# Geometry and Complexity Theory 

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## Preface

This book describes recent applications of algebraic geometry and representation theory to complexity theory. I focus on two central problems: the complexity of matrix multiplication and Valiant's algebraic variants of $\mathbf{P} v$. NP.

I have attempted to make this book accessible to both computer scientists and geometers, and the exposition as self-contained as possible. Two goals are to convince computer scientists of the utility of techniques from algebraic geometry and representation theory, and to show geometers beautiful, interesting, and important geometry questions arising in complexity theory.

Computer scientists have made extensive use combinatorics, graph theory, probability, and linear algebra. I hope to show that even elementary techniques from algebraic geometry and representation theory can substantially advance the search for lower, and even upper bounds in complexity theory. I believe such additional mathematics will be necessary for further advances on questions discussed in this book as well as related complexity problems. Techniques are introduced as needed to deal with concrete problems.

For geometers, I expect that complexity theory will be as good a source for questions in algebraic geometry as has been modern physics. Recent work has indicated that subjects such as Fulton-McPherson intersection theory, the Hilbert scheme of points, and the Kempf-Weyman method for computing syzygies all have something to add to complexity theory. In addition, complexity theory has a way of rejuvenating old questions that had been nearly forgotten but remain beautiful and intriguing: questions of Hadamard, Darboux, Lüroth, and the classical Italian school. At the same time, complexity
theory has brought different areas of mathematics together in new ways: for instance combinatorics, representation theory and algebraic geometry all play a role in understanding the coordinate ring of the orbit closure of the determinant.

This book evolved from several classes I have given on the subject: a spring 2013 semester course at Texas A\&M, summer courses at: Scuola Matematica Inter-universitaria, Cortona (July 2012), CIRM, Trento (June 2014), U. Chicago (IMA sponsored) (July 2014), KAIST, Deajeon (August 2015), and Obergurgul, Austria (September 2016), a fall 2016 semester course at Texas A\&M, and most importantly, a fall 2014 semester course at UC Berkeley as part of the semester-long program, Algorithms and Complexity in Algebraic Geometry, at the Simons Institute for the Theory of Computing.

Since I began writing this book, even since the first draft was completed in fall 2014, the research landscape has shifted considerably: the two paths towards Valiant's conjecture that had been considered the most viable have been shown to be unworkable, at least as originally proposed. On the other hand, there have been significant advances in our understanding of the matrix multiplication tensor. The contents of this book are the state of the art as of January 2017.

Prerequisites. Chapters 1-8 only require a solid background in linear algebra and a willingness to accept several basic results from algebraic geometry that are stated as needed. Nothing beyond [Sha13a] is used in these chapters. Because of the text [Lan12], I am sometimes terse regarding basic properties of tensors and multi-linear algebra. Chapters 9 and 10 contain several sections requiring further background.

Layout. All theorems, propositions, remarks, examples, etc., are numbered together within each section; for example, Theorem 1.3.2 is the second numbered item in Section 1.3. Equations are numbered sequentially within each Chapter. I have included hints for selected exercises, those marked with the symbol © at the end, which is meant to be suggestive of a life preserver. Exercises are marked with (1),(2), or (3), indicating the level of difficulty. Important exercises are also marked with an exclamation mark, sometimes even two, e.g., (1!!) is an exercise that is easy and very important.

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## Introduction

A dramatic leap for signal processing occurred in the 1960s with the implementation of the fast Fourier transform, an algorithm that surprised the engineering community with its efficiency. ${ }^{1}$ Is there a way to predict the existence of such fast unexpected algorithms? Can we prove when they do not exist? Complexity theory addresses these questions.

This book is concerned with the use of geometry towards these goals. I focus primarily on two central questions: the complexity of matrix multiplication and algebraic variants of the famous $\mathbf{P}$ versus $\mathbf{N P}$ problem. In the first case, a surprising algorithm exists and it is conjectured that even better algorithms exist. In the second case, it is conjectured that no surprising algorithm exists.

In this chapter I introduce the main open questions discussed in this book, establish notation that will be used throughout the book, and introduce fundamental geometric notions.

### 1.1. Matrix multiplication

Much of scientific computation amounts to linear algebra, and the basic operation of linear algebra is matrix multiplication. All operations of linear algebra- solving systems of linear equations, computing determinants, etc.use matrix multiplication.

[^0]1.1.1. The standard algorithm. The standard algorithm for multiplying matrices is row-column multiplication: Let $A, B$ be $2 \times 2$ matrices
\[

A=\left($$
\begin{array}{ll}
a_{1}^{1} & a_{2}^{1} \\
a_{1}^{2} & a_{2}^{2}
\end{array}
$$\right), \quad B=\left($$
\begin{array}{ll}
b_{1}^{1} & b_{2}^{1} \\
b_{1}^{2} & b_{2}^{2}
\end{array}
$$\right) .
\]

Remark 1.1.1.1. While computer scientists generally keep all indices down (to distinguish from powers), I use the convention from differential geometry that in a matrix $X$, the entry in the $i$-th row and $j$-th column is labeled $x_{j}^{i}$.

The usual algorithm to calculate the matrix product $C=A B$ is

$$
\begin{aligned}
& c_{1}^{1}=a_{1}^{1} b_{1}^{1}+a_{2}^{1} b_{1}^{2}, \\
& c_{2}^{1}=a_{1}^{1} b_{2}^{1}+a_{2}^{1} b_{2}^{2}, \\
& c_{1}^{2}=a_{1}^{2} b_{1}^{1}+a_{2}^{2} b_{1}^{2}, \\
& c_{2}^{2}=a_{1}^{2} b_{2}^{1}+a_{2}^{2} b_{2}^{2} .
\end{aligned}
$$

It requires 8 multiplications and 4 additions to execute, and applied to $\mathbf{n} \times \mathbf{n}$ matrices, it uses $\mathbf{n}^{3}$ multiplications and $\mathbf{n}^{3}-\mathbf{n}^{2}$ additions.

This algorithm has been around for about two centuries.
In 1968, V. Strassen set out to prove the standard algorithm was optimal in the sense that no algorithm using fewer multiplications exists (personal communication). Since that might be difficult to prove, he set out to show it was true at least for $2 \times 2$ matrices - at least over $\mathbb{Z}_{2}$. His spectacular failure opened up a whole new area of research:
1.1.2. Strassen's algorithm for multiplying $2 \times 2$ matrices using 7 scalar multiplications [Str69]. Set

$$
\begin{align*}
I & =\left(a_{1}^{1}+a_{2}^{2}\right)\left(b_{1}^{1}+b_{2}^{2}\right),  \tag{1.1.1}\\
I I & =\left(a_{1}^{2}+a_{2}^{2}\right) b_{1}^{1}, \\
I I I & =a_{1}^{1}\left(b_{2}^{1}-b_{2}^{2}\right) \\
I V & =a_{2}^{2}\left(-b_{1}^{1}+b_{1}^{2}\right) \\
V & =\left(a_{1}^{1}+a_{2}^{1}\right) b_{2}^{2} \\
V I & =\left(-a_{1}^{1}+a_{1}^{2}\right)\left(b_{1}^{1}+b_{2}^{1}\right), \\
V I I & =\left(a_{2}^{1}-a_{2}^{2}\right)\left(b_{1}^{2}+b_{2}^{2}\right),
\end{align*}
$$

Exercise 1.1.2.1: (1) Show that if $C=A B$, then

$$
\begin{aligned}
& c_{1}^{1}=I+I V-V+V I I, \\
& c_{1}^{2}=I I+I V, \\
& c_{2}^{1}=I I I+V, \\
& c_{2}^{2}=I+I I I-I I+V I .
\end{aligned}
$$

This raises questions:
(1) Can one find an algorithm that uses just six multiplications?
(2) Could Strassen's algorithm have been predicted in advance?
(3) Since it uses more additions, is it actually better in practice?
(4) This algorithm was found by accident and looks ad-hoc. Is there any way to make sense of it? E.g., is there any way to see that it is correct other than a brute force calculation?
(5) What about algorithms for $\mathbf{n} \times \mathbf{n}$ matrices?

I address question (4) in §1.1.15, and the others below, with the last question first:
1.1.3. Fast multiplication of $\mathbf{n} \times \mathbf{n}$ matrices. In Strassen's algorithm, the entries of the matrices need not be scalars - they could themselves be matrices. Let $A, B$ be $4 \times 4$ matrices, and write

$$
A=\left(\begin{array}{ll}
a_{1}^{1} & a_{2}^{1} \\
a_{1}^{2} & a_{2}^{2}
\end{array}\right), \quad B=\left(\begin{array}{ll}
b_{1}^{1} & b_{2}^{1} \\
b_{1}^{2} & b_{2}^{2}
\end{array}\right) .
$$

where $a_{j}^{i}, b_{j}^{i}$ are $2 \times 2$ matrices. One may apply Strassen's algorithm to get the blocks of $C=A B$ in terms of the blocks of $A, B$ performing 7 multiplications of $2 \times 2$ matrices. Since one can apply Strassen's algorithm to each block, one can multiply $4 \times 4$ matrices using $7^{2}=49$ multiplications instead of the usual $4^{3}=64$.

If $A, B$ are $2^{k} \times 2^{k}$ matrices, one may multiply them using $7^{k}$ multiplications instead of the usual $8^{k}$. If $\mathbf{n}$ is not a power of two, enlarge the matrices with blocks of zeros to obtain matrices whose size is a power of two. Asymptotically, by recursion and block multiplication one can multiply $\mathbf{n} \times \mathbf{n}$ matrices using approximately $\mathbf{n}^{\log _{2}(7)} \simeq \mathbf{n}^{2.81}$ multiplications. To see this, let $\mathbf{n}=2^{k}$ and write $7^{k}=\left(2^{k}\right)^{a}$, so $a=\log _{2} 7$.
1.1.4. Regarding the number of additions. The number of additions in Strassen's algorithm also grows like $\mathbf{n}^{2.81}$, so this algorithm is more efficient in practice when the matrices are large. For any efficient algorithm for matrix multiplication, the total complexity is governed by the number of multiplications; see [BCS97, Prop. 15.1]. This is fortuitous because there
is a geometric object, tensor rank, discussed in $\S 1.1 .11$ below, that counts the number of multiplications in an optimal algorithm (within a factor of two), and thus provides a geometric measure of the complexity of matrix multiplication.

Just how large a matrix must be in order to obtain a substantial savings with Strassen's algorithm (size about two thousand suffice) and other practical matters are addressed in $[\mathbf{B B}]$.
1.1.5. An even better algorithm? Regarding question (1) above, one cannot improve upon Strassen's algorithm for $2 \times 2$ matrices. This was first shown in [Win71]. I will give a proof, using geometry and representation theory, of a stronger statement in $\S 8.3 .2$. However for $\mathbf{n}>2$, very little is known, as discussed below and in Chapters 2-5. What is known is that better algorithms than Strassen's exist for $\mathbf{n} \times \mathbf{n}$ matrices when $\mathbf{n}$ is large.
1.1.6. How to predict in advance? The answer to question (2) is yes! In fact it could have been predicted 100 years ago.

Had someone asked Terracini in 1913, he would have been able to predict the existence of something like Strassen's algorithm from geometric considerations alone. Matrix multiplication is a bilinear map (see §1.1.9). Terracini would have been able to tell you, thanks to a simple parameter count (see §2.1.6), that even a general bilinear map $\mathbb{C}^{4} \times \mathbb{C}^{4} \rightarrow \mathbb{C}^{4}$ can be executed using seven multiplications, and thus, fixing any $\epsilon>0$, one can perform any bilinear map $\mathbb{C}^{4} \times \mathbb{C}^{4} \rightarrow \mathbb{C}^{4}$ within an error of $\epsilon$ using seven multiplications.
1.1.7. Big/Little $\mathbf{O}$ etc. notation. For functions $f, g$ of a real variable (or integer) $x$ :
$f(x)=O(g(x))$ if there exists a constant $C>0$ and $x_{0}$ such that $|f(x)| \leq C|g(x)|$ for all $x \geq x_{0}$,
$f(x)=o(g(x))$ if $\lim _{x \rightarrow \infty} \frac{|f(x)|}{|g(x)|}=0$,
$f(x)=\Omega(g(x))$ if there exists a constant $C>0$ and $x_{0}$ such that $C|f(x)| \geq|g(x)|$ for all $x \geq x_{0}$,
$f(x)=\omega(g(x))$ if if $\lim _{x \rightarrow \infty} \frac{|g(x)|}{|f(x)|}=0$, and
$f(x)=\Theta(g(x))$ if $f(x)=O(g(x))$ and $f(x)=\Omega(g(x))$.
1.1.8. The exponent of matrix multiplication. The following quantity is the standard measure of the complexity of matrix multiplication:

Definition 1.1.8.1. The exponent $\omega$ of matrix multiplication is

$$
\begin{gathered}
\omega:=\inf \{h \in \mathbb{R} \mid \mathbf{n} \times \mathbf{n} \text { matrices can be multiplied using } \\
\\
\left.O\left(\mathbf{n}^{h}\right) \text { arithmetic operations }\right\}
\end{gathered}
$$

where inf denotes the infimum.
By Theorem 1.1.11.3 below, Strassen's algorithm shows $\omega \leq \log _{2}(7)<$ 2.81, and it is easy to prove $\omega \geq 2$. Determining $\omega$ is a central open problem in complexity theory. After Strassen's work it was shown in 1978 that $\omega \leq$ 2.7962 [Pan78], then $\omega \leq 2.7799$ [Bin80] in 1979, then $\omega \leq 2.55$ [ $\mathbf{S c h} 81$ ] in 1981, then $\omega \leq 2.48$ [Str87] in 1987, and then $\omega \leq 2.38$ [CW90] in 1989, which might have led people in 1990 to think a resolution was near. However, then nothing happened for over twenty years, and the current "world record" of $\omega<2.373$ [Wil, Gal, Sto] is not much of an improvement since 1990. These results are the topic of Chapter 3.

If one is interested in multiplying matrices of reasonable size, only the algorithm in [Pan78] is known to beat Strassen's. This "practical" exponent is discussed in Chapter 4.

The above work has led to the following astounding conjecture:
Conjecture 1.1.8.2. $\omega=2$.
That is, it is conjectured that asymptotically, it is nearly just as easy to multiply matrices as it is to add them!

Although I am unaware of anyone taking responsibility for the conjecture, most computer scientists I have discussed it with expect it to be true. (For example, multiplying $\mathbf{n}$-bit integers is possible in near linear time $O(\mathbf{n} \log (\mathbf{n}))$, which is almost as efficient as adding them.)

I have no opinion on whether the conjecture should be true or false and thus discuss both upper and lower bounds for the complexity of matrix multiplication, focusing on the role of geometry. Chapters 2 and 5 are dedicated to lower bounds and Chapters 3 and 4 to upper bounds.
1.1.9. Matrix multiplication as a bilinear map. I will use the notation

$$
M_{\langle\mathbf{n}, \mathbf{m}, \mathbf{l}\rangle}: \mathbb{C}^{\mathbf{n} \times \mathbf{m}} \times \mathbb{C}^{\mathbf{m} \times \mathbf{l}} \rightarrow \mathbb{C}^{\mathbf{n} \times \mathbf{1}}
$$

for matrix multiplication of an $\mathbf{n} \times \mathbf{m}$ matrix with an $\mathbf{m} \times \mathbf{l}$ matrix, and write $M_{\langle\mathbf{n}\rangle}=M_{\langle\mathbf{n}, \mathbf{n}, \mathbf{n}\rangle}$.

Matrix multiplication is a bilinear map, that is, for all $X_{j}, X \in \mathbb{C}^{\mathbf{n} \times \mathbf{m}}$, $Y_{j}, Y \in \mathbb{C}^{\mathbf{m} \times 1}$, and $a_{j}, b_{j} \in \mathbb{C}$,
$M_{\langle\mathbf{n}, \mathbf{m}, \mathbf{1}\rangle}\left(a_{1} X_{1}+a_{2} X_{2}, Y\right)=a_{1} M_{\langle\mathbf{n}, \mathbf{m}, \mathbf{1}\rangle}\left(X_{1}, Y\right)+a_{2} M_{\langle\mathbf{n}, \mathbf{m}, \mathbf{l}\rangle}\left(X_{2}, Y\right)$, and
$M_{\langle\mathbf{n}, \mathbf{m}, \mathbf{l}\rangle}\left(X, b_{1} Y_{1}+b_{2} Y_{2}\right)=b_{1} M_{\langle\mathbf{n}, \mathbf{m}, \mathbf{l}\rangle}\left(X, Y_{1}\right)+b_{2} M_{\langle\mathbf{n}, \mathbf{m}, \mathbf{l}\rangle}\left(X, Y_{2}\right)$.

The set of all bilinear maps $\mathbb{C}^{\mathbf{a}} \times \mathbb{C}^{\mathbf{b}} \rightarrow \mathbb{C}^{\mathbf{c}}$ is a vector space. (In our case $\mathbf{a}=\mathbf{n m}, \mathbf{b}=\mathbf{m l}$, and $\mathbf{c}=\mathbf{l n}$.) Write $a_{1}, \ldots, a_{\mathbf{a}}$ for a basis of $\mathbb{C}^{\mathbf{a}}$ and similarly for $\mathbb{C}^{\mathbf{b}}, \mathbb{C}^{\mathbf{c}}$. Then $T: \mathbb{C}^{\mathbf{a}} \times \mathbb{C}^{\mathbf{b}} \rightarrow \mathbb{C}^{\mathbf{c}}$ is uniquely determined by its action on basis vectors:

$$
\begin{equation*}
T\left(a_{i}, b_{j}\right)=\sum_{k=1}^{\mathbf{c}} t^{i j k} c_{k} \tag{1.1.2}
\end{equation*}
$$

That is, the vector space of bilinear maps $\mathbb{C}^{\mathbf{a}} \times \mathbb{C}^{\mathbf{b}} \rightarrow \mathbb{C}^{\mathbf{c}}$, which I will denote by $\mathbb{C}^{\mathbf{a} *} \otimes \mathbb{C}^{\mathbf{b} *} \otimes \mathbb{C}^{\mathbf{c}}$, has dimension abc. (The notation $\mathbb{C}^{\mathbf{a} *} \otimes \mathbb{C}^{\mathbf{b} *} \otimes \mathbb{C}^{\mathbf{c}}$ is motivated in §2.1.) If we represent a bilinear map by a three-dimensional matrix, it may be thought of as taking two column vectors and returning a third column vector.
1.1.10. Ranks of linear maps. I use the notation $\mathbb{C}^{\mathbf{a}}$ for the column vectors of height a and $\mathbb{C}^{\mathbf{a} *}$ for the row vectors.

Definition 1.1.10.1. A linear map $f: \mathbb{C}^{\mathbf{a}} \rightarrow \mathbb{C}^{\mathbf{b}}$ has rank one if there exist $\alpha \in \mathbb{C}^{\mathbf{a} *}$ and $w \in \mathbb{C}^{\mathbf{b}}$ such that $f(v)=\alpha(v) w$. (In other words, every rank one matrix is the product of a row vector with a column vector.) In this case I write $f=\alpha \otimes w$. The rank of a linear map $h: \mathbb{C}^{\mathbf{a}} \rightarrow \mathbb{C}^{\mathbf{b}}$ is the smallest $r$ such that $h$ may be expressed as a sum of $r$ rank one linear maps.

Given an $\mathbf{a} \times \mathbf{b}$ matrix $X$, one can always change bases, i.e., multiply $X$ on the left by an invertible $\mathbf{a} \times \mathbf{a}$ matrix and on the right by an invertible $\mathbf{b} \times \mathbf{b}$ matrix to obtain a matrix with some number of 1's along the diagonal and zeros elsewhere. The number of 1's appearing is called the rank of the matrix and it is the rank of the linear map $X$ determines. In other words, the only property of a linear map $\mathbb{C}^{\mathbf{a}} \rightarrow \mathbb{C}^{\mathbf{b}}$ that is invariant under changes of bases is its rank, and for each rank we have a normal form. This is not surprising because the dimension of the space of such linear maps is $\mathbf{a b}$, we have $\mathbf{a}^{2}$ parameters of changes of bases in $\mathbb{C}^{\mathbf{a}}$ that we can make in a matrix representing the map, and $\mathbf{a}^{2}+\mathbf{b}^{2}>\mathbf{a b}$.
1.1.11. Tensor rank. For bilinear maps $\mathbb{C}^{\mathbf{a}} \times \mathbb{C}^{\mathbf{b}} \rightarrow \mathbb{C}^{\mathbf{c}}$ we are not so lucky as with linear maps, as usually $\mathbf{a b c}>\mathbf{a}^{2}+\mathbf{b}^{2}+\mathbf{c}^{2}$, i.e., there are fewer free parameters of changes of bases than the number of parameters needed to describe the map. This indicates why the study of bilinear maps is vastly more complicated than the study of linear maps.

Nonetheless, there are properties of a bilinear map that will not change under a change of basis. The most important properties for complexity are tensor rank and tensor border rank. Tensor border rank is defined in §1.1.12 below. Tensor rank is a generalization of the rank of a linear map. Tensor rank is defined properly in $\S 2.1 .3$. Informally, a bilinear map $T$ has tensor
rank one if it can be computed with one multiplication. More precisely, $T$ has tensor rank one if in some coordinate system the multi-dimensional matrix representing it has exactly one nonzero entry. This may be expressed without coordinates:

Definition 1.1.11.1. $T \in \mathbb{C}^{\mathbf{a} *} \otimes \mathbb{C}^{\mathbf{b} *} \otimes \mathbb{C}^{\mathbf{c}}$ has tensor rank one if there exist row vectors $\alpha \in \mathbb{C}^{\mathbf{a} *}, \beta \in \mathbb{C}^{\mathbf{b} *}$ and a column vector $w \in \mathbb{C}^{\mathbf{c}}$ such that $T(u, v)=\alpha(u) \beta(v) w . T$ has tensor rank $r$ if it can be written as the sum of $r$ rank one tensors but no fewer, in which case we write $\mathbf{R}(T)=r$. Let $\hat{\sigma}_{r}^{0}=\hat{\sigma}_{r, \mathbf{a}, \mathbf{b}, \mathbf{c}}^{0}$ denote the set of bilinear maps in $\mathbb{C}^{\mathbf{a} *} \otimes \mathbb{C}^{\mathbf{b} *} \otimes \mathbb{C}^{\mathbf{c}}$ of tensor rank at most $r$.

Remark 1.1.11.2. The peculiar notation $\hat{\sigma}_{r}^{0}$ will be explained in §4.7.1. For now, to give an idea where it comes from: $\sigma_{r}=\sigma_{r}\left(\operatorname{Seg}\left(\mathbb{P}^{\mathbf{a}-1} \times \mathbb{P}^{\mathbf{b}-1} \times \mathbb{P}^{\mathbf{c}-1}\right)\right)$ is standard notation in algebraic geometry for the $r$-th secant variety of the Segre variety, which is the object we will study. The hatted object $\hat{\sigma}_{r}$ denotes its cone in affine space and the 0 indicates the subset of this set consisting of tensors of rank at most $r$.

The following theorem shows that tensor rank is a legitimate measure of complexity:
Theorem 1.1.11.3. (Strassen [Str69], also see [BCS97, §15.1])

$$
\omega=\inf \left\{\tau \in \mathbb{R} \mid \mathbf{R}\left(M_{\langle\mathbf{n}\rangle}\right)=O\left(\mathbf{n}^{\tau}\right)\right\} .
$$

That is, $\mathbf{n} \times \mathbf{n}$ matrices may be multiplied using $O\left(\mathbf{n}^{\omega+\epsilon}\right)$ arithmetic operations if and only if the tensor rank of $M_{\langle\mathbf{n}\rangle}$ is $O\left(\mathbf{n}^{\omega+\epsilon}\right)$.

Our goal is thus to determine, for a given $r$, whether or not matrix multiplication lies in $\hat{\sigma}_{r}^{0}$.
1.1.12. How to use algebraic geometry to prove lower bounds for the complexity of matrix multiplication? Algebraic geometry deals with the study of zero sets of polynomials. By a polynomial on the space of bilinear maps $\mathbb{C}^{\mathbf{a}^{*}} \otimes \mathbb{C}^{\mathbf{b} *} \otimes \mathbb{C}^{\mathbf{c}}$, I mean a polynomial in the coefficients $t^{i j k}$, i.e., in abc variables. In $\S 1.1 .14$ I describe a plan to use algebraic geometry to prove upper complexity bounds. A plan to use algebraic geometry for lower bounds is:

Plan to show $M_{\langle\mathbf{n}, \mathbf{m}, \mathbf{1}\rangle} \notin \hat{\sigma}_{r}^{0}$ via algebraic geometry.

- Find a polynomial $P$ on the space of bilinear maps $\mathbb{C}^{\mathrm{nm}} \times \mathbb{C}^{\mathrm{ml}} \rightarrow$ $\mathbb{C}^{\text {nl }}$, such that $P(T)=0$ for all $T \in \hat{\sigma}_{r}^{0}$.
- Show that $P\left(M_{\langle\mathbf{n}, \mathbf{m}, \mathbf{l}\rangle}\right) \neq 0$.

Chapters 2 and 5 discuss techniques for finding such polynomials, using algebraic geometry and representation theory, the study of symmetry in linear algebra.
1.1.13. Representation theory. Representation theory is the systematic study of symmetry. We will primarily be concerned with properties of bilinear maps, tensors, polynomials, etc. that are invariant under changes of bases. Representation theory will facilitate the study of these properties. It has been essential for proving lower bounds for the complexity of $M_{\langle\mathbf{n}\rangle}$.

Let $V$ be a complex vector space of dimension $\mathbf{v}$. (I reserve the notation $\mathbb{C}^{\mathbf{v}}$ for the column vectors with their standard basis.) Let $G L(V)$ denote the group of invertible linear maps $V \rightarrow V$, and I write $G L_{\mathbf{v}}$ for $G L\left(\mathbb{C}^{\mathbf{v}}\right)$. If we have fixed a basis of $V$, this is the group of invertible $\mathbf{v} \times \mathbf{v}$ matrices. If $G$ is a group and $\mu: G \rightarrow G L(V)$ is a group homomorphism, we will say $G$ acts on $V$ and that $V$ is a $G$-module. The image of $\mu$ is called a representation of $G$.

For example, the permutation group on $n$ elements $\mathfrak{S}_{n}$ acts on $\mathbb{C}^{n}$ by

$$
\sigma\left(\begin{array}{c}
v_{1} \\
\vdots \\
v_{n}
\end{array}\right)=\left(\begin{array}{c}
v_{\sigma^{-1}(1)} \\
\vdots \\
v_{\sigma^{-1}(n)}
\end{array}\right),
$$

where $\sigma \in \mathfrak{S}_{n}$ is a a permutation. That is, the image of $\mathfrak{S}_{n}$ in $G L_{n}$ is the set of permutation matrices. (The inverse is used so that for a vector $\bar{v}$, $\sigma(\tau \bar{v})=(\sigma \tau) \bar{v}$.)

A group action is irreducible if there does not exist a proper subspace $U \subset V$ such that $\mu(g) u \in U$ for all $u \in U$ and $g \in G$.

The action of $\mathfrak{S}_{n}$ on $\mathbb{C}^{n}$ is not irreducible since the line spanned by $e_{1}+\cdots+e_{n}$ is preserved by $\mathfrak{S}_{n}$. Note that the subspace spanned by $e_{1}-$ $e_{2}, \ldots, e_{1}-e_{n}$ is also preserved by $\mathfrak{S}_{n}$. Both these subspaces are irreducible $\mathfrak{S}_{n}$-modules.

The essential point is: the sets $X$, such as $X=\hat{\sigma}_{r}^{0} \subset \mathbb{C}^{\text {abc }}$, for which we want polynomials that vanish at the points of $X$, are invariant under the action of groups:

Definition 1.1.13.1. A set $X \subset V$ is invariant under a group $G \subset G L(V)$ if for all $x \in X$ and all $g \in G, g(x) \in X$. Let $G_{X} \subset G L(V)$ denote the group preserving $X$, the largest subgroup of $G L(V)$ under which $X$ is invariant.

When one says that an object has symmetry, it means the object is invariant under the action of a group.

In the case at hand, $X=\hat{\sigma}_{r}^{0} \subset V=A \otimes B \otimes C$. Then $\hat{\sigma}_{r}^{0}$ is invariant under the action of the group $G L(A) \times G L(B) \times G L(C)$ in $G L(V)$, i.e., this image lies in $G_{\hat{\sigma}_{r}^{0}}$.

Recall that an ideal $I$ in a ring $R$ is a vector subspace such that for all $P \in I$ and $Q \in R, P Q \in I$.

Definition 1.1.13.2. For a set $X \subset V$, we will say a polynomial $P$ vanishes on $X$, if $P(x)=0$ for all $x \in X$. The set of all polynomials vanishing on $X$ forms an ideal in the space of polynomials on $V$, called the ideal of $X$ and denoted $I(X)$.

If a polynomial $P$ is in the ideal of $X$, then the polynomial $g(P)$ will also be in the ideal of $X$ for all $g \in G_{X}$. That is:

The ideal of polynomials vanishing on $X$ is a $G_{X}$-module.
The systematic exploitation of symmetry is used throughout this book: to study the ideals of varieties such as $\hat{\sigma}_{r}$ via their irreducible components in Chapter 2, to find new decompositions of the matrix multiplication tensor in Chapter 4, to find normal forms e.g., to prove the state of the art lower bound for the complexity of matrix multiplication in Chapter 5, and to define the only restricted model where an exponential separation of the permanent from the determinant is known in Chapter 7. Chapter 8 is dedicated to representation theory, and Chapters 9 and 10 approach problems in algebraic geometry using representation theory.
1.1.14. How to use algebraic geometry to prove upper bounds for the complexity of matrix multiplication? Based on the above discussion, one could try:
Plan to show $M_{\langle\mathbf{n}, \mathbf{m}, \mathbf{l}\rangle} \in \hat{\sigma}_{r}^{0}$ with algebraic geometry.

- Find a set of polynomials $\left\{P_{j}\right\}$ on the space of bilinear maps $\mathbb{C}^{\mathrm{nm}} \times$ $\mathbb{C}^{\mathrm{ml}} \rightarrow \mathbb{C}^{\mathrm{nl}}$ such that $T \in \hat{\sigma}_{r}^{0}$ if and only if $P_{j}(T)=0$ for all $j$.
- Show that $P_{j}\left(M_{\langle\mathbf{n}, \mathbf{m}, \mathbf{l}\rangle}\right)=0$ for all $j$.

This plan has a problem: Consider the set $S=\left\{(w, z) \in \mathbb{C}^{2} \mid z=0, w \neq\right.$ $0\}$, whose real picture looks like the $z$-axis with the origin removed:


Any polynomial $P \in I(S)$, i.e., any $P$ that evaluates to zero at all points of $S$, will also be zero at the origin.


Exercise 1.1.14.1: (1!) Prove the above assertion.
Just as in this example, the zero set of the polynomials vanishing on $\hat{\sigma}_{r}^{0}$ is larger than $\hat{\sigma}_{r}^{0}$ when $r>1$ (see $\S 2.1 .5$ ) so one cannot certify membership in $\hat{\sigma}_{r}^{0}$ via polynomials, but rather its Zariski closure which I now define:

Definition 1.1.14.2. The Zariski closure of a set $S \subset V$, denoted $\bar{S}$, is the set of $u \in V$ such that $P(u)=0$ for all $P \in I(S)$. A set $S$ is said to be Zariski closed or an algebraic variety if $S=\bar{S}$, i.e., $S$ is the common zero set of a collection of polynomials.

In the example above, $\bar{S}=\left\{(w, z) \in \mathbb{C}^{2} \mid z=0\right\}$.
When $U=\mathbb{C}^{\mathbf{a} *} \otimes \mathbb{C}^{\mathbf{b} *} \otimes \mathbb{C}^{\mathbf{c}}$, let $\hat{\sigma}_{r}:=\overline{\hat{\sigma}_{r}^{0}}$ denote the Zariski closure of the set of bilinear maps of tensor rank at most $r$.

We will see that for almost all $\mathbf{a}, \mathbf{b}, \mathbf{c}$, and $r, \hat{\sigma}_{r}^{0} \subsetneq \hat{\sigma}_{r}$. The problem with the above plan is that it would only show $M_{\langle\mathbf{n}\rangle} \in \hat{\sigma}_{r}$.

Definition 1.1.14.3. $T \in \mathbb{C}^{\mathbf{a}} \otimes \mathbb{C}^{\mathbf{b}} \otimes \mathbb{C}^{\mathbf{c}}$ has border rank $r$ if $T \in \hat{\sigma}_{r}$ and $T \notin \hat{\sigma}_{r-1}$. In this case we write $\underline{\mathbf{R}}(T)=r$.

For the study of the exponent of matrix multiplication, we have good luck:
Theorem 1.1.14.4 (Bini $[\operatorname{Bin} 80]$, see $\S 3.2$ ).

$$
\omega=\inf \left\{\tau \in \mathbb{R} \mid \underline{\mathbf{R}}\left(M_{\langle\mathbf{n}\rangle}\right)=O\left(n^{\tau}\right)\right\} .
$$

That is, although we may have $\underline{\mathbf{R}}\left(M_{\langle\mathbf{n}\rangle}\right)<\mathbf{R}\left(M_{\langle\mathbf{n}\rangle}\right)$, they are not different enough to effect the exponent. In other words, as far as the exponent is concerned, the plan does not have a problem.

For $\mathbf{n}=2$, we will see that $\underline{\mathbf{R}}\left(M_{\langle 2\rangle}\right)=\mathbf{R}\left(M_{\langle 2\rangle}\right)=7$. It is expected that for $\mathbf{n}>2, \underline{\mathbf{R}}\left(M_{\langle\mathbf{n}\rangle}\right)<\mathbf{R}\left(M_{\langle\mathbf{n}\rangle}\right)$. For $\mathbf{n}=3$, we only know $15 \leq \underline{\mathbf{R}}\left(M_{\langle 3\rangle}\right) \leq$ 20 and $19 \leq \mathbf{R}\left(M_{\langle 3\rangle}\right) \leq 23$. In general, we know $\mathbf{R}\left(M_{\langle\mathbf{n}\rangle}\right) \geq 3 \mathbf{n}^{2}-o(\mathbf{n})$, see $\S 2.6$, and $\underline{\mathbf{R}}\left(M_{\langle\mathbf{n}\rangle}\right) \geq 2 \mathbf{n}^{2}-\left\lceil\log _{2}(\mathbf{n})\right\rceil-1$, see $\S 5.4 .5$.
1.1.15. Symmetry and algorithms. In this subsection I mention three uses of symmetry groups in the study of algorithms.

I first address the question raised in §1.1.2: Can we make sense of Strassen's algorithm (1.1.1)? Just as the set $\hat{\sigma}_{r}$ has a symmetry group, the point $M_{\langle\mathbf{1}, \mathbf{m}, \mathbf{n}\rangle}$ also has a symmetry group that includes $G L_{\mathbf{1}} \times G L_{\mathbf{m}} \times G L_{\mathbf{n}}$. (Do not confuse this with $G L_{\mathbf{l m}} \times G L_{\mathbf{m n}} \times G L_{\mathbf{n l}}$ acting on $\mathbb{C}^{\operatorname{lm}} \otimes \mathbb{C}^{\mathbf{m n}} \otimes \mathbb{C}^{\mathbf{n l}}$ which preserves $\hat{\sigma}_{r}^{0}$.) If we let this group act on Strassen's algorithm for $M_{\langle 2\rangle}$, in general we get a new algorithm that also computes $M_{\langle 2\rangle}$. But perhaps the algorithm itself has symmetry.

It does, and the first step to seeing the symmetry is to put all three vector spaces on an equal footing. A linear map $f: A \rightarrow B$ determines a bilinear form $A \times B^{*} \rightarrow \mathbb{C}$ by $(a, \beta) \mapsto \beta(f(a))$. Similarly, a bilinear map $A \times B \rightarrow C$ determines a trilinear form $A \times B \times C^{*} \rightarrow \mathbb{C}$.
Exercise 1.1.15.1: (2!) Show that $M_{\langle\mathbf{n}\rangle}$, considered as a trilinear form, is $(X, Y, Z) \mapsto \operatorname{trace}(X Y Z) \odot$

Since $\operatorname{trace}(X Y Z)=\operatorname{trace}(Y Z X)$, we see that $G_{M_{\langle\mathbf{n}\rangle}}$ also includes a cyclic $\mathbb{Z}_{3}$-symmetry. In Chapter 4 we will see that Strassen's algorithm is invariant under this $\mathbb{Z}_{3}$-symmetry!

This hints that we might be able to use geometry to help find algorithms. This is the topic of Chapter 4.

For tensors or polynomials with continuous symmetry, their algorithms come in families. So to prove lower bounds, i.e., non-existence of a family of algorithms, one can just prove non-existence of a special member of the family. This idea is used to prove to the state of the art lower bound for matrix multiplication presented in §5.4.5.

### 1.2. Separation of algebraic complexity classes

In 1955, John Nash (see [NR16, Chap. 1]) sent a letter to the NSA regarding cryptography, conjecturing an exponential increase in mean key computation length with respect to the length of the key. In a 1956 letter to von Neumann (see [Sip92, Appendix]) Gödel tried to quantify the apparent difference between intuition and systematic problem solving. Around the same time, researchers in the Soviet Union were trying to determine if "brute force search" was avoidable in solving problems such as the famous traveling salesman problem where there seems to be no fast way to find a solution, but any proposed solution can be easily checked, see [Tra84]. (The problem is to determine if there exists a way to visit, say, twenty cities traveling less than a thousand miles. If I claim to have an algorithm to do so, you just need to look at my plan and check the distances.) These discussions eventually gave rise to the complexity classes $\mathbf{P}$, which models problems admitting a fast algorithm to produce a solution, and NP which models problems admitting a fast algorithm to verify a proposed solution. The famous conjecture $\mathbf{P} \neq \mathbf{N P}$ of Cook, Karp, and Levin is that these two classes are distinct. They also showed that many important problems are complete in $\mathbf{N P}$, and hence that resolving the $\mathbf{P}$ versus $\mathbf{N P}$ question has practical importance for understanding whether these problems can be routinely computed. See [Sip92] for a history of the problem and [NR16, Chap. 1] for an up to date survey.

The transformation of this conjecture to a conjecture in geometry goes via algebra:
1.2.1. From complexity to algebra. The $\mathbf{P}$ v. $\mathbf{N P}$ conjecture is generally believed to be out of reach at the moment, so there have been weaker conjectures proposed that might be more tractable. One such comes from a standard counting problem discussed in §6.1.1. This variant has the advantage that it admits a clean algebraic formulation that I now discuss.
L. Valiant [Val79] conjectured that a sequence of polynomials for which there exists an "easy" recipe to write down its coefficients should not necessarily admit a fast evaluation. He defined algebraic complexity classes that are now called VP and VNP, respectively the sequences of polynomials that are "easy" to evaluate, and the sequences whose coefficients are "easy" to write down (see $\S 6.1 .3$ for their definitions), and conjectured:
Conjecture 1.2.1.1 (Valiant [Val79]). VP $\neq$ VNP.
For the precise relationship between this conjecture and the $\mathbf{P} \neq \mathbf{N P}$ conjecture, see [BCS97, Chap. 21]. Analogously with the original conjecture, many natural polynomials are complete in VNP and hence resolving VP versus VNP is important for understanding the computability of these natural polynomials in practice.

Many problems from graph theory, combinatorics, and statistical physics (partition functions) are in VNP. A good way to think of VNP is as the class of sequences of polynomials that can be written down "explicitly".

Most problems from linear algebra (e.g., inverting a matrix, computing its determinant, multiplying matrices) are in VP.

Valiant also showed that a particular polynomial sequence, the permanent $\left(\operatorname{perm}_{n}\right)$, is complete for the class VNP in the sense that VP $\neq \mathbf{V N P}$ if and only if $\left(\operatorname{perm}_{n}\right) \notin \mathbf{V P}$. As explained in $\S 6.1 .1$, the permanent is natural for computer science. Although it is not immediately clear, the permanent is also natural to geometry, see $\S 6.6 .2$. The formula for the permanent of an $n \times n$ matrix $x=\left(x_{j}^{i}\right)$ is:

$$
\begin{equation*}
\operatorname{perm}_{n}(x):=\sum_{\sigma \in \mathfrak{S}_{n}} x_{\sigma(1)}^{1} \cdots x_{\sigma(n)}^{n} . \tag{1.2.1}
\end{equation*}
$$

Here $\mathfrak{S}_{n}$ denotes the group of permutations of $\{1, \ldots, n\}$.
How would one show there is no fast algorithm for the permanent? First we need a precise class of algorithms to consider. To this end, in $\S 6.1 .3$ I define algebraic circuits, which is the standard class of algorithms for computing a polynomial studied in algebraic complexity theory, and their size, which
is a measure of the complexity of the algorithm. Let circuit-size $\left(\right.$ perm $\left._{n}\right)$ denote the size of the smallest algebraic circuit computing perm $n_{n}$. Valiant's hypothesis 1.2.1.1 may be rephrased as:
Conjecture 1.2.1.2 (Valiant [Val79]). circuit-size( perm $_{n}$ ) grows faster than any polynomial in $n$.
1.2.2. From algebra to algebraic geometry. As with our earlier discussion, to prove lower complexity bounds for the permanent, one could work as follows:

Let $S^{n} \mathbb{C}^{N}$ denote the vector space of all homogeneous polynomials of degree $n$ in $N$ variables, so perm $n$ is a point of the vector space $S^{n} \mathbb{C}^{n^{2}}$. If we write an element of $S^{n} \mathbb{C}^{N}$ as $p\left(y_{1}, \ldots, y_{N}\right)=\sum_{1 \leq i_{1} \leq \cdots \leq i_{n} \leq N} c^{i_{1}, \ldots, i_{n}} y_{i_{1}} \cdots y_{i_{n}}$, then we may view the coefficients $c^{i_{1}, \ldots, i_{n}}$ as coordinates on the vector space $S^{n} \mathbb{C}^{N}$. We will look for polynomials on our space of polynomials, that is, polynomials in the coefficients $c^{i_{1}, \ldots, i_{n}}$.
Plan to show $\left(\right.$ perm $\left._{n}\right) \notin \mathbf{V P}$, or at least bound its circuit size by $r$, with algebraic geometry.

- Find a polynomial $P$ on the space $S^{n} \mathbb{C}^{n^{2}}$ such that $P(p)=0$ for all $p \in S^{n} \mathbb{C}^{n^{2}}$ with circuit-size $(p) \leq r$.
- Show that $P\left(\operatorname{perm}_{n}\right) \neq 0$.

By the discussion above on Zariski closure, this may be a more difficult problem than Valiant's original hypothesis: we are not just trying to exclude perm $_{n}$ from having a circuit, but we are also requiring it not be "near" to having a small circuit. I return to this issue in $\S 1.2 .5$ below.
1.2.3. Benchmarks and restricted models. Valiant's hypothesis is expected to be extremely difficult, so it is reasonable to work towards partial results. Two types of partial results are as follows: First, one could attempt to prove the conjecture under additional hypotheses. In the complexity literature, a conjecture with supplementary hypotheses is called a restricted model. For an example of a restricted model, one could restrict to circuits that are formulas (the underlying graph is a formula, see Remark 6.1.5.2). The definition of a formula coincides with our usual notion of a formula. Restricted models are discussed in Chapter 7. Second, one can fix a complexity measure, e.g., circuit-size $\left(\operatorname{perm}_{n}\right)$, and prove lower bounds for it. I will refer to such progress as improving benchmarks.
1.2.4. Another path to algebraic geometry. The permanent resembles one of the most, perhaps the most, studied polynomial, the determinant of
an $n \times n$ matrix $x=\left(x_{j}^{i}\right)$ :

$$
\begin{equation*}
\operatorname{det}_{n}(x):=\sum_{\sigma \in \mathfrak{S}_{n}} \operatorname{sgn}(\sigma) x_{\sigma(1)}^{1} \cdots x_{\sigma(n)}^{n} \tag{1.2.2}
\end{equation*}
$$

Here $\operatorname{sgn}(\sigma)$ denotes the sign of the permutation $\sigma$. The determinant, despite its enormous formula of $n$ ! terms, can be computed very quickly, e.g., by Gaussian elimination. (See $\S 6.1 .3$ for a division-free algorithm.) In particular $\left(\operatorname{det}_{n}\right) \in \mathbf{V P}$. It is not known if $\operatorname{det}_{n}$ is complete for $\mathbf{V P}$, that is, whether or not a sequence of polynomials is in VP if and only if it can be reduced to the determinant in the sense made precise below.

Although

$$
\operatorname{perm}_{2}\left(\begin{array}{ll}
a & b \\
c & d
\end{array}\right)=\operatorname{det}_{2}\left(\begin{array}{cc}
a & -b \\
c & d
\end{array}\right)
$$

Marcus and Minc [MM61], building on work of Pólya and Szegö (see [Gat87]), proved that one could not express $\operatorname{perm}_{m}(y)$ as a size $m$ determinant of a matrix whose entries are affine linear functions of the $x_{j}^{i}$ when $m>2$. This raised the question that perhaps the permanent of an $m \times m$ matrix could be expressed as a slightly larger determinant, which would imply VP $=\mathbf{V N P}$. More precisely, we say $p\left(y^{1}, \ldots, y^{M}\right)$ is an affine linear projection of $q\left(x^{1}, \ldots, x^{N}\right)$, if there exist affine linear functions $x^{\alpha}(y)=x^{\alpha}\left(y^{1}, \ldots, y^{M}\right)$ such that $p(y)=q(x(y))$. For example,

$$
\operatorname{perm}_{3}(y)=\operatorname{det}_{7}\left(\begin{array}{ccccccc}
0 & 0 & 0 & 0 & y_{3}^{3} & y_{2}^{3} & y_{1}^{3}  \tag{1.2.3}\\
y_{1}^{1} & 1 & & & & & \\
y_{2}^{1} & & 1 & & & & \\
y_{3}^{1} & & & 1 & & & \\
& y_{2}^{2} & y_{1}^{2} & 0 & 1 & & \\
& y_{3}^{2} & 0 & y_{1}^{2} & & 1 & \\
& 0 & y_{3}^{2} & y_{2}^{2} & & & 1
\end{array}\right)
$$

This formula is due to B. Grenet [Gre11], who also generalized it to express $\operatorname{perm}_{m}$ as a determinant of size $2^{m}-1$, see $\S 6.6 .3$.

Valiant conjectured that one cannot do much better than this:
Definition 1.2.4.1. Let $p$ be a polynomial. Define the determinantal complexity of $p$, denoted $\operatorname{dc}(p)$, to be the smallest $n$ such that $p$ is an affine linear projection of the determinant.

Valiant showed that for any polynomial $p, \operatorname{dc}(p)$ is finite but possibly larger than circuit-size $(p)$, so the following conjecture is possibly weaker than Conjecture 1.2.1.2.
Conjecture 1.2.4.2 (Valiant [Val79]). $\mathrm{dc}\left(\operatorname{perm}_{m}\right)$ grows faster than any polynomial in $m$.

The state of the art, obtained with classical differential geometry, is $\mathrm{dc}\left(\operatorname{perm}_{m}\right) \geq \frac{m^{2}}{2}$, due to Mignon and Ressayre [MR04]. An exposition of their result is given in $\S 6.4$.
1.2.5. Geometric Complexity Theory. The "Zariski closed" version of Conjecture 1.2.4.2 is the flagship conjecture of Geometric Complexity Theory (GCT) and is discussed in Chapters 6 and 8. To state it in a useful form, first rephrase Valiant's hypothesis as follows:

Let $\operatorname{End}\left(\mathbb{C}^{n^{2}}\right)$ denote the space of all linear maps $\mathbb{C}^{n^{2}} \rightarrow \mathbb{C}^{n^{2}}$, which acts on $S^{n} \mathbb{C}^{n^{2}}$ under the action $L \cdot p(x):=p\left(L^{T}(x)\right)$, where $x$ is viewed as a column vector of size $n^{2}, L$ is an $n^{2} \times n^{2}$ matrix, and $T$ denotes transpose. (The transpose is used so that $L_{1} \cdot\left(L_{2} \cdot p\right)=\left(L_{1} L_{2}\right) \cdot p$.) Let

$$
\operatorname{End}\left(\mathbb{C}^{n^{2}}\right) \cdot p=\left\{L \cdot p \mid L \in \operatorname{End}\left(\mathbb{C}^{n^{2}}\right)\right\}
$$

Define an auxiliary variable $\ell \in \mathbb{C}^{1}$ so $\ell^{n-m} \operatorname{perm}_{m} \in S^{n} \mathbb{C}^{m^{2}+1}$. Consider any linear inclusion $\mathbb{C}^{m^{2}+1} \rightarrow \mathbb{C}^{n^{2}}$ (e.g., with the $M a t_{m \times m}$ in the upper left hand corner and $\ell$ in the $(m+1) \times(m+1)$ slot and zeros elsewhere in the space of $n \times n$ matrices), so we may consider $\ell^{n-m} \operatorname{perm}_{m} \in S^{n} \mathbb{C}^{n^{2}}$. Then

$$
\begin{equation*}
\operatorname{dc}\left(\operatorname{perm}_{m}\right) \leq n \Longleftrightarrow \ell^{n-m} \operatorname{perm}_{m} \in \operatorname{End}\left(\mathbb{C}^{n^{2}}\right) \cdot \operatorname{det}_{n} \tag{1.2.4}
\end{equation*}
$$

This situation begins to resemble our matrix multiplication problem: we have an ambient space $S^{n} \mathbb{C}^{n^{2}}$ (resp. $\left(\mathbb{C}^{n^{2}}\right)^{\otimes 3}$ for matrix multiplication), a subset $\operatorname{End}\left(\mathbb{C}^{n^{2}}\right) \cdot \operatorname{det}_{n}$ (resp. $\hat{\sigma}_{r}^{0}$, the tensors of rank at most $r$ ), and a point $\ell^{n-m} \operatorname{perm}_{m}\left(\operatorname{resp} . M_{\langle\mathbf{n}\rangle}\right)$ and we want to show the point is not in the subset. Note one difference here: the dimension of the ambient space is exponentially large with respect to the dimension of our subset. As before, if we want to separate the point from the subset with polynomials, we are attempting to prove a stronger statement.

Definition 1.2.5.1. For $p \in S^{d} \mathbb{C}^{M}$, let $\overline{\mathrm{dc}}(p)$ denote the smallest $n$ such that $\ell^{n-d} p \in \overline{\operatorname{End}\left(\mathbb{C}^{n^{2}}\right) \cdot \operatorname{det}_{n}}$, the Zariski closure of $\operatorname{End}\left(\mathbb{C}^{n^{2}}\right) \cdot \operatorname{det}_{n}$. Call $\overline{\mathrm{dc}}$ the border determinantal complexity of $p$.

Conjecture 1.2.5.2. $[\mathrm{MS01}] \overline{d c}\left(\right.$ perm $\left._{m}\right)$ grows faster than any polynomial in $m$.

For this problem, we do not have an analog of Bini's theorem 1.1.14.4 that promises similar asymptotics for the two complexity measures. In this situation Mulmuley [Mul] conjectures that there exist sequences of polynomials $\left(p_{m}\right)$ such that $\overline{\mathrm{dc}}\left(p_{m}\right)$ grows like a polynomial in $m$ but dc $\left(p_{m}\right)$ grows faster than any polynomial. Moreover, he speculates that this gap explains why Valiant's hypothesis is so difficult.

Representation theory indicates a path towards solving Conjecture 1.2.5.2. To explain the path, I introduce the following terminology:
Definition 1.2.5.3. A polynomial $p \in S^{n} \mathbb{C}^{N}$ is characterized by its symmetries if, letting $G_{p}:=\left\{g \in G L_{N} \mid g \cdot p=p\right\}$, for any $q \in S^{n} \mathbb{C}^{N}$ with $G_{q} \supseteq G_{p}$, one has $p=\lambda q$ for some $\lambda \in \mathbb{C}$.

There are two essential observations:

- $\overline{\operatorname{End}\left(\mathbb{C}^{n^{2}}\right) \cdot \operatorname{det}_{n}}=\overline{G L_{n^{2}} \cdot \operatorname{det}_{n}}$, that is, the variety $\overline{\operatorname{End}\left(\mathbb{C}^{n^{2}}\right) \cdot \operatorname{det}_{n}}$ is an orbit closure.
- $\operatorname{det}_{n}$ and perm $_{n}$ are characterized by their symmetries.

In principle representation theory (more precisely, the Peter-Weyl Theorem, see §8.6) gives a description of the polynomials vanishing on an orbit closure modulo the effect of the boundary. (It describes the ring of regular functions on the orbit.) Unfortunately for the problem at hand, the approach to Valiant's conjecture via the Peter-Weyl theorem, outlined in [MS01, MS08], was recently shown [IP15, BIP16] to be not viable as proposed. Nevertheless, the approach suggests several alternative paths that could be viable. For this reason, I explain the approach and the proof of its non-viability in Chapter 8.

Unlike matrix multiplication, progress on Valiant's hypothesis and its variants is in its infancy. To gain insight as to what techniques might work, it will be useful to examine "toy" versions of the problem - these questions are of mathematical significance in their own right and lead to interesting connections between combinatorics, representation theory, and geometry. Chapter 9 is dedicated to one such problem, dating back to Hermite and Hadamard, to determine the ideal of the Chow variety of polynomials that decompose into a product of linear forms.

### 1.3. How to find hay in a haystack: the problem of explicitness

A "random" bilinear map $b: \mathbb{C}^{\mathrm{m}} \times \mathbb{C}^{\mathrm{m}} \rightarrow \mathbb{C}^{\mathrm{m}}$ will have tensor rank around $\frac{\mathrm{m}^{2}}{3}$, see $\S 4.7$ for the precise rank. (In particular, the standard algorithm for matrix multiplication already shows that it is pathological as a tensor as $\mathbf{n}^{3} \ll \frac{\left(\mathbf{n}^{2}\right)^{2}}{3}$.) On the other hand, how would one find an explicit tensor of tensor rank around $\frac{\mathrm{m}^{2}}{3}$ ? This is the problem of finding hay in a haystack ${ }^{2}$. Our state of the art for this question is so dismal that there is no known explicit bilinear map of tensor rank $3 \mathbf{m}$, in fact the highest rank of an explicit tensor known (modulo the error term) is for matrix multiplication [Lan14]:

[^1]$\mathbf{R}\left(M_{\langle\mathbf{n}\rangle}\right) \geq 3 \mathbf{n}^{2}-o\left(\mathbf{n}^{2}\right)$. Other explicit sequences $T_{\mathbf{m}}: \mathbb{C}^{\mathbf{m}} \times \mathbb{C}^{\mathbf{m}} \rightarrow \mathbb{C}^{\mathbf{m}}$ with $\mathbf{R}\left(T_{\mathbf{m}}\right) \geq 3 \mathbf{m}-o(\mathbf{m})$ were found in [Zui15] and the largest known rank tensor, from [AFT11], has $\mathbf{R}\left(T_{\mathbf{m}}\right) \geq 3 \mathbf{m}-o(\log (\mathbf{m}))$. It is a frequently stated open problem to find explicit bilinear maps $T_{\mathbf{m}}: \mathbb{C}^{\mathbf{m}} \times \mathbb{C}^{\mathbf{m}} \rightarrow \mathbb{C}^{\mathbf{m}}$ with $\mathbf{R}\left(T_{m}\right) \geq(3+\epsilon) \mathbf{m}$. In Chapter 5, I discuss the state of the art of this problem and the related border rank problem, where no explicit tensor $T \in \mathbb{C}^{\mathbf{m}} \otimes \mathbb{C}^{\mathbf{m}} \otimes \mathbb{C}^{\mathbf{m}}$ with $\underline{\mathbf{R}}(T)>2 \mathbf{m}$ is known. Valiant's hypothesis may also be phrased in this language: exhibiting an explicit polynomial sequence that is provably difficult to compute would be sufficient to prove Valiant's hypothesis (a random sequence is provably difficult).

### 1.4. The role of algebraic geometry

Recent research (e.g., [Gal16, BB14]) has shown that in order to prove super-linear lower bounds on tensor rank or border rank, thanks to the cactus variety, one must deal with subtle questions regarding zero dimensional schemes. The work [GKKS13a] indicates that questions regarding the geometry of syzygies could play a role in the resolution of Valiant's hypothesis. Chapter 10 introduces these topics and others from algebraic geometry and representation theory, and explains their role in complexity theory. It is written as an invitation to algebraic geometers with expertise in these areas to work on questions in complexity theory.

## Chapter 2

## The complexity of Matrix multiplication I: first lower bounds

In this chapter I discuss lower complexity bounds for tensors in general and matrix multiplication in particular. The two basic measures of complexity are rank and border rank. I begin, in $\S 2.1$, by defining tensors and their rank. I motivate the definition of border rank with the story of the discovery by Bini et. al. of approximate algorithms for a reduced matrix multiplication tensor and then give its definition. Next, in $\S 2.2$ I present Strassen's equations. In order to generalize them, I present elementary definitions and results from mutli-linear algebra and representation theory in $\S 2.3$, including the essential Schur's lemma. I then, in $\S 2.4$ give Ottaviani's derivation of Strassen's equations that generalizes to Koszul flattenings, which are also derived. In $\S 2.5$, I show a $2 \mathbf{n}^{2}-\mathbf{n}$ lower bound for the border rank of $M_{\langle\mathbf{n}\rangle}$. This border rank lower bound is exploited to prove a $3 \mathbf{n}^{2}-o\left(\mathbf{n}^{2}\right)$ rank lower bound for $M_{\langle\mathbf{n}\rangle}$ in $\S 2.6$. The current state of the art is a $2 \mathbf{n}^{2}-\left\lceil\log _{2}(\mathbf{n})\right\rceil-1$ lower bound for the border rank of $M_{\langle\mathbf{n}\rangle}$, which is presented in §5.4.5, as it requires more geometry and representation theory than what is covered in this chapter.

### 2.1. Matrix multiplication and multi-linear algebra

To better understand matrix multiplication as a bilinear map, I first review basic facts from multi-linear algebra. For more details on this topic, see [Lan12, Chap. 2].
2.1.1. Linear algebra without coordinates. In what follows it will be essential to work without bases, so instead of writing $\mathbb{C}^{\mathbf{v}}$, I use $V$ to denote a complex vector space of dimension $\mathbf{v}$.

The dual space $V^{*}$ to a vector space $V$ is the vector space whose elements are linear maps from $V$ to $\mathbb{C}$ :

$$
V^{*}:=\{\alpha: V \rightarrow \mathbb{C} \mid \alpha \text { is linear }\} .
$$

This notation is consistent with the notation of $\mathbb{C}^{\mathbf{v}}$ for column vectors and $\mathbb{C}^{\mathbf{v} *}$ for row vectors because if in bases elements of $V$ are represented by column vectors, then elements of $V^{*}$ are naturally represented by row vectors and the map $v \mapsto \alpha(v)$ is just row-column matrix multiplication. Given a basis $v_{1}, \ldots, v_{\mathbf{v}}$ of $V$, it determines a basis $\alpha^{1}, \ldots, \alpha^{\mathbf{v}}$ of $V^{*}$ by $\alpha^{i}\left(v_{j}\right)=\delta_{i j}$, called the dual basis.

Let $V^{*} \otimes W$ denote the vector space of all linear maps $V \rightarrow W$. Given $\alpha \in V^{*}$ and $w \in W$ define a linear map $\alpha \otimes w: V \rightarrow W$ by $\alpha \otimes w(v):=\alpha(v) w$. In bases, if $\alpha$ is represented by a row vector and $w$ by a column vector, $\alpha \otimes w$ will be represented by the matrix $w \alpha$. Such a linear map is said to have rank one. Define the rank of an element $f \in V^{*} \otimes W$ to be the smallest $r$ such $f$ may be expressed as a sum of $r$ rank one linear maps.

Recall from Definition 1.1.14.2, that a variety is the common zero set of a collection of polynomials.

Definition 2.1.1.1. A variety $Z \subset V$ is reducible if it is possible to write $Z=Z_{1} \cup Z_{2}$ with $Z_{1}, Z_{2}$ nonempty varieties. Otherwise it is irreducible.

Definition 2.1.1.2. A property of points of an irreducible variety $Z \subset W$ is general or holds generally if the property holds on the complement of a proper subvariety of $Z$.

A general point of a variety $Z \subset V$ is a point not lying on some explicit Zariski closed subset of $Z$. This subset is often understood from the context and so not mentioned.

The complement to the zero set of any polynomial over the complex numbers has full measure, so properties that hold at general points hold with probability one for a randomly chosen point in $Z$.
Theorem 2.1.1.3 (Fundamental theorem of linear algebra). Let $V, W$ be finite dimensional vector spaces, let $f: V \rightarrow W$ be a linear map, and let $A_{f}$ be a matrix representing $f$. Then
(1)

$$
\begin{aligned}
\operatorname{rank}(f) & =\operatorname{dim} f(V) \\
& =\operatorname{dim}\left(\operatorname{span}\left\{\text { columns of } A_{f}\right\}\right) \\
& =\operatorname{dim}\left(\operatorname{span}\left\{\text { rows of } A_{f}\right\}\right) \\
& =\operatorname{dim} V-\operatorname{dim} \operatorname{ker} f .
\end{aligned}
$$

In particular $\operatorname{rank}(f) \leq \min \{\operatorname{dim} V, \operatorname{dim} W\}$.
(2) For general $f \in V^{*} \otimes W, \operatorname{rank}(f)=\min \{\operatorname{dim} V, \operatorname{dim} W\}$.
(3) If a sequence of linear maps $f_{t}$ of rank $r$ has a limit $f_{0}$, then $\operatorname{rank}\left(f_{0}\right) \leq r$.
(4) $\operatorname{rank}(f) \leq r$ if and only if, in any choice of bases, the determinants of all size $r+1$ submatrices of the matrix representing $f$ are zero.

Note that assertion 4) shows that the set of linear maps of rank at most $r$ forms an algebraic variety. Although we take it for granted, it is really miraculous that the fundamental theorem of linear algebra is true. I explain why in §2.1.5.
Exercise 2.1.1.4: (1!) Prove the theorem. ©
Exercise 2.1.1.5: (1) Assuming $V$ is finite dimensional, write down a canonical isomorphism $V \rightarrow\left(V^{*}\right)^{*}$. ©

Many standard notions from linear algebra have coordinate free definitions. For example: A linear map $f: V \rightarrow W$ determines a linear map $f^{T}: W^{*} \rightarrow V^{*}$ defined by $f^{T}(\beta)(v):=\beta(f(v))$ for all $v \in V$ and $\beta \in W^{*}$. Note that this is consistent with the notation $V^{*} \otimes W \simeq W \otimes V^{*}$, being interpreted as the space of all linear maps $\left(W^{*}\right)^{*} \rightarrow V^{*}$, that is, the order we write the factors does not matter. If we work in bases and insist that all vectors are column vectors, the matrix of $f^{T}$ is just the transpose of the matrix of $f$.
Exercise 2.1.1.6: (1) Show that we may also consider an element $f \in$ $V^{*} \otimes W$ as a bilinear map $b_{f}: V \times W^{*} \rightarrow \mathbb{C}$ defined by $b_{f}(v, \beta):=\beta(f(v))$.
2.1.2. Multi-linear maps and tensors. The space $V \otimes W$ is called the tensor product of $V$ with $W$. More generally, for vector spaces $A_{1}, \ldots, A_{n}$ define their tensor product $A_{1} \otimes \cdots \otimes A_{n}$ to be the space of $n$-linear maps $A_{1}^{*} \times \cdots \times A_{n}^{*} \rightarrow \mathbb{C}$, equivalently the space of $(n-1)$-linear maps $A_{1}^{*} \times \cdots \times$ $A_{n-1}^{*} \rightarrow A_{n}$ etc.. When $A_{1}=\cdots=A_{n}=V$, write $V^{\otimes n}=V \otimes \cdots \otimes V$.

Let $a_{j} \in A_{j}$ and define an element $a_{1} \otimes \cdots \otimes a_{n} \in A_{1} \otimes \cdots \otimes A_{n}$ to be the $n$-linear map

$$
a_{1} \otimes \cdots \otimes a_{n}\left(\alpha^{1}, \ldots, \alpha^{n}\right):=\alpha^{1}\left(a_{1}\right) \cdots \alpha^{n}\left(a_{n}\right) .
$$

Exercise 2.1.2.1: (1) Show that if $\left\{a_{j}^{s_{j}} \mid 1 \leq s_{j} \leq \mathbf{a}_{j}\right\}$, is a basis of $A_{j}$, then $\left\{a_{1}^{s_{1}} \otimes \cdots \otimes a_{n}^{s_{n}} \mid 1 \leq s_{j} \leq \mathbf{a}_{j}\right\}$ is a basis of $A_{1} \otimes \cdots \otimes A_{n}$. In particular $\operatorname{dim}\left(A_{1} \otimes \cdots \otimes A_{n}\right)=\mathbf{a}_{1} \cdots \mathbf{a}_{n}$. ©

Remark 2.1.2.2. One may identify $A_{1} \otimes \cdots \otimes A_{n}$ with any re-ordering of the factors. When I need to be explicit about this, I will call the identification the re-ordering isomorphism.

Example 2.1.2.3 (Matrix multiplication). Let $x_{\alpha}^{i}, y_{u}^{\alpha}, z_{i}^{u}$ respectively be bases of $A=\mathbb{C}^{\mathrm{nm}}, B=\mathbb{C}^{\mathrm{ml}}, C=\mathbb{C}^{\mathrm{ln}}$, then the standard expression of matrix multiplication as a tensor is

$$
\begin{equation*}
M_{\langle\mathbf{1}, \mathbf{m}, \mathbf{n}\rangle}=\sum_{i=1}^{\mathbf{n}} \sum_{\alpha=1}^{\mathbf{m}} \sum_{u=1}^{\mathbf{l}} x_{\alpha}^{i} \otimes y_{u}^{\alpha} \otimes z_{i}^{u} \tag{2.1.1}
\end{equation*}
$$

Exercise 2.1.2.4: (2) Write Strassen's algorithm out as a tensor. ©
2.1.3. Tensor rank. An element $T \in A_{1} \otimes \cdots \otimes A_{n}$ is said to have rank one if there exist $a_{j} \in A_{j}$ such that $T=a_{1} \otimes \cdots \otimes a_{n}$.

I will use the following measure of complexity:
Definition 2.1.3.1. Let $T \in A_{1} \otimes \cdots \otimes A_{n}$. Define the rank (or tensor rank) of $T$ to be the smallest $r$ such that $T$ may be written as the sum of $r$ rank one tensors. Write $\mathbf{R}(T)=r$. Let $\hat{\sigma}_{r}^{0} \subset A_{1} \otimes \cdots \otimes A_{n}$ denote the set of tensors of rank at most $r$.

For bilinear maps, tensor rank is comparable to all other standard measures of complexity on the space of bilinear maps, see, e.g., [BCS97, §14.1].

By (2.1.1) we conclude $\mathbf{R}\left(M_{\langle\mathbf{n}, \mathbf{m}, \mathbf{l}\rangle}\right) \leq \mathbf{n m l}$. Strassen's algorithm shows $\mathbf{R}\left(M_{\langle 2,2,2\rangle}\right) \leq 7$. Shortly afterwards, Winograd [Win71] showed $\mathbf{R}\left(M_{\langle 2,2,2\rangle}\right)=$ 7.
2.1.4. Another spectacular failure. After Strassen's failure to prove the standard algorithm for matrix multiplication was optimal, Bini et. al. [BLR80] considered the reduced matrix multiplication operator

$$
\begin{aligned}
M_{\langle 2\rangle}^{r e d}:= & x_{1}^{1} \otimes\left(y_{1}^{1} \otimes z_{1}^{1}+y_{2}^{1} \otimes z_{1}^{2}\right)+x_{2}^{1} \otimes\left(y_{1}^{2} \otimes z_{1}^{1}+y_{2}^{2} \otimes z_{1}^{2}\right)+x_{1}^{2} \otimes\left(y_{1}^{1} \otimes z_{2}^{1}+y_{2}^{1} \otimes z_{2}^{2}\right) \\
& \in \mathbb{C}^{3} \otimes \mathbb{C}^{4} \otimes \mathbb{C}^{4}
\end{aligned}
$$

obtained by setting the $x_{2}^{2}$ entry for $M_{\langle 2\rangle}$ to zero. The standard presentation shows $\mathbf{R}\left(M_{\langle 2\rangle}^{r e d}\right) \leq 6$. Bini et. al. attempted to find a rank five expression for $M_{\langle 2\rangle}^{r e d}$. They searched for such an expression by computer. Their method was to minimize the norm of $M_{\langle 2\rangle}^{r e d}$ minus a rank five tensor that varied (see $\S 4.6$ for a description of the method), and their computer kept on producing rank five tensors with the norm of the difference getting smaller and smaller,
but with larger and larger coefficients. Bini (personal communication) told me about how he lost sleep trying to understand what was wrong with his computer code. This went on for some time, when finally he realized there was nothing wrong with the code: the output it produced was a manifestation of the phenomenon Bini named border rank [Bin80], which was mentioned in the introduction in the context of finding polynomials for upper rank bounds.

The expression for the tensor $M_{\langle 2\rangle}^{r e d}$ that their computer search found was essentially

$$
\begin{align*}
& M_{\langle 2\rangle}^{r e d}=\lim _{t \rightarrow 0} \frac{1}{t}\left[\left(x_{2}^{1}+t x_{1}^{1}\right) \otimes\left(y_{2}^{1}+t y_{2}^{2}\right) \otimes z_{1}^{2}\right.  \tag{2.1.2}\\
&+\left(x_{1}^{2}+t x_{1}^{1}\right) \otimes y_{1}^{1} \otimes\left(z_{1}^{1}+t z_{2}^{1}\right) \\
&-x_{2}^{1} \otimes y_{2}^{1} \otimes\left(\left(z_{1}^{1}+z_{1}^{2}\right)+t z_{2}^{2}\right) \\
&-x_{1}^{2} \otimes\left(\left(y_{1}^{1}+y_{2}^{1}\right)+t y_{1}^{2}\right) \otimes z_{1}^{1} \\
&\left.+\left(x_{2}^{1}+x_{1}^{2}\right) \otimes\left(y_{2}^{1}+t y_{1}^{2}\right) \otimes\left(z_{1}^{1}+t z_{2}^{2}\right)\right] .
\end{align*}
$$

The rank five tensors found by Bini et. al. were the right hand side of (2.1.2) (without the limit) for particular small values of $t$.

In what follows I first explain why border rank is needed in the study of tensors and then properly define it.
2.1.5. The Fundamental theorem of linear algebra is false for tensors. Recall the fundamental theorem of linear algebra from §2.1.1.3.
Theorem 2.1.5.1. If $T \in \mathbb{C}^{\mathrm{m}} \otimes \mathbb{C}^{\mathrm{m}} \otimes \mathbb{C}^{\mathrm{m}}$ is general, i.e., outside the zero set of a certain finite collection of polynomials (in particular outside a certain set of measure zero), then $\mathbf{R}(T) \geq\left\lceil\frac{\mathrm{m}^{3}-1}{3 \mathrm{~m}-2}\right\rceil$.

Tensor rank can jump up (or down) under limits.
The first assertion is proved in $\S 4.7 .1$. To see the second assertion, at least when $r=2$, consider

$$
T(t):=\frac{1}{t}\left[a_{1} \otimes b_{1} \otimes c_{1}-\left(a_{1}+t a_{2}\right) \otimes\left(b_{1}+t b_{2}\right) \otimes\left(c_{1}+t c_{2}\right)\right]
$$

and note that

$$
\lim _{t \rightarrow 0} T(t)=a_{1} \otimes b_{1} \otimes c_{2}+a_{1} \otimes b_{2} \otimes c_{1}+a_{2} \otimes b_{1} \otimes c_{1}
$$

which has rank three.
Exercise 2.1.5.2: (1) Prove $\mathbf{R}\left(a_{1} \otimes b_{1} \otimes c_{2}+a_{1} \otimes b_{2} \otimes c_{1}+a_{2} \otimes b_{1} \otimes c_{1}\right)=3$. ©
Remark 2.1.5.3. Physicists call the tensor $a_{1} \otimes b_{1} \otimes c_{2}+a_{1} \otimes b_{2} \otimes c_{1}+a_{2} \otimes b_{1} \otimes c_{1}$ the $W$-state so I will sometimes denote it $T_{W \text { State }}$.

To visualize why rank can jump up while taking limits, consider the following picture, where the curve represents the points of $\hat{\sigma}_{1}^{0}$. Points of $\hat{\sigma}_{2}^{0}$ (e.g., the dots limiting to the dot labeled $T$ ) are those on a secant line to $\hat{\sigma}_{1}^{0}$, and the points where the rank jumps up, such at the dot labeled $T$, are those that lie on a tangent line to $\hat{\sigma}_{1}^{0}$. This phenomena fails to occur for matrices because for matrices, every point on a tangent line is also on an honest secant line. Thus in some sense it is a miracle that rank is semi-continuous for matrices.


Our situation regarding tensor rank may be summarized as follows:

- The set $\hat{\sigma}_{r}^{0}$ is not closed under taking limits. I will say a set that is closed under taking limits is Euclidean closed.
- It is also not Zariski closed, i.e., the zero set of all polynomials vanishing on $\hat{\sigma}_{r}^{0}$ includes tensors that are of rank greater than $r$.
Exercise 2.1.5.4: (2) Show that the Euclidean closure (i.e., closure under taking limits) of a set is always contained in its Zariski closure. ©

The tensors that are honestly "close" to tensors of rank $r$ would be the Euclidean closure, but to deal with polynomials as proposed in §1.1.121.1.14, we need to work with the potentially larger Zariski closure.

Often the Zariski closure is much larger than the Euclidean closure. For example, the Zariski closure of $\mathbb{Z} \subset \mathbb{C}$ is $\mathbb{C}$, while $\mathbb{Z}$ is already closed in the Euclidean topology.

For the purposes of proving lower bounds, none of this is an issue, but when we discuss upper bounds, we will need to deal with these problems. For now, I mention that with $\hat{\sigma}_{r}^{0}$ we have good luck: the Zariski and Euclidean closures of $\hat{\sigma}_{r}^{0}$ coincide, so our apparently different informal uses of the term border rank coincide. I present the proof in §3.1.6.

Remark 2.1.5.5. This coincidence is a consequence of a standard result in algebraic geometry that the computer science community was unaware of. As a result, it ended up being re-proven in special cases, e.g., in [Lic84, ?].
2.1.6. Border rank. Generalizing the discussion in §1.1.11, $\hat{\sigma}_{r}=\hat{\sigma}_{r, A_{1} \otimes \ldots \otimes A_{n}}$ denotes the Zariski (and by the above discussion Euclidean) closure of $\hat{\sigma}_{r}^{0}$, and the border rank of $T \in A_{1} \otimes \cdots \otimes A_{n}$, denoted $\underline{\mathbf{R}}(T)$, is the smallest $r$ such that $T \in \hat{\sigma}_{r}$. By the above discussion, border rank is semi-continuous.
Exercise 2.1.6.1: (1) Write down an explicit tensor of border rank $r$ in $\mathbb{C}^{r} \otimes \mathbb{C}^{r} \otimes \mathbb{C}^{r}$ with rank greater than $r$. ©

Border rank is easier to work with than rank for several reasons. For example, the maximal rank of a tensor in $\mathbb{C}^{m} \otimes \mathbb{C}^{m} \otimes \mathbb{C}^{m}$ is not known in general. In contrast, the maximal border rank is known to be $\left\lceil\frac{m^{3}-1}{3 m-2}\right\rceil$ for all $m \neq 3$, and is 5 when $m=3$ [Lic85]. In particular Strassen's algorithm could have been predicted in advance with this knowledge. The method of proof is a differential-geometric calculation that dates back to Terracini in the 1900's [Ter11], see $\S 4.7 .1$ for a discussion.
Exercise 2.1.6.2: (1) Prove that if $T \in A \otimes B \otimes C$ and $T^{\prime}:=\left.T\right|_{A^{\prime} \times B^{\prime} \times C^{\prime}}$ (here $T$ is being considered as a trilinear form) for some $A^{\prime} \subseteq A^{*}, B^{\prime} \subseteq B^{*}$, $C^{\prime} \subseteq C^{*}$, then $\mathbf{R}(T) \geq \mathbf{R}\left(T^{\prime}\right)$ and $\underline{\mathbf{R}}(T) \geq \underline{\mathbf{R}}\left(T^{\prime}\right)$. ©
Exercise 2.1.6.3: (1) Let $T_{j} \in A_{j} \otimes B_{j} \otimes C_{j}, 1 \leq j, k, l \leq s$. Consider $T_{1} \oplus \cdots \oplus T_{s} \in\left(\oplus_{j} A_{j}\right) \otimes\left(\oplus_{k} B_{k}\right) \otimes\left(\oplus_{l} C_{l}\right)$ Show that $\mathbf{R}\left(\oplus_{j} T_{j}\right) \leq \sum_{i=1}^{s} \mathbf{R}\left(T_{i}\right)$ and $\underline{\mathbf{R}}\left(\oplus_{j} T_{j}\right) \leq \sum_{i=1}^{s} \underline{\mathbf{R}}\left(T_{i}\right)$.
Exercise 2.1.6.4: (1) Let $T_{j} \in A_{j} \otimes B_{j} \otimes C_{j}, 1 \leq j, k, l \leq s$. Let $A=\otimes_{j} A_{j}$, $B=\otimes_{k} B_{k}$, and $C=\otimes_{l} C_{l}$, consider $T_{1} \otimes \cdots \otimes T_{s} \in A \otimes B \otimes C$. Show that $\mathbf{R}\left(\otimes_{i=1}^{s} T_{i}\right) \leq \Pi_{i=1}^{s} \mathbf{R}\left(T_{i}\right)$, and $\underline{\mathbf{R}}\left(\otimes_{i=1}^{s} T_{i}\right) \leq \Pi_{i=1}^{s} \underline{\mathbf{R}}\left(T_{i}\right)$
2.1.7. Our first lower bound. Given $T \in A \otimes B \otimes C$, write $T \in A \otimes(B \otimes C)$ and think of $T$ as a linear map $T_{A}: A^{*} \rightarrow B \otimes C$. I will write $T\left(A^{*}\right) \subset B \otimes C$ for the image.
$\operatorname{Proposition~2.1.7.1.~} \underline{\mathbf{R}}(T) \geq \operatorname{rank}\left(T_{A}\right)$.
Exercise 2.1.7.2: (1!) Prove Proposition 2.1.7.1. ©
Say dimensions $\mathbf{a}, \mathbf{b}, \mathbf{c}$ are unbalanced if any of the inequalities $\mathbf{a}>\mathbf{b c}$, $\mathbf{b}>\mathbf{a c}, \mathbf{c}>\mathbf{a b}$ hold, and otherwise that they are balanced.

Permuting the three factors, and assuming the dimensions are balanced, we have equations for $\hat{\sigma}_{r, A \otimes B \otimes C}$ for $r \leq \max \{\mathbf{a}-1, \mathbf{b}-1, \mathbf{c}-1\}$, namely the size $r+1$ minors of the linear maps $T_{A}, T_{B}, T_{C}$.

Definition 2.1.7.3. A tensor $T \in A \otimes B \otimes C$ is concise if the maps $T_{A}, T_{B}$ and $T_{C}$ are all injective.

Exercise 2.1.7.4: (2!) Find a choice of bases such that

$$
M_{\langle\mathbf{n}\rangle_{A}}\left(A^{*}\right)=\left(\begin{array}{lll}
x & & \\
& \ddots & \\
& & x
\end{array}\right)
$$

where $x=\left(x_{j}^{i}\right)$ is $\mathbf{n} \times \mathbf{n}$, i.e., the image in the space of $\mathbf{n}^{2} \times \mathbf{n}^{2}$ matrices is block diagonal with all blocks the same.
Exercise 2.1.7.5: (1) Show that $\underline{\mathbf{R}}\left(M_{\langle\mathbf{n}\rangle}\right) \geq \mathbf{n}^{2}$.
Exercise 2.1.7.6: (1) Show $\underline{\mathbf{R}}\left(M_{\langle\mathbf{m}, \mathbf{n}, 1\rangle}\right)=\mathbf{R}\left(M_{\langle\mathbf{m}, \mathbf{n}, 1\rangle}\right)=\mathbf{m n}$ and $\underline{\mathbf{R}}\left(M_{\langle\mathbf{m}, 1,1\rangle}\right)=$ $\mathbf{R}\left(M_{\langle\mathbf{m}, 1,1\rangle}\right)=\mathbf{m}$.
Exercise 2.1.7.7: (1!) Let $\mathbf{b}=\mathbf{c}$ and assume $T_{A}$ is injective. Show that if $T\left(A^{*}\right)$ is simultaneously diagonalizable under the action of $G L(B) \times G L(C)$ (i.e., there exists $g \in G L(B) \times G L(C)$ such that for any basis $\alpha^{1}, \ldots, \alpha^{\text {a }}$ of $A^{*}$, the elements $g \cdot T\left(\alpha^{1}\right), \ldots, g \cdot T\left(\alpha^{\mathbf{a}}\right)$ are all diagonal) then $\mathbf{R}(T) \leq \mathbf{b}$, and therefore if $T\left(A^{*}\right)$ is the limit of simultaneously diagonalizable subspaces then $\underline{\mathbf{R}}(T) \leq \mathbf{b}$.

### 2.2. Strassen's equations

It wasn't until 1983 [ $\mathbf{S t r 8 3}$ ] that the first non-classical equations were found for tensor border rank. These equations had been found in the related settings of partially symmetric tensors in 1877 by Fram-Toeplitz and 1977 by Barth [Toe77, Bar77], and in the completely symmetric case in 1858 by Aronhold $[$ Aro58]. See $[\mathbf{O t t 0 7}]$ for a history. Here they are:
2.2.1. A test beyond the classical equations. The classical equations just used that $B \otimes C$ is a vector space. To extract more information from $T_{A}$, we examine its image in $B \otimes C$, which we will view as a space of linear maps $C^{*} \rightarrow B$. If dimensions are balanced, $T$ is concise and has minimal border rank $\max \{\mathbf{a}, \mathbf{b}, \mathbf{c}\}$, the image should be special in some way - how? Assume $\mathbf{b}=\mathbf{c}$ so the image is a space of linear maps between two vector spaces of the same dimension. (If $\mathbf{b}<\mathbf{c}$, just restrict to some $\mathbb{C}^{\mathbf{b}} \subset C^{*}$.) If $\mathbf{R}(T)=\mathbf{b}$, then $T\left(A^{*}\right)$ will be spanned by $\mathbf{b}$ rank one linear maps.
Lemma 2.2.1.1. If $\mathbf{a}=\mathbf{b}=\mathbf{c}$ and $T_{A}$ is injective, then $\mathbf{R}(T)=\mathbf{a}$ if and only if $T\left(A^{*}\right)$ is spanned by a rank one linear maps.
Exercise 2.2.1.2: (2!) Prove Lemma 2.2.1.1. ©

How can we test if the image is spanned by $\mathbf{b}$ rank one linear maps? If $T=a_{1} \otimes b_{1} \otimes c_{1}+\cdots+a_{\mathbf{a}} \otimes b_{\mathbf{a}} \otimes c_{\mathbf{a}}$ with each set of vectors a basis, then

$$
T\left(A^{*}\right)=\left\{\left.\left(\begin{array}{llll}
x_{1} & & & \\
& x_{2} & & \\
& & \ddots & \\
& & & x_{\mathbf{a}}
\end{array}\right) \right\rvert\, x_{j} \in \mathbb{C}\right\}
$$

and this is the case for a general rank a tensor in $\mathbb{C}^{\mathbf{a}} \otimes \mathbb{C}^{\mathbf{a}} \otimes \mathbb{C}^{\mathbf{a}}$. That is, the space $T\left(A^{*}\right) \subset B \otimes C$, when $T$ has border rank a, lies in the Zariski closure of the subspaces that, under the action of $G L(B) \times G L(C)$ are simultaneously diagonalizable in the sense of Exercise 2.1.7.7. From this perspective our problem becomes: determine polynomials on $A \otimes B \otimes C$ that vanish of the set of $T$ such that $T\left(A^{*}\right)$ is diagonalizable. (The problem is more naturally defined using the Grassmanian of Definition 2.3.3.1 below.)

A set of equations whose zero set is exactly the Zariski closure of the set of tensors giving rise to diagonalizable spaces of matrices is not known! What follows are some equations. (More are given in Chapter 5.) Recall that $B \otimes C=\operatorname{Hom}\left(C^{*}, B\right)$, the space of linear maps from $C^{*}$ to $B$. If instead we had $\operatorname{Hom}(B, B)=\operatorname{End}(B)$, the space of linear maps from $B$ to itself, a necessary condition for endomorphisms to be simultaneously diagonalizable is that they must commute, and the algebraic test for a subspace $U \subset$ $\operatorname{End}(B)$ to be abelian is simple: the commutators $\left[X_{i}, X_{j}\right]:=X_{i} X_{j}-X_{j} X_{i}$ must vanish on a basis $X_{1}, \ldots, X_{\mathbf{u}}$ of $U$. (I emphasize that commutators only make sense for maps from a vector space to itself.) These degree two equations exactly characterize abelian subspaces. We do not have maps from a vector space to itself, but we can fix the situation if there exists $\alpha \in A^{*}$ such that $T(\alpha): C^{*} \rightarrow B$ is invertible, as then we could test if the commutators $\left[T\left(\alpha_{1}\right) T(\alpha)^{-1}, T\left(\alpha_{2}\right) T(\alpha)^{-1}\right]$ are zero. So we now have a test, but it is not expressed in terms of polynomials on $A \otimes B \otimes C$, and we cannot apply it to all tensors. These problems are fixed in $\S 2.4 .1$. For now I record what we have so far:
Proposition 2.2.1.3. Let $\mathbf{b}=\mathbf{c}$ and let $T \in A \otimes B \otimes C$ be such that there exists $\alpha \in A^{*}$ with $\operatorname{rank}(T(\alpha))=\mathbf{b}$, so $\underline{\mathbf{R}}(T) \geq \mathbf{b}$. If $\underline{\mathbf{R}}(T)=\mathbf{b}$, then for all $X_{1}, X_{2} \in T\left(A^{*}\right) T(\alpha)^{-1} \subset \operatorname{End}(B),\left[X_{1}, X_{2}\right]=0$.
2.2.2. Strassen's equations: original formulation. If $T \in A \otimes B \otimes C$ is "close to" having rank $\mathbf{a}=\mathbf{b}=\mathbf{c}$, one expects, using $\alpha$ with $T(\alpha)$ invertible, that $T\left(A^{*}\right) T(\alpha)^{-1} \subset \operatorname{End}(B)$ will be "close to" being abelian. The following theorem makes this precise:
Theorem 2.2.2.1 (Strassen). [Str83] Let $T \in A \otimes B \otimes C$ and assume $\mathbf{b}=\mathbf{c}$. Assume that there exists $\alpha \in A^{*}$ such that $\operatorname{rank}(T(\alpha))=\mathbf{b}$. Then for all

$$
\begin{aligned}
& X_{1}, X_{2} \in T\left(A^{*}\right) T(\alpha)^{-1} \subset \operatorname{End}(B), \\
& \underline{\mathbf{R}}(T) \geq \frac{1}{2} \operatorname{rank}\left(\left[X_{1}, X_{2}\right]\right)+\mathbf{b} .
\end{aligned}
$$

I prove Theorem 2.2.2.1 for the case of the determinant of $\left[X_{1}, X_{2}\right.$ ] in §2.4.1 below and in general in §5.2.2.

We now have potential tests for border rank for tensors in $\mathbb{C}^{\mathbf{m}} \otimes \mathbb{C}^{\mathbf{m}} \otimes \mathbb{C}^{\mathbf{m}}$ up to $r=\frac{3}{2} \mathbf{m}$, in fact tests for border rank for tensors in $\mathbb{C}^{3} \otimes \mathbb{C}^{\mathbf{m}} \otimes \mathbb{C}^{\mathbf{m}}$ up to $r=\frac{3}{2} \mathbf{m}$, since our test only used three vectors from $A^{*}$. (I write "potential tests" rather than "polynomial tests" because to write down the commutator we must be able to find an invertible element in $T\left(A^{*}\right)$.)

Strassen uses Theorem 2.2.2.1 to show that $\underline{\mathbf{R}}\left(M_{\langle\mathbf{n}\rangle}\right) \geq \frac{3}{2} \mathbf{n}^{2}$ :
Exercise 2.2.2.2: (2!) Prove $\underline{\mathbf{R}}\left(M_{\langle\mathbf{n}\rangle}\right) \geq \frac{3}{2} \mathbf{n}^{2}$. ©
Exercise 2.2.2.3: (2) Show that $\underline{\mathbf{R}}\left(M_{\langle 2\rangle}^{r e d}\right)=5$ and for $\mathbf{m}>2$ that $\underline{\mathbf{R}}\left(M_{\langle\mathbf{m}, 2,2\rangle}^{\text {red }}\right) \geq$ $3 \mathbf{m}-1$, where $M_{\langle\mathbf{m}, 2,2\rangle}^{r e d}$ is $M_{\langle\mathbf{m}, 2,2\rangle}$ with $x_{1}^{1}$ set to zero.

A natural question arises: exchanging the roles of $A, B, C$ we obtain three sets of such equations - are the three sets of equations the same or different? We should have already asked this question for the three types of usual flattenings: are the equations coming from the minors of $T_{A}, T_{B}, T_{C}$ the same or different? It is easy to write down tensors where $\operatorname{rank}\left(T_{A}\right), \operatorname{rank}\left(T_{B}\right), \operatorname{rank}\left(T_{C}\right)$ are distinct, however for $2 \times 2$ minors, two sets of them vanishing implies the third does as well, see, §8.3.1, where these questions are answered with the help of representation theory.

One can generalize Strassen's equations by taking higher order commutators, see [LM08b]. These generalizations do give new equations, but they do not give equations for border rank beyond the $\frac{3}{2} \mathbf{b}$ of Strassen's equations.

An extensive discussion of Strassen's equations and generalizations appears in [Lan12, §7.6].

### 2.2.3. Coming attractions: border rank bounds beyond Strassen's

 equations. The following more complicated expression gives equations for $\hat{\sigma}_{r}$ for $r>\frac{3}{2} \mathbf{b}$ :Let $T \in \mathbb{C}^{5} \otimes \mathbb{C}^{\mathbf{b}} \otimes \mathbb{C}^{\mathbf{b}}$, write $T=a_{0} \otimes X_{0}+\cdots a_{4} \otimes X_{4}$ with $X_{j} \in B \otimes C$. Assume that $\operatorname{rank}\left(X_{0}\right)=\mathbf{b}$ and choose bases such that $X_{0}=\mathrm{Id}$. Consider the following $5 \mathbf{b} \times 5 \mathbf{b}$ matrix:

$$
T_{A}^{\wedge}=\left(\begin{array}{cccc}
0 & {\left[X_{1}, X_{2}\right]} & {\left[X_{1}, X_{3}\right]} & {\left[X_{1}, X_{4}\right]}  \tag{2.2.1}\\
{\left[X_{2}, X_{1}\right]} & 0 & {\left[X_{2}, X_{3}\right]} & {\left[X_{2}, X_{4}\right]} \\
{\left[X_{3}, X_{1}\right]} & {\left[X_{3}, X_{2}\right]} & 0 & {\left[X_{3}, X_{4}\right]} \\
{\left[X_{4}, X_{1}\right]} & {\left[X_{4}, X_{2}\right]} & {\left[X_{4}, X_{3}\right]} & 0
\end{array}\right) .
$$

The name $T_{A}^{\wedge^{2}}$ is explained in $\S 2.4 .2$ where the proof of the following proposition also appears.
Proposition 2.2.3.1. [LO15] Let $T \in \mathbb{C}^{5} \otimes \mathbb{C}^{\mathbf{b}} \otimes \mathbb{C}^{\mathbf{b}}$ be as written above. Then $\underline{\mathbf{R}}(T) \geq \frac{\operatorname{rank} T_{A}{ }^{2}}{3}$. If $T \in A \otimes \mathbb{C}^{\mathbf{b}} \otimes \mathbb{C}^{\mathbf{b}}$ with $\mathbf{a}>5$, one obtains the same result for all restrictions of $T$ to $\mathbb{C}^{5} \otimes \mathbb{C}^{\mathbf{b}} \otimes \mathbb{C}^{\mathbf{b}}$ for any $\mathbb{C}^{5} \subset A^{*}$.

In particular the minors of (2.2.1) give equations up to border rank $\frac{5}{3} \mathbf{b}$ for tensors in $\mathbb{C}^{\mathbf{a}} \otimes \mathbb{C}^{\mathbf{b}} \otimes \mathbb{C}^{\mathbf{c}}$ for $\mathbf{a} \geq 5$ and $\mathbf{b} \leq \mathbf{c}$.

I do not know how anyone would have found (2.2.1) without using the theory discussed in the next section. Hopefully this will motivate the theoryadverse reader to persevere through it.

### 2.3. Theory needed for the generalization of Strassen's equations

The matrices $\left[X_{1}, X_{2}\right]$ and the right hand side of (2.2.1) are part of a sequence of constructions giving better lower bounds for border rank for tensors. The limits of this method are lower bounds of $2 \mathbf{b}-3$. To describe them, we will need more language from multi-linear algebra. Our first task will be to generalize the space of skew-symmetric matrices. It will be convenient to generalize symmetric matrices at the same time. Before that I present a fundamental result in representation theory.
2.3.1. Schur's lemma. I take a short detour into elementary representation theory to prove a lemma everyone should know. Recall the definition of a $G$-module from §1.1.13.
Definition 2.3.1.1. Let $W_{1}, W_{2}$ be vector spaces, let $G$ be a group, and let $\rho_{j}: G \rightarrow G L\left(W_{j}\right), j=1,2$ be representations. A $G$-module homomorphism, or $G$-module map, is a linear map $f: W_{1} \rightarrow W_{2}$ such that $f\left(\rho_{1}(g) \cdot v\right)=$ $\rho_{2}(g) \cdot f(v)$ for all $v \in W_{1}$ and $g \in G$. One also says that $f$ is $G$-equivariant. For a group $G$ and $G$-modules $V$ and $W$, let $\operatorname{Hom}_{G}(V, W) \subset V^{*} \otimes W$ denote the vector space of $G$-module homomorphisms $V \rightarrow W$.

One says $W_{1}$ and $W_{2}$ are isomorphic $G$-modules if there exists a $G$ module homomorphism $W_{1} \rightarrow W_{2}$ that is a linear isomorphism.

Exercise 2.3.1.2: (1!!) Show that the image and kernel of a $G$-module homomorphism are $G$-modules.

The following easy lemma is central to representation theory:
Lemma 2.3.1.3 (Schur's Lemma). Let $G$ be a group, let $V$ and $W$ be irreducible $G$-modules and let $f: V \rightarrow W$ be a $G$-module homomorphism. Then either $f=0$ or $f$ is an isomorphism. If further $V=W$, then $f=\lambda \mathrm{Id}_{V}$ for some constant $\lambda$.
Exercise 2.3.1.4: (1!!) Prove Schur's Lemma.
We will see numerous examples illustrating the utility of Schur's Lemma. I cannot over-emphasize the importance of this simple Lemma. I use it every day of my mathematical life.

For any group $G, G$-module $M$, and irreducible $G$-module $V$, the isotypic component of $V$ in $M$ is the largest subspace of $M$ isomorphic to $V^{\oplus m_{V}}$ for some $m_{V}$. The integer $m_{V}$ is called the multiplicity of $V$ in $M$.

### 2.3.2. Symmetric and skew-symmetric tensors.

Exercise 2.3.2.1: (1) Let $X$ be a matrix representing a bilinear form on $\mathbb{C}^{\mathrm{m}}$, by $X(v, w)=v^{T} X w$. Show that if $X$ is a symmetric matrix, then $X(v, w)=X(w, v)$ and if $X$ is a skew-symmetric matrix, then $X(v, w)=$ $-X(w, v)$.

Recall that $\mathfrak{S}_{d}$ denotes the permutation group on $d$ elements.
Definition 2.3.2.2. A tensor $T \in V^{\otimes d}$ is said to be symmetric if $T\left(\alpha_{1}, \ldots, \alpha_{d}\right)=$ $T\left(\alpha_{\sigma(1)}, \ldots, \alpha_{\sigma(d)}\right)$ for all $\alpha_{1}, \ldots, \alpha_{d} \in V^{*}$ and all permutations $\sigma \in \mathfrak{S}_{d}$, and skew-symmetric if $T\left(\alpha_{1}, \ldots, \alpha_{d}\right)=\operatorname{sgn}(\sigma) T\left(\alpha_{\sigma(1)}, \ldots, \alpha_{\sigma(d)}\right)$ for all $\alpha_{1}, \ldots, \alpha_{d} \in$ $V^{*}$ and all $\sigma \in \mathfrak{S}_{d}$. Let $S^{d} V \subset V^{\otimes d}$ (resp. $\Lambda^{d} V \subset V^{\otimes d}$ ) denote the space of symmetric (resp. skew-symmetric) tensors.

The spaces $\Lambda^{d} V$ and $S^{d} V$ are independent of a choice of basis in $V$. In particular, the splitting

$$
\begin{equation*}
V^{\otimes 2}=S^{2} V \oplus \Lambda^{2} V \tag{2.3.1}
\end{equation*}
$$

of the space of matrices into the direct sum of symmetric and skew symmetric matrices is invariant under the action of $G L(V)$ given by: for $g \in G L(V)$ and $v \otimes w \in V \otimes V, v \otimes w \mapsto g v \otimes g w$.

Introduce the notations:

$$
x_{1} x_{2} \cdots x_{k}:=\sum_{\sigma \in \mathfrak{G}_{k}} x_{\sigma(1)} \otimes x_{\sigma(2)} \otimes \cdots \otimes x_{\sigma(k)} \in S^{k} V,
$$

and

$$
x_{1} \wedge x_{2} \wedge \cdots \wedge x_{k}:=\sum_{\sigma \in \mathfrak{S}_{k}} \operatorname{sgn}(\sigma) x_{\sigma(1)} \otimes x_{\sigma(2)} \otimes \cdots \otimes x_{\sigma(k)} \in \Lambda^{k} V,
$$

respectively called the symmetric product (or simply product) of $x_{1}, \ldots, x_{k}$ and the wedge product of $x_{1}, \ldots, x_{k}$.

The space $S^{k} V^{*}$ may be thought of as the space of homogeneous polynomials of degree $k$ on $V$ (to a symmetric tensor $T$ associate the polynomial $P_{T}$ where $\left.P_{T}(v):=T(v, \ldots, v)\right)$. Thus $x_{1} \cdots x_{k}$ may also be read as the multiplication of $x_{1}, \ldots, x_{k}$.

If $v_{1}, \ldots, v_{\mathbf{v}}$ is a basis of $V$, then $v_{i_{1}} \otimes \cdots \otimes v_{i_{d}}$ with $i_{j} \in[\mathbf{v}]:=\{1, \ldots, \mathbf{v}\}$ is a basis of $V^{\otimes d}, v_{i_{1}} \cdots v_{i_{d}}$ with $1 \leq i_{1} \leq \cdots \leq i_{d} \leq \mathbf{v}$ is a basis of $S^{d} V$ and $v_{i_{1}} \wedge \cdots \wedge v_{i_{d}}$ with $1 \leq i_{1}<\cdots<i_{d} \leq \mathbf{v}$ is a basis of $\Lambda^{d} V$. Call these bases induced bases. If $x_{j}=\left(x_{j}^{1}, \ldots, x_{j}^{\mathbf{V}}\right)^{T}$ in the basis $v_{1}, \ldots, v_{\mathbf{v}}$, then the expression of $x_{1} \wedge \cdots \wedge x_{k}$ in the induced basis is such that the coefficient of $v_{i_{1}} \wedge \cdots \wedge v_{i_{k}}$ is

$$
\operatorname{det}\left(\begin{array}{ccc}
x_{1}^{i_{1}} & \cdots & x_{1}^{i_{k}} \\
& \vdots & \\
x_{k}^{i_{1}} & \cdots & x_{k}^{i_{k}}
\end{array}\right) .
$$

For example, if $V=\mathbb{C}^{4}$ with basis $e_{1}, \ldots, e_{4}$, then $\Lambda^{2} V$ inherits a basis $e_{1} \wedge e_{2}, \ldots, e_{3} \wedge e_{4}$. If

$$
v=\left(\begin{array}{l}
v_{1} \\
v_{2} \\
v_{3} \\
v_{4}
\end{array}\right), w=\left(\begin{array}{l}
w_{1} \\
w_{2} \\
w_{3} \\
w_{4}
\end{array}\right), \text { then } v \wedge w=\left(\begin{array}{l}
v_{1} w_{2}-v_{2} w_{1} \\
v_{1} w_{3}-v_{3} w_{1} \\
v_{1} w_{4}-v_{4} w_{1} \\
v_{2} w_{3}-v_{3} w_{2} \\
v_{2} w_{4}-v_{4} w_{2} \\
v_{3} w_{4}-v_{4} w_{3}
\end{array}\right) .
$$

Exercise 2.3.2.3: (1) Show that there is a $G L(V)$-module map $\Lambda^{k} V \otimes V \rightarrow$ $\Lambda^{k+1} V$, and more generally there are $G L(V)$-module maps $\Lambda^{k} V \otimes \Lambda^{t} V \rightarrow$ $\Lambda^{k+t} V$ and $S^{k} V \otimes S^{t} V \rightarrow S^{k+t} V$, the latter of which may be interpreted as multiplication of polynomials.

Exercise 2.3.2.4: (1) Let $k \geq t$ and show that there is a $G L(V)$-module map $S^{k} V^{*} \otimes S^{t} V \rightarrow S^{k-t} V^{*}$. This map has the following interpretation: $S^{t} V$ may be interpreted as the homogeneous linear differential operators of order $t$ on the space of polynomials $S^{k} V^{*}$. The map is then $P \otimes D \mapsto D(P)$. Sometimes $D(P)$ is denoted $D\lrcorner P$.

Exercise 2.3.2.5: (1) Show that for $k<l$ there is a $G L(V)$-module map, $\Lambda^{k} V^{*} \otimes \Lambda^{l} V \rightarrow \Lambda^{l-k} V$ that commutes with the action of $G L(V)$. This map is often denoted $\beta \otimes Y \mapsto \beta\lrcorner Y$

Exercise 2.3.2.6: (1) Let $\operatorname{Sym}(V)=\oplus_{j=0}^{\infty} S^{j} V, \Lambda^{\bullet} V=\oplus_{j=0}^{\mathbf{v}} \Lambda^{j} V$ and $V^{\otimes \bullet}=\oplus_{j=0}^{\infty} V^{\otimes j}$. Show that these spaces are all naturally algebras with
the above defined products, respectively called the symmetric, exterior and tensor algebras.
2.3.3. The Grassmannian. Before returning to border rank, I define an important algebraic variety that we will need for the proof of tensor rank lower bounds:

Definition 2.3.3.1. The Grassmannian of $k$-planes through the origin in $V$ is
$G(k, V):=\mathbb{P}\left\{T \in \Lambda^{k} V \mid \exists v_{1}, \ldots, v_{k} \in V\right.$ such that $\left.T=v_{1} \wedge \cdots \wedge v_{k}\right\} \subset \mathbb{P} \Lambda^{k} V$.
The most important special case of a Grassmannian is projective space $\mathbb{P} V=G(1, V) . \mathbb{P} V:=(V \backslash 0) / \sim$ where $v \sim w$ if and only if $v=\lambda w$ for some $\lambda \in \mathbb{C} \backslash 0$.

The interpretation of the Grassmannian as the space parameterizing the $k$-planes through the origin in $V$ is via the correspondence $\left[v_{1} \wedge \cdots \wedge v_{k}\right] \leftrightarrow$ $\operatorname{span}\left\{v_{1}, \ldots, v_{k}\right\}$.

The following exercise shows that the Grassmannian is indeed an algebraic variety. It can be safely skipped on a first reading.
Exercise 2.3.3.2: (3) The Grassmannian is the zero set of equations parametrized by $\Lambda^{k-2 j} V^{*} \otimes \Lambda^{k+2 j} V^{*}$ for $1 \leq j \leq \min \left\{\left\lfloor\frac{\mathbf{v}-k}{2}\right\rfloor,\left\lfloor\frac{k}{2}\right\rfloor\right\}$ as follows: for $\mu \in$ $\Lambda^{k-2 j} V^{*}$ and $\zeta \in \Lambda^{k+2 j} V^{*}$, recall Exercise 2.3.2.5, and consider $\left.T\right\lrcorner \zeta \in \Lambda^{2 j} V^{*}$ and $\mu \omega T \in \Lambda^{2 j} V$. Define $\left.P_{\mu \otimes \zeta}(T):=\langle T\lrcorner \zeta, \mu-T\right\rangle$, the evaluation of an element of $\Lambda^{2 j} V^{*}$ on an element of $\Lambda^{2 j} V$. Note that these are quadratic equations in the coefficients of $T$. Show that the zero set of these equations is the Grassmannian. ©

### 2.4. Koszul flattenings

2.4.1. Reformulation and proof of Strassen's equations. Augment the linear map $T_{B}: B^{*} \rightarrow A \otimes C$ by tensoring it with $\mathrm{Id}_{A}$, to obtain a linear map

$$
\mathrm{Id}_{A} \otimes T_{B}: A \otimes B^{*} \rightarrow A \otimes A \otimes C
$$

So far this is not interesting, but by (2.3.1) the target of this map decomposes as a $G L(A) \times G L(C)$-module as $\left(\Lambda^{2} A \otimes C\right) \oplus\left(S^{2} A \otimes C\right)$, and we may project onto these factors. Write the projections as:

$$
\begin{equation*}
T_{B A}^{\wedge}=T_{A}^{\wedge}: A \otimes B^{*} \rightarrow \Lambda^{2} A \otimes C \text { and } T_{B A}^{\circ}: A \otimes B^{*} \rightarrow S^{2} A \otimes C \tag{2.4.1}
\end{equation*}
$$

Exercise 2.4.1.1: (1) Show that if $T=a \otimes b \otimes c$ is a rank one tensor, then $\operatorname{rank}\left(T_{A}^{\wedge}\right)=\mathbf{a}-1$ and $\operatorname{rank}\left(T_{B A}^{\circ}\right)=\mathbf{a}$.

Exercise 2.4.1.1 implies:

Proposition 2.4.1.2. If $\underline{\mathbf{R}}(T) \leq r$, than $\operatorname{rank}\left(T_{A}^{\wedge}\right) \leq r(\mathbf{a}-1)$ and $\operatorname{rank}\left(T_{B A}^{\circ}\right) \leq$ ra.

The second map will not give border rank lower bounds better than the classical equations, but the first, e.g., when $\mathbf{a}=3$, is a map from a $2 \mathbf{b}$ dimensional vector space to a $2 \mathbf{c}$ dimensional vector space, so if $\mathbf{b} \leq \mathbf{c}$ we can get border rank bounds up to $\frac{3}{2} \mathbf{b}$.

The first set is equivalent to Strassen's equations, as I now show. If a $>3$, one can choose a three dimensional subspace $A^{\prime} \subset A^{*}$ and consider $T$ restricted to $A^{\prime} \times B^{*} \times C^{*}$ to obtain equations. (This is what we did in the case of Strassen's equations where $A^{\prime}$ was spanned by $\alpha, \alpha^{\prime}, \alpha^{\prime \prime}$.)

Let $a_{1}, a_{2}, a_{3}$ be a basis of $A$, with dual basis $\alpha^{1}, \alpha^{2}, \alpha^{3}$ of $A^{*}$ so $T \in$ $A \otimes B \otimes C$ may be written as $T=a_{1} \otimes X_{1}+a_{2} \otimes X_{2}+a_{3} \otimes X_{3}$, where $X_{j}=$ $T\left(\alpha_{j}\right)$. Then $T_{A}^{\wedge}$ will be expressed by a $3 \mathbf{b} \times 3 \mathbf{b}$ matrix. Ordering the basis of $A \otimes B^{*}$ by $a_{3} \otimes \beta^{1}, \ldots, a_{3} \otimes \beta^{\mathbf{b}}, a_{2} \otimes \beta^{1}, \ldots, a_{2} \otimes \beta^{\mathbf{b}}, a_{1} \otimes \beta^{1}, \ldots, a_{1} \otimes \beta^{\mathbf{b}}$, and that of $\Lambda^{2} A \otimes C$ by $\left(a_{1} \wedge a_{2}\right) \otimes c_{1}, \ldots,\left(a_{1} \wedge a_{2}\right) \otimes c_{\mathbf{b}},\left(a_{1} \wedge a_{3}\right) \otimes c_{1}, \ldots,\left(a_{1} \wedge\right.$ $\left.a_{3}\right) \otimes c_{\mathbf{b}},\left(a_{2} \wedge a_{3}\right) \otimes c_{1}, \ldots,\left(a_{2} \wedge a_{3}\right) \otimes c_{\mathbf{b}}$, we obtain the block matrix

$$
T_{A}^{\wedge}=\left(\begin{array}{ccc}
0 & X_{1} & -X_{2}  \tag{2.4.2}\\
X_{2} & X_{3} & 0 \\
X_{1} & 0 & X_{3}
\end{array}\right) .
$$

Recall the following basic identity about determinants of blocked matrices (see, e.g., [Pra94, Thm. 3.1.1]), assuming the block $W$ is invertible:

$$
\operatorname{det}\left(\begin{array}{cc}
X & Y  \tag{2.4.3}\\
Z & W
\end{array}\right)=\operatorname{det}(W) \operatorname{det}\left(X-Y W^{-1} Z\right)
$$

Block (2.4.2) $X=0, Y=\left(X_{1},-X_{2}\right), Z=\binom{X_{2}}{X_{1}}, W=\left(\begin{array}{cc}X_{3} & 0 \\ 0 & X_{3}\end{array}\right)$. Assume $X_{3}=T\left(\alpha^{3}\right)$ is invertible to obtain

$$
\begin{equation*}
\operatorname{det} T_{A}^{\wedge}=\operatorname{det}\left(X_{3}\right)^{2} \operatorname{det}\left(X_{1} X_{3}^{-1} X_{2}-X_{2} X_{3}^{-1} X_{1}\right) \tag{2.4.4}
\end{equation*}
$$

Equation (2.4.4) shows the new formulation is equivalent to the old, at least in the case of maximal rank. (We are only interested in the non-vanishing of the polynomial, not its values, so we can multiply the inner matrix on the right by $X_{3}{ }^{-1}$.) Equation (2.4.4) combined with Proposition 2.4.1.2 proves Theorem 2.2.2.1 in this case.

Note that here we have actual polynomials on $A \otimes B \otimes C$ (the minors of (2.4.2)), whereas in our original formulation of Strassen's equations we did not. To obtain polynomials in the original formulation one uses the adjugate matrix instead of the inverse, see [Lan12, §3.8].

Remark 2.4.1.3. Both the classical equations and Strassen's equations are obtained by taking minors of a matrix whose entries are linear combinations of the coefficients of our tensor. Such constructions are part of a long tradition of finding determinantal equations for algebraic varieties discussed further in Chapters 8 and 10. For the experts, given a variety $X$ and a subvariety $Y \subset X$, one way to find defining equations for $Y$ is to find vector bundles $E, F$ over $X$ and a vector bundle map $\phi: E \rightarrow F$ such that $Y$ is realized as the degeneracy locus of $\phi$, that is, the set of points $x \in X$ such that $\phi_{x}$ drops rank. Strassen's equations in the partially symmetric case had been discovered by Barth $[\operatorname{Bar} 77]$ in this context.
Remark 2.4.1.4. In $\S 8.2$ and $\S 8.3 .1$, we will see two different ways of deriving Strassen's equations via representation theory.
2.4.2. Definition of Koszul flattenings. The reformulation of Strassen's equations suggests the following generalization: let $\operatorname{dim} A=2 p+1$ and consider

$$
\begin{equation*}
T_{A}^{\wedge p}: B^{*} \otimes \Lambda^{p} A \rightarrow \Lambda^{p+1} A \otimes C \tag{2.4.5}
\end{equation*}
$$

obtained by first taking $T_{B} \otimes \operatorname{Id}_{\Lambda^{p}} A: B^{*} \otimes \Lambda^{p} A \rightarrow \Lambda^{p} A \otimes A \otimes C$, and then projecting to $\Lambda^{p+1} A \otimes C$ as in Exercise 2.3.2.3.

If $\left\{a_{i}\right\},\left\{b_{j}\right\},\left\{c_{k}\right\}$ are bases of $A, B, C$ and $T=\sum_{i, j, k} t^{i j k} a_{i} \otimes b_{j} \otimes c_{k}$, then

$$
\begin{equation*}
T_{A}^{\wedge p}\left(\beta \otimes f_{1} \wedge \cdots \wedge f_{p}\right)=\sum_{i, j, k} t^{i j k} \beta\left(b_{j}\right) a_{i} \wedge f_{1} \wedge \cdots \wedge f_{p} \otimes c_{k} \tag{2.4.6}
\end{equation*}
$$

The map $T_{A}^{\wedge p}$ is called a Koszul flattening. Note that if $T=a \otimes b \otimes c$ has rank one, then $\operatorname{rank}\left(T_{A}^{\wedge p}\right)=\binom{2 p}{p}$ as the image is $a \wedge \Lambda^{p} A \otimes c$. By linearity of the map $T \mapsto T_{A}^{\wedge p}$ we conclude:
Proposition 2.4.2.1. [LO15] Let $T \in A \otimes B \otimes C$ with $\operatorname{dim} A=2 p+1$. Then

$$
\underline{\mathbf{R}}(T) \geq \frac{\operatorname{rank}\left(T_{A}^{\wedge p}\right)}{\binom{2 p}{p}}
$$

Since the source (resp. target) has dimension $\binom{2 p+1}{p} \mathbf{b}\left(\right.$ resp. $\left.\binom{2 p+1}{p+1} \mathbf{c}\right)$, assuming $\mathbf{b} \leq \mathbf{c}$, we potentially obtain equations for $\hat{\sigma}_{r}$ up to

$$
r=\frac{\binom{2 p+1}{p} \mathbf{b}}{\binom{2 p}{p}}-1=\frac{2 p+1}{p+1} \mathbf{b}-1 .
$$

Just as with Strassen's equations (case $p=1$ ), if $\operatorname{dim} A>2 p+1$, one obtains the best bound for these equations by restricting to subspaces of $A^{*}$ of dimension $2 p+1$.
Exercise 2.4.2.2: (2) Show that if $T_{A}^{\wedge p}: \Lambda^{p} A \otimes B^{*} \rightarrow \Lambda^{p+1} A \otimes C$ is injective, then $T_{A}^{\wedge q}: \Lambda^{q} A \otimes B^{*} \rightarrow \Lambda^{q+1} A \otimes C$ is injective for all $q<p$. ©
2.4.3. Koszul flattenings in coordinates. To prove lower bounds on the rank of matrix multiplication, and to facilitate a comparison with Griesser's equations discussed in $\S 5.2 .2$, it will be useful to view $T_{A}^{\wedge p}$ in coordinates. Let $\operatorname{dim} A=2 p+1$. Write $T=a_{0} \otimes X_{0}+\cdots+a_{2 p} \otimes X_{2 p}$ where $a_{j}$ is a basis of $A$ with dual basis $\alpha^{j}$ and $X_{j}=T\left(\alpha^{j}\right)$. An expression of $T_{A}^{\wedge p}$ in bases is as follows: write $a_{I}:=a_{i_{1}} \wedge \cdots \wedge a_{i_{p}}$ for the induced basis elements of $\Lambda^{p} A$, require that the first $\binom{2 p}{p-1}$ basis vectors of $\Lambda^{p} A$ have $i_{1}=0$, that the second $\binom{2 p}{p}$ do not, and call these multi-indices $0 J$ and $K$. Order the bases of $\Lambda^{p+1} A$ such that the first $\binom{2 p}{p+1}$ multi-indices do not have 0 , and the second $\binom{2 p}{p}$ do, and furthermore that the second set of indices is ordered the same way as $K$ is ordered, only we write $0 K$ since a zero index is included. The resulting matrix is of the form

$$
\left(\begin{array}{cc}
0 & Q  \tag{2.4.7}\\
\tilde{Q} & R
\end{array}\right)
$$

where this matrix is blocked $\left(\binom{2 p}{p+1} \mathbf{b},\binom{2 p}{p} \mathbf{b}\right) \times\left(\binom{2 p}{p+1} \mathbf{b},\binom{2 p}{p} \mathbf{b}\right)$,

$$
R=\left(\begin{array}{lll}
X_{0} & & \\
& \ddots & \\
& & X_{0}
\end{array}\right)
$$

and $Q, \tilde{Q}$ have entries in blocks consisting of $X_{1}, \ldots, X_{2 p}$ and zero. Thus if $X_{0}$ is of full rank and we change coordinates such that it is the identity matrix, so is $R$ and the determinant equals the determinant of $Q \tilde{Q}$ by (2.4.3). If we order the appearances of the $K$ multi-indices such that the $j$-th $K$ is the complement of the $j$-th $J$ in $[2 p]$, then $Q \tilde{Q}$ will be block skew-symmetric. When $p=1, Q \tilde{Q}=\left[X_{1}, X_{2}\right]$, and when $p=2$ we recover the matrix (2.2.1).

In general $Q \tilde{Q}$ is a block skew-symmetric $\binom{2 p}{p-1} \mathbf{b} \times\binom{ 2 p}{p-1} \mathbf{b}$ matrix whose block entries are either zero or commutators $\left[X_{i}, X_{j}\right]$. Each $\left[X_{i}, X_{j}\right]$ appears (up to sign) $\binom{2 p-1}{2}$ times, and each block row and column contain exactly $\binom{2 p-1}{2}$ nonzero blocks, so the resulting matrix is very sparse.

### 2.5. Matrix multiplication and Koszul flattenings

We would like to apply our new equations to matrix multiplication. In order to do so, we first must understand the matrix multiplication tensor better from a geometric perspective.
2.5.1. The matrix multiplication tensor from an invariant perspective. In the vector space $V^{*} \otimes V$ there is a unique line such that every vector on the line has the same matrix representative for any choice of basis (and corresponding choice of dual basis). This line is of course $\mathbb{C}\left\{\operatorname{Id}_{V}\right\}$, the scalar
multiples of the identity map. We say $\mathbb{C}\left\{\operatorname{Id}_{V}\right\}$ is the unique line in $V^{*} \otimes V$ invariant under the action of $G L(V)$.

We have

$$
M_{\langle U, V, W\rangle} \in\left(U^{*} \otimes V\right) \otimes\left(V^{*} \otimes W\right) \otimes\left(W^{*} \otimes U\right) \simeq U^{*} \otimes U \otimes V^{*} \otimes V \otimes W^{*} \otimes W
$$

Moreover, we know $M_{\langle U, V, W\rangle}$ is invariant under the action of $G L(U) \times$ $G L(V) \times G L(W)$. The only element of $U^{*} \otimes U \otimes V^{*} \otimes V \otimes W^{*} \otimes W$ that is invariant under $G L(U) \times G L(V) \times G L(W)$ is up to scale $\operatorname{Id}_{U} \otimes \operatorname{Id}_{V} \otimes \mathrm{Id}_{W}$. Checking the scale, we conclude:
Proposition 2.5.1.1. $M_{\langle U, V, W\rangle}$, after applying the re-ordering isomorphism, is $\operatorname{Id}_{U} \otimes \operatorname{Id}_{V} \otimes \operatorname{Id}_{W}$.
Exercise 2.5.1.2: (1) If $v_{1}, \ldots, v_{\mathbf{v}}$ is a basis of $V$ and $\alpha^{1}, \ldots, \alpha^{\mathbf{v}}$ is the dual basis of $V^{*}$, show that the identity map on $V$ is $\mathrm{Id}_{V}=\sum_{j} \alpha^{j} \otimes v_{j}$.
Exercise 2.5.1.3: (1) Use Exercise 2.5.1.2 and the coordinate presentation of matrix multiplication to get a second proof of Proposition 2.5.1.1. This proof also shows that $M_{\langle U, V, W\rangle}$ is invariant under the action of the image of $G L(U) \times G L(V) \times G L(W)$ in $G L(A) \times G L(B) \times G L(C)$.
Exercise 2.5.1.4: (1) Show that there is a canonical isomorphism $\left(V^{*} \otimes W\right)^{*} \rightarrow$ $V \otimes W^{*}$ where $\alpha \otimes w(v \otimes \beta):=\alpha(v) \beta(w)$. Now let $V=W$ and let $\mathrm{Id}_{V} \in$ $V^{*} \otimes V \simeq\left(V^{*} \otimes V\right)^{*}$ denote the identity map. What is $\operatorname{Id}_{V}(f)$ for $f \in V^{*} \otimes V$ ? ©
Exercise 2.5.1.5: (1!) Show that $M_{\langle U, V, W\rangle}$ when viewed as a trilinear map

$$
M_{\langle U, V, W\rangle}:\left(U^{*} \otimes V\right)^{*} \times\left(V^{*} \otimes W\right)^{*} \times\left(W^{*} \otimes U\right)^{*} \rightarrow \mathbb{C}
$$

is $(X, Y, Z) \mapsto \operatorname{trace}(X Y Z)$. ©
Exercise 2.5.1.6: (1!) Using Exercise 2.5.1.5, show that $M_{\langle\mathbf{n}\rangle} \in \mathbb{C}^{\mathbf{n}^{2}} \otimes \mathbb{C}^{\mathbf{n}^{2}} \otimes \mathbb{C}^{\mathbf{n}^{2}}$ is preserved by the cyclic permutation of the factors.
Exercise 2.5.1.7: (1!) Using Exercise 2.5.1.5, show that $M_{\langle\mathbf{n}\rangle} \in \mathbb{C}^{\mathbf{n}^{2}} \otimes \mathbb{C}^{\mathbf{n}^{2}} \otimes \mathbb{C}^{\mathbf{n}^{2}}$ is preserved by action $x \otimes y \otimes z \mapsto x^{T} \otimes z^{T} \otimes y^{T}$, where $x^{T}$ is the transpose of the $\mathbf{n} \times \mathbf{n}$ matrix $x$.
Exercise 2.5.1.8: (1) Show that $\mathrm{Id}_{V} \otimes \mathrm{Id}_{W} \in V \otimes V^{*} \otimes W \otimes W^{*}=(V \otimes W) \otimes(V \otimes W)^{*}$, after re-ordering, equals $\operatorname{Id}_{V \otimes W}$.
Exercise 2.5.1.9: (1!) Using Exercise 2.5.1.8, show that $M_{\langle\mathbf{n}, \mathbf{m}, \mathbf{1}\rangle} \otimes M_{\left\langle\mathbf{n}^{\prime}, \mathbf{m}^{\prime}, \mathbf{l}^{\prime}\right\rangle}=$ $M_{\left\langle\mathbf{n n}^{\prime}, \mathbf{m m}^{\prime}, \mathbf{I I}^{\prime}\right\rangle}$.

A fancy proof that $\underline{\mathbf{R}}\left(M_{\langle\mathbf{n}\rangle}\right) \geq \mathbf{n}^{2}$, which will be useful for proving further lower bounds, is as follows: Write $A=U^{*} \otimes V, B=V^{*} \otimes W, C=$ $W^{*} \otimes U$, so $\left(M_{\langle\mathbf{n}\rangle}\right)_{A}: A^{*} \rightarrow B \otimes C$ is a map $U \otimes V^{*} \rightarrow V^{*} \otimes W \otimes W^{*} \otimes U$. This map is, for $f \in A^{*}, f \mapsto f \otimes \operatorname{Id}_{W}$, and thus is clearly injective. In other
words, the map is $u \otimes \nu \mapsto \sum_{k}\left(\nu \otimes w_{k}\right) \otimes\left(w^{k} \otimes u\right)$, where $w_{1}, \ldots, w_{\mathbf{w}}$ is a basis of $W$ with dual basis $w^{1}, \ldots, w^{\mathrm{w}}$.
2.5.2. Koszul flattenings and matrix multiplication. When $T=M_{\langle U, V, W\rangle}$, the Koszul flattening map is

$$
\left(M_{\langle U, V, W\rangle}\right)_{A}^{\wedge p}: V \otimes W^{*} \otimes \Lambda^{p}\left(U^{*} \otimes V\right) \rightarrow \Lambda^{p+1}\left(U^{*} \otimes V\right) \otimes\left(W^{*} \otimes U\right) .
$$

The presence of $\operatorname{Id}_{W}=\operatorname{Id}_{W^{*}}$ implies the map factors as $\left(M_{\langle U, V, W\rangle}\right)_{A}^{\wedge p}=$ $\left(M_{\langle\mathbf{u}, \mathbf{v}, 1\rangle}\right\rangle_{A}^{\wedge p} \otimes \mathrm{Id}_{W^{*}}$, where

$$
\begin{align*}
\left(M_{\langle\mathbf{u}, \mathbf{v}, 1\rangle}\right)_{A}^{\wedge p}: V \otimes \Lambda^{p}\left(U^{*} \otimes V\right) & \rightarrow \Lambda^{p+1}\left(U^{*} \otimes V\right) \otimes U .  \tag{2.5.1}\\
v \otimes\left(\xi^{1} \otimes e_{1}\right) \wedge \cdots \wedge\left(\xi^{p} \otimes e_{p}\right) & \mapsto \sum_{s=1}^{\mathrm{u}} u_{s} \otimes\left(u^{s} \otimes v\right) \wedge\left(\xi^{1} \otimes e_{1}\right) \wedge \cdots \wedge\left(\xi^{p} \otimes e_{p}\right) .
\end{align*}
$$

where $u_{1}, \ldots, u_{\mathbf{u}}$ is a basis of $U$ with dual basis $u^{1}, \ldots, u^{\mathbf{u}}$ of $U^{*}$, so $\operatorname{Id}_{U}=$ $\sum_{s=1}^{\mathbf{u}} u^{s} \otimes u_{s}$.

As discussed above, Koszul flattenings could potentially prove a border rank lower bound of $2 \mathbf{n}^{2}-3$ for $M_{\langle\mathbf{n}\rangle}$. However this does not happen, as there is a large kernel for the maps $M_{\langle\mathbf{n}\rangle}^{\wedge p}$ when $p \geq \mathbf{n}$ : Let $\mathbf{u}=\mathbf{v}=\mathbf{n}$. and let $p=\mathbf{n}$. Then

$$
v \otimes\left(u^{1} \otimes v\right) \otimes \cdots \otimes\left(u^{\mathbf{n}} \otimes v\right) \mapsto \sum_{j}\left(u^{j} \otimes v\right) \wedge\left(u^{1} \otimes v\right) \otimes \cdots \otimes\left(u^{\mathbf{n}} \otimes v\right) \otimes u_{j}=0,
$$

so $M_{\langle\mathbf{n}\rangle}^{\wedge \mathbf{n}}$ is not injective. Since $\left.M_{\langle\mathbf{u}, \mathbf{v}, 1\rangle}\right)_{A}^{\wedge p}$ is a $G L(U) \times G L(V)$-module map, by Schur's lemma 2.3.1.3, $\operatorname{ker}\left(M_{\langle\mathbf{n}\rangle}^{\wedge \mathbf{n}}\right) \subset V \otimes \Lambda^{\mathbf{n}}\left(U^{*} \otimes V\right) \subset V^{\otimes \mathbf{n}+1} \otimes U^{* \otimes \mathbf{n}}$ must be a submodule. It is clearly symmetric in $V$ and skew in $U^{*}$, so the kernel must contain the irreducible submodule $\Lambda^{\mathbf{n}} U^{*} \otimes S^{\mathbf{n}+1} V$.

Now consider the case $p=\mathbf{n}-1$. I claim $\left(M_{\langle\mathbf{n}, \mathbf{n}, 1\rangle}\right)_{A}^{\wedge \mathbf{n}-1}$ is injective. The following argument is due to L. Manivel. Say $X_{1} \otimes v_{1}+\cdots+X_{\mathbf{n}} \otimes v_{\mathbf{n}} \in$ $\operatorname{ker}\left(M_{\langle\mathbf{n}, \mathbf{n}, 1\rangle}\right)_{A}^{\wedge \mathbf{n}-1}$, i.e.,

$$
\sum_{s}\left[X_{1} \wedge\left(u^{s} \otimes v_{1}\right)+\cdots+X_{\mathbf{n}} \wedge\left(u^{s} \otimes v_{\mathbf{n}}\right)\right] \otimes u_{s}=0
$$

Then for each $s$, each term in the brackets must be zero.
Lemma 2.5.2.1. Let $A$ be a vector space, let $X_{1}, \ldots, X_{k} \in \Lambda^{q} A$, and let $a_{1}, \ldots, a_{k} \in A$ be linearly independent. Then if $X_{1} \wedge a_{1}+\cdots+X_{k} \wedge a_{k}=0$, we may write each $X_{j}=\sum_{i=1}^{k} Y_{i j} \wedge a_{i}$ for some $Y_{i j} \in \Lambda^{q-1} A$.
Exercise 2.5.2.2: (2) Prove Lemma 2.5.2.1.○
Exercise 2.5.2.3: (2) Show that $\operatorname{ker}\left(M_{\langle\mathbf{n}, \mathbf{n}, 1\rangle}\right)_{A}^{\wedge \mathbf{n}}=\Lambda^{\mathbf{n}} U^{*} \otimes S^{\mathbf{n}+1} V$. ©

Remark 2.5.2.4. This is a special case of the generalized Cartan Lemma, see [IL16b, §A.1]. With the aid of representation theory one can more precisely describe the $Y_{j i}$. (For those familiar with the notation, use the sequence $0 \rightarrow S_{2,1^{q-1}} A \rightarrow \Lambda^{q} A \otimes A \rightarrow \Lambda^{q+1} A \rightarrow 0$.)

Returning to the proof of injectivity when $p=\mathbf{n}-1$, taking $s=1$, we have $X_{j}=\sum Y_{j,(1, i)} \wedge\left(u^{1} \otimes a_{i}\right)$, so each term in $X_{j}$ is divisible by $\left(u^{1} \otimes a_{i}\right)$ for some $i$, but then taking $s=2$, each term in $X_{j}$ is divisible by ( $u^{2} \otimes a_{l}$ ) for some $l$. Continuing, if $p<\mathbf{n}$ we run out of factors, so there cannot be a kernel. In summary:
Proposition 2.5.2.5. When $p<\mathbf{n}$, the map $\left(M_{\langle\mathbf{n}, \mathbf{n}, 1\rangle}\right)_{A}^{\wedge p}$ is injective.
At this point one would like to say that if some $T^{\wedge p}$ is injective, then restricting to a generic $A^{\prime} \subset A^{*}$, the map $\left.T^{\wedge p}\right|_{\Lambda^{p} A^{\prime} \otimes B^{*}}: \Lambda^{p} A^{\prime} \otimes B^{*} \rightarrow$ $\Lambda^{p+1} A^{\prime} \otimes C$ would still be injective. Unfortunately I do not know how to prove this, because a priori $\left.T^{\wedge p}\right|_{\Lambda^{p} A^{\prime} \otimes B^{*}}$ injects into $\left[\Lambda^{p+1} A^{\prime} \otimes C\right] \oplus\left[\Lambda^{p} A^{\prime} \otimes\left(A / A^{\prime}\right) \otimes C\right]$, and it is not clear to me whether for generic $A^{\prime}$ it must remain injective when one projects to the first factor. What follows are two proofs that this is indeed the case for $\left(M_{\langle\mathbf{n}, \mathbf{n}, 1\rangle}\right)_{A}^{\wedge \mathbf{n}-1}$. The first is combinatorial. It has the advantages that it is elementary and will be used to prove the $2 \mathbf{n}^{2}-\left\lceil\log _{2} \mathbf{n}\right\rceil-1$ lower bound of $\S 5.4 .5$. The second is geometrical. It has the advantage of being shorter and more elegant.
Theorem 2.5.2.6. [LO15] Let $\mathbf{n} \leq \mathbf{m}$. Then

$$
\underline{\mathbf{R}}\left(M_{\langle\mathbf{m}, \mathbf{n}, \mathbf{l}\rangle}\right) \geq \frac{\mathbf{n l}(\mathbf{n}+\mathbf{m}-1)}{\mathbf{m}} .
$$

In particular $\underline{\mathbf{R}}\left(M_{\langle\mathbf{n}\rangle}\right) \geq 2 \mathbf{n}^{2}-\mathbf{n}$.
I prove the case $\mathbf{n}=\mathbf{m}$ and leave the general case to the reader. We need to find $A^{\prime} \subset A^{*}$ of dimension $2 \mathbf{n}-1$ such that, setting $\tilde{A}=A / A^{\prime \perp} \simeq A^{\prime *}$, $\left(\left.M_{\langle\mathbf{n}, \mathbf{n}, 1\rangle}\right|_{A^{\prime} \otimes B^{*} \otimes C^{*}}\right)_{\tilde{A}}^{\wedge \mathbf{n}-1}$ is injective.

First proof. Define the projection

$$
\begin{align*}
\phi: A & \rightarrow \mathbb{C}^{2 \mathbf{n}-1}  \tag{2.5.2}\\
x_{j}^{i} & \mapsto e_{i+j-1} . \tag{2.5.3}
\end{align*}
$$

Let $e_{S}:=e_{s_{1}} \wedge \cdots \wedge e_{s_{\mathbf{n}-1}}$, where $S=\left\{s_{1}, \ldots, s_{\mathbf{n}-1}\right\} \subset[2 \mathbf{n}-1]$. The $\operatorname{map}\left(\left.M_{\langle\mathbf{n}, \mathbf{n}, 1\rangle}\right|_{A^{\prime} \otimes B^{*} \otimes C^{*}}\right)_{\tilde{A}}^{\wedge \mathbf{n}-1}$ is

$$
e_{S} \otimes v_{k} \mapsto \sum_{j} \phi\left(u^{j} \otimes v_{k}\right) \wedge e_{S} \otimes u_{j}=\sum_{j} e_{j+k-1} \wedge e_{S} \otimes u_{j} .
$$

Index a basis of the source by pairs $(S, k)$, with $k \in[\mathbf{n}]$, and the target by $(P, l)$ where $P \subset[2 \mathbf{n}-1]$ has cardinality $\mathbf{n}$ and $l \in[\mathbf{n}]$.

What follows is an ordering of the basis vectors in the target such that the resulting matrix is upper-triangular. Then we just need to show that each diagonal element of the matrix is nonzero to conclude. Unfortunately the order on $(P, l)$ is a little complicated because e.g., if the $l$ 's are ordered sequentially, then to get a diagonal matrix, the $P$ 's must be given an order in the opposite direction.

Define an order relation on the target basis vectors as follows: For $\left(P_{1}, l_{1}\right)$ and $\left(P_{2}, l_{2}\right)$, set $l=\min \left\{l_{1}, l_{2}\right\}$, and declare $\left(P_{1}, l_{1}\right)<\left(P_{2}, l_{2}\right)$ if and only if
(1) In lexicographic order, the set of $l$ minimal elements of $P_{1}$ is strictly after the set of $l$ minimal elements of $P_{2}$ (i.e. the smallest element of $P_{2}$ is smaller than the smallest of $P_{1}$ or they are equal and the second smallest of $P_{2}$ is smaller or equal etc. up to $l$-th), or
(2) the $l$ minimal elements in $P_{1}$ and $P_{2}$ are the same, and $l_{1}<l_{2}$.
(3) the $l$ minimal elements in $P_{1}$ and $P_{2}$ are the same, $l_{1}=l_{2}$, and the set of $\mathbf{n}-l$ tail elements of $P_{1}$ are after the set of $\mathbf{n}-l$ tail elements of $P_{2}$.

The third condition is irrelevant - any breaking of a tie for the first two will lead to an upper-triangular matrix. Note that $(\{\mathbf{n}, \ldots, 2 \mathbf{n}-1\}, 1)$ is the minimal element for this order and $([\mathbf{n}], \mathbf{n})$ is the maximal element. Note further that

$$
e_{\mathbf{n}+1} \wedge \cdots \wedge e_{2 \mathbf{n}-1} \otimes u_{\mathbf{n}} \mapsto e_{\mathbf{n}} \wedge \cdots \wedge e_{2 \mathbf{n}-1} \otimes v_{1}
$$

i.e., that

$$
(\{\mathbf{n}+1, \ldots, 2 \mathbf{n}-1\}, \mathbf{n}) \mapsto(\{\mathbf{n}, \ldots, 2 \mathbf{n}-1\}, 1),
$$

so $(\{\mathbf{n}+1, \ldots, 2 \mathbf{n}-1\}, \mathbf{n})$ will be our first basis element for the source. The order for the source is implicitly described in the proof.

Work by induction: the base case that $(\{\mathbf{n}, \ldots, 2 \mathbf{n}-1\}, 1)$ is in the image has been established. Let $(P, l)$ be any basis element, and assume all $\left(P^{\prime}, l^{\prime}\right)$ with $\left(P^{\prime}, l^{\prime}\right)<(P, l)$ have been shown to be in the image. Write $P=\left(p_{1}, \ldots, p_{\mathbf{n}}\right)$ with $p_{i}<p_{i+1}$. Consider the image of $\left(P \backslash\left\{p_{l}\right\}, 1+p_{l}-l\right)$ which is

$$
\sum_{j} \phi\left(u^{j} \otimes v_{1+p_{l}-l}\right) \wedge e_{P \backslash\left\{p_{l}\right\}} \otimes u_{j}=\sum_{\left\{j \mid j-l+p_{l} \notin P \backslash\left\{p_{l}\right\}\right\}} e_{p_{l}-l+j} \wedge e_{P \backslash\left\{p_{l}\right\}} \otimes u_{j} .
$$

Taking $j=l$ we see $(P, l)$ is among the summands. If $j<l$, the contribution to the summand is a $\left(P^{\prime}, j\right)$ where the first $j$ terms of $P^{\prime}$ equal the first of $P$, so by condition $(2),\left(P^{\prime}, j\right)<(P, l)$. If $j>l$, the summand is a $\left(P^{\prime \prime}, j\right)$ where the first $l-1$ terms of $P$ and $P^{\prime \prime}$ agree, and the $l$-th terms are respectively $p_{l}$ and $p_{l}-l+j$ so by condition (1) $\left(P^{\prime \prime}, j\right)<(P, l)$.

To illustrate, consider the first seven terms when $\mathbf{n}=3$ :

$$
(345,1),(345,2),(345,3),(245,1),(235,1),(234,1),(245,2),
$$

where the order did not matter for the triple $(245,1),(235,1),(234,1)$. We have

$$
\begin{aligned}
& (45,3) \mapsto(345,1) \\
& (35,2) \mapsto(345,2) \\
& (34,3) \mapsto(345,3) \\
& (45,2) \mapsto(245,1)+(345,2) \\
& (35,2) \mapsto(235,1)+(345,3) \\
& (34,2) \mapsto(234,1) \\
& (25,3) \mapsto(245,2) .
\end{aligned}
$$

Second proof. For this proof take $\mathbf{u}=\mathbf{n} \leq \mathbf{v}=\mathbf{m}$. Take a vector space $E$ of dimension 2, and fix isomorphisms $U \simeq S^{\mathbf{n}-1} E, V \simeq S^{\mathbf{m}-1} E^{*}$. Let $A^{\prime}=S^{\mathbf{m}+\mathbf{n}-2} E^{*} \subset S^{\mathbf{n}-1} E^{*} \otimes S^{\mathbf{m}-1} E^{*}=U \otimes V^{*}$, and set $\tilde{A}=A / A^{\prime \perp}$. This turns out to be the same projection operator as in the previous proof.

Our map is

$$
\begin{aligned}
& \Lambda^{\mathbf{n}-1}\left(S^{\mathbf{m}+\mathbf{n}-2} E\right) \otimes S^{\mathbf{n}-1} E \rightarrow \Lambda^{\mathbf{n}}\left(S^{\mathbf{m}+\mathbf{n}-2} E\right) \otimes S^{\mathbf{m}-1} E^{*} \\
& \quad Q_{1} \wedge \cdots \wedge Q_{\mathbf{n}-1} \otimes f \mapsto \sum_{j=0}^{\mathbf{m}-1}\left(f h^{j}\right) \wedge Q_{1} \wedge \cdots \wedge Q_{\mathbf{n}-1} \otimes h_{j}
\end{aligned}
$$

where $h^{j}=x^{j} y^{\mathbf{m}-j-1}$ and $h_{j}$ is the dual basis vector.
Recall the contraction map from Exercise 2.3.2.4, for $\alpha \geq \beta$ :

$$
\begin{aligned}
S^{\alpha} E \times S^{\beta} E^{*} & \rightarrow S^{\alpha-\beta} E \\
(f, g) & \mapsto g\lrcorner f .
\end{aligned}
$$

In the case $f=l^{\alpha}$ for some $l \in E$, then $\left.g\right\lrcorner l^{\alpha}=g(l) l^{\alpha-\beta}$ (here $g(l)$ denotes $g$, considered as a polynomial, evaluated at the point $l$ ), so that $g\lrcorner l^{\alpha}=0$ if and only if $l$ is a root of $g$.

Consider the transposed map, and relabeling $E$ as $E^{*}$ (they are isomorphic as $S L(E) \simeq S L_{2}$ modules):

$$
\begin{aligned}
& \left(\left(\left.M_{\langle 1, \mathbf{m}, \mathbf{n}\rangle}\right|_{A^{\prime} \otimes U^{*} \otimes V^{*}}\right)_{\tilde{A}}^{\wedge p}\right)^{T}: \\
& S^{\mathbf{m}-1} E^{*} \otimes \Lambda^{\mathbf{n}} S^{\mathbf{m}+\mathbf{n}-2} E \rightarrow S^{\mathbf{n}-1} E \otimes \Lambda^{\mathbf{n}-1} S^{\mathbf{m}+\mathbf{n}-2} E \\
& \left.g \otimes\left(f_{1} \wedge \cdots \wedge f_{\mathbf{n}}\right) \mapsto \sum_{i=1}^{\mathbf{n}}(-1)^{i-1}(g\lrcorner f_{i}\right) \otimes f_{1} \wedge \cdots \hat{f}_{i} \cdots \wedge f_{\mathbf{n}} .
\end{aligned}
$$

The map $\left(\left(\left.M_{\langle 1, \mathbf{m}, \mathbf{n}\rangle}\right|_{A^{\prime} \otimes U^{*} \otimes V^{*}}\right)_{\tilde{A}}^{\wedge p}\right)^{T}$ is surjective: Let $l^{\mathbf{n}-1} \otimes\left(l_{1}^{\mathbf{m}+\mathbf{n}-2} \wedge\right.$ $\left.\cdots \wedge l_{\mathbf{n}-1}^{\mathbf{m}+\mathbf{n}-2}\right) \in S^{\mathbf{n}-1} E \otimes \Lambda^{\mathbf{n}-1} S^{\mathbf{m}+\mathbf{n}-2} E$ with $l, l_{i} \in E$. Such elements span the target so it will be sufficient to show any such element is in the image. Assume first that $l$ is distinct from the $l_{i}$. Since $\mathbf{n} \leq \mathbf{m}$, there is a polynomial $g \in S^{\mathbf{m}-1} E^{*}$ which vanishes on $l_{1}, \ldots, l_{\mathbf{n}-1}$ and is nonzero on $l$. Then, up to a nonzero scalar, $g \otimes\left(l_{1}^{\mathbf{m}+\mathbf{n}-2} \wedge \cdots \wedge l_{\mathbf{n}-1}^{\mathbf{m}+\mathbf{n}-2} \wedge l^{\mathbf{m}+\mathbf{n}-2}\right)$ maps to our element.

The condition that $l$ is distinct from the $l_{i}$ may be removed by taking limits, as the image of a linear map is closed.

In $\S 2.6 .2$ we will need the following extension:
Proposition 2.5.2.7. For $2 p<\mathbf{n}-1$ there exist $A^{\prime} \subset U \otimes V^{*}$ of dimension $2 p+1$ such that, setting $\tilde{A}=A /\left(A^{\prime}\right)^{\perp}$,

$$
\left(\left.M_{\langle\mathbf{n}, \mathbf{n}, 1\rangle}\right|_{A^{\prime} \otimes V \otimes U^{*}}\right)_{\tilde{A}}^{\wedge p}: V \otimes \Lambda^{p} \tilde{A} \rightarrow \Lambda^{p+1} \tilde{A} \otimes U
$$

is injective. A general choice of $A^{\prime}$ will have this property.

Proof. Consider $A^{\prime}$ as a subspace of $S^{2 \mathbf{n}-2} E \subset A^{*}$ as in the proof above. Take $A^{\prime}$ spanned by $\ell_{1}^{2 \mathbf{n}-2-\alpha} m_{1}^{\alpha}, \ldots, \ell_{2 p+1}^{2 \mathbf{n}-2-\alpha} m_{2 p+1}^{\alpha}$, where all the $4 p+2$ points $\ell_{k}, m_{j}$ are in general position, and $\alpha<\mathbf{n}-1$ will be chosen below. I show the transposed map is surjective. The target of the transposed map is spanned by vectors of the form $h \otimes \ell_{s_{1}}^{2 \mathbf{n}-2-\alpha} m_{s_{1}}^{\alpha} \wedge \cdots \wedge \ell_{s_{p}}^{2 \mathbf{n}-2-\alpha} m_{s_{p}}^{\alpha}$ where $\left\{s_{1}, \ldots, s_{p}\right\}=S \subset[2 p+1]$. The kernel of the map $\left(\ell_{s_{i}}^{2 \mathbf{n}-2-\alpha} m_{s_{i}}^{\alpha}\right)_{\mathbf{n}-1, \mathbf{n}-1}:$ $S^{\mathbf{n - 1}} E^{*} \rightarrow S^{\mathbf{n}-1} E$ has dimension $\mathbf{n}-\alpha-1$. Since the points were chosen in general linear position, the intersection of the $p$ kernels will have codimension $p(\alpha+1)$. In order to imitate the proof above, we need this intersection to be non-empty, so require $p(\alpha+1)<\mathbf{n}$. Now consider some $\left(\ell_{j}^{2 \mathbf{n}-2-\alpha} m_{j}^{\alpha}\right)_{\mathbf{n}-1, \mathbf{n}-1}$ for $j \notin S$ restricted to the intersection of the kernels. Again since the points were chosen in general linear position, it will be injective, so its image will have dimension $\mathbf{n}-p(\alpha+1)$. We have $p+1$ such maps, and again by general position arguments, the images will be transverse. Thus as long as $(p+1)(\mathbf{n}-p(\alpha+1)) \geq \mathbf{n}$, the span of these $p+1$ images will be all of $S^{\mathbf{n}} E$. Thanks to the hypothesis on $p$, the three inequalities on $\alpha$ are compatible, and we can select any $\alpha$ in the admissible range. Thus every $h \otimes \ell_{s_{1}}^{2 \mathrm{n}-2-\alpha} m_{s_{1}}^{\alpha} \wedge$ $\cdots \wedge \ell_{s_{1}}^{2 n-2-\alpha} m_{s_{1}}^{\alpha}$ will be the image under $\left(\left.M_{\langle\mathbf{n}, \mathbf{n}, 1\rangle}\right|_{A^{\prime} \otimes V \otimes U^{*}}\right)_{\tilde{A}}^{\wedge p}$ of

$$
\sum_{j \notin S} g_{j} \otimes \ell_{j}^{2 \mathbf{n}-2-\alpha} m_{j}^{\alpha} \wedge \ell_{s_{1}}^{2 \mathbf{n}-2-\alpha} m_{s_{1}}^{\alpha} \wedge \cdots \wedge \ell_{s_{1}}^{2 \mathbf{n}-2-\alpha} m_{s_{1}}^{\alpha}
$$

for some $g_{j} \in S^{\mathbf{n}-1} E^{*}$.

Write $\tilde{A}=A / A^{\prime \perp}$. Define

$$
\begin{align*}
P_{2 p+1}: G\left(2 p+1, A^{*}\right) & \rightarrow \mathbb{C}  \tag{2.5.4}\\
A^{\prime} & \mapsto \operatorname{det}\left(\left(\left.M_{\langle\mathbf{n}, \mathbf{n}, \mathbf{m}\rangle}\right|_{A^{\prime} \otimes B^{*} \otimes C^{*}}\right)_{\tilde{A}}^{\wedge p}: \Lambda^{p} \tilde{A} \otimes B^{*} \rightarrow \Lambda^{p+1} \tilde{A} \otimes C\right)
\end{align*}
$$

The above argument shows that $P_{2 p+1}$ is not identically zero for all $2 p \leq$ $\mathbf{n}-1$, but since it is a polynomial, it is not zero on a general $A^{\prime}$.
2.5.3. Why didn't we get a better bound? The above result begs the question: did we fail to get a better bound because this is the best bound Koszul flattenings can give, or is there something pathological about matrix multiplication that prevented the full power of Koszul flattenings? That is, perhaps the Koszul flattenings for $\mathbb{C}^{\mathbf{m}} \otimes \mathbb{C}^{\mathbf{m}} \otimes \mathbb{C}^{\mathbf{m}}$ could be trivial beyond border rank $2 \mathbf{m}-\sqrt{\mathbf{m}}$. This is not the case:
Theorem 2.5.3.1. [Lan15b] The maximal minors of the Koszul flattening $T_{A}^{\wedge p}: \Lambda^{p} \mathbb{C}^{2 p+1} \otimes\left(\mathbb{C}^{2 p+2}\right)^{*} \rightarrow \Lambda^{p+1} \mathbb{C}^{2 p+1} \otimes \mathbb{C}^{2 p+2}$ give nontrivial equations for $\hat{\sigma}_{r} \subset \mathbb{C}^{2 p+1} \otimes \mathbb{C}^{2 p+2} \otimes \mathbb{C}^{2 p+2}$, the tensors of border rank at most $r$ in $\mathbb{C}^{2 p+1} \otimes \mathbb{C}^{2 p+2} \otimes \mathbb{C}^{2 p+2}$, up to $r=4 p+1$.

For $\mathbb{C}^{\mathbf{m}} \otimes \mathbb{C}^{\mathbf{m}} \otimes \mathbb{C}^{\mathbf{m}}$, this implies that when $\mathbf{m}$ is even (resp. odd), the equations are nontrivial up to $r=2 \mathbf{m}-3$ (resp. $r=2 \mathbf{m}-5$ ).
Exercise 2.5.3.2: (2!) Prove the theorem. ©

### 2.6. Lower bounds for the rank of matrix multiplication

2.6.1. The results. Most tensors have rank equal to border rank, in the sense that the set of tensors of rank greater than $r$ in $\hat{\sigma}_{r}$ is a proper subvariety, in particular, a set of measure zero in $\hat{\sigma}_{r}$. I expect matrix multiplication to have larger rank than border rank when $\mathbf{n}>2$ because of its enormous symmetry group, as explained in Chapter 4.

The key to the rank lower bound is that our proof of the border rank lower bound used equations of relatively low degree because of the factorization $\left(M_{\langle\mathbf{n}\rangle}\right)_{A}^{\wedge p}=\left(M_{\langle\mathbf{n}, \mathbf{n}, 1\rangle}\right)_{A}^{\wedge p} \otimes \mathrm{Id}_{W}$, so we were considering minors of a $\operatorname{size}\binom{2 \mathbf{n}-1}{\mathbf{n}} \mathbf{n}$ matrix instead of a size $\binom{2 \mathbf{n}-1}{\mathbf{n}} \mathbf{n}^{2}$ matrix. I will show that if a low degree polynomial is nonzero on $M_{\langle\mathbf{n}\rangle}$, and $M_{\langle\mathbf{n}\rangle}$ has an optimal rank decomposition $M_{\langle\mathbf{n}\rangle}=\sum_{j=1}^{r} a_{j} \otimes b_{j} \otimes c_{j}$, then the polynomial is already zero on a subset of the summands. This is a variant of the substitution method discussed in §5.3.
Theorem 2.6.1.1. [MR13] Let $p \leq \mathbf{n}$ be a natural number. Then

$$
\begin{equation*}
\mathbf{R}\left(M_{\mathbf{n}, \mathbf{n}, \mathbf{m}}\right) \geq\left(1+\frac{p}{p+1}\right) \mathbf{n m}+\mathbf{n}^{2}-\left(2\binom{2 p}{p+1}-\binom{2 p-2}{p-1}+2\right) \mathbf{n} \tag{2.6.1}
\end{equation*}
$$

When $\mathbf{n}=\mathbf{m}$,

$$
\begin{equation*}
\mathbf{R}\left(M_{\langle\mathbf{n}\rangle}\right) \geq\left(3-\frac{1}{p+1}\right) \mathbf{n}^{2}-\left(2\binom{2 p}{p+1}-\binom{2 p-2}{p-1}+2\right) \mathbf{n} . \tag{2.6.2}
\end{equation*}
$$

For example, when $p=1$ one recovers Bläser's bound of $\frac{5}{2} \mathbf{n}^{2}-3 \mathbf{n}$. When $p=3$, the bound (2.6.2) becomes $\frac{11}{4} \mathbf{n}^{2}-26 \mathbf{n}$, which improves Bläser's for $\mathbf{n} \geq 132$. A modification of the method also yields $\mathbf{R}\left(M_{\langle\mathbf{n}\rangle}\right) \geq \frac{8}{3} \mathbf{n}^{2}-7 \mathbf{n}$. See [MR13, Lan14] for proofs of the modifications of the error terms.

I give a proof of a $3 \mathbf{n}^{2}-o\left(\mathbf{n}^{2}\right)$ lower bound for $\mathbf{R}\left(M_{\langle\mathbf{n}\rangle}\right)$ :
Theorem 2.6.1.2. [Lan14] Let $2 p<\mathbf{n}-1$. Then

$$
\mathbf{R}\left(M_{\langle\mathbf{n}, \mathbf{n}, \mathbf{m}\rangle}\right) \geq \frac{2 p+1}{p+1} \mathbf{n m}+\mathbf{n}^{2}-(2 p+1)\binom{2 p+1}{p} \mathbf{n} .
$$

To see this implies $\mathbf{R}\left(M_{\langle\mathbf{n}\rangle}\right) \geq 3 \mathbf{n}^{2}-o\left(\mathbf{n}^{2}\right)$, take $p=\log (\log (\mathbf{n}))$.
2.6.2. Proof of Theorem 2.6.1.2. We will need a few facts from algebraic geometry before the proof.

The following standard Lemma, also used in [Blä03], appears in this form in [Lan12, Lemma 11.5.0.2]:
Lemma 2.6.2.1. Given a polynomial $P$ of degree $d$ on $\mathbb{C}^{\mathbf{a}}$, there exists a subset of basis vectors $\left\{e_{i_{1}}, \ldots, e_{i_{d}}\right\}$ such that $\left.P\right|_{\left\langle e_{i_{1}}, \ldots, e_{i_{d}}\right\rangle}$ is not identically zero.

In other words, there exists a coordinate subspace $\mathbb{C}^{d} \subset \mathbb{C}^{\mathbf{a}}$ such that $\mathbb{C}^{d} \not \subset \mathrm{Zeros}(P)$.

The lemma follows by simply choosing the basis vectors from a degree $d$ monomial that appears in $P$. For example, Lemma 2.6.2.1 implies that a surface in $\mathbb{P}^{3}$ defined by a degree two equation cannot contain six lines whose pairwise intersections span $\mathbb{P}^{3}$.

Recall the Grassmannian $G(k, A)$ from Definition 2.3.3.1.
Lemma 2.6.2.2. Let $A$ be given a basis. For $k, d$ satisfying $d k<\operatorname{dim} A$ and a nonzero homogeneous polynomial $P$ of degree $d$ on $\Lambda^{k} A$ that is not in $I(G(k, A))$, there exist $d k$ basis vectors of $A$ such that, denoting their $d k$-dimensional span by $\tilde{A}, P$ restricted to $G(k, \tilde{A})$ is not identically zero.

Proof. Consider the map $f: A^{\times k} \rightarrow \hat{G}(k, A)$ given by $\left(a_{1}, \ldots, a_{k}\right) \mapsto a_{1} \wedge$ $\cdots \wedge a_{k}$. Then $f$ is surjective. Take the polynomial $P$ and pull it back by $f$. Here the pullback $f^{*}(P)$ is defined by $f^{*}(P)\left(a_{1}, \ldots, a_{k}\right):=P\left(f\left(a_{1}, \ldots, a_{k}\right)\right)$. The pullback is of degree $d$ in each copy of $A$. (I.e., fixing $k-1$ of the $a_{j}$, it becomes a degree $d$ polynomial in the $k$-th.) Now apply Lemma 2.6.2.1 $k$ times to obtain $d k$ basis vectors such that the pulled back polynomial is not
identically zero restricted to their span $\tilde{A}$, and thus $P$ restricted to $\hat{G}(k, \tilde{A})$ is not identically zero.

Remark 2.6.2.3. The bound in Lemma 2.6.2.2 is sharp, as give $A$ a basis $a_{1}, \ldots, a_{\mathbf{a}}$ and consider the polynomial on $\Lambda^{k} A$ with coordinates $x^{I}=$ $x^{i_{1}} \cdots x^{i_{k}}$ corresponding to the vector $\sum_{I} x^{I} a_{i_{1}} \wedge \cdots \wedge a_{i_{k}}$ :

$$
P=x^{1, \ldots, k} x^{k+1, \ldots, 2 k} \cdots x^{(d-1) k+1, \ldots, d k}
$$

Then $P$ restricted to $G\left(k,\left\langle a_{1}, \ldots, a_{d k}\right\rangle\right)$ is non-vanishing but there is no smaller subspace spanned by basis vectors on which it is non-vanishing.

Proof of Theorem 2.6.1.2. Say $\mathbf{R}\left(M_{\langle\mathbf{n}, \mathbf{n}, \mathbf{m}\rangle}\right)=r$ and write an optimal expression

$$
\begin{equation*}
M_{\langle\mathbf{n}, \mathbf{n}, \mathbf{m}\rangle}=\sum_{j=1}^{r} a_{j} \otimes b_{j} \otimes c_{j} . \tag{2.6.3}
\end{equation*}
$$

I will show that the Koszul-flattening equation is already nonzero restricted to a subset of this expression for a judicious choice of $\tilde{A} \subset A$ of dimension $2 p+1$ with $p<\mathbf{n}-1$. Then the rank will be at least the border rank bound plus the number of terms not in the subset. Here are the details:

Recall the polynomial $P_{2 p+1}$ from (2.5.4). It is a polynomial of degree $\binom{2 p+1}{p} \mathbf{n m}>\mathbf{n m}$, so at first sight, e.g., when $\mathbf{m} \sim \mathbf{n}$, Lemma 2.6.2.2 will be of no help because $d k>\operatorname{dim} A=\mathbf{n}^{2}$, but since

$$
\left(\left.M_{\langle\mathbf{n}, \mathbf{n}, \mathbf{m}\rangle}\right|_{A^{\prime} \otimes B^{*} \otimes C^{*}}\right)_{\tilde{A}}^{\wedge p}=\left(\left.M_{\langle\mathbf{n}, \mathbf{n}, 1\rangle}\right|_{A^{\prime} \otimes V \otimes U^{*}}\right)_{\tilde{A}}^{\wedge p} \otimes \operatorname{Id}_{W^{*}},
$$

we actually have $P=\tilde{P}^{\mathrm{m}}$, where

$$
\begin{aligned}
\tilde{P}: G(2 p+1, A) & \rightarrow \mathbb{C} \\
\tilde{A} & \mapsto \operatorname{det}\left(\left(\left.M_{\langle\mathbf{n}, \mathbf{n}, 1\rangle}\right|_{A^{\prime} \otimes V \otimes U^{*}}\right)_{\tilde{A}}^{\wedge p}: \Lambda^{p} \tilde{A} \otimes V \rightarrow \Lambda^{p+1} \tilde{A} \otimes U\right) .
\end{aligned}
$$

Hence we may work with $\tilde{P}$ which is of degree $\binom{2 p+1}{p} \mathbf{n}$ which will be less than $\mathbf{n}^{2}$ if $p$ is sufficiently small. Since $\left(M_{\langle\mathbf{n}, \mathbf{n}, \mathbf{m}\rangle}\right)_{A}: A^{*} \rightarrow B \otimes C$ is injective, some subset of the $a_{j}$ forms a basis of $A$. Lemma 2.6.2.2. implies that there exists a subset of those basis vectors of size $d k=\binom{2 p+1}{p} \mathbf{n}(2 p+1)$, such that if we restrict to terms of the expression (2.6.3) that use only $a_{j}$ whose expansion in the fixed basis has nonzero terms from that subset of $d k$ basis vectors, calling the sum of these terms $M^{\prime}$, we have $\underline{\mathbf{R}}\left(M^{\prime}\right) \geq \frac{2 p+1}{p+1} \mathbf{n m}$. Let $M^{\prime \prime}$ be the sum of the remaining terms in the expression. There are at least $\mathbf{a}-d k=\mathbf{n}^{2}-\binom{2 p+1}{p} \mathbf{n}(2 p+1)$ of the $a_{j}$ appearing in $M^{\prime \prime}$ (the terms corresponding to the complementary basis vectors). Since we assumed we
had an optimal expression for $M_{\langle\mathbf{n}, \mathbf{n}, \mathbf{m}\rangle}$, we have

$$
\begin{aligned}
\mathbf{R}\left(M_{\langle\mathbf{n}, \mathbf{n}, \mathbf{m}\rangle}\right) & =\mathbf{R}\left(M^{\prime}\right)+\mathbf{R}\left(M^{\prime \prime}\right) \\
& \geq \frac{2 p+1}{p+1} \mathbf{n m}+\left[\mathbf{n}^{2}-(2 p+1)\binom{2 p+1}{p} \mathbf{n}\right] .
\end{aligned}
$$

The further lower bounds are obtained by lowering the degree of the polynomial by localizing the equations. An easy such localization is to set $X_{0}=$ Id which reduces the determinant of (2.4.7) to that of $(2.2 .1)$ when $p=$ 2 and yields a similar reduction of degree in general. Further localizations both reduce the degree and the size of the Grassmannian, both of which improve the error term.

## Chapter 3

## The complexity of Matrix Multiplication II: asymptotic upper bounds

This chapter discusses progress towards the astounding conjecture that asymptotically, the complexity of multiplying two $\mathbf{n} \times \mathbf{n}$ matrices is nearly the same as the complexity of adding them. I cover the main advances in upper bounds for the exponent of matrix multiplication beyond Strassen's original discovery in 1969: the 1979 upper bound $\omega<2.78$ of Bini et. al., the 1981 bound $\omega \leq 2.55$ of Schönhage, the 1987 bound $\omega<2.48$ of Strassen, and the Coppersmith-Winograd 1990 bound $\omega<2.38$, emphasizing a geometric perspective. I mention recent "explanations" as to why progress essentially stopped in 1990 from [AFLG15]. In Chapter 4, I discuss other potential paths for upper bounds, and present Pan's $1978 \omega<2.79$ [Pan78], which was the first bound to beat Strassen's and is still (along with its slight modifications) the only decomposition other than Strassen's to be implementable in practice.

The exponent $\omega$ of matrix multiplication is naturally defined in terms of tensor rank:

$$
\omega:=\inf \left\{\tau \in \mathbb{R} \mid \mathbf{R}\left(M_{\langle\mathbf{n}\rangle}\right)=O\left(\mathbf{n}^{\tau}\right)\right\} .
$$

See [BCS97, §15.1] for a the proof that tensor rank yields the same exponent as other complexity measures.

The above-mentioned conjecture is that $\omega=2$. One does not need to work asymptotically to get upper bounds on $\omega$ : Proposition 3.2.1.1 states that for all $\mathbf{n}, \mathbf{R}\left(M_{\langle\mathbf{n}\rangle}\right) \geq \mathbf{n}^{\omega}$. The only methods for proving upper bounds on $\mathbf{R}\left(M_{\langle\mathbf{n}\rangle}\right)$ for any fixed $\mathbf{n}$ that have been used effectively are to find explicit rank decompositions, and very few of these are known.

As I explain in $\S 3.2$, Bini et. al. showed that one may also define the exponent in terms of border rank, namely (see Proposition 3.2.1.10)

$$
\omega=\inf \left\{\tau \in \mathbb{R} \mid \underline{\mathbf{R}}\left(M_{\langle\mathbf{n}\rangle}\right)=O\left(\mathbf{n}^{\tau}\right)\right\} .
$$

Again, we do not need to work asymptotically to get upper bounds on $\omega$ using border rank. Theorem 3.2.1.10 states that for all $\mathbf{n}, \underline{\mathbf{R}}\left(M_{\langle\mathbf{n}\rangle}\right) \geq \mathbf{n}^{\omega}$. In order to make the transition from rank to border rank, we will need a basic result in algebraic geometry. Because of this, I begin, in $\S 3.1$ with some basic facts from the subject. However, the only methods for proving upper bounds on $\underline{\mathbf{R}}\left(M_{\langle\mathbf{n}\rangle}\right)$ for any fixed $\mathbf{n}$ that have been used effectively are to find explicit border rank decompositions, and very few of these are known.

A small help is that we may also use rectangular matrix multiplication to prove upper bounds on $\omega$ : Proposition 3.2.1.10 states that for all $\mathbf{l}, \mathbf{m}, \mathbf{n}$,

$$
\underline{\mathbf{R}}\left(M_{\langle\mathbf{m}, \mathbf{n}, \mathbf{l}\rangle}\right) \geq(\mathbf{l m n})^{\frac{\omega}{3}} .
$$

But again, our knowledge of border rank is scant.
To improve the situation, one needs techniques that enable one to avoid dealing with tensors beyond the small range we have results in. After the work of Bini et. al., all upper bounds on $\omega$ are obtained via tensors other than $M_{\langle 1, \mathbf{m}, \mathbf{n}\rangle}$.

The next advance in upper bounds, due to Schönhage (Theorem 3.3.3.1) and described in $\S 3.3$, is more involved: it says it is sufficient to prove upper bounds on sums of disjoint matrix multiplications.

To go beyond this, Strassen had the idea to looks for a tensor $T \in$ $A \otimes B \otimes C$, that has special combinatorial structure rendering it easy to study, that can be degenerated into a collection of disjoint matrix multiplications.

The inequalities regarding $\omega$ above are strict, e.g., there does not exist $\mathbf{n}$ with $\underline{\mathbf{R}}\left(M_{\langle\mathbf{n}\rangle}\right)$ equal to $\mathbf{n}^{\omega}$. (This does not rule out $\underline{\mathbf{R}}\left(M_{\langle\mathbf{n}\rangle}\right)$ equal to $2 \mathbf{n}^{\omega}$ for all n.) Strassen looked for sequences $T_{N} \in A_{N} \otimes B_{N} \otimes C_{N}$ that could be degenerated into sums $\bigoplus_{i=1}^{s(N)} M_{\left\langle\mathbf{1}_{i}(N), \mathbf{m}_{i}(N) \mathbf{n}_{i}(N)\right\rangle}$ with the border rank of the sums giving upper bounds on $\omega$. This is Strassen's "laser method" described in §3.4.

More precisely, to obtain a sequence of disjoint matrix multiplication tensors, one takes a base tensor $T$ and degenerates the tensor powers $T^{\otimes N} \in$
$\left(A^{\otimes N}\right) \otimes\left(B^{\otimes N}\right) \otimes\left(C^{\otimes N}\right)$. Strassen's degeneration is in the sense of points in the $G L\left(A^{\otimes N}\right) \times G L\left(B^{\otimes N}\right) \times G L\left(C^{\otimes N}\right)$-orbit closure of $T^{\otimes N}$.

After Strassen, all other subsequent upper bounds on $\omega$ use what I will call combinatorial restrictions of $T^{\otimes N}$ for some "simple" tensor $T$, where entries of a coordinate presentation of $T^{\otimes N}$ are just set equal to zero. The choice of entries to zero out is subtle. I describe these developments in $\S 3.4$.

In addition to combinatorial restrictions, Cohn et. al. exploit a geometric change of basis when a tensor is the multiplication tensor of an algebra (or even more general structures). They use the discrete Fourier transform for finite groups (and more general structures) to show that the multiplication tensor in the Fourier basis (and thus in any basis) has "low" rank, but nevertheless in the standard basis admits a combinatorial restriction to a "large" sum of matrix multiplication tensors. I discuss this approach in §3.5.

The proofs in this chapter make essential use of the property from Exercise 2.5.1.9:

$$
\begin{equation*}
M_{\langle\mathbf{l}, \mathbf{m}, \mathbf{n}\rangle} \otimes M_{\left\langle\mathbf{l}^{\prime}, \mathbf{m}^{\prime}, \mathbf{n}^{\prime}\right\rangle}=M_{\left\langle\mathbf{l}^{\prime}, \mathbf{m m}^{\prime}, \mathbf{\mathbf { n } ^ { \prime }}\right\rangle} \tag{3.0.1}
\end{equation*}
$$

where for tensors $T \in A \otimes B \otimes C$ and $T^{\prime} \in A^{\prime} \otimes B^{\prime} \otimes C^{\prime}, T \otimes T^{\prime}$ is considered as a tensor in the triple tensor product $\left(A \otimes A^{\prime}\right) \otimes\left(B \otimes B^{\prime}\right) \otimes\left(C \otimes C^{\prime}\right)$.

### 3.1. Facts and definitions from algebraic geometry

Standard references for this material are [Har95, Mum95, Sha13a]. The first is very good for examples, the second and third have clean proofs, with the proofs in the second more concise.

Several results from this section will be used repeatedly in this book: that the linear projection of a projective variety is a projective variety (Theorem 3.1.4.1), that projective varieties of complementary dimension must intersect (Theorem 3.1.5.1), and that the Zariski and Euclidean closures of certain sets agree (Theorem 3.1.6.1).
3.1.1. Projective varieties. Varieties in a vector space $V$ defined by homogeneous polynomials are invariant under rescaling. For this, and other reasons, it will be convenient to work in projective space (Definition 2.3.3.1). Write $\pi: V \backslash 0 \rightarrow \mathbb{P} V$ for the projection map. For $X \subset \mathbb{P} V$, write $\pi^{-1}(X) \cup$ $\{0\}=: \hat{X} \subset V$, and $\pi(y)=[y]$. If $\hat{X} \subset V$ is a variety, I will also refer to $X \subset \mathbb{P} V$ as a variety. The zero set in $V$ of a collection of polynomials on $V$ is called an affine variety and the image in $\mathbb{P} V$ of the zero set of a collection of homogeneous polynomials on $V$ is called a projective variety. For subsets $Z \subset V, \mathbb{P} Z \subset \mathbb{P} V$ denotes its image under $\pi$. If $P \in S^{d} V^{*}$
is an irreducible polynomial, then its zero set $\operatorname{Zeros}(P) \subset \mathbb{P} V$ is an irreducible variety, called a hypersurface of degree $d$. For a variety $X \subset \mathbb{P} V$, $I_{d}(X):=\left\{P \in S^{d} V^{*} \mid X \subset \operatorname{Zeros}(P)\right\}$ denotes the ideal of $X$ in degree $d$, and $I(X)=\oplus_{d} I_{d}(X) \subset \operatorname{Sym}\left(V^{*}\right)$ is the ideal of $X$.

We will be mostly concerned with varieties in spaces of tensors (for the study of matrix multiplication) and spaces of polynomials (for geometric complexity theory).

### 3.1.2. Examples of varieties.

(1) Projective space $\mathbb{P} V \subseteq \mathbb{P} V$.
(2) The Segre variety of rank one tensors

$$
\begin{aligned}
& \sigma_{1}=\operatorname{Seg}\left(\mathbb{P} A_{1} \times \cdots \times \mathbb{P} A_{n}\right) \\
& \quad:=\mathbb{P}\left\{T \in A_{1} \otimes \cdots \otimes A_{n} \mid \exists a_{j} \in A_{j} \text { such that } T=a_{1} \otimes \cdots \otimes a_{n}\right\} \subset \mathbb{P}\left(A_{1} \otimes \cdots \otimes A_{n}\right) .
\end{aligned}
$$

(3) The Veronese variety

$$
v_{d}(\mathbb{P} V)=\mathbb{P}\left\{P \in S^{d} V \mid P=x^{d} \text { for some } x \in V\right\} \subset \mathbb{P} S^{d} V .
$$

(4) The Grassmannian
$G(k, V):=\mathbb{P}\left\{T \in \Lambda^{k} V \mid \exists v_{1}, \ldots, v_{k} \in V\right.$ such that $\left.T=v_{1} \wedge \cdots \wedge v_{k}\right\} \subset \mathbb{P} \Lambda^{k} V$.
(5) The Chow variety

$$
C h_{d}(V):=\mathbb{P} \overline{\left\{P \in S^{d} V \mid \exists v_{1}, \ldots, v_{d} \in V \text { such that } P=v_{1} \cdots v_{d}\right\}} \subset \mathbb{P} S^{d} V
$$

By definition, projective space is a variety (the zero set of no equations).
Exercise 3.1.2.1: (2) Show that $\operatorname{Seg}\left(\mathbb{P} A_{1} \times \cdots \times \mathbb{P} A_{n}\right)$ is the zero set of the size two minors of the flattenings $A_{j}^{*} \rightarrow A_{1} \otimes \cdots \otimes \hat{A}_{j} \otimes \cdots \otimes A_{n}$, for $1 \leq$ $j \leq n$.

To get equations for $v_{d}(\mathbb{P} V)$, given $P \in S^{d} V$, consider the flattening $P_{1, d-1}: V^{*} \rightarrow S^{d-1} V$ defined by $\frac{\partial}{\partial v} \mapsto \frac{\partial P}{\partial v}$. For example when $d=4, \mathbf{v}=2$ and $P=\sum_{i=0}^{4} p_{i} x^{i} y^{4-i}$, the matrix representing $P_{1,3}$ is

$$
\left(\begin{array}{llll}
p_{4} & p_{3} & p_{2} & p_{1}  \tag{3.1.1}\\
p_{3} & p_{2} & p_{1} & p_{0}
\end{array}\right)
$$

and $v_{4}\left(\mathbb{P}^{1}\right)$ is the zero set of the 6 size two minors of this matrix.
Exercise 3.1.2.2: (1) Show that $v_{d}(\mathbb{P} V)$ is the zero set of the size two minors of the flattening $V^{*} \rightarrow S^{d-1} V$.

We saw equations for the Grassmannian in §2.6.2.
Exercise 3.1.4.2 will show that it is not necessary to take the Zariski closure when defining the Chow variety. Equations for the Chow variety are
known, see $\S 9.6$. However generators of the ideal of the Chow variety are not known explicitly- what is known is presented in Chapter 9.
3.1.3. Dimension via tangent spaces. Informally, the dimension of a variety is the number of parameters needed to describe it locally. For example, the dimension of $\mathbb{P} V$ is $\mathbf{v}-1$ because in coordinates on the open neighborhood where $x_{1} \neq 0$, points of $\mathbb{P} V$ have a unique expression as $\left[1, x_{2}, \ldots, x_{\mathbf{v}}\right]$, where $x_{2}, \ldots, x_{\mathbf{v}}$ are free parameters.

I first define dimension of a variety via dimensions of vector spaces. Define the affine tangent space to $X \subset \mathbb{P} V$ at $[x] \in X, \hat{T}_{x} \hat{X}=\hat{T}_{[x]} X \subset V$, to be the span of the tangent vectors $x^{\prime}(0)$ to analytic curves $x(t)$ on $\hat{X}$ with $x(0)=x$, and note that this is independent of the choice of (nonzero) $x \in[x]$. A point $x \in \hat{X}$ is defined to be a smooth point if $\operatorname{dim} \hat{T}_{y} \hat{X}$ is constant for all $y$ in some neighborhood of $x$.

The dimension of an irreducible variety $\hat{X} \subset V$ is the dimension of the tangent space at a smooth point of $\hat{X}$. If $x$ is a smooth point, $\operatorname{dim} X=$ $\operatorname{dim} \hat{X}-1=\operatorname{dim} \hat{T}_{x} \hat{X}-1$. If $x$ is not a smooth point, it is called a singular point and we let $X_{\text {sing }} \subset X$ denote the singular points of $X$. A variety of dimension one is called a curve.
Remark 3.1.3.1. The above definitions of smooth points and dimension implicitly assumes that $X$ is a reduced variety. A hypersurface $\{P=0\}$ is reduced if when one decomposes $P$ into irreducible factors $P=p_{1}^{a_{1}} \cdots p_{r}^{a_{r}}$, that all $a_{j}=1$. For example $\left\{\ell^{n-m}\right.$ perm $\left._{m}=0\right\}$ is not reduced when $n-m>1$. The definition of dimension in $\S 3.1 .5$ below avoids this problem. For a definition of singular points that avoids this problem, see §6.3.1.

Exercise 3.1.3.2: (2) Show that $\operatorname{dim}\left\{\operatorname{det}_{n}=0\right\}_{\text {sing }}=n^{2}-4$.
If a Zariski open subset of a variety is given parametrically, then one can calculate the tangent space to the variety via the parameter space. For example $\hat{S} e g(\mathbb{P} A \times \mathbb{P} B \times \mathbb{P} C)$ may be thought of as the image of the map

$$
\begin{aligned}
A \times B \times C & \rightarrow A \otimes B \otimes C \\
(a, b, c) & \mapsto a \otimes b \otimes c,
\end{aligned}
$$

so to compute $\hat{T}_{[a \otimes b \otimes c]} S e g(\mathbb{P} A \times \mathbb{P} B \times \mathbb{P} C)$, take curves $a(t) \subset A$ with $a(0)=a$ and similarly for $B, C$, then $\left.\frac{d}{d t}\right|_{t=0} a(t) \otimes b(t) \otimes c(t)=a^{\prime} \otimes b \otimes c+$ $a \otimes b^{\prime} \otimes c+a \otimes b \otimes c^{\prime}$ by the Leibnitz rule. Since $a^{\prime}$ can be any vector in $A$ and similarly for $b^{\prime}, c^{\prime}$ we conclude

$$
\hat{T}_{[a \otimes b \otimes c]} S e g(\mathbb{P} A \times \mathbb{P} B \times \mathbb{P} C)=A \otimes b \otimes c+a \otimes B \otimes c+a \otimes b \otimes C .
$$

The right hand side spans a space of dimension $\mathbf{a}+\mathbf{b}+\mathbf{c}-2$, so $\operatorname{dim}(\operatorname{Seg}(\mathbb{P} A \times$ $\mathbb{P} B \times \mathbb{P} C))=\mathbf{a}+\mathbf{b}+\mathbf{c}-3$.

I can now pay off two debts: in $\S 2.1 .1$, I asserted that the fundamental Theorem of linear algebra is something of a miracle, and in Theorem 2.1.5.1 I asserted that a general tensor in $\mathbb{C}^{\mathbf{m}} \otimes \mathbb{C}^{\mathrm{m}} \otimes \mathbb{C}^{\mathrm{m}}$ has tensor rank around $\frac{\mathrm{m}^{2}}{3}$.

A general point of $\sigma_{2}$ is of the form $\left[a_{1} \otimes b_{1} \otimes c_{1}+a_{2} \otimes b_{2} \otimes c_{2}\right]$, and a general tangent vector at that point is of the form $a_{1} \otimes b_{1} \otimes c_{1}^{\prime}+a_{1} \otimes b_{1}^{\prime} \otimes c_{1}+$ $a_{1}^{\prime} \otimes b_{1} \otimes c_{1}+a_{2} \otimes b_{2} \otimes c_{2}^{\prime}+a_{2} \otimes b_{2}^{\prime} \otimes c_{2}+a_{2}^{\prime} \otimes b_{2} \otimes c_{2}$, hence

$$
\begin{aligned}
& \hat{T}_{\left[a_{1} \otimes b_{1} \otimes c_{1}+a_{2} \otimes b_{2} \otimes c_{2}\right]} \sigma_{2}= \\
& a_{1} \otimes b_{1} \otimes C+a_{1} \otimes B \otimes c_{1}+A \otimes b_{1} \otimes c_{1}+a_{2} \otimes b_{2} \otimes C+a_{2} \otimes B \otimes c_{2}+A \otimes b_{2} \otimes c_{2}
\end{aligned}
$$

so that $\operatorname{dim} \sigma_{2} \leq 2(\operatorname{dim}(\operatorname{Seg}(\mathbb{P} A \times \mathbb{P} B \times \mathbb{P} C))+2-1$ (and equality clearly holds if a, $\mathbf{b}, \mathbf{c} \geq 3)$ and similarly $\operatorname{dim} \sigma_{r} \leq r(\operatorname{dim}(S e g(\mathbb{P} A \times \mathbb{P} B \times \mathbb{P} C))+r-1$. The first chance this has to be the entire ambient space is when this number is $\mathbf{a b c}-1$. When $\mathbf{a}=\mathbf{b}=\mathbf{c}=\mathbf{m}$, this means $r \geq \frac{\mathbf{m}^{3}}{3 \mathbf{m}-2}$, paying the second debt.

For the first,

$$
\begin{aligned}
\hat{T}_{\left[a_{1} \otimes b_{1}+a_{2} \otimes b_{2}\right]} \sigma_{2, A \otimes B} & =\operatorname{span}\left\{a_{1} \otimes b_{1}^{\prime}+a_{1}^{\prime} \otimes b_{1}+a_{2} \otimes b_{2}^{\prime}+a_{2}^{\prime} \otimes b_{2}\right\} \\
& =A \otimes \operatorname{span}\left\{b_{1}, b_{2}\right\}+\operatorname{span}\left\{a_{1}, a_{2}\right\} \otimes B
\end{aligned}
$$

and this space has dimension $2 \operatorname{dim} \operatorname{Seg}(\mathbb{P} A \times \mathbb{P} B)$, instead of the expected $2 \operatorname{dim} \operatorname{Seg}(\mathbb{P} A \times \mathbb{P} B)+1$. This accounts for the upper semi-continuity of matrix rank which fails for tensor rank: any point on a tangent line, i.e., a point of the form $a^{\prime} \otimes b+a \otimes b^{\prime}$ is also transparently on a secant line, i.e., the sum of two rank one matrices.
Exercise 3.1.3.3: (1) Compute $\hat{T}_{\left[x^{d}\right]} v_{d}(\mathbb{P} V)$.
3.1.4. Noether normalization. Consider the curve $\{x y=1\} \subset \mathbb{C}^{2}$ :


If we project the curve onto the $x$-axis, we get the set $\{x \in \mathbb{C} \mid x \neq 0\}$, which, as was discussed in §1.1.14, is not Zariski closed.

One of the many wonderful things about projective space is that the projection of an algebraic variety to a hyperplane is still an algebraic variety.


I remind the reader that unless mentioned otherwise, I work exclusively over the complex numbers, because the next theorem is false over $\mathbb{R}$ :
Theorem 3.1.4.1. If $X \subset \mathbb{P} W$ is a variety, $L \subset W$ is a subspace with $\mathbb{P} L \cap X=\emptyset$, and one considers the projection map $p: W \rightarrow W / L$, then $\mathbb{P} p(\hat{X}) \subset \mathbb{P}(W / L)$ is also a variety.

Theorem 3.1.4.1 is part of the Noether normalization theorem (see, e.g., [Sha13a, §1.5.4] or [Mum95, $\S 2 \mathrm{C}]$ ). It is proved via elimination theory. In addition to failing in affine space, this projection property fails over $\mathbb{R}$ : the curve in $\mathbb{R P}^{2}$ given by $x^{2}+z^{2}-y^{2}=0$ when projected from $[1,0,0]$ is not a real algebraic variety. (It consists of $\mathbb{R P}^{1} \backslash\{[0,1]\}$.)
Exercise 3.1.4.2: (1) Show that if $W=V^{\otimes d}$ and $L$ is the $G L(V)$-complement to $S^{d} V$ in $V^{\otimes d}$, taking $p: V^{\otimes d} \rightarrow V^{\otimes d} / L \simeq S^{d} V$, then $p(\operatorname{Seg}(\mathbb{P} V \times \cdots \times$ $\mathbb{P} V))=C h_{d}(V)$. Conclude the closure is not needed in the definition of the Chow variety. ©

The ideal of the projection of a variety from a coordinate point is obtained by eliminating that coordinate from the equations in the ideal. For example, give $S^{4} \mathbb{C}^{2}$ coordinates $\left(p_{4}, p_{3}, p_{2}, p_{1}, p_{0}\right)$ as above and project from $p_{2}$. Eliminating $p_{2}$ from the equations

$$
p_{4} p_{2}-p_{3}^{2}, p_{4} p_{1}-p_{2} p_{3}, p_{4} p_{0}-p_{1} p_{3}, p_{3} p_{1}-p_{2}^{2}, p_{2} p_{0}-p_{1}^{2}
$$

gives the ideal generated by

$$
p_{4} p_{0}-p_{1} p_{3}, p_{3}^{3}-p_{4}^{2} p_{1}, p_{1}^{3}-p_{0}^{2} p_{3} .
$$

Exercise 3.1.4.3: (2) What equations does one get when projecting from $p_{3}$ ? Give a geometric explanation why the answer is different. (A complete answer to this question is beyond what we have covered, I am just asking for some equations.) ©

Remark 3.1.4.4. Since elimination theory doesn't care which point one projects from, one can even project from a point on a variety. The resulting "map" is not defined at the point one projects from, but the Zariski closure of the image of the points where it is defined at is well defined. This is an example of a rational map.
Exercise 3.1.4.5: (2) What ideal does one get when projecting $v_{4}\left(\mathbb{P}^{1}\right)$ from $p_{4}$ ? (A complete answer to this question is beyond what we have covered, I am just asking for some equations.) ©

As long as $X$ does not surject onto $\mathbb{P} V / L$, we can continue projecting it to smaller and smaller projective spaces.

If $X \subset \mathbb{P} V$ is a projective variety and $f: X \rightarrow Y \subset \mathbb{P}^{N}$ is given by $N+1$ homogeneous polynomials on $V$, then $f$ is an example of a regular map. If $X \subset \mathbb{C}^{M}$ and $Y \subset \mathbb{C}^{N}$ are affine varieties, a regular map $f: X \rightarrow Y$ is one given by $N$ polynomials $p_{1}, \ldots, p_{N}$ on $\mathbb{C}^{M}$, such that $\left(p_{1}(x), \ldots, p_{N}(x)\right) \in Y$ for all $x \in X$. For the definition of a regular map, see see, e.g. [Sha13a, $\S 1.2 .3]$. If $X \subset \mathbb{C}^{N}$ is an affine variety, $\mathbb{C}[X]:=\mathbb{C}\left[x_{1}, \ldots, x_{N}\right] / I(X)$ denotes the space of regular functions on $X$.
Exercise 3.1.4.6: (1) If $X, Y$ are affine varieties and $f: X \rightarrow Y$ is a regular map, show that one gets a map $f^{*}: \mathbb{C}[Y] \rightarrow \mathbb{C}[X]$, called the induced pullback map, and that $f^{*}$ is injective if $f$ is surjective.

Theorem 3.1.4.1 generalizes to:
Theorem 3.1.4.7. (see, e.g., [Sha13a, §5.2, Thm. 1.10]) If $X$ is a projective variety and $f: X \rightarrow Y$ is a regular map, then $f(X)$ is Zariski closed.
Exercise 3.1.4.8: (1) Show that if $X$ is irreducible and $f: X \rightarrow Y$ is regular, then $f(X)$ is irreducible. ©
3.1.5. Dimension via projection. The dimension of $X \subset \mathbb{P} V$ is also the largest integer $n$ such that there exists a surjective linear projection onto a $\mathbb{P}^{n}$. In this case the surjective projection $X \rightarrow \mathbb{P}\left(V / \mathbb{C}^{c}\right)$ may be chosen to be finite to one. The integer $c=\mathbf{v}-1-n$ is called the codimension of $X$ in $\mathbb{P} V$. Noether normalization implies that a general linear space $\mathbb{P} L$ will satisfy $\operatorname{dim}(X \cap \mathbb{P} L)=\mathbf{v}-1-n-\operatorname{dim} \mathbb{P} L$. In particular, the intersection of $X$ with a general linear space of dimension $c+1$ will be a finite number of points. This number of points is called the degree of $X$.

A consequence of this more algebraic definition of dimension is the following result:
Theorem 3.1.5.1. Let $X, Y \subset \mathbb{P}^{N}$ (resp. $X, Y \subset \mathbb{C}^{N}$ ) be irreducible projective (resp. affine) varieties.

Then any non-empty component $Z$ of $X \cap Y$ has $\operatorname{dim} Z \geq \operatorname{dim} X+$ $\operatorname{dim} Y-N$.

Moreover, in the projective case, if $\operatorname{dim} X+\operatorname{dim} Y-N>0$, then $X \cap Y \neq$ $\emptyset$.

For the proof, see, e.g., [Sha13a, §1.6.4].
3.1.6. Zariski and Euclidean closure. Recall from §1.1.14.2 that the Zariski closure of a set can be larger than the Euclidean closure. Nevertheless, the following theorem, proved using Noether normalization, shows that in our situation, the two closures agree:
Theorem 3.1.6.1. Let $Z \subset \mathbb{P} V$ be a subset. Then the Euclidean closure of $Z$ is contained in the Zariski closure of $Z$. If $Z$ contains a Zariski open subset
of its Zariski closure, then the two closures coincide. The same assertions hold for subsets $Z \subset V$.

A proof that uses nothing but Noether normalization is given in [Mum95, Thm. 2.33]. I present a proof using the following basic fact: for every irreducible algebraic curve $C \subset \mathbb{P} V$ there exists a smooth algebraic curve $\tilde{C}$ and a surjective algebraic map $\pi: \tilde{C} \rightarrow C$ that is one-to-one over the smooth points of $C$. (More precisely, $\pi$ is a finite map as defined in §9.5.1.) See, e.g., [Sha13a, $\S 1.2 .5 .3]$ for a proof. The curve $\tilde{C}$ is called the normalization of $C$.

The theorem will follow immediately from the following Lemma:
Lemma 3.1.6.2. Let $Z \subset \mathbb{P} V$ be an irreducible variety and let $Z^{0} \subset Z$ be a Zariski open subset. Let $p \in Z \backslash Z^{0}$. Then there exists an analytic curve $C(t)$ such that $C(t) \in Z^{0}$ for all $t \neq 0$ and $\lim _{t \rightarrow 0} C(t)=p$.


Proof. Let $c$ be the codimension of $Z$ and take a general linear space $\mathbb{P} L \subset$ $\mathbb{P} V$ of dimension $c+1$ that contains $p$. Then $\mathbb{P} L \cap Z$ will be a possibly reducible algebraic curve containing $p$. Take a component $C$ of the curve that contains $p$. If $p$ is a smooth point of the curve we are done, as we can expand a Taylor series about $p$. Otherwise take the the normalization $\pi: \tilde{C} \rightarrow C$ and a point of $\pi^{-1}(p)$, expand a Taylor series about that point and compose with $\pi$ to obtain the desired analytic curve.

### 3.2. The upper bounds of Bini, Capovani, Lotti, and Romani

### 3.2.1. Rank, border rank, and the exponent of matrix multiplica-

 tion.Proposition 3.2.1.1. [Bin80] For all $\mathbf{n}, \mathbf{R}\left(M_{\langle\mathbf{n}\rangle}\right) \geq \mathbf{n}^{\omega}$, i.e., $\omega \leq \frac{\log \mathbf{R}\left(M_{\langle\mathbf{n}\rangle}\right)}{\log (\mathbf{n})}$.

Proof. By the definitions of the exponent and $O$, there exists a constant $C$, such that $C \mathbf{R}\left(M_{\langle\mathbf{n}\rangle}\right) \geq \mathbf{n}^{\omega}$ for all $\mathbf{n}$. By (3.0.1) and Exercise 2.1.6.3, $\mathbf{R}\left(M_{\left\langle\mathbf{n}^{k}\right\rangle}\right) \leq \mathbf{R}\left(M_{\langle\mathbf{n}\rangle}\right)^{k}$. Say $\mathbf{R}\left(M_{\langle\mathbf{n}\rangle}\right)=r$. Then $C r^{k} \geq\left(\mathbf{n}^{k}\right)^{\omega}$, i.e., $C^{\frac{1}{k}} r \geq$ $\mathbf{n}^{\omega}$. Now let $k$ go to infinity, we get $r \geq \mathbf{n}^{\omega}$.
Remark 3.2.1.2. The calculation in the proof of Proposition 3.2.1.1 is typical in the upper bound literature and will show up several times in this chapter: one has an initially hazardous constant (in this case $C$ ) that gets washed out asymptotically by taking high tensor powers of $M_{\langle\mathbf{n}\rangle}$.
Proposition 3.2.1.3. For all $\mathbf{1}, \mathbf{m}, \mathbf{n},(\mathbf{l m n})^{\frac{\omega}{3}} \leq \mathbf{R}\left(M_{\langle\mathbf{m}, \mathbf{n}, 1\rangle}\right)$, i.e., $\omega \leq$ $\frac{3 \log \mathbf{R}\left(M_{\langle\mathbf{m}, \mathbf{n}, 1\rangle}\right)}{\log (\mathbf{m n l})}$.
Exercise 3.2.1.4: (2) Prove Proposition 3.2.1.3. ©
Remark 3.2.1.5. The inequalities in Propositions 3.2.1.1 and 3.2.1.3 are strict, see Theorem 3.3.3.5.

To show that $\omega$ may also be defined in terms of border rank, introduce a sequence of ranks that interpolate between rank and border rank.

We say $\mathbf{R}_{h}(T) \leq r$ if there exists an expression

$$
\begin{equation*}
T=\lim _{\epsilon \rightarrow 0} \frac{1}{\epsilon^{h}}\left(a_{1}(\epsilon) \otimes b_{1}(\epsilon) \otimes c_{1}(\epsilon)+\cdots+a_{r}(\epsilon) \otimes b_{r}(\epsilon) \otimes c_{r}(\epsilon)\right) \tag{3.2.1}
\end{equation*}
$$

where $a_{j}(\epsilon), b_{j}(\epsilon), c_{j}(\epsilon)$ are analytic functions of $\epsilon$.
Proposition 3.2.1.6. $\underline{\mathbf{R}}(T) \leq r$ if and only if there exists an $h$ such that $\mathbf{R}_{h}(T) \leq r$.

Proof. We need to show $\underline{\mathbf{R}}(T) \leq r$ implies there exists an $h$ with $\mathbf{R}_{h}(T) \leq r$. Since $\operatorname{Seg}(\mathbb{P} A \times \mathbb{P} B \times \mathbb{P} C)$ is just the product of three projective spaces, every curve in $\operatorname{Seg}(\mathbb{P} A \times \mathbb{P} B \times \mathbb{P} C)$ is of the form $[a(t) \otimes b(t) \otimes c(t)]$ for some curves $a(t) \subset A$ etc., and if the curve is analytic, the functions $a(t), b(t), c(t)$ can be taken to be analytic as well. Thus every analytic curve in $\sigma_{r}^{0}(S e g(\mathbb{P} A \times \mathbb{P} B \times \mathbb{P} C))$ may be written as [ $\sum_{j=1}^{r} a_{j}(t) \otimes b_{j}(t) \otimes c_{j}(t)$ ] for some analytic curves $a_{j}(t) \subset A$ etc. Since the Euclidean and Zariski closures of $\hat{\sigma}_{r}^{0}$ agree by Theorem 3.1.6.1, we conclude that if $T \in \hat{\sigma}_{r}$, then $\mathbf{R}_{h}(T) \leq r$ for $h$ equal to the order of first nonzero term in the Taylor expansion of $\sum_{j=1}^{r} a_{j}(t) \otimes b_{j}(t) \otimes c_{j}(t)$.
Proposition 3.2.1.7. If $\mathbf{R}_{h}(T) \leq r$, then $\mathbf{R}(T) \leq r\binom{h+2}{2}<r h^{2}$.
Proof. Write $T$ as in (3.2.1). Then $T$ is the coefficient of the $\epsilon^{h}$ term of the expression in parentheses. For each summand, there is a contribution of $\sum_{\alpha+\beta+\gamma=h}\left(\epsilon^{\alpha} a_{\alpha}\right) \otimes\left(\epsilon^{\beta} b_{\beta}\right) \otimes\left(\epsilon^{\gamma} c_{\gamma}\right)$ which consists of $\binom{h+2}{2}$ terms.
Remark 3.2.1.8. In fact $\mathbf{R}(T) \leq r(h+1)$, see Exercise 3.5.3.3.

Exercise 3.2.1.9: (1) Show that for $T \in A \otimes B \otimes C$, if $\mathbf{R}_{h}(T) \leq r$, then $\mathbf{R}_{N h}\left(T^{\otimes N}\right) \leq r^{N}$ where $T^{\otimes N}$ is considered as an element of the triple tensor product $\left(A^{\otimes N}\right) \otimes\left(B^{\otimes N}\right) \otimes\left(C^{\otimes N}\right)$.
Theorem 3.2.1.10. [Bini, $[\operatorname{Bin} 80]]$ For all $\mathbf{l}, \mathbf{m}, \mathbf{n}, \omega \leq \frac{3 \log \mathbf{R}\left(M_{\langle\mathbf{m}, \mathbf{n}, \mathbf{l})}\right)}{\log (\mathbf{m n l})}$.
Proof. Write $r=\underline{\mathbf{R}}\left(M_{\langle\mathbf{m}, \mathbf{n}, 1\rangle}\right)$. Set $N=\mathbf{m n l}$. We have $\mathbf{R}_{h}\left(M_{\langle N\rangle}\right) \leq r^{3}$ for some $h$ and thus $\mathbf{R}\left(M_{\left\langle N^{k}\right\rangle}\right) \leq r^{3 k}(h k)^{2}$, which implies

$$
\left(N^{k}\right)^{\omega} \leq r^{3 k}(h k)^{2},
$$

so

$$
N^{\omega} \leq r^{3}(h k)^{\frac{2}{k}} .
$$

Letting $k \rightarrow \infty$ gives the result.
3.2.2. Bini et. al's algorithm. Recall from $\S 2.1 .4$ that $\underline{\mathbf{R}}\left(M_{\langle 2\rangle}^{\text {red }}\right) \leq 5$.

Exercise 3.2.2.1: (1) Use that $\underline{\mathbf{R}}\left(M_{\langle 2\rangle}^{\text {red }}\right) \leq 5$ to show $\underline{\mathbf{R}}\left(M_{\langle 2,2,3\rangle}\right) \leq 10$. More generally, show that if $\underline{\mathbf{R}}\left(M_{\langle\mathbf{m}, 2,2\rangle}^{\text {red }}\right)=r$ and $\underline{\mathbf{R}}\left(M_{\left\langle\mathbf{m}^{\prime}, 2,2\right\rangle}^{r e d}\right)=r^{\prime}$, then setting $n=m+m^{\prime}-1, \underline{\mathbf{R}}\left(M_{\langle n, 2,2\rangle}\right) \leq r+r^{\prime}$.©

Using Proposition 3.2.1.10 we conclude:
Theorem 3.2.2.2. [BCRL79] $\omega<2.78$.

### 3.3. Schönhage's upper bounds

The next contribution to upper bounds for the exponent of matrix multiplication was Schönhage's discovery that the border rank of the sum of two tensors in disjoint spaces can be smaller than the sum of the border ranks, and that this failure could be exploited to prove further upper bounds on the exponent. This result enables one to prove upper bounds with tensors that are easier to analyze because of their low border rank. Before giving Schönhage's bounds, I begin with geometric preliminaries on orbit closures.
3.3.1. Orbit closures. Orbit closures will play a central role in our study of GCT. They also play a role in the work of Schönhage and Strassen on matrix multiplication.

When $r \leq \mathbf{a}_{i}$ for $1 \leq i \leq n, \sigma_{r}\left(\operatorname{Seg}\left(\mathbb{P} A_{1} \times \cdots \times \mathbb{P} A_{n}\right)\right)$ is an orbit closure: Let $a_{j}^{\alpha_{j}}, 1 \leq \alpha_{j} \leq \mathbf{a}_{j}$, be a basis of $A_{j}$, then

$$
\begin{aligned}
& \sigma_{r}\left(S \operatorname{Seg}\left(\mathbb{P} A_{1} \times \cdots \times \mathbb{P} A_{n}\right)\right) \\
& =\overline{G L\left(A_{1}\right) \times \cdots \times G L\left(A_{n}\right) \cdot\left[a_{1}^{1} \otimes \cdots \otimes a_{n}^{1}+\cdots+a_{1}^{r} \otimes \cdots \otimes a_{n}^{r}\right]} \subset \mathbb{P}\left(A_{1} \otimes \cdots \otimes A_{n}\right) .
\end{aligned}
$$

Write $M_{\langle 1\rangle}^{\oplus r}=\sum_{j=1}^{r} a_{j} \otimes b_{j} \otimes c_{j} \in \mathbb{C}^{r} \otimes \mathbb{C}^{r} \otimes \mathbb{C}^{r}$ where $\left\{a_{j}\right\},\left\{b_{j}\right\},\left\{c_{j}\right\}$ are bases. This tensor is sometimes called the unit tensor. Then

$$
\begin{equation*}
\sigma_{r}\left(S e g\left(\mathbb{P}^{r-1} \times \mathbb{P}^{r-1} \times \mathbb{P}^{r-1}\right)\right)=\overline{G L_{r} \times G L_{r} \times G L_{r} \cdot\left[M_{\langle 1\rangle}^{\oplus r}\right]} . \tag{3.3.1}
\end{equation*}
$$

Exercise 3.3.1.1: (2) Let $V$ be a $G$-module and let $v, w \in V$. Show that $w \in \overline{G \cdot v}$ if and only if $\overline{G \cdot w} \subseteq \overline{G \cdot v}$.
Proposition 3.3.1.2. If $T^{\prime} \in \overline{G L(A) \times G L(B) \times G L(C) \cdot T} \subset A \otimes B \otimes C$, then $\underline{\mathbf{R}}\left(T^{\prime}\right) \leq \underline{\mathbf{R}}(T)$.
Exercise 3.3.1.3: (1) Prove Proposition 3.3.1.2. ©
Definition 3.3.1.4. If $T^{\prime} \in \overline{G L(A) \times G L(B) \times G L(C) \cdot T} \subset A \otimes B \otimes C$, we say $T^{\prime}$ is a degeneration of $T$.

Consider the orbit closure of the matrix multiplication tensor

$$
\overline{G L(A) \times G L(B) \times G L(C) \cdot\left[M_{\langle U, V, W\rangle}\right]} \subset \mathbb{P}(A \otimes B \otimes C)
$$

By Exercise 3.3.1.1, we may rephrase our characterization of border rank as, taking inclusions $A, B, C \subset \mathbb{C}^{r}$,

$$
\begin{aligned}
\underline{\mathbf{R}}\left(M_{\langle\mathbf{n}\rangle}\right) \leq r & \Leftrightarrow\left[M_{\langle\mathbf{n}\rangle}\right] \in \sigma_{r}(\operatorname{Seg}(\mathbb{P} A \times \mathbb{P} B \times \mathbb{P} C)) \\
& \Leftrightarrow \overline{G L_{r} \times G L_{r} \times G L_{r} \cdot\left[M_{\langle\mathbf{n}\rangle}\right]} \subset \sigma_{r}\left(\operatorname{Seg}\left(\mathbb{P}^{r-1} \times \mathbb{P}^{r-1} \times \mathbb{P}^{r-1}\right)\right) \\
& \Leftrightarrow \overline{G L_{r} \times G L_{r} \times G L_{r} \cdot\left[M_{\langle\mathbf{n}\rangle}\right]} \subset \overline{G L_{r} \times G L_{r} \times G L_{r} \cdot\left[M_{\langle 1\rangle}^{\oplus r}\right]} .
\end{aligned}
$$

3.3.2. Schönhage's example. Recall from Exercise 2.1.7.6 that $\underline{\mathbf{R}}\left(M_{\langle 1, \mathbf{m}, \mathbf{n}\rangle}\right)=$ $\mathbf{m n}$ and $\underline{\mathbf{R}}\left(M_{\langle N, 1,1\rangle}\right)=N$. Recall the notation from $\S 2.1 .6$ that if $T_{1} \in$ $A_{1} \otimes B_{1} \otimes C_{1}$ and $T_{2} \in A_{2} \otimes B_{2} \otimes C_{2}$, we define the tensor $T_{1} \oplus T_{2} \in\left(A_{1} \oplus\right.$ $\left.A_{2}\right) \otimes\left(B_{1} \oplus B_{2}\right) \otimes\left(C_{1} \oplus C_{2}\right)$. (In Exercise 5.3.1.6 you will show that $\mathbf{R}\left(M_{\langle 1, \mathbf{m}, \mathbf{n}\rangle} \oplus\right.$ $\left.\left.M_{\langle N, 1,1\rangle}\right)=\mathbf{m n}+N.\right)$
Theorem 3.3.2.1 (Schönhage $[\mathbf{S c h} 81])$. Set $N=(\mathbf{n}-1)(\mathbf{m}-1)$. Then

$$
\underline{\mathbf{R}}\left(M_{\langle 1, \mathbf{m}, \mathbf{n}\rangle} \oplus M_{\langle N, 1,1\rangle}\right)=\mathbf{m n}+1
$$

Proof. By conciseness, we only need to show $\underline{\mathbf{R}}\left(M_{\langle 1, \mathbf{m}, \mathbf{n}\rangle} \oplus M_{\langle N, 1,1\rangle}\right) \leq$ $\mathbf{m n}+1$. Write

$$
\begin{aligned}
& M_{\langle 1, \mathbf{m}, \mathbf{n}\rangle}=\sum_{i=1}^{\mathbf{m}} \sum_{j=1}^{\mathbf{n}} x_{i} \otimes y_{j} \otimes z_{i, j}, \\
& M_{\langle N, 1,1\rangle}=\sum_{u=1}^{\mathbf{m}-1} \sum_{v=1}^{\mathbf{n}-1} x_{u, v} \otimes y_{u, v} \otimes z .
\end{aligned}
$$

Then

$$
\begin{aligned}
M_{\langle 1, \mathbf{m}, \mathbf{n}\rangle} \oplus M_{\langle N, 1,1\rangle}=\lim _{t \rightarrow 0} \frac{1}{t^{2}} & {\left[\sum_{u=1}^{\mathbf{m}-1} \sum_{v=1}^{\mathbf{n}-1}\left(x_{u}+t x_{u v}\right) \otimes\left(y_{v}+t y_{u v}\right) \otimes\left(z+t^{2} z_{u v}\right)\right.} \\
& +\sum_{u=1}^{\mathbf{m}-1} x_{u} \otimes\left(y_{\mathbf{n}}+t\left(-\sum_{v} y_{u v}\right)\right) \otimes\left(z+t^{2} z_{u \mathbf{n}}\right) \\
& +\sum_{v=1}^{\mathbf{n}-1}\left(x_{\mathbf{m}}+t\left(-\sum_{u} x_{u v}\right)\right) \otimes y_{v} \otimes\left(z+t^{2} z_{\mathbf{m} v}\right) \\
& \left.+x_{\mathbf{m}} \otimes y_{\mathbf{n}} \otimes\left(z+t^{2} z_{\mathbf{m n}}\right)-\left(\sum_{i} x_{i}\right) \otimes\left(\sum_{s} y_{s}\right) \otimes z\right] .
\end{aligned}
$$

For a discussion of the geometry of this limit, see [Lan12, §11.2.2].
3.3.3. Schönhage's asymptotic sum inequality. To develop intuition how an upper bound on a sum of matrix multiplications could give an upper bound on a single matrix multiplication, say we knew $\underline{\mathbf{R}}\left(M_{\langle\mathbf{n}\rangle}^{\oplus s}\right) \leq r$ with $s \leq \mathbf{n}^{3}$. Then to compute $M_{\left\langle\mathbf{n}^{2}\right\rangle}$ we could write $M_{\left\langle\mathbf{n}^{2}\right\rangle}=M_{\langle\mathbf{n}\rangle} \otimes M_{\langle\mathbf{n}\rangle}$. At worst this is evaluating $\mathbf{n}^{3}$ disjoint copies of $M_{\langle\mathbf{n}\rangle}$. Now group these $\mathbf{n}^{3}$ disjoint copies in groups of $s$ and apply the bound to obtain a savings.

Here is the precise statement:
Theorem 3.3.3.1. [Sch81] [Schönhage's asymptotic sum inequality] For all $\mathbf{l}_{i}, \mathbf{m}_{i}, \mathbf{n}_{i}$, with $1 \leq i \leq s$ :

$$
\sum_{i=1}^{s}\left(\mathbf{m}_{i} \mathbf{n}_{i} \mathbf{l}_{i}\right)^{\frac{\omega}{3}} \leq \underline{\mathbf{R}}\left(\bigoplus_{i=1}^{s} M_{\left\langle\mathbf{m}_{i}, \mathbf{n}_{i}, \mathbf{l}_{i}\right\rangle}\right) .
$$

The main step of the proof, and an outline of the rest of the argument is given below.

Remark 3.3.3.2. A similar result (also proven in [Sch81]) holds for the border rank of the multiplication of matrices with some entries equal to zero, where the product $\mathbf{m}_{i} \mathbf{n}_{i} \mathbf{l}_{i}$ is replaced by the number of multiplications in the naïve algorithm for the matrices with zeros.

Here is a special case that isolates the new ingredient (following [Blä13]):

## Lemma 3.3.3.3.

$$
\mathbf{n}^{\omega} \leq\left\lceil\frac{\underline{\mathbf{R}}\left(M_{\langle\mathbf{n}\rangle}^{\oplus s}\right)}{s}\right\rceil .
$$

In particular, $s \mathbf{n}^{\omega} \leq \underline{\mathbf{R}}\left(M_{\langle\mathbf{n}\rangle}^{\oplus s}\right)$.

Proof. Let $r=\underline{\mathbf{R}}\left(M_{\langle\mathbf{n}\rangle}^{\oplus s}\right)$. It is sufficient to show that for all $N$,

$$
\begin{equation*}
\underline{\mathbf{R}}\left(M_{\left\langle\mathbf{n}^{N}\right\rangle}^{\oplus s}\right) \leq\left\lceil\frac{r}{s}\right\rceil^{N} s \tag{3.3.2}
\end{equation*}
$$

as then, since trivially $\underline{\mathbf{R}}\left(M_{\left\langle\mathbf{n}^{N}\right\rangle}^{\oplus s}\right) \geq \underline{\mathbf{R}}\left(M_{\left\langle\mathbf{n}^{N}\right\rangle}\right) \geq\left(\mathbf{n}^{N}\right)^{\omega}$, we have

$$
\left(\mathbf{n}^{N}\right)^{\omega} \leq\left\lceil\frac{r}{s}\right\rceil^{N} s
$$

i.e.,

$$
\mathbf{n}^{\omega} \leq\left\lceil\frac{r}{s}\right\rceil s^{\frac{1}{N}}
$$

and the result follows letting $N \rightarrow \infty$.
I prove (3.3.2) by induction on $N$. The hypothesis is the case $N=1$. Assume (3.3.2) holds up to $N$ and observe that

$$
M_{\left\langle\mathbf{n}^{N+1}\right\rangle}^{\oplus s}=M_{\langle\mathbf{n}\rangle}^{\oplus s} \otimes M_{\left\langle\mathbf{n}^{N}\right\rangle} .
$$

Now $\underline{\mathbf{R}}\left(M_{\langle\mathbf{n}\rangle}^{\oplus s}\right) \leq r$ implies $M_{\langle\mathbf{n}\rangle}^{\oplus s} \in \overline{G L_{r}^{\times 3} \cdot M_{\langle 1\rangle}^{\oplus r}}$ by Equation (3.3.1), so $M_{\langle\mathbf{n}\rangle}^{\oplus s} \otimes M_{\left\langle\mathbf{n}^{N}\right\rangle} \in \overline{G L_{r}^{\times 3} \cdot M_{\langle 1\rangle}^{\oplus r} \otimes M_{\left\langle\mathbf{n}^{N}\right\rangle}}$. Thus $\underline{\mathbf{R}}\left(M_{\left\langle\mathbf{n}^{N+1}\right.}^{\oplus s}\right) \leq \underline{\mathbf{R}}\left(M_{\langle 1\rangle}^{\oplus r} \otimes M_{\left\langle\mathbf{n}^{N}\right\rangle}\right)$. Recall that $M_{\langle 1\rangle}^{\oplus t} \otimes M_{\left\langle\mathbf{n}^{N}\right\rangle}=M_{\left\langle\mathbf{n}^{N}\right\rangle}^{\oplus t}$. Now

$$
\begin{aligned}
\underline{\mathbf{R}}\left(M_{\left\langle\mathbf{n}^{N+1}\right\rangle}^{\oplus s}\right) & \leq \underline{\mathbf{R}}\left(M_{\left\langle\mathbf{n}^{N}\right\rangle}^{\oplus r}\right) \\
& \leq \underline{\mathbf{R}}\left(M_{\left\langle\mathbf{n}^{N}\right\rangle}^{\left.\oplus \Gamma \frac{r}{s}\right] s}\right) \\
& \leq \underline{\mathbf{R}}\left(M_{\langle 1\rangle}^{\oplus \Gamma \frac{r}{s}} \otimes M_{\left\langle\mathbf{n}^{N}\right\rangle}^{\oplus s}\right) \\
& \leq \underline{\mathbf{R}}\left(M_{\langle 1\rangle}^{\oplus\left[\frac{r}{s}\right]}\right) \underline{\mathbf{R}}\left(M_{\left\langle\mathbf{n}^{N}\right\rangle}^{\oplus s}\right) \\
& \leq\left\lceil\frac{r}{s}\right\rceil\left(\left\lceil\frac{r}{s}\right\rceil^{N} s\right)
\end{aligned}
$$

where the last inequality follows from the induction hypothesis.
The general case of Theorem 3.3.3.1 essentially follows from the above lemma and arguments used previously: one first takes a high tensor power of the sum, then switches to rank at the price of introducing an $h$ that washes out in the end. The new tensor is a sum of products of matrix multiplications that one converts to a sum of matrix multiplications. One then takes the worst term in the summation and estimates with respect to it (multiplying by the number of terms in the summation), and applies the lemma to conclude.
Corollary 3.3.3.4. [Sch81] $\omega<2.55$.
Proof. Applying Theorem 3.3.3.1 to $\underline{\mathbf{R}}\left(M_{\langle 1, \mathbf{m}, \mathbf{n}\rangle} \oplus M_{\langle(\mathbf{m}-1)(\mathbf{n}-1), 1,1\rangle}\right)=\mathbf{m n}+$ 1 gives

$$
(\mathbf{m n})^{\frac{\omega}{3}}+((\mathbf{m}-1)(\mathbf{n}-1))^{\frac{\omega}{3}} \leq \mathbf{m n}+1
$$

and taking $\mathbf{m}=\mathbf{n}=4$ gives the result.
In [CW82] they prove that for any tensor $T$ that is a direct sum of disjoint matrix multiplications, if $\mathbf{R}(T) \leq r$, then there exists $N$ such that $\underline{\mathbf{R}}\left(T \oplus M_{\langle N, 1,1\rangle}\right) \leq r+1$. This, combined with our earlier arguments using $\mathbf{R}_{h}$ to bridge the gap between rank and border rank asymptotically, implies the inequality in Theorem 3.3.3.1 is strict:
Theorem 3.3.3.5. [CW82] For all $\mathbf{l}_{i}, \mathbf{m}_{i}, \mathbf{n}_{i}$, with $1 \leq i \leq s$ :

$$
\sum_{i=1}^{s}\left(\mathbf{m}_{i} \mathbf{n}_{i} \mathbf{l}_{i}\right)^{\frac{\omega}{3}}<\underline{\mathbf{R}}\left(\bigoplus_{i=1}^{s} M_{\left\langle\mathbf{m}_{i}, \mathbf{n}_{i}, \mathbf{l}_{i}\right\rangle}\right)
$$

In particular, for all $\mathbf{n}, \underline{\mathbf{R}}\left(M_{\langle\mathbf{n}\rangle}\right)>\mathbf{n}^{\omega}$, so one cannot determine $\omega$ from $M_{\langle\mathbf{n}\rangle}$ for any fixed $\mathbf{n}$.

### 3.4. Strassen's laser method

3.4.1. Introduction. Recall our situation: we don't understand rank or even border rank in the range we would need to prove upper bounds on $\omega$ via $M_{\langle\mathbf{n}\rangle}$, so we showed upper bounds on $\omega$ could be proved first with rectangular matrix multiplication, then with sums of disjoint matrix multiplications which had the property that the border rank of the sum was less than the sum of the border ranks, and the border rank in each case was determined via an explicit decomposition.

We also saw that to determine the exponent by such methods, one would need to deal with sequences of tensors. Strassen's laser method is based on taking high tensor powers of a fixed tensor, and then degenerating it to a disjoint sum of matrix multiplication tensors. Because it deals with sequences, there is no known obstruction to determining $\omega$ exactly via Strassen's method.

Starting with Strassen's method, all attempts to determine $\omega$ aim at best for a Pyrrhic victory in the sense that even if $\omega$ were determined by these methods, they would not give any indication as to what would be optimally fast matrix multiplication for any given size matrix.

### 3.4.2. Strassen's tensor. Consider the following tensor

$$
\begin{equation*}
T_{S T R}=\sum_{j=1}^{q} a_{0} \otimes b_{j} \otimes c_{j}+a_{j} \otimes b_{0} \otimes c_{j} \in \mathbb{C}^{q+1} \otimes \mathbb{C}^{q+1} \otimes \mathbb{C}^{q} \tag{3.4.1}
\end{equation*}
$$

Exercise 5.3.1.7 will show that $\mathbf{R}\left(T_{S T R}\right)=2 q$, so (3.4.1) is an optimal rank expression. Nevertheless, $\underline{\mathbf{R}}\left(T_{S T R}\right)=q+1$. To see why one could expect
this, consider the $q$ points $a_{0} \otimes b_{0} \otimes c_{j}$. The tensor $T_{S T R}$ is a sum of tangent vectors to these $q$ points:

$$
T_{S T R}=\sum_{j=1}^{q} \lim _{t \rightarrow 0} \frac{1}{t}\left[\left(a_{0}+t a_{j}\right) \otimes\left(b_{0}+t b_{j}\right) \otimes c_{j}-a_{0} \otimes b_{0} \otimes c_{j}\right]
$$

Note that the sum $\sum_{j} a_{0} \otimes b_{0} \otimes c_{j}$ is also a rank one tensor, which leads one to the expression:

$$
\lim _{t \rightarrow 0} \frac{1}{t}\left[\sum_{j=1}^{q}\left(a_{0}+t a_{j}\right) \otimes\left(b_{0}+t b_{j}\right) \otimes c_{j}-a_{0} \otimes b_{0} \otimes\left(c_{1}+\cdots+c_{q}\right)\right]
$$

showing the border rank is at most $q+1$, but since the tensor is concise, we obtain equality. Geometrically, the original $q$ points all lie on the linear space $\left[a_{0} \otimes b_{0} \otimes \mathbb{C}^{q}\right] \subset S e g(\mathbb{P A} \times \mathbb{P} B \times \mathbb{P} C)$.

Now consider $\tilde{T}_{S T R}:=T_{S T R} \otimes \sigma\left(T_{S T R}\right) \otimes \sigma^{2}\left(T_{S T R}\right)$ where $\sigma$ is a cyclic permutation of the three factors. Group triples of spaces together to consider $\tilde{T}_{S T R} \in \mathbb{C}^{q(q+1)^{2}} \otimes \mathbb{C}^{q(q+1)^{2}} \otimes \mathbb{C}^{q(q+1)^{2}}$. We have the upper bound $\underline{\mathbf{R}}\left(\tilde{T}_{S T R}\right) \leq$ $(q+1)^{3}$.

Write $a_{\alpha \beta \gamma}:=a_{\alpha} \otimes a_{\beta} \otimes a_{\gamma}$ and similarly for $b$ 's and $c$ 's. Then, omitting the $\otimes$ 's:

$$
\begin{align*}
\tilde{T}_{S T R}=\sum_{i, j, k=1}^{q} & \left(a_{i j 0} b_{0 j k} c_{i 0 k}+a_{i j k} b_{0 j k} c_{i 00}+a_{i j 0} b_{00 k} c_{i j k}+a_{i j k} b_{00 k} c_{i j 0}\right.  \tag{3.4.2}\\
& \left.+a_{0 j 0} b_{i j k} c_{i 0 k}+a_{0 j k} b_{i j k} c_{i 00}+a_{0 j 0} b_{i 0 k} c_{i j k}+a_{0 j k} b_{i 0 k} c_{i j 0}\right)
\end{align*}
$$

We may think of $\tilde{T}_{S T R}$ as a sum of eight terms, each of which is a $M_{\langle 1, \mathbf{m}, \mathbf{n}\rangle}$ with $\operatorname{lmn}=q^{3}$, e.g., the first is $\sum_{i, j, k=1}^{q} a_{i j 0} b_{0 j k} c_{i 0 k}=M_{\langle q, q, q\rangle}$, the second $M_{\left\langle q^{2}, q, 1\right\rangle}$ etc.. (I will say terms of volume $q^{3}$.) Were they all disjoint expressions, we could use the asymptotic sum inequality to conclude $8 q^{\omega} \leq(q+1)^{3}$ and for small $q$ we would see $\omega<2$. Of course this is not the case, but we can try to zero out some of the variables to keep as many of these eight terms as possible. For example if we set $c_{i 00}, b_{00 k}, b_{i j k}, c_{i j k}$ all to zero, we are left with two disjoint matrix multiplications and we conclude $2 q^{\omega} \leq(q+1)^{3}$. This is best when $q=15$, giving $\omega<2.816$, which is not so interesting.

At this point enters a new idea: since we are dealing with border rank, we have greater flexibility in degeneration than simply zero-ing out terms. By taking limits, we will be able to keep three terms! To explain this, I need to take another detour regarding orbit closures.

### 3.4.3. All tensors of border rank $\left\lfloor\frac{3}{4} \mathbf{n}^{2}\right\rfloor$ are degenerations $M_{\langle\mathbf{n}\rangle}$.

Theorem 3.4.3.1 (Strassen [Str87]). Set $r=\left\lfloor\frac{3}{4} \mathbf{n}^{2}\right\rfloor$ and choose a linear embedding $\mathbb{C}^{r} \subset \mathbb{C}^{\mathbf{n}^{2}}$. Then

$$
\sigma_{r}\left(S e g\left(\mathbb{P}^{r-1} \times \mathbb{P}^{r-1} \times \mathbb{P}^{r-1}\right)\right) \subset \overline{G L_{\mathbf{n}^{2}} \times G L_{\mathbf{n}^{2}} \times G L_{\mathbf{n}^{2}} \cdot\left[M_{\langle\mathbf{n}\rangle}\right]},
$$

i.e.,

$$
\overline{G L_{r} \times G L_{r} \times G L_{r} \cdot\left[M_{\langle 1\rangle}^{\oplus r}\right]} \subset \overline{G L_{\mathbf{n}^{2}} \times G L_{\mathbf{n}^{2}} \times G L_{\mathbf{n}^{2}} \cdot\left[M_{\langle\mathbf{n}\rangle}\right]} .
$$

Proof. The proof will be by a very simple degeneration: let $T^{A} \subset G L(A)=$ $G L_{\mathbf{n}^{2}}$ denote the diagonal $\mathbf{n}^{2} \times \mathbf{n}^{2}$ matrices. I will show

$$
M_{\langle 1\rangle}^{\oplus r} \subset \overline{T^{A} \times T^{B} \times T^{C} \cdot M_{\langle\mathbf{n}\rangle}} .
$$

Write $x_{i j}$ for a basis of $A$ etc., so $M_{\langle\mathbf{n}\rangle}=\sum_{i, j, k} x_{i j} \otimes y_{j k} \otimes z_{k i}$. We want to kill off as few terms as possible such that in the remaining terms, each basis vector appears in at most one monomial. That is if we have $x_{i j}$ appearing, then there should be a unique $k_{0}=k(i, j)$, such that the only term surviving in $\sum_{k} x_{i j} \otimes y_{j k} \otimes z_{k i}$ is $x_{i j} \otimes y_{j k_{0}} \otimes z_{k_{0} i}$. We should view this more symmetrically, fixing some integer $h$ and requiring that the only terms appearing are of the form $x_{i j} \otimes y_{j k} \otimes z_{k i}$ where $i+j+k=h$. To do this, look for curves

$$
\begin{aligned}
x_{i j} & \mapsto t^{\alpha(i, j)} x_{i j} \\
y_{j k} & \mapsto t^{\beta(j, k)} y_{j k} \\
z_{k i} & \mapsto t^{\gamma(k, i)} z_{k i}
\end{aligned}
$$

so that $\alpha+\beta+\gamma=0$ when $i+j+k=h$ and $\alpha+\beta+\gamma>0$ when $i+j+k \neq h$, as then

$$
\lim _{t \rightarrow 0} \sum_{i, j, k=1}^{\mathbf{n}} t^{\alpha(i, j)+\beta(j, k)+\gamma(k, i)} x_{i j} \otimes y_{j k} \otimes z_{k i}=\sum_{i+j+k=h} x_{i j} \otimes y_{j k} \otimes z_{k i} .
$$

Set $\lambda=i+j+k$. We could satisfy the requirements on $\alpha, \beta, \gamma$ by requiring

$$
\alpha+\beta+\gamma=(h-\lambda)^{2}=h^{2}-2 \lambda h+\lambda^{2} .
$$

Take

$$
\begin{aligned}
\alpha & =\frac{1}{2}\left(i^{2}+j^{2}\right)+2 i j+\left(\frac{h}{3}-i-j\right) h \\
\beta & =\frac{1}{2}\left(k^{2}+j^{2}\right)+2 k j+\left(\frac{h}{3}-k-j\right) h \\
\gamma & =\frac{1}{2}\left(i^{2}+k^{2}\right)+2 i k+\left(\frac{h}{3}-i-k\right) h .
\end{aligned}
$$

Exercise 3.4.3.2: (1) Verify that $\alpha+\beta+\gamma=(h-\lambda)^{2}$.
Exercise 3.4.3.3: (2) Show that the best value of $h$ is $h=\left\lceil\frac{3 \mathbf{n}}{2}\right\rceil+1$ which yields $r=\left\lfloor\frac{3}{4} \mathbf{n}^{2}\right\rfloor$ to finish the proof.

Remark 3.4.3.4. This degeneration is more complicated than setting linear combinations of variables to zero because there are values of $i, j, k$ where one of $\alpha, \beta, \gamma$ is negative. To avoid negative terms for the curves in $A, B, C$, we could add $r$ to each of $\alpha, \beta, \gamma$ and then divide the entire entire expression by $t^{3 r}$.

Call degenerations that only use the diagonal matrices toric degenerations.
Corollary 3.4.3.5. Every tensor in $\mathbb{C}^{\frac{3}{2} \mathbf{n}} \otimes \mathbb{C}^{\frac{3}{2} \mathbf{n}} \otimes \mathbb{C}^{\frac{3}{2} \mathbf{n}}$ arises as a toric degeneration of $M_{\langle\mathbf{n}\rangle}$.

Proof. As mentioned in §2.1.6, the maximum border rank of any tensor in $\mathbb{C}^{\frac{3}{2}} \mathbf{n} \otimes \mathbb{C}^{\frac{3}{2}} \mathbf{n} \otimes \mathbb{C}^{\frac{3}{2}} \mathbf{n}$ is at most $\frac{3}{4} \mathbf{n}^{2}$, and any tensor of border rank $\frac{3}{4} \mathbf{n}^{2}$ is a degeneration of $M_{\langle\mathbf{n}\rangle}$.
Remark 3.4.3.6. Theorem 3.4.3.1 may be interpreted as saying that one can degenerate $M_{\langle\mathbf{n}\rangle}$ to a tensor that computes $\left\lfloor\frac{3}{4} \mathbf{n}^{2}\right\rfloor$ independent scalar multiplications. If we have any tensor realized as $M_{\langle\mathbf{n}\rangle} \otimes T$, the same degeneration procedure works to degenerate it to $M_{\langle 1\rangle}^{\oplus\left\lfloor\frac{3}{4} \mathbf{n}^{2}\right\rfloor} \otimes T$.
3.4.4. A better bound using the toric degeneration. Now we return to the expression (3.4.2). There are four kinds of $A$-indices, $i j 0, i j k, 0 j 0$ and $0 j k$. To emphasize this, and to suggest what kind of degeneration to perform, label these with superscripts [11], [21], [12] and [22]. Label each of the $B$ and $C$ indices (which come in four types as well) similarly to obtain:

$$
\begin{aligned}
\tilde{T}_{S T R}= & \sum_{i, j, k=1}^{q}\left(a_{i j 0}^{[11]} b_{0 j k}^{[11]} c_{i 0 k}^{[11]}+a_{i j k}^{[21]} b_{0 j k}^{[11]} c_{i 00}^{[12]}+a_{i j 0}^{[11]} b_{00 k}^{[12]} c_{i j k}^{[21]}+a_{i j k}^{[21]} b_{00 k}^{[12]} c_{i j 0}^{[22]}\right. \\
& \left.+a_{0 j 0}^{[12]} b_{i j k}^{[21]} c_{i 0 k}^{[11]}+a_{0 j k}^{[22]} b_{i j k}^{[21]} c_{i 00}^{[12]}+a_{0 j 0}^{[12]} b_{i 0 k}^{[22]} c_{i j k}^{[21]}+a_{0 j k}^{[22]} b_{i 0 k}^{[22]} c_{i j 0}^{[22]}\right)
\end{aligned}
$$

This expression has the structure of block $2 \times 2$ matrix multiplication. Think of it as a sum of $q^{3} 2 \times 2$ matrix multiplications. Now use Theorem 3.4.3.1 to degenerate each $2 \times 2$ matrix multiplication to a sum of 3 disjoint terms. Namely, following the recipe that the three indices must add to 4 , we keep all terms $a^{[s, t]} b^{[t, u]} c^{[u, s]}$ where $s+t+u=4$, namely we degenerate $\tilde{T}_{S T R}$ to

$$
\sum_{i, j, k=1}^{q} a_{i j k}^{[21]} b_{0 j k}^{[11]} c_{i 00}^{[12]}+a_{i j 0}^{[11]} b_{00 k}^{[12]} c_{i j k}^{[21]}+a_{0 j 0}^{[12]} b_{i j k}^{[21]} c_{i 0 k}^{[11]} .
$$

The asymptotic sum inequality implies $3 q^{\omega} \leq(q+1)^{3}$, which gives the best bound on $\omega$ when $q=7$, namely $\omega<2.642$, which is still not as good as Schönhage's bound.
3.4.5. Strassen's bound. We do better by using the standard trick of this chapter: taking a high tensor power of $\tilde{T}_{S T R}$, as $\tilde{T}_{S T R}^{\otimes N}$ contains $\left(2^{N}\right)^{2}$ matrix multiplications $M_{\langle\mathbf{l}, \mathbf{m}, \mathbf{n}\rangle}$, all with $\mathbf{l m n}=q^{3 N}$, and again by Theorem 3.4.3.1 we may keep $\frac{3}{4} 2^{2 N}$ of them. The asymptotic sum inequality applied to the degenerated tensor gives

$$
\frac{3}{4} 2^{2 N} q^{N \omega} \leq(q+1)^{3 N}
$$

Taking $N$-th roots and letting $N$ tend to infinity, the $\frac{3}{4}$ goes away and we obtain

$$
2^{2} q^{\omega} \leq(q+1)^{3} .
$$

Finally, the case $q=5$ implies:
Theorem 3.4.5.1. [Str87] $\omega<2.48$.
3.4.6. Asymptotic rank. The above discussion suggests the introduction of yet another complexity measure for tensors: given $T \in A \otimes B \otimes C$, we can consider $T^{\otimes N} \in A^{\otimes N} \otimes B^{\otimes N} \otimes C^{\otimes N}$ and this construction played a central role in Strassen's laser method to prove upper bounds for the complexity of matrix multiplication via auxiliary tensors.

Definition 3.4.6.1. The asymptotic rank $\tilde{\mathbf{R}}(T)$ of a tensor $T \in A \otimes B \otimes C$, is

$$
\tilde{\mathbf{R}}(T):=\inf _{N}\left[\mathbf{R}\left(T^{\otimes N}\right)\right]^{\frac{1}{N}} .
$$

Exercise 3.4.6.2: (1) Show that in the definition, one can replace the infimum by $\lim _{N \rightarrow \infty}$ by using Lemma 3.4.7.2 below.
Exercise 3.4.6.3: (2) Show that $\tilde{\mathbf{R}}(T) \leq \underline{\mathbf{R}}(T)$. ©
Since $M_{\langle 2\rangle}^{\otimes k}=M_{\left\langle 2^{k}\right\rangle}$, we have $\tilde{\mathbf{R}}\left(M_{\langle 2\rangle}\right)=2^{\omega}$.
Conjecture 3.4.6.4. $[\mathbf{S t r} 91]$ Let $T \in \mathbb{C}^{\mathbf{m}} \otimes \mathbb{C}^{\mathbf{m}} \otimes \mathbb{C}^{\mathrm{m}}$ be concise. Then $\tilde{\mathbf{R}}(T)=\mathbf{m}$.

Note that, If Conjecture 3.4.6.4 holds for $T=M_{\langle 2\rangle}$, this would imply $\omega=2$.

More subtly, if the conjecture holds for $T_{c w, 2}$ introduced in $\S 3.4 .9$ below, then $\omega=2$, see [BCS97, Rem. 15.44].
3.4.7. Degeneracy value. I now formalize what we did to get Strassen's bound. The starting point is if a tensor $T$ degenerates to $\bigoplus_{i=1}^{s} M_{\left\langle\mathbf{l}_{i}, \mathbf{m}_{i}, \mathbf{n}_{i}\right\rangle}$, then $\sum_{i=1}^{s}\left(\mathbf{l}_{i} \mathbf{m}_{i} \mathbf{n}_{i}\right)^{\frac{\omega}{3}} \leq \underline{\mathbf{R}}(T)$, and more generally we worked with degenerations of $T^{\otimes N}$ as well. Informally define the degeneracy value of $T$ to be the best upper bound on $\omega$ we can get in this manner. More precisely:

Definition 3.4.7.1. Let $T \in A \otimes B \otimes C$. Fix $N \geq 1$ and $\rho \in[2,3]$. Define $V_{\rho, N}^{\text {degen }}(T)$ to be the maximum of $\sum_{i=1}^{s}\left(\mathbf{l}_{i} \mathbf{m}_{i} \mathbf{n}_{i}\right)^{\frac{\rho}{3}}$ over all degenerations of $T^{\otimes N}$ to $\bigoplus_{i=1}^{s} M_{\left\langle\mathbf{l}_{i}, \mathbf{m}_{i}, \mathbf{n}_{i}\right\rangle}$ over all choices of $s, \mathbf{l}_{i}, \mathbf{m}_{i}, \mathbf{n}_{i}$ and define the degeneracy value of $T$ to be $V_{\rho}^{\text {degen }}(T):=\sup _{N} V_{\rho, N}^{\text {degen }}(T)^{\frac{1}{N}}$.

The asymptotic sum inequality implies $V_{\omega}^{\text {degen }}(T) \leq \underline{\mathbf{R}}(T)$, or in other words, if $V_{\rho}^{\text {degen }}(T) \geq \underline{\mathbf{R}}(T)$, then $\omega \leq \rho$.

The supremum in the definition can be replaced by a limit, thanks to Fekete's lemma, since the sequence $\log \left(V_{\rho, N}^{\text {degen }}(T)\right)$ is super-additive:
Lemma 3.4.7.2 (Fekete's Lemma). For every super-additive sequence $\left\{a_{n}\right\}_{n=1}^{\infty}$ (i.e. $a_{n+m} \geq a_{n}+a_{m}$ ), the limit $\lim _{n \rightarrow \infty} \frac{a_{n}}{n}$ exists (possibly $+\infty$ ) and is equal to $\sup \frac{a_{n}}{n}$.
Exercise 3.4.7.3: (3) Prove Fekete's Lemma.
Fekete's lemma implies $\frac{1}{N} \log V_{\rho, N}^{\text {degen }}(T)$ tends to a limit. See [AFLG15] for details.

There is also an analogue of the asymptotic sum inequality for degeneracy value:
Theorem 3.4.7.4. $\sum_{i=1}^{s} V_{\omega}^{\text {degen }}\left(T_{i}\right) \leq \underline{\mathbf{R}}\left(\oplus_{i=1}^{s} T_{i}\right)$.
The proof is similar to the proof of the asymptotic sum inequality. It is clear that $V_{\omega}^{\text {degen }}\left(T_{1} \otimes T_{2}\right) \geq V_{\omega}^{\text {degen }}\left(T_{1}\right) \otimes V_{\omega}^{\text {degen }}\left(T_{2}\right)$. To show $V_{\omega}^{\text {degen }}\left(T_{1} \oplus\right.$ $\left.T_{2}\right) \geq V_{\omega}^{\text {degen }}\left(T_{1}\right)+V_{\omega}^{\text {degen }}\left(T_{2}\right)$ one expands out $V_{\omega, N}^{\text {degen }}\left(T_{1} \oplus T_{2}\right)$, the result is a sum of products with coefficients, but as with the asymptotic sum inequality, one can essentially just look at the largest term, and as $N$ tends to infinity, the coefficient becomes irrelevant after taking $N$-th roots.

Thus tensors of low border rank with high degeneracy value give upper bounds on $\omega$. The problem is that we have no systematic way of estimating degeneracy value. For an extreme example, for a given $r$ the tensor of border rank $r$ with the highest degeneracy value is $M_{\langle 1\rangle}^{\oplus r}$ as all border rank $r$ tensors are degenerations of it.

In subsequent work, researchers restrict to a special type of value that is possible to estimate.
3.4.8. The value of a tensor. Let $\operatorname{End}(A) \times \operatorname{End}(B) \times \operatorname{End}(C)$ act on $A \otimes B \otimes C$ by the action inherited from the $G L(A) \times G L(B) \times G L(C)$ action (not the Lie algebra action). Then for all $X \in \operatorname{End}(A) \times \operatorname{End}(B) \times \operatorname{End}(C)$ and $T \in A \otimes B \otimes C$, we have $\mathbf{R}(X \cdot T) \leq \mathbf{R}(T)$ and $\underline{\mathbf{R}}(X \cdot T) \leq \underline{\mathbf{R}}(T)$ by Exercise 2.1.6.2.
Definition 3.4.8.1. One says $T$ restricts to $T^{\prime}$ if $T^{\prime} \in \operatorname{End}(A) \times \operatorname{End}(B) \times$ $\operatorname{End}(C) \cdot T$.

Definition 3.4.8.2. For $T \in A \otimes B \otimes C, N \geq 1$ and $\rho \in[2,3]$ define $V_{\rho, N}^{\text {restr }}(T)$ to be the maximum of $\sum_{i=1}^{s}\left(\mathbf{l}_{i} \mathbf{m}_{i} \mathbf{n}_{i}\right)^{\frac{\rho}{3}}$ over all restrictions of $T^{\otimes N}$ to $\oplus_{i=1}^{s} M_{\left\langle\mathbf{l}_{i}, \mathbf{m}_{i}, \mathbf{n}_{i}\right\rangle}$ and define the restriction value of $T$ to be $V_{\rho}^{\text {restr }}(T):=\sup _{N} V_{\rho, N}^{\text {restr }}(T)^{\frac{1}{N}}$.

I emphasize that the degeneration used by Strassen is more general than restriction.

Coppersmith-Winograd and all subsequent work, use only the following type of restriction:

Definition 3.4.8.3. Let $A, B, C$ be given bases, so write them as $\mathbb{C}^{\mathbf{a}}, \mathbb{C}^{\mathbf{b}}, \mathbb{C}^{\mathbf{c}}$. We say $T \in \mathbb{C}^{\mathbf{a}} \otimes \mathbb{C}^{\mathbf{b}} \otimes \mathbb{C}^{\mathbf{c}}$ combinatorially restricts to $T^{\prime}$ if $T$ restricts to $T^{\prime}$ by setting some of the coordinates of $T$ to zero.

The condition that $T \in \mathbb{C}^{\mathbf{a}} \otimes \mathbb{C}^{\mathbf{b}} \otimes \mathbb{C}^{\mathbf{c}}$ admits a combinatorial restriction to the matrix multiplication tensor $M_{\langle\mathbf{l}, \mathbf{m}, \mathbf{n}\rangle}$ may be phrased as follows (following [CU03]): write $a_{\alpha}, b_{\beta}, c_{\gamma}$ for the given bases of $A, B, C$ and write $T=\sum_{\alpha=1}^{\mathbf{a}} \sum_{\beta=1}^{\mathbf{b}} \sum_{\gamma=1}^{\mathbf{c}} t^{\alpha, \beta, \gamma} a_{\alpha} \otimes b_{\beta} \otimes c_{\gamma}$. Then $T \in \mathbb{C}^{\mathbf{a}} \otimes \mathbb{C}^{\mathbf{b}} \otimes \mathbb{C}^{\mathbf{c}}$ combinatorially restricts to $M_{\langle\mathbf{l}, \mathbf{m}, \mathbf{n}\rangle}$ means that there exist injections

$$
\begin{aligned}
\alpha:[\mathbf{l}] \times[\mathbf{m}] & \rightarrow[\mathbf{a}] \\
\beta:[\mathbf{m}] \times[\mathbf{n}] & \rightarrow[\mathbf{b}] \\
\gamma:[\mathbf{n}] \times[\mathbf{l}] & \rightarrow[\mathbf{c}]
\end{aligned}
$$

such that

$$
t^{\alpha\left(i, j^{\prime}\right), \beta\left(j, k^{\prime}\right), \gamma\left(k, i^{\prime}\right)}=\left\{\begin{array}{cc}
1 & \text { if } i=i^{\prime}, j=j^{\prime}, k=k^{\prime}  \tag{3.4.3}\\
0 & \text { otherwise }
\end{array}\right\} .
$$

One can similarly phrase combinatorial restriction to a sum of disjoint matrix multiplication tensors.
Definition 3.4.8.4. For $T \in \mathbb{C}^{\mathbf{a}} \otimes \mathbb{C}^{\mathbf{b}} \otimes \mathbb{C}^{\mathbf{c}}, N \geq 1$ and $\rho \in[2,3]$ define $V_{\rho, N}(T)$ to be the maximum of $\sum_{i=1}^{s}\left(\mathbf{l}_{i} \mathbf{m}_{i} \mathbf{n}_{i}\right)^{\frac{\rho}{3}}$ over all combinatorial restrictions of $T^{\otimes N}$ to $\oplus_{i=1}^{s} M_{\left\langle 1_{i}, \mathbf{m}_{i}, \mathbf{n}_{i}\right\rangle}$ and define the combinatorial value (or value for short, since it is the value used in the literature) of $T$ to be $V_{\rho}(T):=\lim _{N \rightarrow \infty} V_{\rho, N}(T)^{\frac{1}{N}}$. (The limit is shown to exist in [DS13].)

Note that the values satisfy $V_{\rho}^{\text {degen }} \geq V_{\rho}^{\text {restr }} \geq V_{\rho}$. As with all the values we have

- $V_{\rho}(T)$ is a non-decreasing function of $\rho$,
- $V_{\omega}(T) \leq \underline{\mathbf{R}}(T)$.

Thus if $V_{\rho}(T) \geq \underline{\mathbf{R}}(T)$, then $\omega \leq \rho$.
Combinatorial value can be estimated in principle, as for each $N$, there are only a finite number of combinatorial restrictions. In practice, the tensor
is presented in such a way that there are "obvious" combinatorial degenerations to disjoint matrix multiplication tensors and at first, one optimizes just among these obvious combinatorial degenerations. However, it may be that there are matrix multiplication tensors of the form $\sum_{j} a_{0} \otimes b_{j} \otimes c_{j}$ as well as tensors of the form $a_{0} \otimes b_{k} \otimes c_{k}$ where $k$ is not in the range of $j$. Then one can merge these tensors to $a_{0} \otimes\left(\sum_{j} b_{j} \otimes c_{j}+b_{k} \otimes c_{k}\right)$ to increase value because although formally speaking they were not disjoint, they do not interfere with each other. (The value increases as e.g., $q^{\omega}+r^{\omega}<(q+r)^{\omega}$.) So the actual procedure is to optimize among combinatorial restrictions with merged tensors.
3.4.9. The Coppersmith-Winograd tensors. Coppersmith and Winograd apply Strassen's laser method, enhanced with merging, using combinatorial restrictions to the following two tensors:

The "easy Coppersmith-Winograd tensor":

$$
\begin{equation*}
T_{q, c w}:=\sum_{j=1}^{q} a_{0} \otimes b_{j} \otimes c_{j}+a_{j} \otimes b_{0} \otimes c_{j}+a_{j} \otimes b_{j} \otimes c_{0} \in \mathbb{C}^{q+1} \otimes \mathbb{C}^{q+1} \otimes \mathbb{C}^{q+1} \tag{3.4.4}
\end{equation*}
$$

Proposition 5.5.1.1 will imply $\mathbf{R}\left(T_{q, c w}\right)=2 q+1$ so the above expression is not optimal. We also have
Proposition 3.4.9.1. $\underline{\mathbf{R}}\left(T_{q, c w}\right)=q+2$.
Proof. Consider the second derivatives of a curve in the Segre: Let $x(t)=$ $a(t) \otimes b(t) \otimes c(t)$, write $x^{\prime}$ for $x^{\prime}(0)$ and similarly for all derivatives. Then

$$
x^{\prime \prime}=\left(a^{\prime \prime} \otimes b \otimes c+a \otimes b^{\prime \prime} \otimes c+a \otimes b \otimes c^{\prime \prime}\right)+2\left(a^{\prime} \otimes b^{\prime} \otimes c+a^{\prime} \otimes b \otimes c^{\prime}+a \otimes b^{\prime} \otimes c^{\prime}\right)
$$

so if we begin with the base point $a_{0} \otimes b_{0} \otimes c_{0}$, each term in the summand for $T_{q, c w}$ is a term of the second kind. The terms in the first parenthesis are ordinary tangent vectors. Thus take $q$ curves beginning at $a_{0} \otimes b_{0} \otimes c_{0}$, we can cancel out all the terms of the first type with a single vector to obtain the resulting border rank $q+2$ expression:

$$
\begin{aligned}
T_{q, c w}= & \lim _{t \rightarrow 0} \frac{1}{t^{2}}\left[\sum_{j=1}^{q}\left(a_{0}+t a_{j}\right) \otimes\left(b_{0}+t b_{j}\right) \otimes\left(c_{0}+t c_{j}\right)\right] \\
& -\left(a_{0}+t \sum_{j} a_{j}\right) \otimes\left(b_{0}+t \sum_{j} b_{j}\right) \otimes\left(c_{0}+t \sum_{j} c_{j}\right)-(q-1) a_{0} \otimes b_{0} \otimes c_{0} .
\end{aligned}
$$

Exercise 3.4.9.2: (2) Show that $\underline{\mathbf{R}}\left(T_{q, c w}\right) \geq q+2$ so that equality holds.

A slightly more complicated tensor yields even better results: Let

$$
\begin{align*}
T_{q, C W}:= & \sum_{j=1}^{q}\left(a_{0} \otimes b_{j} \otimes c_{j}+a_{j} \otimes b_{0} \otimes c_{j}+a_{j} \otimes b_{j} \otimes c_{0}\right)  \tag{3.4.5}\\
& +a_{0} \otimes b_{0} \otimes c_{q+1}+a_{0} \otimes b_{q+1} \otimes c_{0}+a_{q+1} \otimes b_{0} \otimes c_{0} \in \mathbb{C}^{q+2} \otimes \mathbb{C}^{q+2} \otimes \mathbb{C}^{q+2}
\end{align*}
$$

and call $T_{q, C W}$ the Coppersmith-Winograd tensor.
Exercise 3.4.9.3: (2) Show the Coppersmith-Winograd tensor also has border rank $q+2$ by modifying the curves used to obtain $T_{q, c w}$. ©

Now suggestively re-label $T_{q, C W}$ as we did with Strassen's tensor:

$$
\begin{align*}
T_{q, C W}:= & \sum_{j=1}^{q}\left(a_{0}^{[0]} \otimes b_{j}^{[1]} \otimes c_{j}^{[1]}+a_{j}^{[1]} \otimes b_{0}^{[0]} \otimes c_{j}^{[1]}+a_{j}^{[1]} \otimes b_{j}^{[1]} \otimes c_{0}^{[0]}\right)  \tag{3.4.6}\\
& +a_{0}^{[0]} \otimes b_{0}^{[0]} \otimes c_{q+1}^{[2]}+a_{0}^{[0]} \otimes b_{q+1}^{[2]} \otimes c_{0}^{[0]}+a_{q+1}^{[2]} \otimes b_{0}^{[0]} \otimes c_{0}^{[0]} \in \mathbb{C}^{q+2} \otimes \mathbb{C}^{q+2} \otimes \mathbb{C}^{q+2}
\end{align*}
$$

to see that $T_{q, C W}$ is the sum of 3 matrix multiplications of volume $q^{2}$, and 3 of volume 1 , all non-disjoint. To get more interesting matrix multiplications, consider $T_{q, C W}^{\otimes 2}$, but this time, instead of double superscripts, simply add the superscripts.

$$
\begin{aligned}
T_{q, C W}^{\otimes 2}= & \sum_{i, j=1}^{q}\left[a_{00}^{[0]} \otimes b_{i j}^{[2]} \otimes c_{i j}^{[2]}+a_{0 j}^{[1]} \otimes b_{i 0}^{[1]} \otimes c_{i j}^{[2]}+a_{0 j}^{[1]} \otimes b_{i j}^{[2]} \otimes c_{i 0}^{[1]}+a_{i 0}^{[1]} \otimes b_{0 j}^{[1]} \otimes c_{i j}^{[2]}+a_{i 0}^{[1]} \otimes b_{i j}^{[2]} \otimes c_{0 j}^{[1]}\right. \\
& \left.+a_{i j}^{[2]} \otimes b_{i 0}^{[1]} \otimes c_{0 j}^{[1]}+a_{i j}^{[2]} \otimes b_{00}^{[0]} \otimes c_{i j}^{[2]}+a_{i j}^{[2]} \otimes b_{i j}^{[2]} \otimes c_{00}^{[1]}+a_{i j}^{[2]} \otimes b_{0 j}^{[1]} \otimes c_{i 0}^{[1]}\right] \\
+ & \sum_{j=1}^{q}\left[a_{0, q+1}^{[2]} \otimes b_{j 0}^{[1]} \otimes c_{j 0}^{[1]}+a_{q+1,0}^{[2]} \otimes b_{0 j}^{[1]} \otimes c_{0 j}^{[1]}+a_{q+1, j}^{[3]} \otimes b_{0 j}^{[1]} \otimes c_{00}^{[0]}+a_{j, q+1}^{[3]} \otimes b_{j 0}^{[1]} \otimes c_{00}^{[0]}\right. \\
& \left.+a_{q+1, j}^{[3]} \otimes b_{00}^{[0]} \otimes c_{0 j}^{[1]}+a_{j, q+1}^{[3]} \otimes b_{00}^{[0]} \otimes c_{j 0}^{[1]}\right] \\
+ & a_{q+1, q+1}^{[4]} \otimes b_{00}^{[0]} \otimes c_{00}^{[0]}+a_{00}^{[0]} \otimes b_{q+1, j}^{[3]} \otimes c_{0 j}^{[1]}+a_{00}^{[0]} \otimes b_{0 j}^{[1]} \otimes c_{q+1, j}^{[3]} \\
+ & a_{00}^{[0]} \otimes b_{q+1, q+1}^{[4]} \otimes c_{00}^{[0]}+a_{00}^{[0]} \otimes b_{00}^{[0]} \otimes c_{q+1, q+1}^{[4]} .
\end{aligned}
$$

Now we have non-disjoint matrix multiplications of volumes $q^{2}, q$ and 1 . Thus when we zero-out terms to get disjoint matrix multiplications in $\left(T_{q, C W}^{\otimes 2}\right)^{\otimes N}$, in order to optimize value, we need to weight the $q^{2}$ terms more than the $q$ terms etc..

As mentioned above, one can obtain better upper bounds with merging. One needs to make a choice how to merge. Coppersmith and Winogrand
group the $\mathbb{C}^{\mathbf{a}^{2}}$-variables

$$
\begin{aligned}
\mathcal{A}^{[0]} & =\left\{a_{00}^{[0]}\right\} \\
\mathcal{A}^{[1]} & =\left\{a_{0}^{[1]}, a_{0 j}^{[1]}\right\} \\
\mathcal{A}^{[2]} & =\left\{a_{q+1,0}^{[2]}, a_{i j}^{[2]}, a_{0, q+1}^{[2]}\right\} \\
\mathcal{A}^{[3]} & =\left\{a_{q+1, j}^{33]}, a_{i, q+1}^{3]}\right\} \\
\mathcal{A}^{[4]} & =\left\{a_{q+1, q+1}^{[4]}\right\}
\end{aligned}
$$

and similarly for $b$ 's and $c$ 's. Then

$$
T_{q, C W}^{\otimes 2}=\sum_{I+J+K=4} \mathcal{A}^{[I]} \otimes \mathcal{B}^{[J]} \otimes \mathcal{C}^{[K]}
$$

where e.g., $\mathcal{A}^{[I]}$ is to be interpreted as the sum of all elements of $\mathcal{A}^{[I]}$. Most of these terms are just matrix multiplications, however terms with $1+1+2$ are not:

$$
\begin{aligned}
\mathcal{A}^{[1]} \otimes \mathcal{B}^{[1]} \otimes \mathcal{C}^{[2]}= & \sum_{i=1}^{q} a_{i 0}^{[1]} \otimes b_{i 0}^{[1]} \otimes c_{0, q+1}^{[2]}+\sum_{j=1}^{q} a_{0 j}^{[1]} \otimes b_{0 j}^{[1]} \otimes c_{q+1,0}^{[2]} \\
& +\sum_{i, j=1}^{q}\left[a_{i 0}^{[1]} \otimes b_{0 j}^{[1]} \otimes c_{i j}^{[2]}+a_{0 j}^{[1]} \otimes b_{i 0}^{[1]} \otimes c_{i j}^{[2]}\right] .
\end{aligned}
$$

To this term we estimate value using the laser method, i.e., we degenerate tensor powers of $\mathcal{A}^{[1]} \otimes \mathcal{B}^{[1]} \otimes \mathcal{C}^{[2]}$ to disjoint matrix multiplication tensors. Coppersmith and Winograd show that it has value at least $2^{\frac{2}{3}} q^{\omega}\left(q^{3 \omega}+2\right)^{\frac{1}{3}}$.

Now there is an optimization problem to solve, that I briefly discuss in §3.4.10 below.

Coppersmith and Winograd get their best result of $\omega<2.3755$ by merging $T_{q, C W}^{\otimes 2}$ and then optimizing over the various combinatorial restrictions. In subsequent work Stothers [Sto], resp. Williams [Wil], resp. LeGall [Gal] used merging with $T_{q, C W}^{\otimes 4}$ resp. $T_{q, C W}^{\otimes 8}$, resp. $T_{q, C W}^{\otimes 16}$ and $T_{q, C W}^{\otimes 32}$ leading to the current "world record":
Theorem 3.4.9.4. [Gal] $\omega<2.3728639$.
Ambainis, Filmus and LeGall [AFLG15] showed that taking higher powers of $T_{q, C W}$ when $q \geq 5$ cannot be used to prove $\omega<2.30$ by this method alone. Their argument avoids higher powers by more sophisticated methods to account for when potential merging in higher tensor powers can occur.

Thus one either needs to develop new methods, or find better base tensors.

I discuss the search for better base tensors in Remark 5.5.3.4.
3.4.10. How one optimizes in practice. To get an idea of how the optimization procedure works, start with some base tensor $T$ that contains a collection of matrix multiplication tensors $M_{\left\langle\mathbf{l}_{i}, \mathbf{m}_{i}, \mathbf{n}_{i}\right\rangle}, 1 \leq i \leq x$ that are not disjoint. Then $T^{\otimes N}$ will contain matrix multiplication tensors of the form $M_{\left\langle\mathbf{l}_{\mu}, \mathbf{m}_{\mu}, \mathbf{n}_{\mu}\right\rangle}$ where $\mathbf{l}_{\mu}=\mathbf{l}_{\mu_{1}} \cdots \mathbf{l}_{\mu_{N}}$ and similarly for $\mathbf{m}_{\mu}, \mathbf{n}_{\mu}$, where $\mu_{j} \in[x]$.

Each matrix multiplication tensor will occur with a certain multiplicity and certain variables. The problem becomes to zero out variables in a way that maximizes the value of what remains. More precisely, for large $N$, one wants to maximize the sum $\sum_{j} K_{j}\left(\mathbf{l}_{\mu_{j}} \mathbf{m}_{\mu_{j}} \mathbf{n}_{\mu_{j}}\right)^{\frac{\rho}{3}}$ where the surviving matrix multiplication tensors are $M_{\left\langle\mathbf{1}_{j} \mathbf{m}_{\mu_{j}} \mathbf{n}_{\mu_{j}}\right\rangle}^{\oplus K_{j}}$ and disjoint. One then takes the smallest $\rho$ such that $\sum_{j} K_{j}\left(\mathbf{l}_{\mu_{j}} \mathbf{m}_{\mu_{j}} \mathbf{n}_{\mu_{j}}\right)^{\frac{\rho}{3}} \geq \underline{\mathbf{R}}(T)$ and concludes $\omega \leq \rho$. One ingredient is the Salem-Spencer Theorem:
Theorem 3.4.10.1 (Salem and Spencer [SS42]). Given $\epsilon>0$, there exists $M_{\epsilon} \simeq 2^{\frac{c}{\epsilon^{2}}}$ such that for all $M>M_{\epsilon}$, there is a set $B$ of $M^{\prime}>M^{1-\epsilon}$ distinct integers $0<b_{1}<b_{2}<\cdots<b_{M^{\prime}}<\frac{M}{2}$ with no three terms in an arithmetic progression, i.e., for $b_{i}, b_{j}, b_{k} \in B, b_{i}+b_{j}=2 b_{k}$ if and only if $b_{i}=b_{j}=b_{k}$. In fact no three terms form an arithmetic progression $\bmod M$.

This theorem assures one can get away with only zero-ing out a relatively small number of terms, so in some sense it plays the role of Strassen's degeneration theorem. I state it explicitly to emphasize that it is an existence result, not an algorithm. In the general case one assigns probability distributions and optimizes using techniques from probability to determine what percentage of each type gets zero-ed out. See [CW82] for the basic idea and [AFLG15] for the state of the art regarding this optimization.

### 3.5. The Cohn-Umans program

A conceptually appealing approach to proving upper bounds on $\omega$ was initiated by H. Cohn and C. Umans.

Imagine a tensor that comes presented in two different bases. In one, the cost of the tensor is clear: it may be written as a sum of small disjoint matrix multiplication tensors. On the other hand, in the other its value (in the sense discussed above) is high, because it may be seen to degenerate to good matrix multiplication tensors. Such a situation does arise in practice! It occurs for structure tensors for the group algebra of a finite group, as defined below. In one (the "matrix coefficient basis"), one gets an upper bound on the rank of the tensor, and in the other (the "standard basis") there are many potential combinatorial degenerations and one gets a lower bound on the value.

I state the needed representation theory now, and defer proofs of the statements to $\S 8.6$. I then present their method.
3.5.1. Structure tensor of an algebra. Let $\mathcal{A}$ be a finite dimensional algebra, i.e., a vector space with a multiplication operation, with basis $a_{1}, \ldots, a_{\mathbf{a}}$ and dual basis $\alpha^{1}, \ldots, \alpha^{\mathbf{a}}$. Write $a_{i} a_{j}=\sum_{k} A_{i j}^{k} a_{k}$ for the multiplication in $\mathcal{A}$, where the $A_{i j}^{k}$ are constants. The multiplication $\mathcal{A} \times \mathcal{A} \rightarrow \mathcal{A}$ is bilinear and one defines the corresponding structure tensor of $\mathcal{A}$

$$
\begin{equation*}
T_{\mathcal{A}}:=\sum_{i, j, k} A_{i j}^{k} \alpha^{i} \otimes \alpha^{j} \otimes a_{k} \in \mathcal{A}^{*} \otimes \mathcal{A}^{*} \otimes \mathcal{A} \tag{3.5.1}
\end{equation*}
$$

For example, $M_{\langle\mathbf{n}\rangle}$ is the structure tensor for the algebra of $\mathbf{n} \times \mathbf{n}$-matrices with operation matrix multiplication.
The group algebra of a finite group. Let $G$ be a finite group and let $\mathbb{C}[G]$ denote the vector space of complex-valued functions on $G$, called the group algebra of $G$. The following exercise justifies the name:
Exercise 3.5.1.1: (1) Show that if the elements of $G$ are $g_{1}, \ldots, g_{r}$, then $\mathbb{C}[G]$ has a basis indexed $\delta_{g_{1}}, \ldots, \delta_{g_{r}}$, where $\delta_{g_{i}}\left(g_{j}\right)=\delta_{i j}$. Show that $\mathbb{C}[G]$ may be given the structure of an algebra by defining $\delta_{g_{i}} \delta_{g_{j}}:=\delta_{g_{i} g_{j}}$ and extending linearly.

Thus if $G$ is a finite group, then $T_{\mathbb{C}[G]}=\sum_{g, h \in G} \delta_{g}^{*} \otimes \delta_{h}^{*} \otimes \delta_{g h}$.

## Example 3.5.1.2.

$$
T_{\mathbb{C}\left[\mathbb{Z}_{m}\right]}=\sum_{0 \leq i, j<m} \delta_{i}^{*} \otimes \delta_{j}^{*} \otimes \delta_{i+j \bmod m}
$$

Notice that, introducing coordinates $x_{0}, \ldots, x_{m-1}$ on $\mathbb{C}\left[\mathbb{Z}_{m}\right]$, so $v \in \mathbb{C}\left[\mathbb{Z}_{m}\right]$ may be written $\sum x_{s} \delta_{s}$, one obtains a circulant matrix for $T_{\mathbb{C}\left[\mathbb{Z}_{m}\right]}\left(\mathbb{C}\left[\mathbb{Z}_{m}\right]^{*}\right) \subset$ $\mathbb{C}\left[\mathbb{Z}_{m}\right]^{*} \otimes \mathbb{C}\left[\mathbb{Z}_{m}\right]^{*}:$

$$
T_{\mathbb{C}\left[\mathbb{Z}_{m}\right]}\left(\mathbb{C}\left[\mathbb{Z}_{m}\right]^{*}\right)=\left\{\left.\left(\begin{array}{cccc}
x_{0} & x_{1} & \cdots & x_{m-1}  \tag{3.5.2}\\
x_{m-1} & x_{0} & x_{1} & \cdots \\
\vdots & & \ddots & \\
x_{1} & x_{2} & \cdots & x_{0}
\end{array}\right) \right\rvert\, x_{j} \in \mathbb{C}\right\}
$$

In what follows I slightly abuse notation and write the matrix with entries $x_{j}$ rather than the form above. Note that all entries of the matrix are non-zero and filled with basis vectors. This holds in general for the presentation of $\mathbb{C}[G]$ in the standard basis, which makes it useful for combinatorial restrictions.

What are $\underline{\mathbf{R}}\left(T_{\mathbb{C}\left[\mathbb{Z}_{m}\right]}\right)$ and $\mathbf{R}\left(T_{\mathbb{C}\left[\mathbb{Z}_{m}\right]}\right)$ ? The space of circulant matrices forms an abelian subspace, which indicates the rank and border rank might
be minimal or nearly minimal among concise tensors. We will determine the rank and border rank of $T_{\mathbb{C}\left[\mathbb{Z}_{m}\right]}$ momentarily via the discrete Fourier transform.
3.5.2. The structure theorem of $\mathbb{C}[G]$. I give a proof of the following theorem and an explanation of the $G \times G$-module structure on $\mathbb{C}[G]$ in $\S 8.6 .5$.
Theorem 3.5.2.1. Let $G$ be a finite group, then as a $G \times G$-module,

$$
\begin{equation*}
\mathbb{C}[G]=\bigoplus_{i} V_{i}^{*} \otimes V_{i} \tag{3.5.3}
\end{equation*}
$$

where the sum is over all the distinct irreducible representations of $G$. In particular, if $\operatorname{dim} V_{i}=d_{i}$, then as an algebra,

$$
\begin{equation*}
\mathbb{C}[G] \simeq \bigoplus_{i} M a t_{d_{i} \times d_{i}}(\mathbb{C}) \tag{3.5.4}
\end{equation*}
$$

3.5.3. The (generalized) discrete Fourier transform. We have two natural expressions for $T_{\mathbb{C}[G]}$, the original presentation in terms of the algebra multiplication in terms of delta functions, the standard basis, and the matrix coefficient basis in terms of the entries of the matrices in (3.5.4). The change of basis matrix from the standard basis to the matrix coefficient basis is called the (generalized) Discrete Fourier Transform (DFT).

Example 3.5.3.1. The classical DFT is the case $G=\mathbb{Z}_{m}$. The irreducible representations of $\mathbb{Z}_{m}$ are all one dimensional: $\rho_{k}: \mathbb{Z}_{m} \rightarrow G L_{1}$. Let $\sigma \in \mathbb{Z}_{m}$ be a generator, then $\rho_{k}(\sigma) v=e^{\frac{2 \pi i k}{m}} v$ for $0 \leq k \leq m$. The DFT matrix is

$$
\left(e^{\frac{2 \pi i(j+k)}{m}}\right)_{0 \leq j, k \leq m-1} .
$$

Proposition 3.5.3.2. $\underline{\mathbf{R}}\left(T_{\mathbb{C}\left[\mathbb{Z}_{m}\right]}\right)=\mathbf{R}\left(T_{\mathbb{C}\left[\mathbb{Z}_{m}\right]}\right)=m$.
Proof. Theorem 3.5.2.1 implies $T_{\mathbb{C}\left[\mathbb{Z}_{m}\right]}=M_{\langle 1\rangle}^{\oplus m}$.
Compared with (3.5.2), in the matrix coefficient basis the image $T_{\mathbb{C}\left[\mathbb{Z}_{m}\right]}\left(\mathbb{C}\left[\mathbb{Z}_{m}\right]^{*}\right)$ is the set of diagonal matrices:

$$
T_{\mathbb{C}\left[\mathbb{Z}_{m}\right]}\left(\mathbb{C}\left[\mathbb{Z}_{m}\right]^{*}\right)=\left(\begin{array}{llll}
y_{0} & & & \\
& y_{1} & & \\
& & \ddots & \\
& & & y_{m-1}
\end{array}\right)
$$

Exercise 3.5.3.3: (2) Show that if $T \in \hat{\sigma}_{r}^{0, h}$, then $\mathbf{R}(T) \leq r(h+1)$. ©
Exercise 3.5.3.4: (2) Obtain a fast algorithm for multiplying two polynomials in one variable by the method you used to solve the previous exercise.

Example 3.5.3.5. Consider $\mathfrak{S}_{3}$. In the standard basis,

$$
T_{\mathbb{C}\left[\mathfrak{S}_{3}\right]}\left(\mathbb{C}\left[\mathfrak{S}_{3}\right]^{*}\right)=\left(\begin{array}{llllll}
x_{0} & x_{1} & x_{2} & x_{3} & x_{4} & x_{5} \\
x_{1} & x_{0} & x_{4} & x_{5} & x_{2} & x_{3} \\
x_{2} & x_{5} & x_{0} & x_{4} & x_{3} & x_{1} \\
x_{3} & x_{4} & x_{5} & x_{0} & x_{1} & x_{2} \\
x_{4} & x_{3} & x_{1} & x_{2} & x_{5} & x_{0} \\
x_{5} & x_{2} & x_{3} & x_{1} & x_{0} & x_{4}
\end{array}\right) .
$$

Here I have written an element of $\mathbb{C}\left[\mathfrak{S}_{3}\right]$ as $x_{0} \delta_{\text {Id }}+x_{1} \delta_{(12)}+x_{2} \delta_{(13)}+x_{3} \delta_{(23)}+$ $x_{4} \delta_{(123)}+x_{5} \delta_{(132)}$. The irreducible representations of $\mathfrak{S}_{3}$ are the trivial, denoted [3], the sign, denoted $[1,1,1]$ and the two-dimensional standard representation (the complement of the trivial in $\mathbb{C}^{3}$ ), which is denoted $[2,1]$. (See $\S 8.6 .5$ for an explanation of the notation.) Since $\operatorname{dim}[3]=1, \operatorname{dim}[1,1,1]=1$ and $\operatorname{dim}[2,1]=2$, by Theorem 3.5.2.1 $T_{\mathbb{C}\left[\mathfrak{S}_{3}\right]}=M_{\langle 1\rangle}^{\oplus 2} \oplus M_{\langle 2\rangle}$, and in the matrix coefficient basis:

$$
T_{\mathbb{C}\left[\mathfrak{S}_{3}\right]}\left(\mathbb{C}\left[\mathfrak{S}_{3}\right]^{*}\right)=\left(\begin{array}{llllll}
y_{0} & & & & & \\
& y_{1} & & & & \\
& & y_{2} & y_{3} & & \\
& & y_{4} & y_{5} & & \\
& & & & y_{2} & y_{3} \\
& & & & y_{4} & y_{5}
\end{array}\right)
$$

where the blank entries are zero. We conclude $\mathbf{R}\left(T_{\mathbb{C}\left[\mathfrak{G}_{3}\right]}\right) \leq 1+1+7=9$.
3.5.4. Upper bounds via finite groups. Here is the main idea:

Use the standard basis to get a lower bound on the value of $T_{\mathbb{C}[G]}$ and the matrix coefficient basis to get an upper bound on its cost.

Say $T_{\mathbb{C}[G]}$ expressed in its standard basis combinatorially restricts to a sum of matrix multiplications, say $\oplus_{j=1}^{s} M_{\left\langle\mathbf{l}_{j}, \mathbf{m}_{j}, \mathbf{n}_{j}\right\rangle}$. The standard basis is particularly well suited to combinatorial restrictions because all the coefficients of the tensor in this basis are zero or one, and all the entries of the matrix $T_{\mathbb{C}[G]}\left(\mathbb{C}[G]^{*}\right)$ are nonzero and coordinate elements. (Recall that all the entries of the matrix $M_{\langle 1, \mathbf{m}, \mathbf{n}\rangle}\left(A^{*}\right)$ are either zero or coordinate elements.) Using the matrix coefficient basis, we see $T_{\mathbb{C}[G]}=$ $\oplus_{u=1}^{q} M_{\left\langle d_{u}\right\rangle}$, where $d_{u}$ is the dimension of the $u$-th irreducible representation of $G$. Thus $\underline{\mathbf{R}}\left(\oplus_{j=1}^{s} M_{\left\langle\mathbf{1}_{j}, \mathbf{m}_{j}, \mathbf{n}_{j}\right\rangle}\right) \leq \underline{\mathbf{R}}\left(\oplus_{u=1}^{q} M_{\left\langle d_{u}\right\rangle}\right)$ and $\mathbf{R}\left(\oplus_{j=1}^{s} M_{\left\langle\mathbf{1}_{j}, \mathbf{m}_{j}, \mathbf{n}_{j}\right\rangle}\right) \leq$ $\mathbf{R}\left(\oplus_{u=1}^{q} M_{\left\langle d_{u}\right\rangle}\right)$.

The asymptotic sum inequality implies:
Proposition 3.5.4.1. [CU03, CU13] If $T_{\mathbb{C}[G]}$ degenerates to $\oplus_{j=1}^{s} M_{\left\langle\mathbf{1}_{j}, \mathbf{m}_{j}, \mathbf{n}_{j}\right\rangle}$ and $d_{u}$ are the dimensions of the irreducible representations of $G$, then $\sum_{j=1}^{s}\left(\mathbf{l}_{j} \mathbf{m}_{j} \mathbf{n}_{j}\right)^{\frac{\omega}{3}} \leq \mathbf{R}\left(\oplus_{u=1}^{q} M_{\left\langle d_{u}\right\rangle}\right) \leq \sum d_{u}^{3}$. In fact, $\sum_{j=1}^{s}\left(\mathbf{l}_{j} \mathbf{m}_{j} \mathbf{n}_{j}\right)^{\frac{\omega}{3}} \leq$ $\sum d_{u}^{\omega}$.

In this section I will denote the standard basis for $\mathbb{C}[G]$ given by the group elements (which I have been denoting $\delta_{g_{i}}$ ) simply by $g_{i}$.

Basis elements of $\mathbb{C}[G]$ are indexed by elements of $G$, so our sought-after combinatorial restriction is of the form:

$$
\begin{array}{r}
\alpha:[\mathbf{l}] \times[\mathbf{m}] \rightarrow G \\
\beta:[\mathbf{m}] \times[\mathbf{n}] \rightarrow G \\
\gamma:[\mathbf{n}] \times[\mathbf{l}] \rightarrow G .
\end{array}
$$

Recall the requirement that $t^{\alpha\left(i, j^{\prime}\right), \beta\left(j, k^{\prime}\right), \gamma\left(k, i^{\prime}\right)}$ is one if and only if $i=i^{\prime}$, $j=j^{\prime}, k=k^{\prime}$, and is otherwise zero. Here, when considering $T_{\mathbb{C}[G]}$ as a trilinear map, we have

$$
t^{\alpha, \beta, \gamma}= \begin{cases}1 & \alpha \beta \gamma=\mathrm{Id} \\ 0 & \text { otherwise }\end{cases}
$$

We want that $\alpha\left(i, j^{\prime}\right) \beta\left(j, k^{\prime}\right) \gamma\left(k, i^{\prime}\right)=$ Id if and only if $i=i^{\prime}, j=j^{\prime}, k=$ $k^{\prime}$. To simplify the requirement, assume the maps factor to $s_{1}:[1] \rightarrow G$, $s_{2}:[\mathbf{m}] \rightarrow G, s_{3}:[\mathbf{n}] \rightarrow G$, and that $\alpha\left(i, j^{\prime}\right)=s_{1}^{-1}(i) s_{2}\left(j^{\prime}\right), \beta\left(j, k^{\prime}\right)=$ $s_{2}^{-1}(j) s_{3}\left(k^{\prime}\right)$ and $\gamma\left(k, i^{\prime}\right)=s_{3}^{-1}(k) s_{1}\left(i^{\prime}\right)$. Our requirement becomes

$$
s_{1}^{-1}(i) s_{2}\left(j^{\prime}\right) s_{2}^{-1}(j) s_{3}\left(k^{\prime}\right) s_{3}^{-1}(k) s_{1}\left(i^{\prime}\right)=\operatorname{Id} \Leftrightarrow i=i^{\prime}, j=j^{\prime}, k=k^{\prime} .
$$

Let $S_{j}$ denote the image of $s_{j}$. Our requirement is summarized in the following definition:

Definition 3.5.4.2. [CU03] A triple of subsets $S_{1}, S_{2}, S_{3} \subset G$ satisfies the triple product property if for any $s_{j}, s_{j}^{\prime} \in S_{j}, s_{1}^{\prime} s_{1}^{-1} s_{2}^{\prime} s_{2}^{-1} s_{3}^{\prime} s_{3}^{-1}=\mathrm{Id}$ implies $s_{1}^{\prime}=s_{1}, s_{2}^{\prime}=s_{2}, s_{3}^{\prime}=s_{3}$.

There is a corresponding simultaneous triple product property when there is a combinatorial restriction to a collection of disjoint matrix multiplication tensors.

Example 3.5.4.3. [CKSU05] Let $G=\left(\mathbb{Z}_{N}^{\times 3} \times \mathbb{Z}_{N}^{\times 3}\right) \rtimes \mathbb{Z}_{2}$ where $\mathbb{Z}_{2}$ acts by switching the two factors, so $|G|=2 N^{6}$. Write elements of $G$ as $\left[\left(\omega^{i}, \omega^{j}, \omega^{k}\right)\left(\omega^{l}, \omega^{s}, \omega^{t}\right) \tau^{\epsilon}\right]$ where $0 \leq i, j, k, s, t, u \leq N-1, \omega$ is a primitive $N$-th root of unity, $\tau$ is a generator of $\mathbb{Z}_{2}$, and $\epsilon \in\{0,1\}$. Set $\mathbf{l}=\mathbf{m}=\mathbf{n}=2 N(N-1)$. Label the elements of $[\mathbf{n}]=[2 N(N-1)]$ by a triple $(a, b, \epsilon)$ where $1 \leq a \leq N-1,0 \leq b \leq N-1$ and $\epsilon \in\{0,1\}$, and
define

$$
\begin{aligned}
s_{1}:[\mathbf{l}] & \rightarrow G \\
(a, b, \epsilon) & \mapsto\left[\left(\omega^{a}, 1,1\right)\left(1, \omega^{b}, 1\right) \tau^{\epsilon}\right] \\
s_{2}:[\mathbf{m}] & \rightarrow G \\
(a, b, \epsilon) & \mapsto\left[\left(1, \omega^{a}, 1\right)\left(1,1, \omega^{b}\right) \tau^{\epsilon}\right] \\
s_{3}:[\mathbf{n}] & \rightarrow G \\
(a, b, \epsilon) & \mapsto\left[\left(1,1, \omega^{a}\right)\left(\omega^{b}, 1,1\right) \tau^{\epsilon}\right] .
\end{aligned}
$$

As explained in [CKSU05], the triple product property indeed holds (there are several cases), so $T_{\mathbb{C}[G]}$ combinatorially restricts to $M_{\langle 2 N(N-1)\rangle}$. Now $G$ has $2 N^{3}$ irreducible one dimensional representations and $\binom{N^{3}}{2}$ irreducible two dimensional representations (see [CKSU05]). Thus $\mathbf{R}\left(M_{\langle 2 N(N-1)\rangle}\right) \leq$ $2 N^{3}+8\binom{N^{3}}{2}$, which is less than $\mathbf{n}^{3}=[2 N(N-1)]^{3}$ for all $N \geq 5$. Asymptotically this is about $\frac{7}{16} \mathbf{n}^{3}$. If one applies Proposition 3.5.4.1 with $N=17$ (which is optimal), one obtains $\omega<2.9088$. Note that this does not even exploit Strassen's algorithm, so one actually has $\mathbf{R}\left(M_{\langle\mathbf{n}\rangle}\right) \leq 2 N^{3}+7\binom{N^{3}}{2}$, however this does not effect the asymptotics. If one could use the failure of additivity for border rank one potentially could do better.

While this is worse than what one would obtain just using Strassen's algorithm (writing $40=32+8$ and using Strassen in blocks), the algorithm is different. In [CKSU05] they obtain a bound of $\omega<2.41$ by such methods, but key lemmas in their proof are almost the same as the key lemmas used by Coopersmith-Winograd in their optimizations.
3.5.5. Further ideas towards upper bounds. The structure tensor of $\mathbb{C}[G]$ had the convenient property that in the standard basis all the coefficients of the tensor are zero or one, and all entries of the matrix $T_{\mathbb{C}[G]}\left(\mathbb{C}[G]^{*}\right)$ are basis vectors. In [CU13] they propose looking at combinatorial restrictions of more general structure tensors, where the coefficients can be more general, but vestiges of these properties are preserved. They make the following definition, which is very particular to matrix multiplication:

Definition 3.5.5.1. We say $T \in A \otimes B \otimes C$, given in bases $a_{\alpha}, b_{\beta}, c_{\gamma}$ of $A, B, C$, combinatorially supports $M_{\langle\mathbf{l}, \mathbf{m}, \mathbf{n}\rangle}$, if such that, writing $T=\sum t^{\alpha, \beta, \gamma} a_{\alpha} \otimes b_{\beta} \otimes c_{\gamma}$,
there exist injections

$$
\begin{array}{r}
\alpha:[\mathbf{l}] \times[\mathbf{m}] \rightarrow[\mathbf{a}] \\
\beta:[\mathbf{m}] \times[\mathbf{n}] \rightarrow[\mathbf{b}] \\
\gamma:[\mathbf{n}] \times[\mathbf{l}] \rightarrow[\mathbf{c}]
\end{array}
$$

such that $t^{\alpha\left(i, j^{\prime}\right), \beta\left(j, k^{\prime}\right) \gamma\left(k, i^{\prime}\right)} \neq 0$ if and only if $i=i^{\prime}, j=j^{\prime}$ and $k=k^{\prime}$. (Recall that $T$ combinatorially restricts to $M_{\langle\mathbf{l}, \mathbf{m}, \mathbf{n}\rangle}$ if moreover $t^{\alpha(i, j), \beta(j, k) \gamma(k, i)}=$ 1 for all $i, j, k$.)
$T$ combinatorially supports $M_{\langle\mathbf{m}, \mathbf{n}, \mathbf{1}\rangle}$ if there exists a coordinate expression of $T$ such that, upon setting some of the coefficients in the multidimensional matrix representing $T$ to zero, one obtains mnl nonzero entries such that in that coordinate system, matrix multiplication is supported on exactly those mnl entries. They then proceed to define the $s$-rank of a tensor $T^{\prime}$, which is the lowest rank of a tensor $T$ that combinatorially supports it. This is a strange concept because the s-rank of a generic tensor is one: a generic tensor is combinatorially supported by $T=\left(\sum_{j} a_{j}\right) \otimes\left(\sum_{k} b_{k}\right) \otimes\left(\sum_{l} c_{l}\right)$ where $\left\{a_{j}\right\}$ is a basis of $A$ etc..

Despite this, they show that $\omega \leq \frac{3}{2} \omega_{s}-1$ where $\omega_{s}$ is the analog of the exponent of matrix multiplication for $s$-rank. In particular, $\omega_{s}=2$ would imply $\omega=2$. The idea of the proof is that if $T$ combinatorially supports $M_{\langle\mathbf{n}\rangle}$, then $T^{\otimes 3}$ combinatorially degenerates to $M_{\langle\mathbf{n}\rangle}^{\oplus t}$ with $t=O\left(\mathbf{n}^{2-o(1)}\right)$. Compare this with the situation when $T$ combinatorially restricts to $M_{\langle\mathbf{n}\rangle}$, then $T^{\otimes 3}$ combinatorially restricts to $M_{\langle\mathbf{n}\rangle} \otimes M_{\left\langle\mathbf{n}^{2}\right\rangle}$ and thus toric degenerates to $M_{\langle\mathbf{n}\rangle}^{\oplus\left\lfloor\frac{3}{4} \mathbf{n}^{2}\right\rfloor}$ by Theorem 3.4.3.1.

## Chapter 4

## The complexity of Matrix multiplication III: explicit decompositions via geometry

One might argue that the exponent of matrix multiplication is unimportant for the world we live in, since $\omega$ might not be relevant until the sizes of the matrices are on the order of number of atoms in the known universe. For implementation, it is more important to develop explicit decompositions that provide a savings for matrices of sizes that need to be multiplied in practice. One purpose of this chapter is to discuss such decompositions. Another is to gain insight into the asymptotic situation by studying the symmetry groups that occur in the known decompositions of $M_{\langle\mathbf{n}\rangle}$. I begin, in §4.1, by discussing generalities about decompositions: the generalized Comon conjecture positing that optimal decompositions with symmetry exist, a review of Strassen's original decomposition of $M_{\langle 2\rangle}$ that hints that this is indeed the case, and defining symmetry groups of decompositions. In particular, I point out that decompositions come in families essentially parametrized by $G_{M_{\langle\mathbf{n}\rangle}}$, and one gains insight studying the entire family rather than individual members. In $\S 4.2$, I describe two decompositions of $M_{\langle\mathbf{n}\rangle}$ that have appeared in the literature, a recent one by Grochow-Moore, and Pan's 1978
decomposition that still holds the world record for practical matrix multiplication in a sense I now make precise.

Introduce $\omega_{\text {prac }, k}$ to be the smallest $\tau$ such that there exists $\mathbf{n} \leq k$ with $\mathbf{R}\left(M_{\langle\mathbf{n}\rangle}\right) \leq \mathbf{n}^{\tau}$. In contrast to the exponent there is no hidden constant. By definition $\omega_{\text {prac }, k} \geq \omega_{\text {prac }, k^{\prime}}$ for all $k^{\prime}>k$ and for all $k, \omega_{p r a c, k}>\omega$. If we decide that we want to multiply unstructured matrices of size, say 10,000 but no larger, then $\omega_{\text {prac, } 10,000}$ will be a more useful quantity than $\omega$. In this regard, the best result is Pan's decomposition (Theorem 4.2.1.1) which implies $\omega_{\text {prac }, 70} \leq 2.79512$. In comparison, using Schönhage's order two border rank 21 decomposition of $M_{\langle 3\rangle}$, converted to a rank decomposition of a $M_{\left\langle 3^{k}\right\rangle}$ (as discussed in $\S 3.2 .1$ ), on needs matrices on the order of $10^{35}$ before one beats Strassen's 2.81 . Using Bini et. al.'s order one border rank 10 decomposition for $M_{\langle 2,2,3\rangle}$ converted to a rank decomposition of $M_{\left\langle 12^{k}\right\rangle}$, one needs matrices of size on the order of $10^{40}$. In order to make e.g., Coppersmith-Winograd's method viable, one needs matrices of size larger than the number of atoms in the known universe (larger than $10^{81}$ ).

Problem 4.0.0.1. Prove upper bounds on $\omega_{\text {prac }, 1,000}$ or $\omega_{\text {prac }, 10,000}$.

This is currently an active area of research.
In §4.3, I revisit Strassen's decomposition and give a proof of Burichenko's theorem [Bur14] that its symmetry group is as large as one could naïvely hope it to be. In order to determine symmetry groups and determine if different decompositions are in the same family, one needs invariants of decompositions. These are studied in $\S 4.4$. Two interesting examples of decompositions of $M_{\langle 3\rangle}$ are given in $\S 4.5$, a variant of Laderman's decomposition and decomposition with $\mathbb{Z}_{4} \times \mathbb{Z}_{3}$-symmetry from [BILR]. In $\S 4.6$ I briefly describe the alternating least squares method that has been used to find decompositions numerically. Border rank decompositions also have geometry associated with them. In order to describe the geometry, I give some geometric preliminaries, including the definition of secant varieties in §4.7. I conclude with two examples of border rank decompositions and their geometry in $\S 4.8$ from [LR0].

### 4.1. Symmetry and decompositions

4.1.1. Warm-up: Strassen's decomposition. Strassen's algorithm, written as a tensor, is

$$
\begin{align*}
& M_{\langle 2\rangle}=\left(\begin{array}{ll}
1 & 0 \\
0 & 1
\end{array}\right)^{\otimes 3}  \tag{4.1.1}\\
& +\left(\begin{array}{ll}
1 & 0 \\
0 & 0
\end{array}\right) \otimes\left(\begin{array}{cc}
0 & 0 \\
1 & -1
\end{array}\right) \otimes\left(\begin{array}{ll}
0 & 1 \\
0 & 1
\end{array}\right)+\left(\begin{array}{cc}
0 & 0 \\
1 & -1
\end{array}\right) \otimes\left(\begin{array}{ll}
0 & 1 \\
0 & 1
\end{array}\right) \otimes\left(\begin{array}{ll}
1 & 0 \\
0 & 0
\end{array}\right)+\left(\begin{array}{ll}
0 & 1 \\
0 & 1
\end{array}\right) \otimes\left(\begin{array}{ll}
1 & 0 \\
0 & 0
\end{array}\right) \otimes\left(\begin{array}{cc}
0 & 0 \\
1 & -1
\end{array}\right) \\
& -\left(\begin{array}{ll}
0 & 0 \\
0 & 1
\end{array}\right) \otimes\left(\begin{array}{cc}
1 & -1 \\
0 & 0
\end{array}\right) \otimes\left(\begin{array}{ll}
1 & 0 \\
1 & 0
\end{array}\right)-\left(\begin{array}{cc}
1 & -1 \\
0 & 0
\end{array}\right) \otimes\left(\begin{array}{ll}
1 & 0 \\
1 & 0
\end{array}\right) \otimes\left(\begin{array}{ll}
0 & 0 \\
0 & 1
\end{array}\right)-\left(\begin{array}{ll}
1 & 0 \\
1 & 0
\end{array}\right) \otimes\left(\begin{array}{ll}
0 & 0 \\
0 & 1
\end{array}\right) \otimes\left(\begin{array}{cc}
1 & -1 \\
0 & 0
\end{array}\right) .
\end{align*}
$$

A first observation is the $\mathbb{Z}_{3}$-symmetry of $M_{\langle 2\rangle}$ (see Exercise 2.5.1.6), which I will call the standard cyclic symmetry, also occurs in Strassen's decomposition: the $\mathbb{Z}_{3}$ action fixes the first term, and permutes the other two triples of terms. This motivates the study of symmetry groups of rank decompositions.
Exercise 4.1.1.1: (2) Show that if we change bases by
$g_{U}=\left(\begin{array}{ll}1 & -1 \\ 0 & -1\end{array}\right) \in G L(U), g_{V}=\left(\begin{array}{ll}-1 & 0 \\ -1 & 1\end{array}\right) \in G L(V), g_{W}=\left(\begin{array}{ll}0 & 1 \\ 1 & 0\end{array}\right) \in G L(W)$,
then the new decomposition of $M_{\langle 2\rangle}$ has four terms fixed by the standard cyclic $\mathbb{Z}_{3}$. ©
4.1.2. Symmetry groups of tensors and their rank decompositions. Consider $\operatorname{Seg}\left(\mathbb{P} A_{1} \times \cdots \times \mathbb{P} A_{d}\right) \subset \mathbb{P}\left(A_{1} \otimes \cdots \otimes A_{d}\right)$. If all the vector spaces have different dimensions, consider the symmetry group of the cone over the Segre as a subgroup of $G L\left(A_{1}\right) \times \cdots \times G L\left(A_{d}\right)$ (more precisely of $G L\left(A_{1}\right) \times$ $\cdots \times G L\left(A_{d}\right) /\left(\mathbb{C}^{*}\right)^{d-1}$, because if $\lambda_{1} \cdots \lambda_{d}=1$, then $\left(\lambda_{1} \operatorname{Id}_{A_{1}}, \ldots, \lambda_{d} \operatorname{Id}_{A_{d}}\right) \in$ $G L\left(A_{1}\right) \times \cdots \times G L\left(A_{d}\right)$ acts trivially). If all dimensions are the same, consider the symmetry group as a subgroup of $\left(G L\left(A_{1}\right) \times \cdots \times G L\left(A_{d}\right) /\left(\mathbb{C}^{*}\right)^{\times d-1}\right) \rtimes$ $\mathfrak{S}_{d}$, where the $\mathfrak{S}_{d}$ acts by permuting the factors after some isomorphism of the $A_{j}$ has been chosen. One can also consider intermediate cases. For $T \in\left(\mathbb{C}^{N}\right)^{\otimes d}$, let

$$
G_{T}:=\left\{g \in\left(G L_{N}^{\times d} /\left(\mathbb{C}^{*}\right)^{\times d-1}\right) \rtimes \mathfrak{S}_{d} \mid g T=T\right\},
$$

and for $T \in A_{1} \otimes \cdots \otimes A_{d}$ with different dimensions, define

$$
G_{T}:=\left\{g \in G L\left(A_{1}\right) \times \cdots \times G L\left(A_{d}\right) /\left(\mathbb{C}^{*}\right)^{\times d-1} \mid g T=T\right\} .
$$

For a polynomial $P \in S^{d} V$, write

$$
G_{P}:=\{g \in G L(V) \mid g P=P\},
$$

For a rank decomposition $T=\sum_{j=1}^{r} t_{j}$, define the set $\mathcal{S}:=\left\{t_{1}, \ldots, t_{r}\right\}$, which I also call the decomposition. If $T$ has a rank decomposition $\mathcal{S}$ and a nontrivial symmetry group $G_{T}$, then given $g \in G_{T}, g \cdot \mathcal{S}:=\left\{g t_{1}, \ldots, g t_{r}\right\}$ is also a rank decomposition of $T$.

Definition 4.1.2.1. The symmetry group of a decomposition $\mathcal{S}$ is $\Gamma_{\mathcal{S}}:=$ $\left\{g \in G_{T} \mid g \cdot \mathcal{S}=\mathcal{S}\right\}$. Let $\Gamma_{\mathcal{S}}^{\prime}=\Gamma_{\mathcal{S}} \cap\left(\Pi_{j} G L\left(A_{j}\right)\right)$.

A guiding principle of this chapter (for which there is no theoretical justification, but holds in several situations, see $\S 7.1 .2$ and $\S 6.6 .3$ ) is that if $T$ has a large symmetry group, then there will exist optimal decompositions of $T$ with symmetry. This even extends to border rank decompositions, as we will see in §4.7.4.

Naïvely, one might think that some decompositions in a family have better symmetry groups than others. Strictly speaking this is not correct:
Proposition 4.1.2.2. [CILO16] For $g \in G_{T}, \Gamma_{g \cdot \mathcal{S}}=g \Gamma_{\mathcal{S}} g^{-1}$.
Proof. Let $h \in \Gamma_{\mathcal{S}}$, then $g h g^{-1}\left(g t_{j}\right)=g\left(h t_{j}\right) \in g \cdot \mathcal{S}$ so $\Gamma_{g \cdot \mathcal{S}} \subseteq g \Gamma_{\mathcal{S}_{t}} g^{-1}$, but the construction is symmetric in $\Gamma_{g \cdot \mathcal{S}}$ and $\Gamma_{\mathcal{S}}$.

As explained below, there may be preferred decompositions in a family where certain symmetries take a particularly transparent form.

For a polynomial $P \in S^{d} V$ and a symmetric rank decomposition $P=$ $\ell_{1}^{d}+\cdots+\ell_{r}^{d}$ for some $\ell_{j} \in V$ (also called a Waring decomposition), and $g \in G_{P} \subset G L(V)$, the same result holds with $\mathcal{S}=\left\{\ell_{1}^{d}, \ldots, \ell_{r}^{d}\right\}$.

In summary, decompositions come in $\operatorname{dim}\left(G_{T}\right)$-dimensional families, and each member of the family has the same abstract symmetry group.
4.1.3. Symmetries of $M_{\langle\mathbf{n}\rangle}$. Let $P G L(U)$ denote $G L(U) / \mathbb{C}^{*}$, where $\mathbb{C}^{*}=$ $\left\{\lambda \operatorname{Id}_{U} \mid \lambda \in \mathbb{C}^{*}\right\}$. This group acts on $\mathbb{P} U$, as well as on $U^{*} \otimes U$. The first action is clear, the second because the action of $G L(U)$ on $\alpha \otimes u$ is $\alpha g^{-1} \otimes g u$ so the scalars times the identity will act trivially.

In $\S 2.5 .1$ we saw that $P G L_{\mathbf{n}}^{\times 3} \rtimes\left(\mathbb{Z}_{3} \rtimes \mathbb{Z}_{2}\right) \subseteq G_{M_{\langle\mathbf{n}\rangle}}$. I emphasize that this $\mathbb{Z}_{2}$ is not contained in either the $\mathfrak{S}_{3}$ permuting the factors or the $P G L(A) \times$ $P G L(B) \times P G L(C)$ acting on them.
Proposition 4.1.3.1. [dG78, Thms. 3.3,3.4] $G_{M_{\langle\mathbf{n}\rangle}}=P G L_{\mathbf{n}}^{\times 3} \rtimes\left(\mathbb{Z}_{3} \rtimes \mathbb{Z}_{2}\right)$.
A proof is given in $\S 8.12 .4$.

### 4.1.4. The Comon conjecture and its generalization.

Conjecture 4.1.4.1 (P. Comon [Com02]). If $T \in S^{d} \mathbb{C}^{N} \subset\left(\mathbb{C}^{N}\right)^{\otimes d}$, then there exists an optimal rank decomposition of $T$ made from symmetric tensors.

After being initially greeted with skepticism by algebraic geometers (Comon is an engineer), the community has now embraced this conjecture and generalized it.

Question 4.1.4.2. [Generalized Comon Conjecture] [BILR] Let $T \in\left(\mathbb{C}^{N}\right)^{\otimes d}$ be invariant under some $\Gamma \subset \mathfrak{S}_{d}$. Does there exist an optimal rank decomposition $\mathcal{S}$ of $T$ satisfying $\Gamma \subseteq \Gamma_{\mathcal{S}}$ ?

I use the following special case as a working hypothesis:
Conjecture 4.1.4.3. [BILR] If $\mathbf{R}\left(M_{\langle\mathbf{n}\rangle}\right)=r$, then there exists a rank $r$ decomposition of $M_{\langle\mathbf{n}\rangle}$ that has standard cyclic symmetry.
4.1.5. Decomposition of $A^{\otimes 3}$ under $\mathbb{Z}_{3}$. In order to search for standard cyclic $\mathbb{Z}_{3}$ decompositions of $M_{\langle\mathbf{n}\rangle}$ we need to understand the $G L(A)$ decomposition of $A^{\otimes 3}$.
Exercise 4.1.5.1: (1!) Verify that the cyclic $\mathbb{Z}_{3}$ acts trivially on both $S^{3} A$ and $\Lambda^{3} A$.

Proposition 4.1.5.2. Let $\mathbb{Z}_{3} \subset \mathfrak{S}_{3}$ act on $A^{\otimes 3}$ by cyclically permuting factors. Then

$$
\left(A^{\otimes 3}\right)^{\mathbb{Z}_{3}}=S^{3} A \oplus \Lambda^{3} A
$$

Proposition 4.1.5.2 is proved in Exercise 8.7.2.4.
Thus if we are searching for cyclic $\mathbb{Z}_{3}$-invariant decompositions for $M_{\langle\mathbf{n}\rangle}$, the size of our search space is cut down from $\mathbf{n}^{6}$ dimensions to $\frac{\mathbf{n}^{6}+2 \mathbf{n}^{2}}{3}$ dimensions.

It is easy to write down the decomposition of $M_{\langle\mathbf{n}\rangle} \in S^{3} A \oplus \Lambda^{3} A$ into its symmetric and skew-symmetric components:

$$
\begin{aligned}
\operatorname{trace}(X Y Z) & =\frac{1}{2}[\operatorname{trace}(X Y Z)+\operatorname{trace}(Y X Z)]+\frac{1}{2}[\operatorname{trace}(X Y Z)-\operatorname{trace}(Y X Z)] \\
& =: M_{\langle\mathbf{n}\rangle}^{S}(X, Y, Z)+M_{\langle\mathbf{n}\rangle}^{\Lambda}(X, Y, Z)
\end{aligned}
$$

Exercise 4.1.5.3: (1) Verify that the first term in brackets lives in $S^{3} A$ and second lives in $\Lambda^{3} A$.

Remark 4.1.5.4. In $\left[\mathbf{C H I}^{+}\right]$we show that the exponent of $M_{\langle\mathbf{n}\rangle}^{S}$ is the same as that of $M_{\langle\mathbf{n}\rangle}$. Since $M_{\langle\mathbf{n}\rangle}^{S}$ is a polynomial, this suggests one can use further tools from algebraic geometry (study of cubic hypersurfaces) in the attempt to determine the exponent.

### 4.2. Two decomposition families of $M_{\langle\mathbf{n}\rangle}$ of rank $<\mathbf{n}^{3}$

Call a subset of points $\left\{\left[a_{1}\right], \ldots,\left[a_{r}\right]\right\}$ of $\mathbb{P} A$ a pinning set if the stabilizer of this set in $P G L(A)$ is finite and no subset of the points has a finite stabilizer.

If we choose vector representatives for the $\left[a_{j}\right]$ call it a framed pinning. For example, if the subset contains a collection of $\mathbf{a}+1$ elements in general linear position, it is a pinning set, and $\mathbf{a}+1$ is the minimal cardinality of a pinning set. Call such a pinning set a standard pinning.

A standard pinning determines $\binom{$ a }{2} points in $\mathbb{P} A^{*}$ obtained by intersecting sets of $\mathbf{a}-2$ hyperplanes coming from the standard pinning points.
4.2.1. Pan's decomposition family. Pan's 1978 decomposition still holds "world record" for practical matrix multiplication.

Let $\mathbf{n}=2 \mathbf{m}$. Let $\mathbb{Z}_{3}$ denote the standard cyclic permutation of factors. Introduce the notation $\overline{\mathrm{I}}=i+\mathbf{m}, \overline{\mathrm{J}}=j+\mathbf{m}$. Write $x_{j}^{i}=u^{i} \otimes v_{j}, y_{j}^{i}=v^{i} \otimes w_{j}$, $z_{j}^{i}=w^{i} \otimes u_{j}$. Let $\mathbb{Z}_{2}^{U}$ be generated by $\sigma_{U}$ which is the exchange $u^{i} \leftrightarrow u^{\bar{i}}$, (which also sends $u_{i} \leftrightarrow u_{\overline{1}}$ ) and define $\mathbb{Z}_{2}^{V}$ and $\mathbb{Z}_{2}^{W}$ similarly, with generators $\sigma_{V}, \sigma_{W}$. Let $\mathbb{Z}_{2}^{\sigma}$ be generated by the product of the generators, so $\sigma$ acts by: $x_{j}^{i} \leftrightarrow x_{\mathrm{J}}^{\overline{\mathrm{I}}}, y_{j}^{i} \leftrightarrow y_{\mathrm{J}}^{\overline{\mathrm{I}}}, z_{j}^{i} \leftrightarrow z_{\mathrm{J}}^{\overline{\mathrm{I}}}$.

Because of the cyclic $\mathbb{Z}_{3}$ symmetry, it will be convenient to identify the three spaces and I will use $x_{j}^{i}$ for all three. In what follows, indices are to be considered $\bmod \mathbf{n}$.

For a finite group $\Gamma \subset G L_{N} \times G L_{N} \times G L_{N} \rtimes \mathfrak{S}_{3}$, introduce the notation

$$
\begin{equation*}
\langle x \otimes y \otimes z\rangle_{\Gamma}:=\sum_{g \in \Gamma} g \cdot(x \otimes y \otimes z) . \tag{4.2.1}
\end{equation*}
$$

Let $\mathbb{Z}_{2}^{\tau}$ denote the standard transpose $x \otimes y \otimes z \mapsto y^{T} \otimes x^{T} \otimes z^{T}$ and let $\mathbb{Z}_{2}^{\tau^{\prime}}$ denote the transpose-like symmetry obtained by composing the standard transpose symmetry with $\sigma_{V}$.

Theorem 4.2.1.1. [Pan78] With notations as above $M_{\langle\mathbf{n}\rangle}$ equals

$$
\begin{equation*}
\sum_{(i, j, k) \mid 0 \leq i \leq j<k \leq \mathbf{m}-1}\left\langle\left(x_{j}^{i}+x_{k}^{j}+x_{i}^{k}\right)^{\otimes 3}\right\rangle_{\mathbb{Z}_{2}^{\sigma}} \tag{4.2.2}
\end{equation*}
$$

$$
\begin{align*}
& -\sum_{(i, j, k) \mid 0 \leq i \leq j<k \leq \mathbf{m}-1}\left\langle\left(x_{j}^{i}-x_{k}^{\bar{j}}+x_{\overline{1}}^{k}\right) \otimes\left(-x_{j}^{\overline{1}}+x_{\bar{k}}^{j}+x_{i}^{k}\right) \otimes\left(x_{\bar{J}}^{i}+x_{k}^{j}-x_{i}^{\bar{k}}\right)\right\rangle_{\mathbb{Z}_{2}^{\sigma} \times \mathbb{Z}_{3}}  \tag{4.2.3}\\
& (4.2 .4) \\
& +\sum_{i, j=0}^{\mathbf{m}-1}\left\langle x_{j}^{i} \otimes x_{j}^{i} \otimes\left[\left(\mathbf{m}-\delta_{i j}\right) x_{j}^{i}+\sum_{k=0}^{\mathbf{m}-1}\left[\left(x_{i}^{k}+x_{k}^{j}\right)-\delta_{i j}\left(x_{i}^{k}+x_{k}^{i}\right)\right]\right]\right\rangle_{\mathbb{Z}_{2}^{\sigma} \times \mathbb{Z}_{3} \times \mathbb{Z}_{2}^{r_{2}^{\prime}}} \\
& (4.2 .5)  \tag{4.2.5}\\
& +\sum_{i, j=0}^{\mathbf{m}-1}\left\langle x_{j}^{i} \otimes x_{j}^{\overline{1}} \otimes\left[\left(\mathbf{m}-\delta_{i j}\right) x_{\overline{\mathrm{J}}}^{i}+\sum_{k=0}^{\mathbf{m}-1}\left[\left(x_{i}^{\bar{k}}+x_{k}^{j}\right)-\delta_{i j}\left(x_{i}^{\bar{k}}+x_{k}^{i}\right)\right]\right]\right\rangle_{\mathbb{Z}_{2}^{\sigma} \times \mathbb{Z}_{3} \times \mathbb{Z}_{2}^{\tau}} .
\end{align*}
$$

Note that the terms (4.2.2),(4.2.4),(4.2.5) are $M_{\langle\mathbf{n}\rangle}$ plus "garbage" terms. The second summation eliminates the garbage terms. Call the decomposition $\mathcal{S}_{\text {Pan }}$.

Remark 4.2.1.2. According to Burichenko (announced in [Bur15, Thm 1.1]) $\Gamma_{\mathcal{S}_{\text {Pan }}}=\mathfrak{S}_{\mathbf{m}} \times \mathbb{Z}_{2} \times \mathfrak{S}_{3}$.

Exercise 4.2.1.3: (2) Show that the number of triples $(i, j, k)$ with $0 \leq i \leq$ $j<k \leq \mathbf{m}-1$ is $\frac{2}{3}\left(\mathbf{m}^{3}-\mathbf{m}\right)$ and conclude that Pan's decomposition is of rank $\frac{1}{3} \mathbf{n}^{3}+6 \mathbf{n}^{2}-\frac{4}{3} \mathbf{n}$.
Exercise 4.2.1.4: (1) Show that when $\mathbf{n}=70$, Pan's decomposition has rank 143, 240 and conclude that $\omega_{\text {prac, } 70} \leq 2.79512$.
4.2.2. The Grochow-Moore decompositions. The group $\mathfrak{S}_{\mathbf{n}+1}$ acts irreducibly on $\mathbb{C}^{\mathbf{n}}$ (see $\S 1.1 .13$ for the action and $\S 8.7 .2$ for the proof), and the induced action on $\mathbb{C}^{\mathbf{n} *} \otimes \mathbb{C}^{\mathbf{n}}$ has a unique trivial representation, namely $\mathrm{Id}_{\mathbb{C}^{\mathrm{n}}}$, see Exercise 8.6.8.3.
Exercise 4.2.2.1: (1) Show any $T \in\left(U^{*} \otimes V\right) \otimes\left(V^{*} \otimes W\right) \otimes\left(W^{*} \otimes U\right)$ that is acted on trivially by $\mathfrak{S}_{\mathbf{n}+1}^{\times 3}$, where the first copy acts on $U, U^{*}$ the second on $V, V^{*}$ and the third on $W, W^{*}$, is up to scale $M_{\langle\mathbf{n}\rangle}$.

Let $u^{1}, \ldots, u^{\mathbf{n + 1}} \in U^{*}$ be a framed pinning normalized so that $u^{1}+$ $\cdots+u^{\mathbf{n + 1}}=0$, with $\binom{\mathbf{n}+1}{2}$ induced points $u_{i j}:=u_{[\mathbf{n}+1] \backslash\{i, j\}}$ for $i<j$, with normalizations: $\sum_{j} u^{j}=0$, and $u^{i}\left(u_{i k}\right)=1, u^{i}\left(u_{k i}\right)=-1$. Adopt the notation $u_{j i}:=-u_{i j}$, so $u_{i i}=0$.

Given a framed pinning, define a dual framed pinning $u_{1}, \ldots, u_{\mathbf{n}+1}$ of $U$ by requiring

$$
u^{i}\left(u_{j}\right)=\left\{\begin{array}{cc}
1 & i=j \\
-\frac{1}{\mathbf{n}} & i \neq j
\end{array}\right\} .
$$

Exercise 4.2.2.2: (1) When $\mathbf{n}=2$ compute the dual pinning to $\binom{1}{0},\binom{0}{1},\binom{-1}{-1}$. ©

Exercise 4.2.2.3: (1) Show that with the above normalizations $u_{i j}=u_{i}-$ $u_{j}$.
Exercise 4.2.2.4: (1) Show that with the above normalizations $\operatorname{Id}_{U}=$ $\frac{\mathbf{n}}{\mathbf{n}+1} \sum_{i=1}^{\mathbf{n + 1}} u^{i} \otimes u_{i}$.
Proposition 4.2.2.5. [GM16] Notations as above. Then

$$
M_{\langle\mathbf{n}\rangle}=\left(\frac{\mathbf{n}}{\mathbf{n}+1}\right)^{3} \sum_{i, j, k=1}^{\mathbf{n}+1} u^{i} v_{j} \otimes v^{j} w_{k} \otimes w^{k} u_{i}
$$

Proof. Note that the right hand side is invariant under $\mathfrak{S}_{\mathbf{n}+1}^{\times 3}$ so it is some constant times $M_{\langle\mathbf{n}\rangle}$. To check the constant is correct, evaluate the right hand side on, e.g., $\mathrm{Id}_{\mathbf{n}}^{\otimes 3}$.

Proposition 4.2.2.5 gives a rank $(\mathbf{n}+1)^{3}$ decomposition of $M_{\langle\mathbf{n}\rangle}$, so at first glance it does not appear interesting. However, it is used to prove the following theorem:
Theorem 4.2.2.6. [GM16] Let $u^{1}, \ldots, u^{\mathbf{n}+1} \in U^{*}$ be a framed pinning with induced vectors $u_{i j} \in U$ as above, and choose identifications $U \simeq V \simeq$ $W$ to obtain inherited pinnings and induced vectors. The following is a rank $\mathbf{n}^{3}-\mathbf{n}+1$ decomposition of $M_{\langle\mathbf{n}\rangle}$, call it $\mathcal{S}_{G M}$, with $\Gamma_{\mathcal{S}_{G M}} \supset \mathfrak{S}_{\mathbf{n}+1} \rtimes \mathbb{Z}_{3}$.

$$
M_{\langle\mathbf{n}\rangle}=\operatorname{Id}_{\mathbf{n}}^{\otimes 3}-\left(\frac{\mathbf{n}}{\mathbf{n}+1}\right)^{3} \sum_{i, j, k \in[\mathbf{n}+1]} \text { and distinct } u^{i} v_{i j} \otimes v^{j} w_{j k} \otimes w^{k} u_{k i}
$$

Proof. First, notice that

$$
\sum_{i, j, k \in[\mathbf{n}+1] \text { and distinct }} u^{i} v_{i j} \otimes v^{j} w_{j k} \otimes w^{k} u_{k i}=\sum_{i, j, k \in[\mathbf{n}+1]} u^{i} v_{i j} \otimes v^{j} w_{j k} \otimes w^{k} u_{k i}
$$

because $v_{i i}=0$. By Exercise 4.2.2.3 we may write $v_{i j}=v_{i}-v_{j}$ One then expands out, using Exercise 4.2.2.4 and Proposition 4.2.2.5 to conclude.

Theorem 4.2.2.6 gives another perspective on Strassen's decomposition family for $M_{\langle 2\rangle}$.

### 4.3. Strassen's decomposition revisited

Let $\mathcal{S t r}$ denote the Strassen decomposition of $M_{\langle 2\rangle}$.
4.3.1. The Strassen family. As discussed above, decompositions are best studied in families. In the case of $M_{\langle 2\rangle}$, there is a unique family:
Theorem 4.3.1.1. [dG78] The set of rank seven decompositions of $M_{\langle 2\rangle}$ is the orbit $G_{M_{\langle 2\rangle}} \cdot \mathcal{S t r}$.

The proof follows from a careful analysis of every possible decomposition, taking into account that an element $a \otimes b \otimes c$ is not just a triple of vectors, but a triple of endomorphisms $\mathbb{C}^{2} \rightarrow \mathbb{C}^{2}$, and the analysis is via the possible triples of ranks that can appear.

In preparation for studying the Strassen family of decompositions, write

$$
\begin{equation*}
u_{1}=\binom{1}{0}, u_{2}=\binom{0}{1}, u^{1}=(1,0), u^{2}=(0,1) \tag{4.3.1}
\end{equation*}
$$

and set $v_{j}=w_{j}=u_{j}$ and $v^{j}=w^{j}=u^{j}$.
Strassen's decomposition becomes

$$
\begin{align*}
M_{\langle 2\rangle}= & \left(v_{1} u^{1}+v_{2} u^{2}\right) \otimes\left(w_{1} v^{1}+w_{2} v^{2}\right) \otimes\left(u_{1} w^{1}+u_{2} w^{2}\right)  \tag{4.3.2}\\
& +\left\langle v_{1} u^{1} \otimes w_{2}\left(v^{1}-v^{2}\right) \otimes\left(u_{1}+u_{2}\right) w^{2}\right\rangle_{\mathbb{Z}_{3}} \\
& +\left\langle v_{2} u^{2} \otimes w_{1}\left(v^{2}-v^{1}\right) \otimes\left(u_{1}+u_{2}\right) w^{1}\right\rangle_{\mathbb{Z}_{3}} .
\end{align*}
$$

From this presentation we transparently recover much of the entire Strassen family, namely by letting $u_{1}, u_{2}, v_{1}, v_{2}$, and $w_{1}, w_{2}$ be arbitrary bases, with dual basis vectors denoted with superscripts. We obtain a family parametrized by $P G L(U) \times P G L(V) \times P G L(W)$, and since the decomposition (4.3.2) is manifestly $\mathbb{Z}_{3}$-invariant, the only potential additional decompositions arise from applying a transpose symmetry such as $x \otimes y \otimes z \mapsto x^{T} \otimes z^{T} \otimes y^{T}$. Call such a transpose symmetry convenient.
Exercise 4.3.1.2: (1) Show that if we set $u_{3}=\binom{-1}{-1}$ and $u^{3}=(1,-1)$ and similarly for $v, w$, then the matrices in Exercise 4.1.1.1 respectively correspond to the permutations $(2,3),(1,3)$ and $(1,2)$. The matrix in the first term of the decomposition that one obtains from Exercise 4.1.1.1 also corresponds to a permutation. Which one?

Exercise 4.3.1.3: (2) Find a change of basis such that the first term in the decomposition of Exercise 4.1.1.1 becomes $\left(\begin{array}{cc}\omega & 0 \\ 0 & \omega^{2}\end{array}\right)^{\otimes 3}$ where $\omega=e^{\frac{2 \pi i}{3}}$ and write out the decomposition in this basis.

Under $x \otimes y \otimes z \mapsto x^{T} \otimes z^{T} \otimes y^{T}$, Strassen's decomposition is mapped to:

$$
\begin{align*}
M_{\langle 2\rangle}= & \left(\begin{array}{ll}
1 & 0 \\
0 & 1
\end{array}\right)^{\otimes 3}  \tag{4.3.3}\\
& +\left\langle\left(\begin{array}{ll}
1 & 0 \\
0 & 0
\end{array}\right) \otimes\left(\begin{array}{ll}
0 & 0 \\
1 & 1
\end{array}\right) \otimes\left(\begin{array}{cc}
0 & 1 \\
0 & -1
\end{array}\right)\right\rangle_{\mathbb{Z}_{3}} \\
& -\left\langle\left(\begin{array}{ll}
0 & 0 \\
0 & 1
\end{array}\right) \otimes\left(\begin{array}{ll}
1 & 1 \\
0 & 0
\end{array}\right) \otimes\left(\begin{array}{cc}
1 & 0 \\
-1 & 0
\end{array}\right)\right\rangle_{\mathbb{Z}_{3}} .
\end{align*}
$$

Notice that this is almost Strassen's decomposition (4.1.1)- just some the signs are wrong. We can "fix" the problem by conjugating all the matrices with

$$
g_{0}:=\left(\begin{array}{cc}
0 & -1 \\
1 & 0
\end{array}\right)
$$

Exercise 4.3.1.4: (1) Verify that acting by $g_{0}^{\times 3} \in P G L(U) \times P G L(V) \times$ $P G L(W)$ takes (4.3.3) to Strassen's decomposition.

Exercise 4.3.1.4 shows that there is a non-standard $\mathbb{Z}_{2} \subset P G L_{2}^{\times 3} \rtimes\left(\mathbb{Z}_{3} \rtimes\right.$ $\mathbb{Z}_{2}$ ) contained in $\Gamma_{\mathcal{S t r}}$, namely the convenient transpose symmetry composed with $g_{0}^{\times 3}$. It also implies a refinement of deGroote's theorem:
Proposition 4.3.1.5. [Bur14, CILO16] The set of rank seven decompositions of $M_{\langle 2\rangle}$ is $P G L_{2}^{\times 3} \cdot \mathcal{S t r}$.

With the expression (4.3.2), notice that if we exchange $u_{1} \leftrightarrow u_{2}$ and $u^{1} \leftrightarrow u^{2}$, the decomposition is also preserved by this $\mathbb{Z}_{2} \subset P G L_{2}^{\times 3}$, with orbits (4.3.2) and the exchange of the triples. So we see $\Gamma_{\mathcal{S} t r} \supseteq \mathbb{Z}_{2} \rtimes\left(\mathbb{Z}_{3} \rtimes \mathbb{Z}_{2}\right)$, where the first $\mathbb{Z}_{2}$ is diagonally embedded in $P G L_{2}^{\times 3}$.

Although the above description of the Strassen family of decompositions for $M_{\langle 2\rangle}$ is satisfying, it becomes even more transparent with a projective perspective. With the projective perspective, we will see that $\Gamma_{\mathcal{S} t r}$ is even larger.
4.3.2. $M_{\langle 2\rangle}$ viewed projectively. That all rank 7 decompositions of $M_{\langle 2\rangle}$ are obtained via $P G L_{2}^{\times 3}$ suggests using a projective perspective. The group $P G L_{2}$ acts simply transitively on triples of distinct points of $\mathbb{P}^{1}$. So to fix a decomposition in the family, select a pinning (triple of points) in each space. I focus on $\mathbb{P} U$. Call the points $\left[u_{1}\right],\left[u_{2}\right],\left[u_{3}\right]$. Then these determine three points in $\mathbb{P} U^{*},\left[u^{1 \perp}\right],\left[u^{2 \perp}\right],\left[u^{3 \perp}\right]$. Choose representatives $u_{1}, u_{2}, u_{3}$ satisfying $u_{1}+u_{2}+u_{3}=0$. I could have taken any linear relation, it just would introduce coefficients in the decomposition. I take the most symmetric relation to keep all three points on an equal footing. Similarly, fix the scales
on the $u^{j \perp}$ by requiring $u^{j \perp}\left(u_{j-1}\right)=1$ and $u^{j \perp}\left(u_{j+1}\right)=-1$, where indices are considered $\bmod \mathbb{Z}_{3}$, so $u_{3+1}=u_{1}$ and $u_{1-1}=u_{3}$.

In comparison with what we had before, letting the old vectors be hatted, $\hat{u}_{1}=u_{1}, \hat{u}_{2}=u_{2}, \hat{u}^{1}=u^{2 \perp}$, and $\hat{u}^{2}=-u^{1 \perp}$. The effect is to make the symmetries of the decomposition more transparent. Our identifications of the ordered triples $\left\{u_{1}, u_{2}, u_{3}\right\}$ and $\left\{v_{1}, v_{2}, v_{3}\right\}$ determines a linear isomorphism $a_{0}: U \rightarrow V$, and similarly for the other pairs of vector spaces. Note that $a_{0}=v_{j} \otimes u^{j+1 \perp}+v_{j+1} \otimes u^{j+2 \perp}$ for any $j=1,2,3$.

Then

$$
\begin{align*}
M_{\langle 2\rangle} & =a_{0} \otimes b_{0} \otimes c_{0}  \tag{4.3.4}\\
& +\left\langle\left(v_{1} u^{2 \perp}\right) \otimes\left(w_{3} v^{1 \perp}\right) \otimes\left(u_{2} w^{3 \perp}\right)\right\rangle_{\mathbb{Z}_{3}} \\
& +\left\langle\left(v_{1} u^{3 \perp}\right) \otimes\left(w_{2} v^{1 \perp}\right) \otimes\left(u_{3} w^{2 \perp}\right)\right\rangle_{\mathbb{Z}_{3}} .
\end{align*}
$$

Here, to make the terms shifted by $\mathbb{Z}_{3}$ live in the proper space, one must act by $a_{0}, b_{0}, c_{0}$ appropriately, e.g., to shift $v_{1} u^{2 \perp}$ to the second slot, one takes $b_{0} v_{1} u^{2 \perp} a_{0}{ }^{-1}$.

With this presentation, taking $a_{0}=b_{0}=c_{0}=\mathrm{Id}$, the diagonally embedded $\mathfrak{S}_{3} \subset P G L_{2}^{\times 3}$ acting by permuting the indices transparently preserves the decomposition, with two orbits, the fixed point $a_{0} \otimes b_{0} \otimes c_{0}$ and the orbit of $\left(v_{1} u^{2 \perp}\right) \otimes\left(w_{3} v^{1 \perp}\right) \otimes\left(u_{2} w^{3 \perp}\right)$. The action on each of $U, V, W$ is the standard irreducible two dimensional representation.

We now see $\Gamma_{\mathcal{S} t r} \supseteq \mathfrak{S}_{3} \rtimes\left(\mathbb{Z}_{3} \rtimes \mathbb{Z}_{2}\right)$, with $\mathfrak{S}_{3} \subset \Gamma_{\mathcal{S t r}}^{\prime}$. With a little more work, one sees that equality holds:
Theorem 4.3.2.1. [Bur14] The symmetry group $\Gamma_{\mathcal{S t r}}$ of Strassen's decomposition of $M_{\langle 2\rangle}$ is $\left(\mathfrak{S}_{3} \times \mathbb{Z}_{3}\right) \rtimes \mathbb{Z}_{2} \subset P G L_{2}^{\times 3} \rtimes\left(\mathbb{Z}_{3} \rtimes \mathbb{Z}_{2}\right)=G_{M_{\langle 2\rangle}}$.

Remark 4.3.2.2. One can prove Strassen's decomposition is indeed matrix multiplication simply by the group invariance, see [CILO16].

### 4.4. Invariants associated to a decomposition of $M_{\langle\mathbf{n}\rangle}$

Given two decompositions of $M_{\langle\mathbf{n}\rangle}$, how can we determine if they are in the same family? Given one, how can we determine its symmetry group? These questions are related, as a necessary condition for two decompositions to be in the same family is that they have isomorphic symmetry groups. I first define invariants $\mathcal{S}_{s, t, u}$ that are subsets of points in $\mathbb{P}(A \otimes B \otimes C)$. Keeping track of the cardinalities of these sets dates at least back to [JM86]. I then further define subsets $\mathcal{S}_{U} \subset \mathbb{P} U, \mathcal{S}_{U^{*}} \subset \mathbb{P} U^{*}$ that give more information. I describe further invariants associated to a decomposition via graphs. I
then discuss the sets $\mathcal{S}_{U}, \mathcal{S}_{U}^{*}$ in more detail: it turns out that the collection of points themselves has geometry that is also useful for distinguishing decompositions and determining symmetry groups.
4.4.1. Invariants of decompositions of $M_{\langle\mathbf{n}\rangle}$. Let $M_{\langle\mathbf{n}\rangle}=\sum_{j=1}^{r} t_{j}$ be a rank decomposition for $M_{\langle\mathbf{n}\rangle}$ and write $t_{j}=a_{j} \otimes b_{j} \otimes c_{j}$. Let $\mathbf{r}_{j}:=\left(\operatorname{rank}\left(a_{j}\right), \operatorname{rank}\left(b_{j}\right), \operatorname{rank}\left(c_{j}\right)\right)$, and let $\tilde{\mathbf{r}}_{j}$ denote the unordered triple. The following proposition is clear:
Proposition 4.4.1.1. [BILR] Let $\mathcal{S}$ be a rank decomposition of $M_{\langle\mathbf{n}\rangle}$. Partition $\mathcal{S}$ by unordered rank triples into disjoint subsets: $\left\{\tilde{\mathcal{S}}_{1,1,1}, \tilde{\mathcal{S}}_{1,1,2}, \ldots, \tilde{\mathcal{S}}_{n, n, n}\right\}$. Write the corresponding ordered triplets as $\left\{\mathcal{S}_{1,1,1}, \mathcal{S}_{1,1,2}, \mathcal{S}_{1,2,1}, \ldots, \mathcal{S}_{n, n, n}\right\}$. Then $\Gamma_{\mathcal{S}}$ preserves each $\tilde{\mathcal{S}}_{s, t, u}$ and $\Gamma_{\mathcal{S}}^{\prime}$ preserves each $\mathcal{S}_{s, t, u}$.

We can say more about rank one elements: If $a \in U^{*} \otimes V$ and $\operatorname{rank}(a)=$ 1 , then there are unique points $[\mu] \in \mathbb{P} U^{*}$ and $[v] \in \mathbb{P} V$ such that $[a]=[\mu \otimes v]$. So given a decomposition $\mathcal{S}$ of $M_{\langle\mathbf{n}\rangle}$, define $\mathcal{S}_{U^{*}} \subset \mathbb{P} U^{*}$ and $\mathcal{S}_{U} \subset \mathbb{P} U$ to correspond to the $U^{*}$ and $U$ elements appearing in $\mathcal{S}_{1,1,1}$. Then $\Gamma_{\mathcal{S}}^{\prime}$ preserves the sets $\mathcal{S}_{U}$ and $\mathcal{S}_{U^{*}}$ up to projective equivalence.

I will say a decomposition has a transpose-like $\mathbb{Z}_{2}$ invariance if it is invariant under a $\mathbb{Z}_{2}$ such as $x \otimes y \otimes z \mapsto x^{T} \otimes z^{T} \otimes y^{T}$ composed with an element of $P G L(U) \times P G L(V) \times P G L(W)$.
Exercise 4.4.1.2: (1) Show that if a decomposition of $M_{\langle\mathbf{n}\rangle}$ is cyclic $\mathbb{Z}_{3^{-}}$ invariant and also has a transpose-like $\mathbb{Z}_{2}$-invariance, then $\mathcal{S}_{U}$ and $\mathcal{S}_{U^{*}}$ have the same cardinality.
4.4.2. A graph. Define a bipartite graph $\mathcal{I G}_{\mathcal{S}}$, the incidence graph where the top vertex set is given by elements in $\mathcal{S}_{U^{*}}$ and the bottom vertex set by elements in $\mathcal{S}_{U}$. Draw an edge between elements $[\mu]$ and $[v]$ if they are incident, i.e., $\mu(v)=0$. Geometrically, $[v]$ belongs to the hyperplane determined by $[\mu]$ (and vice-versa). One can weight the vertices of this graph in several ways, the simplest (and in practice this has been enough) is just by the number of times the element appears in the decomposition. Let $\Gamma_{\mathcal{I} \mathcal{G}_{\mathcal{S}}} \subset G_{M_{\langle\mathbf{n}\rangle}}$ denote the automorphism group of $\mathcal{I}_{\mathcal{S}}$, so $\Gamma_{\mathcal{S}}^{\prime} \subseteq \Gamma_{\mathcal{I} \mathcal{G}_{\mathcal{S}}}$, and if we take the triple of incidence graphs, we get a similar inclusion for $\Gamma_{\mathcal{S}}$. See the examples in $\S 4.5 .1$ and $\S 4.5 .2$.

If a decomposition is $\mathbb{Z}_{3}$ invariant, the incidence graphs form $V, V^{*}$ and from $W, W^{*}$ are isomorphic, and otherwise they give additional information.

Given a $\mathbb{Z}_{3}$-invariant decomposition, a necessary condition for it to also have a transpose-like $\mathbb{Z}_{2}$ symmetry is that there is an isomorphism of the bipartite graph swapping the sets of (weighted) vertices.

In practice (see the examples below) the incidence graph has been enough to determine the symmetry group $\Gamma_{\mathcal{S}}$, in the sense that it cuts the possible
size of the group down and it becomes straight-forward to determine $\Gamma_{\mathcal{S}}$ from $\Gamma_{\mathcal{I} \mathcal{G}_{S}}$.

Remark 4.4.2.1. In [BILR] a second graph, called the pairing graph is defined that gives further information about $\Gamma_{\mathcal{S}}^{\prime}$.
4.4.3. Configurations of points in projective space. In practice, perhaps because of the numerical methods used, the sets $\mathcal{S}_{U}$, and $\mathcal{S}_{U^{*}}$ have been relatively small. It is not surprising that they each are spanning sets. Usually they have come from configurations in a sense I now describe. For $\mathbb{P}^{1}$, a configuration is simply a triple of points and the triple of points they determine in the dual vector space. For example Strassen's decomposition is built from a configuration. The higher dimensional analog of such pairs of triples is more complicated.

I emphasize that the decompositions of [BILR] were found by numerical searches, without distinguishing any configurations. However in most cases, we were able to give a simple description of the vectors appearing in the decomposition in terms of a configuration. This bodes well for future work.

I restrict the discussion to $\mathbb{P}^{2}$, see $[\mathbf{B I L R}]$ for the general case. The group $P G L_{3}$ acts simply transitively on the set of 4-ples of points in general linear position (i.e., such that any three of them span $\mathbb{P}^{2}$ ).

Start with any 4 -ple of points in general linear position. In the decomposition, actual vectors will appear. Even in the decomposition, since what will appear are vectors tensored with each other, there is only a "global scale" for each term. Take the simplest (to write down) 4-ple, choosing the fourth vector in order to have the linear relation $u_{1}+u_{1}+u_{3}+u_{4}=0$. I'll call this the default configuration. That is, the default configuration starts with

$$
u_{1}=\left(\begin{array}{l}
1 \\
0 \\
0
\end{array}\right), u_{2}=\left(\begin{array}{l}
0 \\
1 \\
0
\end{array}\right), u_{3}=\left(\begin{array}{l}
0 \\
0 \\
1
\end{array}\right), u_{4}=\left(\begin{array}{l}
-1 \\
-1 \\
-1
\end{array}\right) .
$$

The $\left\{\left[u_{j}\right]\right\}$ determine points in the dual space by taking pairwise intersections of the lines (hyperplanes) that they determine in $\mathbb{P} U^{*}$.

$$
\begin{aligned}
& v_{12}=(0,0,1), \quad v_{13}=(0,1,0), \quad v_{14}=(0,1,-1), \\
& v_{23}=(-1,0,0), \quad v_{24}=(-1,0,1), \quad v_{34}=(1,-1,0) .
\end{aligned}
$$

Here $\left[v_{i j}\right]$ is the line in $\mathbb{P}^{2}$ (considered as a point in the dual space $\mathbb{P}^{2 *}$ ) through the points $\left[u_{i}\right]$ and $\left[u_{j}\right]$ in $\mathbb{P}^{2}$ (or dually, the point of intersection of the two lines $\left[u_{i}\right],\left[u_{j}\right]$ in $\mathbb{P}^{2 *}$ ). Here choices of representatives are being made. I have made choices that will be useful for the decomposition $\mathcal{S}_{B I L R, \mathbb{Z}_{4} \times \mathbb{Z}_{3}}$ of $\S 4.5 .1$ below.

The $v_{i, j}$ in turn determine new points of intersection:

$$
u_{12,34}=\left(\begin{array}{l}
1 \\
1 \\
0
\end{array}\right), u_{13,24}=\left(\begin{array}{l}
1 \\
0 \\
1
\end{array}\right), u_{14,23}=\left(\begin{array}{l}
0 \\
1 \\
1
\end{array}\right) .
$$

which determine new points
$v_{(12,34),(13,24)}=(-1,1,1), v_{(12,34),(14,23)}=(1,-1,1), v_{(13,24),(14,23)}=(1,1,-1)$,
which determine new points in $U$, etc.., see [BILR] for details. In practice only vectors from the first three sets of a configuration ( 7 for $U, 6$ for $V$ or vice-versa) have been useful.

### 4.5. Cyclic $\mathbb{Z}_{3}$-invariant rank 23 decompositions of $M_{\langle 3\rangle}$

In [BILR] five new standard cyclic families of decompositions were found, as well as a standard cyclic variant of Laderman's decomposition. What follows is one of the new decompositions and the standard cyclic variant of Laderman's decomposition.
4.5.1. A rank 23 decomposition of $M_{\langle 3\rangle}$ with $\mathbb{Z}_{4} \times \mathbb{Z}_{3}$ symmetry. Take a configuration and let $a_{0}: U \rightarrow V$ send $u_{j}$ to $v_{j+1}$. In the default configuration

$$
a_{0}=\left(\begin{array}{lll}
0 & 0 & -1 \\
1 & 0 & -1 \\
0 & 1 & -1
\end{array}\right)
$$

corresponds to the generator of $\mathbb{Z}_{4}$ that cyclically permutes indices.
Theorem 4.5.1.1. [BILR] Let $u_{i j}, v_{i}, v_{i j \mid k l}$ be as in $\S$ 4.4.3. Then

$$
\begin{aligned}
M_{\langle 3\rangle}= & -a_{0}^{\otimes 3} \\
& \left\langle\left(u_{24} v_{12 \mid 34}\right)^{\otimes 3}\right\rangle_{\mathbb{Z}_{2} \subset \mathbb{Z}_{4}} \\
& \left\langle-\left[u_{24} v_{4}+u_{12} v_{3}\right]^{\otimes 3}\right\rangle_{\mathbb{Z}_{4}} \\
& \left\langle\left(u_{12} v_{3}\right)^{\otimes 3}\right\rangle_{\mathbb{Z}_{4}} \\
& \left\langle\left(u_{12} v_{1}\right) \otimes\left(u_{23} v_{3}\right) \otimes\left(u_{24} v_{4}\right)\right\rangle_{\mathbb{Z}_{4} \times \mathbb{Z}_{3}} .
\end{aligned}
$$

Here is the incidence graph:


Given the distribution of the frequencies of the points: $(4,4,4,4,1,1)$ in $V,(3,3,3,3,3,3)$ in $U^{*}$, a transpose-like symmetry is not possible. Moreover, it is clear one cannot upgrade the $\mathbb{Z}_{4}$ to $\mathfrak{S}_{4}$ since only two of the three $v_{i j \mid k l}$ appear in the decomposition: $v_{12 \mid 34}, v_{14 \mid 23}\left(v_{13 \mid 24}\right.$ is omitted). So, e.g., the transposition $(2,3)$ takes $\mathcal{S}_{B I L R, \mathbb{Z}_{4} \times \mathbb{Z}_{3}}$ to a different decomposition in the family.
Proposition 4.5.1.2. [BILR] $\Gamma_{\mathcal{S}_{B I L R, \mathbb{Z}_{4} \times \mathbb{Z}_{3}}}=\mathbb{Z}_{4} \times \mathbb{Z}_{3}$
Exercise 4.5.1.3: (2) Use the incidence graph to prove Proposition 4.5.1.2.
4.5.2. Laderman's decomposition. I now discuss a variant of Laderman's rank 23 decomposition of $M_{\langle 3\rangle}$, which I denote $\mathcal{L} a d$. According to Burichenko [Bur15], one has a $\mathbb{Z}_{2} \times \mathbb{Z}_{2} \subset P G L(U) \times P G L(V) \times P G L(W)$ contained in $\Gamma_{\mathcal{L a d}}$ and the full cyclic permutation and a transpose-like $\mathbb{Z}_{3} \rtimes \mathbb{Z}_{2}$ also in $\Gamma_{\mathcal{L a d}}$, acting in a twisted way. Thanks to the transpose-like symmetry, it is better to label points in the dual space by their image under the transpose-like symmetry rather than annihilators, to make the symmetry more transparent. Here it is:

Points:

$$
\begin{gathered}
u_{1}=\left(\begin{array}{l}
1 \\
0 \\
0
\end{array}\right), u_{2}=\left(\begin{array}{l}
0 \\
1 \\
0
\end{array}\right), u_{3}=\left(\begin{array}{l}
0 \\
0 \\
1
\end{array}\right), u_{12}=\left(\begin{array}{c}
1 \\
-1 \\
0
\end{array}\right), u_{23}=\left(\begin{array}{c}
0 \\
1 \\
-1
\end{array}\right) . \\
v_{1}=(1,0,0), v_{2}=(0,1,0), v_{3}=(0,0,1), \\
v_{12}=(1,1,0), v_{23}=(0,1,1) .
\end{gathered}
$$

Note that the configuration of points in $\mathbb{P} U$ is
Exercise 4.5.2.1: (1) Determine the subgroup of $P G L_{3}$ fixing the configuration of two lines in the plane. ©


Figure 4.5.1. Configuration from the symmetric Laderman decomposition

Exercise 4.5.2.2: (2) What is the subgroup of your answer to Exercise 4.5.2.1 that preserves the full configuration in $\mathbb{P}^{2}$ (i.e., two lines, intersecting in a point, each with two additional marked points).
Theorem 4.5.2.3. [BILR, Lad76] Notations as above. Then

$$
\begin{align*}
M_{\langle 3\rangle}= & \left(u_{2} v_{2}\right)^{\otimes 3}  \tag{4.5.1}\\
& \left(u_{3} v_{3}\right)^{\otimes 3}  \tag{4.5.2}\\
& \left(u_{12} v_{1}\right)^{\otimes 3}  \tag{4.5.3}\\
& \left(u_{1} v_{12}\right)^{\otimes 3}  \tag{4.5.4}\\
& \left(u_{2} v_{1}-u_{1} v_{12}\right)^{\otimes 3}  \tag{4.5.5}\\
& \left\langle\left(u_{1} v_{3}\right) \otimes\left(u_{3} v_{1}\right) \otimes\left(u_{1} v_{1}\right)\right\rangle_{\mathbb{Z}_{3}}  \tag{4.5.6}\\
& \left\langle\left(u_{23} v_{1}\right) \otimes\left(u_{12} v_{3}\right) \otimes\left(u_{23} v_{3}\right)\right\rangle_{\mathbb{Z}_{3}}  \tag{4.5.7}\\
& \left\langle\left(u_{3} v_{12}\right) \otimes\left(u_{1} v_{23}\right) \otimes\left(u_{3} v_{23}\right)\right\rangle_{\mathbb{Z}_{3}}  \tag{4.5.8}\\
& \left\langle\left(u_{2} v_{3}-u_{23} v_{1}\right) \otimes\left(u_{1} v_{2}-u_{12} v_{3}\right) \otimes\left(u_{3} v_{2}-u_{23} v_{3}\right)\right\rangle_{\mathbb{Z}_{3}}  \tag{4.5.9}\\
& \left\langle\left(u_{23} v_{12}+u_{2} v_{3}-u_{1} v_{23}\right) \otimes\left(u_{2} v_{3}\right) \otimes\left(u_{3} v_{2}\right)\right\rangle_{\mathbb{Z}_{3}}  \tag{4.5.10}\\
& \left\langle\left(u_{12} v_{12}+u_{2} v_{3}-u_{3} v_{2}\right) \otimes\left(u_{2} v_{1}\right) \otimes\left(u_{1} v_{2}\right)\right\rangle_{\mathbb{Z}_{3}} . \tag{4.5.11}
\end{align*}
$$

The transpose-like $\mathbb{Z}_{2}$ is $x \otimes y \otimes z \mapsto\left(\epsilon_{2} y \epsilon_{2}\right)^{T} \otimes\left(\epsilon_{2} x \epsilon_{2}\right)^{T} \otimes\left(\epsilon_{2} z \epsilon_{2}\right)^{T}$, where $\epsilon_{2}=\left(\begin{array}{lll}1 & & \\ & -1 & \\ & & 1\end{array}\right)$. (Note the similarities with Strassen's decomposition.)
In other words send $u_{1} \leftrightarrow v_{1}, u_{2} \leftrightarrow-v_{2}, u_{3} \leftrightarrow v_{3}$ and then switch the first two factors in $A \otimes B \otimes C$. This action performs the exchanges (4.5.3) $\leftrightarrow$ (4.5.4) and (4.5.7) $\leftrightarrow(4.5 .8)$, and fixes all other terms in the decomposition.

Here is the incidence graph:


### 4.6. Alternating least squares (ALS) method for decompositions

I now explain the method used to find decompositions numerically.
Let $A, B, C$ respectively have bases $\left\{e_{i}\right\},\left\{f_{j}\right\},\left\{g_{k}\right\}$. Given a tensor $T=$ $\sum_{i=1}^{\mathrm{a}} \sum_{j=1}^{\mathrm{b}} \sum_{k=1}^{\mathrm{c}} t^{i j k} e_{i} \otimes f_{j} \otimes g_{k} \in A \otimes B \otimes C$, say we have reason to believe it has rank at most $r$. To find a rank $r$ expression we could work as follows: For $1 \leq u \leq r$, write $a_{u}=\sum_{i} X_{u}^{i} e_{i}, b_{u}=\sum_{j} Y_{u}^{j} f_{j}$, and $c_{u}=\sum_{k} Z_{u}^{k} g_{k}$ where the $X_{u}^{i}, Y_{u}^{j}, Z_{u}^{k}$ are constants to be determined. We want $\sum_{u=1}^{r} a_{u} \otimes b_{u} \otimes c_{u}=T$, i.e.,

$$
\begin{equation*}
\sum_{u=1}^{r} X_{u}^{i} Y_{u}^{j} Z_{u}^{k}=t^{i j k} \tag{4.6.1}
\end{equation*}
$$

for all $i, j, k$. If we restrict ourselves to real coefficients, we want

$$
\begin{equation*}
\operatorname{objfn}_{1}:=\sum_{i, j, k}\left(\sum_{u=1}^{r} X_{u}^{i} Y_{u}^{j} Z_{u}^{k}-t^{i j k}\right)^{2}, \tag{4.6.2}
\end{equation*}
$$

called the objective function, to be zero. (One can obtain a similar equation for complex coefficients by splitting all complex numbers into their real and imaginary parts. I stick to the real presentation for simplicity of exposition.) Now (4.6.2) is a degree six polynomial, but it is quadratic in each of the unknown quantities. To solve in practice, one begins with an initial "guess" of the $X_{u}^{i}, Y_{u}^{j}, Z_{u}^{k}$, e.g., chosen at random. Then one tries to minimize (4.6.2) e.g., as a function of the $X_{u}^{i}$ while holding the $Y_{u}^{j}, Z_{u}^{k}$ fixed. This is a linear problem. Once one obtains a solution, one starts again, holding the $X_{u}^{i}$ and $Z_{u}^{k}$ fixed and solving for the $Y_{u}^{j}$. Then one repeats, minimizing for the $Z_{u}^{k}$, and then cycling around again and again until the result converges (or fails
to, in which case one can start again with different initial points). This algorithm was first written down in [Bre70].

Now this procedure could "attempt" to find a border rank solution, that is, the coefficients could go off to infinity. If one wants a rank decomposition, one can add a penalty term to (4.6.2), instead minimizing

$$
\begin{equation*}
\operatorname{objfn}_{2}:=\sum_{i, j, k}\left(\sum_{u=1}^{r} X_{u}^{i} Y_{u}^{j} Z_{u}^{k}-t^{i j k}\right)^{2}+\epsilon\left(\sum_{u, i, j, k}^{r}\left(X_{u}^{i}\right)^{2}+\left(Y_{u}^{j}\right)^{2}+\left(Z_{u}^{k}\right)^{2}\right) \tag{4.6.3}
\end{equation*}
$$

for some $\epsilon$ that in practice is found by trial and error.
In the literature (e.g. [Lad76, JM86, Smi13, AS13]) they prefer coefficient values to be from a small list of numbers, ideally confined to something like $0, \pm 1$ or $0, \pm 1, \pm \frac{1}{2}$. If the tensor in question has a large symmetry group (as does matrix multiplication), one can use the group action to fix some of the coefficients to these desired values.

According to Smirnov, in $[\mathbf{S m i 1 3}]$, for $T=M_{\langle\mathbf{n}\rangle}$ (but not rectangular matrix multiplication) the critical points of objfn ${ }_{1}$ are integers in practice, although he does not give an explanation why one would expect this to be the case. Thus, by these heuristics, if one can obtain a decomposition with objfn $_{1}<1$, then it will converge to zero by the ALS process, producing either a decomposition or limit to a border rank decomposition.

### 4.7. Secant varieties and additional geometric language

To better discuss border rank decompositions in $\S 4.8$, I now introduce the language of secant varieties. This language will also enable us to discuss rank decompositions in a larger context and will arise in the study of Valiant's conjecture and its variants.
4.7.1. Secant Varieties. Given a variety $X \subset \mathbb{P} V$, define the $X$-rank of $[p] \in \mathbb{P} V, \mathbf{R}_{X}([p])$, to be the smallest $r$ such that there exist $x_{1}, \ldots, x_{r} \in \hat{X}$ such that $p$ is in the span of $x_{1}, \ldots, x_{r}$, and the $X$-border rank $\underline{\mathbf{R}}_{X}([p])$ is defined to be the smallest $r$ such that there exist curves $x_{1}(t), \ldots, x_{r}(t) \in \hat{X}$ such that $p$ is in the span of the limiting plane $\lim _{t \rightarrow 0}\left\langle x_{1}(t), \ldots, x_{r}(t)\right\rangle$, where $\left\langle x_{1}(t), \ldots, x_{r}(t)\right\rangle \subset G(r, V)$ is viewed as a curve the Grassmannian. Here and in what follows, I am assuming that for $t \neq 0, x_{1}(t), \ldots, x_{r}(t)$ are linearly independent (otherwise we are really dealing with a decomposition of lower border rank).

Let $\sigma_{r}(X) \subset \mathbb{P} V$ denote the set of points of $X$-border rank at most $r$, called the $r$-th secant variety of $X$. (Theorem 3.1.6.1 assures us that $\sigma_{r}(X)$
is indeed a variety.) In other words

$$
\sigma_{r}(X)=\overline{\bigcup_{x_{1}, \ldots, x_{r} \in X}\left\langle x_{1}, \ldots, x_{r}\right\rangle}
$$

where $\left\langle x_{1}, \ldots, x_{r}\right\rangle$ denotes the linear span in projective space and the overline denotes Zariski closure. The notation is such that $\sigma_{1}(X)=X$. When $X=\operatorname{Seg}\left(\mathbb{P} A_{1} \times \cdots \times \mathbb{P} A_{n}\right)$ is the set of rank one tensors, $\sigma_{r}(X)=\sigma_{r}$.

Let $X \subset \mathbb{P} V$ be a smooth variety, and let $p \in \sigma_{2}(X)$. If $p$ is not a point of $X$, nor a point on an honest secant line, then $p$ must line on some tangent line to $X$, where here I take the naïve definition of tangent line, namely a point on a limit of secant lines.

Terracini's lemma (see, e.g., [Lan12, §5.3]) generalizes our caculation of $\hat{T}_{\left[a_{1} \otimes b_{1} \otimes c_{1}+a_{2} \otimes b_{2} \otimes c_{2}\right]} \operatorname{Seg}(\mathbb{P} A \times \mathbb{P} B \times \mathbb{P} C)$ of $\S 3.1 .3$ : if $z=\left[x_{1}+\cdots+x_{r}\right]$ with $\left[x_{j}\right] \in X$ general points, then $\hat{T}_{z} \sigma_{r}(X)=\sum_{j=1}^{r} \hat{T}_{\left[x_{j}\right]} X$. In particular $\operatorname{dim} \sigma_{r}(X) \leq r \operatorname{dim} X+r-1$.

Thus $\operatorname{dim} \sigma_{r}(X) \leq \min \{r \operatorname{dim} X+r-1, \mathbf{v}-1\}$, and when equality holds we will say $\sigma_{r}(X)$ is of the expected dimension. The expected dimension is indeed what occurs "most" of the time. For example, $\operatorname{dim} \sigma_{r}\left(\mathbb{P}^{N} \times \mathbb{P}^{N} \times \mathbb{P}^{N}\right)$ is the expected dimension $\min \left\{3 N r+r-1, N^{3}-1\right\}$ for all $(r, N)$ except $(r, N)=(4,2)[\mathbf{L i c} 85]$.
4.7.2. Homogeneous varieties, orbit closures, and $G$-varieties. The Segre, Veronese and Grassmannian of $\S 3.1 .2$ are examples of homogeneous varieties:

Definition 4.7.2.1. A subvariety $X \subset \mathbb{P} V$, is homogeneous if it is a closed orbit of some point $x \in \mathbb{P} V$ under the action of some group $G \subset G L(V)$. If $P \subset G$ is the subgroup fixing $x$, write $X=G / P$.

A variety $X \subset \mathbb{P} V$ is called a $G$-variety for a group $G \subset G L(V)$, if for all $g \in G$ and $x \in X, g \cdot x \in X$.

Orbit closures (see $\S 3.3 .1$ ) and homogeneous varieties are $G$-varieties.
Exercise 4.7.2.2: (1) What are the points in $\overline{G L_{n} \cdot\left(x_{1} \cdots x_{n}\right)}$ that are not in $G L_{n} \cdot\left(x_{1} \cdots x_{n}\right)$ ?
4.7.3. The abstract secant variety. Given projective varieties $Y_{j} \subset \mathbb{P} V_{j}$, one can define their Segre product $Y_{1} \times \cdots \times Y_{r} \subset \operatorname{Seg}\left(\mathbb{P} V_{1} \times \cdots \times \mathbb{P} V_{r}\right) \subset$ $\mathbb{P}\left(V_{1} \otimes \cdots \otimes V_{r}\right)$. Let $X \subset \mathbb{P} V$ be a variety. Consider the set

$$
\begin{aligned}
S_{r}(X)^{0}:= & \left\{\left(x_{1}, \ldots, x_{r}, z\right) \in X^{\times r} \times \mathbb{P} V \mid z \in \operatorname{span}\left\{x_{1}, \ldots, x_{r}\right\}\right\} \\
& \subset \operatorname{Seg}\left(X^{\times r} \times \mathbb{P} V\right) \subset \mathbb{P} V^{\otimes r+1}
\end{aligned}
$$

and let $S_{r}(X):=\overline{S_{r}(X)^{0}}$ denote its Zariski closure. (For those familiar with quotients, it would be more convenient to deal with $X^{(\times r)}:=X^{\times r} / \mathfrak{S}_{r}$.) We have a map $\pi^{0}: S_{r}(X)^{0} \rightarrow \mathbb{P} V$, extending to a map $\pi: S_{r}(X) \rightarrow \mathbb{P} V$, given by projection onto the last factor and the image is $\sigma_{r}^{0}(X)$ (resp. $\sigma_{r}(X)$ ). Call $S_{r}(X)$ the abstract $r$-th secant variety of $X$. As long as $r<\mathbf{v}$ and $X$ is not contained in a linear subspace of $\mathbb{P} V, \operatorname{dim} S_{r}(X)=r \operatorname{dim} X+r-1$ because $\operatorname{dim} X^{\times r}=r \operatorname{dim} X$ and a general set of $r$ points on $X$ will span a $\mathbb{P}^{r-1}$.

If $\sigma_{r}(X)$ is of the expected dimension and is not all of $\mathbb{P} V$, so its dimension equals that of $S_{r}(X)$, then for general points $z \in \sigma_{r}(X)^{0},\left(\pi^{0}\right)^{-1}(z)$ will consist of a finite number of points and each point will correspond to a decomposition $\bar{z}=\overline{x_{1}}+\cdots+\overline{x_{r}}$ for $\overline{x_{j}} \in \hat{x}_{j}, \bar{z} \in \hat{z}$. In summary:
Proposition 4.7.3.1. If $X^{n} \subset \mathbb{P}^{N}$ and $\sigma_{r}(X)$ is of (the expected) dimension $r n+r-1<N$, then each of the points of a Zariski dense subset of $\sigma_{r}(X)$ has a finite number of decompositions into a sum of $r$ elements of $X$.

If the fiber of $\pi^{0}$ over $z \in \sigma_{r}^{0}(X)$ is $k$-dimensional, then there is a $k$ parameter family of decompositions of $z$ as a sum of $r$ rank one tensors. This occurs, for example if $z \in \sigma_{r-1}^{0}(X)$, but it can also occur for points in $\sigma_{r}(X) \backslash \sigma_{r-1}(X)$.

For example, every point of $\sigma_{7}\left(\operatorname{Seg}\left(\mathbb{P}^{3} \times \mathbb{P}^{3} \times \mathbb{P}^{3}\right)\right)=\mathbb{P}^{63}$ has a 5 dimensional family of points in the fiber, but $M_{\langle 2\rangle}$ has a nine dimensional family. A general point of $\sigma_{23}\left(\operatorname{Seg}\left(\mathbb{P}^{8} \times \mathbb{P}^{8} \times \mathbb{P}^{8}\right)\right)$ will have a finite number of points in the fiber, but $M_{\langle 3\rangle}$ has at least a 24 -dimensional fiber, in fact by [JM86], at least a 27 -dimensional fiber.

If $X$ is a $G$-variety, then $\sigma_{r}(X)$ is also a $G$-variety, and if $z \in \sigma_{r}^{0}(X)$ is fixed by $G_{z} \subset G$, then $G_{z}$ will act (possibly trivially) on $\left(\pi^{0}\right)^{-1}(z)$, and every distinct (up to re-ordering if one is not working with $X^{(\times r)}$ ) point in its orbit will correspond to a distinct decomposition of $z$. Let $q \in\left(\pi^{0}\right)^{-1}(x)$. If $\operatorname{dim}\left(G_{z} \cdot q\right)=d_{z}$, then there is at least a $d_{z}$ parameter family of decompositions of $z$ as a sum of $r$ elements of $X$.

Remark 4.7.3.2. Note that $\operatorname{codim}\left(S_{r-1}(X), S_{r}(X)\right) \leq \operatorname{dim} X-1$, where the inclusion is just by adding any point of $X$ to a border rank $r-1$ decomposition. In particular, in the case of the Segre relevant for matrix multiplication, this codimension is at most $3\left(\mathbf{n}^{2}-1\right)$. On the other hand $\operatorname{dim} G_{M_{\langle\mathbf{n}\rangle}}=3\left(\mathbf{n}^{2}-1\right)$, so by a dimension count, one might "expect" $\pi_{r}^{-1}\left(M_{\langle\mathbf{n}\rangle}\right)$ to intersect $S_{r-1}(X)$, meaning that we could keep reducing the border rank of $M_{\langle\mathbf{n}\rangle}$ all the way down to one. Of course since $S_{r}(\operatorname{Seg}(\mathbb{P} A \times \mathbb{P} B \times \mathbb{P} C))$ is not a projective space, Theorem 3.1.5.1 does not apply, but this dimension count illustrates the pathology of the tensor $M_{\langle\mathbf{n}\rangle}$.
4.7.4. What is a border rank decomposition? Usually an $X$-border rank decomposition of some $v \in V$ is presented as $v=\lim _{t \rightarrow 0}\left(x_{1}(t)+\right.$ $\left.\cdots+x_{r}(t)\right)$ where $\left[x_{j}(t)\right]$ are curves in $X$. In order to discuss border rank decompositions geometrically, it will be useful to study the corresponding curve in the Grassmannian $\left\langle x_{1}(t), \ldots, x_{r}(t)\right\rangle \subset G(r, V)$. The geometry of the intersection of the limiting $r$ plane that contains $v$ with $X$ has useful information.

To better understand this geometry, consider
$\tilde{S}_{r}^{0}(X):=\left\{\left([v],\left(\left[x_{1}\right], \ldots,\left[x_{r}\right]\right), E\right) \mid v \in\left\langle x_{1}, \ldots, x_{r}\right\rangle \subseteq E\right\} \subset \mathbb{P} V \times X^{\times r} \times G(r, V)$
and $\tilde{S}_{r}(X):=\overline{\tilde{S}_{r}^{0}(X)}$.
We can stratify $\sigma_{r}(X)$ and $\tilde{S}_{r}(X)$ by the $h$ 's of the intermediate ranks $\mathbf{R}_{h}$ of $\S 3.2 .1$. The case $h=0$ is rank. The next case $h=1$ has a straight-forward geometry.

To understand the $h=1$ case, first consider the case $r=2$, so $v=$ $\lim _{t \rightarrow 0} \frac{1}{t}\left(x_{1}(t)+x_{2}(t)\right)$ for curves $\left[x_{j}(t)\right] \subset X$. Then we must have $\lim _{t \rightarrow 0}\left[x_{1}(t)\right]=$ $\lim _{t \rightarrow 0}\left[x_{2}(t)\right]$, letting $[x]$ denote this limiting point, we obtain an element of $\hat{T}_{x} X$. In the case of $\sigma_{r}(X)$, one needs $r$ curves such that the points are linearly independent for $t \neq 0$ and such that they become dependent when $t=0$. This is most interesting when no subset of $r-1$ points becomes linearly dependent. Then one may obtain an arbitrary point of $\hat{T}_{x_{1}} X+\cdots+\hat{T}_{x_{r}} X$ (see [Lan12, §10.8.1]). For some varieties there may not exist $r$ distinct points on them that are linearly dependent (e.g., $v_{d}\left(\mathbb{P}^{1}\right)$ when $d>r$ ). An easy way for such sets of points to exist is if there is a $\mathbb{P}^{r-1}$ on the variety, as was the case for $T_{S T R}$ of $\S 5.6$. The decompositions for $M_{\langle\mathbf{m}, 2,2\rangle}^{r e d}$ I discuss in the next section are not quite from such simple configurations, but nearly are. Because of this I next discuss the geometry of linear spaces on the Segre.
4.7.5. Lines on Segre varieties. There are three types of lines on $\operatorname{Seg}(\mathbb{P A} A$ $\mathbb{P} B \times \mathbb{P} C): \alpha$-lines, which are of the form $\mathbb{P}\left(\left\langle a_{1}, a_{2}\right\rangle \otimes b \otimes c\right)$ for some $a_{j} \in A$, $b \in B, c \in C$, and the other two types are defined similarly and called $\beta$ and $\gamma$ lines.
Exercise 4.7.5.1: (2) Show that all lines on $\operatorname{Seg}(\mathbb{P} A \times \mathbb{P} B \times \mathbb{P} C)$ are one of these types. ©

Given two lines $L_{\beta}, L_{\gamma} \subset \operatorname{Seg}(\mathbb{P} A \times \mathbb{P} B \times \mathbb{P} C)$ respectively of type $\beta, \gamma$, if they do not intersect, then $\left\langle L_{\beta}, L_{\gamma}\right\rangle=\mathbb{P}^{3}$ and if the lines are general, furthermore $\left\langle L_{\beta}, L_{\gamma}\right\rangle \cap \operatorname{Seg}(\mathbb{P} A \times \mathbb{P} B \times \mathbb{P} C)=L_{\beta} \sqcup L_{\gamma}$.

However if $L_{\beta}=\mathbb{P}\left(a \otimes\left\langle b_{1}, b_{2}\right\rangle \otimes c\right)$ and $L_{\gamma}=\mathbb{P}\left(a^{\prime} \otimes b \otimes\left\langle c_{1}, c_{2}\right\rangle\right)$ with $b \in$ $\left\langle b_{1}, b_{2}\right\rangle$ and $c \in\left\langle c_{1}, c_{2}\right\rangle$, then they still span a $\mathbb{P}^{3}$ but $\left\langle L_{\beta}, L_{\gamma}\right\rangle \cap \operatorname{Seg}(\mathbb{P} A \times$
$\mathbb{P} B \times \mathbb{P} C)=L_{\beta} \sqcup L_{\gamma} \sqcup L_{\alpha}$, where $L_{\alpha}=\mathbb{P}\left(\left\langle a, a^{\prime}\right\rangle \otimes b \otimes c\right)$, and $L_{\alpha}$ intersects both $L_{\beta}$ and $L_{\gamma}$.

Let $x, y, z \in \operatorname{Seg}(\mathbb{P} A \times \mathbb{P} B \times \mathbb{P} C)$ be distinct points that all lie on a line $L \subset S e g(\mathbb{P} A \times \mathbb{P} B \times \mathbb{P} C)$. Then
$\hat{T}_{x} S e g(\mathbb{P} A \times \mathbb{P} B \times \mathbb{P} C) \subset\left\langle\hat{T}_{y} S e g(\mathbb{P} A \times \mathbb{P} B \times \mathbb{P} C), \hat{T}_{z} S e g(\mathbb{P} A \times \mathbb{P} B \times \mathbb{P} C)\right\rangle$.
The analogous statement is true for lines on any cominuscule variety, see [BL14, Lemma 3.3]. Because of this, it will be more geometrical to refer to $\hat{T}_{L} S e g(\mathbb{P} A \times \mathbb{P} B \times \mathbb{P} C):=\left\langle\hat{T}_{y} S e g(\mathbb{P} A \times \mathbb{P} B \times \mathbb{P} C), \hat{T}_{z} S e g(\mathbb{P} A \times \mathbb{P} B \times \mathbb{P} C)\right\rangle$, as the choice of $y, z \in L$ is irrelevant.
Exercise 4.7.5.2: (1) Verify (4.7.1).
The matrix multiplication tensor $M_{\langle U, V, W\rangle}$ endows $A, B, C$ with additional structure, e.g., $B=V^{*} \otimes W$, so there are two types of distinguished $\beta$-lines (corresponding to lines of rank one matrices), call them $\left(\beta, \nu^{*}\right)$-lines and $(\beta, \omega)$-lines, where, e.g., a $\nu^{*}$-line is of the form $\mathbb{P}\left(a \otimes\left(\left\langle v^{1}, v^{2}\right\rangle \otimes w\right) \otimes c\right)$, and among such lines there are further distinguished ones where moreover both $a$ and $c$ also have rank one. Call such further distinguished lines special $\left(\beta, \nu^{*}\right)$-lines.

### 4.8. Border rank decompositions

4.8.1. $M_{\langle 2\rangle}^{\text {red }}$. Here $A \subset U^{*} \otimes V$ has dimension three.

What follows is a slight modification of the decomposition of $M_{\langle 2\rangle}^{r e d}$ from [BCRL79] that appeared in [LR0]. Call it the $B C L R$-decomposition. I label the points such that $x_{1}^{1}$ is set equal to zero. The main difference is that in the original all five points moved, but here one is stationary.

$$
\begin{aligned}
& p_{1}(t)=x_{2}^{1} \otimes\left(y_{2}^{2}+y_{1}^{2}\right) \otimes\left(z_{2}^{2}+t z_{1}^{1}\right) \\
& p_{2}(t)=-\left(x_{2}^{1}-t x_{2}^{2}\right) \otimes y_{2}^{2} \otimes\left(z_{2}^{2}+t\left(z_{1}^{1}+z_{1}^{2}\right)\right) \\
& p_{3}(t)=x_{1}^{2} \otimes\left(y_{1}^{2}+t y_{2}^{1}\right) \otimes\left(z_{2}^{2}+z_{2}^{1}\right) \\
& p_{4}(t)=\left(x_{1}^{2}-t x_{2}^{2}\right) \otimes\left(-y_{1}^{2}+t\left(y_{1}^{1}-y_{2}^{1}\right)\right) \otimes z_{2}^{1} \\
& p_{5}(t)=-\left(x_{1}^{2}+x_{2}^{1}\right) \otimes y_{1}^{2} \otimes z_{2}^{2}
\end{aligned}
$$

and

$$
\begin{equation*}
M_{\langle 2\rangle}^{r e d}=\lim _{t \rightarrow 0} \frac{1}{t}\left[p_{1}(t)+\cdots+p_{5}(t)\right] \tag{4.8.1}
\end{equation*}
$$

Use the notation $x_{j}^{i}=u^{i} \otimes v_{j}, y_{k}^{j}=v^{j} \otimes w_{k}$ and $z_{i}^{k}=w^{k} \otimes u_{i}$.

Theorem 4.8.1.1. [LR0] Let $E^{B C L R}=\lim _{t \rightarrow 0}\left\langle p_{1}(t), \ldots, p_{5}(t)\right\rangle \in G(5, A \otimes B \otimes C)$. Then $E^{B C L R} \cap \operatorname{Seg}(\mathbb{P} A \times \mathbb{P} B \times \mathbb{P} C)$ is the union of three lines:

$$
\begin{aligned}
& L_{12,(\beta, \omega)}=x_{2}^{1} \otimes\left(v^{2} \otimes W\right) \otimes z_{2}^{1} \\
& L_{21,\left(\gamma, \omega^{*}\right)}=x_{1}^{2} \otimes y_{2}^{2} \otimes\left(W^{*} \otimes u_{2}\right) \\
& L_{\alpha}=\left\langle x_{1}^{2}, x_{2}^{1}\right\rangle \otimes y_{2}^{2} \otimes z_{2}^{1} .
\end{aligned}
$$



Here $L_{12,(\beta, \omega)}$ is a special $(\beta, \omega)$-line, $L_{21,\left(\gamma, \omega^{*}\right)}$, is a special $\left(\gamma, \omega^{*}\right)$-line, and $L_{\alpha}$, is an $\alpha$-line with rank one $B$ and $C$ points. Moreover, the $C$-point of $L_{12,(\beta, \omega)}$ lies in the $\omega^{*}$-line of $L_{21,\left(\gamma, \omega^{*}\right)}$, the $B$-point of $L_{21,\left(\gamma, \omega^{*}\right)}$ lies in the $\omega$-line of $L_{12,(\beta, \omega)}$ and $L_{\alpha}$ is the unique line on the Segre intersecting $L_{12,(\beta, \omega)}$ and $L_{21,\left(\gamma, \omega^{*}\right)}$ (and thus it is contained in their span).

Furthermore, $E^{B C L R}=\left\langle M_{\langle 2\rangle}^{r e d}, L_{12,(\beta, \omega)}, L_{21,\left(\gamma, \omega^{*}\right)}\right\rangle$ and

$$
M_{\langle 2\rangle}^{r e d} \in\left\langle\hat{T}_{L_{12,(\beta, \omega)}} S e g(\mathbb{P} A \times \mathbb{P} B \times \mathbb{P} C), \hat{T}_{L_{21,\left(\gamma, \omega^{*}\right)}} S e g(\mathbb{P} A \times \mathbb{P} B \times \mathbb{P} C)\right\rangle
$$

Proof. Write $p_{j}=p_{j}(0)$. Then (up to sign, which is irrelevant for geometric considerations)

$$
\begin{aligned}
& p_{1}=x_{2}^{1} \otimes\left(y_{2}^{2}+y_{1}^{2}\right) \otimes z_{2}^{2} \\
& p_{2}=x_{2}^{1} \otimes y_{2}^{2} \otimes z_{2}^{2} \\
& p_{3}=x_{1}^{2} \otimes y_{1}^{2} \otimes\left(z_{2}^{2}+z_{2}^{1}\right) \\
& p_{4}=x_{1}^{2} \otimes y_{1}^{2} \otimes z_{2}^{1} \\
& p_{5}=\left(x_{1}^{2}+x_{2}^{1}\right) \otimes y_{1}^{2} \otimes z_{2}^{2} .
\end{aligned}
$$

Then $L_{12,(\beta, \omega)}=\left\langle p_{1}, p_{2}\right\rangle, L_{21,\left(\gamma, \omega^{*}\right)}=\left\langle p_{3}, p_{4}\right\rangle$, and $p_{5} \in L_{\alpha}$.
To see there are no other points in $E^{B C L R} \cap \operatorname{Seg}(\mathbb{P} A \times \mathbb{P} B \times \mathbb{P} C)$, first note that any such point would have to lie on $\operatorname{Seg}\left(\mathbb{P}\left\langle x_{2}^{1}, x_{1}^{2}\right\rangle \times \mathbb{P}\left\langle y_{1}^{2}, y_{2}^{2}\right\rangle \times \mathbb{P}\left\langle z_{2}^{1}, z_{2}^{2}\right\rangle\right)$ because there is no way to eliminate the rank two $x_{2}^{2} \otimes\left(y_{1}^{2} \otimes z_{2}^{1}+y_{2}^{2} \otimes z_{2}^{2}\right)$ term in $M_{\langle 2\rangle}^{r e d}$ with a linear combination of $p_{1}, \ldots, p_{4}$. Let $\left[\left(s x_{2}^{1}+t x_{1}^{2}\right) \otimes\left(u y_{2}^{2}+\right.\right.$ $\left.\left.v y_{1}^{2}\right) \otimes\left(p z_{2}^{2}+q z_{2}^{1}\right)\right]$ be an arbitrary point on this variety. To have it be in the span of $p_{1}, \ldots, p_{4}$ it must satisfy the equations $s u q=0, s v q=0, t u q=0$, tup $=0$. Keeping in mind that one cannot have $(s, t)=(0,0),(u, v)=(0,0)$,
or $(p, q)=(0,0)$, we conclude the only solutions are the three lines already exhibited.

We have

$$
\begin{aligned}
& p_{1}(0)^{\prime}=x_{2}^{1} \otimes\left(y_{2}^{2}+y_{1}^{2}\right) \otimes z_{1}^{1} \\
& p_{2}(0)^{\prime}=x_{2}^{2} \otimes y_{2}^{2} \otimes z_{2}^{2}-x_{2}^{1} \otimes y_{2}^{2} \otimes\left(-z_{1}^{2}+z_{1}^{1}\right) \\
& p_{3}(0)^{\prime}=x_{1}^{2} \otimes y_{2}^{1} \otimes\left(z_{2}^{2}+z_{2}^{1}\right) \\
& p_{4}(0)^{\prime}=x_{2}^{2} \otimes y_{1}^{2} \otimes z_{2}^{1}+x_{1}^{2} \otimes\left(y_{1}^{1}-y_{2}^{1}\right) \otimes z_{2}^{1} \\
& p_{5}(0)^{\prime}=0 .
\end{aligned}
$$

Then $M_{\langle 2\rangle}^{r e d}=\left(p_{1}^{\prime}+p_{2}^{\prime}\right)+\left(p_{3}^{\prime}+p_{4}^{\prime}\right)$ where $p_{1}^{\prime}+p_{2}^{\prime} \in T_{L_{12,(\beta, \omega)}} \operatorname{Seg}(\mathbb{P} A \times$ $\mathbb{P} B \times \mathbb{P} C)$ and $p_{3}^{\prime}+p_{4}^{\prime} \in T_{L_{21,\left(\gamma, \omega^{*}\right)}} S e g(\mathbb{P} A \times \mathbb{P} B \times \mathbb{P} C)$.

Remark 4.8.1.2. By removing $x_{1}^{1}$ from our tensor, we lose the cyclic $\mathbb{Z}_{3}$ symmetry but retain a standard transpose symmetry $x \otimes y \otimes z \mapsto x^{T} \otimes z^{T} \otimes y^{T}$. Similarly we lose the $G L(U) \times G L(V)$ symmetry but retain the $G L(W)$ action. By composing the standard transpose symmetry with another $\mathbb{Z}_{2}$ action which switches the basis vectors of $W$, the action swaps $p_{1}(t)+p_{2}(t)$ with $p_{3}(t)+p_{4}(t)$ and $L_{12,(\beta, \omega)}$ with $L_{21,\left(\gamma, \omega^{*}\right)}$. This action fixes $p_{5}$.

Remark 4.8.1.3. Note that it is important that $p_{5}$ lies neither on $L_{12,(\beta, \omega)}$ nor on $L_{21,\left(\gamma, \omega^{*}\right)}$, so that no subset of the five points lies in a linearly degenerate position to enable us to have tangent vectors coming from all five points, but I emphasize that any point on the line $L_{\alpha}$ not on the original lines would have worked equally well, so the geometric object is this configuration of lines.
4.8.2. $M_{\langle 3,2,2\rangle}^{r e d}$. Here is the decomposition in [AS13, Thm. 2] due to Alexeev and Smirnov, only changing the element set to zero in their decomposition to $x_{1}^{1}$. The decomposition is order two and the only nonzero coefficients appearing are $\pm 1, \pm \frac{1}{2}$.

$$
\begin{aligned}
& p_{1}(t)=\left(\frac{-1}{2} t^{2} x_{2}^{3}-\frac{1}{2} t x_{1}^{2}+x_{1}^{2}\right) \otimes\left(-y_{1}^{2}+y_{2}^{2}+t y_{1}^{1}\right) \otimes\left(z_{3}^{1}+t z_{2}^{1}\right) \\
& p_{2}(t)=\left(x_{1}^{2}+\frac{1}{2} x_{2}^{1}\right) \otimes\left(y_{1}^{2}-y_{2}^{2}\right) \otimes\left(z_{3}^{1}+z_{3}^{2}+t z_{2}^{1}+t z_{2}^{2}\right) \\
& p_{3}(t)=\left(t^{2} x_{2}^{3}+t x_{1}^{3}-\frac{1}{2} t x_{2}^{2}-x_{1}^{2}\right) \otimes\left(y_{1}^{2}+y_{2}^{2}+t y_{2}^{1}\right) \otimes z_{3}^{2} \\
& p_{4}(t)=\left(\frac{1}{2} t^{2} x_{2}^{3}-t x_{1}^{3}-\frac{1}{2} t x_{2}^{2}+x_{1}^{2}\right) \otimes\left(y_{1}^{2}+y_{2}^{2}-t y_{1}^{1}\right) \otimes z_{3}^{1} \\
& p_{5}(t)=\left(-t^{2} x_{2}^{3}+t x_{2}^{2}-x_{2}^{1}\right) \otimes y_{1}^{2} \otimes\left(z_{3}^{2}+\frac{1}{2} t z_{2}^{1}+\frac{1}{2} t z_{2}^{2}-t^{2} z_{1}^{1}\right) \\
& p_{6}(t)=\left(\frac{1}{2} t x_{2}^{2}+x_{1}^{2}\right) \otimes\left(-y_{1}^{2}+y_{2}^{2}+t y_{2}^{1}\right) \otimes\left(z_{3}^{2}+t z_{2}^{2}\right) \\
& p_{7}(t)=\left(-t x_{1}^{3}+x_{1}^{2}+\frac{1}{2} x_{2}^{1}\right) \otimes\left(y_{1}^{2}+y_{2}^{2}\right) \otimes\left(-z_{3}^{1}+z_{3}^{2}\right) \\
& p_{8}(t)=\left(t x_{2}^{2}+x_{2}^{1}\right) \otimes y_{2}^{2} \otimes\left(z_{3}^{1}+\frac{1}{2} t z_{2}^{1}+\frac{1}{2} t z_{2}^{2}+t^{2} z_{1}^{2}\right) .
\end{aligned}
$$

Then

$$
M_{\langle 3,2,2\rangle}^{r e d}=\frac{1}{t^{2}}\left[p_{1}(t)+\cdots+p_{8}(t)\right] .
$$

Remark 4.8.2.1. In [BDHM15] they prove $\underline{\mathbf{R}}\left(M_{\langle 3,2,2\rangle}^{r e d}\right)=8$.
Theorem 4.8.2.2. [LR0] Let $E^{A S, 3}=\lim _{t \rightarrow 0}\left\langle p_{1}(t), \ldots, p_{8}(t)\right\rangle \in G(8, A \otimes B \otimes C)$. Then $E^{A S, 3} \cap S e g(\mathbb{P} A \times \mathbb{P} B \times \mathbb{P} C)$ is the union of two irreducible algebraic surfaces, both abstractly isomorphic to $\mathbb{P}^{1} \times \mathbb{P}^{1}$ : The first is a sub-Segre variety:

$$
\operatorname{Seg}_{21,(\beta, \omega),\left(\gamma, \omega^{*}\right)}:=\left[x_{1}^{2}\right] \times \mathbb{P}\left(v^{2} \otimes W\right) \times \mathbb{P}\left(W^{*} \otimes u_{3}\right),
$$

The second, $\mathbb{L}_{\alpha}$ is a one-parameter family of lines passing through a parametrized curve in $\operatorname{Seg}_{21,(\beta, \omega),\left(\gamma, \omega^{*}\right)}$ and the plane conic curve (which has the same parametrization):

$$
C_{12,(\beta, \omega),\left(\gamma, \omega^{*}\right)}:=\mathbb{P}\left(\cup_{[s, t] \in \mathbb{P}^{1}} x_{2}^{1} \otimes\left(s y_{1}^{2}-t y_{2}^{2}\right) \otimes\left(s z_{3}^{2}+t z_{3}^{1}\right)\right) .
$$

The three varieties $C_{12,(\beta, \omega),\left(\gamma, \omega^{*}\right)}, \operatorname{Seg}_{21,(\beta, \omega),\left(\gamma, \omega^{*}\right)}$, and $\mathbb{L}_{\alpha}$ respectively play roles analogous to the lines $L_{12,(\beta, \omega)}, L_{21,\left(\gamma, \omega^{*}\right)}$, and $L_{\alpha}$, as described below.


Figure 4.8.1. The curve $C_{12,(\beta, \omega),\left(\gamma, \omega^{*}\right)}$ with its four points, the surface $S_{21,(\beta, \omega),\left(\gamma, \omega^{*}\right)}$, with its four points (only two of which are visible), and the surface $\mathbb{L}_{\alpha}$ with its two points which don't lie on either the curve or surface $\operatorname{Seg}_{21,(\beta, \omega),\left(\gamma, \omega^{*}\right)}$.

Proof. The limit points are (up to sign):

$$
\begin{aligned}
& p_{1}=x_{1}^{2} \otimes\left(y_{1}^{2}-y_{2}^{2}\right) \otimes z_{3}^{1} \\
& p_{3}=x_{1}^{2} \otimes\left(y_{1}^{2}+y_{2}^{2}\right) \otimes z_{3}^{2} \\
& p_{4}=x_{1}^{2} \otimes\left(y_{1}^{2}+y_{2}^{2}\right) \otimes z_{3}^{1} \\
& p_{6}=x_{1}^{2} \otimes\left(y_{1}^{2}-y_{2}^{2}\right) \otimes z_{3}^{2} \\
& p_{5}=x_{2}^{1} \otimes y_{1}^{2} \otimes z_{3}^{2} \\
& p_{8}=x_{2}^{1} \otimes y_{2}^{2} \otimes z_{3}^{1} \\
& p_{2}=\left(x_{1}^{2}+\frac{1}{2} x_{2}^{1}\right) \otimes\left(y_{1}^{2}-y_{2}^{2}\right) \otimes\left(z_{3}^{1}+z_{3}^{2}\right) \\
& p_{7}=\left(x_{1}^{2}+\frac{1}{2} x_{2}^{1}\right) \otimes\left(y_{1}^{2}+y_{2}^{2}\right) \otimes\left(z_{3}^{1}-z_{3}^{2}\right)
\end{aligned}
$$

Just as with $M_{\langle 2\rangle}^{r e d}$, the limit points all lie on a $\operatorname{Seg}\left(\mathbb{P}^{1} \times \mathbb{P}^{1} \times \mathbb{P}^{1}\right)$, in fact the "same" $\operatorname{Seg}\left(\mathbb{P}^{1} \times \mathbb{P}^{1} \times \mathbb{P}^{1}\right)$. Pictorially the Segres are:

$$
\left(\begin{array}{ll}
0 & * \\
* &
\end{array}\right) \times\left(\begin{array}{ll}
* & \\
* & *
\end{array}\right) \times\binom{ *}{*}
$$

for $M_{\langle 2,2,2\rangle}^{r e d}$ and

$$
\left(\begin{array}{ll}
0 & * \\
* &
\end{array}\right) \times\left(\begin{array}{cc}
* & \\
* & *
\end{array}\right) \times\left(\begin{array}{ll}
* & \\
*
\end{array}\right)
$$

for $M_{\langle 3,2,2\rangle}^{\text {red }}$. Here $E^{A S, 3} \cap \operatorname{Seg}(\mathbb{P} A \times \mathbb{P} B \times \mathbb{P} C)$ is the union of a oneparameter family of lines $\mathbb{L}_{\alpha}$ passing through a plane conic and a special $\mathbb{P}^{1} \times \mathbb{P}^{1}: \operatorname{Seg}_{21,(\beta, \omega),\left(\gamma, \omega^{*}\right)}:=\left[x_{1}^{2}\right] \times \mathbb{P}\left(v^{2} \otimes W\right) \times \mathbb{P}\left(W^{*} \otimes u_{3}\right)$ (which contains $p_{1}, p_{3}, p_{4}, p_{6}$ ). To define the family and make the similarity with the BCLR case clearer, first define the plane conic curve

$$
C_{12,(\beta, \omega),\left(\gamma, \omega^{*}\right)}:=\mathbb{P}\left(\cup_{[s, t] \in \mathbb{P}^{1}} x_{2}^{1} \otimes\left(s y_{1}^{2}-t y_{2}^{2}\right) \otimes\left(s z_{3}^{2}+t z_{3}^{1}\right)\right)
$$

The points $p_{5}, p_{8}$ lie on this conic (respectively the values $(s, t)=(1,0)$ and $(s, t)=(0,1))$. Then define the variety

$$
\mathbb{L}_{\alpha}:=\mathbb{P}\left(\cup_{[\sigma, \tau] \in \mathbb{P}^{1}} \cup_{[s, t] \in \mathbb{P}^{1}}\left(\sigma x_{2}^{1}+\tau x_{1}^{2}\right) \otimes\left(s y_{1}^{2}-t y_{2}^{2}\right) \otimes\left(s z_{3}^{2}+t z_{3}^{1}\right)\right),
$$

which is a one-parameter family of lines intersecting the conic and the special $\mathbb{P}^{1} \times \mathbb{P}^{1}$. The points $p_{2}, p_{7}$ lie on $\mathbb{L}_{\alpha}$ but not on the conic. Explicitly $p_{2}$ (resp. $\left.p_{7}\right)$ is the point corresponding to the values $(\sigma, \tau)=\left(1, \frac{1}{2}\right)$ and $(s, t)=(1,1)$ (resp. $(s, t)=(1,-1)$ ).

The analog of $L_{\alpha}$ in the $M_{\langle 2\rangle}^{r e d}$ decomposition is $\mathbb{L}_{\alpha}$, and $C_{12,(\beta, \omega),\left(\gamma, \omega^{*}\right)}$ and $\operatorname{Seg}_{21,(\beta, \omega),\left(\gamma, \omega^{*}\right)}$ are the analogs of the lines $L_{12,(\beta, \omega)}, L_{21,\left(\gamma, \omega^{*}\right)}$. (A difference here is that $C_{12,(\beta, \omega),\left(\gamma, \omega^{*}\right)} \subset \mathbb{L}_{\alpha}$.)

The span of the configuration is the span of a $\mathbb{P}^{2}$ (the span of the conic) and a $\mathbb{P}^{3}$ (the span of the $\mathbb{P}^{1} \times \mathbb{P}^{1}$ ), i.e., a $\mathbb{P}^{6}$.

The proof that these are the only points in the intersection is similar to the BCLR case.

More decompositions are described geometrically in [LR0].
It would be reasonable to expect that the BCLR and Alekseev-Smirnov decompositions generalize to all $\mathbf{m}$, so that $\underline{\mathbf{R}}\left(M_{\langle\mathbf{m}, 2,2\rangle}^{\text {red }}\right) \leq 3 \mathbf{m}-1$, which would imply that $\underline{\mathbf{R}}\left(M_{\langle\mathbf{n}, 2,2\rangle}\right) \leq 3 \mathbf{n}+1$ for all $\mathbf{n}$.

# The complexity of Matrix multiplication IV: The complexity of tensors and more lower bounds 

In Chapter 2 we developed equations to test the border rank of tensors. In this chapter I explain further techniques for proving lower and upper bounds for border rank and rank. I also discuss geometric properties that could be useful for future investigations.

I begin, in $\S 5.1$ by making explicit the dictionary between ( $1_{A^{\prime}}$-generic) tensors in $\mathbb{C}^{\mathbf{a}} \otimes \mathbb{C}^{\mathbf{m}} \otimes \mathbb{C}^{\mathbf{m}}$ and linear subspaces of $\operatorname{End}\left(\mathbb{C}^{\mathbf{m}}\right)$. This enables one to both find new ways to bound rank and border rank via linear algebra, and to use knowledge of tensors to make progress on classical questions in linear algebra.

While up until now I have emphasized the use of explicit polynomials to test membership in varieties, sometimes varieties satisfy Zariski closed conditions that are easy to describe but difficult to write as polynomials. Some such are discussed in $\S 5.1$. Two more such conditions are discussed in §5.2. One particularly useful such technique, the border substitution method is discussed in detail in $\S 5.4$. In particular, it enables the $2 \mathbf{n}^{2}-\log _{2}(\mathbf{n})-1$ lower bound for $\underline{\mathbf{R}}\left(M_{\langle\mathbf{n}\rangle}\right)$ presented in §5.4.5.

Regarding tensor rank, the only general method for proving tensor rank lower bounds I am aware of is the substitution method discussed in §5.3.

The best upper bounds for the exponent $\omega$ were obtained with $T_{S T R}, T_{c w, q}$, and $T_{C W, q}$. What makes these tensors special? It is clear they have nice combinatorial properties, but do they have distinguishing geometric features? I discuss several such geometric properties in $\S 5.5$. If such features could be identified, one could in principle look for other tensors with the same properties and to apply the laser method to those tensors, as was proposed in
[AFLG15].
Several tensors that have been studied arise naturally as structure tensors of algebras. I discuss rank and border rank lower bounds for structure tensors of algebras in §5.6. In particular I present Bläser's and Zuiddam's sequences of tensors with rank to border rank ratio approaching three.

### 5.1. Tensors and classical linear algebra

This section follows [LM15].
5.1.1. 1-genericity. How good are Strassen's equations? We have seen that unless there exists $\alpha \in A^{*}$ with $T(\alpha) \subset B \otimes C$ of maximal rank (or $\beta \in B^{*}$, resp. $\gamma \in C^{*}$ with $T(\beta)$, resp. $T(\gamma)$, of maximal rank), they are essentially useless. The following definition names the class of tensors they are useful for.

Definition 5.1.1.1. A tensor $T \in A \otimes B \otimes C$ is $1_{A}$-generic if there exists $\alpha \in A^{*}$ with $T(\alpha) \subset B \otimes C$ of maximal rank, and $T$ is 1 -generic if it is $1_{A}, 1_{B}$ and $1_{C}$-generic.

Fortunately $M_{\langle\mathbf{n}\rangle}$ and all tensors used to study the exponent of matrix multiplication are 1-generic.

The 1-genericity of $M_{\langle\mathbf{n}\rangle}$ has the consequence that for the purpose of proving $\underline{\mathbf{R}}\left(M_{\langle\mathbf{n}\rangle}\right) \leq r$, it would be sufficient to find a collection of polynomials such that their common zero set simply contains $\sigma_{r}\left(\operatorname{Seg}\left(\mathbb{P}^{\mathbf{n}^{2}-1} \times \mathbb{P}^{\mathbf{n}^{2}-1} \times\right.\right.$ $\left.\mathbb{P}^{\mathbf{n}^{2}-1}\right)$ ) as an irreducible component, as long as all other components of the zero set are contained in the set of non-1-generic tensors.

Say a tensor $T$ is $1_{A}$-generic, $\mathbf{b}=\mathbf{c}$ and Strassen's commutators are identically zero- can we conclude $\underline{\mathbf{R}}(T)=\mathbf{b}$ ?

I address this question in this section and the next. I first show that the properties of tensor rank and border rank of tensors in $A \otimes B \otimes C$ can be studied as properties of a-dimensional linear subspaces of $B \otimes C$.
5.1.2. The dictionary. The following standard result shows that the rank and border rank of a tensor $T \in A \otimes B \otimes C$, may be recovered from the subspace $T\left(A^{*}\right) \subset B \otimes C$. I present a version of it from [LM15].
Proposition 5.1.2.1. For a tensor $T \in A \otimes B \otimes C, \mathbf{R}(T)$ equals the minimal number of rank one elements of $B \otimes C$ needed to span (a space containing) $T\left(A^{*}\right)$, and similarly for the permuted statements.

Say $\operatorname{dim} T\left(A^{*}\right)=k$. Let $Z_{r} \subset G(k, B \otimes C)$ denote the set of $k$-planes in $B \otimes C$ that are contained in the span of $r$ rank one elements, so $\mathbf{R}(T) \leq r$ if and only if $T\left(A^{*}\right) \in Z_{r}$. Then $\underline{\mathbf{R}}(T) \leq r$ if and only if $T\left(A^{*}\right) \in \overline{Z_{r}}$.

Proof. Let $T$ have rank $r$ so there is an expression $T=\sum_{i=1}^{r} a_{i} \otimes b_{i} \otimes c_{i}$. (The vectors $a_{i}$ need not be linearly independent, and similarly for the $b_{i}$ and $c_{i}$.) Then $T\left(A^{*}\right) \subseteq\left\langle b_{1} \otimes c_{1}, \ldots, b_{r} \otimes c_{r}\right\rangle$ shows that the number of rank one matrices needed to span $T\left(A^{*}\right) \subset B \otimes C$ is at most $\mathbf{R}(T)$.

For the other inequality, say $T\left(A^{*}\right)$ is contained in the span of rank one elements $b_{1} \otimes c_{1}, \ldots, b_{r} \otimes c_{r}$. Let $\alpha^{1}, \ldots, \alpha^{\text {a }}$ be a basis of $A^{*}$, with dual basis $e_{1}, \ldots, e_{\mathbf{a}}$ of $A$. Then $T\left(\alpha^{i}\right)=\sum_{s=1}^{r} x_{s}^{i} b_{s} \otimes c_{s}$ for some constants $x_{s}^{i}$. But then $T=\sum_{s, i} e_{i} \otimes\left(x_{s}^{i} b_{s} \otimes c_{s}\right)=\sum_{s=1}^{r}\left(\sum_{i} x_{s}^{i} e_{i}\right) \otimes b_{s} \otimes c_{s}$ proving $\mathbf{R}(T)$ is at most the number of rank one matrices needed to span $T\left(A^{*}\right) \subset B \otimes C$.
Exercise 5.1.2.2: (1) Prove the border rank assertion.
5.1.3. Equations via linear algebra. All the equations we have seen so far arise as Koszul flattenings, which all vanish if Strassen's equations for minimal border rank are zero, as can be seen by the coordinate expressions (2.2.1) and the discussion in $\S 2.4 .3$. Thus we have robust equations only if $T$ is $1_{A}, 1_{B}$ or $1_{C}$-generic, because otherwise the presence of $T(\alpha)^{\wedge \mathbf{a}-1}$ in the expressions make them likely to vanish. When $T$ is $1_{A}$-generic, the Koszul flattenings $T_{A}^{\wedge p}: \Lambda^{p} A \otimes B^{*} \rightarrow \Lambda^{p+1} A \otimes C$ provide measures of the failure of $T\left(A^{*}\right) T(\alpha)^{-1} \subset \operatorname{End}(B)$ to be an abelian subspace.

A first concern is that perhaps the choice of $\alpha \in A^{*}$ effects this failure. The following lemma addresses that concern, at least in the case of minimal border rank:
Lemma 5.1.3.1. [LM15] Let $T \in A \otimes B \otimes C=\mathbb{C}^{\mathbf{a}} \otimes \mathbb{C}^{\mathbf{a}} \otimes \mathbb{C}^{\mathbf{a}}$ be $1_{A}$-generic and assume $\operatorname{rank}\left(T\left(\alpha_{0}\right)\right)=\mathbf{a}$. If $T\left(A^{*}\right) T\left(\alpha_{0}\right)^{-1}$ is abelian then $T\left(A^{*}\right) T\left(\alpha_{0}^{\prime}\right)^{-1}$ is abelian for any $\alpha_{0}^{\prime} \in A^{*}$ such that $\operatorname{rank}\left(T\left(\alpha_{0}^{\prime}\right)\right)=\mathbf{a}$.

Proof. Say $T\left(A^{*}\right) T\left(\alpha_{0}\right)^{-1}$ is abelian, and set $X_{i}=T\left(\alpha_{i}\right) T\left(\alpha_{0}\right)^{-1}$, so $\left[X_{1}, X_{2}\right]=$ 0 . Set $X_{i}^{\prime}=T\left(\alpha_{i}\right) T\left(\alpha_{0}\right)^{-1}$ and $X^{\prime}=T\left(\alpha_{0}^{\prime}\right) T\left(\alpha_{0}\right)^{-1}$, so $\left[X_{i}, X^{\prime}\right]=0$ as well, which implies $\left[X_{i},\left(X^{\prime}\right)^{-1}\right]=0$. We want to show $\left[X_{1}^{\prime}, X_{2}^{\prime}\right]=0$. But

$$
\begin{aligned}
& X_{j}^{\prime}=X_{j}\left(X^{\prime}\right)^{-1}, \text { so } \\
& \qquad \begin{aligned}
X_{1}^{\prime} X_{2}^{\prime}-X_{2}^{\prime} X_{1}^{\prime} & =X_{1}\left(X^{\prime}\right)^{-1} X_{2}\left(X^{\prime}\right)^{-1}-X_{2}\left(X^{\prime}\right)^{-1} X_{1}\left(X^{\prime}\right)^{-1} \\
& =X_{1} X_{2}\left(X^{\prime}\right)^{-1}\left(X^{\prime}\right)^{-1}-X_{2} X_{1}\left(X^{\prime}\right)^{-1}\left(X^{\prime}\right)^{-1} \\
& =\left[X_{1}, X_{2}\right]\left(X^{\prime}\right)^{-1}\left(X^{\prime}\right)^{-1} \\
& =0 .
\end{aligned}
\end{aligned}
$$

Definition 5.1.3.2. Let $\mathbf{a}=\mathbf{b}=\mathbf{c}$ and let $\mathrm{Abel}_{A} \subset A \otimes B \otimes C$ denote the set of concise, $1_{A}$-generic tensors such that for some (and hence any) $\alpha \in A^{*}$ with $T(\alpha)$ of maximal rank, $T\left(A^{*}\right) T(\alpha)^{-1} \subset \operatorname{End}(B)$ is abelian. Note that Abel $_{A}$ is not Zariski closed.

Let $\operatorname{Diag}_{\operatorname{End}(B)}^{0} \subset G(\mathbf{b}, \operatorname{End}(B))$ denote the set of $\mathbf{b}$-dimensional subspaces that are simultaneously diagonalizable under the action of $G L(B)$ and let $\operatorname{Diag}_{\operatorname{End}(B)}=\overline{\operatorname{Diag}_{\operatorname{End}(B)}^{0}}$ denote its Zariski closure. Let $\alpha \in A^{*}$ be such that $T(\alpha)$ is of maximal rank, and let

$$
\operatorname{Diag}_{A}:=\overline{\left\{T \in \operatorname{Abel}_{A} \mid T\left(A^{*}\right) T(\alpha)^{-1} \in \operatorname{Diag}_{\operatorname{End}(B)}\right\}} \cap \operatorname{Abel}_{A} .
$$

By definition, $\operatorname{Diag}_{A} \subseteq \mathrm{Abel}_{A}$. To what extent does equality hold? The following proposition gives a necessary algebraic condition to be in $\mathrm{Diag}_{A}$ :
Proposition 5.1.3.3. [Ger61] The set

$$
\{U \in G(\mathbf{a}, \operatorname{End}(B)) \mid U \text { is closed under composition }\}
$$

is Zariski closed.
In particular, if $T \in A \otimes B \otimes C=\mathbb{C}^{\mathbf{a}} \otimes \mathbb{C}^{\mathbf{a}} \otimes \mathbb{C}^{\mathbf{a}}$ is $1_{A}$-generic with $\underline{\mathbf{R}}(T)=$ a, then for all $\alpha \in A^{*}$ with $T(\alpha)$ invertible, $T\left(A^{*}\right) T(\alpha)^{-1}$ is closed under composition.

Proof. If $u_{1}, \ldots, u_{\mathbf{a}}$ is a basis of $U$, then $U$ is closed under composition if and only if for all $u \in U$,

$$
\left(u u_{j}\right) \wedge u_{1} \wedge \cdots \wedge u_{\mathbf{a}}=0 \forall 1 \leq j \leq \mathbf{a} .
$$

Let $\left(\text { Abel }_{A} \times A^{*}\right)^{0}=\{(T, \alpha) \mid \operatorname{rank}(T(\alpha))=\mathbf{b}\}$, and note that the map $\left(\text { Abel }_{A} \times A^{*}\right)^{0} \rightarrow G(\mathbf{a}, \operatorname{End}(B))$, given by $(T, \alpha) \mapsto T\left(A^{*}\right) T(\alpha)^{-1}$ is continuous. The "in particular" assertion follows from this continuity because if $U \in \operatorname{Diag}_{\operatorname{End}(B)}^{0}$, then $U$ is closed under composition.

Exercise 5.1.3.4: (2) Show that if $T(\alpha), T\left(\alpha^{\prime}\right)$ are invertible and $T\left(A^{*}\right) T(\alpha)^{-1}$ is closed under composition, then $T\left(A^{*}\right) T\left(\alpha^{\prime}\right)^{-1}$ is closed under composition.

Let End $\mathrm{Abel}_{A} \subseteq \mathrm{Abel}_{A}$ denote the subset of tensors with $T(A) T(\alpha)^{-1}$ closed under composition for some (and hence all) $\alpha \in A^{*}$ with $T(\alpha)$ invertible. We have

$$
\begin{equation*}
\operatorname{Diag}_{A} \subseteq \operatorname{End} \operatorname{Abel}_{A} \subseteq \operatorname{Abel}_{A}, \tag{5.1.1}
\end{equation*}
$$

where the first inclusion is Proposition 5.1.3.3 and the second is by definition. Are these containments strict?

A classical theorem states that when $\mathbf{a}=3$ the three sets are equal. Moreover:
Theorem 5.1.3.5. [IM05] When $\mathbf{a} \leq 4, \operatorname{Diag}_{A}=\operatorname{End}_{\mathrm{Abel}}^{A}{ }^{\prime}=\mathrm{Abel}_{A}$.
See [IM05] for the proof, which has numerous cases.
What happens when $\mathbf{a}=5$ ?
Proposition 5.1.3.6. [Lei16] Let $T_{\text {Leit,5 }}=a_{1} \otimes\left(b_{1} \otimes c_{1}+b_{2} \otimes c_{2}+b_{3} \otimes c_{3}+\right.$ $\left.b_{4} \otimes c_{4}+b_{5} \otimes c_{5}\right)+a_{2} \otimes\left(b_{1} \otimes c_{3}+b_{3} \otimes c_{5}\right)+a_{3} \otimes b_{1} \otimes c_{4}+a_{4} \otimes b_{2} \otimes c_{4}+a_{5} \otimes b_{2} \otimes c_{5}$, which gives rise to the linear space

$$
T_{\text {Leit,5 }}\left(A^{*}\right)=\left(\begin{array}{lllll}
x_{1} & & & &  \tag{5.1.2}\\
& x_{1} & & & \\
x_{2} & & x_{1} & & \\
x_{3} & x_{4} & & x_{1} & \\
& x_{5} & x_{2} & & x_{1}
\end{array}\right)
$$

Then $T_{\text {Leit }, 5}\left(A^{*}\right) T\left(\alpha^{1}\right)^{-1}$ is an abelian Lie algebra, but not End-closed. I.e., $T_{\text {Leit }, 5} \in$ Abel $_{A}$ but $T_{\text {Leit }, 5} \notin$ End Abel $_{A}$.

Throughout this chapter, an expression of the form (5.1.2) is to be read as $T\left(x_{1} \alpha^{1}+\cdots x_{\mathbf{a}} \alpha^{\mathbf{a}}\right)$ where $\alpha^{1}, \ldots, \alpha^{\mathbf{a}}$ is a basis of $A^{*}$.
Exercise 5.1.3.7: (1) Verify that $T_{\text {Leit,5 }}\left(A^{*}\right) T\left(\alpha^{1}\right)^{-1}$ is not closed under composition.

Thus when $\mathbf{a} \geq 5$, End $\mathrm{Abel}_{A} \subsetneq \mathrm{Abel}_{A}$. The following proposition shows that the first containment in (5.1.1) is also strict when $\mathbf{a} \geq 7$ :
Proposition 5.1.3.8. [LM15] The tensor corresponding to

$$
T_{\text {end }, 7}\left(A^{*}\right)=\left(\begin{array}{ccccccc}
x_{1} & & & & & & \\
& x_{1} & & & & & \\
& & x_{1} & & & & \\
& & & x_{1} & & & \\
& x_{2}+x_{7} & x_{3} & x_{4} & x_{1} & & \\
x_{2} & x_{3} & x_{5} & x_{6} & & x_{1} & \\
x_{4} & x_{5} & x_{6} & x_{7} & & & x_{1}
\end{array}\right)
$$

is in End $\mathrm{Abel}_{A}$, but has border rank at least 8.

The proof is given in §5.2.1.
We have seen that set-theoretic equations for End $\mathrm{Abel}_{A}$ are easy, whereas set-theoretic equations for $\operatorname{Diag}_{A}$ are not known. One might hope that if $T \in \operatorname{End}_{\mathrm{Abel}_{A}}$, that at least $\underline{\mathbf{R}}(T)$ should be close to a. This hope fails miserably:
Proposition 5.1.3.9. [LM15] There exist $1_{A}$-generic tensors in $\mathbb{C}^{\mathbf{a}} \otimes \mathbb{C}^{\mathbf{a}} \otimes \mathbb{C}^{\mathbf{a}}$ in End $\mathrm{Abel}_{A}$ of border rank greater than $\frac{\mathbf{a}^{2}}{8}$.

Proof. Consider $T$ such that

$$
T\left(A^{*}\right) \subset\left(\begin{array}{cccccc}
x_{1} & & & & &  \tag{5.1.3}\\
& \ddots & & & & \\
& & x_{1} & & & \\
* & \cdots & * & x_{1} & & \\
\vdots & \vdots & \vdots & & \ddots & \\
* & \cdots & * & & & x_{1}
\end{array}\right) \text {. }
$$

and set $x_{1}=0$. We obtain a generic tensor in $\mathbb{C}^{\mathbf{a}-1} \otimes \mathbb{C}^{\left\lfloor\frac{2}{2}\right\rfloor} \otimes \mathbb{C}^{\left\lceil\frac{a}{2}\right\rceil}$, which will have border greater than $\frac{a^{2}}{8}$. Conclude by applying Exercise 2.1.6.2.

Tensors of the form (5.1.3) expose a weakness of Strassen's equations that I discuss further in §5.4.2. Variants of the tensors of the form (5.1.3) are 1-generic and still exhibit the same behavior.
5.1.4. Sufficient conditions for a concise tensor to be of minimal border rank. A classical result in linear algebra says a subspace $U \subset$ $\operatorname{End}(B)$ is diagonalizable if and only if $U$ is abelian and every $x \in U$ (or equivalently for each $x_{j}$ in a basis of $U$ ), is diagonalizable. This implies:
Proposition 5.1.4.1. A necessary and sufficient condition for a concise $1_{A}$-generic tensor $T \in A \otimes B \otimes C$ with $\mathbf{a}=\mathbf{b}=\mathbf{c}$ to be of minimal rank $\mathbf{a}$ is that for some basis $\alpha_{1}, \ldots, a_{\mathbf{a}}$ of $A^{*}$ with $\operatorname{rank}\left(T\left(\alpha_{1}\right)\right)=\mathbf{b}$, the space $T(A) T\left(\alpha_{1}\right)^{-1} \subset \operatorname{End}(B)$ is abelian and each $T\left(\alpha_{j}\right) T\left(\alpha_{1}\right)^{-1}$ is diagonalizable.

Although we have seen several necessary conditions to be of minimal border rank, the question is open in general:

Problem 5.1.4.2. [BCS97, Prob. 15.2] Classify concise tensors of minimal border rank.

Below is a sufficient condition to be of minimal border rank.
For $x \in \operatorname{End}(B)$, define the centralizer of $x$, denoted $C(x)$, by

$$
C(x):=\{y \in \operatorname{End}(B) \mid[y, x]=0\} .
$$

Definition 5.1.4.3. An element $x \in \operatorname{End}(B)$ is regular if $\operatorname{dim} C(x)=\mathbf{b}$, and it is regular semi-simple if $x$ is diagonalizable with distinct eigenvalues.

Exercise 5.1.4.4: (2) An $\mathbf{m} \times \mathbf{m}$ matrix is regular nilpotent if it is zero except for the super diagonal where the entries are all 1's. Show that a regular nilpotent element is indeed regular, and that its centralizer is the space of upper-triangular matrices where the entries on each (upper) diagonal are the same, e.g., when $\mathbf{m}=3$ the centralizer is

$$
\left\{\left.\left(\begin{array}{lll}
x & y & z \\
& x & y \\
& & x
\end{array}\right) \right\rvert\, x, y, z \in \mathbb{C}\right\} .
$$

Exercise 5.1.4.5: (2) Show that $\operatorname{dim} C(x) \geq \mathbf{b}$, with equality if and only if the minimal polynomial of $x$ equals the characteristic polynomial. ©

Note that $x$ is regular semi-simple if and only if $C(x) \subset \operatorname{End}(B)$ is a diagonalizable subspace. In this case the eigenvalues of $x$ are distinct.
Proposition 5.1.4.6. (L. Manivel, [LM15]) Let $U \subset \operatorname{End}(B)$ be an abelian subspace of dimension $\mathbf{b}$ such that there exists $x \in U$ that is regular. Then $U \in \operatorname{Diag}_{\operatorname{End}(B)} \subset G(\mathbf{b}, \operatorname{End}(B))$.

Proof. Since the Zariski closure of the set of regular semi-simple elements is all of $\operatorname{End}(B)$, for any $x \in \operatorname{End}(B)$, there exists a curve $x_{t}$ of regular semi-simple elements with $\lim _{t \rightarrow 0} x_{t}=x$. Consider the induced curve in the Grassmannian $C\left(x_{t}\right) \subset G(\mathbf{b}, \operatorname{End}(B))$. Then $C_{0}:=\lim _{t \rightarrow 0} C\left(x_{t}\right)$ exists and is contained in $C(x) \subset \operatorname{End}(B)$ and since $U$ is abelian, we also have $U \subseteq$ $C(x)$. But if $x$ is regular, then $\operatorname{dim} C(x)=\operatorname{dim}(U)=\mathbf{b}$, so $\lim _{t \rightarrow 0} C\left(x_{t}\right), C_{0}$ and $U$ must all be equal and thus $U$ is a limit of diagonalizable subspaces.

Proposition 5.1.4.6 applied to $T(A) T(\alpha)^{-1}$ provides a sufficient condition for a concise $1_{A}$-generic tensor $T \in A \otimes B \otimes C$ to be of minimal border rank. The condition is not necessary, even for 1 -generic tensors, e.g., the Coppersmith-Winograd tensor $T_{q, C W}$ of (3.4.5), is 1-generic of minimal border rank but $T_{q, C W}\left(A^{*}\right) T_{q, C W}(\alpha)^{-1}$ does not contain a regular element for any $\alpha \in A^{*}$.
Exercise 5.1.4.7: (2) Show that the centralizer of $T_{\mathbb{C}\left[\mathbb{Z}_{m}\right]}\left(x_{1}\right)$ from Example 3.5.1.2 is $T_{\mathbb{C}\left[\mathbb{Z}_{m}\right]}\left(\mathbb{C}\left[\mathbb{Z}_{m}\right]\right)$ to obtain a second proof that $\underline{\mathbf{R}}\left(T_{\mathbb{C}\left[\mathbb{Z}_{m}\right]}\right)=m$.

Problem 5.1.4.8. Determine a criterion for $U \in G(\mathbf{b}, \operatorname{End}(B))$ to be in the closure of the diagonalizable b-planes, when $U$ does not contain a regular element.

### 5.1.5. Strassen's equations and symmetric tensors.

Proposition 5.1.5.1. [LM15] Let $T \in A \otimes B \otimes C=\mathbb{C}^{\mathrm{m}} \otimes \mathbb{C}^{\mathrm{m}} \otimes \mathbb{C}^{\mathrm{m}}$ be $1_{A}$ and $1_{B}$ generic and satisfy the $A$-Strassen equations. Then, after a suitable choice of identification of $A$ with $B$ via bases, $T$ is isomorphic to a tensor in $S^{2} A \otimes C$.

In particular:
(1) After making choices of general $\alpha \in A^{*}$ and $\beta \in B^{*}, T\left(A^{*}\right)$ and $T\left(B^{*}\right)$ are $G L_{m}$-isomorphic subspaces of $\operatorname{End}\left(\mathbb{C}^{\mathbf{m}}\right)$.
(2) If $T$ is 1-generic, then $T$ is isomorphic to a tensor in $S^{3} \mathbb{C}^{\mathrm{m}}$.

Proof. Let $\left\{a_{i}\right\},\left\{b_{j}\right\},\left\{c_{k}\right\}$ respectively be bases of $A, B, C$, with dual bases $\left\{\alpha_{i}\right\},\left\{\beta_{j}\right\},\left\{\gamma_{k}\right\}$. Write $T=\sum t^{i j k} a_{i} \otimes b_{j} \otimes c_{k}$. After a change of basis in $A$ so that $\operatorname{rank}\left(T\left(\alpha_{1}\right)\right)=\mathbf{m}$ and in $B, C$, so that it is the identity matrix, we may assume $t^{1 j k}=\delta_{j k}$ and after a change of basis $B$ so that $T\left(\beta_{1}\right)$ is of full rank and further changes of bases in $A, B, C$, we may assume $t^{i 1 k}=\delta_{i k}$ as well. (To obtain $t^{i 1 k}=\delta_{i k}$ only requires changes of bases in $A, C$, but a further change in $B$ may be needed to preserve $t^{1 j k}=\delta_{j k}$.) Identify $T\left(A^{*}\right) \subset \operatorname{End}\left(\mathbb{C}^{\mathbf{m}}\right)$ via $\alpha^{1}$. Strassen's $A$-equations then say

$$
0=\left[T\left(\alpha^{i_{1}}\right), T\left(\alpha^{i_{2}}\right)\right]_{(j, k)}=\sum_{l} t^{i_{1} j l} t^{i_{2} l k}-t^{i_{2 j} j l} t^{i_{1} l k} \forall i_{1}, i_{2}, j, k .
$$

Consider when $j=1$ :

$$
0=\sum_{l} t^{i_{1} 1 l} t^{i_{2} l k}-t^{i_{2} 1 l} t^{i_{1} l k}=t^{i_{2} i_{1} k}-t^{i_{1} i_{2} k} \forall i_{1}, i_{2}, k,
$$

because $t^{i l l}=\delta_{i, l}$. But this says $T \in S^{2} \mathbb{C}^{\mathbf{m}} \otimes \mathbb{C}^{\mathbf{m}}$.
For the last assertion, say $L_{B}: B \rightarrow A$ is such that $\operatorname{Id}_{A} \otimes L_{B} \otimes \operatorname{Id}_{C}(T) \in$ $S^{2} A \otimes C$ and $L_{C}: C \rightarrow A$ is such that $\operatorname{Id}_{A} \otimes \operatorname{Id}_{B} \otimes L_{C} \in S^{2} A \otimes B$. Then $\operatorname{Id}_{A} \otimes L_{B} \otimes L_{C}(T)$ is in $A^{\otimes 3}$, symmetric in the first and second factors as well as the first and third. But $\mathfrak{S}_{3}$ is generated by two transpositions, so $\mathrm{Id}_{A} \otimes L_{B} \otimes L_{C}(T) \in S^{3} A$.

Thus the $A, B$, and $C$-Strassen equations for minimal border rank, despite being non-isomorphic modules (see [LM08a]), when restricted to 1generic tensors, all have the same zero sets.

### 5.2. Indirectly defined equations

This section and §5.4.1 discuss Zariski closed conditions that in principle give rise to equations, but they are difficult to write down explicitly- to do so systematically one would need to use elimination theory which is impossible to implement in practice other than in very small cases. Nonetheless, for certain tensors these conditions can be used to prove lower bounds on border
rank, e.g., the lower bound on $\underline{\mathbf{R}}\left(M_{\langle\mathbf{n}\rangle}\right)$ via Griesser's equations in $\S 5.2 .2$ and the state of the art lower bound on $\underline{\mathbf{R}}\left(M_{\langle\mathbf{n}\rangle}\right)$ of Theorem 5.4.5.1.

### 5.2.1. Intersection properties.

Exercise 5.2.1.1: (2) [BCS97, Ex. 15.14] Given $T \in \mathbb{C}^{\mathbf{a}} \otimes \mathbb{C}^{\mathbf{a}} \otimes \mathbb{C}^{\mathbf{a}}=$ $A \otimes B \otimes C$ that is concise, show that $\mathbb{P} T\left(A^{*}\right) \cap \operatorname{Seg}(\mathbb{P} B \times \mathbb{P} C)=\emptyset$ implies $\underline{\mathbf{R}}(T)>\mathbf{a}$. ©

Proof of Proposition 5.1.3.8. The fact that $T_{\text {end, } 7}\left(A^{*}\right)$ is End-closed follows by inspection. The tensor has border rank at least 8 by Exercise 5.2.1.1 as $T_{\text {end, } 7}\left(A^{*}\right)$ does not intersect the Segre. Indeed, if it intersected Segre, the vanishing of size two minors implies $x_{1}=x_{4}=0,\left(x_{2}+x_{7}\right) x_{2}=0$ and $\left(x_{2}+x_{7}\right) x_{7}=0$. If $x_{2}+x_{7}=0$ then $x_{3}=0$, and $x_{7}^{2}=\left(x_{2}+x_{7}\right) x_{7}=0$ and hence $x_{2}=0$ as well and we are done. If $x_{2}=0$ analogously we obtain $x_{7}=0$ and $x_{3}=x_{5}=x_{6}=0$.

A complete flag in a vector space $V$ is a sequence of subspaces $0 \subset V_{1} \subset$ $V_{2} \subset \cdots \subset V_{\mathbf{v}}$ with $\operatorname{dim} V_{j}=j$.
Proposition 5.2.1.2. [Lei16, LM15] Let $T \in \mathbb{C}^{\mathbf{a}} \otimes \mathbb{C}^{\mathbf{a}} \otimes \mathbb{C}^{\mathbf{a}}=A \otimes B \otimes C$ be concise. If $\underline{\mathbf{R}}(T)=\mathbf{a}$, then there exists a complete flag $A_{1} \subset \cdots \subset A_{\mathbf{a}-1} \subset$ $A_{\mathbf{a}}=A^{*}$, with $\operatorname{dim} A_{j}=j$, such that $\mathbb{P} T\left(A_{j}\right) \subset \sigma_{j}(\operatorname{Seg}(\mathbb{P} B \times \mathbb{P} C))$.

Proof. Write $T=\lim _{t \rightarrow 0} \sum_{j=1}^{\mathbf{a}} a_{j}(t) \otimes X_{j}(t)$ where $X_{j}(t) \in B \otimes C$ have rank one. Since $T$ is concise, we may assume without loss of generality that $a_{1}(t), \ldots, a_{\mathbf{a}}(t)$ is a basis of $A$ for $t \neq 0$. Let $\alpha^{1}(t), \ldots, \alpha^{\mathbf{a}}(t) \in A^{*}$ be the dual basis. Then take $A_{k}(t)=\operatorname{span}\left\{\alpha^{1}(t), \ldots, \alpha^{k}(t)\right\} \in G\left(k, A^{*}\right)$ and $A_{k}=\lim _{t \rightarrow 0} A_{k}(t)$. Since $\mathbb{P} T^{*}\left(A_{k}(t)\right) \subset \sigma_{k}(\operatorname{Seg}(\mathbb{P} B \times \mathbb{P} C))$ the same must be true in the limit.

One can say even more. For example:
Proposition 5.2.1.3. [LM15] Let $T \in \mathbb{C}^{\mathbf{a}} \otimes \mathbb{C}^{\mathbf{a}} \otimes \mathbb{C}^{\mathbf{a}}=A \otimes B \otimes C$. If $\underline{\mathbf{R}}(T)=$ a and $T\left(A^{*}\right) \cap \operatorname{Seg}(\mathbb{P} B \times \mathbb{P} C)=\left[X_{0}\right]$ is a single point, then $\mathbb{P}\left(T\left(A^{*}\right) \cap\right.$ $\left.\hat{T}_{\left[X_{0}\right]} \operatorname{Seg}(\mathbb{P} B \times \mathbb{P} C)\right)$ must contain a $\mathbb{P}^{1}$.

Proof. Say $T\left(A^{*}\right)$ were the limit of $\operatorname{span}\left\{X_{1}(t), \ldots, X_{\mathbf{a}}(t)\right\}$ with each $X_{j}(t)$ of rank one. Then since $\mathbb{P} T\left(A^{*}\right) \cap S e g(\mathbb{P} B \times \mathbb{P} C)=\left[X_{0}\right]$, we must have each $X_{j}(t)$ limiting to $X_{0}$. But then $\lim _{t \rightarrow 0} \operatorname{span}\left\{X_{1}(t), X_{2}(t)\right\}$, which must be two-dimensional, must be contained in $\hat{T}_{\left[X_{0}\right]} S e g(\mathbb{P} B \times \mathbb{P} C)$ and $T\left(A^{*}\right)$.
5.2.2. Griesser's equations. The following theorem describes potential equations for $\sigma_{r}(\operatorname{Seg}(\mathbb{P} A \times \mathbb{P} B \times \mathbb{P} C))$ in the range $\mathbf{b}<r \leq 2 \mathbf{b}-1$.
Theorem 5.2.2.1. [Gri86] Let $\mathbf{b}=\mathbf{c}$. Given a $1_{A}$-generic tensor $T \in$ $A \otimes B \otimes C$ with $\underline{\mathbf{R}}(T) \leq r$, let $\alpha_{0} \in A^{*}$ be such that $T\left(\alpha_{0}\right)$ is invertible. For
$\alpha^{\prime} \in A^{*}$, let $X\left(\alpha^{\prime}\right)=T\left(\alpha^{\prime}\right) T\left(\alpha_{0}\right)^{-1} \in \operatorname{End}(B)$. Fix $\alpha_{1} \in A^{*}$. Consider the space of endomorphisms $U:=\left\{\left[X\left(\alpha_{1}\right), X\left(\alpha^{\prime}\right)\right]: B \rightarrow B \mid \alpha^{\prime} \in A^{*}\right\} \subset \mathfrak{s l}(B)$. Then there exists $E \in G(2 \mathbf{b}-r, B)$ such that $\operatorname{dim}(U . E) \leq r-\mathbf{b}$.

Remark 5.2.2.2. Compared with the minors of $T_{A}^{\wedge p}$, here one is just examining the first block column of the matrix appearing in the expression $Q \tilde{Q}$ in (2.4.7), but one is apparently extracting more refined information from it.

Proof. For the moment assume $\mathbf{R}(T)=r$ and $T=\sum_{j=1}^{r} a_{j} \otimes b_{j} \otimes c_{j}$. Let $\hat{B}=\mathbb{C}^{r}$ be equipped with basis $e_{1}, \ldots, e_{r}$. Define $\pi: \hat{B} \rightarrow B$ by $\pi\left(e_{j}\right)=b_{j}$. Let $i: B \rightarrow \hat{B}$ be such that $\pi \circ i=\operatorname{Id}_{B}$. Choose $B^{\prime} \subset \hat{B}$ of dimension $r-\mathbf{b}$ such that $\hat{B}=i(B) \oplus B^{\prime}$, and denote the inclusion and projection respectively $i^{\prime}: B^{\prime} \rightarrow \hat{B}$ and $\pi^{\prime}: \hat{B} \rightarrow B^{\prime}$. Pictorially:

\[

\]

Let $\alpha_{0}, \alpha_{1}, \ldots, \alpha_{\mathbf{a}-1}$ be a basis of $A^{*}$. Let $\hat{T}=\sum_{j=1}^{r} a_{j} \otimes e_{j} \otimes e_{j}^{*} \in A \otimes \hat{B} \otimes \hat{B}^{*}$ and let $\hat{X}_{j}:=\hat{T}\left(\alpha_{j}\right) \hat{T}\left(\alpha_{0}\right)^{\wedge r-1}$. (Recall that the matrix of $\hat{T}\left(\alpha_{0}\right)^{\wedge r-1}$ is the cofactor matrix of $\hat{T}\left(\alpha_{0}\right)$.) Now in $\operatorname{End}(\hat{B})$ all the commutators $\left[\hat{X}_{i}, \hat{X}_{j}\right]$ are zero because $\mathbf{R}(\hat{T})=r$. For all $2 \leq s \leq \mathbf{a}-1,\left[\hat{X}_{1}, \hat{X}_{s}\right]=0$ implies

$$
\begin{align*}
0 & =\pi\left[\hat{X}_{1}, \hat{X}_{s}\right] i \\
& =\left[X_{1}, X_{s}\right]+\left(\pi \hat{X}_{1} i^{\prime}\right)\left(\pi^{\prime} \hat{X}_{s} i\right)-\left(\pi \hat{X}_{s} i^{\prime}\right)\left(\pi^{\prime} \hat{X}_{1} i\right) \tag{5.2.1}
\end{align*}
$$

Now take $E \subseteq \operatorname{ker} \pi^{\prime} \hat{X}_{1} i \subset B$ of dimension $2 \mathbf{b}-r$. Then for all $s,\left[X_{1}, X_{s}\right]$. $E \subset$ Image $\pi \hat{X}_{1} i^{\prime}$, which has dimension at most $r-\mathbf{b}$ because $\pi \hat{X}_{1} i^{\prime}: B^{\prime} \rightarrow B$ and $\operatorname{dim} B^{\prime}=r-\mathbf{b}$. The general case follows because these conditions are all Zariski closed.

Proof of Theorem 2.2.2.1. Here there is just one commutator [ $X_{1}, X_{2}$ ] and its rank is at most the sum of the ranks of the other two terms in (5.2.1). But each of the other two terms is a composition of linear maps including $i^{\prime}$ which can have rank at most $r-\mathbf{b}$, so their sum can have rank at most $2(r-\mathbf{b})$.
Remark 5.2.2.3. It is not known to what extent Griesser's equations are non-trivial. Proving non-triviality of equations, even when the equations can be written down explicitly, is often more difficult than finding the equations. For example, it took several years after Koszul-flattenings were discovered to prove they were non-trivial to almost the full extent possible. Regarding Griesser's equations, it is known they are non-trivial up to $r \leq \frac{3}{2} \mathbf{m}+\frac{\sqrt{\mathbf{m}}}{2}-2$
when $\mathbf{m}$ is odd, and a similar, slightly smaller bound when $\mathbf{m}$ is even by Proposition 5.2.2.5 below. On the other hand the equations are trivial when $r=2 \mathbf{b}-1$ and all $\mathbf{a}$, and when $r=2 \mathbf{b}-2$, and $\mathbf{a} \leq \frac{\mathbf{b}}{2}+2$, in particular $\mathbf{a}=\mathbf{b}=4$ by $[\mathbf{L a n 1 5 b}]$. I do not know whether or not the equations are trivial for $r=2 \mathbf{b}-2, \mathbf{a}=\mathbf{b}$ and $\mathbf{b}>4$.

Griesser's equations are most robust when $T\left(\alpha_{1}\right) T\left(\alpha_{0}\right)^{-1}$ is a generic endomorphism, which motivates the following definition:

Definition 5.2.2.4. For a $1_{A}$-generic tensor $T \in A \otimes B \otimes C$, define $T$ to be $2_{A}$-generic if there exist $\alpha \in A^{*}$ such that $T(\alpha): C^{*} \rightarrow B$ is of maximal rank and $\alpha^{\prime} \in A^{*}$ such that $T\left(\alpha^{\prime}\right) T(\alpha)^{-1}: B \rightarrow B$ is regular semi-simple.

Proposition 5.1.4.6 implies that when $T \in \mathbb{C}^{\mathbf{m}} \otimes \mathbb{C}^{\mathbf{m}} \otimes \mathbb{C}^{\mathrm{m}}$ is concise, $2_{A^{-}}$ generic and satisfies Strassen's equations, then $\underline{\mathbf{R}}(T)=\mathbf{m}$.

Unfortunately for proving lower bounds, $M_{\langle\mathbf{n}\rangle}$ is not $2_{A^{-}}$generic. The equations coming from Koszul flattenings, and even more so Griesser's equations, are less robust for tensors that fail to be $2_{A}$-generic. This partially explains why $M_{\langle\mathbf{n}\rangle}$ satisfies some of the Koszul flattening equations and Griesser's equations (as shown below). Thus an important problem is to identify modules of equations for $\sigma_{r}$ that are robust for non-2-generic tensors.
Proposition 5.2.2.5. [Lan15b] Matrix multiplication $M_{\langle\mathbf{n}\rangle}$ fails to satisfy Griesser's equations for $r \leq \frac{3}{2} \mathbf{n}^{2}-1$ when $\mathbf{n}$ is even and $r \leq \frac{3}{2} \mathbf{n}^{2}+\frac{\mathbf{n}}{2}-2$ when $\mathbf{n}$ is odd, and satisfies the equations for all larger $r$.

Proof. Consider matrix multiplication $M_{\langle\mathbf{n}\rangle} \in \mathbb{C}^{\mathbf{n}^{2}} \otimes \mathbb{C}^{\mathbf{n}^{2}} \otimes \mathbb{C}^{\mathbf{n}^{2}}=A \otimes B \otimes C$. Recall from Exercise 2.1.7.4 that with a judicious ordering of bases, $M_{\langle\mathbf{n}\rangle}\left(A^{*}\right)$ is block diagonal

$$
\left(\begin{array}{lll}
x & &  \tag{5.2.2}\\
& \ddots & \\
& & x
\end{array}\right)
$$

where $x=\left(x_{j}^{i}\right)$ is $\mathbf{n} \times \mathbf{n}$. In particular, the image is closed under brackets. Choose $X_{0} \in M_{\langle\mathbf{n}\rangle}\left(A^{*}\right)$ to be the identity. It is not possible to have $X_{1} \in$ $M_{\langle\mathbf{n}\rangle}\left(A^{*}\right)$ diagonal with distinct entries on the diagonal, the most generic choice for $X_{1}$ is to be block diagonal with each block having the same $\mathbf{n}$ distinct entries. For a subspace $E$ of dimension $2 \mathbf{n}^{2}-r=d \mathbf{n}+e$ with $0 \leq e \leq \mathbf{n}-1$, the image of a generic choice of $\left[X_{1}, X_{2}\right], \ldots,\left[X_{1}, X_{\mathbf{n}^{2}-1}\right]$ applied to $E$ is of dimension at least $(d+1) \mathbf{n}$ if $e \geq 2$, at least $(d+1) \mathbf{n}-1$ if $e=1$ and $d \mathbf{n}$ if $e=0$, and equality will hold if we choose $E$ to be, e.g., the span of the first $2 \mathbf{n}^{2}-r$ basis vectors of $B$. (This is because the [ $X_{1}, X_{s}$ ] will span the entries of type (5.2.2) with zeros on the diagonal.) If
$\mathbf{n}$ is even, taking $2 \mathbf{n}^{2}-r=\frac{\mathbf{n}^{2}}{2}+1$, so $r=\frac{3 \mathbf{n}^{2}}{2}-1$, the image occupies a space of dimension $\frac{\mathbf{n}^{2}}{2}+\mathbf{n}-1>\frac{\mathbf{n}^{2}}{2}-1=r-\mathbf{n}^{2}$. If one takes $2 \mathbf{n}^{2}-r=$ $\frac{\mathbf{n}^{2}}{2}$, so $r=\frac{3 \mathbf{n}^{2}}{2}$, the image occupies a space of dimension $\frac{\mathbf{n}^{2}}{2}=r-\mathbf{n}^{2}$, showing Griesser's equations cannot do better for $\mathbf{n}$ even. If $\mathbf{n}$ is odd, taking $2 \mathbf{n}^{2}-r=\frac{\mathbf{n}^{2}}{2}-\frac{\mathbf{n}}{2}+2$, so $r=\frac{3 \mathbf{n}^{2}}{2}+\frac{\mathbf{n}}{2}-2$, the image will have dimension $\frac{\mathbf{n}^{2}}{2}+\frac{\mathbf{n}}{2}>r-\mathbf{n}^{2}=\frac{\mathbf{n}^{2}}{2}+\frac{\mathbf{n}}{2}-1$, and taking $2 \mathbf{n}^{2}-r=\frac{\mathbf{n}^{2}}{2}-\frac{\mathbf{n}}{2}+1$ the image can have dimension $\frac{\mathbf{n}^{2}}{2}-\frac{\mathbf{n}}{2}+(\mathbf{n}-1)=r-\mathbf{n}^{2}$, so the equations vanish for this and all larger $r$. Thus Griesser's equations for $\mathbf{n}$ odd give Lickteig's bound $\underline{\mathbf{R}}\left(M_{\langle\mathbf{n}\rangle}\right) \geq \frac{3 \mathbf{n}^{2}}{2}+\frac{\mathbf{n}}{2}-1$.

### 5.3. The substitution method

The following method has a long history dating back to [Pan66], see [BCS97, Chap. 6] and [Blä14, Chapter 6] for a history and many applications. It is the only general technique available for proving lower bounds on tensor rank that I am aware of. However, limit of the method is at most tensor rank lower bounds of $3 \mathbf{m}-1$ in $\mathbb{C}^{\mathbf{m}} \otimes \mathbb{C}^{\mathbf{m}} \otimes \mathbb{C}^{\mathbf{m}}$. (In $\S 10.1$ I will describe a powerful method for proving lower bounds on symmetric rank.)
5.3.1. Lower bounds on tensor rank via the substitution method.

Proposition 5.3.1.1. [AFT11, Appendix B] Let $T \in A \otimes B \otimes C$. Fix a basis $a_{1}, \ldots, a_{\mathbf{a}}$ of $A$, with dual basis $\alpha^{1}, \ldots, \alpha^{\mathbf{a}}$. Write $T=\sum_{i=1}^{\mathbf{a}} a_{i} \otimes M_{i}$, where $M_{i} \in B \otimes C$. Let $\mathbf{R}(T)=r$ and $M_{1} \neq 0$. Then there exist constants $\lambda_{2}, \ldots, \lambda_{\mathbf{a}}$, such that the tensor

$$
\tilde{T}:=\sum_{j=2}^{\mathbf{a}} a_{j} \otimes\left(M_{j}-\lambda_{j} M_{1}\right) \in \operatorname{span}\left\{a_{2}, \ldots, a_{\mathbf{a}}\right\} \otimes B \otimes C,
$$

has rank at most $r-1$. Moreover, if $\operatorname{rank}\left(M_{1}\right)=1$ then for any choice of $\lambda_{j}, \mathbf{R}(\tilde{T})$ is either $r$ or $r-1$.

The same assertions hold exchanging the role of $A$ with that of $B$ or $C$.
Proof. (Following [LM15].) By Proposition 5.1.2.1 there exist $X_{1}, \ldots, X_{r} \in$ $\hat{S} e g(\mathbb{P} B \times \mathbb{P} C)$ and scalars $d_{j}^{i}$ such that:

$$
M_{j}=\sum_{i=1}^{r} d_{j}^{i} X_{i} .
$$

Since $M_{1} \neq 0$ we may assume $d_{1}^{1} \neq 0$ and define $\lambda_{j}=\frac{d_{j}^{1}}{d_{1}^{1}}$. Then the subspace $\tilde{T}\left(\left\langle\alpha^{2}, \ldots, \alpha^{\mathbf{a}}\right\rangle\right)$ is spanned by $X_{2}, \ldots, X_{r}$ so Proposition 5.1.2.1 implies $\mathbf{R}(\tilde{T}) \leq r-1$. The last assertion holds because if $\operatorname{rank}\left(M_{1}\right)=1$ then we may assume $X_{1}=M_{1}$, so we cannot lower the rank by more than one.

In practice, the method is used iteratively, with each of $A, B, C$ playing the role of $A$ above, to reduce $T$ to a smaller and smaller tensor, at each step gaining one in the lower bound for the rank of $T$. At some steps one may project $T$ to a smaller space to simplify the calculation.

Example 5.3.1.2. [AFT11] Let $T_{a f t, 3} \in A \otimes B \otimes C$ have an expression in bases such that, letting the columns of the following matrix correspond to $B$-basis vectors and the rows to $C$ basis vectors,

$$
T_{a f t, 3}\left(A^{*}\right)=\left(\begin{array}{llllllll}
x_{1} & & & & & & & \\
& x_{1} & & & & & & \\
& & x_{1} & & & & & \\
& & & x_{1} & & & & \\
x_{2} & & & & x_{1} & & & \\
& x_{2} & & & & x_{1} & & \\
x_{3} & & x_{2} & & & & x_{1} & \\
x_{4} & x_{3} & & x_{2} & & & & x_{1}
\end{array}\right) \text {. }
$$

For the first iteration of the substitution method, start with $b_{8} \in B$ in the role of $a_{1}$ in the proposition. Write

$$
\begin{aligned}
T_{a f t, 3}= & b_{1} \otimes\left(a_{1} \otimes c_{1}+a_{2} \otimes c_{5}+a_{3} \otimes c_{7}+a_{4} \otimes c_{8}\right)+b_{2} \otimes\left(a_{1} \otimes c_{2}+a_{2} \otimes c_{6}+a_{3} \otimes c_{8}\right) \\
& +b_{3} \otimes\left(a_{1} \otimes c_{3}+a_{2} \otimes c_{7}\right)+b_{4} \otimes\left(a_{1} \otimes c_{4}+a_{2} \otimes c_{8}\right) \\
& +b_{5} \otimes a_{1} \otimes c_{5}+b_{6} \otimes a_{1} \otimes c_{6}+b_{6} \otimes a_{1} \otimes c_{6}+b_{7} \otimes a_{1} \otimes c_{7}+b_{8} \otimes a_{1} \otimes c_{8} .
\end{aligned}
$$

Then there exist $\lambda_{1}, \ldots, \lambda_{7}$ and a new tensor $T^{\prime} \in A \otimes \mathbb{C}^{7} \otimes C$ with $\mathbf{R}(T) \geq$ $\mathbf{R}\left(T^{\prime}\right)+1$ where

$$
T^{\prime}\left(A^{*}\right)=\left(\begin{array}{lllllll}
x_{1} & & & & & & \\
& x_{1} & & & & & \\
& & x_{1} & & & & \\
& & & x_{1} & & & \\
x_{2} & & & & x_{1} & & \\
& x_{2} & & & & x_{1} & \\
x_{3} & & x_{2} & & & & x_{1} \\
x_{4} & x_{3} & & x_{2} & & &
\end{array}\right)+\left(\begin{array}{lllll} 
& & & \\
& & & & \\
& & & & \\
\lambda_{1} x_{1} & \lambda_{2} x_{1} & \cdots & & \lambda_{7} x_{1}
\end{array}\right) .
$$

Continue removing the last three columns until we get a tensor $T^{\prime \prime} \in A \otimes \mathbb{C}^{4} \otimes C$ with

$$
T^{\prime \prime}\left(A^{*}\right)=\left(\begin{array}{lllll}
x_{1} & & & \\
& x_{1} & & \\
& & x_{1} & \\
& & & x_{1} \\
x_{2} & & & \\
& x_{2} & & \\
x_{3} & & x_{2} & \\
x_{4} & x_{3} & & x_{2}
\end{array}\right)+\left(\begin{array}{llll} 
& & & \\
\mu_{1,1} x_{1} & \mu_{2,1} x_{1} & \mu_{3,1} x_{1} & \mu_{4,1} x_{1} \\
\mu_{1,2} x_{1} & \mu_{2,2} x_{1} & \mu_{3,2} x_{1} & \mu_{4,2} x_{1} \\
\mu_{1,3} x_{1} & \mu_{2,3} x_{1} & \mu_{3,3} x_{1} & \mu_{4,3} x_{1} \\
\mu_{1,4} x_{1} & \mu_{2,4} x_{1} & \mu_{3,4} x_{1} & \mu_{4,4} x_{1}
\end{array}\right) .
$$

Now apply the method successively to $c_{1}, \ldots, c_{4}$ to obtain a tensor $T^{\prime \prime \prime}$ with $T^{\prime \prime \prime}\left(A^{*}\right) \in \mathbb{C}^{4} \otimes \mathbb{C}^{4}$ such that $\mathbf{R}\left(T_{a f t, 3}\right) \geq 8+\mathbf{R}\left(T^{\prime \prime \prime}\right)$. Now project $T^{\prime \prime \prime}$ to the space given by $x_{1}=0$, so all the unknown constants disappear. The new tensor cannot have rank or border rank greater than that of $T^{\prime \prime \prime}$. Iterate the method with the projection of $T^{\prime \prime \prime}$ until one arrives at $\tilde{T}\left(A^{*}\right) \in \mathbb{C}^{1} \otimes \mathbb{C}^{1}$ and the bound $\mathbf{R}\left(T_{a f t, 3}\right) \geq 8+4+2+1=15$. In fact $\mathbf{R}\left(T_{a f t, 3}\right)=15$ : observe that $T_{a f t, 3}\left(A^{*}\right) T_{a f t, 3}\left(\alpha^{1}\right)^{-1}$ is a projection of the centralizer of a regular nilpotent element as in Exercise 5.3.1.8 below, which implies $\mathbf{R}\left(T_{a f t, 3}\right) \leq 15$.

On the other hand $\underline{\mathbf{R}}\left(T_{a f t, 3}\right)=8$, again because $T_{\text {aft }, 3}\left(A^{*}\right) T_{a f t, 3}\left(\alpha^{1}\right)^{-1}$ is a projection of the centralizer of a regular nilpotent element, so Proposition 5.1.4. 6 applies.

This example generalizes to $T_{a f t, k} \in \mathbb{C}^{k+1} \otimes \mathbb{C}^{2^{k}} \otimes \mathbb{C}^{2^{k}}$ of rank $2 \cdot 2^{k}-1$ and border rank $2^{k}$. The tensor $T^{\prime \prime \prime}$ above is $T_{\text {aft }, 2}$.

Example 5.3.1.3. [AFT11] Let $T_{A F T, 3}=a_{1} \otimes\left(b_{1} \otimes c_{1}+\cdots+b_{8} \otimes c_{8}\right)+$ $a_{2} \otimes\left(b_{1} \otimes c_{5}+b_{2} \otimes c_{6}+b_{3} \otimes c_{7}+b_{4} \otimes c_{8}\right)+a_{3} \otimes\left(b_{1} \otimes c_{7}+b_{2} \otimes c_{8}\right)+a_{4} \otimes b_{1} \otimes c_{8}+$ $a_{5} \otimes b_{8} \otimes c_{1}+a_{6} \otimes b_{8} \otimes c_{2}+a_{7} \otimes b_{8} \otimes c_{3}+a_{8} \otimes b_{8} \otimes c_{4}$, so

$$
T_{A F T, 3}\left(A^{*}\right)=\left(\begin{array}{cccccccc}
x_{1} & & & & & & & x_{5} \\
& x_{1} & & & & & & x_{6} \\
& & x_{1} & & & & & x_{7} \\
& & & x_{1} & & & & x_{8} \\
x_{2} & & & & x_{1} & & & \\
& x_{2} & & & & x_{1} & & \\
x_{3} & & x_{2} & & & & x_{1} & \\
x_{4} & x_{3} & & x_{2} & & & & x_{1}
\end{array}\right)
$$

Begin the substitution method with $b_{8}$ in the role of $a_{1}$ in the proposition, then project to $\left\langle\alpha^{8}, \ldots, \alpha^{5}\right\rangle^{\perp}$ to obtain a tensor $\tilde{T}$ represented by the matrix

$$
\left(\begin{array}{lllllll}
x_{1} & & & & & & \\
& x_{1} & & & & & \\
& & x_{1} & & & & \\
& & & x_{1} & & & \\
x_{2} & & & & x_{1} & & \\
& x_{2} & & & & x_{1} & \\
x_{3} & & x_{2} & & & & x_{1} \\
x_{4} & x_{3} & & x_{2} & & &
\end{array}\right),
$$

and $\mathbf{R}\left(T_{A F T, 3}\right) \geq 4+\mathbf{R}(\tilde{T})$. The substitution method then gives $\mathbf{R}(\tilde{T}) \geq 14$ by Example 5.3.1.2 and thus $\mathbf{R}\left(T_{A F T, 3}\right) \geq 18$. This example generalizes to $T_{A F T, k} \in \mathbb{C}^{2^{k}+1} \otimes \mathbb{C}^{2^{k}} \otimes \mathbb{C}^{2^{k}+1}$ of rank at least $3\left(2^{k}+1\right)-k-4$. In fact, equality holds: in the case of $T_{A F T, 3}$, it is enough to consider 17 matrices with just one nonzero entry corresponding to all nonzero entries of $T_{A F T, 3}\left(A^{*}\right)$, apart from the top left and bottom right corner and one matrix with 1 at each corner and all other entries equal to 0 . Moreover, as observed in [Lan15b], for these tensors $\left(2^{k}+1\right)+1 \leq \underline{\mathbf{R}}\left(T_{A F T, k}\right) \leq 2^{k+1}-k$.
Exercise 5.3.1.4: (2) Prove $\left(2^{k}+1\right)+1 \leq \underline{\mathbf{R}}\left(T_{A F T, k}\right) \leq 2^{k+1}-k$. ©
In summary:
Proposition 5.3.1.5. The tensors $T_{A F T, k} \in \mathbb{C}^{2^{k}+1} \otimes \mathbb{C}^{2^{k}} \otimes \mathbb{C}^{2^{k}+1}$ of [AFT11] satisfy $\left(2^{k}+1\right)+1 \leq \underline{\mathbf{R}}\left(T_{A F T, k}\right) \leq 2\left(2^{k}+1\right)-2-k<3\left(2^{k}+1\right)-k-4=$ $\mathbf{R}\left(T_{A F T, k}\right)$.
Exercise 5.3.1.6: (2) Show that for all $\mathbf{m}, \mathbf{n}, N, \mathbf{R}\left(M_{\langle 1, \mathbf{m}, \mathbf{n}\rangle} \oplus M_{\langle N, 1,1\rangle}\right)=$ $\mathbf{m n}+N$.

Exercise 5.3.1.7: (2) Show that Strassen's tensor from §5.6, $T_{S T R, q}=$ $\sum_{j=1}^{q}\left(a_{0} \otimes b_{j} \otimes c_{j}+a_{j} \otimes b_{0} \otimes c_{j}\right) \in \mathbb{C}^{q+1} \otimes \mathbb{C}^{q+1} \otimes \mathbb{C}^{q}$ satisfies $\mathbf{R}\left(T_{S T R, q}\right)=2 q$.
Exercise 5.3.1.8: (3) Show that a tensor $T \in \mathbb{C}^{\mathbf{m}} \otimes \mathbb{C}^{\mathbf{m}} \otimes \mathbb{C}^{\mathbf{m}}$ corresponding to the centralizer of a regular nilpotent element satisfies $\mathbf{R}(T)=2 \mathbf{m}-1$. ©

To date, $T_{A F T, k}$ and its cousins are the only known examples of explicit tensors $T \in \mathbb{C}^{\mathbf{m}} \otimes \mathbb{C}^{\mathbf{m}} \otimes \mathbb{C}^{\mathbf{m}}$ satisfying $\mathbf{R}(T) \geq 3 \mathbf{m}-O(\log (\mathbf{m}))$. There are several known to satisfy $\mathbf{R}(T) \geq 3 \mathbf{m}-O(\mathbf{m})$, e.g., $M_{\langle\mathbf{n}\rangle}$, as was shown in $\S 2.6$, and $T_{W \text { State }}^{\otimes n} \in \mathbb{C}^{2^{n}} \otimes \mathbb{C}^{2^{n}} \otimes \mathbb{C}^{2^{n}}$ discussed in $\S 5.6$.
Problem 5.3.1.9. [Blä14] Find an explicit tensor $T \in \mathbb{C}^{\mathbf{m}} \otimes \mathbb{C}^{\mathbf{m}} \otimes \mathbb{C}^{\mathbf{m}}$ satisfying $\mathbf{R}(T) \geq(3+\epsilon) \mathbf{m}$ for any $\epsilon>0$.
Remark 5.3.1.10. Proposition 5.3.1.1 holds with any choice of basis, so we get to pick $\left[\alpha^{1}\right] \in \mathbb{P} A^{*}$, as long as $M_{1} \neq 0$ (which is automatic if $T$ is
$A$-concise). On the other hand, there is no choice of the $\lambda_{j}$, so when dealing with $\tilde{T}$, one has to assume the $\lambda_{j}$ are as bad as possible for proving lower bounds. For this reason, it is easier to implement this method on tensors with simple combinatorial structure or tensors that are sparse in some basis.

From a geometric perspective, we are restricting $T$, considered as a trilinear form $A^{*} \times B^{*} \times C^{*} \rightarrow \mathbb{C}$, to the hyperplane $A^{\prime} \subset A^{*}$ defined by $\alpha^{1}+$ $\sum_{j=2}^{\mathrm{a}} \lambda_{j} \alpha^{j}=0$ and our condition is that $\mathbf{R}\left(\left.T\right|_{A^{\prime} \otimes B^{*} \otimes C^{*}}\right) \leq \mathbf{R}(T)-1$. Our freedom is the choice of $\left\langle a_{2}, \ldots, a_{\mathbf{a}}\right\rangle \subset A$, and then $A^{\prime}$ (which we do not get to choose) is any hyperplane satisfying the open condition $\left\langle a_{2}, \ldots, a_{\mathbf{a}}\right\rangle^{\perp} \not \subset A^{\prime}$.
5.3.2. Strassen's additivity conjecture. Given $T_{1} \in A_{1} \otimes B_{1} \otimes C_{1}$ and $T_{2} \in A_{2} \otimes B_{2} \otimes C_{2}$, if one considers $T_{1}+T_{2} \in\left(A_{1} \oplus A_{2}\right) \otimes\left(B_{1} \oplus B_{2}\right) \otimes\left(C_{1} \oplus C_{2}\right)$, where each $A_{j} \otimes B_{j} \otimes C_{j}$ is naturally included in $\left(A_{1} \oplus A_{2}\right) \otimes\left(B_{1} \oplus B_{2}\right) \otimes\left(C_{1} \oplus\right.$ $C_{2}$ ), we saw that $\mathbf{R}\left(T_{1}+T_{2}\right) \leq \mathbf{R}\left(T_{1}\right)+\mathbf{R}\left(T_{2}\right)$. Also recall Schönhage's example $\S 3.3 .2$ that $\underline{\mathbf{R}}\left(M_{\langle 1, \mathbf{m}, \mathbf{n}\rangle} \oplus M_{\langle(\mathbf{n}-1)(\mathbf{m}-1), 1,1\rangle}\right)=\mathbf{m n}+1<2 \mathbf{m n}-$ $\mathbf{m}-\mathbf{n}+1=\underline{\mathbf{R}}\left(M_{\langle 1, \mathbf{m}, \mathbf{n}\rangle}\right)+\underline{\mathbf{R}}\left(M_{\langle(\mathbf{n}-1)(\mathbf{m}-1), 1,1\rangle}\right)$. Before this example was known, Strassen made the following conjecture:
Conjecture 5.3.2.1. [Str73] With the above notation, $\mathbf{R}\left(T_{1}+T_{2}\right)=\mathbf{R}\left(T_{1}\right)+$ $\mathbf{R}\left(T_{2}\right)$.

Exercise 5.3.1.6 shows that despite the failure of a border rank analog of the conjecture for $M_{\langle 1, \mathbf{m}, \mathbf{n}\rangle} \oplus M_{\langle(\mathbf{n}-1)(\mathbf{m}-1), 1,1\rangle}$, the rank version does hold in this case.

While this conjecture has been studied from several different perspectives, e.g. [FW84, JT86, Bsh98, CCC15b, BGL13], very little is known about it, and experts are divided as to whether it should be true or false.

In many cases of low rank the substitution method provides the correct rank. In light of this, the following theorem indicates why providing a counter-example to Strassen's conjecture would need new techniques for proving rank lower bounds.
Theorem 5.3.2.2. [LM15] Let $T_{1} \in A_{1} \otimes B_{1} \otimes C_{1}$ and $T_{2} \in A_{2} \otimes B_{2} \otimes C_{2}$ be such that that $\mathbf{R}\left(T_{1}\right)$ can be determined by the substitution method applied to two of $A_{1}, B_{1}, C_{1}$. Then Strassen's additivity conjecture holds for $T_{1} \oplus T_{2}$, i.e., $\mathbf{R}\left(T_{1} \oplus T_{2}\right)=\mathbf{R}\left(T_{1}\right)+\mathbf{R}\left(T_{2}\right)$.

Proof. With each application of the substitution method to elements of $A_{1}$, $B_{1}$, and $C_{1}, T_{1}$ is modified to a tensor of lower rank living in a smaller space and $T_{2}$ is unchanged. After all applications, $T_{1}$ has been modified to zero and $T_{2}$ is still unchanged.

The rank of any tensor in $\mathbb{C}^{2} \otimes B \otimes C$ can be computed using the substitution method as follows: by dimension count, we can always find either
$\beta \in B^{*}$ or $\gamma \in C^{*}$, such that $T(\beta)$ or $T(\gamma)$ is a rank one matrix. In particular, Theorem 5.3.2.2 provides an easy proof of Strassen's additivity conjecture if the dimension of any of $A_{1}, B_{1}$ or $C_{1}$ equals 2 . This was first shown in [JT86] by other methods.

### 5.4. The border substitution method

What follows are indirectly defined equations for border rank, in other words, indirectly defined algebraic varieties that contain $\sigma_{r}(\operatorname{Seg}(\mathbb{P} A \times \mathbb{P} B \times \mathbb{P} C))$. While we don't have equations for these varieties, sometimes one can prove membership or non-membership by direct arguments. The method is primarily useful for tensors with symmetry, as there border rank decompositions come in families, and it suffices to prove non-membership for a convenient member of a putative family.
5.4.1. The border substitution method. The substitution method may be restated as follows:
Proposition 5.4.1.1. Let $T \in A \otimes B \otimes C$ be $A$-concise. Fix a' $<\mathbf{a}$ and $\tilde{A} \subset A$ of dimension $\mathbf{a}^{\prime}$. Then

$$
\mathbf{R}(T) \geq \min _{\left\{A^{\prime} \in G\left(\mathbf{a}^{\prime}, A^{*}\right) \mid A^{\prime} \cap \tilde{A}^{\perp}=0\right\}} \mathbf{R}\left(\left.T\right|_{A^{\prime} \otimes B^{*} \otimes C^{*}}\right)+\left(\mathbf{a}-\mathbf{a}^{\prime}\right) .
$$

Here $\tilde{A}$ in the case $\mathbf{a}^{\prime}=\mathbf{a}-1$ plays the role of $\left\langle a_{2}, \ldots, a_{\mathbf{a}}\right\rangle$ in Proposition 5.3.1.1. Recall that $\left.T\right|_{A^{\prime} \otimes B^{*} \otimes C^{*}} \in\left(A /\left(A^{\prime}\right)^{\perp}\right) \otimes B \otimes C$.

More generally,
Proposition 5.4.1.2. Let $T \in A \otimes B \otimes C$ be concise. Fix $\mathbf{a}^{\prime}<\mathbf{a}, \tilde{A} \subset A$, $\tilde{B} \subset B$ and $\tilde{C} \subset C$ respectively of dimensions $\mathbf{a}^{\prime}, \mathbf{b}^{\prime}$, and $\mathbf{c}^{\prime}$. Then

$$
\begin{aligned}
\mathbf{R}(T) \geq & \left(\mathbf{a}-\mathbf{a}^{\prime}\right)+\left(\mathbf{b}-\mathbf{b}^{\prime}\right)+\left(\mathbf{c}-\mathbf{c}^{\prime}\right) \\
& +\min \\
& \left\{\begin{array}{l}
A^{\prime} \in G\left(\mathbf{a}^{\prime}, A^{*}\right) \mid A^{\prime} \cap \tilde{A}^{\perp}=0 \\
B^{\prime} \in G\left(\mathbf{b}^{\prime}, B^{*}\right) \mid B^{\prime} \cap \tilde{B}^{\perp}=0 \\
C^{\prime} \in G\left(\mathbf{c}^{\prime}, C^{*}\right) \mid A^{\prime} \cap \tilde{C}^{\perp}=0
\end{array}\right\} \mathbf{R ( T | _ { A ^ { \prime } \otimes B ^ { * } \otimes C ^ { * } } ) .}
\end{aligned}
$$

A border rank version is as follows:
Proposition 5.4.1.3. [BL16, LM] Let $T \in A \otimes B \otimes C$ be $A$-concise. Fix $\mathbf{a}^{\prime}<\mathbf{a}$. Then

$$
\underline{\mathbf{R}}(T) \geq \min _{A^{\prime} \in G\left(\mathbf{a}^{\prime}, A^{*}\right)} \underline{\mathbf{R}}\left(\left.T\right|_{A^{\prime} \otimes B^{*} \otimes C^{*}}\right)+\left(\mathbf{a}-\mathbf{a}^{\prime}\right) .
$$

Proof. Say $\underline{\mathbf{R}}(T)=r$, so $T=\lim _{t \rightarrow 0} T_{t}$, for some tensors $T_{t}=\sum_{j=1}^{r} a_{j}(t) \otimes b_{j}(t) \otimes c_{j}(t)$. Without loss of generality, we may assume $a_{1}(t), \ldots, a_{\mathbf{a}}(t)$ form a basis of $A$. Let $A_{t}^{\prime}=\left\langle a_{\mathbf{a}^{\prime}+1}, \ldots, a_{\mathbf{a}}\right\rangle^{\perp} \subset A^{*}$. Then $\mathbf{R}\left(\left.T_{t}\right|_{A_{t}^{\prime} \otimes B^{*} \otimes C^{*}}\right) \leq r-$ ( $\mathbf{a}-\mathbf{a}^{\prime}$ ) by Proposition 5.4.1.1. Let $A^{\prime}=\lim _{t \rightarrow 0} A_{t}^{\prime} \in G\left(\mathbf{a}^{\prime}, A^{*}\right)$. Then
$\left.T\right|_{A^{\prime} \otimes B^{*} \otimes C^{*}}=\left.\lim _{t \rightarrow 0} T_{t}\right|_{A_{t}^{\prime} \otimes B^{*} \otimes C^{*}}$ so $\underline{\mathbf{R}}\left(\left.T\right|_{A^{\prime} \otimes B^{*} \otimes C^{*}}\right) \leq r-\left(\mathbf{a}-\mathbf{a}^{\prime}\right)$, i.e., $r \geq \underline{\mathbf{R}}\left(\left.T\right|_{A^{\prime} \otimes B^{*} \otimes C^{*}}\right)+\left(\mathbf{a}-\mathbf{a}^{\prime}\right)$.

Corollary 5.4.1.4. [BL16] Let $T \in A \otimes B \otimes C$ be $A$-concise. Then $\underline{\mathbf{R}}(T) \geq$ $\mathbf{a}-1+\min _{\alpha \in A^{*} \backslash\{0\}} \operatorname{rank}(T(\alpha))$.

The Corollary follows because for matrices, rank equals border rank, and $\mathbb{C}^{1} \otimes B \otimes C=B \otimes C$.

Although our freedom in the substitution method was minor (a restriction to a Zariski open subset of the Grassmannian determined by $\tilde{A}^{\perp}$ ), it is still useful for tensors with simple combinatorial structure. With the border substitution method we have no freedom at all, but nevertheless it will be useful for tensors with symmetry, as the symmetry group will enable us to restrict to special $A^{\prime}$.

As was the case for the substitution method, this procedure can be iterated: write $T_{1}=\left.T\right|_{A^{\prime} \otimes B^{*} \otimes C^{*}}$. If $T_{1}$ is $B$-concise, apply the proposition again with $B$, if not, let $B_{1} \subset B$ be maximal such that $T_{1}$ is $B_{1}$-concise and then apply the proposition. By successive iterations one finds:
Corollary 5.4.1.5. [LM16] If for all $A^{\prime} \subset A^{*}, B^{\prime} \subset B^{*}, C^{\prime} \subset C^{*}$ respectively of dimensions $\mathbf{a}^{\prime}, \mathbf{b}^{\prime}, \mathbf{c}^{\prime}$ one has $\left.T\right|_{A^{\prime} \otimes B^{\prime} \otimes C^{\prime}} \neq 0$, then $\underline{\mathbf{R}}(T)>$ $\mathbf{a}+\mathbf{b}+\mathbf{c}-\left(\mathbf{a}^{\prime}+\mathbf{b}^{\prime}+\mathbf{c}^{\prime}\right)$.

It is obvious this method cannot prove border rank bounds better than $\mathbf{a}+\mathbf{b}+\mathbf{c}-3$. The actual limit of the method is even less, as I now explain.

### 5.4.2. Limits of the border substitution method.

Definition 5.4.2.1. A tensor $T \in A \otimes B \otimes C$ is ( $\mathbf{a}^{\prime}, \mathbf{b}^{\prime}, \mathbf{c}^{\prime}$ )-compressible if there exist subspaces $A^{\prime} \subset A^{*}, B^{\prime} \subset B^{*}, C^{\prime} \subset C^{*}$ of respective dimensions $\mathbf{a}^{\prime}, \mathbf{b}^{\prime}, \mathbf{c}^{\prime}$ such that $\left.T\right|_{A^{\prime} \otimes B^{\prime} \otimes C^{\prime}}=0$, i.e., there exists $\left(A^{\prime}, B^{\prime}, C^{\prime}\right) \in$ $G\left(\mathbf{a}^{\prime}, A^{*}\right) \times G\left(\mathbf{b}^{\prime}, B^{*}\right) \times G\left(\mathbf{c}^{\prime}, C^{*}\right)$, such that $A^{\prime} \otimes B^{\prime} \otimes C^{\prime} \subset T^{\perp}$, where $T^{\perp} \subset$ $(A \otimes B \otimes C)^{*}$ is the hyperplane annihilating $T$. Otherwise one says $T$ is ( $\mathbf{a}^{\prime}, \mathbf{b}^{\prime}, \mathbf{c}^{\prime}$ )-compression generic.

Let $X\left(\mathbf{a}^{\prime}, \mathbf{b}^{\prime}, \mathbf{c}^{\prime}\right)$ be the set of all tensors that are $\left(\mathbf{a}^{\prime}, \mathbf{b}^{\prime}, \mathbf{c}^{\prime}\right)$-compressible.
Corollary 5.4.1.5 may be rephrased as:

$$
\sigma_{\mathbf{a}+\mathbf{b}+\mathbf{c}-\left(\mathbf{a}^{\prime}+\mathbf{b}^{\prime}+\mathbf{c}^{\prime}\right)} S e g(\mathbb{P} A \times \mathbb{P} B \times \mathbb{P} C) \subset X\left(\mathbf{a}^{\prime}, \mathbf{b}^{\prime}, \mathbf{c}^{\prime}\right)
$$

Proposition 5.4.2.2. [LM16] The set $X\left(\mathbf{a}^{\prime}, \mathbf{b}^{\prime}, \mathbf{c}^{\prime}\right) \subseteq \mathbb{P}(A \otimes B \otimes C)$ is Zariski closed of dimension at most

$$
\min \left\{\mathbf{a b c}-1,\left(\mathbf{a b c}-\mathbf{a}^{\prime} \mathbf{b}^{\prime} \mathbf{c}^{\prime}-1\right)+\left(\mathbf{a}-\mathbf{a}^{\prime}\right) \mathbf{a}^{\prime}+\left(\mathbf{b}-\mathbf{b}^{\prime}\right) \mathbf{b}^{\prime}+\left(\mathbf{c}-\mathbf{c}^{\prime}\right) \mathbf{c}^{\prime}\right\} .
$$

In particular, if

$$
\begin{equation*}
\mathbf{a} \mathbf{a}^{\prime}+\mathbf{b} \mathbf{b}^{\prime}+\mathbf{c} \mathbf{c}^{\prime}<\left(\mathbf{a}^{\prime}\right)^{2}+\left(\mathbf{b}^{\prime}\right)^{2}+\left(\mathbf{c}^{\prime}\right)^{2}+\mathbf{a}^{\prime} \mathbf{b}^{\prime} \mathbf{c}^{\prime} \tag{5.4.1}
\end{equation*}
$$

then $X\left(\mathbf{a}^{\prime}, \mathbf{b}^{\prime}, \mathbf{c}^{\prime}\right) \subsetneq \mathbb{P}(A \otimes B \otimes C)$, so in this range the substitution methods may be used to prove nontrivial lower bounds for border rank.

Proof. The following is a standard construction in algebraic geometry called an incidence correspondence (see, e.g., [Har95, $\S 6.12]$ for a discussion): Let

$$
\begin{aligned}
& \mathcal{I}:= \\
& \left\{\left(\left(A^{\prime}, B^{\prime}, C^{\prime}\right),[T]\right) \in\left[G\left(\mathbf{a}^{\prime}, A^{*}\right) \times G\left(\mathbf{b}^{\prime}, B^{*}\right) \times G\left(\mathbf{c}^{\prime}, C^{*}\right)\right] \times \mathbb{P}(A \otimes B \otimes C)\right. \\
& \left.\mid A^{\prime} \otimes B^{\prime} \otimes C^{\prime} \subset T^{\perp}\right\}
\end{aligned}
$$

and note that the projection of $\mathcal{I}$ to $\mathbb{P}(A \otimes B \otimes C)$ has image $X\left(\mathbf{a}^{\prime}, \mathbf{b}^{\prime}, \mathbf{c}^{\prime}\right)$. A fiber of the other projection $\mathcal{I} \rightarrow G\left(\mathbf{a}^{\prime}, A^{*}\right) \times G\left(\mathbf{b}^{\prime}, B^{*}\right) \times G\left(\mathbf{c}^{\prime}, C^{*}\right)$ is $\mathbb{P}\left(\left(A^{\prime} \otimes B^{\prime} \otimes C^{\prime}\right)^{\perp}\right)$, a projective space of dimension $\mathbf{a b c}-\mathbf{a}^{\prime} \mathbf{b}^{\prime} \mathbf{c}^{\prime}-1$. Hence:

$$
\operatorname{dim} \mathcal{I}:=\left(\mathbf{a b c}-\mathbf{a}^{\prime} \mathbf{b}^{\prime} \mathbf{c}^{\prime}-1\right)+\left(\mathbf{a}-\mathbf{a}^{\prime}\right) \mathbf{a}^{\prime}+\left(\mathbf{b}-\mathbf{b}^{\prime}\right) \mathbf{b}^{\prime}+\left(\mathbf{c}-\mathbf{c}^{\prime}\right) \mathbf{c}^{\prime} .
$$

Since the map $\mathcal{I} \rightarrow X$ is surjective, this proves the dimension assertion. Since the projection to $\mathbb{P}(A \otimes B \otimes C)$ is a regular map, the Zariski closed assertion also follows.

The proof and examples show that beyond this bound one expects $X\left(\mathbf{a}^{\prime}, \mathbf{b}^{\prime}, \mathbf{c}^{\prime}\right)=\mathbb{P}(A \otimes B \otimes C)$, so that the method cannot be used. Also note that tensors could be quite compressible and still have near maximal border rank, a weakness we already saw with the tensor of (5.1.3) (which also satisfies Strassen's equations).

The inequality in Proposition 5.4.2.2 may be sharp or nearly so. For tensors in $\mathbb{C}^{\mathbf{m}} \otimes \mathbb{C}^{\mathbf{m}} \otimes \mathbb{C}^{\mathbf{m}}$ the limit of this method alone would be a border rank lower bound of $3\left(\mathbf{m}-\sqrt{3 \mathbf{m}+\frac{9}{4}}+\frac{3}{2}\right)$.
5.4.3. How to exploit symmetry. As mentioned above, the border substitution method is particularly useful for tensors $T$ with a large symmetry group $G_{T}$, as one can replace the unknown $A^{\prime}$ by representatives of the closed $G_{T}$-orbits in the Grassmannian. For matrix multiplication, one obtains:
Theorem 5.4.3.1. [LM]

$$
M_{\langle\mathbf{n}\rangle} \in \sigma_{r}\left(\operatorname{Seg}\left(\mathbb{P}^{\mathbf{n}^{2}-1} \times \mathbb{P}^{\mathbf{n}^{2}-1} \times \mathbb{P}^{\mathbf{n}^{2}-1}\right)\right)
$$

if and only if there exist curves $p_{j}(t) \subset \operatorname{Seg}\left(\mathbb{P}^{\mathbf{n}^{2}-1} \times \mathbb{P}^{\mathbf{n}^{2}-1} \times \mathbb{P}^{\mathbf{n}^{2}-1}\right)$ such that for $2 \leq j \leq r, \lim _{t \rightarrow 0} p_{j}(t)=x_{2}^{1} \otimes y_{2}^{1} \otimes z_{2}^{1}$ and $M_{\langle\mathbf{n}\rangle} \in \lim _{t \rightarrow 0}\left\langle x_{2}^{1} \otimes y_{2}^{1} \otimes z_{2}^{1}, p_{2}(t), \ldots, p_{r}(t)\right\rangle$.

In §5.4.5, Theorem 5.4.3.1 is used to improve the lower bounds for border rank.

In this section and the next, I explain the theory. One can also use these methods when attempting to search for new decompositions to limit one's
searches for decompositions with certain normal forms. In order to discuss these methods, I first develop language to discuss the $G_{T}$ orbit closures in the Grassmannian.

To simplify notation, for a tensor $T \in A_{1} \otimes \ldots \otimes A_{k}$, and $\tilde{A} \subset A_{1}$, write

$$
T / \tilde{A}:=\left.T\right|_{\tilde{A}^{\perp} \otimes A_{2}^{*} \otimes \cdots \otimes A_{k}^{*}} \in\left(A_{1} / \tilde{A}\right) \otimes A_{2} \otimes \ldots \otimes A_{k}
$$

Define

$$
B_{\rho, \mathbf{a}^{\prime}}(T):=\left\{\tilde{A} \in G\left(\mathbf{a}^{\prime}, A_{1}\right) \mid \underline{\mathbf{R}}(T / \tilde{A}) \leq \rho\right\} .
$$

Proposition 5.4.3.2. [LM16] The set $B_{\rho, \mathbf{a}^{\prime}}(T)$ is Zariski closed.
The proof requires some standard notions from geometry and can be skipped on a first reading.

A vector bundle $\mathcal{V}$ on a variety $X$ is a variety $\mathcal{V}$ equipped with a surjective regular map $\pi: \mathcal{V} \rightarrow X$ such that for all $x \in X, \pi^{-1}(x)$ is a vector space of dimension $\mathbf{v}$, satisfying certain compatibility conditions (in particular, local triviality: for all $x \in X$, there exists an open subset $U$ containing $x$ such that $\left.\left.\mathcal{V}\right|_{U} \simeq \mathbb{C}^{\mathbf{v}} \times U\right)$. See [Sha13b, §6.1.2] for an algebraic definition or [Spi79, Chap. 3 p 71$]$ for a differential-geometric definition. A section of $\mathcal{V}$ is a regular map $s: X \rightarrow \mathcal{V}$ such that $\pi \circ s=\operatorname{Id}_{X}$.

Two vector bundles over the Grassmannian $G(k, V)$ are ubiquitous: First the tautological subspace bundle $\pi_{\mathcal{S}}: \mathcal{S} \rightarrow G(k, V)$ where $\pi_{\mathcal{S}}{ }^{-1}(E)=E$. This is a vector subbundle of the trivial bundle with fiber $V$, which I denote $\underline{V}$. The tautological quotient bundle $\pi_{\mathcal{Q}}: \mathcal{Q} \rightarrow G(k, V)$ has fiber $\pi_{\mathcal{Q}}{ }^{-1}(E)=$ $V / E$, i.e., we have an exact sequence of vector bundles

$$
0 \rightarrow \mathcal{S} \rightarrow \underline{V} \rightarrow \mathcal{Q} \rightarrow 0
$$

All three bundles are $G L(V)$-homogeneous. See e.g., [Wey03, $\S 3.3$ ] for more details.

For any vector bundle over a projective variety, the corresponding bundle of projective spaces is a projective variety, and a sub-fiber bundle defined by homogeneous equations is also projective.

Proof. Consider the bundle $\pi: \mathcal{Q} \otimes A_{1} \otimes \cdots \otimes A_{k} \rightarrow G\left(\mathbf{a}^{\prime}, A_{1}\right)$, where $\pi^{-1}(\tilde{A})=$ $\left(A_{1} / \tilde{A}\right) \otimes A_{2} \otimes \cdots \otimes A_{k}$. Given $T$, define a natural section $s_{T}: G\left(\mathbf{a}^{\prime}, A_{1}\right) \rightarrow$ $\mathcal{Q} \otimes A_{1} \otimes \cdots \otimes A_{k}$ by $s_{T}(\tilde{A}):=T / \tilde{A}$. Let $X \subset \mathbb{P}\left(\mathcal{Q} \otimes A_{2} \otimes \cdots \otimes A_{k}\right)$ denote the subvariety (that is also a sub-fiber bundle) defined by $X \cap \mathbb{P}\left(\left(A_{1} / \tilde{A}\right) \otimes A_{2} \otimes \cdots \otimes A_{k}\right)=$ $\sigma_{\rho}\left(\operatorname{Seg}\left(\mathbb{P}\left(\left(A_{1} / \tilde{A}\right) \times \mathbb{P} A_{2} \times \cdots \times \mathbb{P} A_{k}\right)\right)\right.$. By the discussion above, $X$ is realizable as a projective variety. Let $\tilde{\pi}: X \rightarrow G\left(\mathbf{a}^{\prime}, A_{1}\right)$ denote the projectivization of $\pi$ restricted to $X$. Then $B_{\rho, \mathbf{a}^{\prime}}(T)=\tilde{\pi}\left(X \cap \mathbb{P}_{s_{T}}\left(G\left(\mathbf{a}^{\prime}, A_{1}\right)\right)\right)$. Since the intersection of two projective varieties is a projective variety, as is the image of a projective variety under a regular map (see Theorem 3.1.4.7), we conclude.

Lemma 5.4.3.3. [LM16] Let $T \in A_{1} \otimes \ldots \otimes A_{k}$ be a tensor, let $G_{T} \subset$ $G L\left(A_{1}\right) \times \cdots \times G L\left(A_{k}\right)$ denote its stabilizer and let $G_{1} \subset G L\left(A_{1}\right)$ denote its projection to $G L\left(A_{1}\right)$. Then $B_{\rho, \mathbf{a}^{\prime}}(T)$ is a $G_{1}$-variety.

Proof. Let $g=\left(g_{1}, \ldots, g_{k}\right) \in G_{T}$. Then $\underline{\mathbf{R}}(T / \tilde{A})=\underline{\mathbf{R}}(g \cdot T / g \cdot \tilde{A})=$ $\underline{\mathbf{R}}\left(T / g_{1} \tilde{A}\right)$.

Recall the definition of a homogeneous variety $X=G / P \subset \mathbb{P} V$ from Definition 4.7.2.1.
Lemma 5.4.3.4. [BL14] Let $X=G / P \subset \mathbb{P} V$ be a homogeneous variety and let $p \in \sigma_{r}(X)$. Then there exist a point $x_{0} \in \hat{X}$ and $r-1$ curves $z_{j}(t) \in \hat{X}$ such that $p \in \lim _{t \rightarrow 0}\left\langle x_{0}, z_{1}(t), \ldots, z_{r-1}(t)\right\rangle$.

Proof. Since $p \in \sigma_{r}(X)$, there exist $r$ curves $x(t), y_{1}(t), \ldots, y_{r-1}(t) \in \hat{X}$ such that

$$
p \in \lim _{t \rightarrow 0} \mathbb{P}\left\langle x(t), y_{1}(t), \ldots, y_{r-1}(t)\right\rangle .
$$

Choose a curve $g_{t} \in G$, such that $g_{t}(x(t))=x_{0}=x(0)$ for all $t$ and $g_{0}=\mathrm{Id}$. We have

$$
\begin{aligned}
\left\langle x(t), y_{1}(t), \ldots, y_{r-1}(t)\right\rangle & =g_{t}^{-1} \cdot\left\langle x_{0}, g_{t} \cdot y_{1}(t), \ldots, g_{t} \cdot y_{r-1}(t)\right\rangle \text { and } \\
\lim _{t \rightarrow 0}\left\langle x(t), y_{1}(t), \ldots, y_{r-1}(t)\right\rangle & =\lim _{t \rightarrow 0}\left(g_{t}{ }^{-1} \cdot\left\langle x_{0}, g_{t} \cdot y_{1}(t), \ldots, g_{t} \cdot y_{r-1}(t)\right\rangle\right) \\
& =\lim _{t \rightarrow 0}\left\langle x_{0}, g_{t} \cdot y_{1}(t), \ldots, g_{t} \cdot y_{r-1}(t)\right\rangle .
\end{aligned}
$$

Set $z_{j}(t)=g_{t} \cdot y_{j}(t)$ to complete the proof.
Exercise 5.4.3.5: (1) Show that if $X$ is a $G$-variety, then any orbit $G \cdot x$ for $x \in X$ of minimal dimension must be Zariski closed. ©

The following Lemma applies both to $M_{\langle\mathbf{n}\rangle}$ and to the determinant polynomial:
Lemma 5.4.3.6 (Normal form lemma). [LM] Let $X=G / P \subset \mathbb{P} V$ be a homogeneous variety and let $v \in V$ be such that $G_{v}:=\{g \in G \mid g[v]=[v]\}$ has a single closed orbit $\mathcal{O}_{\text {min }}$ in $X$. Then any border rank $r$ decomposition of $v$ may be modified using $G_{v}$ to a border rank $r$ decomposition with limit plane $\lim _{t \rightarrow 0}\left\langle x_{1}(t), \ldots, x_{r}(t)\right\rangle$ such that there is a stationary point $x_{1}(t) \equiv x_{1}$ lying in $\mathcal{O}_{\text {min }}$.

If moreover every orbit of $G_{v} \cap G_{x_{1}}$ contains $x_{1}$ in its closure, we may further assume that all other $x_{j}(t)$ limit to $x_{1}$.

Proof. I prove the second statement. By Lemma 5.4.3.4, it is sufficient to show that we can have all points limiting to the same point $x_{1}(0)$.

Work by induction. Say we have shown that $x_{1}(t), \ldots, x_{q}(t)$ all limit to the same point $x_{1} \in \mathcal{O}_{\text {min }}$. It remains to show that our curves can be modified so that the same holds for $x_{1}(t), \ldots, x_{q+1}(t)$. Take a curve $g_{\epsilon} \in G_{v} \cap G_{x_{1}}$ such that $\lim _{\epsilon \rightarrow 0} g_{\epsilon} x_{q+1}(0)=x_{1}$. For each fixed $\epsilon$, acting on the $x_{j}(t)$ by $g_{\epsilon}$, we obtain a border rank decomposition for which $g_{\epsilon} x_{i}(t) \rightarrow$ $g_{\epsilon} x_{1}(0)=x_{1}(0)$ for $i \leq q$ and $g_{\epsilon} x_{q+1}(t) \rightarrow g_{\epsilon} x_{q+1}(0)$. Fix a sequence $\epsilon_{n} \rightarrow 0$. Claim: we may choose a sequence $t_{n} \rightarrow 0$ such that

- $\lim _{n \rightarrow \infty} g_{\epsilon_{n}} x_{q+1}\left(t_{n}\right)=x_{1}(0)$,
- $\lim _{n \rightarrow \infty}<g_{\epsilon_{n}} x_{1}\left(t_{n}\right), \ldots, g_{\epsilon_{n}} x_{r}\left(t_{n}\right)>$ contains $v$ and
- $\lim _{n \rightarrow \infty} g_{\epsilon_{n}} x_{j}\left(t_{n}\right)=x_{1}(0)$ for $j \leq q$.

The first point holds as $\lim _{\epsilon \rightarrow 0} g_{\epsilon} x_{q+1}(0)=x_{1}$. The second follows as for each fixed $\epsilon_{n}$, taking $t_{n}$ sufficiently small we may assure that a ball of radius $1 / n$ centered at $v$ intersects $<g_{\epsilon_{n}} x_{1}\left(t_{n}\right), \ldots, g_{\epsilon_{n}} x_{r}\left(t_{n}\right)>$. In the same way we may assure that the third point is satisfied. Considering the sequence $\tilde{x}_{i}\left(t_{n}\right):=g_{\epsilon_{n}} x_{i}\left(t_{n}\right)$ we obtain the desired border rank decomposition.

Exercise 5.4.3.7: (1) Write out a proof of the first assertion in the normal form lemma.

Applying the normal form lemma to matrix multiplication, in order to prove $\left[M_{\langle\mathbf{n}\rangle}\right] \notin \sigma_{r}(\operatorname{Seg}(\mathbb{P} A \times \mathbb{P} B \times \mathbb{P} C))$, it is sufficient to prove it is not contained in a smaller variety. This variety, called the greater areole is discussed in the next section.
5.4.4. Larger geometric context. Recall that for $X \subset \mathbb{P} V, \sigma_{r}(X)$ may be written as

$$
\sigma_{r}(X)=\bigcup_{x_{j}(t) \subset X, 1 \leq j \leq r}\left\{z \in \mathbb{P} V \mid z \in \lim _{t \rightarrow 0}\left\langle x_{1}(t), \ldots, x_{r}(t)\right\rangle\right\}
$$

where the union is over all curves $x_{j}(t)$ in $X$, including stationary ones. (One can take algebraic or analytic curves.) Remarkably, for the study of certain points such as $M_{\langle\mathbf{n}\rangle}$ and $\operatorname{det}_{n}$ with large symmetry groups, it is sufficient to consider "local" versions of secant varieties.

It is better to discuss Theorem 5.4.3.1 in the larger context of secant varieties, so make the following definition:

Definition 5.4.4.1 (Greater Areole). $[\mathbf{L M}]$ Let $X \subset \mathbb{P} V$ be a projective variety and let $p \in X$. The $r$-th greater areole at $p$ is

$$
\tilde{\mathfrak{a}}_{r}(X, p):=\bigcup_{\substack{x_{j}(t) \subset X \\ x_{j}(t) \rightarrow p}} \lim _{t \rightarrow 0}\left\langle x_{1}(t), \ldots, x_{r}(t)\right\rangle \subset \mathbb{P} V .
$$

Then Theorem 5.4.3.1 may be restated as:

## Theorem 5.4.4.2. [LM]

$$
M_{\langle\mathbf{n}\rangle} \in \sigma_{r}\left(\operatorname{Seg}\left(\mathbb{P}^{\mathbf{n}^{2}-1} \times \mathbb{P}^{\mathbf{n}^{2}-1} \times \mathbb{P}^{\mathbf{n}^{2}-1}\right)\right)
$$

if and only if

$$
M_{\langle\mathbf{n}\rangle} \in \tilde{\mathfrak{a}}_{r}\left(S e g\left(\mathbb{P}^{\mathbf{n}^{2}-1} \times \mathbb{P}^{\mathbf{n}^{2}-1} \times \mathbb{P}^{\mathbf{n}^{2}-1}\right),\left[x_{2}^{1} \otimes y_{2}^{1} \otimes z_{2}^{1}\right]\right)
$$

Exercise 5.4.4.3: (2) Show that when $G / P=v_{n}\left(\mathbb{P}^{n^{2}-1}\right)$ is the Veronese variety and $v=\operatorname{det}_{n}, \mathcal{O}_{\text {min }}=v_{n}(\operatorname{Seg}(\mathbb{P} E \times \mathbb{P} F))$ is the unique closed $G_{\text {det }_{n}}{ }^{-}$ orbit, and every orbit of $G_{\operatorname{det}_{n},\left(x_{1}^{1}\right)^{n}}$ contains $\left(x_{1}^{1}\right)^{n}$ in its closure, so the normal form lemma applies. ©
Exercise 5.4.4.4: (2) When $G / P=\operatorname{Seg}\left(\mathbb{P}\left(U^{*} \otimes V\right) \times \mathbb{P}\left(V^{*} \otimes W\right) \otimes \mathbb{P}\left(W^{*} \otimes U\right)\right) \subset$ $\mathbb{P}\left(\mathbb{C}^{\mathbf{n}^{2}} \otimes \mathbb{C}^{\mathbf{n}^{2}} \otimes \mathbb{C}^{\mathbf{n}^{2}}\right)$ and $v=M_{\langle\mathbf{n}\rangle}$, let
$\mathcal{K}:=\left\{[\mu \otimes v \otimes \nu \otimes w \otimes \omega \otimes u] \in \operatorname{Seg}\left(\mathbb{P} U^{*} \times \mathbb{P} V \times \mathbb{P} V^{*} \times \mathbb{P} W \times \mathbb{P} W^{*} \times \mathbb{P} U\right) \mid \mu(u)=\omega(w)=\nu(v)=0\right\}$.
Show that $\mathcal{K}$ is the unique closed $G_{M_{\langle U, V, W\rangle}}$-orbit in $\operatorname{Seg}(\mathbb{P} A \times \mathbb{P} B \times \mathbb{P} C)$, and every orbit of $G_{M_{\langle U, V, W\rangle}, x_{2}^{1} \otimes y_{2}^{1} \otimes z_{2}^{1}}$ contains $x_{2}^{1} \otimes y_{2}^{1} \otimes z_{2}^{1}$ in its closure. (Of course the same is true for any $k \in \mathcal{K}$.) ©

### 5.4.5. The border rank bound $\underline{\mathbf{R}}\left(M_{\langle\mathbf{n}\rangle}\right) \geq 2 \mathbf{n}^{2}-\left\lceil\log _{2}(\mathbf{n})\right\rceil-1$.

Theorem 5.4.5.1. [LM16] Let $0<m<\mathbf{n}$. Then

$$
\underline{\mathbf{R}}\left(M_{\langle\mathbf{n}, \mathbf{n}, \mathbf{w}\rangle}\right) \geq 2 \mathbf{n w}-\mathbf{w}+m-\left\lfloor\frac{\mathbf{w}\binom{\mathbf{n}-1+m}{m-1}}{\binom{2 \mathbf{n}-2}{\mathbf{n}-1}}\right\rfloor .
$$

In particular, taking $\mathbf{w}=\mathbf{n}$ and $m=\mathbf{n}-\left\lceil\log _{2}(\mathbf{n})\right\rceil-1$,

$$
\underline{\mathbf{R}}\left(M_{\langle\mathbf{n}\rangle}\right) \geq 2 \mathbf{n}^{2}-\left\lceil\log _{2}(\mathbf{n})\right\rceil-1 .
$$

Proof. First observe that the "In particular" assertion follows from the main assertion because, taking $m=\mathbf{n}-c$, we want $c$ such that

$$
\frac{\mathbf{n}\binom{2 \mathbf{n}-1-c}{\mathbf{n}}}{\binom{2 \mathbf{n}-2}{\mathbf{n}-1}}<1 .
$$

This ratio is

$$
\frac{(\mathbf{n}-1) \cdots(\mathbf{n}-c)}{(2 \mathbf{n}-2)(2 \mathbf{n}-3) \cdots(2 \mathbf{n}-c)}=\frac{\mathbf{n}-c}{2^{c-1}} \frac{\mathbf{n}-1}{\mathbf{n}-\frac{2}{2}} \frac{\mathbf{n}-2}{\mathbf{n}-\frac{3}{2}} \frac{\mathbf{n}-3}{\mathbf{n}-\frac{4}{2}} \cdots \frac{\mathbf{n}-c+1}{\mathbf{n}-\frac{c}{2}}
$$

so if $c-1 \geq \log _{2}(\mathbf{n})$ it is less than one.
For the rest of the proof, introduce the following notation: a Young diagram associated to a partition $\lambda=\left(\lambda_{1}, \ldots, \lambda_{\ell}\right)$ is a collection of left aligned boxes, with $\lambda_{j}$ boxes in the $j$-th row. Label it with the upsidedown convention as representing entries in the south-west corner of an $\mathbf{n} \times$ $\mathbf{n}$ matrix. More precisely for $(i, j) \in \lambda$, number the boxes of $\lambda$ by pairs
(row, column), however, number the rows starting from $\mathbf{n}$, i.e. $i=\mathbf{n}$ is the first row. For example

$$
\begin{equation*}
 \tag{5.4.2}
\end{equation*}
$$

is labeled $x=(\mathbf{n}, 1), y=(\mathbf{n}, 2), z=(\mathbf{n}-1,1), w=(\mathbf{n}-2,1)$. Let $\tilde{A}_{\lambda}:=$ $\operatorname{span}\left\{u^{i} \otimes v_{j} \mid(i, j) \in \lambda\right\}$ and write $M_{\langle\mathbf{n}, \mathbf{n}, \mathbf{w}\rangle}^{\lambda}:=M_{\langle\mathbf{n}, \mathbf{n}, \mathbf{w}\rangle} / \tilde{A}_{\lambda}$.

The proof consists of two parts. The first is to show that for any $k<\mathbf{n}$ there exists a Young diagram $\lambda$ with $k$ boxes such that $\underline{\mathbf{R}}\left(M_{\langle\mathbf{n}, \mathbf{n}, \mathbf{w}\rangle}^{\lambda}\right) \leq$ $\underline{\mathbf{R}}\left(M_{\langle\mathbf{n}, \mathbf{n}, \mathbf{w}\rangle}\right)-k$, and this is done by induction on $k$. The second is to use Koszul flattenings to obtain a lower bound on $\underline{\mathbf{R}}\left(M_{\langle\mathbf{n}, \mathbf{n}, \mathbf{w}\rangle}^{\lambda}\right)$ for any $\lambda$.

As usual, write $M_{\langle\mathbf{n}, \mathbf{n}, \mathbf{w}\rangle} \in A \otimes B \otimes C=\left(U^{*} \otimes V\right) \otimes\left(V^{*} \otimes W\right) \otimes\left(W^{*} \otimes U\right)$ where $\mathbf{u}=\mathbf{v}=\mathbf{n}$.

Part 1) First consider the case $k=1$. By Proposition 5.4.1.3 there exists $[a] \in B_{\underline{\mathbf{R}}\left(M_{\langle\mathbf{n}, \mathbf{n}, \mathbf{w}\rangle}\right)-1, \mathbf{n}^{2}-1}\left(M_{\langle\mathbf{n}, \mathbf{n}, \mathbf{w}\rangle}\right)$ such that the reduced tensor drops border rank. The group $G L(U) \times G L(V) \times G L(W)$ stabilizes $M_{\langle\mathbf{n}, \mathbf{n}, \mathbf{w}\rangle}$. Lemma 5.4.3.3 applies with $G_{1}=G L(U) \times G L(V) \subset G L(A)$. Since the $G L(U) \times G L(V)$-orbit closure of any $[a] \in \mathbb{P} A$ contains $\left[u^{\mathbf{n}} \otimes v_{1}\right]$, we may replace $[a]$ by $\left[u^{\mathbf{n}} \otimes v_{1}\right]$.

Now assume that $\underline{\mathbf{R}}\left(M_{\langle\mathbf{n}, \mathbf{n}, \mathbf{w}\rangle}^{\lambda^{\prime}}\right) \leq \underline{\mathbf{R}}\left(M_{\langle\mathbf{n}, \mathbf{n}, \mathbf{w}\rangle}\right)-k+1$, where $\lambda^{\prime}$ has $k-1$ boxes. Again by Proposition 5.4.1.3 there exists $\left[a^{\prime}\right] \in B_{\underline{\mathbf{R}}\left(M_{\langle\mathbf{n}, \mathbf{n}, \mathbf{w}\rangle}\right)-k, \mathbf{n}^{2}-k}\left(M_{\langle\mathbf{n}, \mathbf{n}, \mathbf{w}\rangle}^{\lambda^{\prime}}\right)$ such that when we reduce by $\left[a^{\prime}\right]$ the border rank of the reduced tensor drops. We no longer have the full action of $G L(U) \times G L(V)$. However, the product of parabolic subgroups of $G L(U) \times G L(V)$, which by definition are the subgroups that stabilize the flags in $U^{*}$ and $V$ induced by $\lambda^{\prime}$, stabilizes $M_{\langle\mathbf{n}, \mathbf{n}, \mathbf{w}\rangle}^{\lambda^{\prime}}$. In particular, all parabolic groups are contained in the Borel subgroup of upper-triangular matrices. By the diagonal (torus) action and Lemma 5.4.3.3 we may assume that $a$ has just one nonzero entry outside of $\lambda$. Further, using the upper-triangular (Borel) action we can move the entry south-west to obtain the Young diagram $\lambda$.

For example, when the Young diagram is (5.4.2) with $\mathbf{n}=4$, and we want to move $x_{4}^{1}$ into the diagram, we may multiply it on the left and right respectively by

$$
\left(\begin{array}{ccccc}
\epsilon & & & \\
1 & 1 & & \\
& & 1 & \\
& & & & 1
\end{array}\right) \text { and }\left(\begin{array}{llll}
\epsilon & & & 1 \\
& \epsilon & & \\
& & 1 & \\
& & & 1
\end{array}\right)
$$

where blank entries are zero. Then $x_{4}^{1} \mapsto \epsilon^{2} x_{4}^{1}+\epsilon\left(x_{4}^{2}+x_{1}^{4}\right)+x_{1}^{2}$ and we let $\epsilon \rightarrow 0$.

Part 2) Recall that for the matrix multiplication operator, the Koszul flattening of $\S 2.4$ factors as $M_{\langle\mathbf{n}, \mathbf{n}, \mathbf{w}\rangle}=M_{\langle\mathbf{n}, \mathbf{n}, 1\rangle} \otimes \mathrm{Id}_{W}$, so it will suffice to apply the Koszul flattening to $M_{\langle\mathbf{n}, \mathbf{n}, 1\rangle}^{\lambda} \in\left[\left(U^{*} \otimes V\right) / A_{\lambda}\right] \otimes V^{*} \otimes U$. We need to show that for all $\lambda$ of size $m$,

$$
\underline{\mathbf{R}}\left(M_{\langle\mathbf{n}, \mathbf{n}, 1\rangle}^{\lambda}\right) \geq 2 \mathbf{n}-1-\frac{\binom{\mathbf{n}-1+m}{m-1}}{\binom{2 \mathbf{n}-1}{\mathbf{n}-1}} .
$$

This will be accomplished by restricting to a suitable $A^{\prime} \subset\left[\left(U^{*} \otimes V\right) / A_{\lambda}\right]^{*}$ of dimension $2 \mathbf{n}-1$, such that, setting $\hat{A}=\left(A^{\prime}\right)^{*}$,

$$
\left.\operatorname{rank}\left(\left.M_{\langle\mathbf{n}, \mathbf{n}, 1\rangle}^{\lambda}\right|_{A^{\prime} \otimes V \otimes U^{*}}\right)_{\hat{A}}^{\wedge \mathbf{n}-1}\right) \geq\binom{ 2 \mathbf{n}-1}{\mathbf{n}-1} \mathbf{n}-\binom{\mathbf{n}-1+m}{m-1},
$$

i.e.,

$$
\left.\operatorname{dim} \operatorname{ker}\left(\left.M_{\langle\mathbf{n}, \mathbf{n}, 1\rangle}^{\lambda}\right|_{A^{\prime} \otimes V \otimes U^{*}}\right)_{\hat{A}}^{\wedge \mathbf{n}-1}\right) \leq\binom{\mathbf{n}-1+m}{m-1}
$$

and applying Proposition 2.4.2.1. Since we are working in bases, we may consider $M_{\langle\mathbf{n}, \mathbf{n}, 1\rangle}^{\lambda} \in\left(A / A_{\lambda}\right) \otimes B \otimes C$ in $A \otimes B \otimes C$, with specific coordinates set equal to 0 .

Recall the map $\phi: A \rightarrow \mathbb{C}^{2 \mathbf{n}-1}=\hat{A}$ given by $u^{i} \otimes v_{j} \mapsto e_{i+j-1}$ from (2.5.2) and the other notations from the proof of Theorem 2.5.2.6. The crucial part is to determine how many zeros are added to the diagonal when the entries of $\lambda$ are set to zero. The map $\left(\left.M_{\langle\mathbf{n}, \mathbf{n}, 1\rangle}^{\lambda}\right|_{A^{\prime} \otimes V \otimes U^{*}} ^{*} \hat{A}^{\wedge \mathbf{n}-1}\right.$ is

$$
(S, j):=e_{s_{1}} \wedge \cdots \wedge e_{s_{\mathbf{n}-1}} \otimes v_{j} \mapsto \sum_{\{k \in[\mathbf{n}] \mid(i, j) \notin \lambda\}} e_{j+i-1} \wedge e_{s_{1}} \wedge \cdots \wedge e_{s_{\mathbf{n}-1}} \otimes u^{i} .
$$

Recall that when working with $M_{\langle\mathbf{n}, \mathbf{n}, 1\rangle}$, the diagonal terms in the matrix were indexed by pairs $\left[(S, j)=\left(P \backslash p_{l}, 1+p_{l}-l\right),(P, l)\right]$, in other words that $\left(P \backslash p_{l}, 1+p_{l}-l\right)$ mapped to ( $P, l$ ) plus terms that are lower in the order. So fix $(i, j) \in \lambda$, we need to count the number of terms $(P, i)$ that will not appear anymore as a result of $(i, j)$ being in $\lambda$. That is, fixing $(i, j)$, we need to count the number of $\left(p_{1}, \ldots, p_{i-1}\right)$ with $p_{1}<\cdots<p_{i-1}<i+j-1$, of which there are $\binom{i+j-2}{i-1}$, and multiply this by the number of $\left(p_{i+1}, \ldots, p_{\mathbf{n}}\right)$ with $i+j-1<p_{i+1}<\cdots<p_{\mathbf{n}} \leq 2 \mathbf{n}-1$, of which there are $\binom{2 \mathbf{n}-1-(i+j-1)}{\mathbf{n}-i}$. In summary, each $(i, j) \in \lambda$ kills $g(i, j):=\binom{i+j-1}{i-1}\binom{2 \mathbf{n}-i-j}{\mathbf{n}-i}$ terms on the diagonal. Hence, it is enough to prove that $\sum_{(i, j) \in \lambda} g(i, j) \leq\binom{\mathbf{n}-1+m}{m-1}$.
Exercise 5.4.5.2: (1) Show that $\sum_{j=1}^{m}\binom{\mathbf{n}+j-2}{j-1}=\binom{m+\mathbf{n}-2}{m-1}$.

By Exercise 5.4.5.2 and a similar calculation, we see $\sum_{i=\mathbf{n}}^{\mathbf{n}-m+1} g(i, 1)=$ $\sum_{j=1}^{m} g(\mathbf{n}, j)=\binom{\mathbf{n}-2+m}{m-1}$. So it remains to prove that the Young diagram that maximizes $f_{\lambda}:=\sum_{(i, j) \in \lambda} g(i, j)$ has one row or column. Use induction on the size of $\lambda$, the case $|\lambda|=1$ being trivial. Note that $g(\mathbf{n}-i, j)=$ $g(\mathbf{n}-j, i)$. Moreover, $g(i, j+1) \geq g(i, j)$.

Say $\lambda=\lambda^{\prime} \cup\{(i, j)\}$. By induction it is sufficient to show that:

$$
\begin{equation*}
g(\mathbf{n}, i j)=\binom{\mathbf{n}-1+i j-1}{\mathbf{n}-1} \geq\binom{\mathbf{n}+i-j-1}{i-1}\binom{\mathbf{n}-i+j}{j-1}=g(i, j), \tag{5.4.3}
\end{equation*}
$$

where $\mathbf{n}>i j$.
Exercise 5.4.5.3: (3) Prove the estimate. ©
5.4.6. The boundary case. The proof of Corollary 5.4.6.1 below uses elementary properties of Chern classes and can be skipped by readers unfamiliar with them. Let $\pi_{A}: G\left(\mathbf{a}^{\prime}, A^{*}\right) \times G\left(\mathbf{b}^{\prime}, B^{*}\right) \times G\left(\mathbf{c}^{\prime}, C^{*}\right) \rightarrow G\left(\mathbf{a}^{\prime}, A^{*}\right)$ denote the projection and similarly for $\pi_{B}, \pi_{C}$. Let $\mathcal{E}=\mathcal{E}\left(\mathbf{a}^{\prime}, \mathbf{b}^{\prime}, \mathbf{c}^{\prime}\right):=$ $\pi_{A}^{*}\left(\mathcal{S}_{A}\right) \otimes \pi_{B}^{*}\left(\mathcal{S}_{B}\right) \otimes \pi_{C}^{*}\left(\mathcal{S}_{C}\right)$ be the vector bundle that is the tensor product of the pullbacks of tautological subspace bundles $\mathcal{S}_{A}, \mathcal{S}_{B}, \mathcal{S}_{C}$. In each particular case it is possible to explicitly compute how many different $A^{\prime} \otimes B^{\prime} \otimes C^{\prime}$ a generic hyperplane may contain as follows:

## Corollary 5.4.6.1. [LM16]

(1) If (5.4.1) holds then a generic tensor is ( $\left.\mathbf{a}^{\prime}, \mathbf{b}^{\prime}, \mathbf{c}^{\prime}\right)$-compression generic.
(2) If (5.4.1) does not hold then $\operatorname{rank} \mathcal{E}^{*} \leq \operatorname{dim}\left(G\left(\mathbf{a}^{\prime}, A^{*}\right) \times G\left(\mathbf{b}^{\prime}, B^{*}\right) \times\right.$ $G\left(\mathbf{c}^{\prime}, C^{*}\right)$ ). If the top Chern class of $\mathcal{E}^{*}$ is nonzero, then no tensor is $\left(\mathbf{a}^{\prime}, \mathbf{b}^{\prime}, \mathbf{c}^{\prime}\right)$-compression generic.

Proof. The first assertion is a restatement of Proposition 5.4.2.2.
For the second, notice that $T$ induces a section $\tilde{T}$ of the vector bundle $\mathcal{E}^{*} \rightarrow G\left(\mathbf{a}^{\prime}, A^{*}\right) \times G\left(\mathbf{b}^{\prime}, B^{*}\right) \times G\left(\mathbf{c}^{\prime}, C^{*}\right)$ defined by $\tilde{T}\left(A^{\prime} \otimes B^{\prime} \otimes C^{\prime}\right)=$ $\left.T\right|_{A^{\prime} \otimes B^{\prime} \otimes C^{\prime}}$. The zero locus of $\tilde{T}$ is $\left\{\left(A^{\prime}, B^{\prime}, C^{\prime}\right) \in G\left(\mathbf{a}^{\prime}, A^{*}\right) \times G\left(\mathbf{b}^{\prime}, B^{*}\right) \times\right.$ $\left.G\left(\mathbf{c}^{\prime}, C^{*}\right) \mid A^{\prime} \otimes B^{\prime} \otimes C^{\prime} \subset T^{\perp}\right\}$. In particular, $\tilde{T}$ is non-vanishing if and only if $T$ is $\left(\mathbf{a}^{\prime}, \mathbf{b}^{\prime}, \mathbf{c}^{\prime}\right)$-compression generic. If the top Chern class is nonzero, there cannot exist a non-vanishing section.

### 5.5. Geometry of the Coppersmith-Winograd tensors

As we saw in Chapter 3, in practice, only tensors of minimal, or near minimal border rank have been used to prove upper bounds on the exponent of matrix multiplication. Call a tensor that gives a "good" upper bound for
the exponent via the methods of [Str87, CW90], of high CoppersmithWinograd value or high $C W$-value for short. Ambainis, Filmus and LeGall [AFLG15] showed that taking higher powers of $T_{C W, q}$ when $q \geq 5$ cannot prove $\omega<2.30$ by this method alone. They posed the problem of finding additional tensors of high value. The work in this section was motivated by their problem - to isolate geometric features of the Coppersmith-Winograd tensors and find other tensors with such features. However, it turned out that the features described here, with the exception of a large rank/border rank ratio, actually characterize them. The study is incomplete because the CW-value of a tensor also depends on its presentation, and in different bases a tensor can have quite different CW-values. Moreover, even determining the value in a given presentation still involves some "art" in the choice of a good decomposition, choosing the correct tensor power, estimating the value and probability of each block [Wil].
5.5.1. The Coppersmith-Winograd tensors. Recall the CoppersmithWinograd tensors

$$
\begin{equation*}
T_{q, c w}:=\sum_{j=1}^{q} a_{0} \otimes b_{j} \otimes c_{j}+a_{j} \otimes b_{0} \otimes c_{j}+a_{j} \otimes b_{j} \otimes c_{0} \in \mathbb{C}^{q+1} \otimes \mathbb{C}^{q+1} \otimes \mathbb{C}^{q+1} \tag{5.5.1}
\end{equation*}
$$

and

$$
\begin{align*}
T_{q, C W}:= & \sum_{j=1}^{q}\left(a_{0} \otimes b_{j} \otimes c_{j}+a_{j} \otimes b_{0} \otimes c_{j}+a_{j} \otimes b_{j} \otimes c_{0}\right)  \tag{5.5.2}\\
& +a_{0} \otimes b_{0} \otimes c_{q+1}+a_{0} \otimes b_{q+1} \otimes c_{0}+a_{q+1} \otimes b_{0} \otimes c_{0} \in \mathbb{C}^{q+2} \otimes \mathbb{C}^{q+2} \otimes \mathbb{C}^{q+2},
\end{align*}
$$

both of which have border rank $q+2$.
Written as symmetric tensors (polynomials): $T_{q, c w}=x_{0}\left(\sum_{j=1}^{q} x_{j}^{2}\right)$ and $T_{q, C W}=x_{0}\left(\sum_{j=1}^{q} x_{j}^{2}+x_{0} x_{q+1}\right)$.
Proposition 5.5.1.1. [LM15] $\mathbf{R}\left(T_{q, c w}\right)=2 q+1, \mathbf{R}\left(T_{q, C W}\right)=2 q+3$.
Proof. I first prove the lower bound for $T_{q, c w}$. Apply Proposition 5.3.1.1 to show that the rank of the tensor is at least $2 q-2$ plus the rank of $a_{0} \otimes b_{1} \otimes c_{1}+$ $a_{1} \otimes b_{0} \otimes c_{1}+a_{1} \otimes b_{1} \otimes c_{0}$, which is 3 . An analogous estimate provides the lower bound for $\mathbf{R}\left(T_{q, C W}\right)$. To show that $\mathbf{R}\left(T_{q, c w}\right) \leq 2 q+1$ consider the following rank 1 matrices, whose span contains $T\left(A^{*}\right)$ :

1) $q+1$ matrices with all entries equal to 0 apart from one entry on the diagonal equal to 1 ,
2) $q$ matrices indexed by $1 \leq j \leq q$, with all entries equal to zero apart from the four entries $(0,0),(0, j),(j, 0),(j, j)$ equal to 1 .

Exercise 5.5.1.2: (2) Using the lower bound for $T_{q, c w}$, prove the lower bound for $T_{q, C W}$.

In $\S 5.6$ we saw that $\underline{\mathbf{R}}\left(T_{S T R, q}\right)=q+1$, and by Exercise 5.3.1.7, $\mathbf{R}\left(T_{S T R, q}\right)=$ $2 q$. Strassen's tensor has rank nearly twice the border rank, like the CoppersmithWinograd tensors. So one potential source of high CW-value tensors are tensors with a large gap between rank and border rank.
5.5.2. Extremal tensors. Let $A, B, C=\mathbb{C}^{\mathbf{a}}$. There are normal forms for curves in $\operatorname{Seg}(\mathbb{P} A \times \mathbb{P} B \times \mathbb{P} C)$ up to order a - 1 , namely
$T_{t}=\left(a_{1}+t a_{2}+\cdots+t^{\mathbf{a}-1} a_{\mathbf{a}}+O\left(t^{\mathbf{a}}\right)\right) \otimes\left(b_{1}+t b_{2}+\cdots+t^{\mathbf{a}-1} b_{\mathbf{a}}+O\left(t^{\mathbf{a}}\right)\right) \otimes\left(c_{1}+t c_{2}+\cdots+t^{\mathbf{a}-1} c_{\mathbf{a}}+O\left(t^{\mathbf{a}}\right)\right)$
and if the $a_{j}, b_{j}, c_{j}$ are each linearly independent sets of vectors, call the curve general to order $\mathbf{a}-1$.
Proposition 5.5.2.1. [LM15] Let $T \in A \otimes B \otimes C=\mathbb{C}^{\mathbf{a}} \otimes \mathbb{C}^{\mathbf{a}} \otimes \mathbb{C}^{\mathbf{a}}$. If there exists a curve $T_{t}$ that is general to order a such that

$$
T\left(A^{*}\right)=\left.\frac{d^{\mathbf{a}-1} T_{t}\left(A^{*}\right)}{(d t)^{\mathbf{a}-1}}\right|_{t=0}
$$

then, for suitably chosen $\alpha \in A^{*}$ and bases, $T\left(A^{*}\right) T(\alpha)^{-1}$ is the centralizer of a regular nilpotent element.

Proof. Note that $\left.\frac{d^{q} T_{t}}{(d t)^{q}}\right|_{t=0}=q!\sum_{i+j+k=q+3} a_{i} \otimes b_{j} \otimes c_{k}$, i.e.,

$$
\left.\frac{d^{q} T_{t}\left(A^{*}\right)}{(d t)^{q}}\right|_{t=0}=\left(\begin{array}{ccccccc}
x_{q-2} & x_{q-3} & \cdots & \cdots & x_{1} & 0 & \cdots \\
x_{q-3} & x_{q-4} & \cdots & x_{1} & 0 & \cdots & \cdots \\
\vdots & & & & & & \\
\vdots & . \cdot & & & & & \\
x_{1} & 0 & \cdots & & & & \\
0 & 0 & \cdots & & & & \\
\vdots & \vdots & & & & & \\
0 & 0 & \cdots & & & &
\end{array}\right) .
$$

In particular, each space contains the previous ones, and the last equals

$$
\left(\begin{array}{ccccc}
x_{\mathbf{a}} & x_{\mathbf{a}-1} & \cdots & & x_{1} \\
x_{\mathbf{a}-1} & x_{\mathbf{a}-2} & \cdots & x_{1} & 0 \\
\vdots & \vdots & \ddots & & \\
\vdots & x_{1} & & & \\
x_{1} & 0 & & &
\end{array}\right)
$$

which is isomorphic to the centralizer of a regular nilpotent element.

This provides another, explicit proof that the centralizer of a regular nilpotent element belongs to the closure of diagonalizable algebras.

Note that the Coppersmith-Winograd tensor $T_{\mathbf{a}-2, C W}$ satisfies $\mathbb{P} T\left(A^{*}\right) \cap$ $\operatorname{Seg}(\mathbb{P} B \times \mathbb{P} C)=[X]$ is a single point, and $\mathbb{P} \hat{T}_{[X]} \operatorname{Seg}(\mathbb{P} B \times \mathbb{P} C) \cap \mathbb{P} T\left(A^{*}\right)$ is a $\mathbb{P}^{\mathbf{a}-2}$. It turns out these properties characterize it among $1_{A}$-generic tensors:
Theorem 5.5.2.2. [LM15] Let $T \in A \otimes B \otimes C=\mathbb{C}^{\mathbf{a}} \otimes \mathbb{C}^{\mathbf{a}} \otimes \mathbb{C}^{\mathbf{a}}$ be of border rank $\mathbf{a}>2$. Assume $\mathbb{P} T\left(A^{*}\right) \cap \operatorname{Seg}(\mathbb{P} B \times \mathbb{P} C)=[X]$ is a single point, and $\mathbb{P} \hat{T}_{[X]} \operatorname{Seg}(\mathbb{P} B \times \mathbb{P} C) \supset \mathbb{P} T\left(A^{*}\right)$. Then $T$ is not $1_{A}$-generic.

If
(i) $\mathbb{P} T\left(A^{*}\right) \cap \operatorname{Seg}(\mathbb{P} B \times \mathbb{P} C)=[X]$ is a single point,
(ii) $\mathbb{P} \hat{T}_{[X]} \operatorname{Seg}(\mathbb{P} B \times \mathbb{P} C) \cap \mathbb{P} T\left(A^{*}\right)$ is a $\mathbb{P}^{\text {a-2 }}$, and
(iii) $T$ is $1_{A}$-generic,
then $T$ is isomorphic to the Coppersmith-Winograd tensor $T_{\mathbf{a}-2, C W}$.
Proof. For the first assertion, no element of $\mathbb{P} \hat{T}_{[X]} S e g(\mathbb{P} B \times \mathbb{P} C)$ has rank greater than two.

For the second, I first show that $T$ is 1-generic. Choose bases such that $X=b_{1} \otimes c_{1}$, then, after modifying the bases, the $\mathbb{P}^{\mathbf{a}-2}$ must be the projectivization of

$$
E:=\left(\begin{array}{ccccc}
x_{1} & x_{2} & \cdots & x_{\mathbf{a}-1} & 0  \tag{5.5.3}\\
x_{2} & & & & \\
\vdots & & & & \\
x_{\mathbf{a}-1} & & & & \\
0 & & & &
\end{array}\right) .
$$

(Rank one tangent vectors cannot appear by property (i).)
Write $T\left(A^{*}\right)=\operatorname{span}\{E, M\}$ for some matrix $M$. As $T$ is $1_{A}$-generic we can assume that $M$ is invertible. In particular, the last row of $M$ must contain a nonzero entry. In the basis order where $M$ corresponds to $T\left(\alpha^{\mathbf{a}}\right)$, the space of matrices $T\left(B^{*}\right)$ has triangular form and contains matrices with nonzero diagonal entries. The proof for $T\left(C^{*}\right)$ is analogous, hence $T$ is 1-generic.

By Proposition 5.1.5.1 we may assume that $T\left(A^{*}\right)$ is contained in the space of symmetric matrices. Hence, we may assume that $E$ is as above and $M$ is a symmetric matrix. By further changing the bases we may assume that $M$ has:
(1) the first row and column equal to zero, apart from their last entries that are nonzero (we may assume they are equal to 1 ),
(2) the last row and column equal to zero apart from their first entries.

Hence the matrix $M$ is determined by a submatrix $M^{\prime}$ of rows and columns 2 to $\mathbf{a}-1$. As $T\left(A^{*}\right)$ contains a matrix of maximal rank, the matrix $M^{\prime}$ must have rank $\mathbf{a}-2$. We can change the basis $\alpha^{2}, \ldots, \alpha^{\mathbf{a}-1}$ in such a way that the quadric corresponding to $M^{\prime}$ equals $x_{2}^{2}+\cdots+x_{\mathbf{a}-1}^{2}$. This will also change the other matrices, which correspond to quadrics $x_{1} x_{i}$ for $1 \leq i \leq \mathbf{a}-1$, but will not change the space that they span. We obtain the tensor $T_{\mathbf{a}-2, C W}$.
5.5.3. Compression extremality. In this subsection I discuss tensors for which the border substitution method fails miserably. In particular, although the usual substitution method correctly determines the rank of the Coppersmith-Winograd tensors, the tensors are special in that they are nearly characterized by the failure of the border substitution method to give lower border rank bounds.

Definition 5.5.3.1. A 1-generic, tensor $T \in A \otimes B \otimes C$ is said to be maximally compressible if there exists hyperplanes $H_{A} \subset A^{*}, H_{B} \subset B^{*}, H_{C} \subset C^{*}$ such that $\left.T\right|_{H_{A} \times H_{B} \times H_{C}}=0$.

If $T \in S^{3} A \subset A \otimes A \otimes A, T$ is maximally symmetric compressible if there exists a hyperplane $H_{A} \subset A^{*}$ such that $\left.T\right|_{H_{A} \times H_{A} \times H_{A}}=0$.

Recall from Proposition 5.1.5.1 that a tensor $T \in \mathbb{C}^{\mathbf{a}} \otimes \mathbb{C}^{\mathbf{a}} \otimes \mathbb{C}^{\mathbf{a}}$ that is 1-generic and satisfies Strassen's equations, with suitable choices of bases becomes a tensor in $S^{3} \mathbb{C}^{\text {a }}$.
Theorem 5.5.3.2. [LM15] Let $T \in S^{3} \mathbb{C}^{\mathbf{a}}$ be 1-generic and maximally symmetric compressible. Then $T$ is one of:
(1) $T_{\mathbf{a}-1, c w}$
(2) $T_{\mathbf{a}-2, C W}$
(3) $T=a_{1}\left(a_{1}^{2}+\cdots a_{\mathbf{a}}^{2}\right)$.

In particular, the only 1-generic, maximally symmetric compressible, minimal border rank tensor in $\mathbb{C}^{\mathbf{a}} \otimes \mathbb{C}^{\mathbf{a}} \otimes \mathbb{C}^{\mathbf{a}}$ is isomorphic to $T_{\mathbf{a}-2, C W}$.

Proof. Let $a_{1}$ be a basis of the line $H_{A}{ }^{\perp} \subset \mathbb{C}^{\text {a }}$. Then $T=a_{1} Q$ for some $Q \in S^{2} \mathbb{C}^{\mathbf{a}}$. By 1-genericity, the rank of $Q$ is either a or $\mathbf{a}-1$. If the rank is $\mathbf{a}$, there are two cases, either the hyperplane $H_{A}$ is tangent to $Q$, or it intersects it transversely. The second is case (3). The first has a normal form $a_{1}\left(a_{1} a_{\mathbf{a}}+a_{2}^{2}+\cdots+a_{\mathbf{a}-1}^{2}\right)$, which, when written as a tensor, is $T_{\mathbf{a}-2, C W}$. If $Q$ has rank a-1, by 1 -genericity, $\operatorname{ker}\left(Q_{1,1}\right)$ must be in $H_{A}$ and thus we may choose coordinates such that $Q=\left(a_{2}^{2}+\cdots+a_{\mathbf{a}}^{2}\right)$, but then $T$, written as a tensor, is $T_{\mathbf{a}-1, c w}$.

Proposition 5.5.3.3. [LM15] The Coppersmith-Winograd tensor $T_{\mathbf{a}-2, C W}$ is the unique up to isomorphism 1-generic tensor in $\mathbb{C}^{\mathbf{a}} \otimes \mathbb{C}^{\mathbf{a}} \otimes \mathbb{C}^{\mathbf{a}}$ that is maximally compressible and satisfies any of the following:
(1) Strassen's equations
(2) cyclic $\mathbb{Z}_{3}$-invariance
(3) has border rank a.

Proof. Let $a_{1}, \ldots, a_{\mathbf{a}}$ be a basis of $A$ with $H_{A}=a_{1}{ }^{\perp}$ and similarly for $H_{B}=b_{1}{ }^{\perp}$ and $H_{C}=c_{1}{ }^{\perp}$. Thus (allowing re-ordering of the factors $A, B, C$ ) $T=a_{1} \otimes X+b_{1} \otimes Y+c_{1} \otimes Z$ where $X \in B \otimes C, Y \in A \otimes C, Z \in A \otimes B$. Now no $\alpha \in H_{A}$ can be such that $T(\alpha)$ is of maximal rank, as for any $\beta_{1}, \beta_{2} \in H_{B}$, $T\left(\alpha, \beta_{j}\right) \subset \mathbb{C}\left\{c_{1}\right\}$. So $T\left(a^{1}\right), T\left(b^{1}\right), T\left(c^{1}\right)$ are all of rank $\mathbf{a}$, where $a^{1}$ is the dual basis vector to $a_{1}$ etc. After a modification, we may assume $X$ has rank a.

Let $(g, h, k) \in G L(A) \times G L(B) \times G L(C)$. We may normalize $X=\mathrm{Id}$, which forces $g=h$. We may then rewrite $X, Y, Z$ such that $Y$ is full rank and renormalize

$$
X=Y=\left(\begin{array}{ll}
\frac{1}{3} & \\
& \mathrm{Id}_{\mathbf{a}-1}
\end{array}\right),
$$

which forces $h=k$ and uses up our normalizations.
Now we use any of the above three properties. The weakest is the second, but by $\mathbb{Z}_{3}$-invariance, if $X=Y$, we must have $Z=X=Y$ as well and $T$ is the Coppersmith-Winograd tensor. The other two imply the second by Proposition 5.1.5.1.

Remark 5.5.3.4. Theorem 5.5.3.2 and Proposition 5.5.3.3 were motivated by the suggestion in [AFLG15] to look for tensors to replace the CoppersmithWinograd tensor in Strassen's laser method. We had hoped to isolate geometric properties of the tensor and then find other tensors with similar properties, to then test the method on. However, the properties we found, with the exception of a large rank/border rank ratio, essentially characterized the tensors.

### 5.6. Ranks and border ranks of Structure tensors of algebras

In this section I discuss ranks and border ranks of a class of tensors that appear to be more tractable than arbitrary tensors: structure tensors of algebras. It turns out this class is larger than appears at first glance: as explained in §5.6.1, all tensors in $A \otimes B \otimes C=\mathbb{C}^{\mathbf{m}} \otimes \mathbb{C}^{\mathbf{m}} \otimes \mathbb{C}^{\mathrm{m}}$ that are $1_{A}$ and $1_{B}$-generic are equivalent to structure tensors of algebras with unit. In $\S 5.6 .2$, I show structure tensors corresponding to algebras of the form $\mathbb{C}\left[x_{1}, \ldots, x_{n}\right] / \mathcal{I}$, where $\mathcal{I}$ is an ideal whose zero set is finite, are equivalent
to symmetric tensors and give several examples. (For those familiar with the language, these are structure tensors of coordinate rings of zero dimensional affine schemes, see $\S 10.1 .1$.) The algebra structure can facilitate the application of the substitution and border substitution methods, as is illustrated in $\S 5.6 .3$ and $\S 5.6 .4$ respectively. In particular, using algebras of the form $\mathbb{C}\left[x_{1}, \ldots, x_{n}\right] / \mathcal{I}$, I present examples of tensors with rank to border rank ratio approaching three. I conclude with Bläser and Lysikov's study of structure tensors of algebras that have minimal border rank.

Throughout this section $\mathcal{A}$ denotes a finite dimensional associative algebra and $T_{\mathcal{A}} \in \mathcal{A}^{*} \otimes \mathcal{A}^{*} \otimes \mathcal{A}$ denotes its structure tensor as discussed in §3.5.1.
5.6.1. Algebras and minimal border rank tensors. The following reduction theorem is due to Bläser and Lysikov:
Theorem 5.6.1.1. [BL16] Let $\mathcal{A}, \mathcal{A}_{1}$ be algebras of dimension $\mathbf{m}$ with structure tensors $T_{\mathcal{A}}, T_{\mathcal{A}_{1}}$. Then $T_{\mathcal{A}} \subset \overline{G L_{\mathbf{m}}^{\times 3} \cdot T_{\mathcal{A}_{1}}}$ if and only if $T_{\mathcal{A}} \subset$ $\overline{G L_{\mathbf{m}} \cdot T_{\mathcal{A}_{1}}}$.

Proof. Write $\mathbb{C}^{\mathrm{m}} \simeq \mathcal{A} \simeq \mathcal{A}_{1}$ as a vector space, so $T_{\mathcal{A}}, T_{\mathcal{A}_{1}} \in \mathbb{C}^{\mathrm{m} *} \otimes \mathbb{C}^{\mathrm{m} *} \otimes \mathbb{C}^{\mathrm{m}}$ Write $T_{\mathcal{A}}=\lim _{t \rightarrow 0} T_{t}$, where $T_{t}:=\left(f_{t}, g_{t}, h_{t}\right) \cdot T_{\mathcal{A}_{1}}$, with $f_{t}, g_{t}, h_{t}$ curves in $G L_{\mathbf{m}}$. Let $e \in \mathcal{A}$ denote the identity element. Then $T_{t}(e, y)=h_{t} T_{\mathcal{A}_{1}}\left(f_{t}^{-1} e, g_{t}{ }^{-1} y\right)=$ $y+O(t)$. Write $L_{f_{t}-1}: \mathcal{A} \rightarrow \mathcal{A}$ for $f_{t}^{-1} e$ considered as a linear map. Then $h_{t} L_{f_{t}-1 e} g_{t}^{-1}=\operatorname{Id}+O(t)$ so we may replace $g_{t}$ by $\tilde{g}_{t}:=h_{t} L_{f_{t}-1 e}$. Similarly, using that $T_{t}(y, e)=y+O(t)$, we may replace $f_{t}$ by $\tilde{f}_{t}:=h_{t} R_{\tilde{g}_{t}-1}$ e where $R$ is used to remind us that it corresponds to right multiplication in the algebra, so our new curve is $T_{\mathcal{A}}=\lim _{t \rightarrow 0}\left(\left(R_{\tilde{g}_{t}-1}\right)^{-1} h_{t}^{-1},\left(L_{f_{t}-1}\right)^{-1} h_{t}{ }^{-1}, h_{t}\right) \cdot T_{\mathcal{A}_{1}}$. Finally, noting that for any linear maps $X, Y \in \operatorname{End}\left(\mathbb{C}^{\mathbf{m}}\right), T_{\mathcal{A}}(X y, Y z)=$ $X T_{\mathcal{A}}(y, z) Y$, and taking $X_{t}=L_{f_{t}-1} e^{-1}, Y_{t}=R_{\tilde{g}_{t}-1} e^{-1}$, our new action is by $h_{t} L_{f_{t}{ }^{-1} e} R_{\tilde{g}_{t}{ }^{-1} e} \in G L_{\mathbf{m}} \subset G L_{\mathbf{m}}^{\times 3}$.

Proposition 5.6.1.2. [BL16] Let $T \in A \otimes B \otimes C=\mathbb{C}^{\mathbf{m}} \otimes \mathbb{C}^{\mathbf{m}} \otimes \mathbb{C}^{\mathbf{m}}$ be $1_{A}$ and $1_{B}$ generic. Then there exists an algebra $\mathcal{A}$ with unit such that $T$ is equivalent to $T_{\mathcal{A}}$, i.e., they are in the same $G L_{\mathbf{m}}^{\times 3}$-orbit.

Proof. Take $\alpha \in A^{*}, \beta \in B^{*}$ with $T(\alpha): B^{*} \rightarrow C$ and $T(\beta): A^{*} \rightarrow C$ of full rank. Give $C$ the algebra structure $c_{1} \cdot c_{2}:=T\left(T(\beta)^{-1} c_{1}, T(\alpha)^{-1} c_{2}\right)$ and note that the structure tensor of this algebra is in the same $G L_{\mathbf{m}}^{\times 3}$-orbit as $T$.

Exercise 5.6.1.3: (1) Verify that the product above indeed gives $C$ the structure of an algebra with unit.

Combining Theorem 5.6.1.1 and Proposition 5.6.1.2, we obtain:

Theorem 5.6.1.4. [BL16] Let $T \in A \otimes B \otimes C=\mathbb{C}^{\mathbf{m}} \otimes \mathbb{C}^{\mathbf{m}} \otimes \mathbb{C}^{\mathrm{m}}$ be $1_{A}$ and $1_{B}$ generic. Take $\alpha \in A^{*}, \beta \in B^{*}$ with $T(\alpha) \in B \otimes C, T(\beta) \in A \otimes C$ of full rank, and use them to construct an equivalent tensor $\tilde{T} \in C^{*} \otimes C^{*} \otimes C$. Then $\underline{\mathbf{R}}(T)=\mathbf{m}$, i.e., $T \in \overline{G L(A) \times G L(B) \times G L(C) \cdot M_{\langle 1\rangle}^{\oplus \mathbf{m}}}$, if and only if $\tilde{T} \in \overline{G L(C) \cdot M_{\langle 1\rangle}^{\oplus \mathbf{m}}}$.

Recall the Comon conjecture from $\S 4.1 .4$ that posits that for symmetric tensors, $\mathbf{R}(T)=\mathbf{R}_{S}(T)$. One can define a border rank version:
Conjecture 5.6.1.5. [Border rank Comon conjecture] [BGL13] Let $T \in$ $S^{3} \mathbb{C}^{\mathbf{m}} \subset\left(\mathbb{C}^{\mathbf{m}}\right)^{\otimes 3}$. Then $\underline{\mathbf{R}}(T)=\underline{\mathbf{R}}_{S}(T)$.

Theorem 5.6.1.4 combined with Proposition 5.1.5.1, which says that minimal border rank 1-generic tensors are symmetric, implies:
Proposition 5.6.1.6. The border rank Comon conjecture holds for 1-generic tensors of minimal border rank.
5.6.2. Structural tensors of algebras of the form $\mathbb{C}\left[x_{1}, \ldots, x_{n}\right] / \mathcal{I}$. Let $\mathcal{I} \subset \mathbb{C}\left[x_{1}, \ldots, x_{n}\right]$ be an ideal whose zero set in affine space is finite, so that $\mathcal{A}_{\mathcal{I}}:=\mathbb{C}\left[x_{1}, \ldots, x_{n}\right] / \mathcal{I}$ is a finite dimensional algebra. Let $\left\{p_{I}\right\}$ be a basis of $\mathcal{A}_{\mathcal{I}}$ with dual basis $\left\{p_{I}^{*}\right\}$ We can write the structural tensor of $\mathcal{A}_{\mathcal{I}}$ as

$$
T_{\mathcal{A}_{\mathcal{I}}}=\sum_{p_{I}, p_{J} \in \mathcal{A}_{\mathcal{I}}} p_{I}^{*} \otimes p_{J}^{*} \otimes\left(p_{I} p_{J} \bmod \mathcal{I}\right) .
$$

This tensor is transparently in $S^{2} \mathcal{A}^{*} \otimes \mathcal{A}$.
Given an algebra $\mathcal{A}=\mathcal{A}_{\mathcal{I}} \in S^{2} \mathcal{A}^{*} \otimes \mathcal{A}$ defined by an ideal as above, note that since $T_{\mathcal{A}}(1, \cdot) \in \operatorname{End}(\mathcal{A})$ and $T_{\mathcal{A}}(\cdot, 1) \in \operatorname{End}(\mathcal{A})$ have full rank and the induced isomorphism $B^{*} \rightarrow C$ is just $\left(\mathcal{A}^{*}\right)^{*} \rightarrow \mathcal{A}$, and similarly for the isomorphism $A^{*} \rightarrow C$, and since the algebra is abelian Strassen's equations are satisfied, so by Proposition 5.1.5.1 there exists a choice of bases such that $T_{\mathcal{A}} \in S^{3} \mathcal{A}$.
Proposition 5.6.2.1. [Michalek and Jelisiejew, personal communication] Structural tensors of algebras of the form $\mathcal{A}=\mathbb{C}\left[x_{1}, \ldots, x_{n}\right] / \mathcal{I}$ are symmetric if either of the following equivalent conditions hold:

- $\mathcal{A}^{*}=\mathcal{A} \cdot f$ for some $f \in \mathcal{A}^{*}$, where for $a, b \in \mathcal{A},(a \cdot f)(b):=f(a b)$.
- $T_{\mathcal{A}}$ is 1-generic.

Proof. We have already seen that if $T_{\mathcal{A}}$ is 1-generic and satisfies Strassen's equations, then $T_{\mathcal{A}}$ is symmetric.

The following are clearly equivalent for an element $f \in \mathcal{A}^{*}$ :

1) $T_{\mathcal{A}}(f) \in \mathcal{A}^{*} \otimes \mathcal{A}^{*}$ is of full rank,
2) the pairing $\mathcal{A} \otimes \mathcal{A} \rightarrow \mathbb{C}$ given by $(a, b) \mapsto f(a b)$ is non-degenerate,
3) $\mathcal{A} f=\mathcal{A}^{*}$.

Remark 5.6.2.2. The condition that $\mathcal{A}^{*}=\mathcal{A} \cdot f$ for some $f \in \mathcal{A}^{*}$ is called Gorenstein. There are numerous definitions of Gorenstein. One that is relevant for Chapter 10 is that $\mathcal{A}$ is Gorenstein if and only if $\mathcal{A}$ is the annhilator of some polynomial $D$ in the dual space, i.e., $D \in \mathbb{C}\left[\frac{\partial}{\partial x_{1}}, \ldots, \frac{\partial}{\partial x_{n}}\right]$.
Example 5.6.2.3. [Zui15] Consider $\mathcal{A}=\mathbb{C}[x] /\left(x^{2}\right)$, with basis 1 , $x$, so

$$
T_{\mathcal{A}}=1^{*} \otimes 1^{*} \otimes 1+x^{*} \otimes 1^{*} \otimes x+1^{*} \otimes x^{*} \otimes x .
$$

Writing $e_{0}=1^{*}, e_{1}=x^{*}$ in the first two factors and $e_{0}=x, e_{1}=1$ in the third,

$$
T_{\mathcal{A}}=e_{0} \otimes e_{0} \otimes e_{1}+e_{1} \otimes e_{0} \otimes e_{0}+e_{0} \otimes e_{1} \otimes e_{0}
$$

That is, $T_{\mathcal{A}}=T_{W \text { State }}$ is a general tangent vector to $\operatorname{Seg}(\mathbb{P} A \times \mathbb{P} B \times \mathbb{P} C)$.
More generally, consider $\mathcal{A}=\mathbb{C}\left[x_{1}, \ldots, x_{n}\right] /\left(x_{1}^{2}, \ldots, x_{n}^{2}\right)$, with basis $x_{I}=$ $x_{i_{1}} \cdots x_{i_{|I|}}$, where $1 \leq i_{1}<\cdots<i_{|I|} \leq n$, and by convention $x_{\emptyset}=1$. Then

$$
T_{\mathcal{A}}=\sum_{I, J \subset[n] \mid I \cap J=\emptyset} x_{I}^{*} \otimes x_{J}^{*} \otimes x_{I \cup J} .
$$

Similar to above, let $e_{I}=x_{I}^{*}$ in the first two factors and $e_{I}=x_{[n] \backslash I}$ in the third, we obtain

$$
T_{\mathcal{A}}=\sum_{\substack{\begin{subarray}{c}{I U J \cup K=[n]] \\
\left\{I, J,\left.K\right|_{|I|+|J|+|K|=n} ^{|l|}\right\}} }}\end{subarray}} e_{I} \otimes e_{J} \otimes e_{K}
$$

so we explicitly see $T_{\mathcal{A}} \in S^{3} \mathbb{C}^{2^{n}}$.
Exercise 5.6.2.4: (2) Show that for $\mathcal{A}=\mathbb{C}\left[x_{1}, \ldots, x_{n}\right] /\left(x_{1}^{2}, \ldots, x_{n}^{2}\right), T_{\mathcal{A}} \simeq$ $T_{W \text { State }}^{\otimes n}$, where for $T \in A \otimes B \otimes C$, consider $T^{\otimes n} \in\left(A^{\otimes n}\right) \otimes\left(B^{\otimes n}\right) \otimes\left(C^{\otimes n}\right)$ as a three-way tensor.
Exercise 5.6.2.5: (2) Let $\mathcal{A}=\mathbb{C}[x] /\left(x^{n}\right)$. Show that $T_{\mathcal{A}}(\mathcal{A}) T_{\mathcal{A}}(1)^{-1} \subset$ $\operatorname{End}(\mathcal{A})$ corresponds to the centralizer of a regular nilpotent element, so in particular $\underline{\mathbf{R}}\left(T_{\mathcal{A}}\right)=n$ and $\mathbf{R}\left(T_{\mathcal{A}}\right)=2 n-1$ by Exercise 5.3.1.8 and Proposition 5.1.4.6.
Exercise 5.6.2.6: (2) Fix natural numbers $d_{1}, \ldots, d_{n}$. Let $\mathcal{A}=\mathbb{C}\left[x_{1}, \ldots, x_{n}\right] /\left(x_{1}^{d_{1}}, \ldots, x_{n}^{d_{n}}\right)$.
Find an explicit identification $\mathcal{A}^{*} \rightarrow \mathcal{A}$ that renders $T_{\mathcal{A}} \in S^{3} \mathcal{A}$. ©
Example 5.6.2.7. [Zui15] Consider the tensor
$T_{\text {WState }, k}=a_{1,0} \otimes \cdots \otimes a_{k-1,0} \otimes a_{k, 1}+a_{1,0} \otimes \cdots \otimes a_{k-2,0} \otimes a_{k-1,1} \otimes a_{k, 0}+\cdots+a_{1,1} \otimes a_{2,0} \otimes \cdots \otimes a_{k, 0}$
that corresponds to a general tangent vector to $\operatorname{Seg}\left(\mathbb{P}^{1} \times \cdots \times \mathbb{P}^{1}\right) \in \mathbb{P}\left(\left(\mathbb{C}^{2}\right)^{\otimes k}\right)$. (Note that $\left.T_{W \text { State }}=T_{W \text { State }, 3}.\right)$ This tensor is called the generalized $W$ state by physicists. Let $\mathcal{A}_{d, N}=\left(\mathbb{C}[x] /\left(x^{d}\right)\right)^{\otimes N} \simeq \mathbb{C}\left[x_{1}, \ldots, x_{N}\right] /\left(x_{1}^{d}, \ldots, x_{N}^{d}\right)$.
Exercise 5.6.2.8: (2) Show that $T_{\mathcal{A}_{d, N}}=\left(T_{W \text { State, } d}\right)^{\otimes N}$.

Example 5.6.2.9 (The Coppersmith-Winograd tensor). [LM, BL16] Consider the algebra

$$
\mathcal{A}_{C W, q}=\mathbb{C}\left[x_{1}, \ldots, x_{q}\right] /\left(x_{i} x_{j}, x_{i}^{2}-x_{j}^{2}, x_{i}^{3}, i \neq j\right)
$$

Let $\left\{1, x_{i},\left[x_{1}^{2}\right]\right\}$ be a basis of $\mathcal{A}$, where $\left[x_{1}^{2}\right]=\left[x_{j}^{2}\right]$ for all $j$. Then

$$
\begin{aligned}
T_{\mathcal{A}_{C W, q}}= & 1^{*} \otimes 1^{*} \otimes 1+\sum_{i=1}^{q}\left(1^{*} \otimes x_{i}^{*} \otimes x_{i}+x_{i}^{*} \otimes 1^{*} \otimes x_{i}\right) \\
& +x_{i}^{*} \otimes x_{i}^{*} \otimes\left[x_{1}^{2}\right]+1^{*} \otimes\left[x_{1}^{2}\right]^{*} \otimes\left[x_{1}^{2}\right]+\left[x_{1}^{2}\right]^{*} \otimes 1^{*} \otimes\left[x_{1}^{2}\right] .
\end{aligned}
$$

Set $e_{0}=1^{*}, e_{i}=x_{i}^{*}, e_{q+1}=\left[x_{1}^{2}\right]^{*}$ in the first two factors and $e_{0}=\left[x_{1}^{2}\right]$, $e_{i}=x_{i}, e_{q+1}=1$ in the third to obtain

$$
\begin{aligned}
T_{\mathcal{A}_{C W, q}}= & T_{C W, q}=e_{0} \otimes e_{0} \otimes e_{q+1}+\sum_{i=1}^{q}\left(e_{0} \otimes e_{i} \otimes e_{i}+e_{i} \otimes e_{0} \otimes e_{i}+e_{i} \otimes e_{i} \otimes e_{0}\right) \\
& +e_{0} \otimes e_{q+1} \otimes e_{0}+e_{q+1} \otimes e_{0} \otimes e_{0}
\end{aligned}
$$

so we indeed obtain the Coppersmith-Winograd tensor.
When is the structure tensor of $\mathcal{A}_{\mathcal{I}}$ of minimal border rank? Note that if $T \in \mathbb{C}^{\mathbf{m}} \otimes \mathbb{C}^{\mathbf{m}} \otimes \mathbb{C}^{\mathbf{m}}$ is the structure tensor of an algebra $\mathcal{A}$ that is a degeneration of $(\mathbb{C}[x] /(x))^{\oplus \mathbf{m}}$ (whose structure tensor is $M_{\langle 1\rangle}^{\oplus \mathbf{m}}$ ), then $\underline{\mathbf{R}}(T)=\mathbf{m}$.
5.6.3. The substitution method applied to structure tensors of algebras. Let $\mathcal{A}$ be a finite dimensional associative algebra. The radical of $\mathcal{A}$ is the intersection of all maximal left ideals and denoted $\operatorname{Rad}(\mathcal{A})$. When $\mathcal{A}$ is abelian, the radical is often call the nilradical.
Exercise 5.6.3.1: (2) Show that every element of $\operatorname{Rad}(\mathcal{A})$ is nilpotent and that if $\mathcal{A}$ is abelian, $\operatorname{Rad}(\mathcal{A})$ consists exactly of the nilpotent elements of $\mathcal{A}$. (This exercise requires knowledge of standard notions from algebra.) ©

Theorem 5.6.3.2. [Blä00, Thm. 7.4] For any integers $p, q \geq 1$,

$$
\mathbf{R}\left(T_{\mathcal{A}}\right) \geq \operatorname{dim}\left(\operatorname{Rad}(\mathcal{A})^{p}\right)+\operatorname{dim}\left(\operatorname{Rad}(\mathcal{A})^{q}\right)+\operatorname{dim} \mathcal{A}-\operatorname{dim}\left(\operatorname{Rad}(\mathcal{A})^{p+q-1}\right) .
$$

For the proof we will need the following Lemma, whose proof I skip:
Lemma 5.6.3.3. [Blä00, Lem. 7.3] Let $\mathcal{A}$ be a finite dimensional algebra, let $U, V \subseteq \mathcal{A}$ be vector subspaces such that $U+\operatorname{Rad}(\mathcal{A})^{p}=\mathcal{A}$ and $V+$ $\operatorname{Rad}(\mathcal{A})^{q}=\mathcal{A}$. Then $\langle U V\rangle+\operatorname{Rad}(\mathcal{A})^{p+q-1}=\mathcal{A}$.

Proof of Theorem 5.6.3.2. Use Proposition 5.4.1.2 with

$$
\begin{aligned}
& \tilde{A}=\left(\operatorname{Rad}(\mathcal{A})^{p}\right)^{\perp} \subset \mathcal{A}^{*}, \\
& \tilde{B}=\left(\operatorname{Rad}(\mathcal{A})^{q}\right)^{\perp} \subset \mathcal{A}^{*}, \text { and } \\
& \tilde{C}=\operatorname{Rad}(\mathcal{A})^{p+q-1} \subset \mathcal{A} .
\end{aligned}
$$

Then observe that any $A^{\prime} \subset \mathcal{A} \backslash \operatorname{Rad}(\mathcal{A})^{p}, B^{\prime} \subset \mathcal{A} \backslash \operatorname{Rad}(\mathcal{A})^{q}$, can play the roles of $U, V$ in the Lemma, so $T_{\mathcal{A}}\left(A^{\prime}, B^{\prime}\right) \not \subset \operatorname{Rad}(\mathcal{A})^{p+q-1}$. Since $C^{\prime} \subset$ $\mathcal{A}^{*} \backslash\left(\operatorname{Rad}(\mathcal{A})^{p+q-1}\right)^{\perp}$, we conclude.

Remark 5.6.3.4. Theorem 5.6.3.2 illustrates the power of the (rank) substitution method over the border substitution method. By merely prohibiting a certain Zariski closed set of degenerations, we can make $T_{\mathcal{A}}$ noncompressible. Without that prohibition, $T_{\mathcal{A}}$ can indeed be compressed in general.

Remark 5.6.3.5. Using similar (but easier) methods, one can show that if $\mathcal{A}$ is simple of dimension a, then $\mathbf{R}\left(T_{\mathcal{A}}\right) \geq 2 \mathbf{a}-1$, see, e.g., [BCS97, Prop. 17.22]. More generally, the Alder-Strassen Theorem [AS81] states that if there are $m$ maximal two-sided ideals in $\mathcal{A}$, then $\mathbf{R}\left(T_{\mathcal{A}}\right) \geq 2 \mathbf{a}-m$

Theorem 5.6.3.6. [Bla01a] Let $\mathcal{A}_{\text {trunc }, d}:=\mathbb{C}\left[x_{1}, \ldots, x_{n}\right] /\left(S^{d} \mathbb{C}^{n}\right)=\bigoplus_{j=0}^{d-1} S^{j} \mathbb{C}^{n}$. Then

$$
\mathbf{R}\left(T_{\mathcal{A}_{\text {trunc }, d}}\right) \geq 3\binom{n+d}{d-1}-\binom{n+\left\lfloor\frac{d}{2}\right\rfloor}{\left\lfloor\frac{d}{2}\right\rfloor-1}-\binom{n+\left\lceil\frac{d}{2}\right\rceil}{\left\lceil\frac{d}{2}\right\rceil-1} .
$$

Proof. Apply Theorem 5.6.3.2. Here $\operatorname{Rad}\left(\mathcal{A}_{\text {trunc, }}\right)$ is a vector space complement to $\{\operatorname{Id}\}$ in $\mathcal{A}_{\text {trunc, },}$, so it has dimension $\binom{n+d}{d-1}-1$ and $\operatorname{Rad}\left(\mathcal{A}_{\text {trunc }, d}\right)^{k}=$ $\sum_{j=k}^{d-1} S^{j} \mathbb{C}^{n}$ which has dimension $\binom{n+d}{d-1}-\binom{n+k}{k-1}$.

In $\S 5.6 .5$ we will see that any algebra $\mathbb{C}\left[x_{1}, \ldots, x_{n}\right] / \mathcal{I}$ where $\mathcal{I}$ is an ideal generated by monomials, gives rise to a tensor of minimal border rank. Thus, as was observed by Bläser:
Corollary 5.6.3.7 (Bläser, personal communication). Let $d=d(n)<n$ be an integer valued function of $n$. Then

$$
\frac{\mathbf{R}\left(T_{\mathcal{A}_{\text {trunc }, d}}\right)}{\underline{\mathbf{R}}\left(T_{\mathcal{A}_{\text {trunc }, d}}\right)} \geq 3-o(n) .
$$

If $d=\left\lfloor\frac{n}{2}\right\rfloor$, then the error term is on the order of $1 / \operatorname{dim} \mathcal{A}_{\text {trunc, } d}$.
Theorem 5.6.3.8. [Zui15] $\mathbf{R}\left(T_{\text {WState }}^{\otimes n}\right)=3 \cdot 2^{n}-o\left(2^{n}\right)$.

Proof. We have $\mathcal{A}=\mathbb{C}\left[x_{1}, \ldots, x_{n}\right] /\left(x_{1}^{2}, \ldots, x_{n}^{2}\right)$, so the degree $n-s$ component of $\mathcal{A}$ is $\mathcal{A}_{s}=\operatorname{span} \bigcup_{S \subset[n]}\left\{x_{1} \cdots \hat{x}_{i_{1}} \cdots \hat{x}_{i_{s}} \cdots x_{n}\right\}=\operatorname{span} \bigcup_{S \subset[n]}\left\{\frac{x_{1} \cdots \cdots x_{n}}{x_{i_{1}} \cdots x_{i_{s}}}\right\}$. In particular $\operatorname{dim} \mathcal{A}_{s}=\binom{n}{s}$.

Note that $\operatorname{Rad}(\mathcal{A})^{m}=\oplus_{j \geq m} \mathcal{A}_{j}$. Recall that $\sum_{j=0}^{n}\binom{n}{j}=2^{n}$. Take $p=q$ in Theorem 5.6.3.2. We have

$$
\begin{aligned}
\mathbf{R}\left(T_{\mathcal{A}}\right) & \geq 2^{n}+2 \sum_{j=p}^{n}\binom{n}{j}-\sum_{k=2 p-1}^{n}\binom{n}{k} \\
& =3 \cdot 2^{n}-2 \sum_{j=0}^{p}\binom{n}{j}-\sum_{k=0}^{n-2 p+1}\binom{n}{k} .
\end{aligned}
$$

Write $p=\epsilon n$, for some $0<\epsilon<1$. Since $\sum_{j=0}^{\epsilon n}\binom{n}{j} \leq 2^{(-\epsilon \log (\epsilon)-(1-\epsilon) \log (1-\epsilon)) n}$ (see §7.5.1), taking, e.g., $\epsilon=\frac{1}{3}$ gives the result.

Corollary 5.6.3.9. $[$ Zui15 $] \frac{\mathbf{R}\left(T_{V \text { State }}^{\otimes n}\right)}{\underline{\underline{R}}\left(T_{\text {WState }}^{\otimes \infty}\right)} \geq 3-o(1)$, where the right hand side is viewed as a function of $n$.

More generally, Zuiddam shows, for $T_{\text {WState, } k}^{\otimes n} \in\left(\mathbb{C}^{n}\right)^{\otimes k}$ :
Theorem 5.6.3.10. [Zui15] $\mathbf{R}\left(T_{W \text { State }, k}^{\otimes n}\right)=k 2^{n}-o\left(2^{n}\right)$.
Regarding the maximum possible ratio for rank to border rank, there is the following theorem applicable even to $X$-rank and $X$-border rank:
Theorem 5.6.3.11. [BT15] Let $X \subset \mathbb{P} V$ be a complex projective variety not contained in a hyperplane. Let $\underline{\mathbf{R}}_{X, \max }$ denote the maximum $X$-border rank of a point in $\mathbb{P} V$ and $\mathbf{R}_{X, \text { max }}$ the maximum possible $X$-rank. Then $\mathbf{R}_{X, \max } \leq 2 \underline{\mathbf{R}}_{X, \max }$.

Proof. Let $U \subset \mathbb{P} V$ be a Zariski dense open subset of points of rank exactly $\mathbf{R}_{X, \max }$. Let $q \in \mathbb{P} V$ be any point and let $p$ be any point in $U$. The line $L$ through $q$ and $p$ intersects $U$ at another point $p$ (in fact, at infinitely many more points). Since $p$ and $p^{\prime}$ span $L, q$ is a linear combination of $p$ and $p^{\prime}$, thus $\mathbf{R}_{X}(q) \leq \mathbf{R}_{X}(p)+\mathbf{R}_{X}\left(p^{\prime}\right)$

Theorem 5.6.3.11 implies that the maximal possible rank of any tensor in $\mathbb{C}^{\mathbf{m}} \otimes \mathbb{C}^{\mathbf{m}} \otimes \mathbb{C}^{\mathbf{m}}$ is at most $2\left\lceil\frac{\mathbf{m}^{3}-1}{3 \mathbf{m}-2}\right\rceil$, so for any concise tensor the maximal rank to border rank ratio is bounded above by approximately $\frac{2 \mathrm{~m}}{3}$, which is likely far from sharp.

### 5.6.4. The border substitution method and tensor powers of $T_{c w, 2}$.

 Lemma 5.6.4.1. [BL16] For any tensor $T_{1} \in A_{1} \otimes B_{1} \otimes C_{1}$, and any $q \geq 2$, $\min _{\alpha \in\left(A \otimes A_{1}\right)^{*} \backslash\{0\}}\left(\left.\operatorname{rank}\left(T_{c w, q} \otimes T_{1}\right)\right|_{\alpha \otimes B^{*} \otimes C^{*}}\right) \geq 2 \min _{\alpha_{1} \in A_{1} \backslash\{0\}}\left(\operatorname{rank}\left(\left.T_{1}\right|_{\alpha_{1} \otimes B_{1}^{*} \otimes C_{1}^{*}}\right)\right)$.Proof. Write $\alpha=1 \otimes \alpha_{0}+\sum_{j=1}^{q} e_{j}^{*} \otimes \alpha_{j} \in\left(A \otimes A_{1}\right)^{*}$ for some $\alpha_{0}, \alpha_{j} \in A_{1}^{*}$. If all the $\alpha_{j}$ are zero for $1 \leq j \leq q$, then $T_{c w, q}\left(e_{0}^{*} \otimes \alpha_{0}\right)$ is the reordering and
grouping of

$$
\sum_{i=1}^{q}\left(e_{i} \otimes e_{i}\right) \otimes T_{1}\left(\alpha_{0}\right)
$$

which has rank (as a linear map) at least $q \cdot \operatorname{rank}\left(T_{1}\left(\alpha_{0}\right)\right)$. Otherwise without loss of generality, assume $\alpha_{1} \neq 0$. Note that $T_{c w, q}\left(e_{1}^{*} \otimes \alpha_{1}\right)$ is the reordering and grouping of

$$
e_{1} \otimes e_{0} \otimes T_{1}\left(\alpha_{1}\right)+e_{0} \otimes e_{1} \otimes T_{1}\left(\alpha_{1}\right)
$$

which has rank two, and is linearly independent of any of the other factors appearing in the image, so the rank is at least $2 \cdot \operatorname{rank}\left(T_{1}\left(\alpha_{0}\right)\right)$.
Theorem 5.6.4.2. [BL16] For all $q \geq 2$, consider $T_{c w, q}^{\otimes n} \in \mathbb{C}^{(q+1)^{n}} \otimes \mathbb{C}^{(q+1)^{n}} \otimes \mathbb{C}^{(q+1)^{n}}$. Then $\underline{\mathbf{R}}\left(T_{c w, q}^{\otimes n}\right) \geq(q+1)^{n}+2^{n}-1$.

Proof. Note that $T_{c u, q}^{\otimes n}=T_{c w, q} \otimes T_{c u, q}^{\otimes(n-1)}$. Apply the Lemma iteratively and use Corollary 5.4.1.4.

Remark 5.6.4.3. As was pointed out in [BCS97, Rem. 15.44] if the asymptotic rank (see Definition 3.4.6.1) of $T_{c w, 2}$ is the minimal 3, then the exponent of matrix multiplication is 2 . The bound in the theorem does not rule this out.
5.6.5. Smoothable ideals and tensors of minimal border rank. In §5.6.1 we saw that classifying $1_{A}$ and $1_{B}$ generic tensors of minimal border rank is equivalent to the potentially simpler problem of classifying algebras in the $G L_{\mathbf{m}}$-orbit closure of $M_{\langle 1\rangle}^{\oplus \mathbf{m}}$. We can translate this further when the algebras are of the form $\mathbb{C}\left[x_{1}, \ldots, x_{N}\right] / \mathcal{I}$ for some ideal $\mathcal{I}$. The question then becomes if $\mathcal{I}$ is a degeneration of an ideal whose zero set consists of $\mathbf{m}$ distinct points (counted with multiplicity).

The degenerations of ideals have been well-studied, and we are interested in the degeneration of the ideal of $\mathbf{m}$ distinct points to other ideals.

For example, the following algebras have the desired property and thus their structure tensors are of minimal border rank (see [CEVV09]):

- $\operatorname{dim}(\mathcal{A}) \leq 7$,
- $\mathcal{A}$ is generated by two elements,
- the radical of $\mathcal{A}$ satisfies $\operatorname{dim}\left(\operatorname{Rad}(\mathcal{A})^{2} / \operatorname{Rad}(\mathcal{A})^{3}\right)=1$,
- the radical of $\mathcal{A}$ satisfies $\operatorname{dim}\left(\operatorname{Rad}(\mathcal{A})^{2} / \operatorname{Rad}(\mathcal{A})^{3}\right)=2, \operatorname{dim} \operatorname{Rad}(\mathcal{A})^{3} \leq$ 2 and $\operatorname{Rad}(\mathcal{A})^{4}=0$.
An ideal $\mathcal{I}$ is a monomial ideal if it is generated by monomials (in some coordinate system). Choose an order on monomials such that if $|I|>|J|$, then $x^{I}<x^{J}$. Given $f \in \mathbb{C}\left[x_{1}, \ldots, x_{n}\right]$, define $i n(f)$ to be the lowest monomial term of $f$, the initial term of $f$. Given an ideal $\mathcal{I}$, define its initial
ideal (with respect to some chosen order) as $(\operatorname{in}(f) \mid f \in \mathcal{I})$. An ideal can be degenerated to its initial ideal.

Proposition 5.6.5.1. [CEVV09] Monomial ideals are smoothable, so if $\mathcal{I}$ is a monomial ideal then the structure tensor of $\mathbb{C}\left[x_{1}, \ldots, x_{n}\right] / \mathcal{I}$ is of minimal border rank.

Proof. Write $\mathcal{I}=\left(x^{I_{1}}, \ldots, x^{I_{s}}\right)$ for the ideal, where $I_{\alpha}=\left(i_{\alpha, 1}, \ldots, i_{\alpha,\left|I_{\alpha}\right|}\right)$, and let $\mathbf{m}=\operatorname{dim} \mathbb{C}\left[x_{1}, \ldots, x_{N}\right] / \mathcal{I}$. Take a sequence $a_{1}, a_{2}, \ldots$ of distinct elements of $\mathbb{C}$. Define

$$
f_{q}:=\Pi_{j=1}^{N}\left(x_{j}-a_{1}\right)\left(x_{j}-a_{2}\right) \cdots\left(x_{j}-a_{i_{s, q}}\right) .
$$

Note that $i n\left(f_{q}\right)=x^{I_{q}}$. Let $\mathcal{J}$ be the ideal generated by the $f_{q}$. Then $\operatorname{in}(\mathcal{J}) \supset\left(\operatorname{in}\left(f_{1}\right), \ldots, \operatorname{in}\left(f_{s}\right)\right)=\mathcal{I}$, so $\operatorname{dim} \mathbb{C}\left[x_{1}, \ldots, x_{N}\right] / \mathcal{J} \leq \mathbf{m}$. But now for any of the $x^{I_{q}} \in \mathcal{I}$, there each $f_{q}$ vanishes at $\left(a_{I_{q, 1}}, \ldots, a_{I_{q, N}}\right) \in \mathbb{C}^{N}$. Thus $\mathcal{J}$ must be the radical ideal vanishing at the $s$ points and have initial ideal $\mathcal{I}$, so $\mathcal{I}$ is smoothable.

## Chapter 6

## Valiant's hypothesis I: permanent $v$. determinant and the complexity of polynomials

Recall from the introduction that for a polynomial $P$, the determinantal complexity of $P$, denoted $\operatorname{dc}(P)$, is the smallest $n$ such that $P$ is an affine linear projection of the determinant, and Valiant's hypothesis 1.2.4.2 that dc $\left(\operatorname{perm}_{m}\right)$ grows faster than any polynomial in $m$. In this chapter I discuss the conjecture, progress towards it, and its Geometric Complexity Theory ( GCT ) variant.

I begin, in $\S 6.1$, with a discussion of circuits, context for Valiant's hypothesis, definitions of the complexity classes VP and VNP, and the strengthening of Valiant's hypothesis of [MS01] that is more natural for algebraic geometry and representation theory. In particular, I explain why it might be considered as an algebraic analog of the famous $\mathbf{P} \neq$ NP conjecture (although there are other conjectures in the Boolean world that are more closely related to it).

Our study of matrix multiplication indicates a strategy for Valiant's hypothesis: look for polynomials on the space of polynomials that vanish on the determinant and not on the permanent, and to look for such polynomials with the aid of geometry and representation theory. Here there is extra
geometry available: a polynomial $P \in S^{d} V$ defines a hypersurface

$$
\operatorname{Zeros}(P):=\left\{[\alpha] \in \mathbb{P} V^{*} \mid P(\alpha)=0\right\} \subset \mathbb{P} V^{*} .
$$

Hypersurfaces in projective space have been studied for hundreds of years and much is known about them.

In $\S 6.2$ I discuss the simplest polynomials on spaces of polynomials, the catalecticants that date back to Sylvester.

One approach to Valiant's hypothesis discussed at several points in this chapter is to look for pathologies of the hypersurface $\operatorname{Zeros}\left(\operatorname{det}_{n}\right)$ that persist under degeneration, and that are not shared by $\operatorname{Zeros}\left(\ell^{n-m} \operatorname{perm}_{m}\right)$. The simplest pathology of a hypersurface is its singular set. I discuss the singular loci of the permanent and determinant, and make general remarks on singularities in §6.3.

I then present the classical and recent lower bounds on dc $\left(\operatorname{perm}_{m}\right)$ of von zur Gathen and Alper-Bogart-Velasco in $\S 6.3 .4$. These lower bounds on $\mathrm{dc}\left(\operatorname{perm}_{m}\right)$ rely on a key regularity result observed by von zur Gathen. These results do not directly extend to the measure $\overline{d c}\left(\operatorname{perm}_{m}\right)$ defined in §6.1.6 because of the regularity result.

The best general lower bound on $\operatorname{dc}\left(\operatorname{perm}_{m}\right)$, namely $\operatorname{dc}\left(\operatorname{perm}_{m}\right) \geq \frac{m^{2}}{2}$, comes from local differential geometry: the study of Gauss maps. It is presented in $\S 6.4$. This bound extends to $\overline{d c}\left(\right.$ perm $\left._{m}\right)$ after some work. The extension is presented in $\S 6.5$. To better utilize geometry and representation theory, I describe the symmetries of the permanent and determinant in §6.6. Given $P \in S^{d} V$, let $G_{P}:=\{g \in G L(V) \mid g \cdot P=P\}$ denote the symmetry group of the polynomial $P$.

Since $\operatorname{det}(A X B)=\operatorname{det}(X)$ if $A, B$ are $n \times n$ matrices with determinant one, and $\operatorname{det}\left(X^{T}\right)=\operatorname{det}(X)$, writing $V=E \otimes F$ with $E, F=\mathbb{C}^{n}$, we have a map

$$
(S L(E) \times S L(F)) \rtimes \mathbb{Z}_{2} \rightarrow G_{\operatorname{det}_{n}}
$$

where the $\mathbb{Z}_{2}$ is transpose and $S L(E)$ is the group of linear maps with determinant equal to one.

Similarly, letting $T_{E}^{S L} \subset S L(E)$ denote the diagonal matrices, we have a map

$$
\left[\left(T_{E}^{S L} \rtimes \mathfrak{S}_{n}\right) \times\left(T_{F}^{S L} \rtimes \mathfrak{S}_{n}\right)\right] \rtimes \mathbb{Z}_{2} \rightarrow G_{\operatorname{perm}_{n}}
$$

In $\S 6.6$, I show that both maps are surjective.
Just as it is interesting and useful to study the difference between rank and border rank, it is worthwhile to study the difference between dc and $\overline{\mathrm{dc}}$, which I discuss in $\S 6.7$.

One situation where there is some understanding of the difference between $\overline{\mathrm{dc}}$ and dc is for cubic surfaces: a smooth cubic polynomial $P$ in three
variables satisfies $\operatorname{dc}(P)=3$, and thus every cubic polynomial $Q$ in three variables satisfies $\overline{\mathrm{dc}}(Q)=3$. I give an outline of the proof in $\S 6.8$. Finally, although it is not strictly related to complexity theory, I cannot resist a brief discussion of determinantal hypersurfaces - those degree $n$ polynomials $P$ with $\operatorname{dc}(P)=n$, which I also discuss in in $\S 6.8$.

In this chapter I emphasize material that is not widely available to computer scientists, and do not present proofs that already have excellent expositions in the literature, such as the completeness of the permanent for VNP.

This chapter may be read mostly independently of chapters 2-5.

### 6.1. Circuits and definitions of VP and VNP

In this section I give definitions of VP and VNP via arithmetic circuits and show $\left(\operatorname{det}_{n}\right) \in \mathbf{V P}$. I also discuss why Valiant's hypothesis is a cousin of $\mathbf{P} \neq \mathbf{N P}$, namely I show that the permanent can compute the number of perfect matchings of a bipartite graph, something considered difficult, while the determinant can be computed by a polynomial size circuit.
6.1.1. The permanent can do things considered difficult. A standard problem in graph theory, for which the only known algorithms are exponential in the size of the graph, is to count the number of perfect matchings of a bipartite graph, that is, a graph with two sets of vertices and edges only joining vertices from one set to the other.


Figure 6.1.1. A bipartite graph, Vertex sets are $\{A, B, C\}$ and $\{\alpha, \beta, \gamma\}$.

A perfect matching is a subset of the edges such that each vertex shares an edge from the subset with exactly one other vertex.

To a bipartite graph one associates an incidence matrix $x_{j}^{i}$, where $x_{j}^{i}=1$ if an edge joins the vertex $i$ above to the vertex $j$ below and is zero otherwise. The graph above has incidence matrix

$$
\left(\begin{array}{lll}
1 & 1 & 0 \\
0 & 1 & 1 \\
0 & 1 & 1
\end{array}\right)
$$



Figure 6.1.2. Two perfect matchings of the graph from Figure 6.1.1.

A perfect matching corresponds to a matrix constructed from the incidence matrix by setting some of the entries to zero so that the resulting matrix has exactly one 1 in each row and column, i.e., is a matrix obtained by applying a permutation to the columns of the identity matrix.
Exercise 6.1.1.1: (1) Show that if $x$ is the incidence matrix of a bipartite graph, then $\operatorname{perm}_{n}(x)$ indeed equals the number of perfect matchings.

For example, $\operatorname{perm}_{3}\left(\begin{array}{lll}1 & 1 & 0 \\ 0 & 1 & 1 \\ 0 & 1 & 1\end{array}\right)=2$.
Thus a classical problem: determine the complexity of counting the number of perfect matchings of a bipartite graph (which is complete for the complexity class $\sharp \mathbf{P}$, see [BCS97, p. 574]), can be studied via algebra determine the complexity of evaluating the permanent.

### 6.1.2. Circuits.

Definition 6.1.2.1. An arithmetic circuit $\mathcal{C}$ is a finite, directed, acyclic graph with vertices of in-degree 0 or 2 and exactly one vertex of out-degree 0 . The vertices of in-degree 0 are labeled by elements of $\mathbb{C} \cup\left\{x_{1}, \ldots, x_{n}\right\}$, and called inputs. Those of in-degree 2 are labeled with + or $*$ and are called gates. If the out-degree of $v$ is 0 , then $v$ is called an output gate. The size of $\mathcal{C}$ is the number of edges.

To each vertex $v$ of a $\operatorname{circuit} \mathcal{C}$, associate the polynomial that is computed at $v$, which will be denoted $\mathcal{C}_{v}$. In particular the polynomial associated with the output gate is called the polynomial computed by $\mathcal{C}$.

At first glance, circuits do not look geometrical, as they depend on a choice of coordinates. While computer scientists always view polynomials as being given in some coordinate expression, in geometry one is interested in properties of objects that are independent of coordinates. These perspectives are compatible because with circuits one is not concerned with the precise size of a circuit, but its size up to, e.g., a polynomial factor. Reducing the size at worst by a polynomial factor, we can think of the inputs to our circuits as arbitrary affine linear or linear functions on a vector space.


Figure 6.1.3. Circuit for $(x+y)^{3}$

### 6.1.3. Arithmetic circuits and complexity classes.

Definition 6.1.3.1. Let $d(n), N(n)$ be polynomials and let $f_{n} \in \mathbb{C}\left[x_{1}, \ldots, x_{N(n)}\right]_{\leq d(n)}$ be a sequence of polynomials. We say $\left(f_{n}\right) \in \mathbf{V P}$ if there exists a sequence of circuits $\mathcal{C}_{n}$ of size polynomial in $n$ computing $f_{n}$.

Often the phrase "there exists a sequence of circuits $\mathcal{C}_{n}$ of size polynomial in $n$ computing $f_{n}$ " is abbreviated "there exists a polynomial sized circuit computing $\left(f_{n}\right)$ ".

The class VNP, which consists of sequences of polynomials whose coefficients are "easily" described, has a more complicated definition:

Definition 6.1.3.2. A sequence $\left(f_{n}\right)$ is in VNP if there exists a polynomial $p$ and a sequence $\left(g_{n}\right) \in \mathbf{V P}$ such that

$$
f_{n}(x)=\sum_{\epsilon \in\{0,1\}^{p(n)}} g_{n}(x, \epsilon) .
$$

One may think of the class VP as a bundle over VNP where elements of VP are thought of as sequences of maps, say $g_{n}: \mathbb{C}^{N(n)} \rightarrow \mathbb{C}$, and elements of VNP are projections of these maps by eliminating some of the variables by averaging or "integration over the fiber". In algebraic geometry, it is well known that projections of varieties can be far more complicated than the original varieties. See [Bas14] for more on this perspective.

The class VNP is sometimes described as the polynomial sequences that can be written down "explicitly". Mathematicians should take note that the computer science definition of explicit is different from what a mathematician might use. For example, as pointed out in $[$ FS13a], roots of unity are not explicit because using them computationally typically requires
expanding them as a decimal with exponential precision, which is inefficient. On the other hand, the lexicographically first function $f:\{0,1\}^{\lfloor\log \log n\rfloor} \rightarrow$ $\{0,1\}$ with the maximum possible circuit complexity among all functions on $\lfloor\log \log n\rfloor$ bits is explicit because, while seemingly unstructured, this function can be writtend down efficiently via brute-force. See [FS13a] for the definition.

Definition 6.1.3.3. One says that a sequence $\left(g_{m}\left(y_{1}, \ldots, y_{M(m)}\right)\right)$ can be polynomially reduced to $\left(f_{n}\left(x_{1}, \ldots, x_{N(n)}\right)\right)$ if there exists a polynomial $n(m)$ and affine linear functions $X_{1}\left(y_{1}, \ldots, y_{M}\right), \ldots, X_{N}\left(y_{1}, \ldots, y_{M}\right)$ such that $g_{m}\left(y_{1}, \ldots, y_{M(m)}\right)=f_{n}\left(X_{1}(y), \ldots, X_{N(n)}(y)\right)$. A sequence $\left(p_{n}\right)$ is hard for a complexity class $\mathbf{C}$ if $\left(p_{n}\right)$ can be reduced to every $\left(f_{m}\right) \in \mathbf{C}$, and it is complete for $\mathbf{C}$ if furthermore $\left(p_{n}\right) \in \mathbf{C}$.

Exercise 6.1.3.4: (1) Show that every polynomial of degree $d$ can be reduced to $x^{d}$.

Theorem 6.1.3.5. [Valiant] [Val79] $\left(\right.$ perm $\left._{m}\right)$ is complete for VNP.
There are many excellent expositions of the proof, see, e.g. [BCS97] or [Gat87].

Thus Conjecture 1.2.1.1 is equivalent to:
Conjecture 6.1.3.6. [Valiant][Val79] There does not exist a polynomial size circuit computing the permanent.

Now for the determinant:
Proposition 6.1.3.7. $\left(\operatorname{det}_{n}\right) \in \mathbf{V P}$.
Remark 6.1.3.8. $\operatorname{det}_{n}$ would be VP complete if $\mathrm{dc}\left(p_{m}\right)$ grew no faster than a polynomial for all sequences $\left(p_{m}\right) \in \mathbf{V P}$.

One can compute the determinant quickly via Gaussian elimination: one uses the group to put a matrix in a form where the determinant is almost effortless to compute (the determinant of an upper triangular matrix is just the product of its diagonal entries). However this algorithm as presented is not a circuit (there are divisions and one needs to check if pivots are zero). After a short detour on symmetric polynomials, I prove Proposition 6.1.3.7 in §6.1.5.
6.1.4. Symmetric polynomials. An ubiquitous class of polynomials are the symmetric polynomials: let $\mathfrak{S}_{N}$ act on $\mathbb{C}^{N}$ by permuting basis elements, which induces an action on the polynomial ring $\mathbb{C}\left[x_{1}, \ldots, x_{N}\right]$. Let $\mathbb{C}\left[x_{1}, \ldots, x_{N}\right]^{\mathfrak{G}_{N}}$ denote the subspace of polynomials invariant under this action. What follows are standard facts and definitions about symmetric functions. For proofs, see, e.g., [Mac95, §I.2].

The elementary symmetric functions (or elementary symmetric polynomials) are

$$
\begin{equation*}
e_{n}=e_{n, N}=e_{n}\left(x_{1}, \ldots, x_{N}\right):=\sum_{J \subset[N]| | J \mid=n} x_{j_{1}} \cdots x_{j_{n}} . \tag{6.1.1}
\end{equation*}
$$

If the number of variables is understood, I write $e_{n}$ for $e_{n, N}$. They generate the ring of symmetric polynomials. They have the generating function

$$
\begin{equation*}
E_{N}(t):=\sum_{k \geq 0} e_{k}\left(x_{1}, \ldots, x_{N}\right) t^{k}=\prod_{i=1}^{N}\left(1+x_{i} t\right) \tag{6.1.2}
\end{equation*}
$$

Exercise 6.1.4.1: (1) Verify the coefficient of $t^{n}$ in $E_{N}(t)$ is $e_{n, N}$.
The power sum symmetric functions are

$$
\begin{equation*}
p_{n}=p_{n, N}=p_{n, N}\left(x_{1}, \ldots, x_{N}\right)=x_{1}^{n}+\cdots+x_{N}^{n} . \tag{6.1.3}
\end{equation*}
$$

They also generate the ring of symmetric polynomials. They have the generating function

$$
\begin{equation*}
P_{N}(t)=\sum_{k \geq 1} p_{k} t^{k-1}=\frac{d}{d t} \ln \left[\prod_{j=1}^{N}\left(1-x_{j} t\right)^{-1}\right] . \tag{6.1.4}
\end{equation*}
$$

Exercise 6.1.4.2: (2) Verify that the coefficient of $t^{n}$ in $P_{N}(t)$ is indeed $p_{n, N}$. ©
Exercise 6.1.4.3: (2) Show that

$$
\begin{equation*}
P_{N}(-t)=-\frac{E_{N}^{\prime}(t)}{E_{N}(t)} \tag{6.1.5}
\end{equation*}
$$

Exercise 6.1.4.3, together with a little more work (see, e.g. [Mac95, p. 28]) shows that

$$
p_{n}=\operatorname{det}_{n}\left(\begin{array}{ccccc}
e_{1} & 1 & 0 & \cdots & 0  \tag{6.1.6}\\
2 e_{2} & e_{1} & 1 & \cdots & 0 \\
\vdots & \vdots & \ddots & \ddots & \vdots \\
\vdots & \vdots & \vdots & & 1 \\
n e_{n} & e_{n-1} & e_{n-2} & \cdots & e_{1}
\end{array}\right) .
$$

Similarly

$$
e_{n}=\frac{1}{n!} \operatorname{det}_{n}\left(\begin{array}{ccccc}
p_{1} & 1 & 0 & \cdots & 0  \tag{6.1.7}\\
p_{2} & p_{1} & 2 & \cdots & 0 \\
\vdots & \vdots & \vdots & \ddots & \vdots \\
p_{n-1} & p_{n-2} & & \cdots & n-1 \\
p_{n} & p_{n-1} & & \cdots & p_{1}
\end{array}\right) .
$$

6.1.5. Proof of Proposition 6.1.3.7. Here is a construction of a small circuit for the determinant that essentially appeared in [Csa76]:

The determinant of a linear map $f: V \rightarrow V$ is the product of its eigenvalues $\lambda_{1}, \ldots, \lambda_{\mathbf{v}}$, i.e., $e_{\mathbf{v}}(\lambda)=\lambda_{1} \cdots \lambda_{\mathbf{v}}$.

On the other hand, trace $(f)$ is the sum of the eigenvalues of $f$, and more generally, letting $f^{k}$ denote the composition of $f$ with itself $k$ times,

$$
\operatorname{trace}\left(f^{k}\right)=p_{k}(\lambda)=\lambda_{1}^{k}+\cdots+\lambda_{\mathbf{v}}^{k} .
$$

The quantities trace $\left(f^{k}\right)$ can be computed with small circuits.
Exercise 6.1.5.1: (2) Write down a circuit for the polynomial $A \mapsto \operatorname{trace}\left(A^{2}\right)$ when $A$ is an $n \times n$ matrix with variable entries.

Thus we can compute $\operatorname{det}_{n}$ via small circuits and (6.1.7). While (6.1.7) is still a determinant, it is almost lower triangular and its naïve computation, e.g., with Laplace expansion, can be done with an $O\left(n^{3}\right)$-size circuit and the full algorithm for computing $\operatorname{det}_{n}$ can be executed with an $O\left(n^{4}\right)$ size circuit.

Remark 6.1.5.2. A more restrictive class of circuits are formulas which are circuits that are trees. Let $\mathbf{V P}_{e}$ denote the sequences of polynomials that admit a polynomial size formula. The circuit in the proof above is not a formula because results from computations are used more than once. It is known that the determinant admits a quasi-polynomial size formula, that is, a formula of size $n^{O(\log n)}$, and it is complete for the complexity class VQP $=$ $\mathbf{V P}{ }_{s}$ consisting of sequences of polynomials admitting a quasi-polynomial size formula see, e.g., [BCS97, §21.5] (or equivalently, a polynomial sized "skew" circuit, see [Tod92]). It is not known whether or not the determinant is complete for VP.
6.1.6. The Geometric Complexity Theory (GCT) variant of Valiant's hypothesis. Recall that when we used polynomials in the study of matrix multiplication, we were proving lower bounds on tensor border rank rather than tensor rank. In the case of matrix multiplication, at least as far as the exponent is concerned, this changed nothing. In the case of determinant versus permanent, it is not known if using polynomial methods leads to a stronger separation of complexity classes. In any case, it will be best to clarify the two different types of lower bounds.

I recall from $\S 1.2$ that using padded polynomials, one can rephrase Valiant's hypothesis as:
Conjecture 6.1.6.1. [Rephrasing of Valiant's hypothesis] Let $\ell$ be a linear coordinate on $\mathbb{C}^{1}$ and consider any linear inclusion $\mathbb{C}^{1} \oplus \mathbb{C}^{m^{2}} \rightarrow \mathbb{C}^{n^{2}}$, so in particular $\ell^{n-m} \operatorname{perm}_{m} \in S^{n} \mathbb{C}^{n^{2}}$. Let $n(m)$ be a polynomial. Then for all
sufficiently large m,

$$
\left[\ell^{n-m} \operatorname{perm}_{m}\right] \notin \operatorname{End}\left(\mathbb{C}^{n^{2}}\right) \cdot\left[\operatorname{det}_{n(m)}\right]
$$

Recall that the formulations are equivalent because if $\operatorname{perm}\left(y_{j}^{i}\right)=\operatorname{det}_{n}(\Lambda+$ $\left.\sum_{i, j} A_{i j} y_{i, j}\right)$, then $\ell^{n-m} \operatorname{perm}_{m}\left(y_{i, j}\right)=\operatorname{det}_{n}\left(\ell \Lambda+\sum_{i, j} A_{i j} y_{i, j}\right)$. Such an expression is equivalent to setting each entry of the $n \times n$ matrix to a linear combination of the variables $\ell, y_{i, j}$, which is precisely what the elements of rank $m^{2}+1$ in $\operatorname{End}\left(\mathbb{C}^{n^{2}}\right)$ can accomplish. Moreover $\ell^{n-m} \operatorname{perm}_{m}=X \cdot \operatorname{det}_{n(m)}$ for some $X \in \operatorname{End}\left(\mathbb{C}^{n^{2}}\right)$ implies $X$ has rank $m^{2}+1$.

Recall the following conjecture, made to facilitate the use of tools from algebraic geometry and representation theory to separate complexity classes:
Conjecture 6.1.6.2. [MS01] Let $\ell$ be a linear coordinate on $\mathbb{C}^{1}$ and consider any linear inclusion $\mathbb{C}^{1} \oplus \mathbb{C}^{m^{2}} \rightarrow \mathbb{C}^{n^{2}}$, so in particular $\ell^{n-m} \operatorname{perm}_{m} \in$ $S^{n} \mathbb{C}^{n^{2}}$. Let $n(m)$ be a polynomial. Then for all sufficiently large $m$,

$$
\left[\ell^{n-m} \operatorname{perm}_{m}\right] \notin \overline{G L_{n^{2}} \cdot\left[\operatorname{det}_{n(m)}\right]}
$$

Note that $\overline{G L_{n^{2}} \cdot\left[\operatorname{det}_{n}\right]}=\overline{\operatorname{End}\left(\mathbb{C}^{n^{2}}\right) \cdot\left[\operatorname{det}_{n}\right]} . \operatorname{In} \S 6.7 .2$ I show $\overline{G L_{n^{2}} \cdot\left[\operatorname{det}_{n}\right]} \supsetneq$ $\operatorname{End}\left(\mathbb{C}^{n^{2}}\right) \cdot\left[\operatorname{det}_{n}\right]$, so Conjecture 6.1.6.2 is a strengthening of Conjecture 6.1.6.1. It will be useful to rephrase the conjecture slightly, to highlight that it is a question about determining whether one orbit closure is contained in another. Let

$$
\mathcal{D e t}_{n}:=\overline{G L_{n^{2}} \cdot\left[\operatorname{det}_{n}\right]}
$$

and let

$$
\mathcal{P e r m} n_{n}^{m}:=\overline{G L_{n^{2}} \cdot\left[\ell^{n-m} \operatorname{perm}_{m}\right]}
$$

Conjecture 6.1.6.3. [MS01] Let $n(m)$ be a polynomial. Then for all sufficiently large $m$,

$$
\mathcal{P e r m}_{n(m)}^{m} \not \subset \mathcal{D e t}_{n(m)}
$$

The equivalence of Conjectures 6.1.6.3 and 6.1.6.2 follows as $\ell^{n-m}$ perm $_{m} \notin$ $\mathcal{D} e t_{n}$ implies $G L_{n^{2}} \cdot \ell^{n-m} \operatorname{perm}_{m} \not \subset \mathcal{D} e t_{n}$, and since $\mathcal{D e t}_{n}$ is closed and both sides are irreducible, there is no harm in taking closure on the left hand side, as you showed in Exercise 3.3.1.1.

Both $\mathcal{P e r m} n_{n}^{m}$ and $\mathcal{D e t} t_{n}$ are invariant under $G L_{n^{2}}$ so their ideals are $G L_{n^{2}}$-modules. To separate them, one may look for a $G L_{n^{2}}$-module $M$ such that $M \subset I\left[\mathcal{D e t}_{n}\right]$ and $M \not \subset I\left[\mathcal{P e r m} m_{n}^{m}\right]$.

In $\S 8.8$ I explain the original program to solve this conjecture. Although that program cannot work as stated, I believe that the re-focusing of a problem of separating complexity classes to questions in algebraic geometry and representation theory could lead to viable paths to resolving Valiant's hypothesis.

### 6.2. Flattenings: our first polynomials on the space of polynomials

In this section I discuss the most classical polynomials on the space of polynomials, which were first introduced by Sylvester in 1852 and called catalecticants by him. They are also called flattenings and in the computer science literature the polynomials induced by the method of partial derivatives.
6.2.1. Three perspectives on $S^{d} \mathbb{C}^{M}$. I review our perspectives on $S^{d} \mathbb{C}^{M}$ from §2.3.2. We have seen $S^{d} \mathbb{C}^{M}$ is the space of symmetric tensors in $\left(\mathbb{C}^{M}\right)^{\otimes d}$. Given a symmetric tensor $T \in S^{d} \mathbb{C}^{M}$, we may form a polynomial $P_{T}$ on $\mathbb{C}^{M *}$ by, for $v \in \mathbb{C}^{M *}, P_{T}(v):=T(v, \ldots, v)$. I use this identification repeatedly without further mention.

One can also recover $T$ from $P_{T}$ via polarization. Then (up to universal constants) $T\left(v_{i_{1}}, \ldots, v_{i_{M}}\right)$ where $1 \leq i_{1} \leq \cdots \leq i_{M}$ is the coefficient of $t_{i_{1}} \cdots t_{i_{M}}$ in $P_{T}\left(t_{1} v_{1}+\cdots+t_{M} v_{M}\right)$. See [Lan12, Chap. 2] for details.

As was mentioned in Exercise 2.3.2.4, we may also think of $S^{d} \mathbb{C}^{M}$ as the space of homogeneous differential operators of order $d$ on $\operatorname{Sym}\left(\mathbb{C}^{M *}\right):=$ $\oplus_{j=0}^{\infty} S^{j} \mathbb{C}^{M *}$.

Thus we may view an element of $S^{d} \mathbb{C}^{M}$ as a homogeneous polynomial of degree $d$ on $\mathbb{C}^{M *}$, a symmetric tensor, and as a homogeneous differential operator of order $d$ on the space of polynomials $\operatorname{Sym}\left(\mathbb{C}^{M *}\right)$.
6.2.2. Catalecticants, a.k.a. the method of partial derivatives. Now would be a good time to read $\S 3.1$ if you have not already done so. I review a few essential points from it.

The simplest polynomials in $S^{n} \mathbb{C}^{N}$ are just the $n$-th powers of linear forms. Their zero set is a hyperplane (counted with multiplicity $n$ ). Let $P \in S^{n} \mathbb{C}^{N}$. How can one test if $P$ is an $n$-th power of a linear form, $P=\ell^{n}$ for some $\ell \in \mathbb{C}^{N}$ ?
Exercise 6.2.2.1: (1!) Show that $P=\ell^{n}$ for some $\ell \in \mathbb{C}^{N}$ if and only if $\operatorname{dim}\left\langle\frac{\partial P}{\partial x^{1}}, \ldots, \frac{\partial P}{\partial x^{N}}\right\rangle=1$, where $x^{1}, \ldots, x^{N}$ are coordinates on $\mathbb{C}^{N}$.

Exercise 6.2 .2 .1 is indeed a polynomial test: The dual space $\mathbb{C}^{N *}$ may be considered as the space of first order homogeneous differential operators on $S^{n} \mathbb{C}^{N}$, and the test is that the $2 \times 2$ minors of the map $P_{1, n-1}: \mathbb{C}^{N *} \rightarrow$ $S^{n-1} \mathbb{C}^{N}$, given by $\frac{\partial}{\partial x^{j}} \mapsto \frac{\partial P}{\partial x^{j}}$ are zero.

Exercise 6.2.2.1 may be phrased without reference to coordinates: recall the inclusion $S^{n} V \subset V \otimes S^{n-1} V=\operatorname{Hom}\left(V^{*}, S^{n-1} V\right)$. For $P \in S^{n} V$, write $P_{1, n-1} \in \operatorname{Hom}\left(V^{*}, S^{n-1} V\right)$.

Definition 6.2.2.2. I will say $P$ is concise if $P_{1, n-1}$ is injective.

In other words, $P$ is concise if every expression of $P$ in coordinates uses all the variables.

Exercise 6.2.2.1 may be rephrased as: $P$ is an $n$-th power of a linear form if and only if $\operatorname{rank}\left(P_{1, n-1}\right)=1$.

Recall that the $n$-th Veronese variety is

$$
v_{n}(\mathbb{P} V):=\left\{[P] \in \mathbb{P} S^{n} V \mid P=\ell^{n} \text { for some } \ell \in V\right\} \subset \mathbb{P}\left(S^{n} V\right)
$$

Exercise 6.2 .2 .1 shows that the Veronese variety is indeed an algebraic variety. It is homogenous, i.e., a single $G L(V)$-orbit.

More generally define the subspace variety

$$
S u b_{k}\left(S^{n} V\right):=\mathbb{P}\left\{P \in S^{n} V \mid \operatorname{rank}\left(P_{1, n-1}\right) \leq k\right\}
$$

Note that $[P] \in S u b_{k}\left(S^{n} V\right)$ if and only if there exists a coordinate system where $P$ can be expressed using only $k$ of the $\operatorname{dim} V$ variables. The subspace variety $S u b_{k}\left(S^{n} V\right) \subset \mathbb{P} S^{n} V$ has the geometric interpretation as the polynomials whose zero sets in projective space are cones with a $\mathbf{v}-k$ dimensional vertex. (In affine space the zero set may look like a cylinder, such as the surface $x^{2}+y^{2}=1$ in $\mathbb{R}^{3}$.) Consider the hypersurface $X_{P} \subset \mathbb{P}^{k-1}$ cut out by restricting $P$ to a subspace $L$ where $\left(\left.P\right|_{L}\right)_{1, n-1}$ is injective. Then points of $\operatorname{Zeros}(P) \subset \mathbb{P} V^{*}$ are of the form $[x+y]$ where $x \in \hat{X}_{P}$ and $y \in \mathbb{P}^{\mathbf{v}-k-1}=\mathbb{P} \operatorname{ker}\left(P_{1, n-1}\right)$. See $\S 6.4 .2$ for more details.

The symmetric rank of $P \in S^{n} V^{*}, \mathbf{R}_{v_{n}(\mathbb{P} V)}(P)=\mathbf{R}_{S}(P)$, is the smallest $r$ such that $P=\ell_{1}^{n}+\cdots+\ell_{r}^{n}$ for $\ell_{j} \in V$. The symmetric border rank of $P$, $\underline{\mathbf{R}}_{v_{n}(\mathbb{P} V)}(P)=\underline{\mathbf{R}}_{S}(P)$, is is the smallest $r$ such that $[P] \in \sigma_{r}\left(v_{n}(\mathbb{P} V)\right)$, the $r$-th secant variety of the Veronese variety (see $\S 4.7 .1$ ). Symmetric rank will appear naturally in the study of Valiant's hypothesis and its variants. In the language of $\S 7.1, \mathbf{R}_{S}(P)$ is essentially the size of the smallest homogeneous $\Sigma \Lambda \Sigma$-circuit computing $P$.

How would one test if $P$ is the sum of two $n$-th powers, $P=\ell_{1}^{n}+\ell_{2}^{n}$ for some $\ell_{1}, \ell_{2} \in \mathbb{C}^{N}$ ?
Exercise 6.2.2.3: (1) Show that $P=\ell_{1}^{n}+\ell_{2}^{n}$ for some $\ell_{j} \in \mathbb{C}^{N}$ implies $\operatorname{dim} \operatorname{span}\left\{\frac{\partial P}{\partial x^{1}}, \ldots, \left.\frac{\partial P}{\partial x^{N}} \right\rvert\, 1 \leq i, j \leq N\right\} \leq 2$.

Exercise 6.2.2.4: (2) Show that any polynomial vanishing on all polynomials of the form $P=\ell_{1}^{n}+\ell_{2}^{n}$ for some $\ell_{j} \in \mathbb{C}^{N}$ also vanishes on $x^{n-1} y$. ©

Exercise 6.2.2.4 reminds us that $\sigma_{2}\left(v_{n}(\mathbb{P} V)\right)$ also includes points on tangent lines.

The condition in Exercise 6.2.2.3 is not sufficient to determine membership in $\sigma_{2}\left(v_{n}(\mathbb{P} V)\right)$, in other words, $\sigma_{2}\left(v_{n}(\mathbb{P} V)\right) \subsetneq S u b_{2}\left(S^{n} V\right)$ : Consider
$P=\ell_{1}^{n-2} \ell_{2}^{2}$. It has $\operatorname{rank}\left(P_{1, n-1}\right)=2$ but $P \notin \sigma_{2}\left(v_{n}(\mathbb{P} V)\right)$ as can be seen by the following exercises:
Exercise 6.2.2.5: (1) Show that $P=\ell_{1}^{n}+\ell_{2}^{n}$ for some $\ell_{j} \in \mathbb{C}^{N}$ implies $\operatorname{dim} \operatorname{span}\left\{\frac{\partial^{2} P}{\partial x^{i} \partial x^{j}}\right\} \leq 2$.
Exercise 6.2.2.6: (1) Show that $P=\ell_{1}^{n-2} \ell_{2}^{2}$ for some distinct $\ell_{j} \in \mathbb{C}^{N}$ implies dim span $\left\{\frac{\partial^{2} P}{\partial x^{i} \partial x^{j}}\right\}>2$.

Let $P_{2, n-2}: S^{2} \mathbb{C}^{N *} \rightarrow S^{n-2} \mathbb{C}^{N}$ denote the map with image $\left\langle\cup_{i, j} \frac{\partial^{2} P}{\partial x^{i} \partial x^{j}}\right\rangle$. Vanishing of the size three minors of $P_{1, n-1}$ and $P_{2, n-2}$ are necessary and sufficient conditions for $P \in \sigma_{2}\left(v_{n}(\mathbb{P} V)\right)$, as was shown in 1886 by Gundelfinger [Gun].

More generally, one can consider the polynomials given by the minors of the maps $S^{k} \mathbb{C}^{N *} \rightarrow S^{n-k} \mathbb{C}^{N}$, given by $D \mapsto D(P)$. Write these maps as $P_{k, n-k}: S^{k} V^{*} \rightarrow S^{n-k} V$. These equations date back to Sylvester [Syl52] and are called the method of partial derivatives in the complexity literature, e.g. [CKW10]. The ranks of these maps gives a complexity measure on polynomials.

Let's give a name to the varieties defined by these polynomials: define Flat ${ }_{r}^{k, d-k}\left(S^{d} V\right):=\left\{P \in S^{d} V \mid \operatorname{rank}\left(P_{k, d-k}\right) \leq r\right\}$.
Exercise 6.2.2.7: (1!) What does the method of partial derivatives tell us about the complexity of $x_{1} \cdots x_{n}$, $\operatorname{det}_{n}$ and perm$n$, e.g., taking $k=\left\lfloor\frac{n}{2}\right\rfloor$ ? ©

Exercise 6.2.2.7 provides an exponential lower bound for the permanent in the complexity measure of symmetric border $\operatorname{rank} \underline{\mathbf{R}}_{S}$, but we obtain the same lower bound for the determinant. Thus this measure will not be useful for separating the permanent from the determinant. It still gives interesting information about other polynomials such as symmetric functions, which we will examine.

The variety of homogeneous polynomials of degree $n$ that are products of linear forms also plays a role in complexity theory. Recall the Chow variety of polynomials that decompose into a product of linear forms from §3.1.2:

$$
C h_{n}(V):=\mathbb{P}\left\{P \in S^{n} V \mid P=\ell_{1} \cdots \ell_{n} \text { for } \ell_{j} \in V\right\} .
$$

One can define a complexity measure for writing a polynomial as a sum of products of linear forms. The "Zariski closed" version of this condition is membership in $\sigma_{r}\left(C h_{n}(V)\right)$. In the language of circuits, $\mathbf{R}_{C h_{n}(V)}(P)$ is (essentially) the size of the smallest homogeneous $\Sigma \Pi \Sigma$ circuit computing a polynomial $P$. I discuss this in $\S 7.5$.

Exercise 6.2.2.7 gives a necessary test for a polynomial $P \in S^{n} \mathbb{C}^{N}$ to be a product of $n$ linear forms, namely $\operatorname{rank}\left(P_{\left\lfloor\frac{n}{2}\right\rfloor,\left\lceil\frac{n}{2}\right\rceil}\right) \leq\binom{ n}{\left\lfloor\frac{n}{2}\right\rfloor}$. A question to
think about: how would one develop a necessary and sufficient condition to show a polynomial $P \in S^{n} \mathbb{C}^{N}$ is a product of $n$ linear forms? See $\S 9.6$ for an answer.

Unfortunately we have very few techniques for finding good spaces of polynomials on polynomials. One such that generalizes flattenings, called Young flattenings is discussed in §8.2.

A natural question is whether or not all flattenings are non-trivial. I address this in $\S 6.2 .4$ below after defining conormal spaces, which will be needed for the proof.
6.2.3. Conormal spaces. Recall the definition of the tangent space to a point on a variety $X \subset \mathbb{P V}$ or $X \subset V, \hat{T}_{x} X \subset V$, from $\S 3.1 .3$. The conormal space $N_{x}^{*} X \subset V^{*}$ is simply defined to be the annihilator of the tangent space: $N_{x}^{*} X=\left(\hat{T}_{x} X\right)^{\perp}$.
Exercise 6.2.3.1: (2!) Show that in $\hat{\sigma}_{r}^{0}\left(\operatorname{Seg}\left(\mathbb{P}^{u-1} \times \mathbb{P}^{v-1}\right)\right)$, the space of $u \times v$ matrices of rank $r$,

$$
\hat{T}_{M} \sigma_{r}^{0}\left(S e g\left(\mathbb{P}^{u-1} \times \mathbb{P}^{v-1}\right)\right)=\left\{X \in M a t_{u \times v} \mid X \operatorname{ker}(M) \subset \operatorname{Image}(M)\right\}
$$

Give a description of $N_{M}^{*} \sigma_{r}^{0}\left(\operatorname{Seg}\left(\mathbb{P}^{u-1} \times \mathbb{P}^{v-1}\right)\right)$.
6.2.4. All flattenings give non-trivial equations. The polynomials obtained from the maximal minors of $P_{i, d-i}$ give nontrivial equations. In other words, let $r_{0}=r_{0}(i, d, \mathbf{v})=\binom{\mathbf{v}+i-1}{i}$. Then I claim that for $i \leq d-i$, Flat $t_{r_{0}-1}^{i, d-i}\left(S^{d} V\right)$ is a proper subvariety of $\mathbb{P} S^{d} V$.

The most natural way to prove the claim would be to exhibit an explict sequences of polynomials with maximal flatting rank. At this writing, $I$ do not know of any such explicit sequence. I give an indirect proof of the claim below.

Problem 6.2.4.1. Find an explicit sequence of polynomials $P_{d, n} \in S^{d} \mathbb{C}^{n}$ with maximal flattening rank. Can one find such an explicit sequence that lies in $\mathbf{V P}, \mathbf{V P}_{s}$ or even $\mathbf{V P}_{e}$ ?

Exercise 6.2.4.2: (1) Show that if $P_{\left\lfloor\frac{d}{2}\right\rfloor,\left\lceil\frac{d}{2}\right\rceil}$ is of maximal rank, then all $P_{k, d-k}$ are of maximal rank.
Theorem 6.2.4.3. [Gre78, IE78] For a general polynomial $P \in S^{d} V$, all the maps $P_{k, d-k}: S^{k} V^{*} \rightarrow S^{d-k} V$ are of maximal rank.

Proof. (Adapted from [IK99].) By Exercise 6.2.4.2 it is sufficient to consider the case $k=\left\lfloor\frac{d}{2}\right\rfloor$. For each $0 \leq t \leq\binom{\mathbf{v}+\left\lfloor\frac{d}{2}\right\rfloor-1}{\left\lfloor\frac{d}{2}\right\rfloor}$, let

$$
\operatorname{Gor}(t):=\left\{P \in S^{d} V \left\lvert\, \operatorname{rank} P_{\left\lfloor\frac{d}{2}\right\rfloor,\left\lceil\frac{d}{2}\right\rceil}=t\right.\right\}
$$

("Gor" is after Gorenstein, see [IK99].) Note that $S^{d} V=\sqcup_{t} \operatorname{Gor}(t)$. Since this is a finite union there must be exactly one $t_{0}$ such that $\overline{\operatorname{Gor}\left(t_{0}\right)}=S^{d} V$. We want to show that $t_{0}=\binom{\mathbf{v}+\left\lfloor\frac{d}{2}\right\rfloor-1}{\left\lfloor\frac{d}{2}\right\rfloor}$. I will do this by computing conormal spaces as $N_{P}^{*} \operatorname{Gor}\left(t_{0}\right)=0$ for $P \in \operatorname{Gor}\left(t_{0}\right)$. Now, for any $t$, the subspace $N_{P}^{*} G o r(t) \subset S^{d} V$ satisfies
$N_{P}^{*} \operatorname{Gor}(t) \subset N_{P_{\left\lfloor\frac{d}{2}\right\rfloor,\left\lceil\frac{d}{2}\right\rceil}^{*}} \sigma_{t}=N_{P_{\left\lfloor\frac{d}{2}\right\rfloor,\left\lceil\frac{d}{2}\right\rceil}^{*}} \sigma_{t}\left(S e g\left(\mathbb{P} S^{\left\lfloor\frac{d}{2}\right\rfloor} V \times \mathbb{P} S^{\left\lceil\frac{d}{2}\right\rceil} V\right)\right) \subset S^{\left\lfloor\frac{d}{2}\right\rfloor} V \otimes S^{\left\lceil\frac{d}{2}\right\rceil} V$,
and $N_{P}^{*} G o r(t)$ is simply the image of $N_{P_{\left\lfloor\frac{d}{2}\right\rfloor,\left\lceil\frac{d}{2}\right\rceil}^{*}} \sigma_{t}$ under the multiplication map $S^{\left\lfloor\frac{d}{2}\right\rfloor} V \otimes S^{\left\lceil\frac{d}{2}\right\rceil} V \rightarrow S^{d} V^{*}$. On the other hand, by Exercise 6.2.3.1,

$$
N_{P_{\left\lfloor\frac{d}{2}\right\rfloor,\left\lceil\frac{d}{2}\right\rceil}^{*}}^{*} \sigma_{t}=\operatorname{ker} P_{\left\lfloor\frac{d}{2}\right\rfloor,\left\lceil\frac{d}{2}\right\rceil} \otimes \operatorname{ker} P_{\left\lceil\frac{d}{2}\right\rceil \backslash\left\lfloor\frac{d}{2}\right\rfloor} .
$$

In order for $N_{P}^{*} \operatorname{Gor}(t)$ to be zero, we need $N_{P_{\left\lfloor\frac{d}{2}\right\rfloor,\left\lceil\frac{d}{2}\right\rceil}^{*}} \sigma_{t}$ to be zero (otherwise there will be something nonzero in the image of the symmetrization map: if $d$ is odd, the two degrees are different and this is clear. If $d$ is even, the conormal space is the tensor product of a vector space with itself), which implies ker $P_{\left\lceil\frac{d}{2},\left\lfloor\left\lfloor\frac{d}{2}\right\rfloor\right.\right.}=0$, and thus $t_{0}=\binom{\mathbf{v}+\left\lfloor\frac{d}{2}\right\rfloor-1}{\left\lfloor\frac{d}{2}\right\rfloor}$.

Note that the maximum symmetric border rank (in all but a few known exceptions) is $\left\lceil\frac{1}{\mathbf{v}}\binom{\mathbf{v}+d-1}{d}\right\rceil$, whereas flattenings only give equations up to symmetric border rank $\binom{\mathbf{v}+\left\lfloor\frac{d}{2}\right\rfloor-1}{\left\lfloor\frac{d}{2}\right\rfloor}$.

### 6.3. Singular loci and Jacobian varieties

As mentioned above, the geometry of the hypersurfaces $Z \operatorname{eros}\left(\operatorname{det}_{n}\right)$ and Zeros $\left(\operatorname{perm}_{m}\right)$ will aid us in comparing the complexity of the determinant and permanent. A simple invariant that will be useful is the dimension of the singular set of a hypersurface. The definition presented in $\S 3.1 .3$ of the singular locus results in a singular locus whose dimension is not upper semi-continuous under degeneration. I first give a new definition that is semi-continuous under degeneration.

### 6.3.1. Definition of the (scheme theoretic) singular locus.

Definition 6.3.1.1. Say a variety $X=\left\{P_{1}=0, \ldots, P_{s}=0\right\} \subset \mathbb{P} V$ has codimension $c$, using the definition of codimension in $\S 3.1 .5$. Then $x \in X$ is a singular point if $d P_{1, x}, \ldots, d P_{s, x}$ fail to span a space of dimension $c$. Let $X_{\text {sing }} \subset X$ denote the singular points of $X$. In particular, if $X=\operatorname{Zeros}(P)$ is a hypersuface and $x \in X$, then $x \in X_{\text {sing }}$ if and only if $d P_{x}=0$. Note that $X_{\text {sing }}$ is also the zero set of a collection of polynomials.

Warning: This definition is a property of the ideal generated by the polynomials $P_{1}, \ldots, P_{s}$, not of $X$ as a set. For example every point of $\left(x_{1}^{2}+\right.$ $\left.\cdots+x_{n}^{2}\right)^{2}=0$ is a singular point. In the language of algebraic geometry, one refers to the singular point of the scheme associated to the ideal generated by $\left\{P_{1}=0, \ldots, P_{s}=0\right\}$.
"Most" hypersurfaces $X \subset \mathbb{P} V$ are smooth, in the sense that $\{P \in$ $\left.\mathbb{P} S^{d} V \mid \operatorname{Zeros}(P)_{\text {sing }} \neq \emptyset\right\} \subset \mathbb{P} S^{d} V$ is a hypersurface, see, e.g., [Lan12, $\S 8.2 .1]$. The dimension of $\operatorname{Zeros}(P)_{\text {sing }}$ is a measure of the pathology of $P$.

Singular loci will also be used in the determination of symmetry groups.
6.3.2. Jacobian varieties. While the ranks of symmetric flattenings are the same for the permanent and determinant, by looking more closely at the maps, we can extract geometric information that distinguishes them.

First, for $P \in S^{n} V$, consider the images $P_{k, n-k}\left(S^{k} V^{*}\right) \subset S^{n-k} V$. This is a space of polynomials and we can consider the ideal they generate, called the $k$-th Jacobian ideal of $P$, and the common zero set of these polynomials is called the $k$-th Jacobian variety of $P$ :

$$
\operatorname{Zeros}(P)_{J a c, k}:=\left\{[\alpha] \in \mathbb{P} V^{*} \mid q(\alpha)=0 \forall q \in P_{k, n-k}\left(S^{k} V^{*}\right)\right\} .
$$

Exercise 6.3.2.1: (1) Show that $\operatorname{Zeros}\left(\operatorname{det}_{n}\right)_{J a c, k}$ is $\sigma_{n-k-1}\left(\operatorname{Seg}\left(\mathbb{P}^{n-1} \times\right.\right.$ $\left.\mathbb{P}^{n-1}\right)$ ), the matrices of rank at most $n-k-1$.

It is not known what the varieties $\mathrm{Zeros}\left(\mathrm{perm}_{m}\right)_{J a c, k}$ are in general. I explicitly determine $\mathrm{Zeros}\left(\operatorname{perm}_{m}\right)_{J a c, m-2}$ in Lemma 6.3.3.4 below as it is used to prove the symmetries of the permanent are what we expect them to be.
6.3.3. Singularities of $\operatorname{Zeros}\left(\operatorname{perm}_{m}\right)$. In contrast to the determinant, the singular set of the permanent is not understood; even its codimension is not known. The problem is more difficult because, unlike in the determinant case, we do not have normal forms for points on $\mathrm{Zeros}\left(\operatorname{perm}_{m}\right)$. In this section I show that codim $\left(\right.$ Zeros $\left.\left(\text { perm }_{m}\right)_{\text {sing }}\right) \geq 5$.
Exercise 6.3.3.1: (1!) Show that the permanent admits a "Laplace type" expansion similar to that of the determinant.

Exercise 6.3.3.2: (2) Show that $\operatorname{Zeros}\left(\text { perm }_{m}\right)_{\text {sing }}$ consists of the $m \times m$ matrices with the property that all size $m-1$ sub-matrices of it have permanent zero.

Exercise 6.3.3.3: (1) Show that $\operatorname{Zeros}\left(\text { perm }_{m}\right)_{\text {sing }}$ has codimension at most $2 m$ in $\mathbb{C}^{m^{2}}$ 。 ©

Since Zeros $\left(\operatorname{perm}_{2}\right)_{\text {sing }}=\emptyset$, let's start with perm ${ }_{3}$. Since we will need it later, I prove a more general result:

Lemma 6.3.3.4. The variety $\operatorname{Zeros}\left(\operatorname{perm}_{m}\right)_{J a c, m-2}$ is the union of the following varieties:
(1) Matrices $A$ with all entries zero except those in a single size 2 submatrix, and that submatrix has zero permanent.
(2) Matrices $A$ with all entries zero except those in the $j$-th row for some $j$.
(3) Matrices $A$ with all entries zero except those in the $j$-th column for some $j$.
In other words, let $X \subset \operatorname{Mat}_{m}(\mathbb{C})$ denote the subvariety of matrices that are zero except in the upper $2 \times 2$ corner and that $2 \times 2$ submatrix has zero permanent, and let $Y$ denote the variety of matrices that are zero except in the first row, then

$$
\begin{equation*}
\operatorname{Zeros}\left(\operatorname{perm}_{m}\right)_{J a c, m-2}=\bigcup_{\sigma \in\left(\mathfrak{S}_{m} \times \mathfrak{S}_{m}\right) \times \mathbb{Z}_{2}} \sigma \cdot X \cup \sigma \cdot Y \tag{6.3.1}
\end{equation*}
$$

Here $\mathfrak{S}_{m} \times \mathfrak{S}_{m}$ acts by left and right multiplication by permutation matrices and the $\mathbb{Z}_{2}$ is generated by sending a matrix to its transpose.

The proof is straight-forward. Here is the main idea: Take a matrix with entries that don't fit that pattern, e.g., one that begins

$$
\begin{array}{lll}
a & b & e \\
* & d & *
\end{array}
$$

and note that it is not possible to fill in the two unknown entries and have all size two sub-permanents, even in this corner, zero. There are just a few such cases since we are free to act by $\left(\mathfrak{S}_{m} \times \mathfrak{S}_{m}\right) \rtimes \mathbb{Z}_{2} \subset G_{\text {perm }_{m}}$.
Corollary 6.3.3.5.

$$
\left\{\operatorname{perm}_{3}=0\right\}_{\text {sing }}=\bigcup_{\sigma \in\left(\mathfrak{(}_{3} \times \mathfrak{G}_{3}\right) \times \mathbb{Z}_{2}} \sigma \cdot X \cup \sigma \cdot Y .
$$

In particular, all the irreducible components of $\left\{\operatorname{perm}_{3}=0\right\}_{\text {sing }}$ have the same dimension and $\operatorname{codim}\left(\left\{\text { perm }_{3}=0\right\}_{\text {sing }}, \mathbb{C}^{9}\right)=6$.

This equi-dimensionality property already fails for perm $_{4}$ : consider

$$
\left\{\left.\left(\begin{array}{cccc}
x_{1}^{1} & x_{2}^{1} & 0 & 0 \\
x_{1}^{2} & x_{2}^{2} & 0 & 0 \\
0 & 0 & x_{3}^{3} & x_{4}^{3} \\
0 & 0 & x_{3}^{4} & x_{4}^{4}
\end{array}\right) \right\rvert\, x_{1}^{1} x_{2}^{2}+x_{1}^{2} x_{2}^{1}=0, x_{3}^{3} x_{4}^{4}+x_{3}^{4} x_{4}^{3}=0\right\} .
$$

This defines a six dimensional irreducible component of $\left\{\operatorname{perm}_{4}=0\right\}_{\text {sing }}$ which is not contained in either a space of matrices with just two nonzero rows (or columns) or the set of matrices that are zero except for in some $3 \times 3$ submatrix which has zero permanent. In [vzG87] von zur Gathen
states that all components of $\left\{\operatorname{perm}_{4}=0\right\}_{\text {sing }}$ are either of dimension six or eight.

Although we do not know the codimension of $\mathrm{Zeros}\left(\operatorname{perm}_{m}\right)_{s i n g}$, the following estimate will suffice for the application of von zur Gathen's regularity theorem 6.3.4.1 below.
Proposition 6.3.3.6 (von zur Gathen [vzG87]).

$$
\operatorname{codim}\left(\operatorname{Zeros}\left(\operatorname{perm}_{m}\right)_{\text {sing }}, \mathbb{P}^{m^{2}-1}\right) \geq 5 .
$$

Proof. I work by induction on $m$. Since $\operatorname{Zeros}\left(\operatorname{perm}_{2}\right)$ is a smooth quadric, the base case is established. Let $I, J$ be multi-indices of the same size and let $s p(I \mid J)$ denote the sub-permanent of the ( $m-|I|, m-|I|$ ) submatrix omitting the index sets $(I, J)$. Let $C \subset$ Zeros $\left(\text { perm }_{m}\right)_{\text {sing }}$ be an irreducible component of the singular set. If $\left.s p\left(i_{1}, i_{2} \mid j_{1}, j_{2}\right)\right|_{C}=0$ for all $\left(i_{1}, i_{2} \mid j_{1}, j_{2}\right)$, we are done by induction as then $C \subset \bigcup_{i, j} \operatorname{Zeros}(s p(i \mid j))_{s i n g}$. So assume there is at least one size $m-2$ subpermanent that is not identically zero on $C$, without loss of generality assume it is $s p(m-1, m \mid m-1, m)$. We have, via permanental Laplace expansions,

$$
\begin{aligned}
0 & =\left.s p(m, m)\right|_{C} \\
& =\sum_{j=1}^{m-2} x_{m-1}^{j} s p(i, m \mid m-1, m)+x_{m-1}^{m-1} s p(m-1, m \mid m-1, m)
\end{aligned}
$$

so on a Zariski open subset of $C, x_{m-1}^{m-1}$ is a function of the $m^{2}-4$ variables $x_{t}^{s},(s, t) \notin\{(m-1, m-1),(m-1, m),(m, m-1),(m, m)\}$, Similar expansions give us $x_{m}^{m-1}, x_{m-1}^{m}$, and $x_{m}^{m}$ as functions of the other variables, so we conclude $\operatorname{dim} C \leq m^{2}-4$. We need to find one more nonzero polynomial that vanishes identically on $C$ that does not involve the variables $x_{m-1}^{m-1}, x_{m-1}^{m}, x_{m}^{m-1}, x_{m}^{m}$ to obtain another relation and to conclude $\operatorname{dim} C \leq$ $m^{2}-5$. Consider

$$
\begin{aligned}
& s p(m-1, m \mid m-1, m) s p(m-2, m)-s p(m-2, m \mid m-1, m) s p(m-1, m) \\
& \quad-s p(m-2, m-1 \mid m-1, m) s p(m, m) \\
& =-2 x_{m-1}^{m-2} s p(m-2, m-1 \mid m-1, m) s p(m-2, m \mid m-1, m) \\
& \quad+\text { terms not involving } x_{m-1}^{m-2},
\end{aligned}
$$

where we obtained the second line by permanental Laplace expansions in the size $m-1$ subpermanents in the expression, and arranged things such that all terms with $x_{m-1}^{m-1}, x_{m-1}^{m}, x_{m}^{m-1}, x_{m}^{m}$ appearing cancel. Since this expression is a sum of terms divisible by size $m-1$ subpermanents, it vanishes identically on $C$. But $2 x_{m-1}^{m-2} s p(m-2, m-1 \mid m-1, m) s p(m-2, m \mid m-1, m)$ is not the zero polynomial, so the whole expression is not the zero polynomial. Thus we obtain another nonzero polynomial that vanishes identically on $C$ and is
not in the ideal generated by the previous four as it does not involve any of $x_{m-1}^{m-1}, x_{m-1}^{m}, x_{m}^{m-1}, x_{m}^{m}$.

It is embarrassing that the following question is still open:
Question 6.3.3.7. What is codim $\left(\mathrm{Zeros}\left(\operatorname{perm}_{m}\right)_{\text {sing }}\right)$ ?

### 6.3.4. von zur Gathen's regularity theorem and its consequences for lower bounds.

Proposition 6.3.4.1 (von zur Gathen [vzG87], also see [ABV15]). Let $M>4$, and let $P \in S^{m} \mathbb{C}^{M}$ be concise and satisfy codim $\left(\{P=0\}_{\text {sing }}, \mathbb{C}^{M}\right) \geq$ 5. If $P=\operatorname{det}_{n} \circ \tilde{A}$, where $\tilde{A}=\Lambda+A: \mathbb{C}^{M} \rightarrow \mathbb{C}^{n^{2}}$ is an affine linear map with $\Lambda$ constant and $A$ linear, then $\operatorname{rank}(\Lambda)=n-1$.

Proof. I first claim that if $\tilde{A}(y) \in \operatorname{Zeros}\left(\operatorname{det}_{n}\right)_{\text {sing }}$ then $y \in \operatorname{Zeros}(P)_{\text {sing }}$. To see this, note that for any $y \in \mathbb{C}^{M}$, the differential of $P$ at $y$ satisfies (by the chain rule)

$$
\left.d P\right|_{y}=\left.d\left(\operatorname{det}_{n} \circ \tilde{A}\right)\right|_{y}=A^{T}\left(\left.d\left(\operatorname{det}_{n}\right)\right|_{\tilde{A}(y)}\right),
$$

where I have used that $\left.d\left(\operatorname{det}_{n}\right)\right|_{\tilde{A}(y)} \in T_{\tilde{A} \tilde{A}(y)}^{*} \mathbb{C}^{n^{2}} \simeq \mathbb{C}^{n^{2} *}$ and $A^{T}: \mathbb{C}^{n^{2} *} \rightarrow$ $\mathbb{C}^{M^{*}}$ is the transpose of the differential of $\tilde{A}$. In particular, if $\left.d\left(\operatorname{det}_{n}\right)\right|_{\tilde{A}(y)}=0$ then $d P_{y}=0$, which is what we needed to show.

Now by Theorem 3.1.5.1, the set

$$
\tilde{A}\left(\mathbb{C}^{M}\right) \cap \operatorname{Zeros}\left(\operatorname{det}_{n}\right)_{\operatorname{sing}} \subset \mathbb{C}^{n^{2}}
$$

is either empty or of dimension at least $\operatorname{dim}\left(\tilde{A}\left(\mathbb{C}^{M}\right)\right)+\operatorname{dim}\left(\operatorname{Zeros}^{\left.\left(\operatorname{det}_{n}\right)_{s i n g}\right)-}\right.$ $n^{2}=M+\left(n^{2}-4\right)-n^{2}=M-4$. (Here $\tilde{A}$ must be injective as $P$ is concise.) The same is true for $\tilde{A}^{-1}\left(\tilde{A}\left(\mathbb{C}^{M}\right) \cap \operatorname{Zeros}\left(\operatorname{det}_{n}\right)_{\text {sing }}\right)$. But this latter set is contained in $\mathrm{Zeros}(P)_{\operatorname{sing}}$, which is of dimension at most $M-5$, so we conclude it is empty.

Thus for all $y \in \mathbb{C}^{M}, \operatorname{rank}(\tilde{A}(y)) \geq n-1$. In particular $\operatorname{rank}(\tilde{A}(0)) \geq$ $n-1$, but $\tilde{A}(0)=\Lambda$. Finally equality holds because if $\Lambda$ had rank $n$, then $\operatorname{det}\left(\tilde{A}\left(\mathbb{C}^{M}\right)\right)$ would have a constant term.
Exercise 6.3.4.2: (1) Prove that any polynomial $P \in S^{d} \mathbb{C}^{M}$ with singular locus of codimension greater than four must have $\mathrm{dc}(P)>d$.
Proposition 6.3.4.3. [Cai90] Let $F \subset \operatorname{Mat}_{n}(\mathbb{C})$ be an affine linear subspace such that for all $X \in F, \operatorname{rank}(F) \geq n-1$. Then $\operatorname{dim} F \leq\binom{ n+1}{2}+1$.

For the proof, see [Cai90]. Note that Proposition 6.3.4.3 is near optimal as consider $F$ the set of upper triangular matrices with 1's on the diagonal, which has dimension $\binom{n}{2}$.
Exercise 6.3.4.4: (2) Use Proposition 6.3.4.3 to show $\operatorname{dc}\left(\operatorname{perm}_{m}\right) \geq \sqrt{2} m$.

Exercise 6.3.4.5: (2) Let $Q \subset \mathbb{P}^{n+1}$ be a smooth quadric hypersurface of dimension $n$. Show that the maximum dimension of a linear projective space contained in $Q$ is $\left\lfloor\frac{n}{2}\right\rfloor$. ©
Theorem 6.3.4.6 (Alper-Bogart-Velasco [ABV15]). Let $P \in S^{d} \mathbb{C}^{M}$ with $d \geq 3$ and such that codim $\left(\operatorname{Zeros}(P)_{\operatorname{sing}}, \mathbb{C}^{M}\right) \geq 5$. Then

$$
\operatorname{dc}(P) \geq \operatorname{codim}\left(\operatorname{Zeros}(P)_{\text {sing }}, \mathbb{C}^{M}\right)+1
$$

Proof. Let $n=\operatorname{dc}(P)$. Say $P=\operatorname{det}_{n} \circ \tilde{A}$, with $\tilde{A}=\Lambda+A$. By Proposition 6.3.4.1, $\operatorname{rank}(\Lambda)=n-1$, and using $G_{\operatorname{det}_{n}}$, we may assume $\Lambda$ is normalized to the matrix that is zero everywhere but the diagonal, where it has one's except in the $(1,1)$-slot where it is zero. Expand $\operatorname{det}(\tilde{A}(y))=p_{0}+p_{1}+\cdots+p_{n}$ as a sum of homogeneous polynomials. Since the right hand side equals $P$, we must have $p_{j}=0$ for $j<d$. Then $p_{0}=\operatorname{det}(\Lambda)=0$ and $p_{1}=A_{1}^{1}$. Now $p_{2}=\sum_{i=2}^{n} A_{i}^{1} A_{1}^{i}=0$ and more generally, each $p_{j}$ is a sum of monomials, each of which contains an element in the first column and an element in the first row of $A$. Each $A_{j}^{i}$ is a linear form on $\mathbb{C}^{M}$ and as such, we can consider the intersection of their kernels. Write $\Gamma=\cap_{i=1}^{n-1}\left(\operatorname{ker} A_{1}^{i}\right) \cap\left(\operatorname{ker} A_{i}^{1}\right)$. Then $\Gamma \subset \operatorname{Zeros}(P)_{\text {sing }}$. Consider the $A_{i}^{1}, A_{1}^{j}$ as coordinates on $\mathbb{C}^{2(n-1)}, p_{2}$ defines a smooth quadric hypersurface in $\mathbb{P}^{2(n-1)-1}$. By Exercise 6.3.4.5, the maximum dimension of a linear space on such a quadric is $n-1$, so the rank of the linear map $\mathbb{C}^{M} \rightarrow \mathbb{C}^{2(n-1)}$ given by $y \mapsto\left(A_{i}^{1}(y), A_{1}^{j}(y)\right)$ is at most $n-1$. But $\Gamma$ is the kernel of this map. We have

$$
n-1 \geq \operatorname{codim}(\Gamma) \geq \operatorname{codim}\left(\operatorname{Zeros}(P)_{\text {sing }}, \mathbb{C}^{M}\right)
$$

and recalling that $n=\operatorname{dc}(P)$ we conclude.
Exercise 6.3.4.7: (2) Prove that $\operatorname{codim}\left(\left(\operatorname{perm}_{m}\right)_{s i n g}\right)=2 m$ when $m=3,4$.
Corollary 6.3.4.8. $[\mathbf{A B V 1 5}] \mathrm{dc}\left(\operatorname{perm}_{3}\right)=7$ and $\mathrm{dc}\left(\operatorname{perm}_{4}\right) \geq 9$.
The upper bound for $\mathrm{dc}\left(\mathrm{perm}_{3}\right)$ is from (1.2.3).
Even if one could prove $\operatorname{codim}\left(\left(\operatorname{perm}_{m}\right)_{\text {sing }}\right)=2 m$ for all $m$, the above theorem would only give a linear bound on $\mathrm{dc}\left(\mathrm{perm}_{m}\right)$. This bound would be obtained from taking one derivative. In the next section, I show that taking two derivatives, one can get a quadratic bound.

### 6.4. Geometry and the state of the art regarding $\operatorname{dc}\left(\operatorname{perm}_{m}\right)$

In mathematics, one often makes transforms to reorganize information, such as the Fourier transform. There are geometric transforms to "reorganize" the information in an algebraic variety. Taking the Gauss image (dual variety) of a hypersurface is one such, as I now describe.
6.4.1. Gauss maps. A classical construction for the geometry of surfaces in 3 -space, is the Gauss map that maps a point of the surface to its unit normal vector on the unit sphere as in Figure 3.


Figure 6.4.1. The shaded area of the surface maps to the shaded area of the sphere.

This Gauss image can be defined for a surface in $\mathbb{P}^{3}$ without the use of a distance function if one instead takes the union of all conormal lines (see $\S 6.2 .3)$ in $\mathbb{P}^{3 *}$. Let $S^{\vee} \subset \mathbb{P}^{3 *}$ denote this Gauss image, also called the dual variety of $S$. One loses qualitative information in this setting, however one still has the information of the dimension of $S^{\vee}$.

This dimension will drop if through all points of the surface there is a curve along which the tangent plane is constant. For example, if $M$ is a cylinder, i.e., the union of lines in three space perpendicular to a plane curve, the Gauss image is a curve:


Figure 6.4.2. Lines on the cylinder are collapsed to a point.

The extreme case is when the surface is a plane, then its Gauss image is just a point.

### 6.4.2. What do surfaces with degenerate Gauss maps "look like"?

 Here is a generalization of the cylinder above: Consider a curve $C \subset \mathbb{P}^{3}$, and a point $p \in \mathbb{P}^{3}$. Define the cone over $C$ with vertex $p$,$$
J(C, p):=\left\{[x] \in \mathbb{P}^{3} \mid x=y+\bar{p} \text { for some } y \in \hat{C}, \bar{p} \in \hat{p}\right\} .
$$



Exercise 6.4.2.1: (1) Show that if $p \neq y, \hat{T}_{[\bar{y}+\bar{p}]} J(C, p)=\operatorname{span}\left\{\hat{T}_{y} C, \hat{p}\right\}$.
Thus the tangent space to the cone is constant along the rulings, and the surface only has a curves worth of tangent (hyper)-planes, so its dual variety is degenerate.
Exercise 6.4.2.2: (2) More generally, let $X \subset \mathbb{P} V$ be an irreducible variety and let $L \subset \mathbb{P} V$ be a linear space. Define $J(X, L)$, the cone over $X$ with vertex $L$ analogously. Show that given $x \in X_{\text {smooth }}$, with $x \notin L$, the tangent space to $J(X, L)^{\vee}$ at $[\bar{x}+\bar{\ell}]$ is constant for all $\ell \in L$.

Here is another type of surface with a degenerate Gauss map: Consider again a curve $C \subset \mathbb{P}^{3}$, and this time let $\tau(C) \subset \mathbb{P}^{3}$ denote the Zariski closure of the union of all points on $\mathbb{P} \hat{T}_{x} C$ as $x$ ranges over the smooth points of $C$. The variety $\tau(C)$ is called the tangential variety to the curve $C$.


Exercise 6.4.2.3: (2) Show that if $y_{1}, y_{2} \in \tau(C)$ are both on a tangent line to $x \in C$, then $\hat{T}_{y_{1}} \tau(C)=\hat{T}_{y_{2}} \tau(C)$, and thus $\tau(C)^{\vee}$ is degenerate. ©

In 1910 C. Segre proved that the above two examples are the only surfaces with degenerate dual varieties:
Theorem 6.4.2.4. [Seg10, p. 105] Let $S^{2} \subset \mathbb{P}^{3}$ be a surface with degenerate Gauss image. Then $S$ is one of the following:
(1) A linearly embedded $\mathbb{P}^{2}$,
(2) A cone over a curve $C$,
(3) A tangential variety to a curve $C$.
(1) is a special case of both (2) and (3) and is the only intersection of the two.

The proof is differential-geometric, see [IL16b, §4.3].
6.4.3. Dual varieties. If $X \subset \mathbb{P} V$ is an irreducible hypersurface, the Zariski closure of its Gauss image will be a projective subvariety of $\mathbb{P} V^{*}$. Gauss images of hypersurfaces are special cases of dual varieties. For an irreducible variety $X \subset \mathbb{P} V$, define $X^{\vee} \subset \mathbb{P} V^{*}$, the dual variety of $X$, by

$$
\begin{aligned}
X^{\vee}: & =\overline{\left\{H \in \mathbb{P}^{*} \mid \exists x \in X_{\text {smooth }}, \hat{T}_{x} X \subseteq \hat{H}\right\}} \\
& =\overline{\left\{H \in \mathbb{P} V^{*} \mid \exists x \in X_{\text {smooth }}, H \in \mathbb{P} N_{x}^{*} X\right\}}
\end{aligned}
$$

Here $H$ refers both to a point in $\mathbb{P} V^{*}$ and the hyperplane in $\mathbb{P} V$ it determines.

That the dual variety is indeed a variety may be seen by considering the following incidence correspondence:

$$
\mathcal{I}:=\overline{\left\{(x, H) \in X_{\text {smooth }} \times \mathbb{P} V^{*} \mid \mathbb{P} \hat{T}_{x} X \subseteq H\right\}} \subset \mathbb{P} V \times \mathbb{P} V^{*}
$$

and note that its image under the projections to $\mathbb{P} V$ and $\mathbb{P} V^{*}$ are respectively $X$ and $X^{\vee}$. When $X$ is smooth, $\mathcal{I}=\mathbb{P} N^{*} X$, the projectivized conormal bundle. Both projections are surjective regular maps, so by Theorem 3.1.4.1, $X^{\vee}$ is an irreducible variety.
Exercise 6.4.3.1: (2) Show

$$
\mathcal{I}=\overline{\left\{(x, H) \in \mathbb{P} V \times\left(X^{\vee}\right)_{\text {smooth }} \mid \mathbb{P} \hat{T}_{H} X^{\vee} \subseteq x\right\}} \subset \mathbb{P} V \times \mathbb{P} V^{*}
$$

and thus $\left(X^{\vee}\right)^{\vee}=X$. (This is called the reflexivity theorem and dates back to C. Segre.) ©

For our purposes, the most important property of dual varieties is that for a smooth hypersurface other than a hyperplane, its dual variety is also a hypersurface. This will be a consequence of the B. Segre dimension formula
6.4.5.1 below. If the dual of $X \subset \mathbb{P} V$ is not a hypersurface, one says that $X$ has a degenerate dual variety. It is a classical problem to study the varieties with degenerate dual varieties.

Exercise 6.4.2.2 shows that higher dimensional cones have degenerate dual varieties. Griffiths and Harris [GH79] vaguely conjectured a higher dimensional generalization of C. Segre's theorem, namely that a variety with a degenerate dual is "built out of" cones and tangent developables. For example, $\operatorname{Zeros}\left(\operatorname{det}_{n}\right)$ may be thought of as the union of tangent lines to tangent lines to $\ldots$ to the Segre variety $\operatorname{Seg}\left(\mathbb{P}^{n-1} \times \mathbb{P}^{n-1}\right)$, and we will see that it indeed has a degenerate dual variety.

Segre's theorem indicates that if we take the Zariski closure in $\mathbb{P} S^{d} V^{*}$ of the set of irreducible hypersurfaces of degree $d$ with degenerate dual varieties, we will obtain a reducible variety. This will complicate the use of dual varieties for Valiant's hypothesis.

For more on dual varieties see [Lan12, §8.2].
6.4.4. $\operatorname{Zeros}\left(\operatorname{det}_{n}\right)_{\text {sing }}$. As far as singularities are concerned, the determinant is quite pathological: Thanks to $G_{\text {det }_{n}}$, the determination of $\operatorname{Zeros}\left(\operatorname{det}_{n}\right)_{s i n g}$ is easy to describe. Any point of $\operatorname{Zeros}\left(\operatorname{det}_{n}\right)$ is in the $G_{\operatorname{det}_{n}}$-orbit of some

$$
p_{r}:=\left(\begin{array}{cc}
\mathrm{Id}_{r} & 0  \tag{6.4.1}\\
0 & 0
\end{array}\right)
$$

where $1 \leq r \leq n-1$ and the blocking is $(r, n-r) \times(r, n-r)$. The nature of the singularity of $x \in \operatorname{Zeros}\left(\operatorname{det}_{n}\right)$ is the same as that of the corresponding $p_{r}$.

Recall that $\sigma_{r}=\sigma_{r}\left(\operatorname{Seg}\left(\mathbb{P}^{n-1} \times \mathbb{P}^{n-1}\right)\right) \subset \mathbb{P}\left(\mathbb{C}^{n} \otimes \mathbb{C}^{n}\right)$ is the set of matrices (up to scale) of rank at most $r$.

The smooth points of $\operatorname{Zeros}\left(\operatorname{det}_{n}\right)=\sigma_{n-1}$ are those in the $G_{\text {det }_{n}}$-orbit of $p_{n-1}$, as shown by the following exercises:
Exercise 6.4.4.1: (1) Show that $d\left(\operatorname{det}_{n}\right)_{p_{n-1}}=d x_{n}^{n}$.
Exercise 6.4.4.2: (1) Show that $\operatorname{Zeros}\left(\operatorname{det}_{n}\right)_{\text {sing }}=\sigma_{n-2}$.
Exercise 6.4.4.3: (1) Show that $\sigma_{r}=\operatorname{Zeros}\left(\operatorname{det}_{n}\right)_{J a c, n-r}$.
Exercise 6.2.3.1 implies $\operatorname{dim} \sigma_{r}\left(\operatorname{Seg}\left(\mathbb{P}^{u-1} \times \mathbb{P}^{v-1}\right)\right)=r(u+v-r)-1$.
6.4.5. What does this have to do with complexity theory? Having a degenerate dual variety is a pathology, and our dimension calculation below will show that if $Q \in S^{m} \mathbb{C}^{M}$ is an irreducible polynomial such that $Q$ is an affine linear degeneration of an irreducible polynomial $P$, then $\operatorname{dim}\left(\operatorname{Zeros}(Q)^{\vee}\right) \leq \operatorname{dim}\left(\operatorname{Zeros}(P)^{\vee}\right)$.

To determine the dual variety of $\operatorname{Zeros}\left(\operatorname{det}_{n}\right) \subset \mathbb{P}(E \otimes F)$, recall that any smooth point of $\mathrm{Zeros}\left(\operatorname{det}_{n}\right)$ is $G_{\text {det }_{n}}$-equivalent to

$$
p_{n-1}=\left(\begin{array}{cccc}
1 & & & \\
& \ddots & & \\
& & 1 & \\
& & & 0
\end{array}\right) \in \operatorname{Zeros}\left(\operatorname{det}_{n}\right) .
$$

and that

$$
N_{p_{n-1}}^{*} \operatorname{Zeros}\left(\operatorname{det}_{n}\right)=\left(\begin{array}{cccc}
0 & 0 & 0 & 0 \\
\vdots & \ddots & \vdots & \vdots \\
0 & 0 & 0 & 0 \\
0 & 0 & 0 & *
\end{array}\right)
$$

Since any smooth point of $\operatorname{Zeros}\left(\operatorname{det}_{n}\right)$ can be moved to $p_{n-1}$ by a change of basis, we conclude that the tangent hyperplanes to $\operatorname{Zeros}\left(\operatorname{det}_{n}\right)$ are parametrized by the rank one matrices $\operatorname{Seg}\left(\mathbb{P} E^{*} \otimes \mathbb{P} F^{*}\right)$, which has dimension $2 n-2$, because they are obtained by multiplying a column vector by a row vector.
Proposition 6.4.5.1 (B. Segre). Let $P \in S^{d} V^{*}$ be irreducible and let $[x] \in$ Zeros $(P)$ be a general point. Then

$$
\begin{equation*}
\operatorname{dim} \operatorname{Zeros}(P)^{\vee}=\operatorname{rank}\left(P_{d-2,2}\left(x^{d-2}\right)\right)-2 \tag{6.4.2}
\end{equation*}
$$

Here $\left(P_{d-2,2}\left(x^{d-2}\right)\right) \in S^{2} V^{*}$, and we are computing the rank of this symmetric matrix. In coordinates, $P_{d-2,2}$ may be written as a symmetric matrix whose entries are polynomials of degree $d-2$ in the coordinates of $x$, and is called the Hesssian.

Proof. Let $x \in \hat{\operatorname{Zeros}}(P) \subset V$ be a general point, so $P(x)=\bar{P}(x, \ldots, x)=0$ and $d P_{x}=\bar{P}(x, \ldots, x, \cdot) \neq 0$ and take $h=d P_{x} \in V^{*}$, so $[h] \in \operatorname{Zeros}(P)^{\vee}$. Now consider a curve $h_{t} \subset \hat{\operatorname{Zeros}}(P)^{\vee}$ with $h_{0}=h$. There must be a corresponding (possibly stationary) curve $x_{t} \in \hat{\operatorname{Zeros}}(P)$ such that $h_{t}=$ $\bar{P}\left(x_{t}, \ldots, x_{t}, \cdot\right)$ and thus $h_{0}^{\prime}=(d-1) \bar{P}\left(x^{d-2}, x_{0}^{\prime}, \cdot\right)$. Thus the dimension of $\hat{T}_{h} \mathrm{Zeros}(P)^{\vee}$ is the rank of $P_{d-2,2}\left(x^{d-2}\right)$ minus one (we subtract one because we are only allowed to feed in vectors $x_{0}^{\prime}$ that are tangent to $\mathrm{Zeros}(P)$ ). Now just recall that $\operatorname{dim} Z=\operatorname{dim} \hat{T}_{z} Z-1$. We needed $x$ to be general to insure that $[h]$ is a smooth point of $\operatorname{Zeros}(P)^{\vee}$.
Exercise 6.4.5.2: (1) Show that if $Q \in S^{m} \mathbb{C}^{M}$ and there exists $\tilde{A}: \mathbb{C}^{M} \rightarrow$ $\mathbb{C}^{N}$ such that $Q(y)=P(\tilde{A}(y))$ for all $y \in \mathbb{C}^{M *}$, then $\operatorname{rank}\left(Q_{m-2,2}(y)\right) \leq$ $\operatorname{rank}\left(P_{m-2, m}(\tilde{A}(y))\right)$.
Exercise 6.4.5.3: (1) Show that every $P \in S u b_{k}\left(S^{d} V\right)$ has dim $\operatorname{Zeros}(P)^{\vee} \leq$ $k-2$.

Exercise 6.4.5.4: (2) Show that $\sigma_{3}\left(C h_{n}\left(\mathbb{C}^{n^{2}}\right)\right) \not \subset \mathcal{D} e t_{n}$.

Exercise 6.4.5.5: (2) Show that $\sigma_{2 n+1}\left(v_{n}\left(\mathbb{P}^{n^{2}-1}\right)\right) \not \subset \mathcal{D e t}_{n}$.
Exercise 6.4.5.6: (2) Show that $\left\{x_{1} \cdots x_{n}+y_{1} \cdots y_{n}=0\right\} \subset \mathbb{P}^{2 n-1}$ is self dual, in the sense that it is isomorphic to its own dual variety.

To show a hypersurface has a nondegenerate dual variety, it suffices to find a point where the Hessian of its defining equation has maximal rank.
6.4.6. Permanent case. Consider the point

$$
y_{0}=\left(\begin{array}{cccc}
1-m & 1 & \cdots & 1 \\
1 & 1 & \cdots & 1 \\
& \vdots & & \\
1 & 1 & \cdots & 1
\end{array}\right)
$$

Exercise 6.4.6.1: (1!) Show $\operatorname{perm}\left(y_{0}\right)=0$. ©
Now compute $\left(\operatorname{perm}_{m}\right)_{m-2,2}\left(y_{0}\right)$ : First note that

$$
\frac{\partial}{\partial y_{j}^{i}} \frac{\partial}{\partial y_{l}^{k}} \operatorname{perm}_{m}(y)=\left\{\begin{array}{cc}
0 & \text { if } i=k \text { or } j=l \\
\operatorname{perm}_{m-2}\left(y_{\hat{j} \hat{l}}^{\hat{i} \hat{l}}\right) & \text { otherwise }
\end{array}\right.
$$

where $y_{\hat{j} \hat{l} \hat{\hat{k}}}^{\hat{\imath}}$ is the size $(m-2) \times(m-2)$ matrix obtained by removing rows $i, k$ and columns $j, l$.
Exercise 6.4.6.2: (2) Show that if we order indices $y_{1}^{1}, \ldots, y_{1}^{m}, y_{2}^{1}, \ldots, y_{2}^{m}, \ldots, y_{m}^{m}$, then the Hessian matrix of the permanent at $y_{0}$ takes the form

$$
\left(\begin{array}{ccccc}
0 & Q & Q & \cdots & Q  \tag{6.4.3}\\
Q & 0 & R & \cdots & R \\
Q & R & 0 & \ddots & \vdots \\
\vdots & \vdots & \ddots & \ddots & R \\
Q & R & \cdots & R & 0
\end{array}\right)
$$

where
$Q=(m-2)\left(\begin{array}{cccc}0 & 1 & \cdots 1 & \\ 1 & 0 & \ddots & \vdots \\ \vdots & \ddots & \ddots & 1 \\ 1 & \cdots & 1 & 0\end{array}\right), R=\left(\begin{array}{ccccc}0 & m-2 & m-2 & \cdots & m-2 \\ m-2 & 0 & -2 & \cdots & -2 \\ m-2 & -2 & 0 & \ddots & \vdots \\ \vdots & \vdots & \ddots & \ddots & -2 \\ m-2 & -2 & \cdots & -2 & 0\end{array}\right)$.
Lemma 6.4.6.3. Let $Q, R$ be invertible $m \times m$ matrices and let $M$ be an $m^{2} \times m^{2}$ matrix of the form (6.4.3). Then $M$ is invertible.

Proof. Without loss of generality, we may assume $Q=\mathrm{Id}_{m}$ by multipling on the left and right by the block diagonal matrix whose block diagonals
are $Q^{-1}, \operatorname{Id}_{m}, \ldots, \operatorname{Id}_{m}$. Let $v=\left(v_{1}, \ldots, v_{m}\right)^{T}$, where $v_{j} \in \mathbb{C}^{m}$, be a vector in the kernel. Then we have the equations

$$
\begin{aligned}
& v_{2}+\cdots+v_{m}=0 \\
& v_{1}+R v_{3}+\cdots+R v_{m}=0, \\
& \vdots \\
& v_{1}+R v_{2}+\cdots+R v_{m-1}=0,
\end{aligned}
$$

i.e.,

$$
\begin{aligned}
& v_{2}+\cdots+v_{m}=0, \\
& v_{1}-R v_{2}=0, \\
& \vdots \\
& v_{1}-R v_{m}=0 .
\end{aligned}
$$

Multiply the first line by $R$ to conclude $(m-1) v_{1}=0$ and hence $v_{1}=0$, and the remaining equations imply the other $v_{j}=0$.

Thus the permanent hypersurface $\operatorname{Zeros}\left(\operatorname{perm}_{m}\right) \subset \mathbb{P}^{m^{2}-1}$ has a nondegenerate Gauss map. When one includes $\mathbb{C}^{m^{2}} \subset \mathbb{C}^{n^{2}}$, so the equation Zeros $\left(\operatorname{perm}_{m}\right)$ becomes an equation in a space of $n^{2}$ variables that only uses $m^{2}$ of the variables, one gets a cone with vertex $\mathbb{P}^{n^{2}-m^{2}-1}$ corresponding to the unused variables, in particular, the Gauss image will have dimension $m^{2}-2$.

If one makes an affine linear substitution $X=X(Y)$, by the chain rule, the Gauss map of $\{\operatorname{det}(X(Y))=0\}$ will be at least as degenerate as the Gauss map of $\{\operatorname{det}(X)=0\}$ by Exercise 6.4.5.2. Using this, one obtains:
Theorem 6.4.6.4 (Mignon-Ressayre [MR04]). If $n(m)<\frac{m^{2}}{2}$, then there do not exist affine linear functions $x_{j}^{i}\left(y_{t}^{s}\right), 1 \leq i, j \leq n, 1 \leq s, t \leq m$ such that $\operatorname{perm}_{m}(Y)=\operatorname{det}_{n}(X(Y))$. I.e., $\operatorname{dc}\left(\operatorname{perm}_{m}\right) \geq \frac{m^{2}}{2}$.

Remark 6.4.6.5. We saw a linear lower bound by taking one derivative and a quadratic lower bound by taking two. Unfortunately it does not appear to be possible to improve the Mignon-Ressayre bound by taking three derivatives.

### 6.5. Extension of the Mignon-Ressayre result to $\overline{\mathrm{dc}}$

To extend the Mignon-Ressayre theorem to $\overline{\mathrm{dc}}$ we will need to find polynomials on $\mathbb{P} S^{n} V$ that vanish on the hypersurfaces with degenerate dual varieties. This was a classically studied question whose answer was known
only in a very few number of small cases. In this section I present an answer to the classical question and its application to Conjecture 1.2.5.2.
6.5.1. First steps towards equations. Let $P \in S^{d} V^{*}$ be irreducible. Segre's formula (6.4.2) may be restated as: $\operatorname{dim} \operatorname{Zeros}(P)^{\vee} \leq k$ if and only if, for all $w \in V$,

$$
\begin{equation*}
P(w)=0 \Rightarrow \operatorname{det}_{k+3}\left(\left.P_{d-2,2}\left(w^{d-2}\right)\right|_{F}\right)=0 \forall F \in G(k+3, V) . \tag{6.5.1}
\end{equation*}
$$

Here $G(k+3, V)$ is the Grassmannian of $(k+3)$-planes through the origin in $V$ (see Definition 2.3.3.1). Equivalently, for any $F \in G(k+3, V)$, the polynomial $P$ must divide $\operatorname{det}_{k+3}\left(P_{d-2,2} \mid F\right) \in S^{(k+3)(d-2)} V^{*}$, where $\operatorname{det}_{k+3}$ is evaluated on the $S^{2} V^{*}$ factor in $S^{2} V^{*} \otimes S^{d-2} V^{*}$.

Thus to find polynomials on $S^{d} V^{*}$ characterizing hypersurfaces with degenerate duals, we need polynomials that detect if a polynomial $P$ divides a polynomial $Q$. Now, $P \in S^{d} V^{*}$ divides $Q \in S^{e} V^{*}$ if and only if $Q \in$ $P \cdot S^{e-d} V^{*}$, i.e.

$$
x^{I_{1}} P \wedge \cdots \wedge x^{I_{D}} P \wedge Q=0
$$

where $x^{I_{j}}$, is a basis of $S^{e-d} V$ (and $D=\binom{\mathbf{v}+e-d-1}{e-d}$ ). Let $\mathcal{D}_{k, d, N} \subset \mathbb{P} S^{d} \mathbb{C}^{N}$ denote the zero set of these equations when $Q=\operatorname{det}_{k+3}\left(\left.P_{d-2,2}\right|_{F}\right)$ as $F$ ranges over $G(k+3, V)$.

Define Dual $_{k, d, N} \subset \mathbb{P}\left(S^{d} V^{*}\right)$ as the Zariski closure of the set of irreducible hypersurfaces of degree $d$ in $\mathbb{P} V \simeq \mathbb{P}^{N-1}$, whose dual variety has dimension at most $k$. Our discussion above implies Dual $_{k, d, N} \subseteq \mathcal{D}_{k, d, N}$.

Note that

$$
\begin{equation*}
\left[\operatorname{det}_{n}\right] \in \operatorname{Dual}_{2 n-2, n, n^{2}} \subseteq \mathcal{D}_{2 n-2, n, n^{2}} \tag{6.5.2}
\end{equation*}
$$

6.5.2. The lower bound on $\overline{d c}\left(\operatorname{perm}_{m}\right)$. The calculation of $\S 6.4 .6$ shows that perm ${ }_{m-2,2}\left(y_{0}^{m-2}\right)$ is of maximal rank. Here we don't have perm ${ }_{m}$, but rather $\ell^{n-m}$ perm $_{m}$.
Proposition 6.5.2.1. Let $U=\mathbb{C}^{M}$, let $R \in S^{m} U^{*}$ be irreducible, let $\ell$ be a coordinate on $L \simeq \mathbb{C}^{1}$ be nonzero, let $U^{*} \oplus L^{*} \subset \mathbb{C}^{N *}$ be a linear inclusion.

If $[R] \in \mathcal{D}_{\kappa, m, M}$ and $[R] \notin \mathcal{D}_{\kappa-1, m, M}$, then $\left[\ell^{d-m} R\right] \in \mathcal{D}_{\kappa, d, N}$ and $\left[\ell^{d-m} R\right] \notin \mathcal{D}_{\kappa-1, d, N}$.

Proof. Let $u_{1}, \ldots, u_{M}, v, w_{M+2}, \ldots, w_{N}$ be a basis of $\mathbb{C}^{N}$ adapted to the inclusions $\mathbb{C}^{M} \subset \mathbb{C}^{M+1} \subset \mathbb{C}^{N}$, so $\left(U^{*}\right)^{\perp}=\left\langle w_{M+2}, \ldots, w_{N}\right\rangle$ and $\left(L^{*}\right)^{\perp}=$
$\left\langle u_{1}, \ldots, u_{M}, w_{M+2}, \ldots, w_{N}\right\rangle$. Let $c=(d-m)(d-m-1)$. In these coordinates, the matrix of $\left(\ell^{d-m} R\right)_{d-2,2}$ in $(M, 1, N-M-1) \times(M, 1, N-M-1)$ block form:

$$
\left(\ell^{d-m} R\right)_{d-2,2}=\left(\begin{array}{ccc}
\ell^{d-m} R_{m-2,2} & \ell^{d-m-1} R_{m-1,1} & 0 \\
\ell^{d-m-1} R_{m-1,1} & c \ell^{d-m-2} R & 0 \\
0 & 0 & 0
\end{array}\right) .
$$

First note that $\operatorname{det}_{M+1}\left(\left.\left(\ell^{d-m} R\right)_{d-2,2}\right|_{F}\right)$ for any $F \in G\left(M+1, \mathbb{C}^{N}\right)$ is either zero or a multiple of $\ell^{d-m} R$. If $\operatorname{dim} \operatorname{Zeros}(R)^{\vee}=M-2$ (the expected dimension), then for a general $F \in G\left(M+1, \mathbb{C}^{N}\right)$, $\operatorname{det}_{M}\left(\left.\left(\ell^{d-m} R\right)_{d-2,2}\right|_{F}\right)$ will not be a multiple of $\left(\ell^{d-m} R\right)_{d-2,2}$, and more generally if $\operatorname{dim} \operatorname{Zeros}(R)^{\vee}=\kappa$, then for a general $F \in G\left(\kappa+2, \mathbb{C}^{N}\right), \operatorname{det}_{\kappa+2}\left(\left.\left(\ell^{d-m} R\right)_{d-2,2}\right|_{F}\right)$ will not be a multiple of $\ell^{d-m} R$ but for any $F \in G\left(\kappa+3, \mathbb{C}^{N}\right)$, $\operatorname{det}_{\kappa+3}\left(\left.\left(\ell^{d-m} R\right)_{d-2,2}\right|_{F}\right)$ will be a multiple of $\ell^{d-m} R$. This shows $[R] \notin \mathcal{D}_{\kappa-1, m, M}$, implies $\left[\ell^{d-m} R\right] \notin$ $\mathcal{D}_{\kappa-1, d, N}$.
Exercise 6.5.2.2: (1) Show that $[R] \in \mathcal{D}_{\kappa, m, M}$, implies $\left[\ell^{d-m} R\right] \in \mathcal{D}_{\kappa, d, N}$. ©

The inclusion (6.5.2) and Proposition 6.5.2.1 imply:
Theorem 6.5.2.3. [LMR13] $\mathcal{P e r m}_{n}^{m} \not \subset \mathcal{D}_{2 n-2, n, n^{2}}$ when $m<\frac{n^{2}}{2}$. In particular, $\overline{d c}\left(\right.$ perm $\left._{m}\right) \geq \frac{m^{2}}{2}$.

On the other hand, by Exercise 6.4.5.3 cones have degenerate duals, so $\ell^{n-m} \operatorname{perm}_{m} \in \mathcal{D}_{2 n-2, n, n^{2}}$ whenever $m \geq \frac{n^{2}}{2}$.

The next step from this perspective would be:
Problem 6.5.2.4. Find equations that distinguish cones (e.g. $\operatorname{Zeros}\left(\ell^{n-m} \operatorname{perm}_{m}\right) \subset$ $\mathbb{P}^{n^{2}-1}$ ) from tangent developables (e.g., $\left.\operatorname{Zeros}\left(\operatorname{det}_{n}\right) \subset \mathbb{P}^{n^{2}-1}\right)$. More precisely, find equations that are zero on tangent developables but nonzero on cones.
6.5.3. A better module of equations. The equations above are of enormous degree. I now derive equations of much lower degree. Since $P \in S^{d} \mathbb{C}^{N}$ divides $Q \in S^{e} \mathbb{C}^{N}$ if and only if for each $L \in G\left(2, \mathbb{C}^{N}\right),\left.P\right|_{L}$ divides $\left.Q\right|_{L}$, it will be sufficient to solve this problem for polynomials on $\mathbb{C}^{2}$. This will have the advantage of producing polynomials of much lower degree.

Let $d \leq e$, let $P \in S^{d} \mathbb{C}^{2}$ and $Q \in S^{e} \mathbb{C}^{2}$. If $P$ divides $Q$ then $S^{e-d} \mathbb{C}^{2} \cdot P$ will contain $Q$. That is,

$$
x^{e-d} P \wedge x^{e-d-1} y P \wedge \cdots \wedge y^{e-d} P \wedge Q=0
$$

Since $\operatorname{dim} S^{e} \mathbb{C}^{2}=e+1$, these potentially give a $\binom{e+1}{e-d+2}$-dimensional vector space of equations, of degree $e-d+1$ in the coefficients of $P$ and linear in the coefficients of $Q$.

By taking our polynomials to be $P=\left.P\right|_{L}$ and $Q=\left.\operatorname{det}_{k+3}\left(\left.P_{n-2,2}\right|_{F}\right)\right|_{L}$ for $F \in G(k+3, V)$ and $L \in G(2, F)$ (or, for those familiar with flag varieties, better to say $\left.(L, F) \in F l a g_{2, k+3}(V)\right)$ we now have equations parametrized by the pairs $(L, F)$. Note that $\operatorname{deg}(Q)=e=(k+3)(d-2)$. These were the polynomials that were used in [LMR13].
Remark 6.5.3.1. More generally, given $P \in S^{d} \mathbb{C}^{2}, Q \in S^{e} \mathbb{C}^{2}$, one can ask if $P, Q$ have at least $r$ roots in common (counting multiplicity). Then $P, Q$ having $r$ points in common says the spaces $S^{e-r} \mathbb{C}^{2} \cdot P$ and $S^{d-r} \mathbb{C}^{2} \cdot Q$ intersect. That is,

$$
x^{e-r} P \wedge x^{e-r-1} y P \wedge \cdots \wedge y^{e-r} P \wedge x^{d-r} Q \wedge x^{d-r-1} y Q \wedge \cdots \wedge y^{d-r} Q=0
$$

In the case $r=1$, we get a single polynomial, called the resultant, which is of central importance. In particular, the proof of Noether normalization from $\S 3.1 .4$, that the projection of a projective variety $X \subset \mathbb{P} W$ from a point $y \in \mathbb{P} W$ with $y \notin X$, to $\mathbb{P}(W / \hat{y})$ is still a projective variety, relies on the resultant to produce equations for the projection.

### 6.6. Symmetries of the determinant and permanent

The permanent and determinant both have the property that they are characterized by their symmetry groups in the sense described in §1.2.5. I expect these symmetry groups to play a central role in the study of Valiant's hypothesis in future work. For example, the only known exponential separation of the permanent from the determinant in any restricted model (as defined in Chapter 7), is the model of equivariant determinantal complexity, which is defined in terms of symmetry groups, see $\S 7.4 .1$.

### 6.6.1. Symmetries of the determinant.

Theorem 6.6.1.1 (Frobenius [Fro97]). Write $\rho: G L_{n^{2}} \rightarrow G L\left(S^{n} \mathbb{C}^{n^{2}}\right)$ for the induced action. Let $\phi \in G L_{n^{2}}$ be such that $\rho(\phi)\left(\operatorname{det}_{n}\right)=\operatorname{det}_{n}$. Then, identifying $\mathbb{C}^{n^{2}}$ with the space of $n \times n$ matrices,

$$
\phi(z)= \begin{cases}g z h, & \text { or } \\ g z^{T} h & \end{cases}
$$

for some $g, h \in G L_{n}$, with $\operatorname{det}_{n}(g) \operatorname{det}_{n}(h)=1$. Here $z^{T}$ denotes the transpose of $z$.

I present the proof from [Die49] below.
Write $\mathbb{C}^{n^{2}}=E \otimes F=\operatorname{Hom}\left(E^{*}, F\right)$ with $E, F=\mathbb{C}^{n}$. Let $\mathbb{Z}_{n}$ denote the cyclic group of order $n$ and consider the inclusion $\mathbb{Z}_{n} \times \mathbb{Z}_{n} \subset G L(E) \times G L(F)$
given by the $n$-th roots of unity times the identity matrix. Let $\mu_{n}$ denote the kernel of the product map $\left(\mathbb{Z}_{n}\right)^{\times 2} \rightarrow \mathbb{Z}_{n}$.
Corollary 6.6.1.2. $G_{\operatorname{det}_{n}}=(S L(E) \times S L(F)) / \mu_{n} \rtimes \mathbb{Z}_{2}$
To prove the Corollary, just note that the $\mathbb{C}^{*}$ corresponding to $\operatorname{det}(g)$ above and $\mu_{n}$ are the kernel of the map $\mathbb{C}^{*} \times S L(E) \times S L(F) \rightarrow G L(E \otimes F)$.
Exercise 6.6.1.3: (2) Prove the $n=2$ case of Theorem 6.6.1.1. ©
Lemma 6.6.1.4. Let $U \subset E \otimes F$ be a linear subspace such that $U \subset$ $Z \operatorname{Zeros}\left(\operatorname{det}_{n}\right)$. Then $\operatorname{dim} U \leq n^{2}-n$. The subvariety of the Grassmannian $G\left(n^{2}-n, E \otimes F\right)$ consisting of maximal linear spaces on $\mathrm{Zeros}\left(\operatorname{det}_{n}\right)$ has two irreducible components, call them $\Sigma_{\alpha}$ and $\Sigma_{\beta}$, where
$\Sigma_{\alpha}=\left\{X \in G\left(n^{2}-n, E \otimes F\right) \mid \operatorname{ker}(X)=\hat{L}\right.$ for some $\left.L \in \mathbb{P} E^{*}\right\} \simeq \mathbb{P} E^{*}$, and

$$
\begin{equation*}
\Sigma_{\beta}=\left\{X \in G\left(n^{2}-n, E \otimes F\right) \mid \operatorname{Image}(X)=\hat{H} \text { for some } H \in \mathbb{P} F^{*}\right\} \simeq \mathbb{P} F^{*} . \tag{6.6.2}
\end{equation*}
$$

Here for $f \in X, f: E^{*} \rightarrow F$ is considered as a linear map, $\operatorname{ker}(X)$ means the intersections of the kernels of all $f \in X$ and Image $(X)$ is the span of all the images.

Moreover, for any two distinct $X_{j} \in \Sigma_{\alpha}, j=1,2$, and $Y_{j} \in \Sigma_{\beta}$ we have

$$
\begin{align*}
\operatorname{dim}\left(X_{1} \cap X_{2}\right) & =\operatorname{dim}\left(Y_{1} \cap Y_{2}\right)=n^{2}-2 n, \text { and }  \tag{6.6.3}\\
\operatorname{dim}\left(X_{i} \cap Y_{j}\right) & =n^{2}-2 n+1 . \tag{6.6.4}
\end{align*}
$$

Exercise 6.6.1.5: (2) Prove Lemma 6.6.1.4. ©
One can say more: each element of $\Sigma_{\alpha}$ corresponds to a left ideal and each element of $\Sigma_{\beta}$ corresponds to a right ideal in the space of $n \times n$ matrices.

Proof of theorem 6.6.1.1. Let $\Sigma=\Sigma_{\alpha} \cup \Sigma_{\beta}$. Then the automorphism of $G\left(n^{2}-n, E \otimes F\right)$ induced by $\phi$ must preserve $\Sigma$. By the conditions (6.6.3),(6.6.4) of Lemma 6.6.1.4, in order to preserve dimensions of intersections, either every $U \in \Sigma_{\alpha}$ must map to a point of $\Sigma_{\alpha}$, in which case every $V \in \Sigma_{\beta}$ must map to a point of $\Sigma_{\beta}$, or, every $U \in \Sigma_{\alpha}$ must map to a point of $\Sigma_{\beta}$, and every $V \in \Sigma_{\beta}$ must map to a point of $\Sigma_{\alpha}$. If we are in the second case, replace $\phi$ by $\phi \circ T$, where $T(z)=z^{T}$, so we may now assume $\phi$ preserves both $\Sigma_{\alpha}$ and $\Sigma_{\beta}$.

Observe that $\phi$ induces an algebraic map $\phi_{E}: \mathbb{P} E^{*} \rightarrow \mathbb{P} E^{*}$.
Exercise 6.6.1.6: (2) Show that $L_{1}, L_{2}, L_{3} \in \mathbb{P} E$ lie on a $\mathbb{P}^{1}$ if and only if then $\operatorname{dim}\left(U_{L_{1}} \cap U_{L_{2}} \cap U_{L_{3}}\right)=n^{2}-2 n$, where $U_{L}=\{X \mid \operatorname{ker}(X)=L\}$.

In order for $\phi$ to preserve $\operatorname{dim}\left(U_{L_{1}} \cap U_{L_{2}} \cap U_{L_{3}}\right)$, the images of the $L_{j}$ under $\phi_{E}$ must also lie on a $\mathbb{P}^{1}$, and thus $\phi_{E}$ must take lines to lines (and
similarly hyperplanes to hyperplanes). But then, (see, e.g., [Har95, §18, p. 229]) $\phi_{E} \in P G L(E)$, and similarly, $\phi_{F} \in P G L(F)$, where $\phi_{F}: \mathbb{P} F^{*} \rightarrow \mathbb{P} F^{*}$ is the corresponding map. Here $P G L(E)$ denotes $G L(E) / \mathbb{C}^{*}$, the image of $G L(E)$ in its action on projective space. Write $\hat{\phi}_{E} \in G L(E), \hat{\phi}_{F} \in G L(F)$ for any choices of lifts.

Consider the map $\tilde{\phi} \in G L(E \otimes F)$ given by $\tilde{\phi}(z)=\hat{\phi}_{E}{ }^{-1} \phi(z) \hat{\phi}_{F}{ }^{-1}$. The map $\tilde{\phi}$ sends each $U \in \Sigma_{\alpha}$ to itself as well as each $V \in \Sigma_{\beta}$, in particular it does the same for all intersections. Hence it preserves $\operatorname{Seg}(\mathbb{P} E \times \mathbb{P} F) \subset$ $\mathbb{P}(E \otimes F)$ point-wise, so it is up to scale the identity map because $E \otimes F$ is spanned by points of $\hat{S} e g(\mathbb{P} E \times \mathbb{P} F)$.
6.6.2. Symmetries of the permanent. Write $\mathbb{C}^{n^{2}}=E \otimes F$. Let $\Gamma_{n}^{E}:=$ $T_{E}^{S L} \rtimes \mathfrak{S}_{n}$, and similarly for $F$. As discussed in the introduction to this chapter, $\left(\Gamma_{n}^{E} \times \Gamma_{n}^{F}\right) \rtimes \mathbb{Z}_{2} \rightarrow G_{\text {perm }_{n}}$, where the nontrivial element of $\mathbb{Z}_{2}$ acts by sending a matrix to its transpose. We would like to show this map is surjective and determine its kernel. However, it is not when $n=2$.
Exercise 6.6.2.1: (1) What is $G_{\text {perm }_{2}}$ ? ©
Theorem 6.6.2.2. $[\mathbf{M M 6 2}]$ For $n \geq 3, G_{\text {perm }_{n}}=\left(\Gamma_{n}^{E} \times \Gamma_{n}^{F}\right) / \mu_{n} \rtimes \mathbb{Z}_{2}$.
Proof. I follow [Ye11]. Recall the description of $\operatorname{Zeros}\left(\operatorname{perm}_{n}\right)_{J a c, n-2}$ from Lemma 6.3.3.4. Any linear transformation preserving the permanent must send a component of $\operatorname{Zeros}\left(\operatorname{perm}_{n}\right)_{J a c, n-2}$ of type (1) to another of type (1). It must send a component $C^{j}$ either to some $C^{k}$ or some $C_{i}$. But if $i \neq j$, $C^{j} \cap C^{i}=0$ and for all $i, j, \operatorname{dim}\left(C^{i} \cap C_{j}\right)=1$. Since intersections must be mapped to intersections, either all components $C^{i}$ are sent to components $C_{k}$ or all are permuted among themselves. By composing with an element of $\mathbb{Z}_{2}$, we may assume all the $C^{i}$ 's are sent to $C^{i}$ 's and the $C_{j}$ 's are sent to $C_{j}$ 's. Similarly, by composing with an element of $\mathfrak{S}_{n} \times \mathfrak{S}_{n}$ we may assume each $C_{i}$ and $C^{j}$ is sent to itself. But then their intersections are sent to themselves. So we have, for all $i, j$,

$$
\begin{equation*}
\left(x_{j}^{i}\right) \mapsto\left(\lambda_{j}^{i} x_{j}^{i}\right) \tag{6.6.5}
\end{equation*}
$$

for some $\lambda_{j}^{i}$ and there is no summation in the expression. Consider the image of a size 2 submatrix, e.g.,

$$
\begin{array}{lll}
x_{1}^{1} & x_{2}^{1}  \tag{6.6.6}\\
x_{1}^{2} & x_{2}^{2}
\end{array}{ }^{\lambda_{1}^{1} x_{1}^{1}} \begin{array}{ll}
\lambda_{2}^{1} x_{2}^{1} \\
\lambda_{1}^{2} x_{1}^{2} & \lambda_{2}^{2} x_{2}^{2} .
\end{array}
$$

In order that the map (6.6.5) is given by an element of $G_{\text {perm }_{n}}$, when $\left(x_{j}^{i}\right) \in$ Zeros $\left(\operatorname{perm}_{n}\right)_{J a c, n-2}$, the permanent of the matrix on the right hand side of (6.6.6) must be zero. Using that $x_{1}^{1} x_{2}^{2}+x_{2}^{1} x_{1}^{2}=0$, the permanent of the right hand side of (6.6.6) is $\lambda_{1}^{1} \lambda_{2}^{2} x_{1}^{1} x_{2}^{2}+\lambda_{1}^{2} \lambda_{2}^{1} x_{2}^{1} x_{1}^{2}=x_{1}^{1} x_{2}^{2}\left(\lambda_{1}^{1} \lambda_{2}^{2}-\lambda_{1}^{2} \lambda_{2}^{1}\right)$ which implies $\lambda_{1}^{1} \lambda_{2}^{2}-\lambda_{2}^{1} \lambda_{1}^{2}=0$, thus all the $2 \times 2$ minors of the matrix $\left(\lambda_{j}^{i}\right)$ are zero,
so it has rank one and is the product of a column vector and a row vector, but then it arises from $x \mapsto t x t^{\prime}$ with $t, t^{\prime}$ diagonal, and for the permanent to be preserved, $\operatorname{det}(t) \operatorname{det}\left(t^{\prime}\right)=1$. Without loss of generality, we may insist both determininants equal one.
6.6.3. Grenet's decomposition: symmetry and the best upper bound on dc $\left(\right.$ perm $\left._{m}\right)$. Recall from Chapter 4 that the symmetries of the matrix multiplication tensor appear in the optimal and conjecturally optimal rank expressions for it. Will the same be true for determinantal expressions of polynomials, in particular of the permanent?

The best known determinantal expression of perm $_{m}$ is of size $2^{m}-1$ and is due to Grenet [Gre11]. (Previously Valiant [Val79] had shown there was an expression of size $4^{m}$.) We saw (Corollary 6.3.4.8) that when $m=3$ this is the best expression. This motivated N. Ressayre and myself to try to understand Grenet's expression. We observed the following equivariance property:

Let $G \subseteq G_{\text {perm }_{m}}$ I will say a determinantal expression for perm ${ }_{m}$ is $G$-equivariant if given $g \in G$, there exist $n \times n$ matrices $B, C$ such that $\tilde{A}_{\text {Grenet }, m}(g \cdot Y)=B \tilde{A}_{\text {Grenet }, m}(Y) C$ or $B \tilde{A}_{\text {Grenet }, m}(Y)^{T} C$. In other words, there exists an injective group homomorphism $\psi: G \rightarrow G_{\operatorname{det}_{n}}$ such that $\tilde{A}_{\text {Grenet }, m}(Y)=\psi(g)\left(\tilde{A}_{\text {Grenet }, m}(g Y)\right)$.
Proposition 6.6.3.1. [LR15] Grenet's expressions $\tilde{A}_{\text {Grenet }}: \operatorname{Mat}_{m}(\mathbb{C}) \rightarrow$ $\operatorname{Mat}_{n}(\mathbb{C})$ such that $\operatorname{perm}_{m}(Y)=\operatorname{det}_{n}\left(\tilde{A}_{\text {Grenet }}(Y)\right)$ are $\Gamma_{m}^{E}$-equivariant.

For example, let

$$
g(t)=\left(\begin{array}{lll}
t_{1} & & \\
& t_{2} & \\
& & t_{3}
\end{array}\right)
$$

Then $A_{\text {Grenet }, 3}(g(t) Y)=B(t) A_{\text {Grenet }, 3}(Y) C(t)$, where

$$
B(t)=\left(\begin{array}{ccccccc}
t_{3} & & & & & & \\
& t_{1} t_{3} & & & & & \\
& & t_{1} t_{3} & & & & \\
& & & t_{1} t_{3} & & & \\
& & & & 1 & & \\
& & & & & 1 & \\
& & & & & & 1
\end{array}\right) \text { and } C(t)=B(t)^{-1} .
$$

Exercise 6.6.3.2: (2) Determine $B(g)$ and $C(g)$ when $g \in \Gamma_{3}^{E}$ is the permutation $(1,2)$.

Via this equivariance, one can give an invariant description of Grenet's expressions:

The space $S^{k} E$ is an irreducible $G L(E)$-module but it is is not in general irreducible as a $\Gamma_{m}^{E}$-module. Let $e_{1}, \ldots, e_{m}$ be a basis of $E$, and let $\left(S^{k} E\right)_{r e g} \subset S^{k} E$ denote the span of $\prod_{i \in I} e_{i}$, for $I \subset[m]$ of cardinality $k$ (the space spanned by the square-free monomials, also known as the space of regular weights): $\left(S^{k} E\right)_{\text {reg }}$ is an irreducible $\Gamma_{m}^{E}$-submodule of $S^{k} E$. Moreover, there exists a unique $\Gamma_{m}^{E}$-equivariant projection $\pi_{k}$ from $S^{k} E$ to $\left(S^{k} E\right)_{\text {reg }}$.

For $v \in E$, define $s_{k}(v):\left(S^{k} E\right)_{r e g} \rightarrow\left(S^{k+1} E\right)_{r e g}$ to be multiplication by $v$ followed by $\pi_{k+1}$. Alternatively, $\left(S^{k+1} E\right)_{r e g}$ is a $\Gamma_{m}^{E}$-submodule of $E \otimes\left(S^{k} E\right)_{\text {reg }}$, and $s_{k}: E \rightarrow\left(S^{k} E\right)_{\text {reg }}^{*} \otimes\left(S^{k+1} E\right)_{\text {reg }}$ is the unique $\Gamma_{m}^{E-}$ equivariant inclusion.

Fix a basis $f_{1}, \ldots, f_{m}$ of $F^{*}$. If $y=\left(y_{1}, \ldots, y_{m}\right) \in E \otimes F$, let $\left(s_{k} \otimes f_{j}\right)(y):=$ $s_{k}\left(y_{j}\right)$.
Proposition 6.6.3.3. [LR15] The following is Grenet's determinantal representation of perm $_{m}$. Let $\mathbb{C}^{n}=\bigoplus_{k=0}^{m-1}\left(S^{k} E\right)_{\text {reg }}$, so $n=2^{m}-1$, and identify $S^{0} E \simeq\left(S^{m} E\right)_{\text {reg }}$ (both are trivial $\Gamma_{m}^{E}$-modules). Set

$$
\Lambda_{0}=\sum_{k=1}^{m-1} \operatorname{Id}_{\left(S^{k} E\right)_{r e g}}
$$

and define

$$
\begin{equation*}
\tilde{A}=\Lambda_{0}+\sum_{k=0}^{m-1} s_{k} \otimes f_{k+1} \tag{6.6.7}
\end{equation*}
$$

Then $(-1)^{m+1} \operatorname{perm}_{m}=\operatorname{det}_{n} \circ \tilde{A}$. To obtain the permanent exactly, replace $\operatorname{Id}_{\left(S^{1} E\right)_{\text {reg }}}$ by $(-1)^{m+1} \operatorname{Id}_{\left(S^{1} E\right)_{\text {reg }}}$ in the formula for $\Lambda_{0}$.

Moreover the map $\tilde{A}$ is $\Gamma_{m}^{E}$-equivariant.
I prove Proposition 6.6.3.3 in §8.11.1.
Remark 6.6.3.4. In bases respecting the block decomposition induced from the direct sum, the linear part, other than the last term which lies in the upper right block, lies just below the diagonal blocks, and all blocks other than the upper right block and the diagonal and sub-diagonal blocks, are zero. This expression is better viewed as an iterated matrix multiplication as in §7.3.1: $\operatorname{perm}(y)=\left(s_{m-1} \otimes f_{m}(y)\right)\left(s_{m-2} \otimes f_{m-1}(y)\right) \cdots\left(s_{0} \otimes f_{1}(y)\right)$.
6.7. dc v. $\overline{\mathrm{dc}}$

Is conjecture 6.1.6.2 really stronger than Valiant's hypothesis 6.1.6.1? That is, do there exist sequences $\left(P_{m}\right)$ of polynomials with $\overline{\mathrm{dc}}\left(P_{m}\right)$ bounded by a polynomial in $m$ but dc $\left(P_{m}\right)$ growing super-polynomially?
K. Mulmuley [Mul] conjectures that this is indeed the case, and that the existence of such sequences "explains" why Valiant's hypothesis is so difficult.

Before addressing this conjecture, one should at least find a sequence $P_{m}$ with dc $\left(P_{m}\right)>\overline{\mathrm{dc}}\left(P_{m}\right)$. I describe one such sequence in §6.7.2.
6.7.1. On the boundary of the orbit of the determinant. Let $W=$ $\mathbb{C}^{n^{2}}=E^{*} \otimes E$ with $E=\mathbb{C}^{n}$, and let $\sigma_{n^{2}-1}\left(S e g\left(\mathbb{P} W^{*} \times \mathbb{P} W\right)\right) \subset \mathbb{P}\left(W^{*} \otimes W\right)$ be the endomorphisms of $W$ of rank at most $n^{2}-1$ An obvious subset of $\partial \mathcal{D e t}_{n}$
 the orbit closure of the determinant restricted to the traceless matrices. This description shows it has codimension one in $\mathcal{D e} t_{n}$ and is irreducible, so it is a component of $\partial \mathcal{D} e t_{n}$.

Other components can be found as follows: Let $U \subset W$ be a subspace such that $\left.\operatorname{det}_{n}\right|_{U}=0$ and let $V$ be a complement. Given a matrix $M$, write $M=M_{U} \oplus M_{V}$. Introduce a parameter $t$ and consider $M \mapsto \operatorname{det}_{n}\left(M_{U}+t M_{V}\right)$ and expand out in a Taylor series. Say the first non-vanishing term is $t^{k}$, then $M \mapsto \overline{\operatorname{det}_{n}}\left(M_{U}, \ldots, M_{U}, M_{V}, \ldots, M_{V}\right)$ where there are $k M_{V}$ 's, is a point of $\mathcal{D e} t_{n}$ and it is "usually" a point of $\partial \mathcal{D} e t_{n}$. One can do more complicated constructions by taking more complicated splittings. In all cases, the first step is to find a subspace $U \subset W$ on which the determinant is zero. It is not hard to see that without loss of generality, one can restrict to $U$ that are unextendable, i.e., there does not exist any $U^{\prime} \supset U$ with $\left.\operatorname{det}_{n}\right|_{U^{\prime}}=0$. For results on such subspaces, see, e.g., [IL99, Atk83, EH88, dSP16, FLR85]. Unfortunately they are little understood in general. The first interesting such example, when $n$ is odd, is the space of skew-symmetric matrices.

When $n=3$, the unextendable subspaces have been classified by Atkinson [Atk83]: There are four such up to $G L_{3} \times G L_{3}$-action, namely

$$
\left(\begin{array}{ccc}
* & * & * \\
* & * & * \\
0 & 0 & 0
\end{array}\right),\left(\begin{array}{lll}
* & * & 0 \\
* & * & 0 \\
* & * & 0
\end{array}\right),\left(\begin{array}{ccc}
* & * & * \\
* & 0 & 0 \\
* & 0 & 0
\end{array}\right),\left\{\left.\left(\begin{array}{ccc}
0 & \alpha & \beta \\
-\alpha & 0 & \gamma \\
-\beta & -\gamma & 0
\end{array}\right) \right\rvert\, \alpha, \beta, \gamma \in \mathbb{C}\right\} .
$$

Another way to study the boundary is to consider the rational map

$$
\begin{align*}
\psi: \mathbb{P}\left(\operatorname{End}\left(\mathbb{C}^{n^{2}}\right)\right) & \rightarrow \operatorname{Det}_{n}  \tag{6.7.1}\\
{[X] } & \mapsto\left[\operatorname{det}_{n} \circ X\right]
\end{align*}
$$

One could hope to understand the components of the boundary by blowing up the indeterminacy locus, which consists of $X \in \operatorname{End}\left(\mathbb{C}^{n^{2}}\right)$ such that $\left.\operatorname{det}_{n}\right|_{\text {Image }(X)}=0$.
6.7.2. A component via the skew-symmetric matrices. The transposition $\tau \in G_{\text {det }_{n}}$ allows us to write $\mathbb{C}^{n^{2}}=E \otimes E=S^{2} E \oplus \Lambda^{2} E$, where the decomposition is into the $\pm 1$ eigenspaces for $\tau$. For $M \in E \otimes E$, write $M=M_{S}+M_{\Lambda}$ reflecting this decomposition.

Define a polynomial $P_{\Lambda} \in S^{n}\left(\mathbb{C}^{n^{2}}\right)^{*}$ by

$$
P_{\Lambda}(M)=\overline{\operatorname{det}}_{n}\left(M_{\Lambda}, \ldots, M_{\Lambda}, M_{S}\right) .
$$

Let $\mathrm{Pf}_{i}\left(M_{\Lambda}\right)$ denote the Pfaffian (see, e.g., [Lan12, §2.7.4] for the definition of the Pfaffian and a discussion of its properties) of the skew-symmetric matrix, obtained from $M_{\Lambda}$ by suppressing its $i$-th row and column. Write $M_{S}=\left(s_{i j}\right)$.
Exercise 6.7.2.1: (2) Show that

$$
P_{\Lambda}(M)=\sum_{i, j} s_{i j} \operatorname{Pf}_{i}\left(M_{\Lambda}\right) \operatorname{Pf}_{j}\left(M_{\Lambda}\right)
$$

In particular, $P_{\Lambda}=0$ if $n$ is even but is not identically zero when $n$ is odd.
Proposition 6.7.2.2. [LMR13] $P_{\Lambda} \in \mathcal{D e t}_{n}$. Moreover, $\overline{G L(W) \cdot P_{\Lambda}}$ is an irreducible codimension one component of the boundary of $\mathcal{D e} t_{n}$, not contained in $\operatorname{End}(W) \cdot\left[\operatorname{det}_{n}\right]$. In particular $\overline{d c}\left(P_{\Lambda}\right)=n<d c\left(P_{\Lambda}\right)$.

The proof of Proposition 6.7.2.2 is given in §8.5.1.
Exercise 6.7.2.3: (3) Show that

$$
\operatorname{Zeros}\left(P_{\Lambda}\right)^{\vee}=\overline{\mathbb{P}\left\{v^{2} \oplus v \wedge w \in S^{2} \mathbb{C}^{n} \oplus \Lambda^{2} \mathbb{C}^{n}, v, w \in \mathbb{C}^{n}\right\}} \subset \mathbb{P}^{n^{2}-1}
$$

As expected, $\operatorname{Zeros}\left(P_{\Lambda}\right)^{\vee}$ resembles $\operatorname{Seg}\left(\mathbb{P}^{n-1} \times \mathbb{P}^{n-1}\right)$.
Remark 6.7.2.4. For those familiar with the notation, $\operatorname{Zeros}\left(P_{\Lambda}\right)$ can be defined as the image of the projective bundle $\pi: \mathbb{P}(\mathcal{E}) \rightarrow \mathbb{P}^{n-1}$, where $\mathcal{E}=\mathcal{O}(-1) \oplus \mathcal{Q}$ is the sum of the tautological and quotient bundles on $\mathbb{P}^{n-1}$, by a sub-linear system of $\mathcal{O}_{E}(1) \otimes \pi^{*} \mathcal{O}(1)$. This sub-linear system contracts the divisor $\mathbb{P}(\mathcal{Q}) \subset \mathbb{P}(\mathcal{E})$ to the Grassmannian $G(2, n) \subset \mathbb{P} \Lambda^{2} \mathbb{C}^{n}$.

For large $n$ I expect there are many components of the boundary, however, for $n=3$, we have:
Theorem 6.7.2.5. [HL16] The boundary $\partial \mathcal{D}$ et $t_{3}$ has exactly two irreducible components: $\overline{G L_{9} \cdot P_{\Lambda}}$ and $\overline{\left.G L_{9} \cdot \operatorname{det}_{3}\right|_{\left(E^{*} \otimes E\right)_{0}}}$.

The proof has two parts: first they resolve (6.7.1), which can be done with one blow-up (so in terms of a limit above, only $\frac{1}{t}$ need show up). They then analyze each component of Atkinson's classification and identify the component of the boundary it lies in.
6.7.3. Mulmuley's conjectures on the wildness of the boundary. There is scant evidence for or against the conjecture of [Mul] mentioned above. In $\S 6.8 .1$ I outline the proof that all $P \in S^{3} \mathbb{C}^{3}$ with smooth zero set have $\operatorname{dc}(P)=3$ and thus for all $Q \in S^{3} \mathbb{C}^{3}, \overline{\operatorname{dc}(Q)}=3$. In this one case, there is a big jump between $\overline{\mathrm{dc}}$ and dc, giving some positive news for the conjecture:
Theorem 6.7.3.1. [ABV15] $\operatorname{dc}\left(x_{1}^{3}+x_{2}^{2} x_{3}+x_{2} x_{4}^{2}\right) \geq 6$, and thus when $n=3, \operatorname{dc}\left(P_{\Lambda}\right) \geq 6$.

The second assertion follows because $x_{1}^{3}+x_{2}^{2} x_{3}+x_{2} x_{4}^{2}$ is the determinant of the degeneration of $P_{\Lambda}$ obtained by taking

$$
M_{\Lambda}=\left(\begin{array}{ccc}
0 & x_{4} & x_{2} \\
-x_{3} & 0 & x_{1} \\
-x_{3} & -x_{1} & 0
\end{array}\right), \quad M_{S}=\left(\begin{array}{ccc}
x_{1} & 0 & 0 \\
0 & x_{4} & 0 \\
0 & 0 & x_{2}
\end{array}\right)
$$

Exercise 6.7.3.2: (1) Using Theorem 6.3.4.6, prove the first assertion of Theorem 6.7.3.1.

### 6.8. Determinantal hypersurfaces

This section uses more results from algebraic geometry that we have not discussed. It is not used elsewhere and can be safely skipped.
6.8.1. All smooth cubic surfaces in $\mathbb{P}^{3}$ are determinantal. Grassmann [Gra55] showed that all smooth cubic surfaces in $\mathbb{P}^{3}$ lie in $\operatorname{End}\left(\mathbb{C}^{9}\right)$. $\operatorname{det}_{3}$, and thus all cubic surfaces in $\mathbb{P}^{3}$ lie in $\mathcal{D e t}_{3}$. I give an outline of the proof from $\left[\mathbf{B K s 0 7}\right.$, Ger89]. Every smooth cubic surface $S \subset \mathbb{P}^{3}$ arises in the following way. Consider $\mathbb{P}^{2}$ and distinguish 6 points not on a conic and with no three colinear. There is a four dimensional space of cubic polynomials, say spanned by $F_{1}, \ldots, F_{4} \in S^{3} \mathbb{C}^{3}$, that vanish on the six points. Consider the rational map $\mathbb{P}^{2} \rightarrow \mathbb{P}^{3}$ defined by these polynomials, i.e. $[y] \mapsto\left[F_{1}(y), \ldots, F_{4}(y)\right]$, where the map is defined on $\mathbb{P}^{2}$ minus the six points and let $S$ denote the closure of the image. (Better, one blows up $\mathbb{P}^{2}$ at the six points to obtain a surface $\tilde{S}$ and $S$ is the image of the corresponding regular map from $\tilde{S}$.) Give $\mathbb{C}^{3}$ coordinates $x_{1}, x_{2}, x_{3}$. By the Hilbert-Burch Theorem (see, e.g., [Eis05, Thm. 3.2]), there exists a $3 \times 4$ matrix $L\left(x_{1}, x_{2}, x_{3}\right)$, linear in $x_{1}, x_{2}, x_{3}$, whose size three minors are the $F_{j}$. Define a $3 \times 3$ matrix $M=M\left(z_{1}, \ldots, z_{4}\right)$ by

$$
M\left(\begin{array}{l}
x_{1} \\
x_{2} \\
x_{3}
\end{array}\right)=L\left(\begin{array}{l}
z_{1} \\
z_{2} \\
z_{3} \\
z_{4}
\end{array}\right)
$$

Then $\operatorname{det}(M)$ is the equation of $S$.

Remark 6.8.1.1. The set of non-equivalent representations of a cubic as a determinant is in one-to-one correspondence with the subsets of 6 (of the 27) lines of $S$ that do not intersect each other, see $[\mathbf{B K s 0 7}]$. In particular there are 72 such representations.
6.8.2. Description of the quartic hypersurfaces in $\mathbb{P}^{3}$ that are determinantal. Classically, there was interest in determining which smooth hypersurfaces of degree $d$ were expressible as a $d \times d$ determinant. The result in the first nontrivial case shows how daunting GCT might be.
Theorem 6.8.2.1 (Letao Zhang and Zhiyuan Li, personal communication). The variety $\mathbb{P}\left\{P \in S^{4} \mathbb{C}^{4} \mid[P] \in \operatorname{Det}_{4}\right\} \subset \mathbb{P} S^{4} \mathbb{C}^{4}$ is a hypersurface of degree 640, 224.

The rest of this subsection uses more advanced language from algebraic geometry and can be safely skipped.

The following "folklore" theorem was made explicit in [Bea00, Cor. 1.12]:

Theorem 6.8.2.2. Let $U=\mathbb{C}^{n+1}$, let $P \in S^{d} U$, and let $Z=\operatorname{Zeros}(P) \subset \mathbb{P}^{n}$ be the corresponding hypersurface of degree $d$. Assume $Z$ is smooth and choose any inclusion $U \subset \mathbb{C}^{d^{2}}$.

If $P \in \operatorname{End}\left(\mathbb{C}^{d^{2}}\right) \cdot\left[\operatorname{det}_{d}\right]$, we may form a map between vector bundles $M: \mathcal{O}_{\mathbb{P}^{n}}(-1)^{d} \rightarrow \mathcal{O}_{\mathbb{P}^{n}}^{d}$ whose cokernel is a line bundle $L \rightarrow Z$ with the properties:
i) $H^{i}(Z, L(j))=0$ for $1 \leq i \leq n-2$ and all $j \in \mathbb{Z}$
ii) $H^{0}(X, L(-1))=H^{n-1}(X, L(j))=0$

Conversely, if there exists $L \rightarrow Z$ satisfying properties i) and ii), then $Z$ is determinantal via a map $M$ as above whose cokernel is $L$.

If we are concerned with the hypersurface being in $\mathcal{D} e t_{n}$, the first case where this is not automatic is for quartic surfaces, where it is a codimension one condition:
Proposition 6.8.2.3. [Bea00, Cor. 6.6] A smooth quartic surface is determinantal if and only if it contains a nonhyperelliptic curve of genus 3 embedded in $\mathbb{P}^{3}$ by a linear system of degree 6 .

Proof of 6.8.2.1. From Proposition 6.8.2.3, the hypersurface is the locus of quartic surfaces containing a (Brill-Noether general) genus 3 curve $C$ of degree six. This translates into the existence of a lattice polarization

|  | $h$ | $C$ |
| :--- | :--- | :--- |
| $h$ | 4 | 6 |
| $C$ | 6 | 4 |

of discriminant $-\left(4^{2}-6^{2}\right)=20$. By the Torelli theorems, the $K 3$ surfaces with such a lattice polarization have codimension one in the moduli space of quartic $K 3$ surfaces.

Let $D_{3,6}$ denote the locus of quartic surfaces containing a genus 3 curve $C$ of degree six in $\mathbb{P}^{34}=\mathbb{P}\left(S^{4} \mathbb{C}^{4}\right)$. It corresponds to the Noether-Lefschetz divisor $N L_{20}$ in the moduli space of the degree four $K 3$ surfaces. Here $N L_{d}$ denotes the Noether-Lefschetz divisor, parameterizing the degree $4 K 3$ surfaces whose Picard lattice has a rank 2 sub-lattice containing $h$ with discriminant $-d$. (h is the polarization of the degree four $K 3$ surface, $h^{2}=$ 4.)

The Noether-Lefschetz number $n_{20}$, which is defined by the intersection number of $N L_{20}$ and a line in the moduli space of degree four $K 3$ surfaces, equals the degree of $D_{3,6}$ in $\mathbb{P}^{34}=\mathbb{P}\left(S^{4} \mathbb{C}^{4}\right)$.

The key fact is that $n_{d}$ can be computed via the modularity of the generating series for any integer $d$. More precisely, the generating series $F(q):=\sum_{d} n_{d} q^{d / 8}$ is a modular form of level 8 , and can be expressed by a polynomial of $A(q)=\sum_{n} q^{n^{2} / 8}$ and $B(q)=\sum_{n}(-1)^{n} q^{n^{2} / 8}$.

The explicit expression of $F(q)$ is in [MP, Thm 2]. As an application, the Noether-Lefschetz number $n_{20}$ is the coefficient of the term $q^{20 / 8}=q^{5 / 2}$, which is 640,224 .

## Chapter 7

## Valiant's hypothesis II: Restricted models and other approaches

This chapter continues the discussion of Valiant's hypothesis and its variants. Chapter 6 described progress via benchmarks such as lower bounds for $\mathrm{dc}\left(\operatorname{perm}_{m}\right)$. Another approach to these problems is to prove complexity lower bounds under supplementary hypotheses, called restricted models in the computer science literature. I begin, in §7.1, with a discussion of the geometry of one of the simplest classes of shallow circuits, the $\Sigma \Lambda \Sigma$-circuits whose complexity essentially measures symmetric tensor rank, and discuss the symmetric tensor rank of the elementary symmetric polynomials. Next, in $\S 7.2$, I discuss $\Sigma \Pi \Sigma$ circuits and their relationship to secant varieties of the Chow variety. There are several complexity measures that are equivalent to determinantal complexity, such as algebraic branching programs and iterated matrix multiplication complexity. These are discussed in §7.3. Additional restricted models are presented in §7.4: Aravind and Joegelkar's rank $k$ determinantal expressions of [AJ15], Shpilka's restricted model [Shp02] of depth-2 symmetric arithmetic circuits, a result of Glynn [Gly13] on a certain class of expressions for the permanent, Nisan's non-commutative ABP's [Nis91], and the equivariant determinantal complexity of [LR15]. Equivariant determinantal complexity is the only known restricted model that gives an exponential separation between the permanent and determinant.

I devote $\S 7.5$ to the restricted models of shallow circuits because there is a path to proving Valiant's hypothesis by proving lower bounds that are stronger than super-polynomial for them. The depth of a circuit $\mathcal{C}$ is the
number of edges in the longest path in $\mathcal{C}$ from an input to its output. If a circuit has small depth, it is called a shallow circuit, and the polynomial it computes can be computed quickly in parallel. The section begins in §7.5.1 with a detour for readers not familiar with big numbers as different levels of super-polynomial growth need to be compared both for statements and proofs. Having already discussed the geometry associated to depth 3 circuits in $\S 7.2$, I explain the geometry associated to the depth 4 and 5 circuits that arise in [GKKS13a] in §7.5.3. I discuss the tantalizing lower bounds of [GKKS13a] in $\S 7.6$, and analyze the method of proof, shifted partial derivatives, in detail. I then show that this method cannot separate the padded permanent from the determinant.

I conclude with a brief discussion of polynomial identity testing (PIT), hitting sets, and effective Noether normalization in §7.7. I believe these topics are potentially of great interest to algebraic geometry.

As pointed out by Shpilka and Yehudayoff in [SY09], restricted circuits of polynomial size only compute polynomials with "simple" structure. Thus to understand them one needs to determine the precise meaning of "simple" for a given restricted class, and then find an "explicit" polynomial without such structure. One could rephrase this geometrically as restricted circuits of a fixed size $s$ define an algebraic variety in $S^{n} \mathbb{C}^{N}$ that is the closure of the set of polynomials computable with a restricted circuit of size $s$. The goal becomes to find an equation of that variety and an explicit polynomial not satisfying that equation.

Recall that computer scientists always work in bases and the inputs to the circuits are constants and variables. For homogeneous circuits, the inputs are simply the variables. The first layer of a $\Sigma \Lambda \Sigma, \Sigma \Pi \Sigma$, or $\Sigma \Lambda \Sigma \Lambda \Sigma$ circuit for a polynomial $P \in S^{d} \mathbb{C}^{N}$ is just to obtain arbitrary linear forms from these variables, so it plays no role in the geometry, and at worst multiplies the circuit size by $N$, and often enlarges it by much less. This fact will be used throughout this chapter.

I continue to work exclusively with homogeneous polynomials and over the complex numbers. In particular, for a $\mathbf{v}$-dimensional complex vector space $V, S^{d} V$ denotes the space of homogeneous polynomials of degree $d$ on $V^{*}$.

### 7.1. Waring rank, depth three powering circuits and symmetric polynomials

Recall from §6.2.2 that the symmetric tensor rank (Waring rank) of a polynomial $P \in S^{d} V$, denoted $\mathbf{R}_{S}(P)$ is the smallest $r$ such that we may write $P=\ell_{1}^{d}+\cdots+\ell_{r}^{d}$ for some $\ell_{j} \in V$. As explained in §7.1.1, such $P$ admit $\Sigma \Lambda \Sigma$ circuits of size at most $r(\mathbf{v}+2)$. Although not directly related to Valiant's
hypothesis, they are a simple enough class of circuits that one can actually prove lower bounds and they are used as the basis for further lower bound results.

Similarly, the class of elementary symmetric polynomials is a class of polynomials simple enough for one to prove complexity results, but rich enough to be of interest. In $\S 7.1 .2$ I discuss the elementary symmetric function $e_{n, n}=x_{1} \cdots x_{n}$, describing its symmetry group and Waring decomposition. In §7.1.3 I discuss the Waring decompositions of elementary symmetric polynomials in general.

Recall the notation $\sigma_{r}\left(v_{d}(\mathbb{P} V)\right)=\overline{\mathbb{P}\left\{P \in S^{d} V \mid P=\ell_{1}^{d}+\cdots \ell_{r}^{d}, \ell_{j} \in V\right\}}$ for the Zariski closure of the set of polynomials in $\mathbb{P} S^{d} V$ of Waring rank at most $r$, called the $r$-th secant variety of the Veronese variety, and that $\underline{\mathbf{R}}_{S}(P)$ denotes the smallest $r$ such that $P \in \sigma_{r}\left(v_{d}(\mathbb{P} V)\right)$.
7.1.1. $\sigma_{r}\left(v_{d}(\mathbb{P} V)\right)$ and $\Sigma \Lambda \Sigma$ circuits. When one studies circuits of bounded depth, one must allow gates to have an arbitrary number of edges coming in to them, which is called unbounded fanin. For such circuits, multiplication by constants is considered free.

A $\Sigma \Lambda^{\delta} \Sigma$ circuit consists of three layers, the first of addition gates, the second of powering gates, that map $\ell \mapsto \ell^{\delta}$ (so each gate has a single input and output), and the third a single addition gate. Such circuits are also called diagonal depth-3 circuits, or depth three powering circuits, see, e.g., [Sax08].
Proposition 7.1.1.1. Say $P \in S^{d} \mathbb{C}^{\mathbf{v}}$ satisfies $\mathbf{R}_{S}(P)=r$. Then $P$ admits a $\Sigma \Lambda^{d} \Sigma$ circuit of size $r(\mathbf{v}+2)$.

Proof. We are given that $P=\ell_{1}^{d}+\cdots+\ell_{r}^{d}$ from some $\ell_{j} \in \mathbb{C}^{\mathbf{v}}$. We need at most $\mathbf{v}$ additions to construct each $\ell_{j}$, of which there are $r$, so $r \mathbf{v}$ edges at the first level. Then there are $r$ powering gates, of one edge each and each of these sends one edge to the final addition gate, for a total of $r \mathbf{v}+r+r$.

The following proposition bounds $\Sigma \Lambda \Sigma$ complexity by $\overline{\mathrm{dc}}$ :
Proposition 7.1.1.2. Let $P \in S^{m} V$ and let $\ell \in V$. Then $\operatorname{dc}(P) \leq$ $m \mathbf{R}_{S}(P)+1$ and $\overline{\mathrm{dc}}(P) \leq m \underline{\mathbf{R}}_{S}(P)+1$.
Exercise 7.1.1.3: (2) Prove Proposition 7.1.1.2. ©
7.1.2. The polynomial $x_{1} \cdots x_{n}$. Consider the polynomial $e_{n, n}:=x_{1} \cdots x_{n} \in$ $S^{n} \mathbb{C}^{n}$ (the $n$-th elementary symmetric function in $n$ variables). This simple polynomial plays a major role in complexity theory and geometry. Its $G L_{n^{-}}$ orbit closure has been studied for over a hundred years and is discussed in Chapter 9. In some sense it is the "weakest" polynomial known that requires an exponential size $\Sigma \Lambda \Sigma$-circuit, which will be important in $\S 7.7$. I
first determine its symmetry group $G_{e_{n, n}}$, which will be used several times in what follows.

It is clear $T_{n}^{S L} \rtimes \mathfrak{S}_{n} \subseteq G_{e_{n, n}}$, where $T_{n}^{S L}$ denotes the diagonal matrices with determinant one (the matrix with $\left(\lambda_{1}, \ldots, \lambda_{n}\right)$ on the diagonal sends $x_{j}$ to $\lambda_{j} x_{j}$ ) and $\mathfrak{S}_{n}$ acts by permuting the basis vectors. We need to determine if the stabilizer is larger. Let $g \in G L_{n}$. Then

$$
g \cdot e_{n, n}=\left(\sum_{j_{1}=1}^{n} g_{1}^{j_{1}} x_{j_{1}}\right) \ldots\left(\sum_{j_{n}=1}^{n} g_{n}^{j_{n}} x_{j_{n}}\right) .
$$

In order that this be equal to $x_{1} \cdots x_{n}$, by unique factorization of polynomials, there must be a permutation $\sigma \in \mathfrak{S}_{n}$ such that for each $k$, we have $\sum_{j} g_{k}^{j} x_{j}=\lambda_{k} x_{\sigma(k)}$ for some $\lambda_{k} \in \mathbb{C}^{*}$. Composing with the inverse of this permutation we have $g_{k}^{j}=\delta_{k}^{j} \lambda_{j}$, and finally we see that we must further have $\lambda_{1} \cdots \lambda_{n}=1$, which means it is an element of $T_{n}^{S L}$, so the original $g$ is an element of $T_{n}^{S L} \rtimes \mathfrak{S}_{n}$. Thus $G_{e_{n, n}}=T_{n}^{S L} \rtimes \mathfrak{S}_{n}$. By the discussion in §4.2, any Waring decomposition of $e_{n, n}$ containing a pinning set can have symmetry group at most $\mathfrak{S}_{n}$.

The optimal Waring decomposition of $x_{1} \cdots x_{n}$ is

$$
\begin{equation*}
x_{1} \cdots x_{n}=\frac{1}{2^{n-1} n!} \sum_{\substack{\epsilon \in\{-1,1\}^{n} \\ \epsilon 1=1}}\left[\left(\Pi_{i=1}^{n} \epsilon_{i} \sum_{j=1}^{n} \epsilon_{j} x_{j}\right)^{n}\right] \tag{7.1.1}
\end{equation*}
$$

a sum with $2^{n-1}$ terms. It is called Fischer's formula in the computer science literature because Fischer wrote it down in 1994 [Fis94]. While similar formulas appeared earlier (e.g. formula (7.1.2) below appeared in 1934), I have not found this precise formula earlier in the literature. I give the proof of its optimality (due to Ranestad and Schreyer [RS11]) in §10.1.2.

This decomposition transparently has an $\mathfrak{S}_{n-1}$-symmetry. Here is a slightly larger expression that transparently has an $\mathfrak{S}_{n}$-symmetry:

$$
\begin{equation*}
x_{1} \cdots x_{n}=\frac{1}{2^{n} n!} \sum_{\epsilon \in\{-1,1\}^{n}}\left[\Pi_{i=1}^{n} \epsilon_{i}\left(\sum_{j=1}^{n} \epsilon_{j} x_{j}\right)^{n}\right] \tag{7.1.2}
\end{equation*}
$$

This formula dates back at least to 1934, where Mazur and Orlicz [MO34] gave it and generalizations.

Remarkably, as was realized by H. Lee [Lee16], Fischer's expression already has an $\mathfrak{S}_{n}$-symmetry when $n$ is odd.

For example:

$$
x y z=\frac{1}{24}\left[(x+y+z)^{3}-(x+y-z)^{3}-(x-y+z)^{3}-(-x+y+z)^{3}\right] .
$$

For an integer set $I$ and an integer $i$, define

$$
\delta(I, i)=\left\{\begin{array}{cc}
-1 & i \in I \\
1 & i \notin I
\end{array}\right.
$$

When $n=2 k+1$ is odd, rewrite Fischer's formula as:
$x_{1} x_{2} \cdots x_{n}=\frac{1}{2^{n-1} n!} \sum_{I \subset[n],|I| \leq k}(-1)^{|I|}\left(\delta(I, 1) x_{1}+\delta(I, 2) x_{2}+\cdots+\delta(I, n) x_{n}\right)^{n}$.
When $n=2 k$ is even, the situation is a little more subtle. One may rewrite Fischer's formula as:

$$
\begin{align*}
x_{1} x_{2} \cdots x_{n}=\frac{1}{2^{n-1} n!} & {\left[\sum_{I \subset[n],|I|<k}(-1)^{|I|}\left(\delta(I, 1) x_{1}+\delta(I, 2) x_{2}+\cdots+\delta(I, n) x_{n}\right)^{n}\right.}  \tag{7.1.4}\\
& \left.+\sum_{I \subset[n],|I|=k, 1 \in I} \frac{(-1)^{k}}{2}\left(\delta(I, 1) x_{1}+\delta(I, 2) x_{2}+\cdots+\delta(I, n) x_{n}\right)^{n}\right] .
\end{align*}
$$

The collection of terms in the second summation is only $\mathfrak{S}_{n}$-invariant up to sign. In the language of Chapter 4 , if we write the decomposition as $\mathcal{S}=\left\{\ell_{1}^{n}, \ldots, \ell_{2^{n-1}}^{n}\right\}$, the decomposition $\mathcal{S}$ is $\mathfrak{S}_{n}$-invariant. Moreover, the set $\left\{\left[\ell_{1}\right], \ldots,\left[\ell_{2^{n-1}}\right]\right\}$ is $\mathfrak{S}_{n}$-invariant, however the set $\left\{\ell_{1}, \ldots, \ell_{2^{n-1}}\right\}$ is only $\mathfrak{S}_{n-1}$-invariant (and $\mathfrak{S}_{n}$-invariant up to sign).

Remark 7.1.2.1. Using the techniques of [RS00], Ranestad (personal communication) has shown that every minimal rank decomposition of $x_{1} \cdots x_{n}$ is in the $T^{S L_{n}}$-orbit of the the right hand side of (7.1.1), so in particular, by Proposition 4.1.2.2 every decomposition has $\mathfrak{S}_{n}$-symmetry.
7.1.3. Symmetric ranks of elementary symmetric polynomials. Here are generalizations of the Waring expressions for $e_{n, n}$ to all symmetric polynomials due to H. Lee:
Theorem 7.1.3.1. [Lee16] Let $d=2 k+1$ and let $N \geq d$. Then
$e_{d, N}=\frac{1}{2^{d-1} d!} \sum_{I \subset[N],|I| \leq k}(-1)^{|I|}\binom{N-k-|I|-1}{k-|I|}\left(\delta(I, 1) x_{1}+\delta(I, 2) x_{2}+\cdots+\delta(I, N) x_{N}\right)^{d}$.
In particular, for $d$ odd, $\mathbf{R}_{S}\left(e_{d, N}\right) \leq \sum_{i=0}^{\left\lfloor\frac{d}{2}\right\rfloor}\binom{N}{i}$.
This formula nearly appeared in [MO34] in 1934, but just as with Fischer's, there was a doubling of size.

Proof. Work by downwards induction, the case $d=N$ is Fischer's formula. Let $d<N$ and let $F_{d, N}$ denote the right hand side of the expression.

Observe that $F_{d, d}=e_{d, d}$ and $F_{d, N-1}=F_{d, N}\left(x_{1}, \ldots, x_{N-1}, 0\right)$ up to a constant. In particular $F_{d, d}=F_{d, N}\left(x_{1}, \ldots, x_{d}, 0, \ldots, 0\right)$ up to a constant. The analogous statement holds setting any subset of the variables to zero. This implies that $F_{d, N}$ is an expression that has all the square-free monomials in $e_{d, N}$ appearing in it, all with the same coefficient. Moreover, there are no other monomials appearing in $F_{d, N}$ as otherwise there would be a monomial involving fewer than $d$ variables that would appear in some specialization to some $e_{d, d}$. One concludes by checking the constant is correct.

Lee gives a similar formula for $d$ even.

### 7.2. Depth three circuits and secant varieties of the Chow variety

In this section I discuss the depth three or $\Sigma \Pi \Sigma$ circuits, which consist of depth three formulas where the first layer of gates consist of additions, the second of multiplications, and the last gate is an addition gate. Remarkably, these circuits are powerful enough to potentially separate VP from VNP, as is explained in §7.5.

There is a subtlety with these circuits: their homogeneous version, used naïvely, lacks computing power. This can be fixed either by allowing inhomogeneous circuits, which is what is done in the computer science literature, or with the help of padding, which I discuss in §7.2.3.
7.2.1. Secant varieties and homogeneous depth three circuits. Recall the Chow variety $C h_{n}(W) \subset \mathbb{P} S^{n} W$. When $\mathbf{w}=\operatorname{dim} W \geq n$, it is the orbit closure $\overline{G L(W) \cdot\left[x_{1} \cdots x_{n}\right]}$. The set of polynomials of the form $\sum_{i=1}^{r} \ell_{i, 1} \cdots \ell_{i, n}$, where $\ell_{i, j} \in W$ (the sum-product polynomial in the computer science literature) is denoted $\sigma_{r}^{0}\left(C h_{n}(W)\right)$, and $\sigma_{r}\left(C h_{n}(W)\right)$ is the Zariski closure in $\mathbb{P} S^{n} W$ of $\sigma_{r}^{0}\left(C h_{n}(W)\right)$, the $r$-th secant variety of the Chow variety.

The relation between secant varieties of Chow varieties and depth three circuits is as follows:
Proposition 7.2.1.1. [Lan15a] A polynomial $P \in S^{n} W$ in $\sigma_{r}^{0}\left(C h_{n}(W)\right)$ is computable by a homogeneous depth three circuit of size $r+n r(1+\mathbf{w})$. If $P \notin \sigma_{r}^{0}\left(C h_{n}(W)\right)$, then $P$ cannot be computed by a homogeneous depth three circuit of size $n(r+1)+r+1$.

Proof. In the first case, $P=\sum_{j=1}^{r}\left(\ell_{1 j} \cdots \ell_{n j}\right)$ for some $\ell_{s j} \in W$. Expressed in terms of a fixed basis of $W$, each $\ell_{s j}$ is a linear combination of at worst $\mathbf{w}$ basis vectors, thus to create the $\ell_{s j}$ requires at worst $n r \mathbf{w}$ additions. Then to multiply them in groups of $n$ is $n r$ multiplications, and finally to add these together is $r$ further additions. In the second case, at best $P$ is in
$\sigma_{r+1}^{0}\left(C h_{n}(W)\right)$, in which case, even if each of the $\ell_{s j}$ 's is a basis vector (so no initial additions are needed), we still must perform $n(r+1)$ multiplications and $r+1$ additions.
7.2.2. Why homogeneous depth three circuits do not appear useful at first glance. Exercise 6.2.2.7 implies that in order that [ $\operatorname{det}_{n}$ ] $\in$ $\sigma_{r}\left(C h_{n}\left(\mathbb{C}^{n^{2}}\right)\right)$ we must have $r\binom{n}{\left\lfloor\frac{n}{2}\right\rfloor} \geq\binom{ n}{\left\lfloor\frac{n}{2}\right\rfloor}^{2}$, i.e., $r>\binom{n}{\left\lfloor\frac{n}{2}\right\rfloor} \sim 2^{n} / n$ (see §7.5.1).

By Proposition 7.2.1.1, we conclude:
Proposition 7.2.2.1. [NW97] The polynomial sequences $\operatorname{det}_{n}$ and perm $_{n}$ do not admit homogeneous depth three circuits of size $2^{n} / n$.

Remark 7.2.2.2. The proof above follows from considering partial derivatives in middle degree. In [NW97] they consider all partial derivatives of all orders simultaneously to improve the lower bound to $2^{n}$.

Thus homogeneous depth three circuits at first sight do not seem that powerful because a homogeneous depth 3 circuit of size $2^{n}$ cannot compute the determinant.

To make matters worse, consider the polynomial corresponding to iterated matrix multiplication of three by three matrices $I M M_{k}^{3} \in S^{k}\left(\mathbb{C}^{9 k}\right)$. It is complete for the class $\mathbf{V P}_{e}$ of sequences with polynomial sized formulas discussed in Remark 6.1.5.2 (see [BOC92]), and also has an exponential lower bound for its Chow border rank:
Exercise 7.2.2.3: (2) Use flattenings to show $I M M_{k}^{3} \notin \sigma_{p o l y(k)}\left(C h_{k}\left(\mathbb{C}^{9 k}\right)\right)$.
By Exercise 7.2.2.3, sequences of polynomials admitting polynomial size formulas do not in general have polynomial size homogeneous depth three circuits.
7.2.3. Homogeneous depth three circuits for padded polynomials. If one works with padded polynomials instead of polynomials (as we did with $\mathcal{D} e t_{n}$ ), the power of homogeneous depth three circuits increases to the power of arbitrary depth three circuits. The following geometric version of a result of Ben-Or and Cleve (presented below as a Corollary) was suggested by K. Efremenko:
Proposition 7.2.3.1. [Lan15a] Let $\mathbb{C}^{m+1}$ have coordinates $\ell, x_{1}, \ldots, x_{m}$ and let $e_{k, m}=e_{k, m}\left(x_{1}, \ldots, x_{m}\right)$ be the $k$-th elementary symmetric polynomial. For all $k \leq m, \ell^{m-k} e_{k, m} \in \sigma_{m}^{0}\left(C h_{m}\left(\mathbb{C}^{m+1}\right)\right)$.

Proof. Fix an integer $u \in \mathbb{Z}$, recall the generating function $E_{m}$ for the elementary symmetric functions from (6.1.2), and define

$$
\begin{aligned}
g_{u}(x, \ell) & =(u \ell)^{m} E_{m}\left(\frac{1}{u \ell}\right) \\
& =\prod_{i=1}^{m}\left(x_{i}+u \ell\right) \\
& =\sum_{k} u^{m-k} e_{k, m}(x) \ell^{m-k} .
\end{aligned}
$$

The second line shows $g_{u}(x, \ell) \in C h_{m}\left(\mathbb{C}^{m+1}\right)$. Letting $u=1, \ldots, m$, we may use the inverse of the Vandermonde matrix to write each $\ell^{m-k} e_{k, m}$ as a sum of $m$ points in $C h_{m}\left(\mathbb{C}^{m+1}\right)$ because

$$
\left(\begin{array}{cccc}
1^{0} & 1^{1} & \cdots & 1^{m} \\
2^{0} & 2^{1} & \cdots & 2^{m} \\
& \vdots & & \\
m^{0} & m^{1} & \cdots & m^{m}
\end{array}\right)\left(\begin{array}{c}
\ell^{m-1} e_{1, m} \\
\ell^{m-2} e_{2, m} \\
\vdots \\
\ell^{0} e_{m, m}
\end{array}\right)=\left(\begin{array}{c}
g_{1}(x, \ell) \\
g_{2}(x, \ell) \\
\vdots \\
g_{m}(x, \ell)
\end{array}\right)
$$

Corollary 7.2.3.2. $[\mathbf{B O C} 92] \ell^{m-k} e_{k, m}$ can be computed by a homogeneous depth three circuit of size $3 m^{2}+m$.

Proof. As remarked above, for any point of $\sigma_{r}\left(C h_{n}\left(\mathbb{C}^{m+1}\right)\right)$ one gets a circuit of size at most $r+n r+r n(m+1)$, but here at the first level all the addition gates have fanin two (i.e., there are two inputs to each addition gate) instead of the possible $m+1$.
Remark 7.2.3.3. The best lower bound for computing the $e_{n}^{k}$ via a $\Sigma \Pi \Sigma$ circuit is $\Omega\left(n^{2}\right)$ [SW01], so Corollary 7.2.3.2 is very close to (and may well be) sharp.

Proposition 7.2.3.4. [Lan15a] Say $P \in S^{m} \mathbb{C}^{M}$ is computable by a depth three circuit of size $s$. Then, for some $n<s+m, \ell^{n-m} P$ is computable by a homogeneous depth three circuit of size $O\left(s^{2}\right)$.

Proof. Start with the inhomogeneous circuit computing $P$. At the first level, add a homogenizing variable $\ell$, so that the affine linear outputs become linear in our original variables plus $\ell$, the product gates will each produce a homogeneous polynomial. While the different product gates may produce polynomials of different degrees, when we add them up what remains must be a sum of homogeneous polynomials, such that when we set $\ell=1$, we obtain the desired homogeneous polynomial. Say the largest power of $\ell$ appearing in this sum is $q$. Note that $q<s$. For each other term there is some other power of $\ell$ appearing, say $q_{i}$ for the $i$-th term. Then to the
original circuit, add $q-q_{i}$ inputs to the $i$-th product gate, where each input is $\ell$. This will not change the size of the circuit by more than $q r<s^{2}$. Our new homogeneous depth three circuit will output $\ell^{q} P$.

### 7.3. Algebraic branching programs

In this section I describe algebraic branching programs, a model of computation with complexity equivalent to that of the determinant, as well as two restrictions of it, one (non-commutative ABP's) that has an exponential lower bound for the permanent (but also for the determinant), and another (read once ABP's) where it is possible to carry out deterministic polynomial identity testing as described in §7.7.

### 7.3.1. Algebraic branching programs and iterated matrix multiplication.

Definition 7.3.1.1 (Nisan [Nis91]). An Algebraic Branching Program (ABP) over $\mathbb{C}$ is a directed acyclic graph $\Gamma$ with a single source $s$ and a single sink $t$. Each edge $e$ is labeled with an affine linear function $\ell_{e}$ in the variables $\left\{y_{i} \mid 1 \leq i \leq M\right\}$. Every directed path $p=e_{1} e_{2} \cdots e_{k}$ computes the product $\Gamma_{p}:=\prod_{j=1}^{k} \ell_{e_{j}}$. For each vertex $v$ the polynomial $\Gamma_{v}$ is defined as $\sum_{p \in \mathcal{P}_{s, v}} \Gamma_{p}$ where $\mathcal{P}_{s, v}$ is the set of paths from $s$ to $v$. We say that $\Gamma_{v}$ is computed by $\Gamma$ at $v$. We also say that $\Gamma_{t}$ is computed by $\Gamma$ or that $\Gamma_{t}$ is the output of $\Gamma$.

The size of $\Gamma$ is the number of vertices. Let $\operatorname{abpc}(P)$ denote the smallest size of an algebraic branching program that computes $P$.

An ABP is layered if we can assign a layer $i \in \mathbb{N}$ to each vertex such that for all $i$, all edges from layer $i$ go to layer $i+1$. An ABP is homogeneous if the polynomials computed at each vertex are all homogeneous. A homogeneous ABP $\Gamma$ is degree layered if $\Gamma$ is layered and the layer of a vertex $v$ coincides with the degree of $v$. For a homogeneous $P$ let dlabpc $(P)$ denote the smallest size of a degree layered algebraic branching program that computes $P$. Of course dlabpc $(P) \geq \operatorname{abpc}(P)$.

Definition 7.3.1.2. The iterated matrix multiplication complexity of a polynomial $P(y)$ in $M$ variables, $\operatorname{immc}(P)$ is the smallest $n$ such that there exists affine linear maps $B_{j}: \mathbb{C}^{M} \rightarrow \operatorname{Mat}_{n}(\mathbb{C}), j=1, \ldots, n$, such that $P(y)=\operatorname{trace}\left(B_{n}(y) \cdots B_{1}(y)\right)$. The homogeneous iterated matrix multiplication complexity of a degree $m$ homogeneous polynomial $P \in S^{m} \mathbb{C}^{M}$, $\operatorname{himmc}(P)$, is the smallest $n$ such that there exist natural numbers $n_{1}, \ldots, n_{m}$ with $1=n_{1}$, and $n=n_{1}+\cdots+n_{m}$, and linear maps $A_{s}: \mathbb{C}^{M} \rightarrow$ Mat $_{n_{s} \times n_{s+1}}$, $1 \leq s \leq m$, with $n_{m+1}=1$, such that $P(y)=A_{m}(y) \cdots A_{1}(y)$.
7.3.2. Determinantal complexity and ABP's. Two complexity measures $m_{1}, m_{2}$ are polynomially related if for any sequence $p_{n}$ of polynomials, there exist constants $C_{1}, C_{2}$ such that for all sufficiently large $n$, $m_{1}\left(p_{n}\right) \leq\left(m_{2}\left(p_{n}\right)\right)^{C_{1}}$ and $m_{2}\left(p_{n}\right) \leq\left(m_{1}\left(p_{n}\right)\right)^{C_{2}}$.

The following folklore theorem was stated explicitly in [IL16a] with precise upper and lower bounds between the various complexity measures:
Theorem 7.3.2.1. [IL16a] The complexity measures dc, abpc, immc, dlabpc and himmc, are all polynomially related.

Additional relations between different models are given in [MP08].
Regarding the geometric search for separating equations, the advantage one gains by removing the padding in the iterated matrix multiplication model is offset by the disadvantage of dealing with the himmc polynomial that for all known equations such as Young flattenings (which includes the method of shifted partial derivatives as a special case) and equations for degenerate dual varieties, behaves far more generically than the determinant.

Work of Mahajan-Vinay [MV97] implies:
Proposition 7.3.2.2. [IL16a] dlabpc $\left(\operatorname{det}_{m}\right) \leq \frac{m^{3}}{3}-\frac{m}{3}+2$ and $\operatorname{himmc}\left(\operatorname{det}_{m}\right) \leq$ $\frac{m^{3}}{3}-\frac{m}{3}+2$.

Remark 7.3.2.3. For $m<7$, the size $2^{m}-1$ Grenet-like expressions from [LR15] for $\operatorname{det}_{m}$ give smaller iterated matrix multiplication expressions. This warns us that small cases can be deceptive.

Remark 7.3.2.4. It is an important and perhaps tractable open problem to prove an $\omega\left(m^{2}\right)$ lower bound for dc $\left(\operatorname{perm}_{m}\right)$. By the more precise version of Theorem 7.3.2.1 in [IL16a], it would suffice to prove an $\omega\left(m^{6}\right)$ lower bound for himmc $\left(\operatorname{perm}_{m}\right)$.

Here are the size $\frac{m^{3}}{3}-\frac{m}{3}+2$ himmc expressions for $\operatorname{det}_{m}$ when $m=3,4,5$ :

$$
\operatorname{det}_{3}(x)=\left(x_{1}^{2}, x_{1}^{3}, x_{2}^{2}, x_{2}^{3}, x_{3}^{3}\right)\left(\begin{array}{ccc}
x_{2}^{2} & x_{2}^{3} & 0 \\
x_{3}^{2} & x_{3}^{3} & 0 \\
-x_{1}^{2} & -x_{1}^{3} & 0 \\
0 & 0 & x_{3}^{2} \\
-x_{1}^{2} & -x_{1}^{3} & -x_{2}^{2}
\end{array}\right)\left(\begin{array}{c}
x_{2}^{1} \\
x_{3}^{1} \\
-x_{1}^{1}
\end{array}\right) .
$$

Let $M_{1}=\left(-x_{1}^{2},-x_{1}^{3},-x_{1}^{4},-x_{2}^{2},-x_{2}^{3},-x_{2}^{4},-x_{3}^{3},-x_{3}^{4},-x_{4}^{4}\right)$. Then

$$
\operatorname{det}_{4}(x)=M_{1}\left(\begin{array}{ccccccc}
x_{2}^{2} & x_{2}^{3} & x_{2}^{4} & 0 & 0 & 0 & 0 \\
x_{3}^{2} & x_{3}^{3} & x_{3}^{4} & 0 & 0 & 0 & 0 \\
x_{4}^{2} & x_{4}^{3} & x_{4}^{4} & 0 & 0 & 0 & 0 \\
-x_{1}^{2} & -x_{1}^{3} & -x_{1}^{4} & 0 & 0 & 0 & 0 \\
0 & 0 & 0 & x_{3}^{2} & x_{3}^{3} & x_{3}^{4} & 0 \\
0 & 0 & 0 & x_{4}^{2} & x_{4}^{3} & x_{4}^{4} & 0 \\
-x_{1}^{2} & -x_{1}^{3} & -x_{1}^{4} & -x_{2}^{2} & -x_{2}^{3} & -x_{2}^{4} & 0 \\
0 & 0 & 0 & 0 & 0 & 0 & x_{4}^{3} \\
-x_{1}^{2} & -x_{1}^{3} & -x_{1}^{4} & -x_{2}^{2} & -x_{2}^{3} & -x_{2}^{4} & -x_{3}^{3}
\end{array}\right)\left(\begin{array}{cccc}
x_{2}^{2} & x_{2}^{3} & x_{2}^{4} & 0 \\
x_{3}^{2} & x_{3}^{3} & x_{3}^{4} & 0 \\
x_{4}^{2} & x_{4}^{3} & x_{4}^{4} & 0 \\
-x_{1}^{2} & -x_{1}^{3} & -x_{1}^{4} & 0 \\
0 & 0 & 0 & x_{3}^{2} \\
0 & 0 & 0 & x_{4}^{2} \\
-x_{1}^{2} & -x_{1}^{3} & -x_{1}^{4} & -x_{2}^{2}
\end{array}\right)\left(\begin{array}{c}
x_{2}^{1} \\
x_{3}^{1} \\
x_{4}^{1} \\
-x_{1}^{1}
\end{array}\right)
$$

Let $M_{1}=\left(x_{1}^{2}, x_{1}^{3}, x_{1}^{4}, x_{1}^{5}, x_{2}^{2}, x_{2}^{3}, x_{2}^{4}, x_{2}^{5}, x_{3}^{3}, x_{3}^{4}, x_{3}^{5}, x_{4}^{4}, x_{4}^{5}, x_{5}^{5}\right)$,

$$
M_{2}=\left(\begin{array}{cccccccccccc}
x_{2}^{2} & x_{2}^{3} & x_{2}^{4} & x_{2}^{5} & 0 & 0 & 0 & 0 & 0 & 0 & 0 & 0 \\
x_{3}^{2} & x_{3}^{3} & x_{3}^{4} & x_{3}^{5} & 0 & 0 & 0 & 0 & 0 & 0 & 0 & 0 \\
x_{4}^{2} & x_{4}^{3} & x_{4}^{4} & x_{4}^{5} & 0 & 0 & 0 & 0 & 0 & 0 & 0 & 0 \\
x_{5}^{2} & x_{5}^{3} & x_{5}^{4} & x_{5}^{5} & 0 & 0 & 0 & 0 & 0 & 0 & 0 & 0 \\
-x_{1}^{2} & -x_{1}^{3} & -x_{1}^{4} & -x_{1}^{5} & 0 & 0 & 0 & 0 & 0 & 0 & 0 & 0 \\
0 & 0 & 0 & 0 & x_{3}^{2} & x_{3}^{3} & x_{3}^{4} & x_{3}^{5} & 0 & 0 & 0 & 0 \\
0 & 0 & 0 & 0 & x_{4}^{2} & x_{4}^{3} & x_{4}^{4} & x_{4}^{5} & 0 & 0 & 0 & 0 \\
0 & 0 & 0 & 0 & x_{5}^{2} & x_{5}^{3} & x_{5}^{4} & x_{5}^{5} & 0 & 0 & 0 & 0 \\
-x_{1}^{2} & -x_{1}^{3} & -x_{1}^{4} & -x_{1}^{5} & -x_{2}^{2} & -x_{2}^{3} & -x_{2}^{4} & -x_{2}^{5} & 0 & 0 & 0 & 0 \\
0 & 0 & 0 & 0 & 0 & 0 & 0 & 0 & x_{4}^{3} & x_{4}^{4} & x_{4}^{5} & 0 \\
0 & 0 & 0 & 0 & 0 & 0 & 0 & 0 & x_{5}^{3} & x_{5}^{4} & x_{5}^{5} & 0 \\
-x_{1}^{2} & -x_{1}^{3} & -x_{1}^{4} & -x_{1}^{5} & -x_{2}^{2} & -x_{2}^{3} & -x_{2}^{4} & -x_{2}^{5} & -x_{3}^{3} & -x_{3}^{4} & -x_{3}^{5} & 0 \\
0 & 0 & 0 & 0 & 0 & 0 & 0 & 0 & 0 & 0 & 0 & x_{5}^{4} \\
-x_{1}^{2} & -x_{1}^{3} & -x_{1}^{4} & -x_{1}^{5} & -x_{2}^{2} & -x_{2}^{3} & -x_{2}^{4} & -x_{2}^{5} & -x_{3}^{3} & -x_{3}^{4} & -x_{3}^{5} & -x_{4}^{4}
\end{array}\right)
$$

$$
M_{4}=\left(\begin{array}{ccccc}
x_{2}^{2} & x_{2}^{3} & x_{2}^{4} & x_{2}^{5} & 0 \\
x_{3}^{2} & x_{3}^{3} & x_{3}^{4} & x_{3}^{5} & 0 \\
x_{4}^{2} & x_{4}^{3} & x_{4}^{4} & x_{4}^{5} & 0 \\
x_{5}^{2} & x_{5}^{3} & x_{5}^{4} & x_{5}^{5} & 0 \\
-x_{1}^{2} & -x_{1}^{3} & -x_{1}^{4} & -x_{1}^{5} & 0 \\
0 & 0 & 0 & 0 & x_{3}^{2} \\
0 & 0 & 0 & 0 & x_{4}^{2} \\
0 & 0 & 0 & 0 & x_{5}^{2} \\
-x_{1}^{2} & -x_{1}^{3} & -x_{1}^{4} & -x_{1}^{5} & -x_{2}^{2}
\end{array}\right), \quad M_{5}=\left(\begin{array}{c}
x_{2}^{1} \\
x_{3}^{1} \\
x_{4}^{1} \\
x_{5}^{1} \\
-x_{1}^{1}
\end{array}\right) .
$$

Then $\operatorname{det}_{5}(x)=M_{1} M_{2} M_{3} M_{4} M_{5}$.

Let


Then $\operatorname{det}_{6}(X)=M_{1} M_{2} M_{3} M_{4} M_{5} M_{6}$.
Compare these with the expression from [LR15]:

$$
\operatorname{det}_{3}(x)=\left(x_{1}^{1}, x_{1}^{2}, x_{1}^{3}\right)\left(\begin{array}{ccc}
x_{2}^{2} & -x_{2}^{3} & 0  \tag{7.3.1}\\
-x_{2}^{1} & 0 & x_{2}^{3} \\
0 & x_{2}^{1} & -x_{2}^{2}
\end{array}\right)\left(\begin{array}{l}
x_{3}^{3} \\
x_{3}^{2} \\
x_{3}^{1}
\end{array}\right) .
$$

and for $\operatorname{det}_{4}$ the sizes of the matrices are $1 \times 4,4 \times 6,6 \times 4,4 \times 1$.
7.3.3. A classical exponential lower bound for the permanent (and determinant). Consider a restricted model where one is not allowed to exploit the commutativity of multiplication. Let $\mathbb{C}\left\{y_{1}, \ldots, y_{N}\right\}$ denote the ring of polynomials in the non-commuting variables $y_{1}, \ldots, y_{N}$. Choose an expression for a polynomial $P$ and consider it in this larger ring. The definition of circuits is the same here, just that we cannot assume $a b=b a$ for expressions $a$ and $b$.
Theorem 7.3.3.1. [Nis91] The degree homogeneous non-commutative ABP complexity of $\operatorname{det}_{m}$ and $\operatorname{perm}_{m}$ are both $2^{m}-1$.

Proof. Choose the representations of the determinant and permanent where the first row comes first, the second comes second, etc. Consider the degree homogeneous ABP $\Gamma$ with $m+1$ layers that computes $\operatorname{det}_{m}\left(\right.$ or $\operatorname{perm}_{m}$ ). Keep the labels from all edges that appear before level $s$ and set the labels on all other layers to constants to see that all terms of the form $\sum_{\sigma \in \mathfrak{S}_{m}} c_{\sigma} y_{\sigma(1)}^{1} \cdots y_{\sigma(s)}^{s}$ can be computed by taking linear combinations of the polynomials $\Gamma_{v}$, where $v$ is a vertex in layer $s$. Since these terms span a vector space of dimension $\binom{m}{s}$ there must be at least $\binom{m}{s}$ linearly independent polynomials $\Gamma_{v}$, so there must be at least $\binom{m}{s}$ vertices on layer $s$. Summing up the binomial coefficients yields the lower bound.

The Grenet determinantal presentation of perm $_{m}[$ Gre11] and the regular determinantal presentation of $\operatorname{det}_{m}$ of [LR15] give rise to columnwise multi-linear iterated matrix multiplication presentations, and thus noncommutative ABP's, of size $2^{m}-1$.

Remark 7.3.3.2. In contrast to ABP's, for general non-commutative circuits, very little is known, see, e.g., [LMS16, HsWY10]. There are exponential bounds for skew circuits in [LMS16] (the class of circuits equivalent in power to the determinant).
7.3.4. Read once ABP's. Another restriction of ABP's is that of read once oblivious ABP's, henceforth ROABP's. Here the ABP is layered. The read-once means that the edges at layer $i$ only use a variable $x_{i}$. On the other hand, the weights are allowed to be low degree polynomials in the $x_{i}$. The word oblivious means additionally that an ordering of the variables is fixed in
advance. I return to this model in $\S 7.7$ because it is restrictive enough that it admits explicit deterministic hitting sets for polynomial identity testing. On the other hand, this model can efficiently simulate depth three powering circuits.

### 7.4. Additional restricted models

The purpose of this section is to survey restricted models that have geometric aspects. Each subsection may be read independently of the others.
7.4.1. Equivariant determinantal complexity. Motivated by the symmetry of Grenet's expressions for the permanent discussed in $\S 6.6 .3$, N . Ressayre and I asked, what happens if one imposes the $\Gamma_{m}^{E}$-equivariance? We found:
Theorem 7.4.1.1. [LR15] Among $\Gamma_{m}^{E}$-equivariant determinantal expressions for perm $_{m}$, Grenet's size $2^{m}-1$ expressions are optimal and unique up to trivialities.

The $\Gamma_{m}^{E}$-equivariance is peculiar as it only makes sense for the permanent. To fix this, we defined a complexity measure that could be applied to all polynomials:

Let $P \in S^{m} \mathbb{C}^{M}$ have symmetry group $G_{P}$, let $A: \mathbb{C}^{M} \rightarrow \mathbb{C}^{n^{2}}$ be the linear part of a determinantal expression of $P$ with constant term $\Lambda$. Let $G_{\operatorname{det}_{n}, \Lambda}=G_{\operatorname{det}_{n}} \cap G_{\Lambda} \subset G L_{n^{2}}$. Note that $G_{P} \times G_{\operatorname{det}_{n}, \Lambda}$ acts on $\mathbb{C}^{M^{*}} \otimes \mathbb{C}^{n^{2}}$ by $(g, h) A(y):=h \cdot A\left(g^{-1} y\right)$.

Definition 7.4.1.2. Define the symmetry group of $\tilde{A}$ to be

$$
G_{\tilde{A}}:=\left\{(g, h) \in G_{P} \times G_{\operatorname{det}_{n}, \Lambda} \mid(g, h) \cdot A=A\right\}
$$

Call $\tilde{A}$ an equivariant determinantal expression for $P$ if the projection from $G_{\tilde{A}}$ to $G_{P}$ is surjective. Define $\operatorname{edc}(P)$ to be the smallest size of an equivariant determinantal expression for $P$.

If $G$ is a subgroup of $G_{P}$, we say that $\tilde{A}$ is $G$-equivariant if $G \subseteq G_{\tilde{A}}$.
Note that if $P$ is a generic polynomial of degree greater than two, $\operatorname{edc}(P)=\operatorname{dc}(P)$ because it will have a trivial symmetry group. One also has $\operatorname{edc}\left(\operatorname{det}_{m}\right)=\operatorname{dc}\left(\operatorname{det}_{m}\right)$ because $A=\operatorname{Id}: \mathbb{C}^{n^{2}} \rightarrow \mathbb{C}^{n^{2}}$ and $\Lambda=0$ is an equivariant expression.
Theorem 7.4.1.3. [LR15] There exists an equivariant determinantal expression for perm $_{m}$ of size $\binom{2 m}{m}-1$.
Theorem 7.4.1.4. [LR15] Among equivariant determinatal expressions for perm $_{m}$, the size $\binom{2 m}{m}-1$ expressions are optimal and unique up to trivialities.

In particular, Valiant's hypothesis holds in the restricted model of equivariant expressions. To my knowledge, equivariant determinantal complexity is the only restricted model with a known exponential separation of the permanent from the determinant.

Proofs are outlined in $\S 8.11 .2$.
Note that $\binom{2 m}{m} \sim 4^{m}$ so the size of the equivariant determinantal expressions are roughly the square of the size of Grenet's expressions. In particular, they are polynomially related in size.

Thus, if one could show either

- there exists an optimal determinantal expression for perm $_{m}$ with some symmetry, or
- there exists an equivariant determinantal expression for perm $_{m}$ of size polynomial in dc $\left(\operatorname{perm}_{m}\right)$,
then one would have proven Valiant's hypothesis. I write "some" symmetry, because as is shown in the proof, full $\Gamma_{m}^{E}$-symmetry is not needed for the exponential lower bound. (I do not know just how large the symmetry group needs to be to obtain an exponential bound.)

Regarding the possibility of proving either of the above, we have seen that the optimal Waring rank expression for $x_{1} \cdots x_{n}$ (and more generally odd degree elementary symmetric functions) have maximal symmetry, as does the optimal rank expression for $M_{\langle 2\rangle}$.
7.4.2. Elementary symmetric polynomial complexity. Let $P \in S^{m} \mathbb{C}^{k}$ and define the elementary symmetric complexity of $P, \operatorname{esc}(P)$, to be the smallest $N$ such that there exists a linear inclusion $\mathbb{C}^{k} \subset \mathbb{C}^{n}$ with $P \in$ $\operatorname{End}\left(\mathbb{C}^{N}\right) \cdot e_{m, N}=: \hat{\mathcal{E}} l e m e n_{m, N}^{0}$, and $\overline{\mathrm{esc}}(P)$ to be the smallest $N$ such that $P \in \overline{\operatorname{End}\left(\mathbb{C}^{N}\right) \cdot e_{m, N}}=\overline{G L_{N} \cdot e_{m, N}}=: \hat{\mathcal{E}}$ lemen $_{m, N}$. A. Shpilka [Shp02] refers to $\operatorname{esc}(P)$ as the "size of the smallest depth two circuit with a symmetric gate at the top and plus gates at the bottom".

For any polynomial $P, \operatorname{esc}(P)$ is finite. More precisely:
Proposition 7.4.2.1. $[\mathbf{S h p 0 2}] \sigma_{r}^{0}\left(v_{m}(\mathbb{P} V)\right) \subset \mathcal{E}$ lemen $n_{m, r m}^{0}$ and $\sigma_{r}\left(v_{m}(\mathbb{P} V)\right) \subset$ $\mathcal{E l e m e n}_{m, r m}$. In other words, if $P \in S^{d} V$ is computable by a $\Sigma \Lambda \Sigma$ circuit of size $r$ then $\operatorname{esc}(P) \leq r m$.

Proof. Without loss of generality, assume $\mathbf{v}=r$ and let $y_{1}, \ldots, y_{r}$ be a basis of $V$. It will be sufficient to show $\sum y_{j}^{m} \in \mathcal{E}$ lemen ${ }_{m, m r}^{0}$. Let $\omega$ be a primitive $m$-th root of unity. Then I claim

$$
\sum y_{j}^{m}=-e_{m, r m}\left(y_{1},-\omega y_{1},-\omega^{2} y_{1}, \ldots, \omega^{m-1} y_{1},-y_{2},-\omega y_{2}, \ldots,-\omega^{m-1} y_{r}\right)
$$

To see this, evaluate the generating function:

$$
\begin{aligned}
& E_{r m}(t)\left(y_{1},-\omega y_{1},-\omega^{2} y_{1}, \ldots, \omega^{m-1} y_{1},-y_{2},-\omega y_{2}, \ldots,-\omega^{m-1} y_{r}\right) \\
& =\prod_{i \in[r]} \prod_{s \in[m]}\left(1-\omega^{s} y_{i}\right) \\
& =\prod_{i \in[r]}\left(1-y_{i}^{m} t^{m}\right) .
\end{aligned}
$$

The coefficient of $t^{m}$ on the last line is $-\sum_{i} y_{i}^{m}$.
Note that $\operatorname{dim}\left(\right.$ Elemen $\left._{m, r m}\right) \leq r^{2} m^{2}$ while $\operatorname{dim}\left(\sigma_{r}\left(v_{m}\left(\mathbb{P}^{r m-1}\right)\right)\right)=r m^{2}-$ 1 , so the dimensions differ only by a factor of $r$. Contrast this with the inclusion implied by Theorem 7.1.3.1 of $\mathcal{E}$ Lemen $_{d, N} \subset \sigma_{q}\left(v_{d}\left(\mathbb{P}^{N-1}\right)\right.$ with $q=\sum_{j=0}^{\left\lfloor\frac{d}{2}\right\rfloor}\binom{N}{j}$ where the second space in general has dimension exponentially larger than the first.

Regarding lower bounds for esc, Corollary 7.2.3.2 implies that esc $(P)$ is at least the square root of the size of the smallest depth three circuit computing $P$.

Shpilka proves lower bounds for esc in the same way the first lower bounds for determinantal complexity were found: by considering linear spaces on the zero set $\operatorname{Zeros}\left(e_{m, N}\right) \subset \mathbb{P}^{N-1}$.
Theorem 7.4.2.2. $[\mathbf{S h p 0 2}]$ Let $L \subset \operatorname{Zeros}\left(e_{m, N}\right) \subset \mathbb{P}^{N-1}$ be a linear space. Then $\operatorname{dim} L \leq \min \left(\max (N-m, m-1), \frac{m+N}{2}\right)-1$.

Proof. The key to the proof is the algebraic independence of the $e_{j, N}$ (see, e.g., [Mac95, §1.2]). Any linear space of dimension $k$ will have an isomorphic projection onto some coordinate $k$-plane. Without loss of generality, assume it has an isomorphic projection onto the span of the first $k$-coordinates, so that $\hat{L} \subset \mathbb{C}^{N}$ has equations $x_{s}=\ell_{s}\left(x_{1}, \ldots, x_{k}\right)$ for $k+1 \leq s \leq N$. We are assuming $\left.e_{m, N}\right|_{\hat{L}}=0$.
Exercise 7.4.2.3: (1) Show that if we have two sets of variables $(x, y)=$ $\left(x_{1}, \ldots, x_{k}, y_{1}, \ldots, y_{N-k}\right)$, then $e_{m, N}(x, y)=\sum_{j=0}^{m} e_{m-j, k}(x) e_{j, N-k}(y)$.

By Exercise 7.4.2.3,

$$
\begin{align*}
0 & =e_{m, N}(x, \ell(x)) \\
& =e_{m, k}(x)+\sum_{j=1}^{m} e_{m-j, k}(x) e_{j, N-k}(\ell(x)) . \tag{7.4.1}
\end{align*}
$$

First assume $k=\operatorname{dim} \hat{L} \geq \max (N-m+1, m)$. Since $e_{k, u}=0$ if $k>u$, if $N-k<m$ the sum in (7.4.1) is from 1 to $N-k$.

Let $\Psi: \mathbb{C}\left[x_{1}, \ldots, x_{k}\right] \rightarrow \mathbb{C}\left[x_{1}, \ldots, x_{k}\right]^{\mathfrak{S}_{k}}$ denote the symmetrization operator. (Sometimes $\Psi$ is called a Reynolds operator.)
Exercise 7.4.2.4: (1) Show that for any functions $f, g$, that $\Psi(f+g)=$ $\Psi(f)+\Psi(g)$.
Exercise 7.4.2.5: (1) Show that if $f$ is a symmetric function and $g$ is a polynomial, then $\Psi(f g)=\Psi(f) \Psi(g)$.

Apply $\Psi$ to (7.4.1) to obtain

$$
0=e_{m, k}(x)+\sum_{j=1}^{N-k} e_{m-j, k}(x) \Psi\left(e_{j}(\ell(x))\right)
$$

but this expresses $e_{m, k}$ as a polynomial in symmetric functions of degree less than $k$, a contradiction.

Now assume $\operatorname{dim} \hat{L} \geq \frac{m+N}{2}$, so

$$
0=e_{m, k}(x)+e_{m, N-k}(\ell(x))+\sum_{j=1}^{m} e_{m-j, k}(x) e_{j}(\ell(x))
$$

The idea is again the same, but we must somehow reduce to a smaller space. If we take $D \in\left\{\ell_{1}, \ldots, \ell_{N-k}\right\}^{\perp} \subset \mathbb{C}^{N}$ and apply it, we can eliminate the $e_{m, N-k}(\ell(x))$ term. But if we take a general such $D$, we will no longer have symmetric functions. However, one can find a $D$ such that, if we restrict to span of the first $m-1$ coordinate vectors, call this space $V_{m-1} \subset \mathbb{C}^{k} \subset \mathbb{C}^{N}$, then $\left.\left(D e_{r, k}\right)\right|_{V_{m-1}}=e_{r-1, m-1}$, see [Shp02]. Unfortunately this is still not good enough, as letting $x^{\prime}=\left(x_{1}, \ldots, x_{m-1}\right)$ we now have

$$
0=e_{m-1, m-1}\left(x^{\prime}\right)+\sum_{j=1}^{m} e_{m-j, k}\left(x^{\prime}\right) e_{j}\left(\ell\left(x^{\prime}\right)\right)
$$

We could argue as before if we could eliminate the $j=1$ term. A modification of $D$ as described in [Shp02] also satisfies $D\left(e_{1, k}(x)\right)=0$.

Thus if $\operatorname{Zeros}(P)$ has large linear spaces on it we obtain lower bounds for $\operatorname{esc}(P)$. Recall that for a projective subspace $L \subset \mathbb{P}^{N-1}$, that $\hat{L} \subset \mathbb{C}^{N}$ denotes the corresponding linear subspace.
Exercise 7.4.2.6: (1) Show $\overline{\mathrm{esc}}\left(\operatorname{det}_{m}\right) \geq 2 m^{2}-3 m$.
Exercise 7.4.2.7: (1) Show that if $m \geq \frac{N+1}{2}$, there exists a linear space of dimension $m-2$ on $\operatorname{Zeros}\left(e_{m, N}\right)$. ©

Say $m$ is odd and $N$ is even. Let

$$
\hat{L}=\operatorname{span}\{(1,-1,0, \ldots, 0),(0,0,1,-1,0, \ldots, 0),, \ldots,(0, \ldots, 0,1,-1)
$$

Notice that all odd power sum functions vanish on $\hat{L}$. When we express $e_{m, N}$ in terms of power sum functions, each term will contain an odd degree power sum so we conclude $\left.e_{m, N}\right|_{\hat{L}}=0$. More generally:
Proposition 7.4.2.8. [Shp02] [attributed to Saks] There exists a $\mathbb{P}^{\left\lfloor\frac{N}{q}\right\rfloor-1} \subset$ Zeros $\left(e_{m, N}\right)$, where $q$ is the smallest integer such that $q$ does not divide $m$.
Exercise 7.4.2.9: (2) Prove Proposition 7.4.2.8. ©
Exercise 7.4.2.7 and Proposition 7.4.2.8 show that Theorem 7.4.2.2 is close to being sharp.

The following conjecture appeared in [Shp02] (phrased differently):
Conjecture 7.4.2.10. [Shp02] There exists a polynomial $r(m)$ such that $\sigma_{r(m)}\left(C h_{m}\left(\mathbb{C}^{m r(m)}\right)\right) \not \subset \mathcal{E}$ lemen ${ }_{m, 2^{m}}$. One might even be able to take $r(m) \equiv$ 2.

The second assertion is quite strong, as when $r=1$ there is containment, and when $r=2$ the left hand side has dimension about $4 m$ and the right hand side has dimension about $4^{m}$.
Exercise 7.4.2.11: (2) Show that $\sigma_{2}\left(C h_{m}\left(\mathbb{C}^{2 m}\right)\right) \not \subset \mathcal{E l e m e n}_{m, \frac{3}{2} m-3}$.
Question 7.4.2.12. [Shp02] What is the maximal dimension of a linear subspace $L \subset \mathbb{P}^{N-1}$ such that $L \subset \operatorname{Zeros}\left(e_{m, N}\right)$ ?
7.4.3. Raz's theorem on tensor rank and formula size. In this section I explain Raz's results that if one considers a tensor as a polynomial, lower bounds on the tensor rank have consequences for the formula size of the corresponding polynomial.

Definition 7.4.3.1. A polynomial $P \in S^{d} V$ is multi-linear if $V=V_{1} \oplus \cdots \oplus$ $V_{d}$ and $P \in V_{1} \otimes \cdots \otimes V_{d} \subset S^{d} V$.

The permanent and determinant may be considered as multi-linear polynomials (in two different ways). In the literature, e.g., [Raz10b], they do not insist on homogeneous polynomials, so they use the term set-multi-linear to describe such polynomials where each monomial appearing is multi-linear (but does not necessarily use variables from each of the $V_{j}$ ).

Given a tensor $T \in A_{1} \otimes \cdots \otimes A_{d}$, by considering $A_{1} \otimes \cdots \otimes A_{d} \subset S^{d}\left(A_{1} \oplus\right.$ $\cdots \oplus A_{d}$ ), we may think of $T$ as defining a multi-linear polynomial. When I want to emphasize $T$ as a multi-linear polynomial, I'll write $P_{T} \in S^{d}\left(A_{1} \oplus\right.$ $\left.\cdots \oplus A_{d}\right)$.

One can compare the tensor rank of $T$ with the circuit complexity of $P_{T}$. Raz compares it with the formula complexity: He shows that superpolynomial lower bounds for multi-linear formulas for polynomial sequences
$P_{n}$ where the degree grows slowly, imply super-polynomial lower bounds for general formulas:
Theorem 7.4.3.2. $\left[\right.$ Raz10b] Let $\operatorname{dim} A_{j}=n$ and let $T_{n} \in A_{1} \otimes \cdots \otimes A_{d}$ be a sequence of tensors with $d=d(n)$ satisfying $d=O\left(\frac{\log (n)}{\log (\log (n))}\right)$. If there exists a formula of size $n^{C}$ for $P_{T_{n}}$, then $\mathbf{R}\left(T_{n}\right) \leq n^{d\left(1-2^{O(C)}\right)}$.
Corollary 7.4.3.3. [Raz10b] Let $\operatorname{dim} A_{j}=n$ and let $T_{n} \in A_{1} \otimes \cdots \otimes A_{d}$ be a sequence of tensors with $d=d(n)$ satisfying $d=O\left(\frac{\log (n)}{\log (\log (n))}\right)$. If $\mathbf{R}\left(T_{n}\right) \geq n^{d(1-o(1))}$, then there is no polynomial size formula for $P_{T}$.

These results were extended in [CKSV16].
Via flattenings, one can exhibit explicit tensors with $\underline{\mathbf{R}}(T) \geq n^{\left\lfloor\frac{d}{2}\right\rfloor}$. Using the substitution method (see $\S 5.3$ ), that was improved for tensor rank to $2 n^{\left\lfloor\frac{d}{2}\right\rfloor}+n-O(d \log (n))$ in [AFT11] by a construction generalizing the one described in $\S 5.3 .1$ for the case $d=3$.

The idea of proof is as follows: A rank decomposition of $T$, viewed as a computation of $P_{T}$, corresponds to a depth-3 multi-linear formula for $P_{T}$. Raz shows that for any polynomial sequence $P_{n}$, if there is a fanin- 2 formula of size $s$ and depth $\delta$ for $P$, then there exists a homogeneous formula of size $O\left(\binom{\delta+d+1}{d} s\right)$ for $P_{n}$. He then shows that for any multi-linear polynomial $P_{n}$, if there exists a fanin- 2 formula of size $s$ and depth $\delta$, then there exists a multi-linear formula of size $O\left((\delta+2)^{d} s\right)$ for $P_{n}$.
7.4.4. Multi-linear formulas. A formula is multi-linear if the polynomial computed by each of its sub-formulas is multi-linear. For example, Ryser's formula for the permanent is multi-linear. On the other hand, the smallest known formula for the determinant is not multi-linear.

In [Raz09], Raz shows that any multi-linear arithmetic formula for perm ${ }_{n}$ or $\operatorname{det}_{n}$ is of size $n^{\Omega(n)}$. The starting point of the proof is the method of partial derivatives. Then Raz makes certain reductions, called random restrictions to reduce to a smaller polynomial that one can estimate more precisely.
7.4.5. Raz's elusive functions and circuit lower bounds. Raz defines the following "hay in a haystack" approach to Valiant's hypothesis. Consider a linear projection of a Veronese proj : $\mathbb{P} S^{r} \mathbb{C}^{s} \rightarrow \mathbb{P}^{m}$, and let $\Gamma_{r, s}:=$ proj $\circ v_{r}: \mathbb{P}^{s-1} \rightarrow \mathbb{P}^{m}$ be the composition of the projection with the Veronese map. A map $f: \mathbb{P}^{n} \rightarrow \mathbb{P}^{m}$ is said to be $(r, s)$-elusive if $f\left(\mathbb{P}^{n}\right)$ is not contained in the image of any such $\Gamma_{r, s}$.

Recall that VNP may be thought of as the set of "explicit" polynomial sequences.

Theorem 7.4.5.1. $[\operatorname{Raz} 10 a]$ Let $m$ be super-polynomial in $n$, and $s \geq m^{\frac{9}{10}}$. If there exists an explicit $(s, 2)$-elusive $f: \mathbb{P}^{n} \rightarrow \mathbb{P}^{m}$, then VP $\neq$ VNP.
Theorem 7.4.5.2. [Raz10a] Let $r(n)=\log (\log (n)), s(n)=n^{\log (\log (\log (n)))}$, $m=n^{r}$, and let $C$ be a constant. If there exists an explicit $(s, r)$-elusive $f: \mathbb{P}^{n} \rightarrow \mathbb{P}^{m}$, then VP $\neq \mathbf{V N P}$.

By a dimension count, a general polynomial in either range will be elusive.

Again, one can approach, e.g., the case where $r=2$, by finding equations for the variety of all images of projections of the quadratic Veronese, and then finding a polynomial (point) not in the zero set.

In the same paper Raz constructs an explicit $f$, whose entries are monomials, that requires circuits of size at least $n^{1+\Omega\left(\frac{1}{r}\right)}$ to compute.
7.4.6. Glynn's Theorem on expressions for the permanent. Recall, for $P \in S^{m} \mathbb{C}^{M}, \mathbf{R}_{C h_{m}\left(\mathbb{C}^{M}\right)}(P)$ is the smallest $r$ such that $P\left(y_{1}, \ldots, y_{M}\right)=$ $\sum_{s=1}^{r} \Pi_{u=1}^{m}\left(\sum_{a=1}^{M} \lambda_{s, u, a} y_{a}\right)$ for some constants $\lambda_{s, u, a}$. This corresponds to the smallest homogeneous $\Sigma^{r} \Pi^{m} \Sigma^{M}$ circuit that computes $P$. If $P$ is multilinear, so $M=m w$ and we may write $y_{a}=\left(y_{i \alpha}\right)$ where $1 \leq i \leq m, 1 \leq$ $\alpha \leq w$, and $P=\sum C_{\alpha} y_{1 \alpha} \cdots y_{m \alpha}$ we could restrict to multi-linear $\Sigma \Pi \Sigma$ circuits (ML- $\Sigma \Pi \Sigma$ circuits), those of the form $\sum_{s=1}^{r} \Pi_{i=1}^{m}\left(\sum_{\alpha=1}^{w} \lambda_{s, \alpha} y_{i \alpha}\right)$. Write $\mathbf{R}_{C h_{m}\left(\mathbb{C}^{M}\right)}^{M L}(P)$ for the smallest multi-linear $\Sigma^{r} \Pi \Sigma^{w}$ circuit for such a $P$. Consider multi-linear $\Sigma \Pi \Sigma$-circuit complexity as a restricted model. In this context, we have the following theorem of D. Glynn:
Theorem 7.4.6.1. [Gly13] $\mathbf{R}_{C h_{m}\left(\mathbb{C}^{M}\right)}^{M L}\left(\operatorname{perm}_{m}\right)=\mathbf{R}_{S}\left(x_{1} \cdots x_{m}\right)=2^{m-1}$.
Moreover, there is a one to one correspondence between Waring decompositions of $x_{1} \cdots x_{m}$ and $M L-\Sigma \Pi \Sigma$ decompositions of perm $_{m}$. The correspondence is as follows: Constants $\lambda_{s, j}, 1 \leq s \leq r, 1 \leq j \leq m$ satisfy

$$
\begin{equation*}
x_{1} \cdots x_{m}=\sum_{s=1}^{r}\left(\sum_{j=1}^{m} \lambda_{s, j} x_{j}\right)^{m} \tag{7.4.2}
\end{equation*}
$$

if and only if

$$
\begin{equation*}
\operatorname{perm}_{m}\left(y_{i j}\right)=m!\sum_{s=1}^{r} \prod_{i=1}^{m}\left(\sum_{j=1}^{m} \lambda_{s, j} y_{i j}\right) . \tag{7.4.3}
\end{equation*}
$$

Proof. Given a Waring decomposition (7.4.2) of $x_{1} \cdots x_{m}$, set $x_{j}=\sum_{k} y_{j k} z_{k}$. The coefficient of $z_{1} \cdots z_{m}$ in the resulting expression on the left hand side is the permanent and the coefficient of $z_{1} \cdots z_{m}$ on the right hand side is the right hand side of (7.4.3).

To see the other direction, given an expression (7.4.3), I will specialize to various matrices to show identities among the $\lambda_{s, j}$ that will imply all coefficients but the desired one on the right hand side of (7.4.2) are zero.

The coefficient of $x_{1}^{b_{1}} \cdots x_{m}^{b_{m}}$, where $b_{1}+\cdots+b_{m}=m$ in (7.4.2) is $\binom{m}{b_{1}, \ldots, b_{m}} \sum_{s} \lambda_{s, 1}^{b_{1}} \cdots \lambda_{s, m}^{b_{m}}$.

Let $y$ be a matrix where there are $b_{j} 1$ 's in column $j$ and zero elsewhere. Then unless each $b_{j}=1$, $\operatorname{perm}(y)=0$. But (7.4.3) says that $0=\operatorname{perm}(y)$ is a nonzero constant times $\sum_{s} \lambda_{s, 1}^{b_{1}} \cdots \lambda_{s, m}^{b_{m}}$. Thus all these terms are zero and the only potential nonzero coefficient in the right hand side of (7.4.2) is the coefficient of $x_{1} \cdots x_{m}$. This coefficient is $m!=\left(\begin{array}{c}m, \ldots, 1\end{array}\right)$ times $\lambda_{s, 1} \cdots \lambda_{s, m}$. Plugging in $y=$ Id shows $1=m!\lambda_{s, 1} \cdots \lambda_{s, m}$.

Remark 7.4.6.2. As mentioned in Remark 7.1.2.1, all rank $2^{m-1}$ expressions for $x_{1} \cdots x_{m}$ come from the $T^{S L_{m}}$ orbit of (7.1.1), so the same holds for size $2^{m-1} M L-\Sigma \Pi \Sigma$ expressions for perm ${ }_{m}$.
7.4.7. Rank $k$ determinantal expressions. Restricted models with a parameter $k$ that converge to the original problem as $k$ grows are particularly appealing, as one can measure progress towards the original conjecture. Here is one such: Given a polynomial $P \in S^{m} \mathbb{C}^{M}$ and determinantal expression $\tilde{A}: \mathbb{C}^{M} \rightarrow \mathbb{C}^{n^{2}}, \tilde{A}(y)=\Lambda+\sum_{j=1}^{M} A_{j} y_{j}$ where $\Lambda, A_{j}$ are matrices, define the rank of $\tilde{A}$ to be the largest rank of the $A_{j}$ 's. Note that this depends on the coordinates up to rescaling them, but for the permanent this is not a problem, as $G_{\text {perm }_{m}}$ defines the coordinates up to scale.

If one could show that perm ${ }_{m}$ did not admit an expression with rank polynomial in $m$, then that would trivially prove Valiant's hypothesis.

The notation of rank of a determinantal expression was introduced in [AJ15], as a generalization of the read of a determinantal expression, which is the maximal number of nonzero entries of the $A_{j}$. As observed by Anderson, Shpilka and Volk (personal communication from Shpilka) as far as complexity is concerned the measures are equivalent: if a polynomial $P$ in $n$ variables admits a rank $k$ determinantal expression of size $s$, then it admits a read- $k$ determinantal expression of size $s+2 n k$.

The state of the art regarding this model is not very impressive:
Theorem 7.4.7.1. [IL16a] The polynomial perm $m_{m}$ does not admit a rank one determinantal expression over $\mathbb{C}$ when $m \geq 3$. In particular, perm $_{m}$ does not admit a read once regular determinantal expression over $\mathbb{C}$ when $m \geq 3$.

### 7.5. Shallow Circuits and Valiant's hypothesis

In this section I discuss three classes of shallow circuits that could be used to prove Valiant's hypothesis. We have already seen the first, the $\Sigma \Pi \Sigma$ circuits. The next is the $\Sigma \Lambda \Sigma \Lambda \Sigma$ circuits, which are depth five circuits where the first layer of gates are additions, the second layer consists of "powering gates", where a powering gate takes $f$ to $f^{\delta}$ for some natural number $\delta$, the third layer addition gates, the fourth layer again powering gates, and the fifth layer is an addition gate. The third is the class of depth four $\Sigma \Pi \Sigma \Pi$ circuits. I describe the associated varieties to these classes of circuits in §7.5.3. A $\Sigma \Lambda^{\alpha} \Sigma \Lambda^{\beta} \Sigma$ circuit means the powers are respectively $\beta$ and $\alpha$, and other superscripts are to be similarly interpreted.
7.5.1. Detour for those not familiar with big numbers. When dealing with shallow circuits, we will have to distinguish between different rates of super-polynomial growth, both in statements and proofs of theorems. This detour is for those readers not used to comparing large numbers.

All these identities follow from (7.5.1), which follows from Stirling's formula, which gives an approximation for the Gamma function, e.g., for $x>0$,

$$
\Gamma(x)=\sqrt{2 \pi} x^{x-\frac{1}{2}} e^{-x} e^{\frac{\theta(x)}{12 x}}
$$

where $0<\theta(x)<1$. Stirling's formula may be proved via complex analysis (estimating a contour integral), see, e.g., [Ah178, §5.2.5]. Let

$$
H_{e}(x):=-x \ln x-(1-x) \ln (1-x)
$$

denote the Shannon entropy.

$$
\begin{align*}
n! & \gtrsim \sqrt{2 \pi n}\left(\frac{n}{e}\right)^{n}  \tag{7.5.1}\\
\ln (n!) & =n \ln (n)-O(\ln (n))  \tag{7.5.2}\\
\binom{2 n}{n} & \gtrsim \frac{4^{n}}{\sqrt{\pi n}}  \tag{7.5.3}\\
\ln \binom{\alpha n}{\beta n} & =\alpha H_{e}\left(\frac{\beta}{\alpha}\right) n-O(\ln n)  \tag{7.5.4}\\
\binom{\alpha n}{\beta n} & =\left[\frac{\alpha^{\alpha}}{\beta^{\beta}(\alpha-\beta)^{\alpha-\beta}}\right]^{n} O\left(\frac{1}{n}\right) . \tag{7.5.5}
\end{align*}
$$

Exercise 7.5.1.1: (1) Show that for $0<x<1,0<H_{e}(x) \leq 1$. For which $x$ is the maximum achieved?
Exercise 7.5.1.2: (1) Show $a^{\log (b)}=b^{\log (a)}$.

Exercise 7.5.1.3: (1!) Consider the following sequences of $n$ :
$\log _{2}(n), n, 100 n, n^{2}, n^{3}, n^{\log _{2}(n)}, 2^{\left[\log _{2}(n)\right]^{2}}, n^{\sqrt{\log _{2}(n)}}, 2^{n},\binom{2 n}{n}, n!, n^{n}$.
In each case, determine for which $n$, the sequence surpasses the number of atoms in the known universe. (It is estimated that there are between $10^{78}$ and $10^{82}$ atoms in the known universe.)
Exercise 7.5.1.4: (1) Compare the growth of $s^{\sqrt{d}}$ and $2^{\sqrt{d \log d s}}$.
Exercise 7.5.1.5: (1!) Compare the growth of $\binom{n^{2}+\frac{n}{2}-1}{\frac{n}{2}}$ and $\binom{n}{\frac{n}{2}}^{2}$. Compare with your answer to Exercise 6.2.2.7.
7.5.2. Depth reduction theorems. A major result in the study of shallow circuits was [VSBR83], where it was shown that if a polynomial of degree $d$ can be computed by a circuit of size $s$, then it can be computed by a circuit of depth $O(\log d \log s)$ and size polynomial in $s$. Since then there has been considerable work on shallow circuits. See, e.g., [GKKS17] for a history.

Here are the results relevant for our discussion. They combine results of [Bre74, GKKS13b, Tav15, Koi, AV08]:
Theorem 7.5.2.1. Let $N=N(d)$ be a polynomial and let $P_{d} \in S^{d} \mathbb{C}^{N}$ be a sequence of polynomials that can be computed by a circuit of polynomial size $s=s(d)$. Let $S(d):=2^{O(\sqrt{d \log (d s) \log (N)})}$.

Then:
(1) $P$ is computable by a homogeneous $\Sigma \Pi \Sigma \Pi$ circuit of size $S(d)$.
(2) $P$ is computable by a $\Sigma \Pi \Sigma$ circuit of size of size $S(d)$.
(3) $P$ is computable, by a homogeneous $\Sigma \Lambda \Sigma \Lambda \Sigma$ circuit of size $S(d)$, and both powering gates of size $O(\sqrt{d})$.

Note that $S(d)$ is approximately $s^{\sqrt{d}}$.
Corollary 7.5.2.2. If $\operatorname{perm}_{m}$ is not computable by one of: a homogeneous $\Sigma \Pi \Sigma \Pi$ circuit, a $\Sigma \Pi \Sigma$ circuit, or a homogeneous $\Sigma \Lambda \Sigma \Lambda \Sigma$ circuit of size $2^{\omega\left(\sqrt{m} \log ^{\frac{3}{2}} m\right)}$, then VP $\neq \mathbf{V N P}$.

Here are ideas towards the proof: In [GKKS13b] they prove upper bounds for the size of a depth three circuit computing a polynomial, in terms of the size of an arbitrary circuit computing the polynomial. They first apply the work of [Koi, AV08], which allows one to reduce an arbitrary circuit of size $s$ computing a polynomial of degree $d$ in $N$ variables to a formula of size $2^{O(\log s \log d)}$ and depth $d$.

The next step is via the iterated matrix multiplication polynomial. By Theorem 7.3.2.1 formula size is at least as large as iterated matrix multiplication complexity. Say we can compute $f \in S^{m} \mathbb{C}^{M}$ via $m$ matrix multiplications of $n \times n$ matrices with linear entries. (Here $n$ will be comparable to $s$.) Group the entries into groups of $\left\lceil\frac{m}{a}\right\rceil$ for some $a$. To simplify the discussion, assume $\frac{m}{a}$ is an integer. Write

$$
X_{1} \cdots X_{m}=\left(X_{1} \cdots X_{\frac{m}{a}}\right)\left(X_{\frac{m}{a}+1} \cdots X_{2 \frac{m}{a}}\right) \cdots\left(X_{m-\frac{m}{a}+1} \cdots X_{m}\right) .
$$

Each term in parenthesis can be computed (using the naïve matrix multiplication algorithm) via a $\Sigma \Pi^{\frac{m}{a}}$-circuit of size $O\left(n^{\frac{m}{a}}\right)$. After getting the resulting matrices, we can compute the rest via a $\Sigma \Pi^{a}$ circuit of size $O\left(n^{a}\right)$. This reduces one to a depth four circuit of size $S=2^{O(\sqrt{d \log d \log s \log n})}$. Then one can get a depth five powering circuit using (7.1.1). (An alternative, perhaps simpler, proof appears in [Sap, Thm. 5.17].)

The new circuit has size $O(S)$ and is of the form $\Sigma \Lambda \Sigma \Lambda \Sigma$. Finally, they use (6.1.6) to convert the power sums to elementary symmetric functions which keeps the size at $O(S)$ and drops the depth to three.
7.5.3. Geometry and shallow circuits. I first rephrase the depth 3 result:
Proposition 7.5.3.1. [Lan15a] Let $d=N^{O(1)}$ and let $P \in S^{d} \mathbb{C}^{N}$ be a polynomial that can be computed by a circuit of size $s$.

Then $\left[\ell^{n-d} P\right] \in \sigma_{r}\left(C h_{n}\left(\mathbb{C}^{N+1}\right)\right)$ with roughly $r n \sim s^{\sqrt{d}}$, more precisely, $r n=2^{O(\sqrt{d \log (N) \log (d s)})}$.
Corollary 7.5.3.2. $[\mathbf{G K K S 1 3 b}]\left[\ell^{n-m} \operatorname{det}_{m}\right] \in \sigma_{r}\left(C h_{n}\left(\mathbb{C}^{m^{2}+1}\right)\right)$ where $r n=$ $2^{O(\sqrt{m} \log m)}$.

Proof. The determinant admits a circuit of size $m^{4}$, so it admits a $\Sigma \Pi \Sigma$ circuit of size

$$
2^{O\left(\sqrt{m \log (m) \log \left(m * m^{4}\right)}\right)}=2^{O(\sqrt{m} \log m)},
$$

so its padded version lies in $\sigma_{r}\left(C h_{n}\left(\mathbb{C}^{m^{2}+1}\right)\right)$ where $r n=2^{O(\sqrt{m} \log m)}$.
Corollary 7.5.3.3. [GKKS13b] If for all but finitely many $m$ and all $r, n$ with $r n=2^{\sqrt{m} \log (m) \omega(1)}$, one has $\left[\ell^{n-m} \operatorname{perm}_{m}\right] \notin \sigma_{r}\left(C h_{n}\left(\mathbb{C}^{m^{2}+1}\right)\right)$, then there is no circuit of polynomial size computing the permanent, i.e., $\mathbf{V P} \neq$ VNP.

Proof. One just needs to observe that the number of edges in the first layer (which are invisible from the geometric perspective) is dominated by the number of edges in the other layers.

I now reformulate the other shallow circuit results in geometric language. I first give a geometric reformulation of homogeneous $\Sigma \Lambda \Sigma \Lambda \Sigma$ circuits As mentioned previously, the first layer just allows one to work with arbitrary linear forms. The second layer of a $\Sigma \Lambda \Sigma \Lambda \Sigma$ circuit sends a linear form $\ell$ to $\ell^{\delta}$, i.e., it outputs points of $v_{\delta}(\mathbb{P} V)$. The next layer consists of addition gates, outputting sums of $d$-th powers, i.e., points of $\sigma_{r}\left(v_{\delta}(\mathbb{P} V)\right)$. The next layer Veronese re-embeds and multiplies (i.e. projects $S^{\delta^{\prime}}\left(S^{\delta} V\right) \rightarrow S^{\delta \delta^{\prime}} V$ ) these secant varieties to obtain points of $\operatorname{mult}\left(v_{\delta^{\prime}}\left(\sigma_{r}\left(v_{\delta}(\mathbb{P} V)\right)\right)\right)$, and the final addition gate outputs a point of $\sigma_{r^{\prime}}\left(\operatorname{mult}\left(\left(v_{\delta^{\prime}}\left(\sigma_{r}\left(v_{\delta}(\mathbb{P} V)\right)\right)\right)\right)\right.$. In what follows I will simply write $\sigma_{r^{\prime}}\left(v_{\delta^{\prime}}\left(\sigma_{r}\left(v_{\delta}(\mathbb{P} V)\right)\right)\right.$ for this variety. Thus we may rephrase Theorem 7.5.2.1(2) of [GKKS13b] as:
Proposition 7.5.3.4. [Lan15a] Let $d=N^{O(1)}$ and let $P_{N} \in S^{d} \mathbb{C}^{N}$ be a polynomial sequence that can be computed by a circuit of size $s$. Then $\left[P_{N}\right] \in \sigma_{r_{1}}\left(v_{\frac{d}{\delta}}\left(\sigma_{r_{2}}\left(v_{\delta}\left(\mathbb{P}^{N-1}\right)\right)\right)\right)$ with roughly $\delta \sim \sqrt{d}$ and $r_{1} r_{2} \sim s^{\sqrt{d}}$, more precisely $r_{1} r_{2} \delta=2^{O(\sqrt{d \log (d s) \log (N)})}$.
Corollary 7.5.3.5. [Lan15a] If for all but finitely many $m, \delta \simeq \sqrt{m}$, and all $r_{1}, r_{2}$ such that $r_{1} r_{2}=2^{\sqrt{m} \log (m) \omega(1)}$, one has $\left[\operatorname{perm}_{m}\right] \notin \sigma_{r_{1}}\left(v_{m / \delta}\left(\sigma_{r_{2}}\left(v_{\delta}\left(\mathbb{P}^{m^{2}-1}\right)\right)\right)\right)$, then there is no circuit of polynomial size computing the permanent, i.e., VP $\neq$ VNP.
Problem 7.5.3.6. Find equations in the ideal of $\sigma_{r_{1}}\left(v_{\delta}\left(\sigma_{r_{2}}\left(v_{\delta}\left(\mathbb{P}^{m^{2}-1}\right)\right)\right)\right)$.
Y. Guan [Gua15b] has compared the flattening rank of a generic polynomial in $\sigma_{r_{1}}\left(v_{\delta}\left(\sigma_{r_{2}}\left(v_{\delta}\left(\mathbb{P}^{m^{2}-1}\right)\right)\right)\right)$ with that of the permanent and showed that

$$
\operatorname{perm}_{n} \notin \sigma_{2 \sqrt{n} \log (n) \omega(1)}\left(v_{\sqrt{n}}\left(\sigma_{2^{2} \sqrt{n}-\log (n) \omega(1)}\left(v_{\sqrt{n}}\left(\mathbb{P}^{n^{2}-1}\right)\right)\right)\right.
$$

Remark 7.5.3.7. The expected dimension of $\sigma_{r}\left(C h_{m}(W)\right)$ is $r m \mathbf{w}+r-$ 1. If we take $n$ and work instead with padded polynomials $\ell^{n-m} P$, the expected dimension of $\sigma_{r}\left(C h_{n}(W)\right)$ is $r n \mathbf{w}+r-1$. In contrast, the expected dimension of $\sigma_{r}\left(v_{d-a}\left(\sigma_{\rho}\left(v_{a}(\mathbb{P} W)\right)\right)\right)$ does not change when one increases the degree, which indicates why padding is so useful for homogeneous depth three circuits but not for $\Sigma \Lambda \Sigma \Lambda \Sigma$ circuits.

I now describe depth four circuits in terms of joins and multiplicative joins. Following [Lan10], for varieties $X \subset \mathbb{P} S^{a} W$ and $Y \subset \mathbb{P} S^{b} W$, define the multiplicative join of $X$ and $Y, M J(X, Y):=\{[x y] \mid[x] \in X,[y] \in$ $Y\} \subset \mathbb{P} S^{a+b} W$, and define $M J\left(X_{1}, \ldots, X_{k}\right)$ similarly. Let $M J^{k}(X)=$ $\operatorname{MJ}\left(X_{1}, \ldots, X_{k}\right)$ when all the $X_{j}=X$, which is a multiplicative analog of the secant variety. Note that $M J^{k}(\mathbb{P} W)=C h_{k}(W)$. The varieties associated to the polynomials computable by depth $k+1$ formulas are of the form $\sigma_{r_{k}}\left(M J^{d_{k-1}}\left(\sigma_{r_{k-2}}\left(\cdots M J^{d_{1}}(\mathbb{P} W) \cdots\right)\right)\right.$, and $M J^{d_{k}}\left(\sigma_{r_{k-1}}\left(M J^{d_{k-2}}\left(\sigma_{r_{k-3}}\left(\cdots M J^{d_{1}}(\mathbb{P} W) \cdots\right)\right)\right)\right.$ ). In particular, a $\Sigma^{r} \Pi^{\alpha} \Sigma^{s} \Pi^{\beta}$ circuit computes (general) points of $\sigma_{r}\left(M J^{\alpha}\left(\sigma_{s}\left(M J^{\beta}(\mathbb{P} W)\right)\right)\right.$.

### 7.6. Hilbert functions of Jacobian ideals (shifted partial derivatives) and VP v. VNP

The paper [GKKS13a], by Gupta, Kamath, Kayal, and Saptharishi (GKKS) won the best paper award at the 2013 Conference on Computational Complexity (CCC) because it came tantalizingly close to proving Valiant's hypothesis by showing that the permanent does not admit a depth four circuit with top fanin $2^{o(\sqrt{m})}$. Compare this with Theorem 7.5.2.1 that implies to prove VP $\neq \mathbf{V N P}$, it would be sufficient to show that perm ${ }_{m}$ is not computable by a homogeneous $\Sigma \Pi^{O(\sqrt{m})} \Sigma \Pi^{O(\sqrt{m})}$ circuit with top fanin $2^{\Omega(\sqrt{m} \log (m))}$.

The caveat is that in the same paper, they proved the same lower bound for the determinant. On the other hand, a key estimate they use (7.6.6) is close to being sharp for the determinant but conjecturally far from being sharp for the permanent.

Their method of proof is via a classical subject in algebraic geometry: the study of Hilbert functions, and opens the way for using techniques from commutative algebra (study of syzygies) in algebraic complexity theory. I begin, in $\S 7.6 .1$, with a general discussion on the growth of Hilbert functions of ideals. In $\S 7.6 .2$, I outline the proof of the above-mentioned GKKS theorem. In §7.6.3, I show that the shifted partial derivative technique alone cannot separated the determinant from the padded permanent. However, more powerful tools from commutative algebra should be useful for future investigations. With this in mind, in $\S 10.4$, I discuss additional information about the permanent and determinant coming from commutative algebra.
7.6.1. Generalities on Hilbert functions. In what follows we will be comparing the sizes of ideals in judiciously chosen degrees. In this section I explain the fastest and slowest possible growth of ideals generated in a given degree.
Theorem 7.6.1.1 (Macaulay, see, e.g., [Gre98]). Let $\mathcal{I} \subset \operatorname{Sym}\left(\mathbb{C}^{N}\right)$ be a homogeneous ideal, and let $d$ be a natural number. Write

$$
\begin{equation*}
\operatorname{dim} S^{d} \mathbb{C}^{N} / \mathcal{I}_{d}=\binom{a_{d}}{d}+\binom{a_{d-1}}{d-1}+\cdots+\binom{a_{\delta}}{\delta} \tag{7.6.1}
\end{equation*}
$$

with $a_{d}>a_{d-1}>\cdots>a_{\delta}$ (such an expression exists and is unique). Then
$\operatorname{dim} \mathcal{I}_{d+\tau} \geq\binom{ N+d+\tau-1}{d+\tau}-\left[\binom{a_{d}+\tau}{d+\tau}+\binom{a_{d-1}+\tau}{d+\tau-1}+\ldots+\binom{a_{\delta}+\tau}{\delta+\tau}\right]$.
See [Gre98] for a proof.

Corollary 7.6.1.2. Let $\mathcal{I}$ be a homogeneous ideal such that $\operatorname{dim} \mathcal{I}_{d} \geq$ $\operatorname{dim} S^{d-q} \mathbb{C}^{N}=\binom{N+d-q-1}{d-q}$ for some $q<d$. Then

$$
\operatorname{dim} \mathcal{I}_{d+\tau} \geq \operatorname{dim} S^{d-q+\tau} \mathbb{C}^{N}=\binom{N+\tau+d-q-1}{\tau+d-q}
$$

Proof of Corollary. First use the identity

$$
\begin{equation*}
\binom{a+b}{b}=\sum_{j=1}^{q}\binom{a+b-j}{b-j+1}+\binom{a+b-q}{b-q} \tag{7.6.3}
\end{equation*}
$$

with $a=N-1, b=d$. Write this as

$$
\binom{N-1+d}{d}=Q_{d}+\binom{N-1+d-q}{d-q}
$$

Set

$$
Q_{d+\tau}:=\sum_{j=1}^{q}\binom{N-1+d+\tau-j}{d+\tau-j+1}
$$

By Macaulay's theorem, any ideal $\mathcal{I}$ with

$$
\operatorname{dim} \mathcal{I}_{d} \geq\binom{ N-1+d-q}{d-q}
$$

must satisfy

$$
\operatorname{dim} \mathcal{I}_{d+\tau} \geq\binom{ N-1+d+\tau}{d+\tau}-Q_{d+\tau}=\binom{N-1+d-q+\tau}{d-q+\tau}
$$

Gotzman [Got78] showed that if $\mathcal{I}$ is generated in degree at most $d$, then equality is achieved for all $\tau$ in (7.6.2) if equality holds for $\tau=1$. This is the slowest possible growth of an ideal. Ideals satisfying this minimal growth exist. For example, lex-segment ideals satisfy this property, see [Gre98]. These are the ideals, say generated by $K$ elements, where the generators are the first $K$ monomials in lexicographic order. For $1 \leq K \leq M$, the generators are $x_{1}^{d}, x_{1}^{d-1} x_{2}, \ldots, x_{1}^{d-1} x_{K}$. For $M+1 \leq K \leq 2 M$, the generators are $x_{1}^{d-1} x_{j}, x_{1}^{d-2} x_{2} x_{s}, 1 \leq j \leq M, 2 \leq s \leq K-M$, etc...

In general, slow growth occurs because there are syzygies among the generators of the ideal, that is there are relations of the form $P_{1} Q_{1}+\cdots+$ $P_{r} Q_{r}=0$, where $P_{j} \in \mathcal{I}$ and the $Q_{j}$ are polynomials of low degree. For any ideal, one has tautological syzygies, called the Koszul syzygies with $r=2$ and $Q_{1}=P_{2}$ and $Q_{2}=-P_{1}$. Ideals which have only these syzygies grow fast. Explicitly, the fastest possible growth of an ideal generated in degree $d$ by $K<N$ generators is like that of a complete intersection: a variety $X \subset \mathbb{P} V$ of codimension $c$ is a complete intersection if its ideal can be generated by
$c$ elements. The degree $D$ component of an ideal generated in degree $d$ by $K$ generators that grows like a complete intersection ideal has dimension

$$
\begin{equation*}
\sum_{j=1}^{K}(-1)^{j+1}\binom{K}{j}\binom{N+D-j d-1}{D-j d} \tag{7.6.4}
\end{equation*}
$$

Fröberg [Frö85] conjectures ideals with this growth exist even when $K>$ $N$ and Iarrobino $[\mathbf{I a r} \mathbf{9 7}]$ conjectures further that the ideal generated by $\ell_{1}^{d}, \ldots, \ell_{K}^{d}$, with the $\ell_{j}$ general, has this growth (this is known for $K \leq N$ ).
Exercise 7.6.1.3: (2) Prove directly that (7.6.4) holds for an ideal generated by $\ell_{1}^{d}, \ell_{2}^{d}$. ©

The study of the growth of ideals is a classical subject in algebraic geometry. The function $\operatorname{HilbF}_{t}(\mathcal{I}):=\operatorname{dim} \mathcal{I}_{t}$ is called the Hilbert function of the ideal $\mathcal{I} \subset \operatorname{Sym}(V)$.

### 7.6.2. Lower complexity bounds for $\operatorname{perm}_{m}$ (and $\operatorname{det}_{n}$ ) for depth four circuits.

Theorem 7.6.2.1. [GKKS13a] Any $\Sigma \Pi^{O(\sqrt{m})} \Sigma \Pi^{O(\sqrt{m})}$ circuit that computes $\operatorname{perm}_{m}$ or $\operatorname{det}_{m}$ must have top fanin at least $2^{\Omega(\sqrt{m})}$.

In other words $\left[\operatorname{perm}_{m}\right] \notin \sigma_{s}\left(M J^{q}\left(\sigma_{t}\left(M J^{m-q}\left(\mathbb{P}^{m^{2}-1}\right)\right)\right)\right)$, for $s=2^{o(\sqrt{m})}$ and $q=O(\sqrt{m})$. In fact they show $\left[\operatorname{perm}_{m}\right] \notin \sigma_{s}\left(M J^{q}\left(\mathbb{P} S^{m-q} \mathbb{C}^{m^{2}}\right)\right)$.

Recall the Jacobian varieties from §6.3.2. The dimension of $\operatorname{Zeros}(P)_{J a c, k}$ is a measure of the nature of the singularities of $\operatorname{Zeros}(P)$. The proof proceeds by comparing the Hilbert functions of Jacobian varieties.

If $P=Q_{1} \cdots Q_{p}$ is the product of $p$ polynomials, and $k \leq p$, then $Z_{J a c, k}$ will be of codimension at most $k+1$ because it contains $\operatorname{Zeros}\left(Q_{i_{1}}\right) \cap \cdots \cap$ $\operatorname{Zeros}\left(Q_{i_{k+1}}\right)$ for all $\left(i_{1}, \ldots, i_{k+1}\right) \subset[p]$.

Now $\sigma_{s}\left(M J^{q}\left(\mathbb{P} S^{m-q} \mathbb{C}^{m^{2}}\right)\right)$ does not consist of polynomials of this form, but sums of such. With the sum of $m$ such, we can arrive at a smooth hypersurface. So the goal is to find a pathology of $Q_{1} \cdots Q_{p}$ that persists even when taking sums. (The goal is to find something that persists even when taking a sum of $2^{\sqrt{m}}$ such!)

In this situation, the dimension of the space of partial derivatives (rank of the flattenings) is not small enough to prove the desired lower bounds. However, the image of the flattening map will be of a pathological nature, in that all the polynomials in the image are in an ideal generated by a small number of lower degree polynomials. To see this, when $P=Q_{1} \cdots Q_{p}$, with $\operatorname{deg}\left(Q_{j}\right)=q$, any first derivative is in $\sum_{j} S^{q-1} V \cdot\left(Q_{1} \cdots \hat{Q}_{j} \cdots Q_{p}\right)$, where the hat denotes omission. The space of $k$-th derivatives, when $k<p$, is in $\sum_{|J|=k} S^{q-k} V \cdot\left(Q_{1} \cdots \hat{Q}_{j_{1}} \cdots \hat{Q}_{j_{k}} \cdots Q_{p}\right)$. In particular, it has dimension at
most

$$
\begin{equation*}
\binom{p}{k} \operatorname{dim} S^{q-k} V=\binom{p}{k}\binom{\mathbf{v}+q-k-1}{q-k} . \tag{7.6.5}
\end{equation*}
$$

More important than its dimension, is its structure: the ideal it generates, in a given degree $D$ "looks like" the polynomials of degree $D-k$ times a small fixed space of dimension $\binom{p}{k}$.

This behavior is similar to the lex-segment ideals. It suggests comparing the Hilbert functions of the ideal generated by a polynomial computable by a "small" depth four circuit, i.e., of the form $\sum_{j=1}^{s} Q_{1 j} \cdots Q_{p j}$ with the ideal generated by the partial derivatives of the permanent, which are just the sub-permanents. As remarked earlier, even the dimension of the zero set of the size $k$ subpermanents is not known in general. Nevertheless, we just need a lower bound on its growth, which we can obtain by degenerating it to an ideal we can estimate.

First we get an upper bound on the growth of the ideal of the Jacobian variety of $Q_{1} \cdots Q_{m}$ : By the discussion above, in degree $m-k+\tau$ it has dimension at most

$$
\binom{p}{k} \operatorname{dim} S^{q-k+\tau} V=\binom{p}{k}\binom{\mathbf{v}+\tau+q-k-1}{q-k} .
$$

To get the lower bound on the growth of the ideal generated by subpermanents we use a crude estimate: given a polynomial $f$ given in coordinates, its leading monomial in some order (say lexicographic), is the monomial in its expression that is highest in the order. So if an ideal is generated by $f_{1}, \ldots, f_{q}$ in degree $d$, then in degree $d+\tau$, it is of dimension at most the number of monomials in degree $d+\tau$ divisible by a leading monomial from one of the $f_{j}$.

If we order the variables in $\mathbb{C}^{m^{2}}$ by $y_{1}^{1}>y_{2}^{1}>\cdots>y_{m}^{1}>y_{1}^{2}>\cdots>y_{m}^{m}$, then the leading monomial of any sub-permanent is the product of the elements on the principal diagonal. Even working with this, the estimate is difficult, so in [GKKS13a] they restrict further to only look at leading monomials among the variables on the diagonal and super diagonal: $\left\{y_{1}^{1}, \ldots, y_{m}^{m}, y_{2}^{1}, y_{3}^{2}, \ldots, y_{m}^{m-1}\right\}$. Among these, they compute that the number of leading monomials of degree $\delta$ is $\binom{2 m-\delta}{\delta}$. In our case, $\delta=m-k$ and $D=\tau+m-k$. Let $I_{d}^{\text {perm }_{m}, k} \subset S^{d} \mathbb{C}^{m^{2}}$ denote the degree $d$ component of the ideal generated by the order $k$ partial derivatives of the permanent, i.e., the $k$-th Jacobian variety of $\operatorname{perm}_{m}$. In [GKKS13a], $I_{d}^{\text {perm }_{m}, k}$ is denoted $\left\langle\partial^{=k} \operatorname{perm}_{m}\right\rangle_{=d-m}$. We have

$$
\begin{equation*}
\operatorname{dim} I_{m-k+\tau}^{\mathrm{perm}_{m}, k} \geq\binom{ m+k}{2 k}\binom{m^{2}+\tau-2 k}{\tau} \tag{7.6.6}
\end{equation*}
$$

and

$$
\begin{equation*}
\operatorname{dim} I_{m-k+\tau}^{\operatorname{det}_{m}, k} \geq\binom{ m+k}{2 k}\binom{m^{2}+\tau-2 k}{\tau} . \tag{7.6.7}
\end{equation*}
$$

Putting the estimates together, if we want to realize the permanent by size $s \Sigma \Pi^{O(\sqrt{m})} \Sigma \Pi^{O(\sqrt{m})}$ circuit, we need

$$
\begin{equation*}
s \geq \frac{\binom{m+k}{2 k}\binom{m^{2}+\tau-2 k}{\tau}}{\binom{c \sqrt{m}+k}{k}\binom{m^{2}+\tau+(\sqrt{m}-1) k}{m^{2}}} \tag{7.6.8}
\end{equation*}
$$

Theorem 7.6.2.1 follows by setting $\tau=m^{\frac{5}{2}}$ and $k=\epsilon m^{\frac{1}{2}}$ where $\epsilon$ is a constant defined below. To see this, one calculates (using the estimates of §7.5.1):

$$
\begin{aligned}
& \ln \frac{\binom{m^{2}+m^{\frac{5}{2}}-2 \epsilon \sqrt{m}}{m^{\frac{5}{2}}}}{\left(\begin{array}{c}
m^{2}+m^{\frac{5}{2}}+(\sqrt{m}-1) \epsilon \sqrt{m} \\
m^{2}
\end{array}\right.}=-2 \epsilon \sqrt{m} \ln \sqrt{m}-\epsilon \sqrt{m} \pm O(1) \\
& \ln \frac{\binom{m^{2}+\epsilon \sqrt{m}}{2 \epsilon \sqrt{m}}}{\binom{(c+\epsilon) \sqrt{m}}{\epsilon \sqrt{m}}}=\sqrt{m} 2 \epsilon \ln \frac{\sqrt{m}}{2 \epsilon}+2 \epsilon \\
& +(c+\epsilon)\left[\frac{\epsilon}{c+\epsilon} \ln \left(\frac{\epsilon}{c+\epsilon}\right)+\left(1-\frac{\epsilon}{c+\epsilon}\right) \ln \left(1-\frac{\epsilon}{c+\epsilon}\right)\right]+O(\ln m)
\end{aligned}
$$

These imply

$$
\ln (s) \geq \epsilon \sqrt{m} \ln \frac{1}{4 \epsilon(c+\epsilon)} \pm O(1)
$$

so choosing $\epsilon$ such that $\frac{1}{4 \epsilon(c+\epsilon)}=e$, yields $\ln (s) \geq \Omega(\sqrt{m})$.
7.6.3. Shifted partial derivatives cannot separate permanent from determinant. Recall the notations for a polynomial $P \in S^{n} V$, that $I_{d}^{P, k}=$ $\left\langle\partial^{=k} P\right\rangle_{=d-n}$ is the degree $d$ component of the ideal generated by the order $k$ partial derivatives of $P$, i.e., the degree $d$ component of the ideal of the $k$-th Jacobian variety of $P$.
Theorem 7.6.3.1. [ELSW16] There exists a constant $M$ such that for all $m>M$, every $n>2 m^{2}+2 m$, any $\tau$, and any $k<n$,

$$
\operatorname{dim} I_{n+\tau}^{\ell^{n-m} \operatorname{perm}_{m}, k}<\operatorname{dim} I_{n+\tau}^{\operatorname{det}_{n}, k} .
$$

In other words

$$
\operatorname{dim}\left\langle\partial^{=k}\left(\ell^{n-m} \operatorname{perm}_{m}\right)\right\rangle_{=\tau}<\operatorname{dim}\left\langle\partial^{=k} \operatorname{det}_{n}\right\rangle_{=\tau} .
$$

The proof of Theorem 7.6.3.1 splits into four cases:

- (C1) Case $k \geq n-\frac{n}{m+1}$. This case has nothing to do with the padded permanent or its derivatives: the estimate is valid for any polynomial in $m^{2}+1$ variables.
- (C2) Case $2 m \leq k \leq n-2 m$. This case uses that when $k<n-m$, the Jacobian ideal of any padded polynomial $\ell^{n-m} P \in S^{n} W$ is contained in the ideal generated in degree $n-m-k$ by $\ell^{n-m-k}$ which has slowest growth by Macaulay's theorem.
- (C3) Case $k<2 m$ and $\tau>\frac{3}{2} n^{2} m$. This case is similar to case C2, only a degeneration of the determinant is used in the comparison.
- (C4) Case $k<2 m$ and $\tau<\frac{n^{3}}{6 m}$. This case uses (7.6.7) and compares it with a very crude upper bound for the dimension of the space of the shifted partial derivatives of the permanent.

Note that C1, C2 overlap when $n>2 m^{2}+2 m$ and C3, C4 overlap when $n>\frac{m^{2}}{4}$, so it suffices to take $n>2 m^{2}+2 m$.
Case C1. The assumption is $(m+1)(n-k) \leq n$. It will be sufficient to show that some $R \in \operatorname{End}(W) \cdot \operatorname{det}_{n}$ satisfies $\operatorname{dim} I_{n-k+\tau}^{\ell^{n-m} \operatorname{perm}_{m}, k}<\operatorname{dim} I_{n-k+\tau}^{R, k}$. Block the matrix $x=\left(x_{u}^{s}\right) \in \mathbb{C}^{n^{2}}$, with $1 \leq s, u \leq n$, as a union of $n-k$ blocks of size $m \times m$ in the upper-left corner plus the remainder, which by our assumption includes at least $n-k$ elements on the diagonal. Set each diagonal block to the matrix $\left(y_{j}^{i}\right)$, with $1 \leq i, j \leq n$, (there are $n-k$ such blocks), fill the remainder of the diagonal with $\ell$ (there are at least $n-k$ such terms), and fill the remainder of the matrix with zeros. Let $R$ be the restriction of the determinant to this subspace. Then the space of partials of $R$ of degree $n-k, I_{n-k}^{R, k} \subset S^{n-k} \mathbb{C}^{n^{2}}$ contains a space of polynomials isomorphic to $S^{n-k} \mathbb{C}^{m^{2}+1}$, and $I_{n-k}^{\ell^{n-m} \text { perm }_{m}, k} \subset S^{n-k} \mathbb{C}^{m^{2}+1}$ so we conclude.

Example 7.6.3.2. Let $m=2, n=6, k=4$. The matrix is

$$
\left(\begin{array}{cccccc}
y_{1}^{1} & y_{2}^{1} & & & & \\
y_{1}^{2} & y_{2}^{2} & & & & \\
& & y_{1}^{1} & y_{2}^{1} & & \\
& & y_{1}^{2} & y_{2}^{2} & & \\
& & & & \ell & \\
& & & & & \ell
\end{array}\right)
$$

The polynomial $\left(y_{1}^{1}\right)^{2}$ is the image of $\frac{\partial^{4}}{\partial x_{2}^{2} \partial x_{4}^{4} \partial x_{5}^{5} \partial x_{6}^{6}}$ and the polynomial $y_{2}^{1} y_{2}^{2}$ is the image of $\frac{\partial^{4}}{\partial x_{1}^{2} \partial x_{3}^{3} \partial x_{5}^{5} \partial x_{6}^{6}}$.

Case C2. As long as $k<n-m, I_{n-k}^{\ell^{n-m} \operatorname{perm}_{m}, k} \subset \ell^{n-m-k} \cdot S^{m} W$, so

$$
\begin{equation*}
\operatorname{dim} I_{n-k+\tau}^{\ell^{n-m} \operatorname{perm}_{m}, k} \leq\binom{ n^{2}+m+\tau-1}{m+\tau} \tag{7.6.9}
\end{equation*}
$$

By Corollary 7.6.1.2, with $N=n^{2}, d=n-k$, and $d-q=m$, it will be sufficient to show that

$$
\begin{equation*}
\operatorname{dim} I_{n-k}^{\operatorname{det}_{n}, k}=\binom{n}{k}^{2} \geq \operatorname{dim} S^{m} W=\binom{n^{2}+m-1}{m} \tag{7.6.10}
\end{equation*}
$$

In the range $2 m \leq k \leq n-2 m$, the quantity $\binom{n}{k}$ is minimized at $k=2 m$ and $k=n-2 m$, so it is enough to show that

$$
\begin{equation*}
\binom{n}{2 m}^{2} \geq\binom{ n^{2}+m-1}{m} \tag{7.6.11}
\end{equation*}
$$

The estimates of $\S 7.5 .1$ show that this holds when $\left(\frac{n}{2 m}-1\right)^{4}>\left(\frac{n^{2}}{m}-\frac{m-1}{m}\right)$ which holds for all sufficiently large $m$ when $n>m^{2}$.
Case C3. Here simply degenerate $\operatorname{det}_{n}$ to $R=\ell_{1}^{n}+\ell_{2}^{n}$ by e.g., setting all diagonal elements to $\ell_{1}$, all the sub-diagonal elements to $\ell_{2}$ as well as the $(1, n)$-entry, and setting all other elements of the matrix to zero. Then $I_{n-k}^{R, k}=\operatorname{span}\left\{\ell_{1}^{n-k}, \ell_{2}^{n-k}\right\}$. Since this is a complete intersection ideal,

$$
\begin{equation*}
\operatorname{dim} I_{n-k+\tau}^{R, k}=2\binom{n^{2}+\tau-1}{\tau}-\binom{n^{2}+\tau-(n-k)-1}{\tau-(n-k)} \tag{7.6.12}
\end{equation*}
$$

Using the estimate (7.6.9) from Case C 2 , it remains to show

$$
2\binom{n^{2}+\tau-1}{\tau}-\binom{n^{2}+\tau+m-1}{\tau+m}-\binom{n^{2}+\tau-(n-k)-1}{\tau-(n-k)}>0
$$

Divide by $\binom{n^{2}+\tau-1}{\tau}$. We need

$$
\begin{align*}
2> & \prod_{j=1}^{m} \frac{n^{2}+\tau+m-j}{\tau+m-j}+\prod_{j=1}^{n-k} \frac{\tau-j}{n^{2}+\tau-j}  \tag{7.6.13}\\
& =\prod_{j=1}^{m}\left(1+\frac{n^{2}}{\tau+m-j}\right)+\prod_{j=1}^{n-k}\left(1-\frac{n^{2}}{n^{2}+\tau-j}\right) . \tag{7.6.14}
\end{align*}
$$

The second line is less than

$$
\begin{equation*}
\left(1+\frac{n^{2}}{\tau}\right)^{m}+\left(1-\frac{n^{2}}{n^{2}+\tau-1}\right)^{n-k} \tag{7.6.15}
\end{equation*}
$$

Consider (7.6.15) as a function of $\tau$. Write $\tau=n^{2} m \delta$, for some constant $\delta$. Then (7.6.15) is bounded above by

$$
e^{\frac{1}{\delta}}+e^{\frac{2}{\delta}-\frac{n}{m \delta}} .
$$

The second term goes to zero for large $m$, so we just need the first term to be less than 2 , so take, e.g., $\delta=\frac{3}{2}$.

Case C4. Compare (7.6.7) with the very crude estimate

$$
\operatorname{dim} I_{n-k+\tau}^{\ell^{n-m} \operatorname{perm}_{m}, k} \leq \sum_{j=0}^{k}\binom{m}{j}^{2}\binom{n^{2}+\tau-1}{\tau}
$$

where $\sum_{j=0}^{k}\binom{m}{j}^{2}$ is the dimension of the space of partials of order $k$ of $\ell^{n-m} \operatorname{perm}_{m}$, and the $\binom{n^{2}+\tau-1}{\tau}$ is what one would have if there were no syzygies. One then concludes using the estimates of $\S 7.5 .1$, although it is necessary to split the estimates into two sub-cases: $k \geq \frac{m}{2}$ and $k<\frac{m}{2}$. See [ELSW16] for details.

### 7.7. Polynomial identity testing, hitting sets and explicit Noether normalization

I give an introduction to the three topics in the section title. Hitting sets are defined with respect to a coordinate system, however they reflect geometry that is independent of coordinates that merits further study.

For simplicity, I work with homogeneous polynomials, and continue to work exclusively over $\mathbb{C}$.
7.7.1. PIT. If someone hands you a homogeneous polynomial, given in terms of a circuit, or in terms of a sequence of symmetrizations and skewsymmetrizations (as often happens in representation theory), how can you test if it is identically zero?

I will only discuss "black box" polynomial identity testing (henceforth PIT), where one is only given the output of the circuit, as opposed to "white box" PIT where the structure of the circuit may also be examined.

Consider the simplest case: say you are told the polynomial in $N$ variables is linear. Then it suffices to test it on $N$ points in general linear position in $\mathbb{P}^{N-1}$. Similarly, if we have a conic the projective plane, six general points suffice to test if the conic is zero (and given six points, it is easy to test if they are general enough).

Any $P \in S^{d} \mathbb{C}^{2}$ vanishing on any $d+1$ distinct points in $\mathbb{P}^{1}$ is identically zero. More generally, for $P \in S^{d} \mathbb{C}^{N},\binom{N+d-1}{d}$ sufficiently general points in $\mathbb{P}^{N-1}$ suffice to test if $P$ is zero. If $N, d$ are large, this is not feasible. Also, it is not clear how to be sure points are sufficiently general. Fortunately, for a small price, we have the following Lemma, which dates back at least to [AT92], addressing the "sufficiently general" issue:
Lemma 7.7.1.1. Let $\Sigma$ be a collection of $d+1$ distinct nonzero complex numbers, let $\Sigma^{N}=\left\{\left(c_{1}, \ldots, c_{N}\right) \mid c_{i} \in \Sigma\right\}$. Then any $P \in S^{d} \mathbb{C}^{N}$ vanishing on $\Sigma^{N}$ is identically zero.

Proof. Work by induction on $N$, the case $N=1$ is clear. Write $P$ as a polynomial in $x_{N}: P\left(x_{1}, \ldots, x_{N}\right)=\sum_{j=0}^{d} P_{j}\left(x_{1}, \ldots, x_{N-1}\right) x_{N}^{j}$ where $P_{j} \in$ $S^{d-j} \mathbb{C}^{N-1}$ and assume $P$ vanishes on $\Sigma^{N}$. For each $\left(c_{1}, \ldots, c_{N-1}\right) \in \Sigma^{N-1}$, $P\left(c_{1}, \ldots, c_{N-1}, x_{N}\right)$ is a polynomial in one variable vanishing on $\Sigma^{1}$, and is therefore identically zero. Thus each $P_{j}$ vanishes identically on $\Sigma^{N-1}$ and by induction is identically zero.

Now say we are given a polynomial with extra structure and we would like to exploit the structure to determine if it is non-zero using a smaller set of points than $\Sigma^{N}$. If we are told it is a $d$-th power (or zero), then its zero set is simply a hyperplane, so $N+1$ points in general linear position again suffice. Now say we are told it has low Waring rank. How could we exploit that information to find a small set of points to test?

### 7.7.2. Hitting sets.

Definition 7.7.2.1. (see, e.g, $[\mathbf{S Y 0 9}, \S 4.1])$ Given a subset $\mathcal{C} \subset \mathbb{C}\left[x_{1}, x_{2}, \ldots, x_{N}\right]$, a finite subset $\mathcal{H} \subset \mathbb{C}^{N}$ is a hitting set for $\mathcal{C}$ if for all nonzero $f \in \mathcal{C}$, there exists $\alpha \in \mathcal{H}$ such that $f(\alpha) \neq 0$.

Lemma 7.7.1.1 provides an explicit, size $(d+1)^{N}$ hitting set for $S^{d} \mathbb{C}^{N}$. Call this the naïve $(d, N)$-hitting set.

In geometric language, a hitting set is a subset $\mathcal{H} \subset \mathbb{C}^{N}$ such that the evaluation map eval $\mathcal{H}_{\mathcal{H}}: \mathbb{C}^{\binom{n+d}{d}} \rightarrow \mathbb{C}^{|\mathcal{H}|}$ satisfies eval $_{\mathcal{H}}{ }^{-1}(0) \cap \mathcal{C}=0$.

Existence of a hitting set implies black box PIT via the evaluation map.
Lemma 7.7.2.2. [HS82] There exist hitting sets for

$$
\mathcal{C}_{s}:=\left\{f \in S^{d} \mathbb{C}^{n} \mid \exists \text { a size } s \text { circuit computing } f\right\},
$$

with size bounded by a polynomial in $s, d$ and $n$.
7.7.3. Waring rank. Returning to the problem of finding an explicit (in the computer science sense, see §6.1.3) hitting set for polynomials of small Waring rank, recall that we do not know defining equations for $\sigma_{r}\left(v_{d}\left(\mathbb{P}^{N-1}\right)\right)$, however, we do have some equations, at least as long as $r<\binom{N+\left\lfloor\frac{d}{2}\right\rfloor-1}{\left\lfloor\frac{d}{2}\right\rfloor}$, namely the flattenings. So it is easier to change the question: we simply look for a hitting set for the larger variety Flat $^{\left\lfloor\frac{d}{2}\right\rfloor,\left\lceil\frac{d}{2}\right\rceil}\left(S^{d} V\right):=\{P \in$ $\left.S^{d} V \left\lvert\, \operatorname{rank}\left(P_{\left\lfloor\frac{d}{2}\right\rfloor,\left\lceil\frac{d}{2}\right\rceil}\right) \leq r\right.\right\}$, where $P_{\left\lfloor\frac{d}{2}\right\rfloor,\left\lceil\frac{d}{2}\right\rceil}: S^{\left\lfloor\frac{d}{2}\right\rfloor} V^{*} \rightarrow S^{\left\lceil\frac{d}{2}\right\rceil} V$ is the partial derivative map. We have a considerable amount of information about Flat ${ }^{\left\lfloor\frac{d}{2}\right\rfloor,\left\lceil\frac{d}{2}\right\rceil}\left(S^{d} V\right)$.

Consider the case $r=2$ : our polynomial is of the form $P=\left(\lambda_{1} x_{1}+\right.$ $\left.\cdots \lambda_{N} x_{N}\right)^{d}+\left(\mu_{1} x_{1}+\cdots \mu_{N} x_{N}\right)^{d}$, for some $\lambda_{i}, \mu_{j} \in \mathbb{C}$. It is no longer sufficient
to check on the $N$ coordinate points, as it could be that $\lambda_{j}^{d}+\mu_{j}^{d}=0$ for all $j$ but $P$ is nonzero. On the other hand, there cannot be too much "interference": restrict to the $\binom{N}{2}$ coordinate $\mathbb{P}^{1}$ 's: it is straightforward to see that if all these restrictions are identically zero, then the polynomial must be zero. Moreover, each of those restrictions can be tested to be zero by just checking on $d+1$ points on a line. Rephrasing geometrically: no point of $\sigma_{2}\left(v_{d}\left(\mathbb{P}^{N-1}\right)\right)$ has a zero set that contains the $\binom{N}{2} \mathbb{P}^{1}$ 's spanned by pairs of points from any collection of $N$ points that span $\mathbb{P}^{N-1}$. (Compare with Lemma 2.6.2.1.) In contrast, consider $\ell_{1} \cdots \ell_{d}=0$ : it contains $d$ hyperplanes!

The general idea is that, if the flattening rank of $P$ is small and $P$ is not identically zero, then for some "reasonably small" $k$, there cannot be a collection of $\binom{N}{k} \mathbb{P}^{k-1}$, s spanned by $k$ subsets of any set of $N$ points spanning $\mathbb{P}^{N-1}$ in the zero set of $P$. In coordinates, this means there is a monomial in the expression for $P$ that involves at most $k$ variables, so it will suffice to restrict $P$ to each of the $\binom{N}{k}$ coordinate subspaces and test these restrictions on a naïve ( $d, k$ )-hitting set.

From the example of $P=\ell_{1} \cdots \ell_{d}$, we see that "small" means at least that $r<\binom{N}{\left.\left\lfloor\frac{d}{2}\right\rfloor\right\rfloor}$. In [FS13a], they show that we may take $k=\log (r)$. (Note that if $r$ is close to $2^{N}$, the assertion becomes vacuous as desired.) Explicitly: Theorem 7.7.3.1. [FS13a] Let $\mathcal{H}$ consist of the $(d+1)^{k}\binom{N}{k}$ points of naïve hitting sets on each coordinate $\mathbb{P}^{k-1}$. Then $\mathcal{H}$ is an explicit hitting set for $\left\{P \in S^{d} \mathbb{C}^{N} \left\lvert\, \operatorname{rank}\left(P_{\left\lfloor\frac{d}{2}\right\rfloor,\left\lceil\frac{d}{2}\right\rceil}\right)<2^{k}\right.\right\}$, in particular for points of $\sigma_{2^{k}}\left(v_{d}\left(\mathbb{P}^{\mathbb{N}-1}\right)\right)$.

An even better hitting set is given in [FSS13].
Recall ROABP's from §7.3.4.
Theorem 7.7.3.2. $[\mathbf{F S 1 3 b}]$ Let $\mathcal{C} \subset \mathbb{C}\left[x_{1}, \ldots, x_{n}\right]$ denote the set of polynomials computable by a depth $n$, width at most $w$, degree at most $r$ ROABP. Then $\mathcal{C}$ has an explicit hitting set $\mathcal{H}$ of size poly $(n, w, r)^{O(\log (n))}$ (quasipolynomial size). Furthermore, one can take $\mathcal{H} \subset \mathbb{Q}^{n}$.
7.7.4. Efficient Noether normalization. One of the difficulties in understanding $\mathcal{D e t}_{n} \subset \mathbb{P} S^{n} \mathbb{C}^{n^{2}}$ is that its codimension is of size exponential in $n$. It would be desirable to have a subvariety of at worst polynomial codimension to work with, as then one could use additional techniques to study its coordinate ring. If one is willing to put up with the destruction of external symmetry, one might simply take a linear projection of $\mathcal{D e} t_{n}$ to a small ambient space. By Noether-Normalization §3.1.4, we know that a "random" linear space of codimension, say $2 n^{4}$ would give rise to an isomorphic projection. However what one would need is an explicit such linear space. In [Mul12] Mulmuley considers this problem of explicitness, in the
context of separating points of an affine variety (described below) via a small subset of its coordinate ring. Call $\mathcal{S} \subset \mathbb{C}[X]$ a separating subset if for all $x, y \in X$, there exists $f \in S$ that distinguishes them, i.e., $f(x) \neq f(y)$.

Remark 7.7.4.1. In [Mul12] the desired linear projection is referred to as a "normalizing map", which is potentially confusing to algebraic geometers because it is not a normalization of the image variety.

Consider $\operatorname{End}\left(\mathbb{C}^{m}\right)^{\oplus r}$, which is an $S L_{m}$-variety under the diagonal action. Write $\bar{A}=A_{1} \oplus \cdots \oplus A_{r} \in \operatorname{End}\left(\mathbb{C}^{m}\right)^{\oplus r}$.
Theorem 7.7.4.2. [Raz74, Pro76] [First fundamental theorem for matrix invariants $] \mathbb{C}\left[\operatorname{End}\left(\mathbb{C}^{m}\right)^{\oplus r}\right]^{S L_{m}}$ is generated by

$$
T_{\alpha}(\bar{A}):=\operatorname{trace}\left(A_{\alpha_{1}} \cdots A_{\alpha_{k}}\right), k \leq m^{2}, \alpha_{1}, \ldots, \alpha_{k} \in[r] .
$$

It is also known that if one takes $k \leq\left\lfloor\frac{m^{2}}{8}\right\rfloor$ one does not obtain generators. In particular, one has an exponentially large (with respect to $m$ ) number of generators.

Put all these polynomials together in a generating function: let $y=\left(y_{j}^{s}\right)$, $1 \leq j \leq m^{2}, 1 \leq s \leq r$ and define

$$
T(y, \bar{A}):=\operatorname{trace}\left[\left(\operatorname{Id}+y_{1}^{1} A_{1}+\cdots+y_{1}^{k} A_{k}\right) \cdots\left(\operatorname{Id}+y_{m^{2}}^{1} A_{1}+\cdots+y_{m^{2}}^{k} A_{k}\right)\right]
$$

The right hand side is an IMM, even a ROABP. Thus all the generating invariants may be read off as the coefficients of the output of an ROABP. This, combined with Theorem 7.7.3.2 implies:
Theorem 7.7.4.3. [FS13a] There exists a poly $(n, r)^{O(\log (n))}$-sized set $\mathcal{H} \subset$ $\mathbb{C}\left[\operatorname{End}\left(\mathbb{C}^{m}\right)^{\oplus r}\right]^{S L_{m}}$ of separating invariants, with poly $(n, r)$-explicit ABP's. In other words, for any $\bar{A}, \bar{B} \in \operatorname{End}\left(\mathbb{C}^{m}\right)^{\oplus r}$, there exists $f \in \mathbb{C}\left[\operatorname{End}\left(\mathbb{C}^{m}\right)^{\oplus r}\right]^{S L_{m}}$ with $f(\bar{A}) \neq f(\bar{B})$ if and only if there exists such $f \in \mathcal{H}$.

Remark 7.7.4.4. A more geometric way of understanding Theorem 7.7.4.3 is to introduce the GIT-quotient $\operatorname{End}\left(\mathbb{C}^{m}\right)^{\oplus r} / / S L_{m}$ (see $\S 9.5 .2$ ), which is an affine algebraic variety whose coordinate ring is $\mathbb{C}\left[\operatorname{End}\left(\mathbb{C}^{m}\right)^{\oplus r}\right]^{S L_{m}}$. Then $\mathcal{H}$ is a subset that separates points of the GIT-quotient $\operatorname{End}\left(\mathbb{C}^{m}\right)^{\oplus r} / / S L_{m}$.

The following conjecture appeared in the 2012 version of [Mul12]:
Conjecture 7.7.4.5. [Mul12] Noether normalization can be performed explicitly for $\operatorname{End}\left(\mathbb{C}^{m}\right)^{\oplus r} / / S L_{m}$ in polynomial time.

Conjecture 7.7.4.5 motivated the work of [FS13a], as Theorem 7.7.4.3 implies:
Corollary 7.7.4.6. [FS13a] Noether normalization can be performed explicitly for $\operatorname{End}\left(\mathbb{C}^{m}\right)^{\oplus r} / / S L_{m}$ in quasi-polynomial time.

Remark 7.7.4.7. The PIT problem is the word problem for the field of rational functions over a set of commuting variables. One can ask the same for the (free) skew field over non-commuting variables. This is answered in [GGOW15] where there are connections to and implications for many areas including PIT, quivers and GIT questions.

# Chapter 8 

## Representation theory and its uses in complexity theory

In this chapter I derive the representation theory of the general linear group $G L(V)$ and give numerous applications to complexity theory. In order to get to the applications as soon as possible, I summarize basic facts about representations of $G L(V)$ in $\S 8.1$. The first application, in $\S 8.2$, explains the theory of Young flattenings underlying the equations that led to the $2 \mathbf{n}^{2}-\mathbf{n}$ lower bound for the border rank of matrix multiplication (Theorem 2.5.2.6). I also explain how the method of shifted partial derivatives may be viewed as a special case of Young flattenings. Next, in $\S 8.3$, I briefly discuss how representation theory has been used to find equations for secant varieties of Segre varieties and other varieties. In $\S 8.4$, I describe severe restrictions on modules of polynomials to be useful for the permanent v . determinant problem. In $\S 8.5$, I give the proofs of several statements about $\mathcal{D e} t_{n}$ from Chapter 7. In $\S 8.6$, I begin to develop representation theory via the double commutant theorem, the algebraic Peter-Weyl theorem and Schur-Weyl duality. The reason for this choice of development is that the (finite) PeterWeyl theorem is the starting point of the Cohn-Umans program of $\S 3.5$ and the algebraic Peter-Weyl theorem was the starting point of the program of [MS01, MS08] described in §8.8. The representations of the general linear group are then derived in $\S 8.7$. In $\S 8.8$ I begin a discussion of the program of [MS01, MS08], as refined in [BLMW11], to separate the permanent from the determinant via representation theory. This is continued in $\S 8.9$, which contains a general discussion of plethysm coefficients, and $\S 8.10$, which
presents results of $[\mathbf{I P} 15]$ and $[\mathbf{B I P} 16]$ that show this program cannot work as stated. I then, in $\S 8.11$ outline the proofs of Theorems 7.4.1.1 and 7.4.1.4 regarding equivariant determinantal expressions for the permanent. I conclude, in $\S 8.12$ with additional theory how to determine symmetry groups of polynomials and illustrate the theory with several examples relevant for complexity theory.

### 8.1. Representation theory of the general linear group

Irreducible representations of $G L(V)$ in $V^{\otimes d}$ are indexed by partitions of $d$ with length at most $\mathbf{v}$, as we will prove in Theorem 8.7.1.2. Let $S_{\pi} V$ denote the isomorphism class of the irreducible representation associated to the partition $\pi$, and let $S_{\bar{\pi}} V$ denote some particular realization of $S_{\pi} V$ in $V^{\otimes d}$. In particular $S_{(d)} V=S^{d} V$ and $S_{(1, \ldots, 1)} V=\Lambda^{d} V$ where there are $d$ 1's. For a partition $\pi=\left(p_{1}, \ldots, p_{k}\right)$, write $|\pi|=p_{1}+\cdots+p_{k}$ and $l(\pi)=k$. If a number is repeated I sometimes use superscripts to record its multiplicity, for example $(2,2,1,1,1)=\left(2^{2}, 1^{3}\right)$.

To visualize $\pi$, define a Young diagram associated to a partition $\pi$ to be a collection of left-aligned boxes with $p_{j}$ boxes in the the $j$-th row, as in Figure 8.1.1.


Figure 8.1.1. Young diagram for $\pi=(4,2,1)$

Define the conjugate partition $\pi^{\prime}$ to $\pi$ to be the partition whose Young diagram is the reflection of the Young diagram of $\pi$ in the north-west to south-east diagonal.


Figure 8.1.2. Young diagram for $\pi^{\prime}=(3,2,1,1)$, the conjugate partition to $\pi=(4,2,1)$.
8.1.1. Lie algebras. Associated to any Lie group $G$ is a Lie algebra $\mathfrak{g}$, which is a vector space that may be identified with $T_{\mathrm{Id}} G$. For basic information on Lie algebras associated to a Lie group, see any of [Spi79, IL16b, Pro07].

When $G=G L(V)$, then $\mathfrak{g}=\mathfrak{g l}(V):=V^{*} \otimes V$. If $G \subseteq G L(V)$, so that $G$ acts on $V^{\otimes d}$, there is an induced action of $\mathfrak{g} \subseteq \mathfrak{g l}(V)$ given by, for $X \in \mathfrak{g}$,

$$
\begin{aligned}
& X .\left(v_{1} \otimes v_{2} \otimes \cdots \otimes v_{d}\right) \\
& \quad=\left(X . v_{1}\right) \otimes v_{2} \otimes \cdots \otimes v_{d}+v_{1} \otimes\left(X . v_{2}\right) \otimes \cdots \otimes v_{d}+\cdots+v_{1} \otimes v_{2} \otimes \cdots \otimes v_{d-1} \otimes\left(X . v_{d}\right) .
\end{aligned}
$$

To see why this is a natural induced action, consider a curve $g(t) \subset G$ with $g(0)=\mathrm{Id}$ and $X=g^{\prime}(0)$ and take

$$
\left.\frac{d}{d t}\right|_{t=0} g(t) \cdot\left(v_{1} \otimes \cdots \otimes v_{d}\right)=\left.\frac{d}{d t}\right|_{t=0}\left(g(t) \cdot v_{1}\right) \otimes \cdots \otimes\left(g(t) \cdot v_{d}\right) .
$$

One concludes by applying the Leibnitz rule.
8.1.2. Weights. Fix a basis $e_{1}, \ldots, e_{\mathbf{v}}$ of $V$, let $T \subset G L(V)$ denote the subgroup of diagonal matrices, called a maximal torus, let $B \subset G L(V)$ be the subgroup of upper triangular matrices, called a Borel subgroup, and let $N \subset B$ be the upper triangular matrices with 1's along the diagonal. The Lie algebra $\mathfrak{n}$ of $N$ consists of nilpotent matrices. Call $z \in V^{\otimes d}$ a weight vector if $T[z]=[z]$. If

$$
\left(\begin{array}{ccc}
x_{1} & & \\
& \ddots & \\
& & x_{\mathbf{v}}
\end{array}\right) z=\left(x_{1}\right)^{p_{1}} \cdots\left(x_{\mathbf{v}}\right)^{p_{\mathbf{v}}} z
$$

we say $z$ has weight $\left(p_{1}, \ldots, p_{\mathbf{v}}\right) \in \mathbb{Z}^{\mathbf{v}}$.
Call $z$ a highest weight vector if $B[z]=[z]$, i.e., if $N z=z$. If $M$ is an irreducible $G L(V)$-module and $z \in M$ is a highest weight vector, call the weight of $z$ the highest weight of $M$. A necessary condition for two irreducible $G L(V)$-modules to be isomorphic is that they have the same highest weight (because they must also be isomorphic $T$-modules). The condition is also sufficient, see $\S 8.7$.
Exercise 8.1.2.1: (1) Show that $z$ is a highest weight vector if and only if $\mathfrak{n} . z=0$.

The elements of $\mathfrak{n}$ are often called raising operators.
Exercise 8.1.2.2: (1) Show that if $z \in V^{\otimes d}$ is a highest weight vector of weight $\left(p_{1}, \ldots, p_{\mathbf{v}}\right)$, then $\left(p_{1}, \ldots, p_{\mathbf{v}}\right)$ is a partition of $d$. ©

When $G=G L\left(A_{1}\right) \times \cdots \times G L\left(A_{n}\right)$, the maximal torus in $G$ is the product of the maximal tori in the $G L\left(A_{j}\right)$, and similarly for the Borel. A weight is then defined to be an $n$-tuple of weights etc...

Because of the relation with weights, it will often be convenient to add a string of zeros to a partition to make it a string of $\mathbf{v}$ integers.
Exercise 8.1.2.3: (1) Show that the space $S^{2}\left(S^{2} \mathbb{C}^{2}\right)$ contains a copy of $S_{(2,2)} \mathbb{C}^{2}$ by showing that $\left(x_{1}^{2}\right)\left(x_{2}^{2}\right)-\left(x_{1} x_{2}\right)\left(x_{1} x_{2}\right) \in S^{2}\left(S^{2} \mathbb{C}^{2}\right)$ is a highest weight vector.
Exercise 8.1.2.4: (1!) Find highest weight vectors in $V, S^{2} V, \Lambda^{2} V, S^{3} V, \Lambda^{3} V$ and the kernels of the symmetrization and skew-symmetrization maps $V \otimes S^{2} V \rightarrow$ $S^{3} V$ and $V \otimes \Lambda^{2} V \rightarrow \Lambda^{3} V$. Show that both of the last two modules have highest weight $(2,1)$, i.e., they are realizations of $S_{(2,1)} V$.

Exercise 8.1.2.5: (2) More generally, find a highest weight vector for the kernel of the symmetrization map $V \otimes S^{d-1} V \rightarrow S^{d} V$ and of the kernel of the "exterior derivative" (or "Koszul") map

$$
\begin{align*}
S^{k} V \otimes \Lambda^{t} V & \rightarrow S^{k-1} V \otimes \Lambda^{t+1} V  \tag{8.1.1}\\
x_{1} \cdots x_{k} \otimes y_{1} \wedge \cdots \wedge y_{t} & \mapsto \sum_{j=1}^{k} x_{1} \cdots \hat{x}_{j} \cdots x_{k} \otimes x_{j} \wedge y_{1} \wedge \cdots \wedge y_{t} .
\end{align*}
$$

Exercise 8.1.2.6: (1!) Let $\pi=\left(p_{1}, \ldots, p_{\ell}\right)$ be a partition with at most $\mathbf{v}$ parts and let $\pi^{\prime}=\left(q_{1}, \ldots, q_{p_{1}}\right)$ denote the conjugate partition. Show that

$$
\begin{equation*}
z_{\pi}:=\left(e_{1} \wedge \cdots \wedge e_{q_{1}}\right) \otimes\left(e_{1} \wedge \cdots \wedge e_{q_{2}}\right) \otimes \cdots \otimes\left(e_{1} \wedge \cdots \wedge e_{q_{p_{1}}}\right) \in V^{\otimes|\pi|} \tag{8.1.2}
\end{equation*}
$$

is a highest weight vector of weight $\pi$.
8.1.3. The Pieri rule. I describe the decomposition of $S_{\pi} V \otimes V$ as a $G L(V)$ module. Write $\pi^{\prime}=\left(q_{1}, \ldots, q_{p_{1}}\right)$ and recall $z_{\pi}$ from (8.1.2). Consider the vectors:

$$
\begin{aligned}
& \left(e_{1} \wedge \cdots \wedge e_{q_{1}} \wedge e_{q_{1}+1}\right) \otimes\left(e_{1} \wedge \cdots \wedge e_{q_{2}}\right) \otimes \cdots \otimes\left(e_{1} \wedge \cdots \wedge e_{q_{p_{1}}}\right) \\
& \vdots \\
& \left(e_{1} \wedge \cdots \wedge e_{q_{1}}\right) \otimes\left(e_{1} \wedge \cdots \wedge e_{q_{2}}\right) \otimes \cdots \otimes\left(e_{1} \wedge \cdots \wedge e_{q_{p_{1}}} \wedge e_{q_{p_{1}}+1}\right) \\
& \left(e_{1} \wedge \cdots \wedge e_{q_{1}}\right) \otimes\left(e_{1} \wedge \cdots \wedge e_{q_{2}}\right) \otimes \cdots \otimes\left(e_{1} \wedge \cdots \wedge e_{q_{p_{1}}}\right) \otimes e_{1} .
\end{aligned}
$$

These are all highest weight vectors obtained by tensoring $z_{\pi}$ with a vector in $V$ and skew-symmetrizing appropriately, so the associated modules are contained in $S_{\pi} V \otimes V$. With a little more work, one can show these are highest weight vectors of all the modules that occur in $S_{\pi} V \otimes V$. If $q_{j}=$ $q_{j+1}$ one gets the same module if one inserts $e_{q_{j}+1}$ into either slot, and its multiplicity in $S_{\pi} V \otimes V$ is one. More generally one obtains:
Theorem 8.1.3.1 (The Pieri formula). The decomposition of $S_{\pi} V \otimes S^{d} V$ is multiplicity free. The partitions corresponding to modules $S_{\mu} V$ that occur
are those obtained from the Young diagram of $\pi$ by adding $d$ boxes to the diagram of $\pi$, with no two boxes added to the same column.

Definition 8.1.3.2. Let $\pi=\left(p_{1}, \ldots, p_{l(\pi)}\right), \mu=\left(m_{1}, \ldots, m_{l(\mu)}\right)$ be partitions with $l(\mu)<l(\pi)$ One says $\mu$ interlaces $\pi$ if $p_{1} \geq m_{1} \geq p_{2} \geq m_{2} \geq \cdots \geq$ $m_{l(\pi)-1} \geq p_{l(\pi)}$.

Exercise 8.1.3.3: (1) Show that $S_{\pi} V \otimes S_{(d)} V$ consists of all the $S_{\mu} V$ such that $|\mu|=|\pi|+d$ and $\pi$ interlaces $\mu$.

Exercise 8.1.3.4: (1) Show that a necessary condition for $S_{\pi} V$ to appear in $S^{d}\left(S^{n} V\right)$ is that $l(\pi) \leq d$.

Although a pictorial proof is possible, the standard proof of the Pieri formula uses a character (see $\S 8.6 .7$ ) calculation, computing $\chi_{\pi} \chi_{(d)}$ as a sum of $\chi_{\mu}$ 's. See, e.g., [Mac95, §I.9]. A different proof, using Schur-Weyl duality is in [GW09, §9.2]. There is an algorithm to compute arbitrary tensor product decompositions called the Littlewood Richardson Rule. See, e.g., [Mac95, §I.9] for details.

Similar considerations give:
Theorem 8.1.3.5. [The skew-Pieri formula] The decomposition of $S_{\pi} V \otimes \Lambda^{k} V$ is multiplicity free. The partitions corresponding to modules $S_{\mu} V$ that occur are those obtained from the Young diagram of $\pi$ by adding $k$ boxes to the diagram of $\pi$, with no two boxes added to the same row.
8.1.4. The $G L(V)$-modules not appearing in the tensor algebra of $V$. The $G L(V)$-module $V^{*}$ does not appear in the tensor algebra of $V$. Nor do the one-dimensional representations for $k>0, \operatorname{det}^{-k}: G L(V) \rightarrow G L\left(\mathbb{C}^{1}\right)$ given by, for $v \in \mathbb{C}^{1}, \operatorname{det}^{-k}(g) v:=\operatorname{det}(g)^{-k} v$.
Exercise 8.1.4.1: (1) Show that if $\pi=\left(p_{1}, \ldots, p_{\mathbf{v}}\right)$ with $p_{\mathbf{v}}>0$, then $\operatorname{det}^{-1} \otimes S_{\pi} V=S_{\left(p_{1}-1, \ldots, p_{\mathbf{v}}-1\right)} V$. ©
Exercise 8.1.4.2: (1) Show that as a $G L(V)$-module, $V^{*}=\Lambda^{\mathbf{v}-1} V \otimes \operatorname{det}^{-1}=$ $S_{1 \mathrm{v}-1} V \otimes \operatorname{det}^{-1}$. ©

Every irreducible $G L(V)$-module is of the form $S_{\pi} V \otimes \operatorname{det}^{-k}$ for some partition $\pi$ and some $k \geq 0$. Thus they may be indexed by non-increasing sequences of integers $\left(p_{1}, \ldots, p_{\mathbf{v}}\right)$ where $p_{1} \geq p_{2} \geq \cdots \geq p_{\mathbf{v}}$. Such a module is isomorphic to $S_{\left(p_{1}-p_{\mathrm{v}}, \ldots, p_{\mathrm{v}-1}-p_{\mathrm{v}}, 0\right)} V \otimes \operatorname{det}^{p_{\mathrm{v}}}$.

Using

$$
S_{\pi} V \otimes V^{*}=S_{\pi} V \otimes \Lambda^{\mathrm{v}-1} V \otimes \operatorname{det}^{-1}
$$

we may compute the decomposition of $S_{\pi} V \otimes V^{*}$ using the skew-symmetric version of the Pieri rule.

Example 8.1.4.3. Let $\mathbf{w}=3$, then

$$
\begin{aligned}
S_{(32)} W \otimes W^{*} & =S_{(43)} W \otimes \operatorname{det}^{-1} \oplus S_{(331)} W \otimes \operatorname{det}^{-1} \oplus S_{(421)} W \otimes \operatorname{det}^{-1} \\
& =S_{(43)} W \otimes \operatorname{det}^{-1} \oplus S_{(22)} W \oplus S_{(31)} W
\end{aligned}
$$

The first module does not occur in the tensor algebra but the rest do.
8.1.5. $S L(V)$-modules in $V^{\otimes d}$. Every $S L(V)$-module is the restriction to $S L(V)$ of some $G L(V)$-module. However distinct $G L(V)$-modules, when restricted to $S L(V)$ can become isomorphic, such as the trivial representation and $\Lambda^{\mathbf{v}} V=S_{(1 \mathbf{v})} V=\operatorname{det}^{1}$.
Proposition 8.1.5.1. Let $\pi=\left(p_{1}, \ldots, p_{\mathbf{v}}\right)$ be a partition. The $S L(V)$ modules in the tensor algebra $V^{\otimes}$ that are isomorphic to $S_{\pi} V$ are $S_{\mu} V$ with $\mu=\left(p_{1}+j, p_{2}+j, \ldots, p_{\mathbf{v}}+j\right)$ for $-p_{\mathbf{v}} \leq j<\infty$.
Exercise 8.1.5.2: (2) Prove Proposition 8.1.5.1. ©
For example, for $S L_{2}$-modules, $S_{p_{1}, p_{2}} \mathbb{C}^{2} \simeq S^{p_{1}-p_{2}} \mathbb{C}^{2}$. We conclude:
Corollary 8.1.5.3. A complete set of the finite dimensional irreducible representations of $S L_{2}$ are the $S^{d} \mathbb{C}^{2}$ with $d \geq 0$.

The $G L(V)$-modules that are $S L(V)$-equivalent to $S_{\pi} V$ may be visualized as being obtained by erasing or adding columns of size $\mathbf{v}$ from the Young diagram of $\pi$, as in Figure 8.1.5.


Figure 8.1.3. Young diagrams for $S L_{3}$-modules equivalent to $S_{421} \mathbb{C}^{3}$

The Lie algebra of $S L(V)$, denoted $\mathfrak{s l}(V)$, is the set of traceless endomorphisms. One can define weights for the Lie algebra of the torus, which are essentially the logs of the corresponding torus in the group. In particular, vectors of $\mathfrak{s l}$-weight zero have $G L(V)$-weight $(d, \ldots, d)=\left(d^{\mathbf{v}}\right)$ for some $d$.
Exercise 8.1.5.4: (1!) Let $T^{S L} \subset S L(V)$ be the diagonal matrices with determinant one. Show that $\left(V^{\otimes d}\right)^{T^{S L}}$ is zero unless $d=\delta \mathbf{v}$ for some natural number $\delta$ and in this case it consists of all vectors of weight $\left(\delta^{\mathbf{v}}\right)$.

### 8.2. Young flattenings

Most known equations for border rank of tensors, i.e., polynomials in the ideal of the variety $\sigma_{r}\left(S e g\left(\mathbb{P} A_{1} \times \cdots \times \mathbb{P} A_{n}\right)\right)$ and symmetric border rank of polynomials, i.e., polynomials in the ideal of the variety $\sigma_{r}\left(v_{d}(\mathbb{P} V)\right)$, are
obtained by taking minors of some auxiliary matrix constructed from the tensor (polynomial). What follows is a general way to use representation theory to find such matrices.
8.2.1. The case of polynomials. Let $P \in S^{d} V$. Recall the flattenings from §6.2: $P_{k, d-k}: S^{k} V^{*} \rightarrow S^{d-k} V$. Flattenings give rise to a $G L(V)$ module of equations because $S^{d} V \subset S^{k} V \otimes S^{d-k} V$. The generalization is similar:
Proposition 8.2.1.1. Given a linear inclusion $S^{d} V \subset U \otimes W$, i.e., $S^{d} V$ is realized as a space of linear maps from $U^{*}$ to $W$, say the rank of the linear map associated to $\ell^{d}$ is $r_{0}$. If the rank of the linear map associated to $P$ is $r$, then $\underline{\mathbf{R}}_{S}(P) \geq \frac{r}{r_{0}}$.
Exercise 8.2.1.2: (1!) Prove Proposition 8.2.1.1. ©
This method works best when $r_{0}$ is small. For example in the classical flattening case $r_{0}=1$.

We will take $U, W$ to be $G L(V)$-modules and the linear inclusion a $G L(V)$-module map because $I\left(\sigma_{r}\left(v_{d}(\mathbb{P} V)\right)\right)$ is a $G L(V)$-module. It turns out that we know all such maps. The Pieri rule $\S 8.1 .3$ says they are all of the form $S^{d} V \subset S_{\pi} V^{*} \otimes S_{\mu} V$ where the Young diagram of $\mu$ is obtained from the Young diagram of $\pi$ by adding $d$ boxes, with no two boxes added to the same column. To make this completely correct, we need to consider sequences with negative integers, where e.g., the Young diagram of $(-d)$ should be thought of as $-d$ boxes in a row. Alternatively, one can work with $S L(V)$-modules, as then $\operatorname{det}^{-d}=S_{(-d)} V=S_{\left(d^{v}-1\right)} V$ as $S L(V)$ modules. For every such pair $(\pi, \mu)$ there is exactly one $G L(V)$-inclusion. Call the resulting linear map a Young-flattening.

The classical case is $\pi=(-k)$ and $\mu=(d-k)$, or in terms of $S L(V)$ modules, $\pi=\left(k^{\mathbf{v}-1}\right)$ and $\mu=\left(k^{\mathbf{v}}, d-k\right)$. The main example in [LO13], called a (polynomial) Koszul flattening was constructed as follows: take the classical flattening $P_{k, d-k}: S^{k} V^{*} \rightarrow S^{d-k} V$ and tensor it with $\operatorname{Id}_{\Lambda^{p} V}$ for some $p$, to get a map $S^{k} V^{*} \otimes \Lambda^{p} V \rightarrow S^{d-k} V \otimes \Lambda^{p} V$. Now include $S^{d-k} V \subset$ $S^{d-k-1} V \otimes V$, to obtain a map $S^{k} V^{*} \otimes \Lambda^{p} V \rightarrow S^{d-k-1} V \otimes V \otimes \Lambda^{p} V$ and finally skew-symmetrize the last two factors to obtain a map

$$
\begin{equation*}
P_{k, d-k}^{\wedge p}: S^{k} V^{*} \otimes \Lambda^{p} V \rightarrow S^{d-k-1} V \otimes \Lambda^{p+1} V \tag{8.2.1}
\end{equation*}
$$

If one views this as a map $S^{d} V \otimes\left(S^{k} V^{*} \otimes \Lambda^{p} V\right) \rightarrow S^{d-k-1} V \otimes \Lambda^{p+1} V$, it is a $G L(V)$-module map. By the Pieri rule,

$$
\left(S^{k} V^{*} \otimes \Lambda^{p} V\right)^{*}=S_{\left(k, 1^{\mathbf{v}-p}\right)} V \otimes \operatorname{det}^{-1} \oplus S_{\left(k+1,1^{\mathbf{v}-p-1)}\right.} V \otimes \operatorname{det}^{-1}
$$

and

$$
S^{d-k-1} V \otimes \Lambda^{p+1} V=S_{\left(d-k-1,1^{p+1}\right)} V \oplus S_{\left(d-k, 1^{p}\right)} V .
$$

Although in practice one usually works with the map (8.2.1), the map is zero except restricted to the map between irreducible modules:

$$
\left[S_{\left(k, 1^{\mathbf{v}-p}\right)} V^{*} \otimes \operatorname{det}^{-1}\right]^{*} \rightarrow S_{\left(d-k, 1^{p}\right)} V
$$

The method of shifted partial derivatives $\S 7.6$ is a type of Young flattening which I will call a Hilbert flattening, because it is via Hilbert functions of Jacobian ideals. It is the symmetric cousin of the Koszul flattening: take the classical flattening $P_{k, d-k}: S^{k} V^{*} \rightarrow S^{d-k} V$ and tensor it with $\operatorname{Id}_{S^{\ell} V}$ for some $\ell$, to get a map $S^{k} V^{*} \otimes S^{\ell} V \rightarrow S^{d-k} V \otimes S^{\ell} V$. Now simply take the projection (multiplication map) $S^{d-k} V \otimes S^{\ell} V \rightarrow S^{d-k+\ell} V$, to obtain a map

$$
\begin{equation*}
P_{k, d-k[\ell]}: S^{k} V^{*} \otimes S^{\ell} V \rightarrow S^{d-k+\ell} V \tag{8.2.2}
\end{equation*}
$$

The target is an irreducible $G L(V)$-module, so the pruning is easier here.
8.2.2. The case of tensors. Young flattenings can also be defined for tensors. For tensors in $A \otimes B \otimes C$, the Koszul flattenings $T_{A}^{\wedge p}: \Lambda^{p} A \otimes B^{*} \rightarrow$ $\Lambda^{p+1} A \otimes C$ used in $\S 2.4$ appear to be the only useful cases.

In principle there are numerous inclusions

$$
A \otimes B \otimes C \subset\left(S_{\pi} A \otimes S_{\mu} B \otimes S_{\nu} C\right)^{*} \otimes\left(S_{\tilde{\pi}} A \otimes S_{\tilde{\mu}} B \otimes S_{\tilde{\nu}} C\right)
$$

where the Young diagram of $\tilde{\pi}$ is obtained from the Young diagram of $\pi$ by adding a box (and similarly for $\mu, \nu$ ), and the case of Koszul flattenings is where (up to permuting the three factors) $\pi=\left(1^{p}\right), \mu=\left(1^{\mathbf{b}-1}\right)$ (so $S_{\mu} B \simeq B^{*}$ as $S L(B)$-modules) and $\nu=\emptyset$.

Exercise 2.4.1.1 already indicates why symmetrization is not useful, and an easy generalization of it proves this to be the case for Young flattenings of tensors. But perhaps additional skew-symmetrization could be useful: Let $T \in A \otimes B \otimes C$ and consider $T \otimes \operatorname{Id}_{\Lambda^{p} A} \otimes \operatorname{Id}_{\Lambda^{q} B} \otimes \operatorname{Id}_{\Lambda^{s} C}$ as a linear map $B^{*} \otimes \Lambda^{q} B^{*} \otimes \Lambda^{p} A \otimes \Lambda^{s} C \rightarrow \Lambda^{q} B^{*} \otimes \Lambda^{p} A \otimes A \otimes \Lambda^{s} C \otimes C$. Now quotient to the exterior powers to get a map:

$$
T_{p, q, s}: \Lambda^{q+1} B^{*} \otimes \Lambda^{p} A \otimes \Lambda^{s} C \rightarrow \Lambda^{q} B^{*} \otimes \Lambda^{p+1} A \otimes \Lambda^{s+1} C
$$

This generalizes the map $T_{A}^{\wedge p}$ which is the case $q=s=0$. Claim: this generalization does not give better lower bounds for border rank than Koszul flattenings when $\mathbf{a}=\mathbf{b}=\mathbf{c}$. (Although it is possible it could give better lower bounds for some particular tensor.) If $T$ has rank one, say $T=a \otimes b \otimes c$, the image of $T_{p, q, s}$ is

$$
\Lambda^{q}\left(b^{\perp}\right) \otimes\left(a \wedge \Lambda^{p} A\right) \otimes\left(c \wedge \Lambda^{s} C\right)
$$

Here $b^{\perp}:=\left\{\beta \in B^{*} \mid \beta(b)=0\right\}$. The image of $(a \otimes b \otimes c)_{p, q, s}$ has dimension

$$
d_{p, q, s}:=\binom{\mathbf{b}-1}{q}\binom{\mathbf{a}-1}{p}\binom{\mathbf{c}-1}{s}
$$

Thus the size $r d_{p, q, s}+1$ minors of $T_{p, q, s}$ potentially give equations for the variety of tensors of border rank at most $r$. We have nontrivial minors as long as

$$
r d_{p, q, s}+1 \leq \min \left\{\operatorname{dim}\left(\Lambda^{q} B^{*} \otimes \Lambda^{p+1} A \otimes \Lambda^{s+1}\right), \operatorname{dim}\left(\Lambda^{q+1} B^{*} \otimes \Lambda^{p} A \otimes \Lambda^{s} C\right)\right\},
$$

i.e., as long as

$$
r<\min \left\{\frac{\binom{\mathbf{b}}{q}\binom{\mathbf{a}}{p+1}\binom{\mathbf{c}}{s+1}}{\binom{\mathbf{b}-1}{q}\binom{\mathbf{a}-1}{p}\binom{\mathbf{c}-1}{s}}, \frac{\binom{\mathbf{b}}{q+1}\binom{\mathbf{a}}{p}\binom{\mathbf{c}}{s}}{\binom{\mathbf{b}-1}{q}\binom{\mathbf{a}-1}{p}\binom{\mathbf{c}-1}{s}}\right\},
$$

i.e.

$$
r<\min \left\{\frac{\mathbf{a b c}}{(\mathbf{b}-q)(p+1)(s+1)}, \frac{\mathbf{a b c}}{(q+1)(\mathbf{a}-p)(\mathbf{c}-s)}\right\} .
$$

Consider the case $q=0$, so we need

$$
r<\min \left\{\frac{\mathbf{a c}}{(p+1)(s+1)}, \frac{\mathbf{a b c}}{(\mathbf{a}-p)(\mathbf{c}-s)}\right\} .
$$

Let's specialize to $\mathbf{a}=\mathbf{c}, p=q$, so we need

$$
r<\min \left\{\frac{\mathbf{a}^{2}}{(p+1)^{2}}, \frac{\mathbf{a}^{2} \mathbf{b}}{(\mathbf{a}-p)^{2}}\right\} .
$$

Consider the case $\mathbf{a}=m p$ for some $m$. Then if $m$ is large, the first term is large, but the second is very close to $\mathbf{b}$. So unless the dimensions are unbalanced, one is unlikely to get any interesting equations out of these Young flattenings.
8.2.3. General perspective. Let $X \subset \mathbb{P} V$ be a $G$-variety for some reductive group $G$, where $V=V_{\lambda}$ is an irreducible $G$-module.
Proposition 8.2.3.1. Given irreducible $G$-modules $V_{\mu}, V_{\nu}$ such that $V_{\lambda} \subset$ $V_{\mu} \otimes V_{\nu}$ and $v \in V$, we obtain a linear map $v_{\mu, \nu}: V_{\mu}^{*} \rightarrow V_{\nu}$. Say the maximum rank of such a linear map for $x \in X$ is $q$, then the size ( $q r+1$ )-minors of $v_{\mu, \nu}$ test membership $\sigma_{r}(X)$.

### 8.3. Additional uses of representation theory to find modules of equations

In this section, I briefly cover additional techniques for finding modules of polynomials in ideals of $G$-varieties. I am brief because either the methods are not used in this book or they are described at length in [Lan12].
8.3.1. A naïve algorithm. Let $X \subset \mathbb{P} W$ be a $G$-variety. We are primarily interested in the cases $X=\sigma_{r}(\operatorname{Seg}(\mathbb{P} A \times \mathbb{P} B \times \mathbb{P} C)) \subset \mathbb{P}(A \otimes B \otimes C)$, where $G=G L(A) \times G L(B) \times G L(C)$ and $X=\mathcal{D e t}_{n} \subset \mathbb{P} S^{n} \mathbb{C}^{n^{2}}$, where $G=G L_{n^{2}}$. Since the ideal of $X$ will be a $G$-module, we can look for irreducible modules in the ideal of $X$ by testing highest weight vectors. If $U \subset S^{d} W^{*}$ is an irreducible $G$-module with highest weight vector $u$, then $U \subset I(X)$ if and only if $u \in I(X)$ because if $u \in I(X)$ then $g(u) \in I(X)$ for all $g \in G$ and such vectors span $U$. Thus in each degree $d$, we can in principle determine $I_{d}(X)$ by a finite calculation. In practice we test each highest weight vector $u$ on a "random" point $[x] \in X$. (If $X$ is an orbit closure it suffices to test on a point in the orbit.) If $u(x) \neq 0$, we know for sure that $U \not \subset I_{d}(X)$. If $u(x)=0$, then with extremely high probability (probability one if the point is truly randomly chosen, and with certainty if dealing with an orbit closure), we have $U \subset I(X)$. After testing several such points, we have high confidence in the result. Once one has a candidate module by such tests, one can often prove it is in the ideal by different methods.

More precisely, if $S^{d} W^{*}$ is multiplicity free as a $G$-module, there are a finite number of highest weight vectors to check. If a given module has multiplicity $m$, then we need to take a basis $u_{1}, \ldots, u_{m}$ of the highest weight space, test on say $x_{1}, \ldots, x_{q}$ with $q \geq m$ if $\sum_{j} y_{j} u_{j}\left(x_{s}\right)=0$ for some constants $y_{1}, \ldots, y_{m}$ and all $1 \leq s \leq q$.

To carry out this procedure in our two cases we would respectively need

- A method to decompose $S^{d}(A \otimes B \otimes C)^{*}\left(\right.$ resp. $\left.S^{d}\left(S^{n} \mathbb{C}^{n^{2}}\right)\right)$ into irreducible submodules.
- A method to explicitly write down highest weight vectors.

There are several systematic techniques for accomplishing both these tasks that work well in small cases, but as cases get larger one needs to introduce additional methods to be able to carry out the calculations in practice. The first task amounts to the well-studied problem of computing Kronecker coefficients defined in $\S 8.8 .2$. I briefly discuss the second task in §8.7.2.
8.3.2. Enhanced search using numerical methods. Rather than discuss the general theory, I outline the method used in [HIL13] to find equations for $\sigma_{6}\left(\operatorname{Seg}\left(\mathbb{P}^{3} \times \mathbb{P}^{3} \times \mathbb{P}^{3}\right)\right)$. First fix a "random" linear space $L \subset \mathbb{P}^{63}$ of dimension 4 (i.e., $\operatorname{codim} \sigma_{6}\left(\operatorname{Seg}\left(\mathbb{P}^{3} \times \mathbb{P}^{3} \times \mathbb{P}^{3}\right)\right)$ ) and consider the finite set $Z:=\sigma_{6}\left(\operatorname{Seg}\left(\mathbb{P}^{3} \times \mathbb{P}^{3} \times \mathbb{P}^{3}\right)\right) \cap L$. The first objective is to compute points in $Z$, with a goal of computing every point in $Z$. To this end, we first computed one point in $Z$ as follows. One first picks a random point $x^{*} \in \sigma_{6}\left(\operatorname{Seg}\left(\mathbb{P}^{3} \times \mathbb{P}^{3} \times \mathbb{P}^{3}\right)\right)$, which is easy since an open dense subset of $\sigma_{6}\left(\operatorname{Seg}\left(\mathbb{P}^{3} \times \mathbb{P}^{3} \times \mathbb{P}^{3}\right)\right)$ is parameterizable. Let $\tilde{L}$ be a system of 59 linear
forms so that $L$ is the zero locus of $\tilde{L}$ and let $L_{t, x^{*}}$ be the zero locus of $L(x)-t \cdot L\left(x^{*}\right)$. Since $x^{*} \in \sigma_{6}\left(\operatorname{Seg}\left(\mathbb{P}^{3} \times \mathbb{P}^{3} \times \mathbb{P}^{3}\right)\right) \cap L_{1, x^{*}}$, a point in $Z$ is the endpoint of the path defined by $\sigma_{6}\left(\operatorname{Seg}\left(\mathbb{P}^{3} \times \mathbb{P}^{3} \times \mathbb{P}^{3}\right)\right) \cap L_{t, x^{*}}$ at $t=0$ starting from $x^{*}$ at $t=1$.

Even though the above process could be repeated for different $x^{*}$ to compute points in $Z$, we instead used monodromy loops [SVW01] for generating more points in $Z$. After performing 21 loops, the number of points in $Z$ that we computed stabilized at 15,456 . The trace test [SVW02] shows that 15,456 is indeed the degree of $\sigma_{6}\left(\operatorname{Seg}\left(\mathbb{P}^{3} \times \mathbb{P}^{3} \times \mathbb{P}^{3}\right)\right)$ thereby showing we had indeed computed $Z$.

From $Z$, we performed two computations. The first was the membership test of $[\mathbf{H S 1 3}]$ for deciding if $M_{\langle 2\rangle} \in \sigma_{6}\left(\operatorname{Seg}\left(\mathbb{P}^{3} \times \mathbb{P}^{3} \times \mathbb{P}^{3}\right)\right)$, which requires tracking 15,456 homotopy paths that start at the points of $Z$ and end on a $\mathbb{P}^{4}$ containing $M_{\langle 2\rangle}$. In this case, each of these 15,456 paths converged to points in $\sigma_{6}\left(\operatorname{Seg}\left(\mathbb{P}^{3} \times \mathbb{P}^{3} \times \mathbb{P}^{3}\right)\right)$ distinct from $M_{\langle 2\rangle}$ providing a numerical proof that $M_{\langle 2\rangle} \notin \sigma_{6}\left(\operatorname{Seg}\left(\mathbb{P}^{3} \times \mathbb{P}^{3} \times \mathbb{P}^{3}\right)\right)$. The second was to compute the minimal degree of nonzero polynomials vanishing on $Z \subset L$. This sequence of polynomial interpolation problems showed that no non-constant polynomials of degree $\leq 18$ vanished on $Z$ and hence on $\sigma_{6}\left(\operatorname{Seg}\left(\mathbb{P}^{3} \times \mathbb{P}^{3} \times \mathbb{P}^{3}\right)\right)$. The $15456 \times 8855$ matrix resulting from polynomial interpolation of homogeneous forms of degree 19 in 5 variables using the approach of [GHPS14] has a 64-dimensional null space. Thus, the minimal degree of nonzero polynomials vanishing on $Z \subset L$ is 19 , showing $\operatorname{dim} I_{19}\left(\sigma_{6}\right) \leq 64$.

The next objective was to verify that the minimal degree of nonzero polynomials vanishing on the curve $C:=\sigma_{6}\left(\operatorname{Seg}\left(\mathbb{P}^{3} \times \mathbb{P}^{3} \times \mathbb{P}^{3}\right)\right) \cap K \subset K$ for a fixed "random" linear space $K \subset \mathbb{P}^{63}$ of dimension 5 was also 19 . We used 50,000 points on $C$ and the $50000 \times 42504$ matrix resulting from polynomial interpolation of homogeneous forms of degree 19 in 6 variables using the approach of [GHPS14] also has a 64 -dimensional null space. With this agreement, we decomposed $S^{19}\left(\mathbb{C}^{4} \otimes \mathbb{C}^{4} \otimes \mathbb{C}^{4}\right)$ and looked for a 64 dimensional submodule. The only reasonable candidate was to take a copy of $S_{5554} \mathbb{C}^{4} \otimes S_{5554} \mathbb{C}^{4} \otimes S_{5554} \mathbb{C}^{4}$. We found a particular copy that was indeed in the ideal and then proved that $M_{\langle 2\rangle}$ is not contained in $\sigma_{6}\left(\operatorname{Seg}\left(\mathbb{P}^{3} \times \mathbb{P}^{3} \times \mathbb{P}^{3}\right)\right)$ by showing a polynomial in this module did not vanish on it. The evaluation was numerical, so the result was:
Theorem 8.3.2.1. [HIL13] With extremely high probability, the ideal of $\sigma_{6}\left(\operatorname{Seg}\left(\mathbb{P}^{3} \times \mathbb{P}^{3} \times \mathbb{P}^{3}\right)\right)$ is generated in degree 19 by the module $S_{5554} \mathbb{C}^{4} \otimes S_{5554} \mathbb{C}^{4} \otimes S_{5554} \mathbb{C}^{4}$. This module does not vanish on $M_{\langle 2\rangle}$.

In the same paper, a copy of the trivial degree twenty module $S_{5555} \mathbb{C}^{4} \otimes S_{5555} \mathbb{C}^{4} \otimes S_{5555} \mathbb{C}^{4}$ is shown to be in the ideal of $\sigma_{6}\left(\operatorname{Seg}\left(\mathbb{P}^{3} \times \mathbb{P}^{3} \times \mathbb{P}^{3}\right)\right)$ by symbolic methods, giving a new proof that:
Theorem 8.3.2.2. [Lan06, HIL13] $\underline{\mathbf{R}}\left(M_{\langle 2\rangle}\right)=7$.
The same methods have shown $I_{45}\left(\sigma_{15}\left(\operatorname{Seg}\left(\mathbb{P}^{3} \times \mathbb{P}^{7} \times \mathbb{P}^{8}\right)\right)=0\right.$ and that $I_{186,999}\left(\sigma_{18}\left(\operatorname{Seg}\left(\mathbb{P}^{6} \times \mathbb{P}^{6} \times \mathbb{P}^{6}\right)\right)=0\right.$ (this variety is a hypersurface), both of which are relevant for determining the border rank of $M_{\langle 3\rangle}$, see [HIL13].
8.3.3. Inheritance. Inheritance is a general technique for studying equations of $G$-varieties that come in series. It is discussed extensively in [Lan12, $\S 7.4, \S 16.4]$.

If $V \subset W$ then $S_{\bar{\pi}} V \subset V^{\otimes d}$ induces a module $S_{\bar{\pi}} W \subset W^{\otimes d}$ by, e.g., choosing a basis of $W$ whose first $\mathbf{v}$ vectors are a basis of $V$. Then the two modules have the same highest weight vector and one obtains the $G L(W)$ module the span of the $G L(W)$-orbit of the highest weight vector.

Because the realizations of $S_{\pi} V$ in $V^{\otimes d}$ do not depend on the dimension of $V$, one can reduce the study of $\sigma_{r}(\operatorname{Seg}(\mathbb{P} A \times \mathbb{P} B \times \mathbb{P} C))$ to that of $\sigma_{r}\left(S e g\left(\mathbb{P}^{r-1} \times \mathbb{P}^{r-1} \times \mathbb{P}^{r-1}\right)\right)$. As discussed in $\S 3.3 .1$ this latter variety is an orbit closure, namely the orbit closure of $M_{\langle 1\rangle}^{\oplus r}$.
Proposition 8.3.3.1. [LM04, Prop. 4.4] For all vector spaces $B_{j}$ with $\operatorname{dim} B_{j}=\mathbf{b}_{j} \geq \operatorname{dim} A_{j}=\mathbf{a}_{j} \geq r$, a module $S_{\bar{\mu}_{1}} B_{1} \otimes \cdots \otimes S_{\bar{\mu}_{n}} B_{n}$ such that $l\left(\mu_{j}\right) \leq \mathbf{a}_{j}$ for all $j$, is in $I_{d}\left(\sigma_{r}\left(\operatorname{Seg}\left(\mathbb{P} B_{1}^{*} \times \cdots \times \mathbb{P} B_{n}^{*}\right)\right)\right)$ if and only if $S_{\bar{\mu}_{1}} A_{1} \otimes \cdots \otimes S_{\bar{\mu}_{n}} A_{n}$ is in $I_{d}\left(\sigma_{r}\left(\operatorname{Seg}\left(\mathbb{P} A_{1}^{*} \times \cdots \times \mathbb{P} A_{n}^{*}\right)\right)\right)$.
Corollary 8.3.3.2. [LM04, AR03] Let $\operatorname{dim} A_{j} \geq r, 1 \leq j \leq n$. The ideal of $\sigma_{r}\left(\operatorname{Seg}\left(\mathbb{P} A_{1} \times \cdots \times \mathbb{P} A_{n}\right)\right)$ is generated by the modules inherited from the ideal of $\sigma_{r}\left(\operatorname{Seg}\left(\mathbb{P}^{r-1} \times \cdots \times \mathbb{P}^{r-1}\right)\right)$ and the modules generating the ideal of $S u b_{r, \ldots, r}\left(A_{1} \otimes \cdots \otimes A_{n}\right)$. The analogous scheme and set-theoretic results hold as well.
8.3.4. Prolongation. Prolongation and multi-prolongation provide a systematic method to find equations for secant varieties that is particularly effective for secant varieties of homogeneous varieties. For a general discussion and proofs see [Lan12, §7.5]. For our purposes, we will need the following:
Proposition 8.3.4.1. Given $X \subset \mathbb{P} V^{*}, I_{r+1}\left(\sigma_{r}(X)\right)=\left(I_{2}(X) \otimes S^{r-1} V\right) \cap$ $S^{r+1} V$.

Proposition 8.3.4.2. [SS09] Let $X \subset \mathbb{P} V$ be a variety with $I_{d-1}(X)=0$.
Then for all $\delta<(d-1) r, I_{\delta}\left(\sigma_{r}(X)\right)=0$.
Corollary 8.3.4.3. $I_{d}\left(\sigma_{d}\left(v_{n}(\mathbb{P} V)\right)=0\right.$.

### 8.4. Necessary conditions for modules of polynomials to be useful for GCT

The polynomial $\ell^{n-m} \operatorname{perm}_{m} \in S^{n} \mathbb{C}^{n^{2}}$ has two properties that can be studied individually: it is padded, i.e., it is divisible by a large power of a linear form, and its zero set is a cone with a $\left(n^{2}-m^{2}-1\right)$-dimensional vertex, that is, it only uses $m^{2}+1$ of the $n^{2}$ variables in an expression in good coordinates. Both of these properties restrict the types of polynomials we should look for. Equipped with the language of representation theory we can give precise descriptions of the modules we should restrict our attention to, which I call GCT useful.

I begin with the study of cones, a classical topic.
8.4.1. Cones. Recall the subspace variety $S u b_{k}\left(S^{d} V\right) \subset \mathbb{P} S^{d} V$ from $\S 6.2 .2$, the polynomials whose associated hypersurfaces are cones with a $\mathbf{v}-k$ dimensional vertex.
Proposition 8.4.1.1. $I_{\delta}\left(S u b_{k}\left(S^{d} V\right)\right)$ consists of the isotypic components of the modules $S_{\pi} V^{*}$ appearing in $S^{\delta}\left(S^{d} V^{*}\right)$ such that $l(\pi)>k$.
Exercise 8.4.1.2: (2!) Prove Proposition 8.4.1.1. ©
With just a little more effort, one can prove the degree $k+1$ equations from Proposition 8.4.1.1 generate the ideal:
Theorem 8.4.1.3. [Wey03, Cor. 7.2.3] The ideal of $\operatorname{Sub}_{k}\left(S^{d} V\right)$ is generated by the image of $\Lambda^{k+1} V^{*} \otimes \Lambda^{k+1} S^{d-1} V^{*} \subset S^{k+1}\left(V^{*} \otimes S^{d-1} V^{*}\right)$ in $S^{k+1}\left(S^{d} V^{*}\right)$, the size $k+1$ minors of the $(k, d-k)$-flattening.
Aside 8.4.1.4. Here is further information about the variety $S u b_{k}\left(S^{d} V\right)$ : It is an example of a variety admitting a Kempf-Weyman desingularization, a type of desingularization that $G$-varieties often admit. Rather than discuss the general theory here (see [Wey03] for a full exposition or [Lan12, Chap. 17] for an elementary introduction), I just explain this example, which gives a proof of Theorem 8.4.1.3, although more elementary proofs are possible. As was mentioned in $\S 5.4 .3$, the Grassmannian $G(k, V)$ has a tautological vector bundle $\pi: \mathcal{S} \rightarrow G(k, V)$, where the fiber over a $k$-plane $E$ is just the $k$-plane itself. The whole bundle is a sub-bundle of the trivial bundle $\underline{V}$ with fiber $V$. Consider the bundle $S^{d} \mathcal{S} \subset S^{d} \underline{V}$. We have a projection map $p: S^{d} \underline{V} \rightarrow S^{d} V$. The image of $S^{d} \mathcal{S}$ under $p$ is $\hat{S} u b_{k}\left(S^{d} V\right)$. Moreover, the map is a desingularization, that is $S^{d} \mathcal{S}$ is smooth, and the map to $\hat{S} u b_{k}\left(S^{d} V\right)$ is generically one to one. In particular, this implies $\operatorname{dim} \hat{S} u b_{k}\left(S^{d} V\right)=\operatorname{dim}\left(S^{d} \mathcal{S}\right)=\binom{k+d-1}{d}+d(\mathbf{v}-k)$. One obtains the entire minimal free resolution of $S u b_{k}\left(S^{d} V\right)$ by "pushing down" a tautological resolution "upstairs". From the minimal free resolution one can read off the generators of the ideal.
8.4.2. The variety of padded polynomials. Define the variety of padded polynomials

$$
\begin{aligned}
& \operatorname{Pad}_{n-m}\left(S^{n} W\right):= \\
& \mathbb{P}\left\{P \in S^{n} W \mid P=\ell^{n-m} h, \text { for some } \ell \in W, h \in S^{m} W\right\} \subset \mathbb{P} S^{n} W
\end{aligned}
$$

Note that $\operatorname{Pad}_{n-m}\left(S^{n} W\right)$ is a $G L(W)$-variety.
Proposition 8.4.2.1. [KL14] Let $\pi=\left(p_{1}, \ldots, p_{\mathbf{w}}\right)$ be a partition of $d n$. If $p_{1}<d(n-m)$, then the isotypic component of $S_{\pi} W^{*}$ in $S^{d}\left(S^{n} W^{*}\right)$ is contained in $I_{d}\left(\operatorname{Pad}_{n-m}\left(S^{n} W\right)\right)$.

Proof. Fix a (weight) basis $e_{1}, \ldots, e_{\mathbf{w}}$ of $W$ with dual basis $x_{1}, \ldots, x_{\mathbf{w}}$ of $W^{*}$. Note any element $\ell^{n-m} h \in \operatorname{Pad}_{n-m}\left(S^{n} W\right)$ is in the $G L(W)$-orbit of $\left(e_{1}\right)^{n-m} \tilde{h}$ for some $\tilde{h}$, so it will be sufficient to show that the ideal in degree $d$ contains the modules vanishing on the orbits of elements of the form $\left(e_{1}\right)^{n-m} h$. The highest weight vector of any copy of $S_{\left(p_{1}, \ldots, p_{\mathbf{w}}\right)} W^{*}$ in $S^{d}\left(S^{n} W^{*}\right)$ will be a linear combination of vectors of the form $m_{I}:=$ $\left(x_{1}^{i_{1}^{1}} \cdots x_{\mathbf{w}}^{i_{\mathbf{w}}^{1}}\right) \cdots\left(x_{1}^{i_{1}^{d}} \cdots x_{\mathbf{w}}^{i_{\mathbf{w}}^{d}}\right)$, where $i_{j}^{1}+\cdots+i_{j}^{d}=p_{j}$ for all $1 \leq j \leq \mathbf{w}$ and $i_{1}^{k}+\cdots+i_{\mathbf{w}}^{k}=n$ for all $1 \leq k \leq d$ as these are all the vectors of weight $\pi$ in $S^{d}\left(S^{n} W\right)$. Each $m_{I}$ vanishes on any $\left(e_{1}\right)^{n-m} h$ unless $p_{1} \geq d(n-m)$. (For a coordinate-free proof, see $[\mathbf{K L 1 4}]$.)

What we really need to study is the variety $\operatorname{Pad}_{n-m}\left(S u b_{k}\left(S^{d} W\right)\right)$ of padded cones.
Proposition 8.4.2.2. $[\mathbf{K L 1 4}] I_{d}\left(\operatorname{Pad}_{n-m}\left(S u b_{k}\left(S^{n} W^{*}\right)\right)\right)$ consists of all modules $S_{\bar{\pi}} W$ such that $S_{\bar{\pi}} \mathbb{C}^{k}$ is in the ideal of $\operatorname{Pad}_{n-m}\left(S^{n} \mathbb{C}^{k^{*}}\right)$ and all modules whose associated partition has length at least $k+1$.
Exercise 8.4.2.3: (2) Prove Proposition 8.4.2.2.
In summary:
Proposition 8.4.2.4. In order for a module $S_{\left(p_{1}, \ldots, p_{l}\right)} W^{*}$, where $\left(p_{1}, \ldots, p_{l}\right)$ is a partition of $d n$ to be GCT-useful for showing $\ell^{n-m} \operatorname{perm}_{m} \notin \overline{G L_{n^{2}} \cdot \operatorname{det}_{n}}$ we must have

- $l \leq m^{2}+1$, and
- $p_{1}>d(n-m)$.


### 8.5. Representation theory and $\mathcal{D e t}_{n}$

8.5.1. Proof of Proposition 6.7.2.2. Recall $P_{\Lambda}(M)=\overline{\operatorname{det}}_{n}\left(M_{\Lambda}, \ldots, M_{\Lambda}, M_{S}\right)$ from $\S 6.7 .1$ where $M=M_{\Lambda}+M_{S}$ is the decomposition of the matrix $M$ into its skew-symmetric and symmetric components. We need to show
$\overline{G L_{n^{2}} \cdot\left[P_{\Lambda}\right]}$ has codimension one in $\mathcal{D e t}_{n}$ and is not contained in $\operatorname{End}\left(\mathbb{C}^{n^{2}}\right)$. [ $\operatorname{det}_{n}$ ]. We compute the stabilizer of $P_{\Lambda}$ inside $G L(E \otimes E)$, where $E=\mathbb{C}^{n}$. The action of $G L(E)$ on $E \otimes E$ by $M \mapsto g M g^{T}$ preserves $P_{\Lambda}$ up to scale, and the Lie algebra of the stabilizer of $\left[P_{\Lambda}\right]$ is a $G L(E)$ submodule of $\operatorname{End}(E \otimes E)$. Note that $\mathfrak{s l}(E)=S_{\left(21^{n-2}\right)} E$ and $\mathfrak{g l}(E)=\mathfrak{s l}(E) \oplus \mathbb{C}$. Decompose End $(E \otimes E)$ as a $S L(E)$-module:

$$
\begin{align*}
& \operatorname{End}(E \otimes E)=\operatorname{End}\left(\Lambda^{2} E\right) \oplus \operatorname{End}\left(S^{2} E\right) \oplus \operatorname{Hom}\left(\Lambda^{2} E, S^{2} E\right) \oplus \operatorname{Hom}\left(S^{2} E, \Lambda^{2} E\right) \\
& \quad=\Lambda^{2} E \otimes \Lambda^{2} E^{*} \oplus S^{2} E \otimes S^{2} E^{*} \oplus \Lambda^{2} E^{*} \otimes S^{2} E \oplus S^{2} E^{*} \otimes \Lambda^{2} E \\
& \quad(8.5 .1)  \tag{8.5.1}\\
& \quad=\left(\mathfrak{g l}(E) \oplus S_{2^{2}, 1^{n-2}} E\right) \oplus\left(\mathfrak{g l}(E) \oplus S_{4,2^{n-1}} E\right) \oplus\left(\mathfrak{s l}(E) \oplus S_{3,1^{n-2}} E\right) \oplus\left(\mathfrak{s l}(E) \oplus S_{3^{2}, 2^{n-2}} E\right)
\end{align*}
$$

By testing highest weight vectors, one concludes the Lie algebra of $G_{P_{\Lambda}}$ is isomorphic to $\mathfrak{g l}(E) \oplus \mathfrak{g l}(E)$, which has dimension $2 n^{2}=\operatorname{dim} G_{\operatorname{det}_{n}}+1$, implying $\overline{G L(W) \cdot P_{\Lambda}}$ has codimension one in $\overline{G L(W) \cdot\left[\operatorname{det}_{n}\right]}$. Since it is not contained in the orbit of the determinant, it must be an irreducible component of its boundary. Since the zero set is not a cone, $P_{\Lambda}$ cannot be in $\operatorname{End}(W) \cdot \operatorname{det}_{n}$ which consists of $G L(W) \cdot \operatorname{det}_{n}$ plus polynomials whose zero sets are cones, as any element of $\operatorname{End}(W)$ either has a kernel or is invertible.
Exercise 8.5.1.1: (3) Verify by testing on highest weight vectors that the only summands in (8.5.1) annihilating $P_{\Lambda}$ are those in $\mathfrak{g l}(E) \oplus \mathfrak{g l}(E)$. Note that as a $\mathfrak{g l}(E)$-module, $\mathfrak{g l}(E)=\mathfrak{s l}(E) \oplus \mathbb{C}$ so one must test the highest weight vector of $\mathfrak{s l}(E)$ and $\mathbb{C}$.

### 8.5.2. The module structure of the equations for hypersurfaces

 with degenerate duals. Recall the equations for $\mathcal{D}_{k, d, N} \subset \mathbb{P}\left(S^{d} \mathbb{C}^{N *}\right)$ that we found in $\S 6.5 .3$. In this subsection I describe the module structure of those equations. It is technical and can be skipped on a first reading.Write $P=\sum_{J} \tilde{P}_{J} x^{J}$ with the sum over $|J|=d$. The weight of a monomial $\tilde{P}_{J_{0}} x^{J_{0}}$ is $J_{0}=\left(j_{1}, \ldots, j_{n}\right)$. Adopt the notation $[i]=(0, \ldots, 0,1,0, \ldots, 0)$ where the 1 is in the $i$-th slot and similarly for $[i, j]$ where there are two 1's. The entries of $P_{d-2,2}$ are, for $i \neq j,\left(P_{d-2,2}\right)_{i, j}=P_{I+[i, j]} x^{I}$, and for $i=j, P_{I+2[i]} x^{I}$, where $|I|=d-2$, and $P_{J}$ is $\tilde{P}_{J}$ with the coefficient adjusted, e.g., $P_{(d, 0, \ldots, 0)}=d(d-1) \tilde{P}_{(d, 0, \ldots, 0)}$ etc.. (This won't matter because we are only concerned with the weights of the coefficients, not their values.) To determine the highest weight vector, take $L=\operatorname{span}\left\{e_{1}, e_{2}\right\}$, $F=\operatorname{span}\left\{e_{1}, \ldots, e_{k+3}\right\}$. The highest weight term of

$$
\left.\left(\left.x_{1}^{e-d} P\right|_{L}\right) \wedge\left(\left.x_{1}^{e-d-1} x_{2} P\right|_{L}\right) \wedge \cdots \wedge\left(\left.x_{2}^{e-d} P\right|_{L}\right) \wedge\left(\operatorname{det}_{k+3}\left(\left.P_{d-2,2}\right|_{F}\right)\right)\right|_{L}
$$

is the coefficient of $x_{1}^{e} \wedge x_{1}^{e-1} x_{2} \wedge \cdots \wedge x_{1}^{e-(e-d+2)} x_{2}^{e-d+2}$. It will not matter how we distribute these for the weight, so take the coefficient of $x_{1}^{e}$ in
$\left.\left(\operatorname{det}_{k+3}\left(P_{d-2,2} \mid F\right)\right)\right|_{L}$. It has leading term

$$
P_{(d, 0, \ldots, 0)} P_{(d-2,2,0, \ldots, 0)} P_{(d-2,0,2,0, \ldots, 0)} \cdots P_{(d-2,0, \ldots, 0,2,0, \ldots, 0)}
$$

which is of weight $\left(d+(k+2)(d-2), 2^{k+2}\right)$. For each $\left(\left.x_{1}^{e-d-s} x_{2}^{s} P\right|_{L}\right)$ take the coefficient of $x_{1}^{e-s-1} x_{2}^{s+1}$ which has the coefficient of $P_{(d-1,1,0, \ldots, 0)}$ each time, to get a total weight contribution of $((e-d+1)(d-1),(e-d+1), 0, \ldots, 0)$ from these terms. Adding the weights together, and recalling that $e=$ $(k+3)(d-2)$ the highest weight is

$$
\left(d^{2} k+2 d^{2}-2 d k-4 d+1, d k+2 d-2 k-3,2^{k+1}\right),
$$

which may be written as

$$
\left((k+2)\left(d^{2}-2 d\right)+1,(k+2)(d-2)+1,2^{k+1}\right) .
$$

In summary:
Theorem 8.5.2.1. [LMR13] The ideal of the variety $\mathcal{D}_{k, d, N} \subset \mathbb{P}\left(S^{d} \mathbb{C}^{N *}\right)$ contains a copy of the $G L_{N}$-module $S_{\pi(k, d)} \mathbb{C}^{N}$, where

$$
\pi(k, d)=\left((k+2)\left(d^{2}-2 d\right)+1, d(k+2)-2 k-3,2^{k+1}\right) .
$$

Since $|\pi|=d(k+2)(d-1)$, these equations have degree $(k+2)(d-1)$.

Observe that the module $\pi(2 n-2, n)$ indeed satisfies the requirements to be ( $m, \frac{m^{2}}{2}$ )-GCT useful, as $p_{1}=2 n^{3}-2 n^{2}+1>n(n-m)$ and $l(\pi(2 n-$ $2, n))=2 n+1$.
8.5.3. Dual $l_{k, d, N}$ v. $\mathcal{D}_{k, d, N}$. Recall that Dual $_{k, d, N} \subset \mathbb{P} S^{d} \mathbb{C}^{N *}$ is the Zariski closure of the irreducible polynomials whose hypersurfaces have $k$-dimensional dual varieties. The following more refined information may be useful for studying permanent v. determinant:
Proposition 8.5.3.1. [LMR13] As subsets of $S^{d} \mathbb{C}^{N^{*}}$, Dual $l_{k, d, N}$ intersected with the irreducible hypersurfaces equals $\mathcal{D}_{k, d, N}$ intersected with the irreducible hypersurfaces.

Proof. Let $P \in \mathcal{D}_{k, d, N}$ be irreducible. For each $(L, F) \in G(2, F) \times G(k+$ $3, V)$ one obtains set-theoretic equations for the condition that $\left.P\right|_{L}$ divides $\left.Q\right|_{L}$, where $Q=\operatorname{det}\left(\left.P_{d-2,2}\right|_{F}\right)$. But $P$ divides $Q$ if and only if restricted to each plane $P$ divides $Q$, so these conditions imply that the dual variety of the irreducible hypersurface $\operatorname{Zeros}(P)$ has dimension at most $k$.

Theorem 8.5.3.2. [LMR13] $\mathcal{D e t}_{n}$ is an irreducible component of $\mathcal{D}_{2 n-2, n, n^{2}}$
The proof of Theorem 8.5.3.2 requires familiarity with Zariski tangent spaces to schemes. Here is an outline: Given two schemes, $X, Y$ with $X$ irreducible and $X \subseteq Y$, an equality of Zariski tangent spaces, $T_{x} X=T_{x} Y$
for some $x \in X_{\text {smooth }}$, implies that $X$ is an irreducible component of $Y$ (and in particular, if $Y$ is irreducible, that $X=Y$ ). The following theorem is a more precise version:
Theorem 8.5.3.3. [LMR13] The scheme $\mathcal{D}_{2 n-2, n, n^{2}}$ is smooth at $\left[\operatorname{det}_{n}\right]$, and $\mathcal{D e t}_{n}$ is an irreducible component of $\mathcal{D}_{2 n-2, n, n^{2}}$.

The idea of the proof is as follows: We need to show $T_{\left[\text {det }_{n}\right]} \mathcal{D}_{n, 2 n-2, n^{2}}=$ $T_{\left[d e t_{n}\right]} \mathcal{D e t}_{n}$. We already know $T_{\left[\operatorname{det}_{n}\right]} \mathcal{D e t}_{n} \subseteq T_{\left[d e t_{n}\right]} \mathcal{D}_{n, 2 n-2, n^{2}}$. Both of these vector spaces are $G_{\operatorname{det}_{n}}$-submodules of $S^{n}(E \otimes F)$. In 8.7.1.3 you will prove the Cauchy formula that $S^{n}(E \otimes F)=\bigoplus_{|\pi|=n} S_{\pi} E \otimes S_{\pi} F$.
Exercise 8.5.3.4: (2) Show that $\left[\operatorname{det}_{n}\right]=S_{1^{n}} E \otimes S_{1^{n}} F$ and $\hat{T}_{\operatorname{det}_{n}} \mathcal{D e t}_{n}=$ $S_{1^{n}} E \otimes S_{1^{n} F} \oplus S_{2,1^{n-1}} E \otimes S_{2,1^{n-1}} F$. ©

So as a $G L(E) \times G L(F)$-module, $T_{\left[d e t_{n}\right]} \mathcal{D} e t_{n}=S_{2,1^{n-2}} E \otimes S_{2,1^{n-2}} F$. The problem now becomes to show that none of the other modules in $S^{n}(E \otimes F)$ are in $T_{\left[d e t_{n}\right]} \mathcal{D}_{n, 2 n-2, n^{2}}$. To do this, it suffices to check a single point in each module. A first guess would be to check highest weight vectors, but these are not so easy to write down in any uniform manner. Fortunately in this case there is another choice, namely the immanants $I M_{\pi}$ defined by Littlewood [Lit06], the unique trivial representation of the diagonal $\mathfrak{S}_{n}$ in the weight $\left(\left(1^{n}\right),\left(1^{n}\right)\right)$ subspace of $S_{\pi} E \otimes S_{\pi} F$, and the proof in [LMR13] proceeds by checking that none of these other than $I M_{2,1^{n-2}}$ are contained in $T_{\left[d e t_{n}\right]} \mathcal{D}_{n, 2 n-2, n^{2}}$.

Theorem 8.5.3.3 implies that the $G L(W)$-module of highest weight $\pi(2 n-$ $2, n)$ given by Theorem 8.5.2.1 gives local equations at [ $\left.\operatorname{det}_{n}\right]$ of $\mathcal{D e t}_{n}$, of degree $2 n(n-1)$. Since $\operatorname{Sub}_{k}\left(S^{n} \mathbb{C}^{N}\right) \subset D u a l_{k, n, N}$, the zero set of the equations is strictly larger than $\mathcal{D e t}_{n}$. Recall that $\operatorname{dim} S u b_{k}\left(S^{n} \mathbb{C}^{n^{2}}\right)=$ $\left({ }_{n}^{k+n+1}\right)+(k+2)(N-k-2)-1$. For $k=2 n-2, N=n^{2}$, this is larger than the dimension of the orbit of $\left[\operatorname{det}_{n}\right]$, and therefore $D u a l_{2 n-2, n, n^{2}}$ is not irreducible.

### 8.6. Double-Commutant and algebraic Peter-Weyl Theorems

I now present the theory that will enable proofs of the statements in $\S 8.1$ and $\S 3.5$.
8.6.1. Algebras and their modules. For an algebra $\mathcal{A}$, and $a \in \mathcal{A}$ the space $\mathcal{A} a$ is a left ideal and a (left) $\mathcal{A}$-module.

Let $G$ be a finite group. Recall from $\S 3.5$. 1 the notation $\mathbb{C}[G]$ for the space of functions on $G$, and $\delta_{g} \in \mathbb{C}[G]$ for the function such that $\delta_{g}(h)=0$ for $h \neq g$ and $\delta_{g}(g)=1$. Define a representation $L: G \rightarrow G L(\mathbb{C}[G])$ by
$L(g) \delta_{h}=\delta_{g h}$ and extending the action linearly. Define a second representation $R: G \rightarrow G L(\mathbb{C}[G])$ by $R(g) \delta_{h}=\delta_{h g^{-1}}$. Thus $\mathbb{C}[G]$ is a $G \times G$-module under the representation $(L, R)$, and for all $c \in \mathbb{C}[G]$, the ideal $\mathbb{C}[G] c$ is a $G$-module under the action $L$.

A representation $\rho: G \rightarrow G L(V)$ induces an algebra homomorphism $\mathbb{C}[G] \rightarrow \operatorname{End}(V)$, and it is equivalent that $V$ is a $G$-module or a left $\mathbb{C}[G]-$ module.

A module $M$ (for a group, ring, or algebra) is simple if it has no proper submodules. The module $M$ is semi-simple if it may be written as the direct sum of simple modules. An algebra is completely reducible if all its modules are semi-simple. For groups alone I will continue to use the terminology irreducible for a simple module, completely reducible for a semi-simple module, and reductive for a group such that all its modules can be decomposed into a direct sum of irreducible modules.
Exercise 8.6.1.1: (2) Show that if $\mathcal{A}$ is completely reducible, $V$ is an $\mathcal{A}$ module with an $\mathcal{A}$-submodule $U \subset V$, then there exists an $\mathcal{A}$-invariant complement to $U$ in $V$ and a projection map $\pi: V \rightarrow U$ that is an $\mathcal{A}$ module map. ©
8.6.2. The double-commutant theorem. Our sought-after decomposition of $V^{\otimes d}$ as a $G L(V)$-module will be obtained by exploiting the fact that the actions of $G L(V)$ and $\mathfrak{S}_{d}$ on $V^{\otimes d}$ commute. In this subsection I discuss commuting actions in general, as this is also the basis of the generalized DFT used in the Cohn-Umans method $\S 3.5$, and the starting point of the program of [MS01, MS08]. References for this section are [Pro07, Chap. 6] and [GW09, §4.1.5]. Let $S \subset \operatorname{End}(V)$ be any subset. Define the centralizer or commutator of $S$ to be

$$
S^{\prime}:=\{X \in \operatorname{End}(V) \mid X s=s X \forall s \in S\}
$$

## Proposition 8.6.2.1.

(1) $S^{\prime} \subset \operatorname{End}(V)$ is a sub-algebra.
(2) $S \subset\left(S^{\prime}\right)^{\prime}$.

Exercise 8.6.2.2: (1!) Prove Proposition 8.6.2.1.
Theorem 8.6.2.3. [Double-Commutant Theorem] Let $\mathcal{A} \subset \operatorname{End}(V)$ be a completely reducible associative algebra. Then $\mathcal{A}^{\prime \prime}=\mathcal{A}$.

There is an ambiguity in the notation $S^{\prime}$ as it makes no reference to $V$, so instead introduce the notation $\operatorname{End}_{S}(V):=S^{\prime}$.

Proof. By Proposition 8.6.2.1, $\mathcal{A} \subseteq \mathcal{A}^{\prime \prime}$. To show the reverse inclusion, say $T \in \mathcal{A}^{\prime \prime}$. Fix a basis $v_{1}, \ldots, v_{\mathbf{v}}$ of $V$. Since the action of $T$ is determined by its action on a basis, we need to find $a \in \mathcal{A}$ such that $a v_{j}=T v_{j}$ for
$j=1, \ldots, \mathbf{v}$. Let $w:=v_{1} \oplus \cdots \oplus v_{\mathbf{v}} \in V^{\oplus \mathbf{v}}$ and consider the submodule $\mathcal{A} w \subseteq V^{\oplus \mathbf{v}}$. By Exercise 8.6.1.1, there exists an $\mathcal{A}$-invariant complement to this submodule and an $\mathcal{A}$-equivariant projection $\pi: V^{\oplus \mathbf{v}} \rightarrow \mathcal{A} w \subset V^{\oplus \mathbf{v}}$, that is, a projection $\pi$ that commutes with the action of $\mathcal{A}$, i.e., $\pi \in \operatorname{End}_{\mathcal{A}}\left(V^{\oplus \mathbf{v}}\right)$. Since $T \in \operatorname{End}_{\mathcal{A}}(V)$ and the action on $V^{\oplus \mathbf{v}}$ is diagonal, $T \in \operatorname{End}_{\mathcal{A}}\left(V^{\oplus \mathbf{v}}\right)$. We have $\pi(T w)=T(\pi(w))$ but $T(\pi(w))=T(w)=T v_{1} \oplus \cdots \oplus T v_{\mathbf{v}}$. But since $\pi(T w) \in \mathcal{A} w$, there must be some $a \in \mathcal{A}$ such that $a w=T(w)$, i.e., $a v_{1} \oplus \cdots \oplus a v_{\mathbf{v}}=T v_{1} \oplus \cdots \oplus T v_{\mathbf{v}}$, i.e., $a v_{j}=T v_{j}$ for $j=1, \ldots, \mathbf{v}$.

Burnside's theorem, stated in $\S 3.5$, has a similar proof:
Theorem 8.6.2.4. [Burnside] Let $\mathcal{A} \subseteq \operatorname{End}(V)$ be a finite dimensional simple sub-algebra of $\operatorname{End}(V)$ (over $\mathbb{C}$ ) acting irreducibly on a finite-dimensional vector space $V$. Then $\mathcal{A}=\operatorname{End}(V)$. More generally, a finite dimensional semi-simple associative algebra $\mathcal{A}$ over $\mathbb{C}$ is isomorphic to a direct sum of matrix algebras:

$$
\mathcal{A} \simeq M a t_{d_{1} \times d_{1}}(\mathbb{C}) \oplus \cdots \oplus M a t_{d_{q} \times d_{q}}(\mathbb{C})
$$

for some $d_{1}, \ldots, d_{q}$.
Proof. For the first assertion, we need to show that given $X \in \operatorname{End}(V)$, there exists $a \in \mathcal{A}$ such that $a v_{j}=X v_{j}$ for $v_{1}, \ldots, v_{\mathbf{v}}$ a basis of $V$. Now just imitate the proof of Theorem 8.6.2.3. For the second assertion, note that $\mathcal{A}$ is a direct sum of simple algebras.

Remark 8.6.2.5. A pessimist could look at this theorem as a disappointment: all kinds of interesting looking algebras over $\mathbb{C}$, such as the group algebra of a finite group, are actually just plain old matrix algebras in disguise. An optimist could view this theorem as stating there is a rich structure hidden in matrix algebras. We will determine the matrix algebra structure explicitly for the group algebra of a finite group.
8.6.3. Consequences for reductive groups. Let $S$ be a group or algebra and let $V, W$ be $S$-modules, adopt the notation $\operatorname{Hom}_{S}(V, W)$ for the space of $S$-module maps $V \rightarrow W$, i.e.,

$$
\begin{aligned}
\operatorname{Hom}_{S}(V, W): & =\{f \in \operatorname{Hom}(V, W) \mid s(f(v))=f(s(v)) \forall s \in S, v \in V\} \\
& =\left(V^{*} \otimes W\right)^{S} .
\end{aligned}
$$

Theorem 8.6.3.1. Let $G$ be a reductive group and let $V$ be a $G$-module. Then
(1) The commutator $\operatorname{End}_{G}(V)$ is a semi-simple algebra.
(2) The isotypic components of $G$ and $\operatorname{End}_{G}(V)$ in $V$ coincide.
(3) Let $U$ be one such isotypic component, say for irreducible representations $A$ of $G$ and $B$ of $\operatorname{End}_{G}(V)$. Then, as a $G \times \operatorname{End}_{G}(V)$-module,

$$
U=A \otimes B
$$

as an $\operatorname{End}_{G}(V)$-module

$$
B=\operatorname{Hom}_{G}(A, U),
$$

and as a $G$-module

$$
A=\operatorname{Hom}_{\operatorname{End}_{G}(V)}(B, U) .
$$

In particular, $\operatorname{mult}(A, V)=\operatorname{dim} B$ and $\operatorname{mult}(B, V)=\operatorname{dim} A$.
Example 8.6.3.2. Below we will see that $\operatorname{End}_{G L(V)}\left(V^{\otimes d}\right)=\mathbb{C}\left[\mathfrak{S}_{d}\right]$. As an $\mathfrak{S}_{3} \times G L(V)$-module, we have the decomposition $V^{\otimes 3}=\left([3] \otimes S^{3} V\right) \oplus$ $\left([2,1] \otimes S_{21} V\right) \oplus\left([1,1,1] \otimes \Lambda^{3} V\right)$ which illustrates Theorem 8.6.3.1.

To prove the theorem, we will need the following lemma:
Lemma 8.6.3.3. For $W \subset V$ a $G$-submodule and $f \in \operatorname{Hom}_{G}(W, V)$, there exists $a \in \operatorname{End}_{G}(V)$ such that $\left.a\right|_{W}=f$.

Proof. Consider the diagram


The vertical arrows are $G$-equivariant projections, and the horizontal arrows are restriction of domain of a linear map. The diagram is commutative. Since the vertical arrows and upper horizontal arrow are surjective, we conclude the lower horizontal arrow is surjective as well.

Proof of Theorem 8.6.3.1. I first prove (3): The space $\operatorname{Hom}_{G}(A, V)$ is an $\operatorname{End}_{G}(V)$-module because for $s \in \operatorname{Hom}_{G}(A, V)$ and $a \in \operatorname{End}_{G}(V)$, the composition as : $A \rightarrow V$ is still a $G$-module map. We need to show (i) that $\operatorname{Hom}_{G}(A, V)$ is an irreducible $\operatorname{End}_{G}(V)$-module and (ii) that the isotypic component of $A$ in $V$ is $A \otimes \operatorname{Hom}_{G}(A, V)$.

To show (i), it is sufficient to show that for all nonzero $s, t \in \operatorname{Hom}_{G}(A, V)$, there exists $a \in \operatorname{End}_{G}(V)$ such that at $=s$. Since $t A$ and $s A$ are isomorphic $G$-modules, by Lemma 8.6.3.3, there exists $a \in \operatorname{End}_{G}(V)$ extending an isomorphism between them, so $a(t A)=s A$, i.e., at : $A \rightarrow s A$ is an isomorphism. Consider the isomorphism $S: A \rightarrow s A$, given by $a \mapsto s a$, so $S^{-1} a t$ is a nonzero scalar $c$ times the identity. Then $\tilde{a}:=\frac{1}{c} a$ has the property that $\tilde{a} t=s$.

To see (ii), let $U$ be the isotypic component of $A$, so $U=A \otimes B$ for some vector space $B$. Let $b \in B$ and define a map $\tilde{b}: A \rightarrow V$ by $a \mapsto a \otimes b$, which
is a $G$-module map where the action of $G$ on the target is just the action on the first factor. Thus $B \subseteq \operatorname{Hom}_{G}(A, V)$. Any $G$-module map $A \rightarrow V$ by definition has image in $U$, so equality holds.
(3) implies (2).

To see (1), note that $\operatorname{End}_{G}(V)$ is semi-simple because if the irreducible $G \times \operatorname{End}_{G}(V)$-components of $V$ are $U_{i}$, then $\operatorname{End}_{G}(V)=\oplus_{i} \operatorname{End}_{G}\left(U_{i}\right)=$ $\oplus_{i} \operatorname{End}_{G}\left(A_{i} \otimes B_{i}\right)=\oplus_{i} \operatorname{End}\left(B_{i}\right)$.
8.6.4. Matrix coefficients. For affine algebraic reductive groups, one can obtain all their (finite dimensional) irreducible representations from the ring of regular functions on $G$, denoted $\mathbb{C}[G]$. Here $G$ is an affine algebraic variety, i.e., a subvariety of $\mathbb{C}^{N}$ for some $N$, so $\mathbb{C}[G]=\mathbb{C}\left[x_{1}, \ldots, x_{N}\right] / I(G)$.
Exercise 8.6.4.1: (1!) Show that $G L_{n}$ is an affine algebraic subvariety of $\mathbb{C}^{n^{2}+1}$ with coordinates $\left(x_{j}^{i}, z\right)$ by considering the polynomial $z \operatorname{det}_{n}\left(x_{j}^{i}\right)-1$.

Thus $\mathbb{C}[G L(W)]$ may be defined to be the restriction of polynomial functions on $\mathbb{C}^{n^{2}+1}$ to the subvariety isomorphic to $G L(W)$. (For a finite group, all complex-valued functions on $G$ are algebraic, so this is consistent with our earlier notation.) If $G \subset G L(W)$ is defined by algebraic equations, this also enables us to define $\mathbb{C}[G]$ because $G \subset G L(W)$ is a subvariety. In this section and the next, we study the structure of $\mathbb{C}[G]$ as a $G$-module.

Let $G$ be an affine algebraic group. Let $\rho: G \rightarrow G L(V)$ be a finite dimensional representation of $G$. Define a map $i_{V}: V^{*} \otimes V \rightarrow \mathbb{C}[G]$ by $i_{V}(\alpha \otimes v)(g):=\alpha(\rho(g) v)$. The space of functions $i_{V}\left(V^{*} \otimes V\right)$ is called the space of matrix coefficients of $V$.
Exercise 8.6.4.2: (1)
i) Show $i_{V}$ is a $G \times G$-module map.
ii) Show that if $V$ is irreducible, $i_{V}$ is injective. ©
iii) If we choose a basis $v_{1}, \ldots, v_{\mathbf{v}}$ of $V$ with dual basis $\alpha^{1}, \ldots, \alpha^{\mathbf{v}}$, then $i_{V}\left(\alpha^{i} \otimes v_{j}\right)(g)$ is the $(i, j)$-th entry of the matrix representing $\rho(g)$ in this basis (which explains the name "matrix coefficients").
iv) Compute the matrix coefficient basis of the three irreducible representations of $\mathfrak{S}_{3}$ in terms of the standard basis $\left\{\delta_{\sigma} \mid \sigma \in \mathfrak{S}_{3}\right\}$.
v) Let $G=G L_{2} \mathbb{C}$, write $g=\left(\begin{array}{ll}a & b \\ c & d\end{array}\right) \in G$, and compute the matrix coefficient basis as functions of $a, b, c, d$ when $V=S^{2} \mathbb{C}^{2}, S^{3} \mathbb{C}^{2}$ and $\Lambda^{2} \mathbb{C}^{2}$.

Theorem 8.6.4.3. Let $G$ be an affine algebraic group and let $V$ be an irreducible $G$-module. Then $i_{V}\left(V^{*} \otimes V\right)$ equals the isotypic component of
type $V$ in $\mathbb{C}[G]$ under the action $L$ and the isotypic component of $V^{*}$ in $\mathbb{C}[G]$ under the action $R$.

Proof. It suffices to prove one of the assertions, consider the action $L$. Let $j: V \rightarrow \mathbb{C}[G]$ be a $G$-module map under the action $L$. We need to show $j(V) \subset i_{V}\left(V^{*} \otimes V\right)$. Define $\alpha \in V^{*}$ by $\alpha(v):=j(v)\left(\operatorname{Id}_{G}\right)$. Then $j(v)=$ $i_{V}(\alpha \otimes v)$, as $j(v) g=j(v)\left(g \cdot \operatorname{Id}_{G}\right)=j(g v)\left(\operatorname{Id}_{G}\right)=\alpha(g v)=i_{V}(\alpha \otimes v) g$.
8.6.5. Application to representations of finite groups. Theorem 8.6.4.3 implies:
Theorem 8.6.5.1. Let $G$ be a finite group, then as a $G \times G$-module under the action $(L, R)$ and as an algebra,

$$
\begin{equation*}
\mathbb{C}[G]=\bigoplus_{i} V_{i} \otimes V_{i}^{*} \tag{8.6.1}
\end{equation*}
$$

where the sum is over all the distinct irreducible representations of $G$.
Exercise 8.6.5.2: (1!) Let $G$ be a finite group and $H$ a subgroup. For the homogeneous space $G / H$, show that $\mathbb{C}[G / H]=\bigoplus_{i} V_{i}^{*} \otimes\left(V_{i}\right)^{H}$ as a $G$ module under the action $L$.
8.6.6. The algebraic Peter-Weyl Theorem. Theorem 8.6.5.1 generalizes to reductive algebraic groups. The proof is unchanged, except that one has an infinite sum:
Theorem 8.6.6.1. Let $G$ be a reductive algebraic group. Then there are only countably many non-isomorphic finite dimensional irreducible G-modules. Let $\Lambda_{G}^{+}$denote a set indexing the irreducible $G$-modules, and for $\lambda \in \Lambda_{G}^{+}$, let $V_{\lambda}$ denote the irreducible module associated to $\lambda$. Then, as a $G \times G$-module

$$
\mathbb{C}[G]=\bigoplus_{\lambda \in \Lambda_{G}^{+}} V_{\lambda} \otimes V_{\lambda}^{*}
$$

Corollary 8.6.6.2. Let $H \subset G$ be a closed subgroup. Then, as a $G$-module, the coordinate ring of the homogeneous space $G / H$ is

$$
\begin{equation*}
\mathbb{C}[G / H]=\mathbb{C}[G]^{H}=\bigoplus_{\lambda \in \Lambda_{G}^{+}} V_{\lambda} \otimes\left(V_{\lambda}^{*}\right)^{H}=\bigoplus_{\lambda \in \Lambda_{G}^{+}} V_{\lambda}^{\oplus \operatorname{dim}\left(V_{\lambda}^{*}\right)^{H}} \tag{8.6.2}
\end{equation*}
$$

Here $G$ acts on the $V_{\lambda}$ and $\left(V_{\lambda}^{*}\right)^{H}$ is just a vector space whose dimension records the multiplicity of $V_{\lambda}$ in $\mathbb{C}[G / H]$.
Exercise 8.6.6.3: (2!) Use Corollary 8.6.6.2 to determine $\mathbb{C}\left[v_{d}(\mathbb{P} V)\right]$ (even if you already know it by a different method).
8.6.7. Characters and representations of finite groups. Let $\rho: G \rightarrow$ $G L(V)$ be a representation. Define a function $\chi_{\rho}: G \rightarrow \mathbb{C}$ by $\chi_{\rho}(g)=$ $\operatorname{trace}(\rho(g))$. The function $\chi_{\rho}$ is called the character of $\rho$.
Exercise 8.6.7.1: (1) Show that $\chi_{\rho}$ is constant on conjugacy classes of $G$.
A function $f: G \rightarrow \mathbb{C}$ such that $f\left(h g h^{-1}\right)=f(g)$ for all $g, h \in G$ is called a class function.
Exercise 8.6.7.2: (1) For representations $\rho_{j}: G \rightarrow G L\left(V_{j}\right)$, show that $\chi_{\rho_{1} \oplus \rho_{2}}=\chi_{\rho_{1}}+\chi_{\rho_{2}}$.
Exercise 8.6.7.3: (1) Given $\rho_{j}: G \rightarrow G L\left(V_{j}\right)$ for $j=1,2$, define $\rho_{1} \otimes \rho_{2}$ : $G \rightarrow G L\left(V_{1} \otimes V_{2}\right)$ by $\rho_{1} \otimes \rho_{2}(g)\left(v_{1} \otimes v_{2}\right)=\rho_{1}(g) v_{1} \otimes \rho_{2}(g) v_{2}$. Show that $\chi_{\rho_{1} \otimes \rho_{2}}=$ $\chi_{\rho_{1}} \chi_{\rho_{2}}$.

Theorem 8.6.5.1 is not yet useful, as we do not yet know what the $V_{i}$ are. Let $\mu_{i}: G \rightarrow G L\left(V_{i}\right)$ denote the representation. It is not difficult to show that the functions $\chi_{\mu_{i}}$ are linearly independent in $\mathbb{C}[G]$. (One uses a $G$-invariant Hermitian inner-product $\left\langle\chi_{V}, \chi_{W}\right\rangle:=\frac{1}{|G|} \sum_{g \in G} \chi_{V}(g) \overline{\chi_{W}(g)}$ and shows that they are orthogonal with respect to this inner-product, see, e.g., [FH91, §2.2].) On the other hand, we have a natural basis of the class functions, namely the $\delta$-functions on each conjugacy class. Let $C_{j}$ be a conjugacy class of $G$ and define $\delta_{C_{j}}:=\sum_{g \in C_{j}} \delta_{g}$. It is straightforward to see, via the DFT (§3.5.1), that the span of the $\delta_{C_{j}}$ 's equals the span of the $\chi_{\mu_{i}}$ 's, that is the number of distinct irreducible representations of $G$ equals the number of conjugacy classes (see, e.g., [FH91, §2.2] for the standard proof using the Hermitian inner-product on class functions and [GW09, §4.4] for a DFT proof).

Remark 8.6.7.4. The classical Heisenberg uncertainty principle from physics, in the language of mathematics, is that it is not possible to localize both a function and its Fourier transform. A discrete analog of this uncertainty principle holds, in that the transforms of the delta functions have large support in terms of matrix coefficients and vice versa. In particular, the relation between these two bases can be complicated.
8.6.8. Representations of $\mathfrak{S}_{d}$. When $G=\mathfrak{S}_{d}$, we get lucky: one may associate irreducible representations directly to conjugacy classes.

The conjugacy class of a permutation is determined by its decomposition into a product of disjoint cycles. The conjugacy classes of $\mathfrak{S}_{d}$ are in 1-1 correspondence with the set of partitions of $d$ : to a partition $\pi=\left(p_{1}, \ldots, p_{r}\right)$ one associates the conjugacy class of an element with disjoint cycles of lengths $p_{1}, \ldots, p_{r}$. Let [ $\pi$ ] denote the isomorphism class of the irreducible $\mathfrak{S}_{d}$-module associated to the partition $\pi$. In summary:

Proposition 8.6.8.1. The irreducible representations of $\mathfrak{S}_{d}$ are indexed by partitions of $d$.

Thus as an $\mathfrak{S}_{d} \times \mathfrak{S}_{d}$ module under the $(L, R)$-action:

$$
\begin{equation*}
\mathbb{C}\left[\mathfrak{S}_{d}\right]=\bigoplus_{|\pi|=d}[\pi]_{L}^{*} \otimes[\pi]_{R} . \tag{8.6.3}
\end{equation*}
$$

We can say even more: as $\mathfrak{S}_{d}$ modules, $[\pi]$ is isomorphic to $[\pi]^{*}$. This is usually proved by first noting that for any finite group $G$, and any irreducible representation $\mu, \chi_{\mu^{*}}=\overline{\chi_{\mu}}$ where the overline denotes complex conjugate and then observing that the characters of $\mathfrak{S}_{d}$ are all real-valued functions. Thus we may rewrite (8.6.3) as

$$
\begin{equation*}
\mathbb{C}\left[\mathfrak{S}_{d}\right]=\bigoplus_{|\pi|=d}[\pi]_{L} \otimes[\pi]_{R} . \tag{8.6.4}
\end{equation*}
$$

Exercise 8.6.8.2: (1) Show $[d] \subset[\pi] \otimes[\mu]$ if and only if $\pi=\mu$. ©
Exercise 8.6.8.3: (1) Show that moreover $[d] \subset[\pi] \otimes[\pi]$ with multiplicity one. ©

### 8.7. Representations of $\mathfrak{S}_{d}$ and $G L(V)$

In this section we finally obtain our goal of the decomposition of $V^{\otimes d}$ as a $G L(V)$-module.
8.7.1. Schur-Weyl duality. We have already seen that the actions of $G L(V)$ and $\mathfrak{S}_{d}$ on $V^{\otimes d}$ commute.
Proposition 8.7.1.1. $\operatorname{End}_{G L(V)}\left(V^{\otimes d}\right)=\mathbb{C}\left[\mathfrak{S}_{d}\right]$.
Proof. We will show that $\operatorname{End}_{\mathbb{C}\left[\mathfrak{S}_{d}\right]}\left(V^{\otimes d}\right)$ is the algebra generated by $G L(V)$ and conclude by the double commutant theorem. Since

$$
\begin{aligned}
\operatorname{End}\left(V^{\otimes d}\right) & =V^{\otimes d} \otimes\left(V^{\otimes d}\right)^{*} \\
& \simeq\left(V \otimes V^{*}\right)^{\otimes d}
\end{aligned}
$$

under the re-ordering isomorphism, $\operatorname{End}\left(V^{\otimes d}\right)$ is spanned by elements of the form $X_{1} \otimes \cdots \otimes X_{d}$ with $X_{j} \in \operatorname{End}(V)$, i.e., elements of $\hat{S e g}(\mathbb{P}(\operatorname{End}(V)) \times$ $\cdots \times \mathbb{P}(\operatorname{End}(V)))$. The action of $X_{1} \otimes \cdots \otimes X_{d}$ on $v_{1} \otimes \cdots \otimes v_{d}$ induced from the $G L(V)^{\times d}$-action is $v_{1} \otimes \cdots \otimes v_{d} \mapsto\left(X_{1} v_{1}\right) \otimes \cdots \otimes\left(X_{d} v_{d}\right)$. Since $g \in$ $G L(V)$ acts by $g \cdot\left(v_{1} \otimes \cdots \otimes v_{d}\right)=g v_{1} \otimes \cdots \otimes g v_{d}$, the image of $G L(V)$ in $\left(V \otimes V^{*}\right)^{\otimes d}$ lies in $S^{d}\left(V \otimes V^{*}\right)$, in fact it is a Zariski open subset of $\hat{v}_{d}\left(\mathbb{P}\left(V \otimes V^{*}\right)\right)$ which spans $S^{d}\left(V \otimes V^{*}\right)$. In other words, the algebra generated by $G L(V)$ is $S^{d}\left(V \otimes V^{*}\right) \subset \operatorname{End}\left(V^{\otimes d}\right)$. But by definition $S^{d}\left(V \otimes V^{*}\right)=\left[\left(V \otimes V^{*}\right)^{\otimes d}\right] \mathfrak{S}_{d}$ and we conclude.

Theorem 8.6.3.1 and Proposition 8.7.1.1 imply:
Theorem 8.7.1.2. [Schur-Weyl duality] The irreducible decomposition of $V^{\otimes d}$ as a $G L(V) \times \mathbb{C}\left[\mathfrak{S}_{d}\right]$-module (equivalently, as a $G L(V) \times \mathfrak{S}_{d}$-module) is

$$
\begin{equation*}
V^{\otimes d}=\underset{|\pi|=d}{\bigoplus} S_{\pi} V \otimes[\pi] \tag{8.7.1}
\end{equation*}
$$

where $S_{\pi} V:=\operatorname{Hom}_{\mathfrak{S}_{d}}\left([\pi], V^{\otimes d}\right)$ is an irreducible $G L(V)$-module.
Note that as far as we know, $S_{\pi} V$ could be zero. (It will be zero whenever $l(\pi) \geq \operatorname{dim} V$.)
Exercise 8.7.1.3: (2) Show that as a $G L(E) \times G L(F)$-module, $S^{d}(E \otimes F)=$ $\bigoplus_{|\pi|=d} S_{\pi} E \otimes S_{\pi} F$. This is called the Cauchy formula. ©
8.7.2. Explicit realizations of representations of $\mathfrak{S}_{d}$ and $G L(V)$. By Theorem 8.6.5.1 we may explicitly realize each irreducible $\mathfrak{S}_{d}$-module via some projection from $\mathbb{C}\left[\mathfrak{S}_{d}\right]$. The question is, which projections?

Given $\pi$ we would like to find elements $c_{\bar{\pi}} \in \mathbb{C}\left[\mathfrak{S}_{d}\right]$ such that $\mathbb{C}\left[\mathfrak{S}_{d}\right] c_{\bar{\pi}}$ is isomorphic to $[\pi]$. I write $\bar{\pi}$ instead of just $\pi$ because the elements are far from unique; there is a vector space of dimension $\operatorname{dim}[\pi]$ of such projection operators by Theorem 8.6.5.1, and the overline signifies a specific realization. In other words, the $\mathfrak{S}_{d}$-module map $R M_{c_{\bar{\pi}}}: \mathbb{C}\left[\mathfrak{S}_{d}\right] \rightarrow \mathbb{C}\left[\mathfrak{S}_{d}\right], f \mapsto f c_{\bar{\pi}}$ should kill all $\mathfrak{S}_{d}^{R}$-modules not isomorphic to $[\pi]_{R}$, and the image should be $[\pi]_{L} \otimes z$ for some $z \in[\pi]_{R}$. If this works, as a bonus, the map $c_{\bar{\pi}}: V^{\otimes d} \rightarrow V^{\otimes d}$ induced from the $\mathfrak{S}_{d^{-}}$-action will have image $S_{\bar{\pi}} V \otimes z \simeq S_{\bar{\pi}} V$ for the same reason, where $S_{\bar{\pi}} V$ is some realization of $S_{\pi} V$ and $z \in[\pi]$.

Here are projection operators for the two representations we understand well:

When $\pi=(d)$, there is a unique up to scale $c_{\overline{(d)}}$ and it is easy to see it must be $c_{\overline{(d)}}:=\sum_{\sigma \in \mathfrak{S}_{d}} \delta_{\sigma}$, as the image of $R M_{c_{\overline{(d)}}}$ is clearly the line through $c_{\overline{(d)}}$ on which $\mathfrak{S}_{d}$ acts trivially. Note further that $c_{\overline{(d)}}\left(V^{\otimes d}\right)=S^{d} V$ as desired.

When $\pi=\left(1^{d}\right)$, again we have a unique up to scale projection, and its clear we should take $c_{\overline{\left(1^{d}\right)}}=\sum_{\sigma \in \mathfrak{S}_{d}} \operatorname{sgn}(\sigma) \delta_{\sigma}$ as the image of any $\delta_{\tau}$ will be $\operatorname{sgn}(\tau) c_{\overline{\left(1^{d}\right)}}$, and $c_{\overline{\left(1^{d}\right)}}\left(V^{\otimes d}\right)=\Lambda^{d} V$.

The only other representation of $\mathfrak{S}_{d}$ that we have a reasonable understanding of is the standard representation $\pi=(d-1,1)$ which corresponds to the complement of the trivial representation in the permutation action on $\mathbb{C}^{d}$. A basis of this space could be given by $e_{1}-e_{d}, e_{2}-e_{d}, \ldots, e_{d-1}-e_{d}$. Note that the roles of $1, \ldots, d-1$ in this basis are the "same" in that if one permutes them, one gets the same basis, and that the role of $d$ with respect
to any of the other $e_{j}$ is "skew" in some sense. To capture this behavior, consider

$$
c_{\overline{(d-1,1)}}:=\left(\delta_{\mathrm{Id}}-\delta_{(1, d)}\right)\left(\sum_{\sigma \in \mathfrak{G}_{d-1}[d-1]} \delta_{\sigma}\right)
$$

where $\mathfrak{S}_{d-1}[d-1] \subset \mathfrak{S}_{d}$ is the subgroup permuting the elements $\{1, \ldots, d-$ $1\}$. Note that $c_{(d-1,1)} \delta_{\tau}=c_{\overline{(d-1,1)}}$ for any $\tau \in \mathfrak{S}_{d-1}[d-1]$ so the image is of dimension at most $d=\operatorname{dim}\left(\mathbb{C}\left[\mathfrak{S}_{d}\right] / \mathbb{C}\left[\mathfrak{S}_{d-1}\right]\right)$.
Exercise 8.7.2.1: (2) Show that the image is $d-1$ dimensional.
Now consider $R M_{c_{(d-1,1)}}\left(V^{\otimes d}\right)$ : after re-orderings, it is the image of the composition of the maps

$$
V^{\otimes d} \rightarrow V^{\otimes d-2} \otimes \Lambda^{2} V \rightarrow S^{d-1} V \otimes V
$$

In particular, in the case $d=3$, it is the image of

$$
V \otimes \Lambda^{2} V \rightarrow S^{2} V \otimes V
$$

which is isomorphic to $S_{21} V$, as was mentioned in in §4.1.5.
Here is the general recipe to construct an $\mathfrak{S}_{d}$-module isomorphic to $[\pi]$ : fill the Young diagram of a partition $\pi$ of $d$ with integers $1, \ldots, d$ from top to bottom and left to right. For example let $\pi=(4,2,1)$ and write:

$$
\begin{equation*}
 \tag{8.7.2}
\end{equation*}
$$

Define $\mathfrak{S}_{\bar{\pi}^{\prime}} \simeq \mathfrak{S}_{q_{1}} \times \cdots \times \mathfrak{S}_{q_{p_{1}}} \subset \mathfrak{S}_{d}$ to be the subgroup that permutes elements in each column and $\mathfrak{S}_{\bar{\pi}}$ is the subgroup of $\mathfrak{S}_{d}$ that permutes the elements in each row.

Explicitly, writing $\pi=\left(p_{1}, \ldots, p_{q_{1}}\right)$ and $\pi^{\prime}=\left(q_{1}, \ldots, q_{p_{1}}\right), \mathfrak{S}_{q_{1}}$ permutes the elements of $\left\{1, \ldots, q_{1}\right\}, \mathfrak{S}_{q_{2}}$ permutes the elements of $\left\{q_{1}+1, \ldots, q_{1}+q_{2}\right\}$ etc.. Similarly, $\mathfrak{S}_{\pi} \simeq \mathfrak{S}_{p_{1}} \times \cdots \times \mathfrak{S}_{p_{\ell}} \subset \mathfrak{S}_{d}$ is the subgroup where $\mathfrak{S}_{p_{1}}$ permutes the elements $\left\{1, q_{1}+1, q_{1}+q_{2}+1, \ldots, q_{1}+\cdots+q_{p_{1}-1}+1\right\}$, $\mathfrak{S}_{p_{2}}$ permutes the elements $\left\{2, q_{1}+2, q_{1}+q_{2}+2, \ldots, q_{1}+\cdots+q_{p_{1}-1}+2\right\}$ etc..

Define two elements of $\mathbb{C}\left[\mathfrak{S}_{d}\right]: s_{\bar{\pi}}:=\sum_{\sigma \in \mathfrak{G}_{\bar{\pi}}} \delta_{\sigma}$ and $a_{\bar{\pi}}:=\sum_{\sigma \in \mathfrak{S}_{\pi^{\prime}}} \operatorname{sgn}(\sigma) \delta_{\sigma}$. Fact: Then $[\pi]$ is the isomorphism class of the $\mathfrak{S}_{d}$-module $\mathbb{C}\left[\mathfrak{S}_{d}\right] a_{\bar{\pi}} S_{\bar{\pi}}$. (It is also the isomorphism class of $\mathbb{C}\left[\mathfrak{S}_{d}\right] s_{\bar{\pi}} a_{\bar{\pi}}$, although these two realizations are generally distinct.)
Exercise 8.7.2.2: (1) Show that $\left[\pi^{\prime}\right]=[\pi] \otimes\left[1^{d}\right]$ as $\mathfrak{S}_{d}$-modules. ©
The action on $V^{\otimes d}$ is first to map it to $\Lambda^{q_{1}} V \otimes \cdots \otimes \Lambda^{q_{p_{1}}} V$, and then the module $S_{\pi} V$ is realized as the image of a map from this space to
$S^{p_{1}} V \otimes \cdots \otimes S^{p_{q_{1}}} V$ obtained by re-ordering then symmetrizing. So despite their original indirect definition, we may realize the modules $S_{\pi} V$ explicitly simply by skew-symmetrizations and symmetrizations.

Other realizations of $S_{\pi} V$ (resp. highest weight vectors for $S_{\pi} V$, in fact a basis of them) can be obtained by letting $\mathfrak{S}_{d}$ act on $R M_{c_{\bar{\pi}}} V^{\otimes d}$ (resp. the highest weight vector of $\left.R M_{c_{\pi}} V^{\otimes d}\right)$.

Example 8.7.2.3. Consider $c_{\overline{(2,2)}}$, associated to

$$
\begin{array}{|l|l|}
\hline 1 & 3  \tag{8.7.3}\\
\hline 2 & 4 \\
\hline
\end{array}
$$

which realizes a copy of $S_{(2,2)} V \subset V^{\otimes 4}$. It first maps $V^{\otimes 4}$ to $\Lambda^{2} V \otimes \Lambda^{2} V$ and then maps that to $S^{2} V \otimes S^{2} V$. Explicitly, the maps are

$$
\begin{aligned}
a \otimes b \otimes c \otimes c & \mapsto(a \otimes b-b \otimes a) \otimes(c \otimes d-d \otimes c)=a \otimes b \otimes c \otimes d-a \otimes b \otimes d \otimes c-b \otimes a \otimes c \otimes d+b \otimes a \otimes d \otimes c \\
& \mapsto(a \otimes b \otimes c \otimes d+c \otimes b \otimes a \otimes d+a \otimes d \otimes c \otimes b+c \otimes d \otimes a \otimes b) \\
& -(a \otimes b \otimes d \otimes c+d \otimes b \otimes a \otimes c+a \otimes c \otimes d \otimes b+d \otimes c \otimes a \otimes b) \\
& -(b \otimes a \otimes c \otimes d+c \otimes a \otimes b \otimes d+b \otimes d \otimes c \otimes a+c \otimes d \otimes b \otimes a) \\
& +(b \otimes a \otimes d \otimes c+d \otimes a \otimes b \otimes c+b \otimes c \otimes d \otimes a+d \otimes c \otimes b \otimes a)
\end{aligned}
$$

Exercise 8.7.2.4: (2) Show that a basis of the highest weight space of $[2,1] \otimes S_{21} V \subset V^{\otimes 3}$ is $v_{1}=e_{1} \wedge e_{2} \otimes e_{1}$ and $v_{2}=e_{1} \otimes e_{1} \wedge e_{2}$. Let $\mathbb{Z}_{3} \subset \mathfrak{S}_{3}$ be the cyclic permutation of the three factors in $V^{\otimes 3}$ and show that $\omega v_{1} \pm \omega^{2} v_{2}$ are eigenvectors for this action with eigenvalues $\omega, \omega^{2}$, where $\omega=e^{\frac{2 \pi i}{3}}$.

### 8.8. The program of [MS01, MS08]

Algebraic geometry was used successfully in [Mul99] to prove lower bounds in the "PRAM model without bit operations" (the model is defined in [Mul99]), and the proof indicated that algebraic geometry, more precisely invariant theory, could be used to resolve the $\mathbf{P}$ v. NC problem (a cousin of permanent v . determinant). This was investigated further in [MS01, MS08] and numerous sequels. In this section I present the program outlined in [MS08], as refined in [BLMW11].

Independent of its viability, I expect the ingredients that went into the program of [MS01, MS08] will play a role in future investigations regarding Valiant's conjecture and thus are still worth studying.
8.8.1. Preliminaries. Let $W=\mathbb{C}^{n^{2}}$. Recall $\mathbb{C}\left[\hat{\mathcal{D e t}} t_{n}\right]:=\operatorname{Sym}\left(S^{n} W^{*}\right) / I\left(\mathcal{D e t}_{n}\right)$, the homogeneous coordinate ring of the (cone over) $\mathcal{D e t}_{n}$. This is the space of polynomial functions on $\hat{\mathcal{D e t}}{ }_{n}$ inherited from polynomials on the ambient space.

Since $I\left(\mathcal{D e t}_{n}\right) \subset \operatorname{Sym}\left(S^{n} W^{*}\right)$ is a $G L(W)$-submodule, and since $G L(W)$ is reductive, we obtain the following splitting as a $G L(W)$-module:

$$
\operatorname{Sym}\left(S^{n} W^{*}\right)=I\left(\mathcal{D e t}_{n}\right) \oplus \mathbb{C}\left[\hat{\mathcal{D e t}_{n}}\right] .
$$

In particular, if a module $S_{\pi} W^{*}$ appears in $\operatorname{Sym}\left(S^{n} W^{*}\right)$ and it does not appear in $\mathbb{C}\left[\hat{\mathcal{D e}_{n}}\right]$, it must appear in $I\left(\mathcal{D e t}_{n}\right)$.

Now consider

$$
\mathbb{C}\left[G L(W) \cdot \operatorname{det}_{n}\right]=\mathbb{C}\left[G L(W) / G_{\operatorname{det}_{n}}\right]=\mathbb{C}[G L(W)]^{G_{\operatorname{det}_{n}}} .
$$

There is an injective map

$$
\mathbb{C}\left[\hat{\mathcal{D e t}_{n}}\right] \rightarrow \mathbb{C}\left[G L(W) \cdot \operatorname{det}_{n}\right]
$$

given by restriction of functions. The map is an injection because any function identically zero on a Zariski open subset of an irreducible variety is identically zero on the variety.

Corollary 8.6.6.2 in principle gives a recipe to determine the modules in $\mathbb{C}\left[G L(W) \cdot \operatorname{det}_{n}\right]$, which motivates the following plan:
Plan : Find a module $S_{\pi} W^{*}$ not appearing in $\mathbb{C}\left[G L(W) / G_{\operatorname{det}_{n}}\right]$ that does appear in $\operatorname{Sym}\left(S^{n} W^{*}\right)$.

By the above discussion such a module must appear in $I\left(\mathcal{D e t}_{n}\right)$.
Definition 8.8.1.1. An irreducible $G L(W)$-module $S_{\pi} W^{*}$ appearing in $\operatorname{Sym}\left(S^{n} W^{*}\right)$ and not appearing in $\mathbb{C}\left[G L(W) / G_{\mathrm{det}_{n}}\right]$ is called an orbit occurrence obstruction.

The precise condition a module must satisfy in order to not occur in $\mathbb{C}\left[G L(W) / G_{\text {det }_{n}}\right]$ is explained in Proposition 8.8.2.2. The discussion in $\S 8.4$ shows that in order to be useful, $\pi$ must have a large first part and few parts.

One might object that the coordinate rings of different orbits could coincide, or at least be very close. Indeed this is the case for generic polynomials, but in GCT one generally restricts to polynomials whose symmetry groups characterize the orbit in the sense of Definition 1.2.5.3. We have seen in $\S 6.6$ that both the determinant and permanent polynomials are characterized by their stabilizers.

Corollary 8.6.6.2 motivates the study of polynomials characterized by their stabilizers: if $P \in V$ is characterized by its stabilizer, then $G \cdot P$ is the unique orbit in $V$ with coordinate ring isomorphic to $\mathbb{C}[G \cdot P]$ as a $G$ module. Thus one can think of polynomial sequences that are complete for their complexity classes and are characterized by their stabilizers as "best" representatives of their class.

Remark 8.8.1.2. All $G L(W)$-modules $S_{\left(p_{1}, \ldots, p_{\mathbf{w}}\right)} W$ may be graded using $p_{1}+\cdots+p_{\mathbf{w}}$ as the grading. One does not have such a grading for $S L(W)-$ modules, which makes their use in GCT more difficult. In [MS01, MS08], it was proposed to use the $S L(W)$-module structure because it had the advantage that the $S L$-orbit of $\operatorname{det}_{n}$ is already closed. The disadvantage from the lack of a grading appears to outweigh this advantage.
8.8.2. The coordinate ring of $G L_{n^{2}} \cdot \operatorname{det}_{n}$. Write $\mathbb{C}^{n^{2}}=E \otimes F$, with $E, F=\mathbb{C}^{n}$. I first compute the $S L(E) \times S L(F)$-invariants in $S_{\pi}(E \otimes F)$ where $|\pi|=d=\delta n$. Recall from $\S 8.7 .1$ that by definition, $S_{\pi} W=\operatorname{Hom}_{\mathfrak{S}_{d}}\left([\pi], W^{\otimes d}\right)$. Thus

$$
\begin{aligned}
S_{\pi}(E \otimes F) & =\operatorname{Hom}_{\mathfrak{S}_{d}}\left([\pi], E^{\otimes d} \otimes F^{\otimes d}\right) \\
& \left.=\operatorname{Hom}_{\mathfrak{S}_{d}}[\pi],\left(\bigoplus_{|\mu|=d}[\mu] \otimes S_{\mu} E\right) \otimes\left(\bigoplus_{|\nu|=d}[\nu] \otimes S_{\nu} F\right)\right) \\
& =\bigoplus_{|\mu|=|\nu|=d} \operatorname{Hom}_{\mathfrak{S}_{d}}([\pi],[\mu] \otimes[\nu]) \otimes S_{\mu} E \otimes S_{\nu} F
\end{aligned}
$$

The vector space $\operatorname{Hom}_{\mathfrak{S}_{d}}([\pi],[\mu] \otimes[\nu])$ simply records the multiplicity of $S_{\mu} E \otimes S_{\nu} F$ in $S_{\pi}(E \otimes F)$. The numbers $k_{\pi, \mu, \nu}=\operatorname{dim} \operatorname{Hom}_{\mathfrak{S}_{d}}([\pi],[\mu] \otimes[\nu])$ are called Kronecker coefficients.
Exercise 8.8.2.1: (2) Show that

$$
k_{\pi, \mu, \nu}=\operatorname{Hom}_{\mathfrak{S}_{d}}([d],[\pi] \otimes[\mu] \otimes[\nu])=\operatorname{mult}\left(S_{\pi} A \otimes S_{\mu} B \otimes S_{\nu} C, S^{d}(A \otimes B \otimes C)\right) .
$$

In particular, $k_{\pi, \mu, \nu}$ is independent of the order of $\pi, \mu, \nu$.
Recall from $\S 8.1 .5$ that $S_{\mu} E$ is a trivial $S L(E)$ module if and only if $\mu=\left(\delta^{n}\right)$ for some $\delta \in \mathbb{Z}$. Thus so far, we are reduced to studying the Kronecker coefficients $k_{\pi, \delta^{n}, \delta^{n}}$. Now take the $\mathbb{Z}_{2}$ action given by exchanging $E$ and $F$ into account. Write $[\mu] \otimes[\mu]=S^{2}[\mu] \oplus \Lambda^{2}[\mu]$. The first module will be invariant under $\mathbb{Z}_{2}=\mathfrak{S}_{2}$, and the second will transform its sign under the transposition. So define the symmetric Kronecker coefficients $s k_{\mu, \mu}^{\pi}:=\operatorname{dim}\left(\operatorname{Hom}_{\mathfrak{S}_{d}}\left([\pi], S^{2}[\mu]\right)\right)$.

We conclude:
Proposition 8.8.2.2. [BLMW11] Let $W=\mathbb{C}^{n^{2}}$. The coordinate ring of the $G L(W)$-orbit of $\operatorname{det}_{n}$ is

$$
\mathbb{C}\left[G L(W) \cdot \operatorname{det}_{n}\right]=\bigoplus_{d \in \mathbb{Z}} \bigoplus_{\pi| | \pi \mid=n d}\left(S_{\pi} W^{*}\right)^{\oplus s k_{d}^{\pi} d^{n}} .
$$

While Kronecker coefficients were studied classically (if not the symmetric version), unfortunately very little is known about them. In the next section I describe a geometric method used to study them.

### 8.9. Plethysm coefficients, Kronecker coefficients and geometry

A basic, if not the basic problem in representation theory is: given a group $G$, an irreducible $G$-module $U$, and a subgroup $H \subset G$, decompose $U$ as an $H$-module. The determination of Kronecker coefficients can be phrased this way with $G=G L(V \otimes W), U=S_{\lambda}(V \otimes W)$ and $H=G L(V) \times G L(W)$. The determination of plethysm coefficients may be phrased as the case $G=$ $G L\left(S^{n} V\right), U=S^{d}\left(S^{n} V\right)$ and $H=G L(V)$.

I briefly discuss a geometric method of L. Manivel and J. Wahl [Wah91, Man97, Man98, Man15b, Man15a] based on the Bott-Borel-Weil theorem that allows one to gain asymptotic information about such decomposition problems.

The Bott-Borel-Weil theorem realizes modules as spaces of sections of vector bundles on homogeneous varieties. The method studies sequences of such sections. It has the properties: (i) the vector bundles come with filtrations that allow one to organize information, (ii) the sections of the associated graded bundles can be computed explicitly, giving one bounds for the coefficients, and (iii) Serre's theorem on the vanishing of sheaf cohomology tells one that the bounds are achieved eventually, and gives an upper bound for when stabilization occurs.

I now discuss the decomposition of $S^{d}\left(S^{n} V\right)$.
8.9.1. Asymptotics of plethysm coefficients. We want to decompose $S^{d}\left(S^{n} V\right)$ as a $G L(V)$-module, or more precisely, to obtain qualitative asymptotic information about this decomposition. Note that $S^{d n} V \subset S^{d}\left(S^{n} V\right)$ with multiplicity one. Beyond that the decomposition gets complicated. Let $x_{1}, \ldots, x_{\mathbf{v}}$ be a basis of $V$, so $\left(\left(x_{1}\right)^{n}\right)^{d}$ is the highest highest weight vector in $S^{d}\left(S^{n} V\right)$.

Define the inner degree lifting map $\mathfrak{m}_{x_{1}}=\mathfrak{m}_{x_{1}}^{d, m, n}: S^{d}\left(S^{m} V\right) \rightarrow S^{d}\left(S^{n} V\right)$ on basis elements by

$$
\begin{align*}
& \left(x_{1}^{i_{1}^{1}} x_{2}^{i_{2}^{1}} \cdots x_{d}^{i_{d}^{1}}\right) \cdots\left(x_{1}^{i_{1}^{d}} \cdots x_{d}^{i_{d}^{d}}\right)  \tag{8.9.1}\\
& \mapsto\left(x_{1}^{i_{1}^{1}+(n-m)} x_{2}^{i_{2}^{1}} \cdots x_{d}^{i_{d}^{1}}\right) \cdots\left(x_{1}^{i_{1}^{d}+(n-m)} \cdots x_{d}^{i_{d}^{d}}\right)
\end{align*}
$$

and extend linearly. Here $i_{1}^{j}+\cdots+i_{d}^{j}=m$ for all $j$.
A vector of weight $\mu=\left(q_{1}, q_{2}, \ldots, q_{d}\right)$ is mapped under $\mathfrak{m}_{x_{1}}$ to a vector of weight $\pi=\left(p_{1}, \ldots, p_{d}\right):=\mu+(d(n-m))=\left(q_{1}+d(n-m), q_{2}, \ldots, q_{d}\right)$ in $S^{d}\left(S^{n} V\right)$.

Define the outer degree lifting map $\mathfrak{o}_{x_{1}}=\mathfrak{o}_{x_{1}}^{\delta, d, n}: S^{\delta}\left(S^{n} V\right) \rightarrow S^{d}\left(S^{n} V\right)$ on basis elements by
$\left(x_{i_{1,1}} \cdots x_{i_{1, n}}\right) \cdots\left(x_{i_{\delta, 1}} \cdots x_{i_{\delta, n}}\right) \mapsto\left(x_{i_{1,1}} \cdots x_{i_{1, n}}\right) \cdots\left(x_{i_{\delta, 1}} \cdots x_{i_{\delta, n}}\right)\left(x_{1}^{n}\right) \cdots\left(x_{1}^{n}\right)$
and extend linearly. A vector of weight $\mu=\left(q_{1}, q_{2}, \ldots, q_{d}\right)$ is mapped under $\mathfrak{o}_{x_{1}}$ to a vector of weight $\pi=\left(p_{1}, \ldots, p_{d}\right):=\mu+((d-\delta) n)=\left(q_{1}+(d-\right.$ $\left.\delta) n, q_{2}, \ldots, q_{d}\right)$ in $S^{d}\left(S^{n} V\right)$.

Both $\mathfrak{m}_{x_{1}}$ and $\mathfrak{o}_{x_{1}}$ take highest weight vectors to highest weight vectors, as Lie algebra raising operators annihilate $x_{1}$.

This already shows qualitative behavior if we allow the first part of a partition to grow. More generally, one has:
Theorem 8.9.1.1. [Man97] Let $\mu$ be a fixed partition. Then mult $\left(S_{(d n-|\mu|, \mu)} V, S^{d}\left(S^{n} V\right)\right)$ is a non-decreasing function of both $d$ and $n$ that is constant as soon as $d \geq|\mu|$ or $n \geq l(\mu)$.

More precisely, the inner and outer degree lifting maps $\mathfrak{m}_{x_{1}}$ and $\mathfrak{o}_{x_{1}}$ are both injective and eventually isomorphisms on highest weight vectors of isotypic components of partitions $\left(p_{1}, \ldots, p_{\mathbf{v}}\right)$ with $\left(p_{2}, \ldots, p_{\mathbf{v}}\right)$ fixed and $p_{1}$ growing.

There are several proofs of the stability. The precise stabilization is proved by computing the space of sections of homogeneous vector bundles on $\mathbb{P} V$ via an elementary application of Bott's theorem (see, e.g., [Wey03, §4.1] for an exposition of Bott's theorem).

One way to view what we just did was to write $V=x_{1} \oplus T$, so

$$
\begin{equation*}
S^{n}\left(x_{1} \oplus T\right)=\bigoplus_{j=0}^{n} x_{1}^{n-j} \otimes S^{j} T \tag{8.9.3}
\end{equation*}
$$

Then decompose the $d$-th symmetric power of $S^{n}\left(x_{1} \oplus T\right)$ and examine the stable behavior as we increase $d$ and $n$. One could think of the decomposition (8.9.3) as the osculating sequence of the $n$-th Veronese embedding of $\mathbb{P} V$ at $\left[x_{1}^{n}\right]$ and the further decomposition as the osculating sequence (see, e.g., [IL16b, Chap. 4]) of the $d$-th Veronese re-embedding of the ambient space refined by (8.9.3).

For Kronecker coefficients and more general decomposition problems the situation is more complicated in that the ambient space is no longer projective space, but a homogeneous variety, and instead of an osculating sequence, one examines jets of sections of a vector bundle.

### 8.9.2. A partial converse to Proposition 8.4.2.1.

Proposition 8.9.2.1. [KL14] Let $\pi=\left(p_{1}, \ldots, p_{\mathbf{w}}\right)$ be a partition of $d n$. If $p_{1} \geq \min \{d(n-1), d n-m\}$, then $I_{d}\left(\operatorname{Pad}_{n-m}\left(S^{n} W\right)\right)$ does not contain a copy of $S_{\pi} W^{*}$.

Proof. The image of the space of highest weight vectors for the isotypic component of $S_{\mu} W^{*}$ in $S^{d}\left(S^{m} W^{*}\right)$ under $\mathfrak{m}_{x_{1}}^{d, m, n}$ will be in $\mathbb{C}\left[\operatorname{Pad}_{n-m}\left(S^{n} W\right)\right]$ because, for example, such a polynomial will not vanish on $\left(e_{1}\right)^{n-m}\left[\left(e_{1}\right)^{i_{1}^{1}} \cdots\left(e_{d}\right)^{i_{d}^{1}}+\right.$ $\left.\cdots+\left(e_{1}\right)^{i_{1}^{d}} \cdots\left(e_{d}\right)^{i_{d}^{d}}\right]$, but if $p_{1} \geq d(n-1)$ we are in the stability range.

For the sufficiency of $p_{1} \geq d n-m$, note that if $p_{1} \geq(d-1) n+(n-m)=$ $d n-m$, then in an element of weight $\pi$, each of the exponents $i_{1}^{1}, \ldots, i_{1}^{d}$ of $x_{1}$ must be at least $n-m$. So there again exists an element of $\operatorname{Pad}_{n-m}\left(S^{n} W\right)$ such that a vector of weight $\pi$ does not vanish on it.

### 8.10. Orbit occurrence obstructions cannot separate $\mathcal{P e r m} n_{n}^{m}$ from $\mathcal{D} e t_{n}$

I present an outline of the proof $[\mathbf{I P 1 5}, \mathbf{B I P} 16]$ that the program of $[\mathbf{M S 0 1}$, MS08] cannot work as originally proposed, or even the refinement discussed in [BLMW11]. Despite this negative news, the program has opened several promising directions, and inspired perspectives that have led to concrete advances such as $[\mathbf{L R 1 5}]$ as described in $\S 7.4 .1$.

Throughout this section, set $W=\mathbb{C}^{n^{2}}$.
8.10.1. Occurrence obstructions cannot separate. The program of [MS01, MS08] proposed to use orbit occurrence obstructions to prove Valiant's conjecture. In [IP15] they show that this cannot work. Furthermore, in [BIP16] they prove that one cannot even use the following relaxation of orbit occurrence obstructions:

Definition 8.10.1.1. An irreducible $G L(W)$-module $S_{\lambda} W^{*}$ appearing in $\operatorname{Sym}\left(S^{n} W^{*}\right)$ and not appearing in $\mathbb{C}\left[\hat{\mathcal{D e t}} t_{n}\right]$ is called an occurrence obstruction.

The extension is all the more remarkable because they essentially prove that occurrence obstructions cannot even be used to separate any degree $m$ polynomial padded by $\ell^{n-m}$ in $m^{2}$ variables from

$$
\begin{equation*}
M J\left(v_{n-k}(\mathbb{P} W), \sigma_{r}\left(v_{k}(\mathbb{P} W)\right)\right)=\overline{G L(W) \cdot\left[\ell^{n-k}\left(x_{1}^{k}+\cdots+x_{r}^{k}\right)\right]} \tag{8.10.1}
\end{equation*}
$$

for certain $k, r$ with $k r \leq n$. Here $M J$ is the multiplicative join of $\S 7.5 .3$.
First I show that the variety $(8.10 .1)$ is contained in $\mathcal{D e t}_{n}$. I will use the following classical result:

Theorem 8.10.1.2. [Valiant [Val79], Liu-Regan [LR06]] Every $f \in \mathbb{C}\left[x_{1}, \ldots, x_{n}\right]$ of formula size $u$ is a projection of $\operatorname{det}_{u+1}$. In other words $f \in \operatorname{End}\left(\mathbb{C}^{(u+1)^{2}}\right)$. $\operatorname{det}_{u+1}$.

Note that the formula size of $x_{1}^{k}+\cdots+x_{r}^{k}$ is at most $r k$.
Corollary 8.10.1.3. [BIP16] If $r k<n$ then $\left[\ell^{n-k}\left(x_{1}^{k}+\cdots+x_{r}^{k}\right)\right] \in \mathcal{D e t}_{n}$ and thus $\overline{G L(W) \cdot\left[\ell^{n-k}\left(x_{1}^{k}+\cdots+x_{r}^{k}\right)\right]} \subset \mathcal{D e t}_{n}$.

Their main theorem is:
Theorem 8.10.1.4. [BIP16] Let $n>m^{25}$. Let $\pi=\left(p_{1}, \ldots, p_{\ell}\right)$ be a partition of dn such that $\ell \leq m^{2}+1$ and $p_{1} \geq d(n-m)$. If a copy of $S_{\pi} W^{*}$ occurs in $S^{d}\left(S^{n} W^{*}\right)$ then a copy also occurs in some $\mathbb{C}\left[\overline{G L(W) \cdot\left[\ell^{n-k}\left(x_{1}^{k}+\cdots+x_{r}^{k}\right)\right]}\right]$ for some $r, k$ with $r k<n$.

By the above discussion, this implies occurrence obstructions cannot be used to separate the permanent from the determinant.

The proof is done by splitting the problem into three cases:
(1) $d \leq \sqrt{\frac{n}{m}}$
(2) $d>\sqrt{\frac{n}{m}}$ and $p_{1}>d n-m^{10}$
(3) $d>\sqrt{\frac{n}{m}}$ and $p_{1} \leq d n-m^{10}$.

The first case is an immediate consequence of the prolongation property §8.3.4: take $r=d$ and $k=m$.

The second reduces to the first by two applications of Manivel's stability theorem:
Proposition 8.10.1.5. [BIP16, Prop. 5.2] Let $|\pi|=d n, l(\pi) \leq m^{2}+1$, $p_{2} \leq k, m^{2} k^{2} \leq n$ and $m^{2} k \leq d$. If a copy of $S_{\pi} W$ occurs in $S^{d}\left(S^{n} W\right)$, then a copy also occurs in $\mathbb{C}\left[\overline{G L(W) \cdot\left[\ell^{n-k}\left(x_{1}^{k}+\cdots+x_{m^{2} k}^{k}\right)\right]}\right]$.

Proof. For a partition $\mu=\left(m_{1}, \ldots, m_{l}\right)$, introduce the notation $\bar{\mu}=\left(m_{2}, \ldots, m_{l}\right)$ First note that the inner degree lifting map (8.9.1) $\mathfrak{m}_{\ell}^{d, k, n}: S^{d}\left(S^{k} W^{*}\right) \rightarrow$ $S^{d}\left(S^{n} W^{*}\right)$ is an isomorphism on highest weight vectors in this range because $d$ is sufficiently large, so there exists $\mu$ with $|\mu|=d k$ and $\bar{\pi}=\bar{\mu}$. Moreover, if $v_{\mu}$ is a highest weight vector of weight $\mu$, then $\mathfrak{m}_{\ell}^{d, k, n}\left(v_{\mu}\right)$ is a highest weight vector of weight $\pi$. Since $m^{2} k$ is sufficiently large, there exists $\nu$ with $|\nu|=$ $m^{2} k^{2}=\left(m^{2} k\right) k$, with $\bar{\nu}=\bar{\mu}$ such that $v_{\mu}=\mathfrak{o}_{x_{1}}\left(w_{\nu}\right)$, where $w_{\nu}$ is a highest weight vector of weight $\nu$ in $S^{m^{2} k}\left(S^{k} W^{*}\right)$. Since $I_{m^{2} k}\left(\sigma_{m^{2} k}\left(v_{k}(\mathbb{P} W)\right)\right)=0$, we conclude that a copy of $S_{\nu} W^{*}$ is in $\mathbb{C}\left[\sigma_{m^{2} k}\left(v_{k}(\mathbb{P} W)\right)\right]$ and then by the discussion above the modules corresponding to $\mu$ and $\pi$ are respectively in the coordinate rings of $M J\left(\left[\ell^{d-m^{2} k}\right], \sigma_{m^{2} k}\left(v_{k}(\mathbb{P} W)\right)\right)$ and $M J\left(\left[\ell^{n-k}\right], \sigma_{m^{2} k}\left(v_{k}(\mathbb{P} W)\right)\right)$. Since $\left(m^{2} k\right) k \leq n$, the result follows by prolongation.

The third case relies on a building block construction made possible by the following exercise:
Exercise 8.10.1.6: (1!) Show that if $V$ is a $G L(W)$-module and $Q \in$ $S_{\lambda} W \subset S^{d} V$ and $R \in S_{\mu} W \subset S^{\delta} V$ are both highest weight vectors, then $Q R \in S_{\lambda+\mu} W \subset S^{d+\delta} V$ is also a highest weight vector.

Exercise 8.10.1.6, combined with the fact that for an irreducible variety $X$, if $Q, R \in \mathbb{C}[X]$, then $Q R \in \mathbb{C}[X]$ enables the building block construction assuming $n>m^{25}$. I will show (Corollary 9.4.1.2) that for $n$ even, there exists a copy of $S_{\left(n^{d}\right)} W$ in $\mathbb{C}\left[\sigma_{d}\left(v_{n}(\mathbb{P} W)\right)\right]$, providing one of the building blocks. The difficulty in their proof lies in establishing the other base building block cases. See [BIP16] for the details.

Remark 8.10.1.7. In [IP15] the outline of the proof is similar, except there is an interesting argument by contradiction: they show that in a certain range of $n$ and $m$, if an orbit occurrence obstruction exists, then the same is true for larger values of $n$ with the same $m$. But this contradicts Valiant's result (see $\S 6.6 .3$ ) that if $n=4^{m}$, then $\ell^{n-m}$ perm $_{m} \in \mathcal{D e t}_{n}$.

It is conceivably possible to carry out a modification of the program, either taking into account information about multiplicities, or with the degree $m$ iterated matrix multiplication polynomial $I M M_{n}^{m}$ in place of the determinant, as the latter can be compared to the permanent without padding.

### 8.11. Equivariant determinantal complexity

The GCT perspective of focusing on symmetry groups led to the discovery of symmetry in Grenet's expression for the permanent, as well as the restricted model of equivariant determinantal complexity. In this section I first give a geometric description of Grenet's expressions in the IMM model, and then outline the proof that the equivariant determinantal complexity of perm $_{m}$ is $\binom{2 m}{m}-1$.
8.11.1. Geometric description of Grenet's expression. I now describe Grenet's size $2^{m}-1$ determinantal expression for perm $_{m}$ from a geometric perspective. The matrix $A_{\text {Grenet }}(y)$ in Grenet's expression is in block format, and taking $\operatorname{det}\left(\tilde{A}_{\text {Grenet }}(y)\right)$ amounts to the matrix multiplication of these blocks (see, e.g., the expression (1.2.3) compared with (7.3.1)), and so are more naturally described as a homogeneous iterated matrix multiplication. Recall that for $P \in S^{m} \mathbb{C}^{N}$, this is a sequence of matrices $M_{1}(y), \ldots, M_{m}(y)$, with $M_{j}$ of size $m_{j-1} \times m_{j}$ and $m_{0}=m_{m}=1$, such that $P(y)=M_{1}(y) \cdots M_{m}(y)$. View this more invariantly as

$$
U_{m}=\mathbb{C} \xrightarrow{M_{m}(y)} U_{m-1} \xrightarrow{M_{m-1}(y)} \cdots \rightarrow U_{1}=\mathbb{C},
$$

where $M_{j}$ is a linear map $\mathbb{C}^{N} \rightarrow U_{j}^{*} \otimes U_{j-1}$. Such a presentation is $G$ equivariant, for some $G \subseteq G_{P}$, if there exist representations $\rho_{j}: G \rightarrow$ $G L\left(U_{j}\right)$, with dual representations $\rho_{j}^{*}: G \rightarrow G L\left(U_{j}^{*}\right)$, such that for all $g \in G,\left(\rho_{j}^{*} \otimes \rho_{j+1}\right)(g) M_{j}(g \cdot y)=M(y)$.

Write $\operatorname{perm}_{m} \in S^{m}(E \otimes F)$. In the case of Grenet's presentation, we need each $U_{j}$ to be a $\Gamma^{E}=\left(T^{S L(E)} \rtimes \mathfrak{S}_{m}\right)$-module and $M_{1}(y) \cdots M_{m}(y)$ to equal the permanent.

Let $\left(S^{k} E\right)_{\text {reg }}$ denote the span of the square free monomials, which I will also call the regular weight space. It is the span of all vectors of weight $\left(a_{1}, \ldots, a_{m}\right)$ with $a_{j} \in\{0,1\}$ and $\sum_{a_{j}}=k$. This is an irreducible $\Gamma^{E}$-module. Note that $\left(S^{m} E\right)_{\text {reg }}$ is a one-dimensional vector space, and $\operatorname{perm}_{m} \in\left(S^{m} E\right)_{\text {reg }} \otimes\left(S^{m} F\right)_{\text {reg }} \subset S^{m} E \otimes S^{m} F \subset S^{m}(E \otimes F)$, which characterizes $\operatorname{perm}_{m}$ up to scale (and the scale can be fixed e.g., by evaluating on the identity matrix).

Note that $E \subset \operatorname{Hom}\left(\left(S^{j} E\right)_{\text {reg }},\left(S^{j+1} E\right)_{\text {reg }}\right)$ via the composition

$$
\begin{equation*}
E \otimes\left(S^{j} E\right)_{\text {reg }} \rightarrow S^{j+1} E \rightarrow\left(S^{j+1} E\right)_{r e g} \tag{8.11.1}
\end{equation*}
$$

where the first map is multiplication and the second projection onto the regular weight space. This inclusion is as a $\Gamma^{E}$-module. Fix a basis $f_{1}, \ldots, f_{m}$ of $F$. Consider the spaces $U_{j}:=\left(S^{j} E\right)_{\text {reg }} \otimes S^{j} F$, and the inclusions $E \otimes f_{j} \subset$ $\operatorname{Hom}\left(\left(S^{j} E\right)_{\text {reg }} \otimes S^{j} F,\left(S^{j+1} E\right)_{\text {reg }} \otimes S^{j+1} F\right)$ where the $E$ side is mapped via (8.11.1) and the $F$ side is multiplied by the vector $f_{j}$.

Taking the chain of maps from $U_{0}$ to $U_{m}$, by construction our output polynomial lies in $\left(S^{m} E\right)_{\text {reg }} \otimes S^{m} F$, but the weight on the second term is $(1, \ldots, 1)$ so it must lie in the one-dimensional space $\left(S^{m} E\right)_{\text {reg }} \otimes\left(S^{m} F\right)_{\text {reg }}$. Finally we check that it is indeed the permanent by evaluating on the identity matrix.

Remark 8.11.1.1. The above construction is a symmetric cousin of a familiar construction in algebra, namely the Koszul maps:

$$
\Lambda^{0} E \xrightarrow{\wedge y_{1}} \Lambda^{1} E \xrightarrow{\wedge y_{2}} \Lambda^{2} E \xrightarrow{\wedge y_{3}} \cdots \xrightarrow{\wedge y_{m}} \Lambda^{m} E .
$$

If we tensor this with exterior multiplication by basis vectors of $F$, we obtain a $S L(E)$-equivariant homogeneous iterated matrix multiplication of $\operatorname{det}_{m} \in$ $\Lambda^{m} E \otimes \Lambda^{m} F$ of size $2^{m}-1$.
(Note that both the Koszul maps and (8.11.1) give rise to complexes, i.e., if we multiply by the same vector twice we get zero.)

This IMM realization of the determinant is related to the IMM version of Grenet's realization of the permanent via the Howe-Young duality functor: The involution on the space of symmetric functions (see [Mac95, §I.2]) that exchanges elementary symmetric functions with complete symmetric functions, (and, for those familiar with the notation, takes the Schur function
$s_{\pi}$ to $s_{\pi^{\prime}}$ ) extends to modules of the general linear group. This functor exchanges symmetrization and skew-symmetrization. For more explanations, see $\S 10.4 .4$, where it plays a central role. I expect it will be useful for future work regarding permanent v . determinant. It allows one to transfer knowledge about the well-studied determinant, to the less understood permanent.

One can have a full $G_{\text {perm }_{m}}$-equivariant expression by considering the inclusions

$$
\begin{equation*}
E \otimes F \subset \operatorname{Hom}\left(\left(S^{j} E\right)_{r e g} \otimes\left(S^{j} F\right)_{r e g},\left(S^{j+1} E\right)_{r e g} \otimes\left(S^{j+1} F\right)_{r e g}\right) \tag{8.11.2}
\end{equation*}
$$

(The transpose invariance is possible because transposing the matrices in the sequence $M_{m-j}$ is sent to $M_{j}^{T}$ and $\left(M_{m-j}\left(y^{T}\right)\right)^{T}=M_{j}(y)$.)
Exercise 8.11.1.2: (1) Show that (8.11.2) gives rise to a size $\binom{2 m}{m}$ IMM expression for perm $_{m}$. ©

Remark 8.11.1.3. One similarly obtains a size $\binom{2 m}{m} G_{\operatorname{det}_{m}}$-equivariant IMM presentation of $\operatorname{det}_{m}$.
8.11.2. Outline of proofs of lower bounds. Recall the lower bound theorems:
Theorem 8.11.2.1. [LR15]Assume $m \geq 3$.

- $\operatorname{edc}\left(\operatorname{perm}_{m}\right)=\binom{2 m}{m}-1$ with equality given by the determinantal expression obtained from (8.11.2).
- The smallest size $\Gamma^{E}$-equivariant determinantal expression for $\operatorname{perm}_{m}$ is $2^{m}-1$ and is given by $\tilde{A}_{\text {Grenet }}$.
Ideas towards the proofs are as follows: Write $\mathbb{C}^{n^{2}}=B \otimes C$. Without loss of generality, one takes the constant part $\Lambda$ of $\tilde{A}$ to be the diagonal matrix with zero in the (1,1)-slot and 1's on the remaining diagonal entries. Then $\Lambda$ determines a splitting $B \otimes C=\left(B_{1} \oplus B_{2}\right) \otimes\left(C_{1} \oplus C_{2}\right)$ with $\operatorname{dim} B_{1}=$ $\operatorname{dim} C_{1}=1$. Consider the linear part of an expression $A: E \otimes F \rightarrow B \otimes C$. We have already seen (in the proof of Theorem 6.3.4.6) the component in $B_{1} \otimes C_{1}$ (i.e., the $(1,1)$ slot in the matrix $\left.A(y)\right)$ must be zero. Thus in order for the expression not to be identically zero, we must have the components of $A(y)$ in $B_{1} \otimes C_{2}$ and $B_{2} \otimes C_{1}$ nonzero (i.e., other entries in the first row and column must be nonzero). Focus on the $\Gamma^{E}$-equivariant case for simplicity of exposition. As a $\Gamma^{E}$-module, $E \otimes F=E^{\oplus m}$. By $\Gamma^{E}$-equivariance, $B_{1} \otimes C_{2}$ must contain at least one copy of $E$, write the submodule as $B_{1} \otimes C_{2,1} \simeq E^{\oplus j}$. For simplicity of discussion, assume $j=1$. Also for simplicity, assume the one-dimensional space $B_{1}$ is a trivial $\Gamma^{E}$-module, so $C_{2,1} \simeq E$ as a $\Gamma^{E_{-}}$ module. Since $\Gamma^{E}$ is reductive, we have a splitting $C_{2}=C_{2,1} \oplus C_{2}^{\prime}$. In order that there is no degree two component appearing, we must have the map to
$C_{2,1}^{*} \otimes B_{1}^{*}$ be zero. The picture of what we have reduced to so far looks like this:

$$
\left(\begin{array}{ccc}
0 & C_{2,1}^{*} \otimes B_{1}^{*} & * \\
B_{1} \otimes C_{2,1} & \operatorname{Id} & * \\
* & ? & *
\end{array}\right) .
$$

Now in order that the determinant is not identically zero, the ? block cannot be identically zero, so there must be some $B_{2,1} \subset B_{2}$, such that $C_{2,1}^{*} \otimes B_{2} \simeq E^{*} \otimes B_{2,1}$ contains a copy of $E$.

Fact: the minimum dimension of a $\Gamma^{E}$-module $M$ such that $E \subset E^{*} \otimes M$ is $\binom{m}{2}$ and the corresponding module is (up to tensoring with a one-dimensional representation) $\left(S^{2} E\right)_{\text {reg }}$.

Remark 8.11.2.2. Were we constructing a $S L(E)$-equivariant regular determinantal presentation of the determinant, we would need an $S L(E)$ module $M$ such that $E \subset E^{*} \otimes M$. By the Pieri rule, the admissible $M$ correspond to Young diagrams with two boxes, i.e., $S^{2} E$ and $\Lambda^{2} E$. Note that $\operatorname{dim}\left(\Lambda^{2} E\right)=\binom{m}{2}$. This "coincidence" of dimensions is attributable to the Howe-Young duality endofunctor.

Continuing, we need some $B_{2,2}$ such that $E \subset\left(S^{2} E\right)_{\text {reg }}^{*} \otimes B_{2,2}$, and the smallest such is $B_{2,2}=\left(S^{3} E\right)_{\text {reg }}$ (just as in the skew case, one needs a Young diagram with three boxes, the smallest such module is $\Lambda^{3} E$ ).

One continues until arriving at $B=\bigoplus_{j=0}^{m-1}\left(S^{j} E\right)_{\text {reg }}$ and one concludes.
Remark 8.11.2.3. In the above discussion I swept two important complications under the rug. First, we don't really have $\Gamma^{E} \subset G_{\operatorname{det}_{n}, \Lambda}$, but rather a group $G \subset G_{\operatorname{det}_{n}, \Lambda}$ that has a surjective map onto $\Gamma^{E}$. This problem is dealt with by observing that the modules for any such $G$ can be labeled using the labels from $\Gamma^{E}$-modules. Second, since $\Gamma^{E}$ is not connected, we need to allow the possibility that the $\mathbb{Z}_{2} \subset G_{\operatorname{det}_{n}, \Lambda}$ is part of the equivariance. This second problem is dealt with by restricting to the alternating group. For details, see [LR15].

### 8.12. Symmetries of additional polynomials relevant for complexity theory

A central insight from GCT is that polynomials that are determined by their symmetry groups should be considered preferred representatives of their complexity classes. This idea has already guided several results: i) the symmetries of the matrix multiplication tensor have given deep insight into its decompositions, ii) these symmetries were critical for proving its border rank lower bounds, and iii) the above results on equivariant determinantal complexity. We have already determined the symmetry groups of the
determinant, permanent, and $x_{1} \cdots x_{n}$. In this section I present the symmetry groups of additional polynomials relevant for complexity theory and techniques for determining them.

Throughout this section $G=G L(V), \operatorname{dim} V=n$, and I use index ranges $1 \leq i, j, k \leq n$.
8.12.1. The Fermat. This example follows [CKW10]. Let fermat ${ }_{n}^{d}:=$ $x_{1}^{d}+\cdots+x_{n}^{d} \in S^{d} \mathbb{C}^{n}$. The $G L_{n}$-orbit closure of [fermat ${ }_{n}^{d}$ ] is the $n$-th secant variety of the Veronese variety $\sigma_{n}\left(v_{d}\left(\mathbb{P}^{n-1}\right)\right) \subset \mathbb{P} S^{d} \mathbb{C}^{n}$. It is clear $\mathfrak{S}_{n} \subset$ $G_{\text {fermat }}$, as well as the diagonal matrices whose entries are $d$-th roots of unity. We need to see if there is anything else. The first idea, to look at the singular locus, does not work, as the zero set is smooth, so consider (fermat $\left.{ }_{n}^{d}\right)_{2, d-2}=$ $x_{1}^{2} \otimes x^{d-2}+\cdots+x_{n}^{2} \otimes x^{d-2}$. Write the further polarization (fermat $\left.{ }_{n}^{d}\right)_{1,1, d-2}$ as a symmetric matrix whose entries are homogeneous polynomials of degree $d-2$ (the Hessian matrix):

$$
\left(\begin{array}{lll}
x_{1}^{d-2} & & \\
& \ddots & \\
& & x_{n}^{d-2}
\end{array}\right)
$$

Were the determinant of this matrix $G L(V)$-invariant, we could proceed as we did with $e_{n, n}$, using unique factorization. Although it is not, it is close enough as follows:

Recall that for a linear map $f: W \rightarrow V$, where $\operatorname{dim} W=\operatorname{dim} V=n$, we have $f^{\wedge n} \in \Lambda^{n} W^{*} \otimes \Lambda^{n} V$ and an element $(h, g) \in G L(W) \times G L(V)$ acts on $f^{\wedge n}$ by $(h, g) \cdot f^{\wedge n}=(\operatorname{det}(h))^{-1}(\operatorname{det}(g)) f^{\wedge n}$. In our case $W=V^{*}$ so $P_{2, d-2}^{\wedge n}(x)=\operatorname{det}(g)^{2} P_{2, d-2}^{\wedge n}(g \cdot x)$, and the polynomial obtained by the determinant of the Hessian matrix is invariant up to scale.

Arguing as in §7.1.2, $\sum_{j}\left(g_{1}^{j_{1}} x_{j_{1}}\right)^{d-2} \cdots\left(g_{n}^{j_{n}} x_{j_{n}}\right)^{d-2}=x_{1}^{d-2} \cdots x_{n}^{d-2}$ and we conclude again by unique factorization that $g$ is in $\mathfrak{S}_{n} \ltimes T_{n}$. Composing with a permutation matrix to make $g \in T$, we see that, by acting on the Fermat itself, that the entries on the diagonal are $d$-th roots of unity.

In summary:
Proposition 8.12.1.1. $G_{x_{1}^{d}+\cdots+x_{n}^{d}}=\mathfrak{S}_{n} \ltimes\left(\mathbb{Z}_{d}\right)^{\times n}$.
Exercise 8.12.1.2: (2) Show that the Fermat is characterized by its symmetries.
8.12.2. The sum-product polynomial. The polynomial

$$
S P_{r}^{n}:=\sum_{i=1}^{r} \Pi_{j=1}^{n} x_{i j} \in S^{n}\left(\mathbb{C}^{n r}\right)
$$

called the sum-product polynomial in the CS literature, was used in our study of depth three circuits. Its $G L(r n)$-orbit closure is the $r$-th secant variety of the Chow variety $\sigma_{r}\left(C h_{n}\left(\mathbb{C}^{n r}\right)\right)$.
Exercise 8.12.2.1: (2) Determine $G_{S P_{r}^{n}}$ and show that $S P_{r}^{n}$ is characterized by its symmetries.
8.12.3. Further Techniques. One technique for determining $G_{P}$ is to form auxiliary objects from $P$ which have a symmetry group $H$ that one can compute, and by construction $H$ contains $G_{P}$. Usually it is easy to find a group $H^{\prime}$ that clearly is contained in $G_{P}$, so if $H=H^{\prime}$, we are done.

Recall that we have already used auxiliary varieties such as $\operatorname{Zeros}(P)_{J a c, k}$ and $\operatorname{Zeros}(P)^{\vee}$ in determining the symmetry groups of $\operatorname{perm}_{n}$ and $\operatorname{det}_{n}$.

One can determine the connected component of the stabilizer by a Lie algebra calculation: If we are concerned with $p \in S^{d} V$, the connected component of the identity of the stabilizer of $p$ in $G L(V)$ is the connected Lie group associated to the Lie subalgebra of $\mathfrak{g l}(V)$ that annihilates $p$. (The analogous statement holds for tensors.) To see this, let $\mathfrak{h} \subset \mathfrak{g l}(V)$ denote the annihilator of $p$ and let $H=\exp (\mathfrak{h}) \subset G L(V)$ the corresponding Lie group. Then it is clear that $H$ is contained in the stabilizer as $h \cdot p=\exp (X) \cdot p=\left(\operatorname{Id}+X+\frac{1}{2} X X+\ldots\right) p$ the first term preserves $p$ and the remaining terms annihilate it. Similarly, if $H$ is the group preserving $p$, taking the derivative of any curve in $H$ through Id at $t=0$ gives $\left.\frac{d}{d t}\right|_{t=0} h(t) \cdot p=0$.

To recover the full stabilizer from knowledge of the connected component of the identity, we have the following observation, the first part comes from [BGL14]:
Proposition 8.12.3.1. Let $V$ be an irreducible $G L(W)$-module. Let $G_{v}^{0}$ be the identity component of the stabilizer $G_{v}$ of some $v \in V$ in $G L(W)$. Then $G_{v}$ is contained in the normalizer $N\left(G_{v}^{0}\right)$ of $G_{v}^{0}$ in $G L(W)$. If $G_{v}^{0}$ is semi-simple and $[v]$ is determined by $G_{v}^{0}$, then up to scalar multiples of the identity in $G L(W), G_{v}$ and $N\left(G_{v}^{0}\right)$ coincide.

Proof. First note that for any group $H$, the full group $H$ normalizes $H^{0}$. (If $h \in H^{0}$, take a curve $h_{t}$ with $h_{0}=\mathrm{Id}$ and $h_{1}=h$, then take any $g \in H$, the curve $g h_{t} g^{-1}$ connects $g h_{1} g^{-1}$ to the identity.) So $G_{v}$ is contained in the normalizer of $G_{v}^{0}$ in $G L(W)$.

For the second assertion, let $h \in N\left(G_{v}^{0}\right)$ be in the normalizer. We have $h^{-1} g h v=g^{\prime} v=v$ for some $g^{\prime} \in G_{v}^{0}$, and thus $g(h v)=(h v)$. But since $[v]$ is the unique line preserved by $G_{v}^{0}$ we conclude $h v=\lambda v$ for some $\lambda \in \mathbb{C}^{*}$.

Here is a lemma for those familiar with roots and weights:

Lemma 8.12.3.2. [BGL14, Prop. 2.2] Let $G^{0}$ be semi-simple and act irreducibly on $V$. Then its normalizer $N\left(G^{0}\right)$ is generated by $G^{0}$, the scalar matrices, and a finite group constructed as follows: Assume we have chosen a Borel for $G^{0}$, and thus have distinguished a set of simple roots $\Delta$ and a group homomorphism $\operatorname{Aut}(\Delta) \rightarrow G L(V)$. Assume $V=V_{\lambda}$ is the irreducible representation with highest weight $\lambda$ of $G^{0}$ and consider the subgroup $\operatorname{Aut}(\Delta, \lambda) \subset \operatorname{Aut}(\Delta)$ that fixes $\lambda$. Then $N\left(G^{0}\right)=\left(\left(\mathbb{C}^{*} \times G^{0}\right) / Z\right) \rtimes \operatorname{Aut}(\Delta, \lambda)$.

For the proof, see $[\mathbf{B G L 1 4}]$.
8.12.4. Iterated matrix multiplication. Let $I M M_{n}^{k} \in S^{n}\left(\mathbb{C}^{k^{2} n}\right)$ denote the iterated matrix multiplication operator for $k \times k$ matrices, $\left(X_{1}, \ldots, X_{n}\right) \mapsto$ $\operatorname{trace}\left(X_{1} \cdots X_{n}\right)$. Letting $V_{j}=\mathbb{C}^{k}$, invariantly

$$
\begin{gathered}
I M M_{n}^{k}=\operatorname{Id}_{V_{1}} \otimes \cdots \otimes \operatorname{Id}_{V_{n}} \in\left(V_{1} \otimes V_{2}^{*}\right) \otimes\left(V_{2} \otimes V_{3}^{*}\right) \otimes \cdots \otimes\left(V_{n-1} \otimes V_{n}^{*}\right) \otimes\left(V_{n} \otimes V_{1}^{*}\right) \\
\subset S^{n}\left(\left(V_{1} \otimes V_{2}^{*}\right) \oplus\left(V_{2} \otimes V_{3}^{*}\right) \oplus \cdots \oplus\left(V_{n-1} \otimes V_{n}^{*}\right) \oplus\left(V_{n} \otimes V_{1}^{*}\right)\right),
\end{gathered}
$$

and the connected component of the identity of $G_{I M M_{n}^{k}} \subset G L\left(\mathbb{C}^{k^{2} n}\right)$ is $G L\left(V_{1}\right) \times \cdots \times G L\left(V_{n}\right)$.

The case of $I M M_{n}^{3}$ is important as this sequence is complete for the complexity class $\mathbf{V P}_{e}$, of sequences of polynomials admitting polynomial size formulas, see $[\mathbf{B O C} 92]$. Moreover $I M M_{n}^{n}$ is complete for the same complexity class as the determinant, namely $\mathbf{V Q P}=\mathbf{V P}_{s}$, see $[\mathbf{B l a ̈ 0 1 b}]$.

The first equality in the following theorem for the case $k=3$ appeared in [dG78, Thms. 3.3,3.4] and [Bur15, Prop. 4.7] with ad-hoc proofs.
Theorem 8.12.4.1. [Ges16] $G_{I M M_{n}^{k}}=\left(G L_{k}^{\times n} / \mathbb{C}^{*}\right) \rtimes D_{n}$, where $D_{n}=$ $\mathbb{Z}_{n} \rtimes \mathbb{Z}_{2}$ is the dihedral group. The $\mathbb{Z}_{n}^{n}$ corresponds to cyclic permutation of factors, and the $\mathbb{Z}_{2}$ is generated by $\left(X_{1}, \ldots, X_{k}\right) \mapsto\left(X_{k}^{T}, \ldots, X_{1}^{T}\right)$.

A "hands on" elementary proof is possible, see, e.g. [Bur15, Prop. 4.7]. Here is an elegant proof for those familiar with Dynkin diagrams from [Ges16] in the special case of $M_{\langle\mathbf{n}\rangle}$, i.e., $k=\mathbf{n}$ and $n=3$.

Proof. It will be sufficient to show the second equality because the $\left(\mathbb{C}^{*}\right)^{\times 2}$ acts trivially on $A \otimes B \otimes C$. For polynomials, the method of [BGL14, Prop. 2.2] adapts to reducible representations. A straight-forward Lie algebra calculation shows the connected component of the identity of $\tilde{G}_{M_{\langle\mathbf{n}\rangle}}$ is $\tilde{G}_{M_{\langle\mathbf{n}\rangle}}^{0}=$ $\left(\mathbb{C}^{*}\right)^{\times 2} \times P G L_{n}^{\times 3}$. As was observed in $[\mathbf{B G L 1 4}]$, the full stabilizer group must be contained in its normalizer $N\left(\tilde{G}_{M_{\langle\mathbf{n}\rangle}}^{0}\right)$, see Proposition 8.12.3.1. But the normalizer of $\tilde{G}_{M_{\langle\mathbf{n}\rangle}}^{0}$ quotiented by $\tilde{G}_{M_{\langle\mathbf{n}\rangle}}^{0}$ is the automorphism group of the marked Dynkin diagram for $A \oplus B \oplus C$, which is


There are three triples of marked diagrams. Call each column consisting of 3 marked diagrams a group. The automorphism group of the picture is $\mathbb{Z}_{3} \rtimes \mathbb{Z}_{2}$, where the $\mathbb{Z}_{2}$ may be seen as flipping each diagram, exchanging the first and third diagram in each group, and exchanging the first and second group. The $\mathbb{Z}_{3}$ may be seen as cyclically permuting each group and the diagrams within each group.
Problem 8.12.4.2. Find equations in the ideal of $\overline{G L_{9 n} \cdot I M M_{n}^{3}}$. Determine lower bounds for the inclusions $\mathcal{P e r m}_{m} \subset \overline{G L_{9 n} \cdot I M M_{n}^{3}}$ and study common geometric properties (and differences) of $\mathcal{D e t}_{n}$ and $\overline{G L_{9 n} \cdot I M M_{n}^{3}}$.
8.12.5. The Pascal determinant. Let $k$ be even, and let $A_{j}=\mathbb{C}^{n}$. Define the $k$-factor Pascal determinant $P D_{k, n}$ to be the unique up to scale element of $\Lambda^{n} A_{1} \otimes \cdots \otimes \Lambda^{n} A_{k} \subset S^{n}\left(A_{1} \otimes \cdots \otimes A_{k}\right)$. Choose the scale such that if $X=\sum x_{i_{1}, \ldots, i_{k}} a_{1, i_{1}} \otimes \cdots \otimes a_{k, i_{k}}$ with $a_{\alpha, j}$ a basis of $A_{\alpha}$, then (8.12.1)

$$
P D_{k, n}(X)=\sum_{\sigma_{2}, \ldots, \sigma_{k} \in \mathfrak{S}_{n}} \operatorname{sgn}\left(\sigma_{2} \cdots \sigma_{k}\right) x_{1, \sigma_{2}(1), \ldots, \sigma_{k}(1)} \cdots x_{n, \sigma_{2}(n), \ldots, \sigma_{k}(n)}
$$

This expression, for fixed $k$, shows that $\left(P D_{k, n}\right) \in$ VNP.
Proposition 8.12.5.1 (Gurvits). The sequence $\left(P D_{4, n}\right)$ is VNP complete.
Proof. It remains to show VNP-hardness. Set $x_{i j k l}=0$ unless $i=j$ and $k=l$. Then $x_{i, \sigma_{2}(i), \sigma_{3}(i), \sigma_{4}(i)}=0$ unless $\sigma_{2}(i)=i$ and $\sigma_{3}(i)=\sigma_{4}(i)$ so the only nonzero monomials are those where $\sigma_{2}=\mathrm{Id}$ and $\sigma_{3}=\sigma_{4}$. Since the $\operatorname{sign}$ of $\sigma_{3}$ is squared, the result is the permanent.

Thus we could just as well work with the sequence $P D_{4, n}$ as the permanent. Since $\operatorname{det}_{n}=P D_{2, n}$, and the symmetry groups superficially resemble each other, it is an initially appealing substitute.

It is clear the identity component of the stabilizer includes $\left(S L\left(A_{1}\right) \times\right.$ $\left.\cdots \times S L\left(A_{k}\right)\right) / \mu_{n, k}$ where $\mu_{n}$ is as in §6.6.1, and a straight-forward Lie algebra calculation confirms this is the entire identity component. (Alternatively, one can use Dynkin's classification [Dyn52] of maximal subalgebras.) It is also clear that $\mathfrak{S}_{k}$ preserves $P D_{n, k}$ by permuting the factors.
Theorem 8.12.5.2 (Garibaldi, personal communication). For all $k$ even

$$
G_{P D_{k, n}}=S L_{n}^{\times k} / \mu_{n, k} \rtimes \mathfrak{S}_{k}
$$

8. Representation theory and its uses in complexity theory

Note that this includes the case of the determinant, and gives a new proof.

The result will follow from the following Lemma and Proposition 8.12.3.1.
Lemma 8.12.5.3. [Garibaldi, personal communication] Let $V=A_{1} \otimes \cdots \otimes A_{k}$. The normalizer of $S L_{n}^{\times k} / \mu_{n}$ in $G L(V)$ is $\left(G L_{n}^{\times k} / \mu_{k}\right) \rtimes \mathfrak{S}_{k}$, where $\mu_{k}$ denotes the kernel of the product map $\left(\mathbb{C}^{*}\right)^{\times k} \rightarrow \mathbb{C}^{*}$.

Proof of Lemma 8.12.5.3. We use Lemma 8.12.3.2. In our case, the Dynkin diagram for $(\Delta, \lambda)$ is and $\operatorname{Aut}(\Delta, \lambda)$ is clearly $\mathfrak{S}_{k}$.


Figure 8.12.1. Marked Dynkin diagram for $V$

The theorem follows.

## The Chow variety of products of linear forms

In the GCT approach to Valiant's conjecture, one wants to understand the $G L_{n^{2}}$-module structure of $\mathbb{C}\left[\overline{G L_{n^{2}} \cdot\left[\operatorname{det}_{n}\right]}\right]$ via $\mathbb{C}\left[G L_{n^{2}} \cdot\left[\operatorname{det}_{n}\right]\right]$. In this chapter I discuss a "toy" problem that turns out to be deep, subtle and have surprising connections with several different areas of mathematics: the orbit closure $\overline{G L_{n} \cdot\left[x_{1} \cdots x_{n}\right]}=C h_{n}\left(\mathbb{C}^{n}\right) \subset \mathbb{P} S^{n} \mathbb{C}^{n}$. This subject has a remarkable history beginning over 100 years ago, with work of Brill, Gordan, Hermite and Hadamard. The history is rife with rediscoveries and errors that only make the subject more intriguing.

I begin, in $\S 9.1$ describing the Hermite-Hadamard-Howe map $h_{n}$ that has been discovered and rediscovered numerous times. Its kernel is the ideal of the Chow variety. I also state the main results regarding this map: the Black-List propagation theorem and Brion's asymptotic surjectivity theorem. In $\S 9.2$ I re-derive the map from a GCT perspective that compares the coordinate ring of the orbit to that of its closure. In $\S 9.3$ I define a map of modules for the permutation group $\mathfrak{S}_{d n}$ that contains equivalent information to the original map. This map was originally defined in a different manner by Black and List as a path to prove a celebrated conjecture of Foulkes that I also explain in the section. Via a variant of this $\mathfrak{S}_{d n}$-map, I give the proof of the Black-List propagation theorem from [Ike15], which is a geometric reinterpretation of the proof in [McK08]. In $\S 9.4$ I illustrate the subtlety of determining the rank of $h_{n}$ by explaining how a very special case of the
problem is equivalent to a famous conjecture in combinatorics due to Alon and Tarsi. In $\S 9.5$, I give yet another derivation of the map $h_{n}$ via algebraic geometry due to Brion. If one is content with set-theoretic equations for the Chow variety, such equations were discovered over a hundred years ago by Brill and Gordan. I give a modern presentation of these equations in §9.6. I conclude in $\S 9.7$ with the proof of Brion's asymptotic surjectivity theorem. This last proof requires more advanced results in algebraic geometry and commutative algebra, and should be skipped by readers unfamiliar with the relevant notions.

### 9.1. The Hermite-Hadamard-Howe map

I begin with the first description of the ideal of $C h_{n}\left(V^{*}\right)$, due to Hadamard (1897).
9.1.1. The Hermite-Hadamard-Howe map and the ideal of the Chow variety. The following linear map was first defined when $\operatorname{dim} V=2$ by Hermite (1854), and in general independently by Hadamard (1897), and Howe (1988).

Definition 9.1.1.1. The Hermite-Hadamard-Howe map $h_{d, n}: S^{d}\left(S^{n} V\right) \rightarrow$ $S^{n}\left(S^{d} V\right)$ is defined as follows: First include $S^{d}\left(S^{n} V\right) \subset V^{\otimes n d}$. Next, reorder the copies of $V$ from $d$ blocks of $n$ to $n$ blocks of $d$ and symmetrize the blocks of $d$ to obtain an element of $\left(S^{d} V\right)^{\otimes n}$. Finally, thinking of $S^{d} V$ as a single vector space, symmetrize the $n$ blocks.

For example, putting subscripts on $V$ to indicate position:

$$
\begin{aligned}
S^{2}\left(S^{3} V\right) \subset V^{\otimes 6} & =V_{1} \otimes V_{2} \otimes V_{3} \otimes V_{4} \otimes V_{5} \otimes V_{6} \\
& \rightarrow\left(V_{1} \otimes V_{4}\right) \otimes\left(V_{2} \otimes V_{5}\right) \otimes\left(V_{3} \otimes V_{6}\right) \\
& \rightarrow S^{2} V \otimes S^{2} V \otimes S^{2} V \\
& \rightarrow S^{3}\left(S^{2} V\right)
\end{aligned}
$$

Note that $h_{d, n}$ is a $G L(V)$-module map.

Example 9.1.1.2. For $(x y)^{2}=(x y)(x y) \in S^{2}\left(S^{2} \mathbb{C}^{2}\right)$, here is $h_{2,2}\left((x y)^{2}\right)$ :

$$
\begin{aligned}
(x y)^{2} & =\frac{1}{4}[(x \otimes y+y \otimes x) \otimes(x \otimes y+y \otimes x)] \\
& =\frac{1}{4}[x \otimes y \otimes x \otimes y+x \otimes y \otimes y \otimes x+y \otimes x \otimes x \otimes y+y \otimes x \otimes y \otimes x] \\
& \mapsto \frac{1}{4}[x \otimes x \otimes y \otimes y+x \otimes y \otimes y \otimes x+y \otimes x \otimes x \otimes y+y \otimes y \otimes x \otimes x] \\
& \mapsto \frac{1}{4}\left[2\left(x^{2}\right) \otimes\left(y^{2}\right)+2(x y) \otimes(x y)\right] \\
& \mapsto \frac{1}{2}\left[\left(x^{2}\right)\left(y^{2}\right)+(x y)(x y)\right] .
\end{aligned}
$$

Exercise 9.1.1.3: (1!) Show that $h_{d, n}\left(\left(x_{1}\right)^{n} \cdots\left(x_{d}\right)^{n}\right)=\left(x_{1} \cdots x_{d}\right)^{n}$.
Theorem 9.1.1.4 (Hadamard [Had97]). $\operatorname{ker} h_{d, n}=I_{d}\left(C h_{n}\left(V^{*}\right)\right)$.
Proof. Given $P \in S^{d}\left(S^{n} V\right)$, we determine if $P$ vanishes on $C h_{n}\left(V^{*}\right)$. Since $\operatorname{Seg}\left(v_{n}(\mathbb{P} V) \times \cdots \times v_{n}(\mathbb{P} V)\right)$ spans $\left(S^{n} V\right)^{\otimes d}$, its projection to $S^{d}\left(S^{n} V\right)$ also spans, so we may write $P=\sum_{j}\left(x_{1 j}\right)^{n} \cdots\left(x_{d j}\right)^{n}$ for some $x_{\alpha, j} \in V$. Let $\ell^{1}, \ldots, \ell^{n} \in V^{*}$. Recall $\bar{P}$ is $P$ considered as a linear form on $\left(S^{n} V^{*}\right)^{\otimes d}$. In what follows I use $\langle-,-\rangle$ to denote the pairing between a vector space and its dual.

$$
\begin{aligned}
P\left(\ell^{1} \cdots \ell^{n}\right) & =\left\langle\bar{P},\left(\ell^{1} \cdots \ell^{n}\right)^{d}\right\rangle \\
& =\sum_{j}\left\langle\left(x_{1 j}\right)^{n} \cdots\left(x_{d j}\right)^{n},\left(\ell^{1} \cdots \ell^{n}\right)^{d}\right\rangle \\
& =\sum_{j}\left\langle\left(x_{1 j}\right)^{n},\left(\ell^{1} \cdots \ell^{n}\right)\right\rangle \cdots\left\langle\left(x_{d j}\right)^{n},\left(\ell^{1} \cdots \ell^{n}\right)\right\rangle \\
& =\sum_{j} \Pi_{s=1}^{n} \Pi_{i=1}^{d} x_{i j}\left(\ell_{s}\right) \\
& =\sum_{j}\left\langle x_{1 j} \cdots x_{d j},\left(\ell^{1}\right)^{d}\right\rangle \cdots\left\langle x_{1 j} \cdots x_{d j},\left(\ell^{n}\right)^{d}\right\rangle \\
& =\left\langle\overline{h_{d, n}(P)},\left(\ell^{1}\right)^{d} \cdots\left(\ell^{n}\right)^{d}\right\rangle .
\end{aligned}
$$

If $h_{d, n}(P)$ is nonzero, there will be some monomial of the form $\left(\ell^{1}\right)^{d} \cdots\left(\ell^{n}\right)^{d}$ it will pair with to be nonzero (using the spanning property in $S^{n}\left(S^{d} V^{*}\right)$ ). On the other hand, if $h_{d, n}(P)=0$, then $\bar{P}$ annihilates all points of $C h_{n}\left(V^{*}\right)$.

### 9.1.2. Information on the rank of $h_{d, n}$.

Exercise 9.1.2.1: (2) Show that $h_{d, n}: S^{d}\left(S^{n} V\right) \rightarrow S^{n}\left(S^{d} V\right)$ is "self-dual" in the sense that $h_{d, n}^{T}=h_{n, d}: S^{n}\left(S^{d} V^{*}\right) \rightarrow S^{d}\left(S^{n} V^{*}\right)$. Conclude that $h_{d, n}$ surjective if and only if $h_{n, d}$ is injective.

Exercise 9.1.2.2: (1) Show that if $h_{d, n}: S^{d}\left(S^{n} \mathbb{C}^{m}\right) \rightarrow S^{n}\left(S^{d} \mathbb{C}^{m}\right)$ is not surjective, then $h_{d, n}: S^{d}\left(S^{n} \mathbb{C}^{k}\right) \rightarrow S^{n}\left(S^{d} \mathbb{C}^{k}\right)$ is not surjective for all $k>m$, and that the partitions corresponding to highest weights of the modules in the kernel are the same in both cases if $d \leq m$. ©
Exercise 9.1.2.3: (1) Show that if $h_{d, n}: S^{d}\left(S^{n} \mathbb{C}^{m}\right) \rightarrow S^{n}\left(S^{d} \mathbb{C}^{m}\right)$ is surjective, then $h_{d, n}: S^{d}\left(S^{n} \mathbb{C}^{k}\right) \rightarrow S^{n}\left(S^{d} \mathbb{C}^{k}\right)$ is surjective for all $k<m$.
Example 9.1.2.4 (The case $\operatorname{dim} V=2$ ). When $\operatorname{dim} V=2$, every polynomial decomposes as a product of linear factors, so the ideal of $C h_{n}\left(\mathbb{C}^{2}\right)$ is zero. We recover the following theorem of Hermite:
Theorem 9.1.2.5 (Hermite reciprocity). The map $h_{d, n}: S^{d}\left(S^{n} \mathbb{C}^{2}\right) \rightarrow$ $S^{n}\left(S^{d} \mathbb{C}^{2}\right)$ is an isomorphism for all $d, n$. In particular $S^{d}\left(S^{n} \mathbb{C}^{2}\right)$ and $S^{n}\left(S^{d} \mathbb{C}^{2}\right)$ are isomorphic $G L_{2}$-modules.

Often in modern textbooks (e.g., [FH91]) only the "In particular" is stated.

Originally Hadamard thought the maps $h_{d, n}$ were always of maximal rank, but later he realized he did not have a proof. In [Had99] he did prove:
Theorem 9.1.2.6 (Hadamard [Had99]). The map $h_{3,3}: S^{3}\left(S^{3} V\right) \rightarrow S^{3}\left(S^{3} V\right)$ is an isomorphism.

Proof. By Exercise 9.1.2.2, we may assume $\mathbf{v}=3$ and $x_{1}, x_{2}, x_{3} \in V^{*}$ are a basis. Say we had $P \in \operatorname{ker}\left(h_{3,3}\right)=I_{3}\left(C h_{3}\left(V^{*}\right)\right)$. Consider $P$ restricted to the line in $\mathbb{P}\left(S^{3} V^{*}\right)$ spanned by $x_{1}^{3}+x_{2}^{3}+x_{3}^{3}$ and $x_{1} x_{2} x_{3}$. Write $P(\mu, \nu):=$ $P\left(\mu\left(x_{1}^{3}+x_{2}^{3}+x_{3}^{3}\right)-\lambda x_{1} x_{2} x_{3}\right)$ as a cubic polynomial on $\mathbb{P}^{1}$ with coordinates $[\mu, \lambda]$. Note that $P(\mu, \nu)$ vanishes at the four points $[0,1],[1,3],[1,3 \omega],\left[1,3 \omega^{2}\right]$ where $\omega$ is a primitive third root of unity. A cubic polynomial on $\mathbb{P}^{1}$ vanishing at four points is identically zero, so the whole line is contained in Zeros $(P)$. In particular, $P(1,0)=0$, i.e., $P$ vanishes on $x_{1}^{3}+x_{2}^{3}+x_{3}^{3}$. Since $\sigma_{3}\left(v_{3}\left(\mathbb{P}^{2}\right)\right)$ is a $G L_{3}$-variety, $P$ must vanish identically on $\sigma_{3}\left(v_{3}\left(\mathbb{P}^{2}\right)\right)$. But $I_{3}\left(\sigma_{3}\left(v_{3}\left(\mathbb{P}^{2}\right)\right)\right)=0$, see, e.g., Corollary 8.3.4.3. (In fact $\sigma_{3}\left(v_{3}\left(\mathbb{P}^{2}\right)\right) \subset \mathbb{P} S^{3} \mathbb{C}^{3}$ is a hypersurface of degree four.)

In the same paper, he posed the question:
Question 9.1.2.7. Is $h_{d, n}$ always of maximal rank?
Howe [How87] also investigated the map $h_{d, n}$ and wrote "it is reasonable to expect" that $h_{d, n}$ is always of maximal rank.

Remark 9.1.2.8. The above proof is due to A. Abdesselam (personal communication). It is a variant of Hadamard's original proof, where instead of $x_{1}^{3}+x_{2}^{3}+x_{3}^{3}$ one uses an arbitrary cubic $f$, and generalizing $x_{1} x_{2} x_{3}$ one uses
the determinant of the Hessian $\operatorname{det}_{3}(H(f))$. Then the curves $f=0$ and $\operatorname{det}_{3}(H(f))=0$ intersect in 9 points (the nine flexes of $f=0$ ) and there are four groups of three lines going through these points, i.e., four places where the polynomial becomes a product of linear forms.

Theorem 9.1.2.9. [BL89] [also see [McK08, Thm. 8.1] and [Ike15]] If $h_{d, n}$ is surjective, then $h_{d^{\prime}, n}$ is surjective for all $d^{\prime}>d$. Equivalently, if $h_{d, n}$ is injective, then $h_{d, n^{\prime}}$ is injective for all $n^{\prime}>n$.

The proof is outlined in §9.3. The following two theorems were shown by a computer calculation:
Theorem 9.1.2.10. [MN05] The map $h_{4,4}$ is an isomorphism.
The results above imply $h_{d, n}$ is of maximal rank for all $n \leq 4$ and all $d$.
Theorem 9.1.2.11. [MN05] The map $h_{5,5}$ is not surjective.
Remark 9.1.2.12. In [MN05] they showed the map $h_{5,5: 0}$ defined in $\S 9.3$ below is not injective. A. Abdessalem realized their computation showed the map $h_{5,5}$ is not injective and pointed this out to them. Evidently there was some miscommunication because in [MN05] they mistakenly say the result comes from [Bri02] rather than their own paper.

The $G L(V)$-module structure of the kernel of $h_{5,5}$ was determined by M-W Cheung, C. Ikenmeyer and S. Mkrtchyan as part of a 2012 AMS MRC program:
Proposition 9.1.2.13. [CIM17] The kernel of $h_{5,5}: S^{5}\left(S^{5} \mathbb{C}^{5}\right) \rightarrow S^{5}\left(S^{5} \mathbb{C}^{5}\right)$ consists of irreducible modules corresponding to the following partitions:

$$
\begin{array}{r}
\{(14,7,2,2),(13,7,2,2,1),(12,7,3,2,1),(12,6,3,2,2) \\
(12,5,4,3,1),(11,5,4,4,1),(10,8,4,2,1),(9,7,6,3)\}
\end{array}
$$

All these occur with multiplicity one in the kernel, but not all occur with multiplicity one in $S^{5}\left(S^{5} \mathbb{C}^{5}\right)$. In particular, the kernel is not a sum of isotypic components.

It would be interesting to understand if there is a pattern to these partitions. Their Young diagrams are:



While the Hermite-Hadamard-Howe map is not always of maximal rank, it is "eventually" of maximal rank:
Theorem 9.1.2.14. [Bri93, Bri97] The Hermite-Hadamard-Howe map

$$
h_{d, n}: S^{d}\left(S^{n} V^{*}\right) \rightarrow S^{n}\left(S^{d} V^{*}\right)
$$

is surjective for $d$ sufficiently large, in fact for $d \gtrsim n^{2}\binom{n+d}{d}$.
I present the proof of Theorem 9.1.2.14 in §9.5.2.
Problem 9.1.2.15 (The Hadamard-Howe Problem). Determine the function $d(n)$ such that $h_{d, n}$ is surjective for all $d \geq d(n)$.

A more ambitious problem would be:
Problem 9.1.2.16. Determine the kernel of $h_{d, n}$.
A less ambitious problem is:
Problem 9.1.2.17. Improve Brion's bound to, say, a polynomial bound in $n$.

Another apparently less ambitious problem is the following conjecture:
Conjecture 9.1.2.18 (Kumar [?]). Let $n$ be even, then $S_{\left(n^{n}\right)} \mathbb{C}^{n} \not \subset$ ker $h_{n, n}$, i.e., $S_{\left(n^{n}\right)} \mathbb{C}^{n} \subset \mathbb{C}\left[C h_{n}\left(\mathbb{C}^{n}\right)\right]$.

Kumar conjectures further that for all $d \leq n, S_{\left(n^{d}\right)} \mathbb{C}^{n} \not \subset \operatorname{ker} h_{d, n}$, i.e., $S_{\left(n^{d}\right)} \mathbb{C}^{n} \subset \mathbb{C}\left[C h_{n}\left(\mathbb{C}^{n}\right)\right]$, but Conjecture 9.1.2.18 is the critical case. By Corollary 9.2.2.2 below, when $n$ is even, the module $S_{\left(n^{d}\right)} \mathbb{C}^{n}$ occurs in $S^{d}\left(S^{n} \mathbb{C}^{n}\right)$ with multiplicity one.

I discuss Conjecture 9.1.2.18 in $\S 9.4$. It turns out to be equivalent to a famous conjecture in combinatorics.

### 9.2. The GCT perspective

In this section, in the spirit of GCT, I compare $\mathbb{C}\left[C h_{n}\left(V^{*}\right)\right]=\mathbb{C}\left[\overline{G L(V) \cdot\left(x_{1} \cdots x_{n}\right)}\right]$ with $\mathbb{C}\left[G L(V) \cdot\left(x_{1} \cdots x_{n}\right)\right]$. Throughout this section, assume $\operatorname{dim} V=n$.
9.2.1. Application of the algebraic Peter-Weyl theorem. Let $x_{1}, \ldots, x_{n} \in$ $V^{*}$ be a basis. Recall from $\S 7.1 .2$ that the symmetry group of $x_{1} \cdots x_{n}$ is $\Gamma_{n}:=T^{S L_{n}} \rtimes \mathfrak{S}_{n}$. Also recall that for any orbit, $G / H$, the algebraic PeterWeyl theorem (see §8.6) implies $\mathbb{C}[G / H]=\bigoplus_{\lambda \in \Lambda_{G}^{+}} V_{\lambda} \otimes\left(V_{\lambda}^{*}\right)^{H}$, so

$$
\begin{equation*}
\mathbb{C}\left[G L(V) \cdot\left(x_{1} \cdots x_{n}\right)\right]=\bigoplus_{l(\pi) \leq n}\left(S_{\pi} V\right)^{\oplus \operatorname{dim}\left(S_{\pi} V^{*}\right)^{\Gamma_{n}}} \tag{9.2.1}
\end{equation*}
$$

where here $\pi=\left(p_{1}, \ldots, p_{n}\right)$ with $p_{j} \in \mathbb{Z}$ satisfies $p_{1} \geq p_{2} \geq \cdots \geq p_{n}$. (Note that the $p_{j}$ are allowed to be negative.) We break up the determination of $\left(S_{\pi} V^{*}\right)^{\Gamma_{n}}$ into two problems: first determine the $T^{S L_{n}-\text { invariants, and }}$ then the $\mathfrak{S}_{n}$ invariants inside the $T^{S L_{n}}$-invariants. By Exercise 8.1.5.4, the
 $|\pi|=s n$ for some $s \in \mathbb{Z}$. Let $\left(S_{\pi} V^{*}\right)_{0}$ denote the space of $T^{S L_{n}}$-invariants. The notation is chosen because this is the $\mathfrak{s l}(V)$-weight zero subspace.

It remains to determine $\left(S_{\pi} V^{*}\right)_{0}^{\mathfrak{S}_{n}}$. This is not known. In the next subsection, I relate it to another unknown quantity. Remarkably, this will enable us to get a satisfactory answer.
9.2.2. Plethysm and the double commutant theorem. The group $\mathfrak{S}_{n}$ acts on the $\mathfrak{s l}$-weight zero subspace by permuting basis vectors. (This is an example of a Weyl group action.) The following theorem is proved using the Double Commutant Theorem 8.6.2.3.
Theorem 9.2.2.1. [Gay76] Let $\mu$ be a partition of $n \delta$ (so that $\left(S_{\mu} V\right)_{0} \neq 0$ ). Suppose that the decomposition of $\left(S_{\mu} V\right)_{0}$ into irreducible $\mathfrak{S}_{n}$-modules is

$$
\left(S_{\mu} V\right)_{0}=\bigoplus_{|\pi|=n}[\pi]^{\oplus s_{\mu} \pi}
$$

Then the decomposition of $S_{\pi}\left(S^{\delta} V\right)$ as a $G L(V)$-module is

$$
S_{\pi}\left(S^{\delta} V\right)=\bigoplus_{|\mu|=\delta n}\left(S_{\mu} V\right)^{\oplus s_{\mu, \pi}}
$$

In particular, for $\delta=1$, i.e., $|\mu|=n,\left(S_{\mu} V\right)_{0}=[\mu]$.
Corollary 9.2 .2 .2 . For any partition $\pi$ of $d n$,

$$
\operatorname{mult}\left(S_{\pi} V, S^{n}\left(S^{d} \mathbb{C}^{n}\right)\right)=\operatorname{mult}\left([n],\left(S_{\pi} \mathbb{C}^{n}\right)_{0}\right)
$$

9.2.3. Back to the coordinate ring. Now specialize to the case of modules appearing in $\operatorname{Sym}\left(S^{n} V\right)$. Introduce the notation $M_{\text {poly }} \subset M$ where $M$ is a $G L(V)$-module and $M_{p o l y}$ is the submodule of the isotypic components of all $S_{\pi} V$ in $M$ where $\pi$ is a partition. (I.e., here I do require the parts of $\pi$ to be non-negative. The notation is used because these are the polynomial $G L(V)$-modules.) If we consider all the $\pi$ 's together, Corollary 9.2.2.2 combined with Equation (9.2.1) implies the following the equality of $G L(V)$-modules:

$$
\mathbb{C}\left[G L(V) \cdot\left(x_{1} \cdots x_{n}\right)\right]_{p o l y}=\bigoplus_{s} S^{n}\left(S^{s} V\right)
$$

In particular, $\bigoplus_{s} S^{n}\left(S^{s} V\right)$ inherits a graded ring structure. We'll return to this in $\S 9.5$.2. If $X$ is an irreducible affine variety and $X^{0} \subset X$ is a Zariski open subset (so $X=\overline{X^{0}}$ ), one has an injection

$$
\mathbb{C}[X] \rightarrow \mathbb{C}\left[X^{0}\right]
$$

by restriction of functions. We thus have a sequence of $G L(V)$-module maps

$$
\mathbb{C}\left[S^{n} V^{*}\right] \rightarrow \mathbb{C}\left[C h_{n}\left(V^{*}\right)\right] \rightarrow \mathbb{C}\left[G L(V) \cdot\left(x_{1} \cdots x_{n}\right)\right]_{\text {poly }}=\bigoplus_{s} S^{n}\left(S^{s} V\right),
$$

with the first surjective and the second injective. Their composition is a map

$$
h_{n}: \operatorname{Sym}\left(S^{n}(V)\right) \rightarrow \mathbb{C}\left[G L(V) \cdot\left(x_{1} \cdots x_{n}\right)\right]_{\text {poly }}
$$

with kernel $I\left(C h_{n}\left(V^{*}\right)\right)$. It should come as no surprise that in degree $d, h_{n}$ is $h_{d, n}$. A proof is given in $\S 9.5$. This gives us a second, GCT interpretation of the Hadamard-Howe map.

## 9.3. $\mathfrak{S}_{d n}$-formulation of the Hadamard-Howe problem

I now give an interpretation of the Hadamard-Howe problem in terms of maps of $\mathfrak{S}_{d n}$-modules.
9.3.1. The Black-List map. The dimension of $V$, as long as it is at least $d$, is irrelevant for the $G L(V)$-module structure of the kernel of $h_{d, n}$. In this section assume $\operatorname{dim} V=d n$.

If one restricts $h_{d, n}$ to the $\mathfrak{s l}(V)$-weight zero subspace, since the permutation of basis vectors commutes with $h_{d, n}$, one obtains a $\mathfrak{S}_{d n}$-module map

$$
\begin{equation*}
h_{d, n: 0}: S^{d}\left(S^{n} V\right)_{0} \rightarrow S^{n}\left(S^{d} V\right)_{0} . \tag{9.3.1}
\end{equation*}
$$

Let $\mathfrak{S}_{n} \backslash \mathfrak{S}_{d} \subset \mathfrak{S}_{d n}$ denote the wreath product, which, by definition, is the normalizer of $\mathfrak{S}_{n}^{\times d}$ in $\mathfrak{S}_{d n}$. It is the semi-direct product of $\mathfrak{S}_{n}^{\times d}$ with $\mathfrak{S}_{d}$, where $\mathfrak{S}_{d}$ acts by permuting the factors of $\mathfrak{S}_{n}^{\times d}$, see e.g., [Mac95, p 158]. The action of the group $\mathfrak{S}_{n} \prec \mathfrak{S}_{d}$ on $V^{\otimes d n}$ induced from the $\mathfrak{S}_{d n}$-action is
as follows: consider $V^{\otimes d n}$ as $\left(V^{\otimes n}\right)^{\otimes d}, d$ blocks of $n$-copies of $V$, permuting the $n$ copies of $V$ within each block as well as permuting the blocks. Thus $S^{d}\left(S^{n} V\right)=\left(V^{\otimes d n}\right)^{\mathfrak{S}_{n} \backslash \mathfrak{S}_{d}}$.

Notice that

$$
\left(V^{\otimes d n}\right)^{\mathfrak{S}_{n} i \mathfrak{S}_{d}}=\left(\bigoplus_{|\pi|=d n}[\pi] \otimes S_{\pi} V\right)^{\mathfrak{S}_{n} \backslash \mathfrak{G}_{d}}=\bigoplus_{|\pi|=d n}[\pi]^{\mathfrak{S}_{n} i \mathfrak{G}_{d}} \otimes S_{\pi} V,
$$

so

$$
\operatorname{mult}\left(S_{\pi} V, S^{d}\left(S^{n} V\right)\right)=\operatorname{dim}[\pi]^{\mathfrak{S}_{n} \backslash \mathfrak{S}_{d}} .
$$

Unfortunately the action of $\mathfrak{S}_{n} 乙 \mathfrak{S}_{d}$ is difficult to analyze.
In other words, recalling the discussion in $\S 9.2 .2$, as a $\mathfrak{S}_{d n}$-module map, (9.3.1) is

$$
\begin{equation*}
h_{d, n: 0}: \operatorname{Ind}_{\mathfrak{S}_{n} \mathfrak{S}_{d}}^{\mathfrak{S}_{d n}} \text { triv } \rightarrow \operatorname{Ind}_{\mathfrak{S}_{d} \mathfrak{S _ { n }}}^{\mathfrak{G}_{n}} \text { triv . } \tag{9.3.2}
\end{equation*}
$$

Call $h_{d, n: 0}$ the Black-List map. Since every irreducible module appearing in $S^{d}\left(S^{n} V\right)$ has a non-zero weight zero subspace, $h_{d, n}$ is the unique $G L(V)$ module extension of $h_{d, n: 0}$.

The above discussion shows that one can deduce the kernel of $h_{d, n}$ from that of $h_{d, n: 0}$ and vice versa. In particular, one is injective if and only if the other is, giving us our third interpretation of the Hadamard-Howe problem.

The map $h_{d, n: 0}$ was defined purely in terms of combinatorics in [BL89] as a path to try to prove the following conjecture of Foulkes:
Conjecture 9.3.1.1. [Fou50] Let $d>n$, let $\pi$ be a partition of $d n$ and let $[\pi]$ denote the corresponding $\mathfrak{S}_{d n}$-module. Then,

$$
\operatorname{mult}\left([\pi], \operatorname{Ind}_{\mathfrak{S}_{n i n} \mathfrak{S}_{d}}^{\mathfrak{S}_{d n}} \text { triv }\right) \geq \operatorname{mult}\left([\pi], \operatorname{Ind}_{\mathfrak{S}_{d} \mathfrak{\mathfrak { S } _ { n }}}^{\mathfrak{S}_{d n}} \text { triv }\right) .
$$

## Equivalently,

$$
\begin{equation*}
\operatorname{mult}\left(S_{\pi} V, S^{d}\left(S^{n} V\right)\right) \geq \operatorname{mult}\left(S_{\pi} V, S^{n}\left(S^{d} V\right)\right) \tag{9.3.3}
\end{equation*}
$$

Theorem 8.9.1.1 shows that equality holds asymptotically in (9.3.3), and Theorem 9.1.2.10 shows it holds for $d \leq 4$. In [CIM17] they show it also holds for $d=5$ by showing $h_{6,5}$ is surjective. Conjecture 9.3.1.1 is still open in general.
9.3.2. Outline of proof of Theorem 9.1.2.9. I prove that if $h_{d, n-1}$ is injective, then $h_{d, n}$ is injective. I follow the proof in [Ike15]. Write $W=$ $E \oplus F$ with $\operatorname{dim} E=d$ and $\operatorname{dim} F=n$. Give $E$ a basis $e_{1}, \ldots, e_{d}$ and $F$ a basis $f_{1}, \ldots, f_{n}$ inducing a basis of $W$ ordered $\left(e_{1}, e_{2}, \ldots, f_{n}\right)$. For a $G L(E) \times G L(F)$-weight $\alpha=\left(a_{1}, \ldots, a_{d}\right), \beta=\left(b_{1}, \ldots, b_{n}\right)$, let $\left(W^{\otimes d n}\right)_{(\alpha, \beta)}$ denote the $(\alpha, \beta)$ weight subspace of $W^{\otimes d n}$. Define the lowering map

$$
\phi_{i, j}:\left(W^{\otimes d n}\right)_{(\alpha, \beta)} \rightarrow\left(W^{\otimes d n}\right)_{\left(a_{1}, \ldots, a_{i-1},\left(a_{i}-1\right), a_{i+1}, \ldots, a_{d}\right), \beta=\left(b_{1}, \ldots,\left(b_{j}+1\right), \ldots, b_{n}\right)}
$$

induced from the map $W \rightarrow W$ that sends $e_{i}$ to $f_{j}$ and maps all other basis vectors to themselves. It is straight-forward to see the $\phi_{i, j}$ commute. Let $\phi_{d \times n}:\left(W^{\otimes d n}\right)_{\left(n^{d},(0)\right)} \rightarrow\left(W^{\otimes d n}\right)_{\left((0), d^{n}\right)}$ denote the composition of $\phi_{1,1} \cdots \phi_{d, b}$ restricted to $\left(W^{\otimes d n}\right)_{\left(n^{d},(0)\right)}$.

Call $\phi_{d \times n}$ the McKay map.
Proposition 9.3.2.1. As $\mathfrak{S}_{d n}$-module maps, $\phi_{d \times n}=h_{d, n ; 0}$, i.e., as maps of $\mathfrak{S}_{d n}$-modules, the McKay map is equivalent to the Black-List map.

The proof is indirect, by showing that the spaces coincide and the kernels are isomorphic $\mathfrak{S}_{d n}$-modules. More precisely:
Proposition 9.3.2.2. [Ikenmeyer, personal communication]
i) $\operatorname{mult}\left([\pi],\left(W^{\otimes d n}\right)_{\left(n^{d},(0)\right)}\right)=\operatorname{mult}\left(S_{\pi} W, S^{d}\left(S^{n} W\right)\right.$
ii) $\operatorname{mult}\left([\pi], \phi_{d \times n}\left(\left(W^{\otimes d n}\right)_{\left(n^{d},(0)\right)}\right)\right)=\operatorname{mult}\left(S_{\pi} W, h_{d, n}\left(S^{d}\left(S^{n} W\right)\right)\right.$.

Ikenmeyer proves Proposition 9.3.2.2 with explicit bases of both spaces defined via tableau. A posteriori this shows $\left(W^{\otimes d n}\right)_{\left(n^{d},(0)\right)}=\operatorname{Ind}_{\mathfrak{G}_{n} \backslash \mathfrak{S}_{d}}^{\mathfrak{G}_{d}}$ triv as $\mathfrak{S}_{d n}$-modules and $h_{d, n ; 0}=\phi_{d \times n}$.

Now for the proof of Theorem 9.1.2.9: We need to show $\phi_{d \times(n-1)}$ injective implies $\phi_{d \times n}$ is injective.

Reorder and decompose

$$
\phi_{d \times n}=\left(\phi_{1,1} \cdots \phi_{1, n-1} \phi_{2,1} \cdots \phi_{d, n-1}\right) \cdot\left(\phi_{1, n} \cdots \phi_{d, n}\right)
$$

and call the first term the left factor and the second the right factor. Each term in the left factor is injective by assumption. It remains to show injectivity of each $\phi_{i, n}$. I will show injectivity of $\phi_{i, n}$ restricted to each $\left(\left((n-1)^{i-1}, n^{d-i}\right),\left(0^{n-1}, i-1\right)\right)$ weight space. Each of these restrictions just involves a raising operator in the $\mathbb{C}^{2}$ with basis $e_{i}, f_{n}$, so we need to see the lowering map $\left(\left(\mathbb{C}^{2}\right)^{\otimes n+i-1}\right)_{(n, i-1)} \rightarrow\left(\left(\mathbb{C}^{2}\right)^{\otimes n+i-1}\right)_{(n-1, i)}$ is injective. Decompose

$$
\left(\mathbb{C}^{2}\right)^{\otimes n+i-1}=\oplus_{p_{2}=0}^{\left\lfloor\frac{n+i-1}{2}\right\rfloor} S_{n+i-1-p_{2}, p_{2}} \mathbb{C}^{2}
$$

The weight $(n-1, i)$ vector in each space may be written as $\left(e_{i} \wedge f_{n}\right)^{\otimes p_{2}} \otimes\left(e_{i}^{n-p_{2}} f_{n}^{i-1-p_{2}}\right)$. The lowering operator is zero on the first factor so this vector maps to $\left(e_{i} \wedge f_{n}\right)^{\otimes p_{2}} \otimes\left(e_{i}^{n-p_{2}-1} f_{n}^{i-p_{2}}\right)$ which is a basis vector in the target.

### 9.4. Conjecture 9.1.2.18 and a conjecture in combinatorics

For any even $n$, the one-dimensional module $S_{\left(n^{d}\right)} \mathbb{C}^{d}$ occurs with multiplicity one in $S^{d}\left(S^{n} \mathbb{C}^{d}\right)\left(\right.$ cf. [How87, Prop. 4.3]). Let $P \in S_{n^{d}}\left(\mathbb{C}^{d}\right) \subset S^{d}\left(S^{n} \mathbb{C}^{d}\right)$ be non-zero. Conjecture 9.1.2.18 and its generalizations may be stated as $P\left(\left(x_{1} \cdots x_{n}\right)^{d}\right) \neq 0$. Our first task is to obtain an expression for $P$.
9.4.1. Realization of the module. Let $V=\mathbb{C}^{d}$. Fix a nonzero basis element $\operatorname{det}_{d} \in \Lambda^{d} V$.
Proposition 9.4.1.1. [KL15] Let $n$ be even. The unique (up to scale) element $P \in S_{\left(n^{d}\right)} V \subset S^{d}\left(S^{n} V\right)$ evaluated on

$$
x=\left(v_{1}^{1} \cdots v_{n}^{1}\right)\left(v_{1}^{2} \cdots v_{n}^{2}\right) \cdots\left(v_{1}^{d} \cdots v_{n}^{d}\right) \in S^{d}\left(S^{n} V^{*}\right), \text { for any } v_{j}^{i} \in V^{*},
$$

is

$$
\begin{equation*}
\langle P, x\rangle=\sum_{\sigma_{1}, \ldots, \sigma_{d} \in \mathfrak{S}_{n}} \operatorname{det}_{d}\left(v_{\sigma_{1}(1)}^{1} \wedge \cdots \wedge v_{\sigma_{d}(1)}^{d}\right) \cdots \operatorname{det}_{d}\left(v_{\sigma_{1}(n)}^{1} \wedge \cdots \wedge v_{\sigma_{d}(n)}^{d}\right) \tag{9.4.1}
\end{equation*}
$$

Proof. Let $\tilde{P}(x)$ denote the right hand side of (9.4.1), so $\tilde{P} \in(V)^{\otimes n d}$. It suffices to check that
(i) $\tilde{P} \in S^{d}\left(S^{n} V\right)$,
(ii) $\tilde{P}$ is $S L(V)$ invariant, and
(iii) $\tilde{P}$ is not identically zero.

Observe that (iii) follows from the identity (9.4.1) by taking $v_{j}^{i}=x_{i}$ where $x_{1}, \ldots, x_{d}$ is a basis of $V^{*}$, and (ii) follows because $S L(V)$ acts trivially on $\operatorname{det}_{d}$.

To prove (i), I show (ia) $\tilde{P} \in S^{d}\left(V^{\otimes n}\right)$ and (ib) $\tilde{P} \in\left(S^{n} V\right)^{\otimes d}$ to conclude. To see (ia), it is sufficient to show that exchanging two adjacent factors in parentheses in the expression of $x$ will not change (9.4.1). Exchange $v_{j}^{1}$ with $v_{j}^{2}$ in the expression for $j=1, \ldots, n$. Then, each individual determinant will change sign, but there are an even number of determinants, so the right hand side of (9.4.1) is unchanged. To see (ib), it is sufficient to show the expression is unchanged if we swap $v_{1}^{1}$ with $v_{2}^{1}$ in (9.4.1). If we multiply by $n$ !, we may assume $\sigma_{1}=\mathrm{Id}$, i.e.,
$\langle\tilde{P}, x\rangle=$
$n!\sum_{\sigma_{2}, \ldots, \sigma_{d} \in \mathfrak{S}_{n}} \operatorname{det}_{d}\left(v_{1}^{1}, v_{\sigma_{2}(1)}^{2}, \ldots, v_{\sigma_{d}(1)}^{d}\right) \operatorname{det}_{d}\left(v_{2}^{1} \wedge v_{\sigma_{2}(2)}^{2} \wedge \cdots \wedge v_{\sigma_{d}(2)}^{d}\right) \cdots \operatorname{det}_{d}\left(v_{n}^{1} \wedge v_{\sigma_{2}(n)}^{2} \wedge \cdots \wedge v_{\sigma_{d}(n)}^{d}\right)$.
With the two elements $v_{1}^{1}$ and $v_{2}^{1}$ swapped, we get
$n!\sum_{\sigma_{2}, \ldots, \sigma_{d} \in \mathfrak{S}_{n}} \operatorname{det}_{d}\left(v_{2}^{1} \wedge v_{\sigma_{2}(1)}^{2} \wedge \cdots \wedge v_{\sigma_{d}(1)}^{d}\right) \operatorname{det}_{d}\left(v_{1}^{1} \wedge v_{\sigma_{2}(2)}^{2} \wedge \cdots \wedge v_{\sigma_{d}(2)}^{d}\right) \cdots \operatorname{det}_{d}\left(v_{n}^{1} \wedge v_{\sigma_{2}(n)}^{2} \wedge \cdots \wedge v_{\sigma_{d}(n)}^{d}\right)$.
Now right compose each $\sigma_{s}$ in (9.4.2) by the transposition $(1,2)$. The expressions become the same.

Corollary 9.4.1.2. The unique (up to scale) polynomial $P \in S_{\left(n^{d}\right)} V \subset$ $S^{d}\left(S^{n} V\right)$ when $n$ is even, is nonzero on $\left(y_{1}\right)^{n}+\cdots+\left(y_{d}\right)^{n}$ if the $y_{j}$ are linearly independent. In particular, $S_{n^{d}} V \subset \mathbb{C}\left[\sigma_{d}\left(v_{n}\left(\mathbb{P} V^{*}\right)\right)\right]$ whenever $\operatorname{dim} V \geq d$.

Proof. The monomial $\left(y_{1}\right)^{n} \cdots\left(y_{d}\right)^{n}$ appears in $\left(\left(y_{1}\right)^{n}+\cdots+\left(y_{d}\right)^{n}\right)^{d}$ and all other monomials appearing pair with $P$ to be zero.

Now specialize to the critical case $d=n$ and evaluate on $\left(x_{1} \cdots x_{n}\right)^{n}$, where $x_{1}, \ldots, x_{n}$ is a basis of $V^{*}$ such that $\operatorname{det}_{n}\left(x_{1} \wedge \cdots \wedge x_{n}\right)=1$.

$$
\begin{equation*}
\left\langle P,\left(x_{1} \cdots x_{n}\right)^{n}\right\rangle=\sum_{\sigma_{1}, \ldots, \sigma_{n} \in \mathfrak{S}_{n}} \operatorname{det}_{d}\left(x_{\sigma_{1}(1)}, \ldots, x_{\sigma_{n}(1)}\right) \cdots \operatorname{det}_{d}\left(x_{\sigma_{1}(n)}, \ldots, x_{\sigma_{n}(n)}\right) . \tag{9.4.3}
\end{equation*}
$$

For a fixed $\left(\sigma_{1}, \ldots, \sigma_{n}\right)$ the contribution will either be 0,1 or -1 . The contribution is zero unless for each $j$, the indices $\sigma_{1}(j), \ldots, \sigma_{n}(j)$ are distinct. Arrange these numbers in an array:

$$
\left(\begin{array}{ccc}
\sigma_{1}(1) & \cdots & \sigma_{n}(1) \\
& \vdots & \\
\sigma_{1}(n) & \cdots & \sigma_{n}(n)
\end{array}\right)
$$

The contribution is zero unless the array is a Latin square, i.e., an $n \times n$ matrix such that each row and column consists of the integers $\{1, \ldots, n\}$. If it is a Latin square, the rows correspond to permutations, and the contribution of the term is the product of the signs of these permutations. Call this the row sign of the Latin square. The products of both the signs of the row permutations and the column permutations is called the sign of the Latin square:
Conjecture 9.4.1.3 (Alon-Tarsi [AT92]). Let $n$ be even. The number of sign -1 Latin squares of size $n$ is not equal to the number of sign +1 Latin squares of size $n$.

Conjecture 9.4.1.3 is known to be true when $n=p \pm 1$, where $p$ is an odd prime; in particular, it is known to be true up to $n=24$ [Gly10, Dri97].

On the other hand, in $[\mathbf{A l p 1 4}, \mathbf{C W 1 6}]$ they show that the ratio of the number of sign -1 Latin squares of size $n$ to the number of sign +1 Latin squares of size $n$ tends to one as $n$ goes to infinity.

In [HR94], Huang and Rota showed:
Theorem 9.4.1.4. [HR94, Identities 8,9] The difference between the number of row even Latin squares of size $n$ and the number of row odd Latin squares of size $n$ equals the difference between the number of even Latin squares of size $n$ and the number of odd Latin squares of size $n$, up to sign.

In particular, the Alon-Tarsi conjecture holds for $n$ if and only if the row-sign Latin square conjecture holds for $n$. Thus
Theorem 9.4.1.5. [KL15] The Alon-Tarsi conjecture holds for $n$ if and only if $S_{\left(n^{n}\right)}\left(\mathbb{C}^{n}\right) \in \mathbb{C}\left[C h_{n}\left(\mathbb{C}^{n}\right)\right]$.

In [KL15] several additional statements equivalent to the conjecture were given. In particular, for those familiar with integration over compact Lie groups, the conjecture holds for $n$ if and only if

$$
\int_{\left(g_{j}^{i}\right) \in S U(n)} \prod_{1 \leq i, j \leq n} g_{j}^{i} d \mu \neq 0
$$

where $d \mu$ is Haar measure.

### 9.5. Algebraic geometry derivation of the Hadamard-Howe map

9.5.1. Preliminaries from algebraic geometry. In modern geometry, one studies a space by the functions on it. The general definition of an affine variety over $\mathbb{C}$ can be made without reference to an ambient space. It corresponds to a finitely generated ring over $\mathbb{C}$ with no nilpotent elements (see, e.g., [Har95, Lect. 5]), as these are the rings that are realizable as the ring of regular functions of a subvariety of affine space. In this section we will deal with two affine varieties that are presented to us in terms of their rings of regular functions, the normalization of an affine variety and the GIT quotient of an affine variety with an algebraic group action.

If $R, S$ are rings with $R \subset S, s \in S$ is integral over $R$ if it satisfies a monic polynomial with coefficients in $R: s^{d}+r_{1} s^{d-1}+\cdots+r_{d}=0$ for some $r_{i} \in R$, and $S$ is integral over $R$ if every element of $S$ is integral over $R$.

A regular map (see §3.1.4) between affine varieties $f: X \rightarrow Y$ such that $f(X)$ is dense in $Y$ is said to be finite if $\mathbb{C}[X]$ is integral over $\mathbb{C}[Y]$ (see, e.g. [Sha13a, §I.5.3]).

An affine variety $Z$ is normal if $\mathbb{C}[Z]$ is integrally closed, that is if every element of $\mathbb{C}(Z)$, the field of fractions of $\mathbb{C}[Z]$, that is integral over $\mathbb{C}[Z]$ is in $\mathbb{C}[Z]$. To every affine variety $Z$ one may associate a unique normal affine variety $\operatorname{Nor}(Z)$, called the normalization of $Z$, such that there is a generically one to one finite map $\pi: \operatorname{Nor}(Z) \rightarrow Z$. If $Z$ is smooth then $\operatorname{Nor}(Z)=Z$, and more generally $\pi$ is one to one over the smooth points of $Z$. For details see [Sha13a, §II.5].
Exercise 9.5.1.1: (1) Show that if $Z$ is a $G$-variety, then $\operatorname{Nor}(Z)$ is too.
Recall from Exercise 3.1.4.6 the inclusion $\mathbb{C}[Z] \rightarrow \mathbb{C}[\operatorname{Nor}(Z)]$ given by pullback of functions. If the non-normal points of $Z$ form a finite set, then the cokernel of this inclusion is finite dimensional.
9.5.2. Coordinate ring of the normalization of the Chow variety. In this section I work in affine space and follow [Bri93]. The normalization of the (cone over the) Chow variety and its coordinate ring have a simple
description that I now explain. The cone $\hat{C} h_{n}\left(V^{*}\right) \subset S^{n} V^{*}$ is the image of the following map:

$$
\begin{aligned}
\phi_{n}: V^{* \times n} & \rightarrow S^{n} V^{*} \\
\left(u_{1}, \ldots, u_{n}\right) & \mapsto u_{1} \cdots u_{n} .
\end{aligned}
$$

Note that $\phi_{n}$ is $G L(V)$-equivariant.
For any affine algebraic group $\Gamma$ and any affine $\Gamma$-variety $Z$, define the GIT quotient $Z / / \Gamma$ to be the affine algebraic variety whose coordinate ring is $\mathbb{C}[Z]^{\Gamma}$. (When $\Gamma$ is finite, this is just the usual set-theoretic quotient. In the general case two $\Gamma$-orbits will be identified under the quotient map $Z \rightarrow Z / / \Gamma$ when there is no $\Gamma$-invariant regular function that can distinguish them.)
Exercise 9.5.2.1: (2!) Consider the space of $n \times n$ matrices $M a t_{n}$ with the action of $G L_{n}$ via conjugation. Give an explicit description of the map $M a t_{n} \rightarrow M a t_{n} / / G L_{n}$. ©
Exercise 9.5.2.2: (2) Show that if $Z$ is normal, then so is $Z / / \Gamma$. ©
In our case $V^{* \times n}$ is an affine $\Gamma_{n}:=T^{S L_{n}} \rtimes \mathfrak{S}_{n}$-variety, where a diagonal matrix in $T^{S L_{n}}$ with entries $\lambda_{j}$ acts on $V^{* \times n}$ by $\left(\alpha_{1}, \ldots, \alpha_{n}\right) \mapsto$ $\left(\lambda_{1} \alpha_{1}, \ldots, \lambda_{n} \alpha_{n}\right)$. The map $\phi_{n}$ factors through the GIT quotient because it is $\Gamma_{n}$-equivariant, giving a map

$$
\begin{equation*}
\psi_{n}: V^{* \times n} / / \Gamma_{n} \rightarrow S^{n} V^{*} \tag{9.5.1}
\end{equation*}
$$

whose image is $\hat{C} h_{n}\left(V^{*}\right)$. By unique factorization, $\psi_{n}$ is generically one to one. Elements of $V^{* \times n}$ of the form $\left(0, u_{2}, \ldots, u_{n}\right)$ cannot be distinguished from $(0, \ldots, 0)$ by $\Gamma_{n}$-invariant functions, so they are identified with $(0, \ldots, 0)$ in the quotient, which is consistent with the fact that $\phi_{n}\left(0, u_{2}, \ldots, u_{n}\right)=$ 0 .

Consider the induced map on coordinate rings:

$$
\psi_{n}^{*}: \mathbb{C}\left[S^{n} V^{*}\right] \rightarrow \mathbb{C}\left[V^{* \times n} / / \Gamma_{n}\right]=\mathbb{C}\left[V^{* \times n}\right]^{\Gamma_{n}} .
$$

For affine varieties, $\mathbb{C}[Y \times Z]=\mathbb{C}[Y] \otimes \mathbb{C}[Z]$ (see e.g., [Sha13a, $\S 2.2$ Ex.1.10]), so

$$
\begin{aligned}
\mathbb{C}\left[V^{* \times n}\right] & =\mathbb{C}\left[V^{*}\right]^{\otimes n} \\
& =\operatorname{Sym}(V) \otimes \cdots \otimes \operatorname{Sym}(V) \\
& =\bigoplus_{i_{1}, \ldots, i_{n} \in \mathbb{Z}_{\geq 0}} S^{i_{1}} V \otimes \cdots \otimes S^{i_{n}} V .
\end{aligned}
$$

Taking $T^{S L_{n}}$ invariants gives

$$
\mathbb{C}\left[V^{* \times n}\right]^{T_{n}^{S L}}=\bigoplus_{i \geq 0} S^{i} V \otimes \cdots \otimes S^{i} V
$$

and finally

$$
\left(\mathbb{C}\left[V^{* \times n}\right]^{T_{n}^{S L}}\right)^{\mathfrak{G}_{n}}=\bigoplus_{i \geq 0} S^{n}\left(S^{i} V\right)
$$

The map

$$
\tilde{h}_{n}:=\psi_{n}^{*}: \operatorname{Sym}\left(S^{n} V\right) \rightarrow \oplus_{i}\left(S^{n}\left(S^{i} V\right)\right),
$$

respects $G L$-degree, so it gives rise to maps $\tilde{h}_{d, n}: S^{d}\left(S^{n} V\right) \rightarrow S^{n}\left(S^{d} V\right)$.
Proposition 9.5.2.3. $\tilde{h}_{d, n}=h_{d, n}$.
Proof. Since elements of the form $x_{1}^{n} \cdots x_{d}^{n}$ span $S^{d}\left(S^{n} V\right)$ it will be sufficient to prove the maps agree on such elements. By Exercise 9.1.1.3, $h_{d, n}\left(x_{1}^{n} \cdots x_{d}^{n}\right)=\left(x_{1} \cdots x_{d}\right)^{n}$. On the other hand, in the algebra $\mathbb{C}\left[V^{*}\right]^{\otimes n}$, the multiplication is $\left(f_{1} \otimes \cdots \otimes f_{n}\right) \odot\left(g_{1} \otimes \cdots \otimes g_{n}\right)=f_{1} g_{1} \otimes \cdots \otimes f_{n} g_{n}$ and this descends to the algebra $\left(\mathbb{C}\left[V^{*}\right]^{\otimes n}\right)^{\Gamma_{n}}$ which is the target of the algebra $\operatorname{map} \tilde{h}_{n}$, i.e.,

$$
\begin{aligned}
\tilde{h}_{d, n}\left(x_{1}^{n} \cdots x_{d}^{n}\right) & =\psi_{n}^{*}\left(x_{1}^{n} \cdots x_{d}^{n}\right) \\
& =\psi_{n}^{*}\left(x_{1}^{n}\right) \odot \cdots \odot \psi_{n}^{*}\left(x_{d}^{n}\right) \\
& =x_{1}^{n} \odot \cdots \odot x_{d}^{n} \\
& =\left(x_{1} \cdots x_{d}\right)^{n} .
\end{aligned}
$$

Proposition 9.5.2.4. $\psi_{n}: V^{* \times n} / / \Gamma_{n} \rightarrow \hat{C} h_{n}\left(V^{*}\right)$ is the normalization of $\hat{C} h_{n}\left(V^{*}\right)$.

I prove Proposition 9.5.2.4 in $\S 9.7$. Thus we get a fourth formulation of the Hadamard-Howe problem: Determine the cokernel of the natural inclusion map

$$
\mathbb{C}\left[\hat{C} h\left(V^{*}\right)\right] \rightarrow \mathbb{C}\left[\operatorname{Nor}\left(\hat{C} h\left(V^{*}\right)\right)\right]
$$

This is equivalent to the other formulations because the cokernel of $\tilde{h}_{n}$ is also the cokernel of the composition $\operatorname{Sym}\left(S^{n} V\right) \rightarrow \mathbb{C}\left[\hat{C} h_{n}\left(V^{*}\right)\right] \rightarrow \mathbb{C}\left[\operatorname{Nor}\left(\hat{C} h\left(V^{*}\right)\right)\right]$. The proof of Proposition 9.5.2.4 and the qualitative assertion of Theorem 9.1.2.14 will hinge on exploiting that the only non-normal point of $\hat{C} h\left(V^{*}\right)$ is the origin. Since it involves more advanced results from algebraic geometry, I postpone the proofs until the end of this chapter.

### 9.6. Brill's equations

Set theoretic equations of $C h_{d}\left(V^{*}\right)$ have been known since 1894. Here is a modern presentation elaborating the presentation in [Lan12, §8.6], which was suggested by E. Briand.
9.6.1. Preliminaries. Our goal is a polynomial test to determine if $f \in$ $S^{d} V^{*}$ is a product of linear factors. We can first try to just determine if $f$ is divisible by a power of a linear form. The discussion in $\S 8.4 .2$ will not be helpful as the conditions there are vacuous when $n-m=1$. We could proceed as in $\S 6.5 .1$ and check if $\ell x^{I_{1}} \wedge \cdots \wedge \ell x^{I_{D}} \wedge f=0$ where the $x^{I_{j}}$ are a basis of $S^{d-1} V^{*}$, but in this case there is a simpler test to see if a given linear form $\ell$ divides $f$ :

Consider the map $\pi_{d, d}: S^{d} V^{*} \otimes S^{d} V^{*} \rightarrow S_{(d, d)} V^{*}$ obtained by projection. (By the Pieri rule 8.1.3.1, $S_{(d, d)} V^{*} \subset S^{d} V^{*} \otimes S^{d} V^{*}$ with multiplicity one.)
Lemma 9.6.1.1. Let $\ell \in V^{*}, f \in S^{d} V^{*}$. Then $f=\ell h$ for some $h \in S^{d-1} V^{*}$ if and only if $\pi_{d, d}\left(f \otimes \ell^{d}\right)=0$.

Proof. Since $\pi_{d, d}$, is linear, it suffices to prove the lemma when $f=\ell_{1} \cdots \ell_{d}$. In that case $\pi_{d, d}\left(f \otimes \ell^{d}\right)$, up to a constant, is $\left(\ell_{1} \wedge \ell\right) \cdots\left(\ell_{d} \wedge \ell\right)$.

We would like a map that sends $\ell_{1} \cdots \ell_{d}$ to $\sum_{j} \ell_{j}^{d} \otimes s t u f f_{j}$, as then we could apply $\pi_{d, d} \otimes \mathrm{Id}$ to $f$ tensored with the result of our desired map to obtain equations. I construct such a map in several steps.

The maps $f \mapsto f_{j, d-j}$ send $\left(\ell_{1} \cdots \ell_{d}\right)$ to $\sum_{|K|=j} \ell_{K} \otimes \ell_{K^{c}}$ where $\ell_{K}=$ $\ell_{k_{1}} \cdots \ell_{k_{j}}$ and $K^{c}$ denotes the complementary index set in $[d]$. The $\ell_{K}$ are monomials appearing in elementary symmetric functions and the idea is to convert this to power sums by the conversion formula obtained from the relation between generating functions (6.1.5):

$$
p_{d}=\mathcal{P}_{d}\left(e_{1}, \ldots, e_{d}\right):=\operatorname{det}\left(\begin{array}{ccccc}
e_{1} & 1 & 0 & \cdots & 0  \tag{9.6.1}\\
2 e_{2} & e_{1} & 1 & \cdots & 0 \\
\vdots & \vdots & \vdots & & \vdots \\
d e_{d} & e_{d-1} & e_{d-2} & \cdots & e_{1}
\end{array}\right)
$$

The desired term comes from the diagonal $e_{1}^{d}$ and the rest of the terms kill off the unwanted terms of $e_{1}^{d}$. This idea almost works- the only problem is that our naïve correction terms have the wrong degree on the right hand side. For example, when $d=3$, naïvely using $p_{3}=e_{1}^{3}-3 e_{1} e_{2}+3 e_{3}$ would give, for the first term, degree $6=2+2+2$ on the right hand side of the tensor product, the second degree $3=2+1$ and the third degree zero. In general, the right hand side of the $e_{1}^{d}$ term would have degree $(d-1)^{d}$, whereas the $d e_{d}$ term would have degree zero. In addition to fixing the degree mismatch, we need to formalize how we will treat the right hand sides.

Define maps

$$
\begin{align*}
E_{j}: S^{\delta} V^{*} & \rightarrow S^{j} V^{*} \otimes S^{\delta-1} V^{*}  \tag{9.6.2}\\
f & \mapsto f_{j, \delta-j} \odot\left(1 \otimes f^{j-1}\right)
\end{align*}
$$

where © is the multiplication introduced in the proof of Proposition 9.5.2.3. The $\left(1 \otimes f^{j-1}\right)$ fixes our degree problem. If $j>\delta$ define $E_{j}(f)=0$.

Our desired map is

$$
\begin{align*}
Q_{d}: S^{d} V^{*} & \rightarrow S^{d} V^{*} \otimes S^{d(d-1)} V^{*}  \tag{9.6.3}\\
f & \mapsto \mathcal{P}_{d}\left(E_{1}(f), \ldots, E_{d}(f)\right) .
\end{align*}
$$

9.6.2. Statement and proof of Brill's equations. Define Brill's map

$$
\begin{align*}
\mathcal{B}: S^{d} V^{*} & \rightarrow S_{d, d} V^{*} \otimes S^{d^{2}-d} V^{*}  \tag{9.6.4}\\
f & \mapsto\left(\pi_{d, d} \otimes \operatorname{Id}_{S^{d^{2}-d} V^{*}}\right)\left[f \otimes Q_{d}(f)\right] .
\end{align*}
$$

Theorem 9.6.2.1 (Brill [Bri93], Gordan [Gor94], Gelfand-Kapranov-Zelevinski [GKZ94], Briand [Bri10]). $[f] \in C h_{d}\left(V^{*}\right)$ if and only if $\mathcal{B}(f)=0$.

The proof will be by an induction argument that will require a generalization of $Q_{d}$. Define

$$
\begin{align*}
Q_{d, \delta}: S^{\delta} V^{*} & \rightarrow S^{d} V^{*} \otimes S^{d(\delta-1)} V^{*}  \tag{9.6.5}\\
f & \mapsto \mathcal{P}_{d}\left(E_{1}(f), \ldots, E_{d}(f)\right) .
\end{align*}
$$

Lemma 9.6.2.2. If $f_{1} \in S^{\delta} V^{*}$ and $f_{2} \in S^{d^{\prime}-\delta} V^{*}$, then

$$
Q_{d, d^{\prime}}\left(f_{1} f_{2}\right)=\left(1 \otimes f_{1}^{d}\right) \odot Q_{d, d^{\prime}-\delta}\left(f_{2}\right)+\left(1 \otimes f_{2}^{d}\right) \odot Q_{d, \delta}\left(f_{1}\right) .
$$

Assume Lemma 9.6.2.2 for the moment.
Proof of Theorem 9.6.2.1. Say $f=\ell_{1} \cdots \ell_{d}$. First note that for $\ell \in$ $V^{*}, E_{j}\left(\ell^{j}\right)=\ell^{j} \otimes \ell^{j-1}$ and $Q_{d, 1}(\ell)=\ell^{d} \otimes 1$. Next, compute $E_{1}\left(\ell_{1} \ell_{2}\right)=$ $\ell_{1} \otimes \ell_{2}+\ell_{2} \otimes \ell_{1}$ and $E_{2}\left(\ell_{1} \ell_{2}\right)=\ell_{1} \ell_{2} \otimes \ell_{1} \ell_{2}$, so $Q_{2,2}\left(\ell_{1} \ell_{2}\right)=\ell_{1}^{2} \otimes \ell_{2}^{2}+\ell_{2}^{2} \otimes \ell_{1}^{2}$. By induction and Lemma 9.6.2.2,

$$
Q_{d, \delta}\left(\ell_{1} \cdots \ell_{\delta}\right)=\sum_{j} \ell_{j}^{d} \otimes\left(\ell_{1}^{d} \cdots \ell_{j-1}^{d} \ell_{j+1}^{d} \cdots \ell_{\delta}^{d}\right)
$$

We conclude $Q_{d}(f)=\sum_{j} \ell_{j}^{d} \otimes\left(\ell_{1}^{d} \cdots \ell_{j-1}^{d} \ell_{j+1}^{d} \cdots \ell_{d}^{d}\right)$ and $\pi_{d, d}\left(\ell_{1} \cdots \ell_{d}, \ell_{j}^{d}\right)=0$ for each $j$ by Lemma 9.6.1.1.

Now assume $\mathcal{B}(f)=0$ and we will see $[f] \in C h_{d}\left(V^{*}\right)$. Compute $Q_{d}(f)=$ $\left(E_{1}(f)\right)^{d}+\sum \mu_{j} \otimes \psi_{j}$ where $\psi_{j} \in S^{d^{2}-d} V^{*}, \mu_{j} \in S^{d} V^{*}$ and $f$ divides $\psi_{j}$ for each $j$ because $E_{1}(f)^{d}$ occurs as a monomial in the determinant (9.6.1) and all the other terms contain an $E_{j}(f)$ with $j>1$, and so are divisible by $f$.

First assume $f$ is reduced, i.e., has no repeated factors, then every component of $\operatorname{Zeros}(f)$ contains a smooth point. Let $z \in \operatorname{Zeros}(f)_{\text {smooth }}$. Thus $\mathcal{B}(f)(\cdot, z)=\pi_{d, d}\left(f \otimes\left(d f_{z}\right)^{d}\right)$ because $E_{1}(f)^{d}=\left(f_{1, d-1}\right)^{d}$ and $f_{1, d-1}(\cdot, z)=$ $d f_{z}$, and all the $\psi_{j}(z)$ are zero. By Lemma 9.6.1.1, $d f_{z}$ divides $f$ for all $z \in \operatorname{Zeros}(f)$. But this implies the tangent space to $f$ is constant in a neighborhood of $z$, i.e., that the component containing $z$ is a linear space. So
when $f$ is reduced, $\operatorname{Zeros}(f)$ is a union of hyperplanes, which is what we set out to prove.

Finally, say $f=g^{k} h$ where $g$ is irreducible of degree $q$ and $h$ is of degree $d-q k$ and is relatively prime to $g$. Apply Lemma 9.6.2.2:

$$
Q_{d}\left(g\left(g^{k-1} h\right)\right)=\left(1 \otimes g^{d}\right) \odot Q_{d, d-q}\left(g^{k-1} h\right)+\left(1 \otimes\left(g^{k-1} h\right)^{d}\right) \odot Q_{d, q}(g) .
$$

A second application gives

$$
\begin{gathered}
Q_{d}\left(g^{k} h\right)=\left(1 \otimes g^{d}\right) \odot\left[\left(1 \otimes g^{d}\right) \odot Q_{d, d-2 q}\left(g^{k-2} h\right)+\left(1 \otimes\left(g^{k-2} h\right)^{d}\right) \odot Q_{d, q}(g)\right. \\
\left.+\left(1 \otimes\left(g^{k-2} h\right)^{d}\right) \odot Q_{d, q}(g)\right] .
\end{gathered}
$$

After $k-1$ applications one obtains:

$$
Q_{d}\left(g^{k} h\right)=\left(1 \otimes g^{d(k-1)}\right) \odot\left[k\left(1 \otimes h^{d}\right) \odot Q_{d, q}(g)+\left(1 \otimes g^{d}\right) \odot Q_{d, d-q k}(h)\right]
$$

and $\left(1 \otimes g^{d(k-1)}\right)$ will also factor out of $\mathcal{B}(f)$. Since $\mathcal{B}(f)$ is identically zero but $g^{d(k-1)}$ is not, we conclude

$$
0=\pi_{d, d} \otimes \operatorname{Id}_{S^{d^{2}-d} V^{*}} f \otimes\left[k\left(1 \otimes h^{d}\right) \odot Q_{d, q}(g)+\left(1 \otimes g^{d}\right) \odot Q_{d, d-q k}(h)\right]
$$

Let $w \in \operatorname{Zeros}(g)$ be a general point, so in particular $h(w) \neq 0$. Evaluating at $(z, w)$ with $z$ arbitrary gives zero on the second term and the first implies $\pi_{d, d} \otimes \operatorname{Id}_{S^{d^{2}-d} V^{*}}\left(f \otimes Q_{d, q}(g)\right)=0$ which implies $d g_{w}$ divides $g$, so $g$ is a linear form. Applying the argument to each non-reduced factor of $f$ we conclude.

Proof of Lemma 9.6.2.2. Define, for $u \in \operatorname{Sym}\left(V^{*}\right) \otimes \operatorname{Sym}\left(V^{*}\right)$,

$$
\begin{aligned}
\Delta_{u}: \operatorname{Sym}\left(V^{*}\right) & \rightarrow \operatorname{Sym}\left(V^{*}\right) \otimes \operatorname{Sym}\left(V^{*}\right) \\
f & \mapsto \sum_{j} u^{j} \odot f_{j, \operatorname{deg}(f)-j} .
\end{aligned}
$$

Exercise 9.6.2.3: (2) Show that $\Delta_{u}(f g)=\left(\Delta_{u} f\right) \odot\left(\Delta_{u} g\right)$, and that the generating series for the $E_{j}(f)$ may be written as

$$
\mathcal{E}_{f}(t)=\frac{1}{1 \otimes f} \odot \Delta_{t(1 \otimes f)} f .
$$

Note that $(1 \otimes f)^{\odot s}=1 \otimes f^{s}$ and $(1 \otimes f g)=(1 \otimes f) \odot(1 \otimes g)$. Thus

$$
\mathcal{E}_{f g}(t)=\left[\frac{1}{1 \otimes f} \odot \Delta_{[t(1 \otimes g)](1 \otimes f)}(f)\right] \odot\left[\frac{1}{1 \otimes g} \odot \Delta_{[t(1 \otimes f)](1 \otimes g)}(g)\right],
$$

and taking the logarithmic derivative (recalling Equation (6.1.5)) we conclude.

Remark 9.6.2.4. There was a gap in the argument in [Gor94], repeated in [GKZ94], when proving the "only if" part of the argument. They assumed that the zero set of $f$ contains a smooth point, i.e., that the differential of
$f$ is not identically zero. This gap was fixed in [Bri10]. In [GKZ94] they use $G_{0}\left(d, \operatorname{dim} V^{*}\right)$ to denote $C h_{d}\left(V^{*}\right)$.
9.6.3. Brill's equations as modules. Brill's equations are of degree $d+1$ on $S^{d} V^{*}$. (The total degree of $S_{d, d} V \otimes S^{d^{2}-d} V$ is $d(d+1)$ which is the total degree of $S^{d+1}\left(S^{d} V\right)$.) Consider the $G L(V)$-module map

$$
\begin{equation*}
S_{(d, d)} V \otimes S^{d^{2}-d} V \rightarrow S^{d+1}\left(S^{d} V\right) \tag{9.6.6}
\end{equation*}
$$

whose image consists of Brill's equations. The components of the target are not known in general and the set of modules present grows extremely fast. One can use the Pieri formula 8.1.3.1 to get the components of the domain. Using the Pieri formula, we conclude:
Proposition 9.6.3.1. As a $G L(V)$-module, Brill's equations are multiplicity free.
Exercise 9.6.3.2: (2) Write out the decomposition and show that only partitions with three parts appear as modules in Brill's equations. ©

Not all partitions with three parts appear:
Theorem 9.6.3.3. [Gua15a] As a $G L(V)$-module, Brill's equations are:

$$
\begin{array}{r}
S_{(732)} V \text { when } d=3, \text { and } \\
\bigoplus_{j=2}^{d} S_{\left(d^{2}-d, d, j\right)} V \text { when } d>3 .
\end{array}
$$

The proof is given by explicitly writing out highest weight vectors and determining their image under (9.6.6).

Remark 9.6.3.4. If $d<\mathbf{v}=\operatorname{dim} V^{*}$, then $C h_{d}\left(V^{*}\right) \subset S u b_{d}\left(S^{d} V^{*}\right)$ so $I\left(C h_{d}\left(V^{*}\right)\right) \supset \Lambda^{d+1} V \otimes \Lambda^{d+1}\left(S^{d-1} V\right)$. J. Weyman (in unpublished notes from 1994) observed that these equations are not in the ideal generated by Brill's equations. More precisely, the ideal generated by Brill's equations does not include modules $S_{\pi} V$ with $l(\pi)>3$ in degree $d+1$, so it does not cut out $C h_{d}\left(V^{*}\right)$ scheme theoretically when $d<\mathbf{v}$. By Theorem 9.1.2.11 the same conclusion holds for $C h_{5}\left(\mathbb{C}^{5}\right)$ and almost certainly holds for all $C h_{n}\left(\mathbb{C}^{n}\right)$ with $n \geq 5$.

### 9.7. Proofs of Proposition 9.5.2.4 and Theorem 9.1.2.14

### 9.7.1. Proof of Proposition 9.5.2.4.

Lemma 9.7.1.1. Let $X, Y$ be affine varieties equipped with polynomial $\mathbb{C}^{*}$ actions with unique fixed points $0_{X} \in X, 0_{Y} \in Y$, and let $f: X \rightarrow Y$ be a $\mathbb{C}^{*}$-equivariant morphism such that as sets, $f^{-1}\left(0_{Y}\right)=\left\{0_{X}\right\}$. Then $f$ is finite.

Assume Lemma 9.7.1.1 for the moment.
Proof of Proposition 9.5.2.4. Since $V^{\times n} / / \Gamma_{n}$ is normal and $\psi_{n}$ of (9.5.1) is regular and generically one to one, it just remains to show $\psi_{n}$ is finite.

Write $[0]=[0, \ldots, 0]$. To show finiteness, by Lemma 9.7.1.1, it is sufficient to show $\psi_{n}^{-1}(0)=[0]$ as a set, as $[0]$ is the unique $\mathbb{C}^{*}$ fixed point in $V^{\times n} / / \Gamma_{n}$, and every $\mathbb{C}^{*}$ orbit closure contains [0]. Now $u_{1} \cdots u_{n}=0$ if and only if some $u_{j}=0$, say $u_{1}=0$. The $T^{S L_{n}}$-orbit closure of $\left(0, u_{2}, \ldots, u_{n}\right)$ contains the origin so $\left[0, u_{2}, \ldots, u_{n}\right]=[0]$.

Sketch of Proof of Lemma 9.7.1.1. $\mathbb{C}[X], \mathbb{C}[Y]$ are $\mathbb{Z}_{\geq 0}$-graded, and the hypothesis $f^{-1}\left(0_{Y}\right)=\left\{0_{X}\right\}$ states that

$$
\mathbb{C}[X] /\left(f^{*}\left(\mathbb{C}[Y]_{>0}\right)\right)
$$

is a finite dimensional vector space. We want to show that $\mathbb{C}[X]$ is integral over $\mathbb{C}[Y]$. This follows from a graded version of Nakayama's Lemma (the algebraic implicit function theorem).

The condition $f^{-1}\left(0_{Y}\right)=\left\{0_{X}\right\}$ as sets in Lemma 9.7.1.1 says that the only maximal ideal of $\mathbb{C}[X]$ containing the ideal generated by $f^{*} \mathbb{C}[Y]_{>0}$ is $\mathbb{C}[X]_{>0}$.

If $R$ is a finitely generated ring with no nilpotents, the points of the associated affine variety are in one to one correspondence with the maximal ideals of $R$ and the prime ideals correspond to the irreducible subvarieties. For any ring $R$ let $\operatorname{Spec}(R)$ denote the set of prime ideals of $R$, called the affine scheme associated to $R$. See [Sha13b, §5.1] for an introduction.

Here are more details for the proof of Lemma 9.7.1.1 (see, e.g. [Kum13, Lemmas 3.1,3.2], or [Eis95, p136, Ex. 4.6a]):
Lemma 9.7.1.2. Let $R, S$ be $\mathbb{Z}_{\geq 0}$-graded, finitely generated domains over $\mathbb{C}$ such that $R_{0}=S_{0}=\mathbb{C}$, and let $f^{*}: R \rightarrow S$ be an injective graded algebra homomorphism. If $S_{>0}$ is the only maximal ideal of $S$ containing the ideal generated by $f^{*}\left(R_{>0}\right)$, then $S$ is a finitely generated $R$-module. In particular, it is integral over $R$.

Proof. Let $\mathfrak{m}$ be the ideal generated by $f^{*}\left(R_{>0}\right)$, so the radical of $\mathfrak{m}$ equals $S_{>0}$, and in particular $S_{>0}^{d}$ must be contained in it for all $d>d_{0}$, for some $d_{0}$. So $S / \mathfrak{m}$ is a finite dimensional vector space, and by the next lemma, $S$ is a finitely generated $R$-module.

Lemma 9.7.1.3. Let $S$ be as above, and let $M$ be a $\mathbb{Z}_{\geq 0}$-graded $S$-module. Assume $M /\left(S_{>0} \cdot M\right)$ is a finite dimensional vector space over $S / S_{>0} \simeq \mathbb{C}$. Then $M$ is a finitely generated $S$-module.

Proof. Choose a set of homogeneous generators $\left\{\bar{x}_{1}, \ldots, \bar{x}_{n}\right\} \subset M /\left(S_{>0} \cdot M\right)$ and let $x_{j} \in M$ be a homogeneous lift of $\bar{x}_{j}$. Let $N \subset M$ be the graded $S$-submodule $S x_{1}+\cdots+S x_{n}$. Then $M=S_{>0} M+N$, as let $a \in M$, consider $\bar{a} \in M /\left(S_{>0} M\right)$ and lift it to some $b \in N$, so $a-b \in S_{>0} M$, and $a=(a-b)+b$. Now quotient by $N$ to obtain

$$
\begin{equation*}
S_{>0} \cdot(M / N)=M / N . \tag{9.7.1}
\end{equation*}
$$

If $M / N \neq 0$, let $d_{0}$ be the smallest degree such that $(M / N)^{d_{0}} \neq 0$. But $S_{>0} \cdot(M / N)^{\geq d_{0}} \subset(M / N)^{\geq d_{0}+1}$ so there is no way to obtain $(M / N)^{d_{0}}$ on the right hand side. Contradiction.

### 9.7.2. Proof of the qualitative assertion in Theorem 9.1.2.14.

Theorem 9.7.2.1. [Bri93] For all $n \geq 1, \psi_{n}$ restricts to a map

$$
\begin{equation*}
\psi_{n}^{o}:\left(V^{* \times n} / / \Gamma_{n}\right) \backslash[0] \rightarrow S^{n} V^{*} \backslash 0 \tag{9.7.2}
\end{equation*}
$$

such that $\psi_{n}^{o *}: \mathbb{C}\left[S^{n} V^{*} \backslash 0\right] \rightarrow \mathbb{C}\left[\left(V^{* \times n} / / \Gamma_{n}\right) \backslash[0]\right]$ is surjective.
Corollary 9.7.2.2. [Bri93] The Hermite-Hadamard-Howe map

$$
h_{d, n}: S^{d}\left(S^{n} V\right) \rightarrow S^{n}\left(S^{d} V\right)
$$

is surjective for $d$ sufficiently large.
Proof of Corollary. Theorem 9.7.2.1 implies $h_{d, n}=\left(\psi_{n}^{*}\right)_{d}$ is surjective for $d$ sufficiently large, because the cokernel of $\psi_{n}^{*}$ is supported at a point and thus must vanish in large degree.

The proof of Theorem 9.7.2.1 will give a second proof that the kernel of $\psi_{n}^{*}$ equals the ideal of $C h_{n}\left(V^{*}\right)$.

Proof of Theorem 9.7.2.1. Since $\psi_{n}$ is $\mathbb{C}^{*}$-equivariant, we can consider the quotient map to projective space

$$
\underline{\psi}_{n}:\left(\left(V^{* \times n} / / \Gamma_{n}\right) \backslash[0]\right) / \mathbb{C}^{*} \rightarrow\left(S^{n} V^{*} \backslash 0\right) / \mathbb{C}^{*}=\mathbb{P} S^{n} V^{*}
$$

and show that $\underline{\psi}_{n}^{*}$ is surjective. Note that $\left(\left(V^{* \times n} / / \Gamma_{n}\right) \backslash[0]\right) / \mathbb{C}^{*}$ is $G L(V)$ isomorphic to $\left(\bar{P}^{n} V^{*}\right)^{\times n} / \mathfrak{S}_{n}$, as

$$
\left(V^{* \times n} / / \Gamma_{n}\right) \backslash[0]=\left(V^{*} \backslash 0\right)^{\times n} / \Gamma_{n}
$$

and $\Gamma_{n} \times \mathbb{C}^{*}=\left(\mathbb{C}^{*}\right)^{\times n} \rtimes \mathfrak{S}_{n}$. So we may write

$$
\underline{\psi}_{n}:\left(\mathbb{P} V^{*}\right)^{\times n} / \mathfrak{S}_{n} \rightarrow \mathbb{P} S^{n} V^{*} .
$$

It will be sufficient to show $\underline{\psi}_{n}^{*}$ is surjective on affine open subsets that cover the source and target. Let $w_{1}, \ldots, w_{\mathbf{v}}$ be a basis of $V^{*}$ and consider the affine open subset of $\mathbb{P} V^{*}$ given by elements where the coordinate on $w_{1}$ is nonzero, and the corresponding induced affine open subsets of $\left(\mathbb{P} V^{*}\right)^{\times n}$ and $\mathbb{P} S^{n} V^{*}$, call these $\left(\mathbb{P} V^{*}\right)_{1}^{\times n}$ and $\left(\mathbb{P} S^{n} V^{*}\right)_{1}$. I will show that the algebra
of $\mathfrak{S}_{n}$-invariant functions on $\left(\mathbb{P} V^{*}\right)_{1}^{\times n}$ is in the image of $\left(\mathbb{P} S^{n} V^{*}\right)_{1}$. The restriction of the quotient by $\mathfrak{S}_{n}$ of $\left(\mathbb{P} V^{*}\right)^{\times n}$ composed with $\underline{\psi}_{n}$ to these open subsets in coordinates is

$$
\left(\left(w_{1}+\sum_{s=2}^{\mathbf{v}} x_{s}^{1} w_{s}\right), \ldots,\left(w_{1}+\sum_{s=2}^{\mathbf{v}} x_{s}^{\mathbf{v}} w_{s}\right)\right) \mapsto \prod_{i=1}^{n}\left(w_{1}+\sum_{s=2}^{\mathbf{v}} x_{s}^{i} w_{s}\right) .
$$

The coefficients appearing on the right hand side are the elementary multisymmetric polynomials (also called the elementary symmetric vector polynomials). These generate the algebra of multi-symmetric polynomials, i.e., the algebra of $\mathfrak{S}_{n}$-invariant functions in the $n$ sets of variables $\left(x_{s}^{i}\right)_{i=1, \ldots, n}$. For a proof see [Wey97, §II.3] or [GKZ94, §4, Thm. 2.4]. The proof proceeds by first showing the power sum multi-symmetric polynomials generate the algebra and then showing one can express the power sum multi-symmetric polynomials in terms of the elementary ones.

Note that the proof also shows that $\left(\mathbb{P} V^{*}\right)^{\times n} / \mathfrak{S}_{n}$ is isomorphic to $C h_{n}\left(V^{*}\right)$ as a projective variety, which is also shown in [GKZ94, §4 Thm. 2.2].

For any orbit closure $\overline{G \cdot v}$ in affine space (where $v \in V$ and $V$ is a $G$ module), we always have an inclusion $\mathbb{C}[\operatorname{Nor}(\overline{G \cdot v})] \subseteq \mathbb{C}[G \cdot v]_{\text {poly }}$ because $G \cdot v$ is also a Zariski open subset of the normalization, as it is contained in the smooth points of $\overline{G \cdot v}$. In our situation $\mathbb{C}\left[\operatorname{Nor}\left(\hat{C} h_{n}\left(V^{*}\right)\right]=\mathbb{C}[G L(V)\right.$. $\left.\left(x_{1} \cdots x_{n}\right)\right]_{\text {poly }}$. This gives a second proof that $\psi_{n}^{*}=h_{n}$.

Remark 9.7.2.3. Given the above situation, I had asked: Under what circumstances is $\mathbb{C}[\operatorname{Nor}(\overline{G L(V) \cdot w})]=\mathbb{C}[G L(V) \cdot w]_{\text {poly }}$ when $W$ is a polynomial $G L(V)$-module, $w \in W$, and $G_{w}$ is reductive? In [H0] Hüttenhain answers this question: A simple necessary condition is that $\operatorname{End}(V) \cdot w=$ $\overline{G L(V) \cdot w}$. A necessary and sufficient condition is that for all $X \in \operatorname{End}(W)$ such that $w \in \operatorname{ker}(X)$, the zero endomorphism lies in $\overline{G_{w} \cdot X} \subset \operatorname{End}(W)$. A simple example where this fails is the cubic $w=x^{3}+y^{3}+z^{3} \in S^{3} \mathbb{C}^{3}$, whose orbit closure is the degree four hypersurface $\sigma_{3}\left(v_{3}\left(\mathbb{P}^{2}\right)\right)$. By this result and the boundary component of $\S 6.7 .1$, we know that the coordinate ring of the normalization of $\mathcal{D} e t_{n}$ and the polynomial part of the coordinate ring of the orbit do not agree.
9.7.3. Proof of Brion's quantitative theorem 9.1.2.14. We have a ring map

$$
h_{n}: \operatorname{Sym}\left(S^{n} V\right) \rightarrow \bigoplus_{i} S^{n}\left(S^{i} V\right)
$$

The proof has three steps:
(1) Show $\mathbb{C}\left[\operatorname{Nor}\left(C h_{n}\left(V^{*}\right)\right)\right]$ is generated in degree at most $(n-1)(\mathbf{v}-1)$ via vanishing of cohomology (Castelnuovo-Mumford regularity, see, e.g., $[$ Eis05, Chap. 4]).
(2) Show that $h_{n}\left(\left(v^{n}\right)^{d(n-1)} \cdot \mathbb{C}\left[\operatorname{Nor}\left(C h_{n}\left(V^{*}\right)\right)\right] \subset \mathbb{C}\left[C h_{n}\left(V^{*}\right)\right]\right.$ via a localization argument to reduce to a question about multi-symmetric functions.
(3) Use that Zariski open subset of the polynomials of degree $n$ in $\mathbf{v}$ variables can be written as a sum of $r_{0} n$-th powers, where $r_{0} \sim \frac{1}{n}\binom{\mathbf{v}+n-1}{n}$ (The Alexander-Hirschowitz theorem [AH95], see [BO08] for a detailed exposition of the proof or [Lan12, Chap. 15]).
Then we conclude that for $d \geq(n-1)(\mathbf{v}-1)\left(r_{0}(n-1)+n\right)$ that $h_{d, n}$ is surjective.

Proof of Step 1. Recall that $\mathbb{C}\left[\operatorname{Nor}\left(C h_{n}\left(V^{*}\right)\right)\right]=\left(\mathbb{C}\left[V^{* \times n}\right]^{T^{S L_{n}}}\right)^{\mathfrak{S}_{n}}$ so it will be sufficient to show that $\mathbb{C}\left[V^{* \times n}\right]^{T^{S L_{n}}}$ is generated in degree at most $(n-1)(\mathbf{v}-1)$. This translates into a sheaf cohomology problem:

$$
\begin{aligned}
\mathbb{C}\left[V^{* \times n}\right]^{T_{n}} & =\bigoplus_{d=0}^{\infty} H^{0}\left(\mathbb{P} V^{* \times n}, \mathcal{O}_{\mathbb{P} V^{*}}(d)^{\times n}\right) \\
& =\bigoplus_{d=0}^{\infty} H^{0}\left(\mathbb{P} S^{n} V^{*}, \operatorname{proj}_{*} \mathcal{O}_{\mathbb{P}^{*}}(d)^{\times n}\right),
\end{aligned}
$$

where proj : $\mathbb{P} V^{* \times n} \rightarrow \mathbb{P}\left(S^{n} V^{*}\right)$ is the projection map. We want an upper bound on the degrees of the generators of the graded $\operatorname{Sym}\left(S^{n} V\right)$ module associated to the sheaf $\operatorname{proj}_{*} \mathcal{O}_{\mathbb{P V}^{*} *}^{\times n}$. Castelnuovo-Mumford regularity [Mum66, Lect. 14] gives a bound in terms of vanishing of sheaf cohomology groups. Here we are dealing with groups we can compute: $H^{j}\left(\mathbb{P} V^{* \times n}, \mathcal{O}(d-j)^{\times n}\right)$, and the result follows from this computation.

Proof of Step 2. Let $v=v_{\mathbf{v}} \in V \backslash 0$, and let $v_{1}, \ldots, v_{\mathbf{v}}$ be a basis of $V$, which may also be considered as linear forms on $V^{*}$, so $x_{i}:=\frac{v_{i}}{v}$ makes sense. Consider the localization of the coordinate ring of the normalization at $v^{n}$, the degree zero elements in the localization of $\mathbb{C}\left[\operatorname{Nor}\left(C h_{n}\left(V^{*}\right)\right)\right]\left[\frac{1}{v^{n}}\right]$ :

$$
\begin{aligned}
\mathbb{C}\left[\operatorname{Nor}\left(C h_{n}\left(V^{*}\right)\right)\right]_{v^{n}}: & =\bigcup_{d \geq 0} S^{n}\left(S^{d} V\right)\left(v^{n}\right)^{-d} \\
& =S^{n}\left(\bigcup_{d \geq 0}\left(S^{d} V\right)\left(v^{n}\right)^{-d}\right. \\
& =S^{n} \mathbb{C}\left[x_{1}, \ldots, x_{\mathbf{v}-1}\right]=: S^{n} \mathbb{C}[\bar{x}] \\
& =\left[(\mathbb{C}[\bar{x}])^{\otimes n}\right]^{\mathfrak{G}_{n}} \\
& =\left(\mathbb{C}\left[\overline{x_{1}}, \ldots, \overline{x_{n}}\right]\right)^{\mathfrak{G}_{n}},
\end{aligned}
$$

where $\overline{x_{j}}=\left(x_{1, j}, \ldots, x_{\mathbf{v}-1, j}\right)$.
Similarly

$$
\begin{aligned}
\operatorname{Sym}\left(S^{n} V\right)_{v^{n}} & =\bigcup_{d \geq 0} S^{d}\left(S^{n} V\right)\left(v^{n}\right)^{-d} \\
& =\operatorname{Sym}\left(S^{n} V / v^{n}\right) \\
& =\operatorname{Sym}\left(\bigoplus_{i=1}^{n} \mathbb{C}[\bar{x}]_{i}\right) .
\end{aligned}
$$

We get a localized graded algebra map $h_{n, v^{n}}$ between these spaces. Hence it is determined in degree one, where the map

$$
\bigoplus_{i=1}^{n} \mathbb{C}[\bar{x}]_{i} \rightarrow \mathbb{C}\left[\overline{x_{1}}, \ldots, \overline{x_{n}}\right]^{\mathfrak{S}_{n}}
$$

takes the degree at most $n$ monomial $x_{1}^{a_{1}} \cdots x_{d-1}^{a_{d_{1}}}$ to the coefficient of $t_{1}^{a_{1}} \cdots t_{d-1}^{a_{d-1}}$ in the expansion of

$$
\Pi_{i=1}^{n}\left(1+\bar{x}_{i} t_{1}+\cdots+\bar{x}_{i d-1} t_{d-1}\right)
$$

Again we obtain elementary multi-symmetric functions which generate the ring of multi-symmetric functions $\mathbb{C}\left[\overline{x_{1}}, \ldots, \overline{x_{n}}\right]^{\mathfrak{G}_{n}}$. Thus $h_{n, v^{n}}$ is surjective.

Moreover, if $f \in \mathbb{C}\left[\overline{x_{1}}, \ldots, \overline{x_{n}}\right]^{\mathfrak{G}_{n}}$ has all its partial degrees at most $d$, then the total degree of $f$ is at most $d n$ in the $\overline{x_{j}}$ 's, so it is a polynomial of degree at most $d n$ in the elementary multi-symmetric functions. In other words, the map

$$
S^{d n}\left(S^{n} V\right)\left(v^{n}\right)^{-d n} \rightarrow S^{n}\left(S^{d} V\right)\left(v^{n}\right)^{-d}
$$

is surjective, so $h_{n}\left(\left(v^{n}\right)^{d(n-1)} \mathbb{C}\left[\operatorname{Nor}\left(C h_{n}(V)\right)\right] \subset \mathbb{C}\left[C h_{n}(V)\right]\right.$.
We conclude by appeal to the Alexander-Hirschowitz theorem [AH95].

## Chapter 10

## Topics using additional algebraic geometry

This chapter covers four mostly independent topics: $\S 10.1$ presents symmetric (Waring) rank lower bounds for polynomials, $\S 10.2$ explains limits of determinantal methods (such as the method of partial derivatives, of shifted partial derivatives and Koszul flattenings) for proving lower complexity bounds for tensors and polynomials, $\S 10.3$ shows that the singularities of the varieties $\mathcal{D} e t_{n}$ and $\mathcal{P} e r m_{n}^{m}$ make their study more difficult (they are not normal varieties), and $\S 10.4$ discusses further commutative algebra results that might be useful in future study of Valiant's hypothesis (syzygies of Jacobian loci of $\operatorname{det}_{n}$ and $\operatorname{perm}_{m}$ ). Other than $\S 10.2$, they can be read independently, and $\S 10.2$ only requires $\S 10.1$.

In $\S 10.1$, I introduce the language of zero dimensional schemes to state and prove the Apolarity Lemma, an important tool for proving symmetric rank lower bounds. This section does not assume any algebraic geometry beyond what was discussed in Chapters 1-9. It will hopefully motivate computer scientists to learn about zero dimensional schemes, and show algebraic geometers interesting applications of the subject to complexity. I introduce the cactus variety, as the apolarity method generally also proves lower bounds on the cactus rank. In $\S 10.2$, the cactus variety is shown to be a major obstruction to proving superlinear border rank lower bounds for tensors: all known equations for lower bounds also give lower bounds for the cactus variety, but tensors in $\mathbb{C}^{\mathbf{m}} \otimes \mathbb{C}^{\mathbf{m}} \otimes \mathbb{C}^{\mathbf{m}}$ never have cactus border rank above $6 \mathbf{m}$.

In $\S 10.3$ I present Kumar's proof that $\mathcal{D e t}_{n}$ is not a normal variety. This section also does not assume anything beyond what was introduced in Chapters 1-9.

I conclude, in §10.4, with Lascoux's derivation [Las78], as made explicit by Weyman [Wey03], of the minimal free resolution of the ideal generated by the size $k$ minors of a matrix of indeterminants, and briefly compare it with the ideal generated by subpermanents. The exposition indicates how tools from commutative algebra might be helpful in proving lower bounds on $\overline{\mathrm{dc}}\left(\operatorname{perm}_{m}\right)$. Parts of this section assume additional background in algebraic geometry and representation theory.

### 10.1. Rank and cactus rank of polynomials

This section and the next deal with two complexity issues using zero dimensional schemes: lower bounds on the symmetric rank of polynomials, and the limits of determinantal methods for proving border rank lower bounds. In this section I discuss lower bounds on rank. I begin, in $\S 10.1 .1$ by introducing the language of affine and projective schemes and defining cactus varieties, a generalization of secant varieties. In §10.1.2 I introduce apolarity as a tool for proving symmetric rank (and cactus rank) lower bounds. The key to using apolarity is Bezout's theorem, a classical theorem in algebraic geometry. I state and prove a version sufficient for our purposes in in $\S 10.1 .3$. With this, the Ranestad-Schreyer results [RS11] on the ranks of monomials follow, in particular, that $\mathbf{R}_{S}\left(x_{1} \cdots x_{n}\right)=2^{n-1}$. This is presented in §10.1.4. Bezout's theorem similarly enables Lee's lower bounds [Lee16] on the symmetric ranks of elementary symmetric functions, which are tight in odd degree and presented in §10.1.5.

To facilitate the distinction between elements of $\operatorname{Sym}(V)$ and $\operatorname{Sym}\left(V^{*}\right)$, I will use lower case letters for elements of $\operatorname{Sym}(V)$ and either upper case letters or differential operator notation for elements of $\operatorname{Sym}\left(V^{*}\right)$, e.g., $f \in$ $S^{d} V, P \in S^{e} V^{*}, \frac{\partial}{\partial x_{n}} \in V^{*}$.
10.1.1. Language from algebraic geometry. As mentioned in §9.7, affine varieties correspond to finitely generated algebras over $\mathbb{C}$ with no nilpotent elements and affine schemes similarly correspond to finitely generated rings. Given a finitely generated algebra, the corresponding scheme is the set of its prime ideals (endowed with a topology that generalizes the Zariski topology for varieties). This enables us to "remember" non-reduced structures. For example, $\mathbb{C}[x] /(x)^{2}$ defines a scheme which we think of as the origin in $\mathbb{A}^{1}$ with multiplicity two. An affine scheme $Z$ is zero dimensional if the corresponding ring (called the ring of regular functions on $Z$ ) is a finite dimensional vector space over $\mathbb{C}$. If a variety is a collection of $d$ distinct points,
its corresponding coordinate ring is a $d$-dimensional vector space. The length or degree of a zero dimensional affine scheme is the dimension of the corresponding ring as a $\mathbb{C}$-vector space. If $Z$ denotes a zero-dimensional affine scheme, where the corresponding ring is a quotient ring of $\mathbb{C}\left[x_{1}, \ldots, x_{n}\right]$, i.e., $\mathbb{C}[Z]=\mathbb{C}\left[x_{1}, \ldots, x_{n}\right] / \mathcal{I}$, we consider $Z$ as a subscheme of the affine space $\mathbb{A}^{n}$.

A homogeneous ideal $I \subset \operatorname{Sym}\left(V^{*}\right)$ is saturated if every $P \in \operatorname{Sym}\left(V^{*}\right)$ such that for all $L \in V^{*}$, there exists a $\delta$ such that $P L^{\delta} \in I$ satisfies $P \in I$. A projective scheme corresponds to a graded ring $S=\operatorname{Sym}\left(V^{*}\right) / I$ with $I$ a saturated homogeneous ideal (not equal to $\operatorname{Sym}\left(V^{*}\right)_{>0}$ ) and the associated projective scheme is the set of homogeneous prime ideals of $S$, excluding $\oplus_{d>0} S_{d}$, see, e.g., [Har77, §II.2]. The corresponding scheme is denoted $\operatorname{Proj}(S)$. (I continue to use the notation $\mathrm{Zeros}(I)$ for the zero set of an ideal $I$.) One says $X=\operatorname{Proj}\left(\operatorname{Sym}\left(V^{*}\right) / I\right)$ is a subscheme of $\mathbb{P} V$ and writes $X \subset \mathbb{P} V$ and $I$ is called the ideal of $X$. More generally, for schemes $X, Y \subset \mathbb{P} V$ defined by ideals $I(X)$ and $I(Y), X \subset Y$ means $I(Y) \subset I(X)$. The support of a scheme determined by an ideal $I$ is $\operatorname{Zeros}(I)$. Define the span of $X,\langle X\rangle$ to be the linear space $\operatorname{Zeros}\left(I(X)_{1}\right) \subset \mathbb{P} V$.

Definition 10.1.1.1. [BB14] Let $X \subset \mathbb{P} V$ be a projective variety and let $y \in \mathbb{P} V$. Define the $X$-cactus rank of $y$ to be the smallest $r$ such that there exists a zero dimensional scheme $Z \subset X$ of length $r$ such that $y \in\langle Z\rangle$. Write $\mathfrak{c r}_{X}(y)=r$. (The usual case of $X$-rank is when $Z$ consists of $r$ distinct points.) Define the $r$-th cactus variety of $X$ to be

$$
\mathfrak{k}_{r}(X):=\overline{\left\{y \in \mathbb{P} V \mid \mathfrak{c r}_{X}(y) \leq r\right\}}
$$

and define the cactus border rank of $y,{\operatorname{cr}_{X}}_{X}(y)$ to be the smallest $r$ such that $y \in \mathfrak{k}_{r}(X)$.

By definition $\mathfrak{c r}_{X}(y) \leq \mathbf{R}_{X}(y)$ and $\underline{\mathfrak{c r}}_{X}(y) \leq \underline{\mathbf{R}}_{X}(y)$, i.e., $\sigma_{r}(X) \subseteq \mathfrak{k}_{r}(X)$, and strict inequality can occur in both cases. The cactus rank was originally defined for polynomials in [Iar95], where it was called scheme length, and in general in [BB14]. We will be mostly concerned with the cactus rank of polynomials with respect to $X=v_{d}(\mathbb{P} V)$. For $f \in S^{d} V$, write $\mathfrak{c r}_{S}(f):=$ $\mathfrak{c r}_{v_{d}(\mathbb{P} V)}(f)$. If $Y \subset v_{d}(\mathbb{P} V)$, then since the Veronese map is an embedding, there exists a subscheme $Z \subset \mathbb{P} V$ such that $Y=v_{d}(Z)$. It will be more convenient to write $v_{d}(Z)$ for a subscheme of $v_{d}(\mathbb{P V})$ in what follows.
Exercise 10.1.1.2: (1!) Show that for any subscheme $Z \subset \mathbb{P} V,(I(Z))_{e}=$ $\left\langle v_{e}(Z)\right\rangle^{\perp} \subset S^{e} V^{*}$, where $\perp$ denotes the annhilator in the dual space. In particular, $I\left(v_{d}(Z)\right)_{1}=I(Z)_{d . \odot}$
10.1.2. The apolar ideal. For $f \in S^{d} V$, recall from $\S 6.2$ the flattening (catalecticant) maps $f_{j, d-j}: S^{j} V^{*} \rightarrow S^{d-j} V$ given by $D \mapsto D(f)$. Define
the annihilator of $f$ or apolar ideal of $f, f^{a n n} \subset \operatorname{Sym}\left(V^{*}\right)$ by

$$
\begin{aligned}
f^{a n n} & :=\bigoplus_{j=1}^{d} \operatorname{ker} f_{j, d-j} \oplus \bigoplus_{k=d+1}^{\infty} S^{k} V^{*} \\
& =\left\{P \in \operatorname{Sym}\left(V^{*}\right) \mid P(f)=0\right\}
\end{aligned}
$$

Given a (not necessarily homogeneous) polynomial $f \in \mathbb{C}\left[x_{1}, \ldots, x_{n}\right]$ define $f^{a n n}:=\left\{\left.P \in \mathbb{C}\left[\frac{\partial}{\partial x_{1}}, \ldots, \frac{\partial}{\partial x_{n}}\right] \right\rvert\, P(f)=0\right\}$.

Recall the notation $\operatorname{Partials}(f):=\left\{P(f) \mid P \in \operatorname{Sym}\left(V^{*}\right)\right\} \subset \operatorname{Sym}(V)$.
Exercise 10.1.2.1: (1) Show that if $e<d$, and $P \in S^{e} V^{*}$ satisfies $P(f) \neq 0$, then there exists $L \in V^{*}$ such that $(L P)(f) \neq 0$. Show more generally that for any $\delta \leq d-e$ there exists $Q \in S^{\delta} V^{*}$ with $Q P(f) \neq 0$.

The following lemma is critical:
Lemma 10.1.2.2. [Apolarity Lemma] A finite subscheme $Z \subset \mathbb{P} V$ satisfies $f \in\left\langle v_{d}(Z)\right\rangle=\operatorname{Zeros}\left(I(Z)_{d}\right)$ if and only if $I(Z) \subseteq f^{a n n}$. In particular, $f \in \operatorname{span}\left\{\ell_{1}^{d}, \ldots, \ell_{r}^{d}\right\} \subset S^{d} V$ if and only if $f^{a n n} \supseteq I\left(\left\{\left[\ell_{1}\right], \ldots,\left[\ell_{r}\right]\right\}\right)$.

If $Z$ satisfies the conditions of the Apolarity Lemma, we say is apolar to $f \in S^{d} V$.

Proof. The "if" is clear. To prove the other direction, consider $I(Z)_{e}$. When $e>d, f_{e}^{a n n}=S^{e} V^{*}$. When $e=d$, it is the hypothesis. For $e<d$, let $P \in I(Z)_{e}$. Then $S^{d-e} V^{*} \cdot P \subset I(Z)_{d} \subset\left(f^{a n n}\right)_{d}$, so $(Q P)(f)=0$ for all $Q \in S^{d-e} V^{*}$, which implies $P(f)=0$ by Exercise 10.1.2.1, which is what we wanted to prove.
10.1.3. Bezout's theorem. Given a graded ring $R$, define its Hilbert function $\operatorname{HilbF}_{k}(R):=\operatorname{dim}\left(R_{k}\right)$. If $R=\operatorname{Sym}\left(V^{*}\right) / I$ for some homogeneous ideal $I \subset \operatorname{Sym}\left(V^{*}\right)$, then there exists a polynomial, called the Hilbert polynomial of $R, \operatorname{HilbP}_{z}(R)$ in a variable $z$ such that for $k$ sufficiently large $\operatorname{HilbF}_{k}(R)=\operatorname{HilbP}_{k}(R)$, see, e.g., [Sha13b, §6.4.2]. For a projective scheme $X$ defined by an ideal $I \subset \operatorname{Sym}\left(V^{*}\right)$, one may define its dimension as $\operatorname{dim}(X):=\operatorname{deg}\left(\operatorname{HilbP}_{z}\left(\operatorname{Sym}\left(V^{*}\right) / I\right)\right)$, which agrees with our previous definition for varieties. One writes $\operatorname{HilbF}_{k}(X)$ and $\operatorname{HilbP}_{z}(X)$ for the Hilbert functions and polynomials of the coordinate ring of $X$.
Exercise 10.1.3.1: (1) Let $X=\mathbb{P}^{\ell-1} \subseteq \mathbb{P} V$ be a linear space. Show $\operatorname{HilbP}_{z}(X)=\frac{1}{(\ell-1)!} z^{\ell-1}+O\left(z^{\ell-2}\right)$.
Exercise 10.1.3.2: (1) Let $\operatorname{Zeros}(P) \subseteq \mathbb{P} V$ be a hypersurface of degree $d$. Show $\operatorname{HilbP}_{z}(\operatorname{Zeros}(P))=\frac{d}{(\mathbf{v}-2)!} z^{\mathbf{v}-3}+O\left(z^{\mathbf{v}-4}\right)$.

The above exercise suggests that one may define $\operatorname{deg}(X)$, the degree of a projective scheme $X$, to be the coefficient of $z^{\operatorname{dim}(X)}$ in $\operatorname{HilbP}_{z}(X)$
times $\operatorname{dim}(X)$ !, which indeed agrees with our previous definition of degree for a projective variety (see, e.g., [Har95, Lect. 14]) and degree of a zero dimensional scheme (in this case the Hilbert polynomial is constant).

Definition 10.1.3.3. Given an ideal $I \subset \operatorname{Sym}\left(V^{*}\right)$ and $G \in S^{e} V^{*}$, say $G$ is transverse to $I$ if the multiplication map

$$
\begin{equation*}
\operatorname{Sym}\left(V^{*}\right) / I \xrightarrow{G} \operatorname{Sym}\left(V^{*}\right) / I \tag{10.1.1}
\end{equation*}
$$

given by $P \mapsto G P$ is injective.
We will need a corollary of the following classical theorem:
Theorem 10.1.3.4 (Bezout's Theorem). Let $I \subset \operatorname{Sym}\left(V^{*}\right)$ be a homogeneous ideal and let $G \in S^{e} V^{*}$ be transverse to $I$. Write $\operatorname{HilbP}_{z}\left(\operatorname{Sym}\left(V^{*}\right) / I\right)=$ $\frac{d}{n!} z^{n}+O\left(z^{n-1}\right)$ and assume $n \geq 1$. Then $\operatorname{HilbP}_{z}\left(\operatorname{Sym}\left(V^{*}\right) /(I+(G))\right)=$ $\frac{e d}{(n-1)!} z^{n-1}+O\left(z^{n-2}\right)$.

Proof. Consider the exact sequence

$$
0 \rightarrow \operatorname{Sym}\left(V^{*}\right) / I \xrightarrow{G} \operatorname{Sym}\left(V^{*}\right) / I \rightarrow \operatorname{Sym}\left(V^{*}\right) /(I+(G)) \rightarrow 0 .
$$

This is an exact sequence of graded $\operatorname{Sym}\left(V^{*}\right)$-modules, i.e., for all $\delta$, the sequence

$$
0 \rightarrow S^{\delta-e} V^{*} / I_{\delta-e} \xrightarrow{G} S^{\delta} V^{*} / I_{\delta} \rightarrow S^{\delta} V^{*} /(I+(G))_{\delta} \rightarrow 0
$$

is exact, so

$$
\operatorname{dim}\left(S^{\delta} V^{*} /(I+(G))_{\delta}\right)=\operatorname{dim}\left(S^{\delta}\left(V^{*}\right) / I_{\delta}\right)-\operatorname{dim}\left(S^{\delta-e}\left(V^{*}\right) / I_{\delta-e}\right)
$$

Thus
$\operatorname{HilbP}_{z}\left(\operatorname{Sym}\left(V^{*}\right) /(I+(G))\right)=\operatorname{HilbP}_{z}\left(\operatorname{Sym}\left(V^{*}\right) / I\right)-\operatorname{HilbP}_{z-e}\left(\operatorname{Sym}\left(V^{*}\right) / I\right)$.
Write $\operatorname{HilbP}_{z}\left(\operatorname{Sym}\left(V^{*}\right) / I\right)=\frac{d}{n!} z^{n}+c z^{n-1}+O\left(z^{n-2}\right)$, so

$$
\operatorname{HilbP}_{z-e}\left(\operatorname{Sym}\left(V^{*}\right) / I\right)=\frac{d}{n!}(z-e)^{n}+c(z-e)^{n-1}+O\left(z^{n-2}\right),
$$

and

$$
\operatorname{HilbP}_{z}\left(\operatorname{Sym}\left(V^{*}\right) /(I+(G))\right)=\frac{e d}{(n-1)!} z^{n-1}+O\left(z^{n-2}\right)
$$

Corollary 10.1.3.5. Let $I \subset \operatorname{Sym}\left(V^{*}\right)$ be a homogeneous ideal defining a zero dimensional scheme in $\mathbb{P} V$, and let $G \in S^{e} V^{*}$ be transverse to $I$. Then $\operatorname{dim}\left(\operatorname{Sym}\left(V^{*}\right) /(I(Z)+(G))\right)=\delta \operatorname{deg}(Z)$,

Proof. First consider the 1-dimensional affine scheme defined by $I(Z)$ in the affine space $V$. To return to the projective setting (so that one can apply Theorem 10.1.3.4), consider $V \subset \mathbb{P}^{\mathbf{v}}$ as an affine open subset, e.g., the subset $\left[1, x_{1}, \ldots, x_{\mathbf{v}}\right]$. Then $I(Z)+(G)$ cuts out a zero dimensional subscheme of $\mathbb{P}^{\mathbf{v}}$ supported at the point $[1,0, \ldots, 0]$, and we can use Bezout's theorem on $\mathbb{P}^{\mathbf{v}}$ to conclude $\operatorname{dim}\left(S^{D}\left(\mathbb{C}^{\mathbf{v}+1 *}\right) /(I(Z)+(G))_{D}\right)=\delta \operatorname{deg}(Z)$ for $D$ sufficiently large (Here the Hilbert polynomial is a constant.) But this implies $\operatorname{dim}\left(\operatorname{Sym}\left(V^{*}\right) /(I(Z)+(G))\right)=\delta \operatorname{deg}(Z)$.

### 10.1.4. Lower bounds on cactus rank.

Theorem 10.1.4.1. [RS11] Let $f \in S^{d} V$ and say $f^{a n n}$ is generated in degrees at most $\delta$. Then

$$
\mathfrak{c r}_{S}(f) \geq \frac{1}{\delta} \operatorname{dim}\left(S y m\left(V^{*}\right) / f^{a n n}\right)=\frac{1}{\delta} \sum_{j=0}^{d} \operatorname{rank} f_{j, d-j}=\frac{1}{\delta} \operatorname{dim}(\operatorname{Partials}(f)) .
$$

Proof. Let $f \in S^{d} V$ and let $Z \subset \mathbb{P} V$ be a zero dimensional projective scheme of minimal degree satisfying $I(Z) \subseteq f^{a n n}$ (i.e., $f \in\left\langle v_{d}(Z)\right\rangle$ ), so $\mathfrak{c r}_{S}(f)=\operatorname{deg}(Z)$. There exists $G \in\left(f^{a n n}\right)_{\delta}$ such that $G$ is transverse to $I(Z)$ because in degrees $e>d,\left(f^{a n n}\right)_{e}=S^{e} V^{*}$, so there must be some polynomial in the ideal transverse to $I(Z)$, and since the ideal is generated in degrees at most $\delta$, there must be such a polynomial in degree $\delta$. By Corollary 10.1.3.5, $\operatorname{dim}\left(\operatorname{Sym}\left(V^{*}\right) /(I(Z)+(G))\right)=\delta \operatorname{deg}(Z)$.

Finally, $\operatorname{dim}\left(\operatorname{Sym}\left(V^{*}\right) /(I(Z)+(G))\right) \geq \operatorname{dim}\left(\operatorname{Sym}\left(V^{*}\right) / f^{a n n}\right)$.
For example, if $f=x y z$, then $\operatorname{rank}\left(f_{3,0}\right)=1, \operatorname{rank}\left(f_{2,1}\right)=3, \operatorname{rank}\left(f_{1,2}\right)=$ $3, \operatorname{rank}\left(f_{0,3}\right)=1$, and $f^{a n n}$ is generated in degree two, so $\mathfrak{c r}_{S}(x y z) \geq 4$. On the other hand $f^{a n n}$ is generated by $\frac{\partial^{2}}{(\partial x)^{2}}, \frac{\partial^{2}}{(\partial y)^{2}}, \frac{\partial^{2}}{(\partial z)^{2}}$, and the scheme with the ideal generated by any two of these has length four, so equality holds.
Exercise 10.1.4.2: (1) Prove that $\mathbf{R}_{S}\left(x_{1} \cdots x_{n}\right)=\mathfrak{c r}_{S}\left(x_{1} \cdots x_{n}\right)=2^{n-1}$, which was first shown in [RS11].
Exercise 10.1.4.3: (2) Prove that more generally, for a monomial $x_{1}^{d_{1}} \cdots x_{n}^{d_{n}}$ with $d_{1} \leq d_{2} \leq \cdots \leq d_{n}$, then $\mathbf{R}_{S}\left(x_{1}^{d_{1}} \cdots x_{n}^{d_{n}}\right) \geq \mathfrak{c r}_{v_{d}(\mathbb{P} V)}\left(x_{1}^{d_{1}} \cdots x_{n}^{d_{n}}\right)=$ $\left(d_{1}+1\right) \cdots\left(d_{n-1}+1\right)$. This was also first shown in [RS11]. ©
Remark 10.1.4.4. L. Oeding [Oed16], using Young flattenings, has shown $\underline{\mathbf{R}}_{S}\left(x_{1} \cdots x_{n}\right)=2^{n-1}$, and using a generalization of them, for $d_{1} \leq d_{2} \leq$ $\cdots \leq d_{n}$, that $\underline{\mathbf{R}}_{S}\left(x_{1}^{d_{1}} \cdots x_{n}^{d_{n}}\right)=\left(d_{1}+1\right) \cdots\left(d_{n-1}+1\right)$.
10.1.5. Waring Ranks of elementary symmetric functions. For ideals $I, J \subset \operatorname{Sym}\left(V^{*}\right)$, introduce the colon ideal

$$
I: J:=\left\{P \in \operatorname{Sym}\left(V^{*}\right) \mid P J \subseteq I\right\} .
$$

Exercise 10.1.5.1: (1) Prove that if $I, J \subset \operatorname{Sym}\left(V^{*}\right)$ are saturated homogeneous ideals such that $\operatorname{Sym}\left(V^{*}\right) / I$ and $\operatorname{Sym}\left(V^{*}\right) / J$ contain no nilpotent elements, then $I: J$ is the ideal of polynomials vanishing on the components of $\operatorname{Zeros}(I)$ that are not contained in $\operatorname{Zeros}(J)$.

Exercise 10.1.5.2: (1) Show that for $D \in \operatorname{Sym}\left(V^{*}\right), f^{a n n}: D=(D(f))^{a n n}$.
Theorem 10.1.5.3. $\left[\mathbf{C C C}^{+} \mathbf{1 5 a}\right]$ For $f \in S^{d} V$, and sufficiently general $P \in S^{e} V^{*}$,

$$
\mathbf{R}_{S}(f) \geq \frac{1}{e} \operatorname{dim}\left(\frac{\operatorname{Sym}\left(V^{*}\right)}{f^{\text {ann }}:(P)+(P)}\right)
$$

Proof. Say $f=\ell_{1}^{d}+\cdots+\ell_{r}^{d}$ is a minimal rank decomposition of $f$, so by the Apolarity Lemma 10.1.2.2 $f^{a n n} \supseteq I\left(\left\{\left[\ell_{1}\right], \ldots,\left[\ell_{r}\right]\right\}\right)$. Take $P$ such that for $1 \leq j \leq r, P\left(\ell_{j}\right) \neq 0$. (Note that a Zariski open subset of $S^{e} V^{*}$ has this property, explaining the "sufficiently general" condition.) Let $J=$ $I\left(\left\{\left[\ell_{1}\right], \ldots,\left[\ell_{r}\right]\right\}\right):(P)$, and let $\operatorname{Proj}(J)$ be the corresponding projective scheme. Our choice of $P$ insures that the multiplication map $\operatorname{Sym}\left(V^{*}\right) / J \xrightarrow{\cdot P}$ $\operatorname{Sym}\left(V^{*}\right) / J$ is injective. To see this, say $H \notin J$ and $H \cdot P \in J$. Then $H \cdot P^{2} \in I\left(\left\{\left[\ell_{1}\right], \ldots,\left[\ell_{r}\right]\right\}\right)$, but $I\left(\left\{\left[\ell_{1}\right], \ldots,\left[\ell_{r}\right]\right\}\right)$ is reduced so $H \cdot P \in$ $I\left(\left\{\left[\ell_{1}\right], \ldots,\left[\ell_{r}\right]\right\}\right)$ which means $H \in J$. Corollary 10.1.3.5 applies to show $e \cdot \operatorname{deg}(\operatorname{Proj}(J))=\operatorname{dim}\left(\operatorname{Sym}\left(V^{*}\right) /(J+(P))\right)$.

The genericity condition also implies $\operatorname{deg}(\operatorname{Proj}(J))=\operatorname{deg}\left(\left\{\left[\ell_{1}\right], \ldots,\left[\ell_{r}\right]\right\}\right)=$ $r=\mathbf{R}_{S}(f)$. Furthermore, since $J \subseteq f^{a n n}:(P)$,

$$
\operatorname{dim}\left(\frac{\operatorname{Sym}\left(V^{*}\right)}{f^{\text {ann }}:(P)+(P)}\right) \leq \operatorname{dim}\left(\frac{\operatorname{Sym}\left(V^{*}\right)}{J+(P)}\right)
$$

and we conclude.
Corollary 10.1.5.4. [Lee16, $\left.\mathbf{C C C}^{+} \mathbf{1 5 a}\right]$ Let $f \in S^{d} V$ be concise, and let $L \in V^{*} \backslash 0$ be arbitrary. Then

$$
\mathbf{R}_{S}(f) \geq \operatorname{dim}\left(\frac{\operatorname{Sym}\left(V^{*}\right)}{f^{\operatorname{ann}}:(L)+(L)}\right)
$$

Proof. The ideal $I\left(\left\{\left[\ell_{1}\right], \ldots,\left[\ell_{r}\right]\right\}\right)$ is reduced and empty in degree one. Thus $L$ is not a zero divisor in $I\left(\left\{\left[\ell_{1}\right], \ldots,\left[\ell_{r}\right]\right\}\right): L$, so Theorem 10.1.5.3 applies with $P=L$.

Let $f=e_{d, n} \in S^{d} \mathbb{C}^{n}$ be the $d$-th elementary symmetric function. Take $L=\frac{\partial}{\partial x_{n}} \in S^{1} \mathbb{C}^{n *}$ and apply Corollary 10.1.5.4. Exercise 10.1.5.2 implies

$$
\begin{aligned}
\left(e_{d, n}\right)^{a n n}: \frac{\partial}{\partial x_{n}}+\left(\frac{\partial}{\partial x_{n}}\right) & =\left(\frac{\partial}{\partial x_{n}} e_{d, n}\right)^{a n n}+\left(\frac{\partial}{\partial x_{n}}\right) \\
& =e_{d-1, n-1}^{a n n, \mathbb{C}^{n}}+\left(\frac{\partial}{\partial x_{n}}\right)
\end{aligned}
$$

where $e_{d-1, n-1}^{a n n, \mathbb{C}^{n}}$ is the annihilator of $e_{d-1, n-1}$ considered as a function of $n$ variables that does not involve $x_{n}$. Now

$$
\frac{\operatorname{Sym}\left(\mathbb{C}^{n *}\right)}{e_{d-1, n-1}^{a n n, \mathbb{C}^{n}}+\left(\frac{\partial}{\partial x_{n}}\right)} \simeq \frac{\operatorname{Sym}\left(\mathbb{C}^{(n-1) *}\right)}{e_{d-1, n-1}^{a n n, \mathbb{C}^{n-1}}}
$$

Exercise 10.1.5.5: (1) Show that for $t \leq \frac{\delta}{2},\left(\operatorname{Sym}\left(\mathbb{C}^{k *}\right) / e_{\delta, k}^{a n n}\right)_{t}$ consists of all square free monomials, so it is of dimension $\binom{k}{t}$. By symmetry, for $\frac{\delta}{2} \leq t \leq \delta, \operatorname{dim}\left(\left(\operatorname{Sym}\left(\mathbb{C}^{k *}\right) / e_{\delta, k}^{a n n}\right)_{t}\right)=\binom{k}{\delta-t}$. Finally, it is zero for $t>\delta$.

Putting it all together:
Theorem 10.1.5.6. [Lee16] For all $n$, and even $d$,

$$
\sum_{j=0}^{\frac{d}{2}}\binom{n}{j}-\binom{n-1}{\frac{d}{2}} \leq \mathbf{R}_{S}\left(e_{d, n}\right) \leq \sum_{j=0}^{\frac{d}{2}}\binom{n}{j}
$$

For all $n$, and odd $d$,

$$
\mathbf{R}_{S}\left(e_{d, n}\right)=\sum_{j=0}^{\frac{d-1}{2}}\binom{n}{j} .
$$

Proof. Let $d=2 k+1$ By Exercise 10.1.5.5 and the discussion above,

$$
\begin{aligned}
\mathbf{R}_{S}\left(e_{d, n}\right) & \geq 2 \sum_{j=0}^{\frac{d-1}{2}}\binom{n-1}{j} \\
& =\binom{n-1}{0}+\sum_{j=1}^{k}\left[\binom{n-1}{j}+\binom{n-1}{j-1}\right] \\
& =1+\sum_{j=1}^{k}\binom{n}{j}
\end{aligned}
$$

But this is the upper bound of Theorem 7.1.3.1. The even case is similar.

### 10.2. Cactus varieties and secant varieties

Recall that $\sigma_{r}(X) \subseteq \mathfrak{k}_{r}(X)$. Cactus border rank might appear to be just a curiosity, but the cactus variety turns out to be the obstruction to proving further lower bounds with current technology for $\underline{\mathbf{R}}\left(M_{\langle\mathbf{n}\rangle}\right)$ and places limits on the utility of lower bound techniques arising from determinantal equations. As I explain in §10.2.1, Almost all the equations discussed in this book
(method of partial derivatives Koszul flattenings, method of shifted partial derivatives, etc.), are equations for the cactus variety. Thus our border rank and symmetric border rank lower bounds are actually cactus border rank lower bounds. The reason this is a problem is explained in §10.2.2: the cactus varieties fill the ambient space much faster than the secant varieties. What is particularly surprising is that the cactus analog of the greater areole of $\S 5.4 .4$ already fills the ambient space.
10.2.1. Young-flattenings and the cactus variety. Recall the variety Flat ${ }_{r}^{i, d-i}(V):=\mathbb{P}\left\{P \in S^{d} V \mid \operatorname{rank}\left(P_{i, d-i}\right) \leq r\right\}$ from §6.2.2.
Proposition 10.2.1.1. [IK99, Thm. 5.3D] (also see [BB14]) For all $r, d, i$, $\mathfrak{k}_{r}\left(v_{d}(\mathbb{P} V)\right) \subseteq$ Flat $_{r}^{i, d-i}(V)$.

Proof. Say $[f] \in\left\langle v_{d}(Z)\right\rangle$ for a zero dimensional subscheme $Z \subset \mathbb{P} V$ of degree at most $r$, i.e., $\hat{f} \subset I(Z)_{d}^{\perp}$. We need to show $[f] \in F_{l a t}^{r}{ }_{r}^{i, d-i}(V)$. We have $\hat{f}\left(S^{i} V^{*}\right) \subset I(Z) \frac{\perp}{d}\left(S^{i} V^{*}\right)$. By the same argument as in the proof of the apolarity Lemma, $I(Z)_{d}^{\perp}\left(S^{i} V^{*}\right) \subset I(Z)_{d-i}^{\perp}$. By Exercise 10.1.1.2, $I(Z) \frac{\perp}{d-i}=\left\langle v_{d-i}(Z)\right\rangle$ and $\operatorname{dim}\left\langle v_{d-i}(Z)\right\rangle \leq \operatorname{deg}(Z)$.

More generally, Galcazka [Gal16] shows that for any variety $X, \mathfrak{k}_{r}(X)$ is in the zero set of any equations for $\sigma_{r}(X)$ arising from minors of a map between vector bundles. In particular:
Theorem 10.2.1.2. [Gal16] The equations from Koszul flattenings for $\sigma_{r}\left(S e g\left(\mathbb{P} A_{1} \times \cdots \times \mathbb{P} A_{k}\right)\right)$ and the equations from Koszul and Hilbert flattenings for $\sigma_{r}\left(v_{d}(\mathbb{P} V)\right)$ are equations for the corresponding cactus varieties.

This "explains" limits of Young flattenings because the cactus varieties fill the ambient space much faster than the secant varieties as I explain in the next section.
10.2.2. Local cactus rank and the local cactus variety. Let $X \subset \mathbb{P} V$ be a variety. Fix $x \in X$, and define the local $X$-cactus rank of $y \in \mathbb{P} V$ based at $x$, denoted $\mathfrak{l c r} r_{X, x}(y)$ to be the smallest $r$ such that there exists a length $r$ zero dimensional scheme $Z \subset X$, with the support of $Z$ equal to $\{x\}$, such that $y \in\langle Z\rangle$.

Define the local cactus variety of $X$ based at $x$ to be

$$
\mathfrak{L e}_{r}(X, x):=\bigcup_{\substack{Z \subset X, \operatorname{deg}(Z) \leq r, \\ \text { support }(Z)=x}}\langle Z\rangle \subset \mathbb{P} V .
$$

Of course $\mathfrak{l}_{r}(X, x) \subseteq \mathfrak{k}_{r}(X)$. Compare the local cactus variety with the greater areole of §5.4.4.

Given $f \in S^{d} \mathbb{C}^{n}$, define the dehomogenization $f_{x_{n}}:=f\left(x_{1}, \ldots, x_{n-1}, 1\right) \in$ $\mathbb{C}\left[x_{1}, \ldots, x_{n-1}\right]$ of $f$ with respect to $x_{n}$.
Theorem 10.2.2.1. [BR13] Let $f \in S^{d} \mathbb{C}^{n}$. Then

$$
\mathfrak{c r}_{S}(f) \leq \mathfrak{l c r}_{S}\left(f,\left[x_{n}^{d}\right]\right) \leq \operatorname{dim}\left(\operatorname{Partials}\left(f_{x_{n}}\right)\right) .
$$

Proof. The ideal defined by $\left(f_{x_{n}}\right)^{\text {ann }}$ in $\mathbb{C}\left[\frac{\partial}{\partial x_{1}}, \ldots, \frac{\partial}{\partial x_{n-1}}\right]$ may be homogenized to define a homogeneous ideal in $\mathbb{C}\left[\frac{\partial}{\partial x_{1}}, \ldots, \frac{\partial}{\partial x_{n}}\right]$. Consider the subscheme of $\mathbb{P}^{n-1}$ defined by this homogeneous ideal. It has degree equal $\operatorname{dim}\left(\operatorname{Partials}\left(f_{x_{n}}\right)\right)$ and support $\{[0, \ldots, 0,1]\}$. Assume $f$ is concise. (If it is not, just start over in the smaller space.) I next show the homogenized $\left(f_{x_{n}}\right)^{a n n}$ is contained in $f^{a n n}$.

Let $G \in\left(f_{x_{n}}\right)^{\text {ann }} \subset \mathbb{C}\left[\frac{\partial}{\partial x_{1}}, \ldots, \frac{\partial}{\partial x_{n-1}}\right]$. Write $G=G_{1}+\cdots+G_{r}$, where $\operatorname{deg}\left(G_{j}\right)=j$. (Note that $G_{0}=0$ in order that $G\left(f_{x_{n}}\right)=0$.) Similarly, write $f_{x_{n}}=f_{0}+\cdots+f_{d}$. Then $G\left(f_{x_{n}}\right)=0$ says that for all $e \in\{0, \ldots, d-$ $1\}, \sum_{j} G_{j}\left(f_{j+e}\right)=0$. Let $G^{h}$ be the homogenization of $G$, i.e., $G^{h}=$ $\left(\frac{\partial}{\partial x_{n}}\right)^{r-1} G_{1}+\left(\frac{\partial}{\partial x_{n}}\right)^{r-2} G_{2}+\cdots+G_{r}$. Then, since $f=x_{n}^{d} f_{0}+x_{n}^{d-1} f_{1}+\cdots+f_{d}$,

$$
G^{h}(f)=0
$$

if and only if

$$
\sum_{e} \sum_{j} x_{n}^{d-r-e} G_{j}\left(f_{e+j}\right)=\sum_{e} x_{n}^{d-r-e} \sum_{j} G_{j}\left(f_{e+j}\right)=0 .
$$

Thus the homogenization of $\left(f_{x_{n}}\right)^{a n n}$ is contained in $f^{a n n}$, and $\mathfrak{l c r}\left(f,\left[x_{n}^{d}\right]\right)$ is at most the degree of the scheme defined by $f_{x_{n}}$, which is $\operatorname{dim}\left(\operatorname{Partials}\left(f_{x_{n}}\right)\right)$.

Corollary 10.2.2.2. Set $N_{n, d}=\sum_{j=0}^{d}\binom{n-1+j-1}{j}$. Then $\mathfrak{F}_{N_{n, d}}\left(v_{d}\left(\mathbb{P}^{n-1}\right),\left[x_{n}^{d}\right]\right)=$ $\mathfrak{k}_{N_{n, d}}\left(v_{d}\left(\mathbb{P}^{n-1}\right)\right)=\mathbb{P} S^{d} \mathbb{C}^{n}$.

Note that $N_{n, d} \sim 2\binom{n+\left[\frac{d}{2}\right\rceil}{ d} \ll \frac{1}{n}\binom{n+d-1}{d}=: r_{n, d}$, the latter being the smallest $r$ such that $\sigma_{r}\left(v_{d}\left(\mathbb{P}^{n-1}\right)\right)=\mathbb{P} S^{d} \mathbb{C}^{n}$ (except for a few exceptions where $r_{n, d}$ is even larger [AH95], see [Lan12, §5.4.1] for a discussion).

So for example, if $n=2 k+1$ is odd then the upper bound for cactus rank in $S^{n} \mathbb{C}^{n^{2}}$ grows as $2\left(\frac{4 k^{2}+5 k}{k}\right) \simeq 2(2 k)^{2 k}$ plus lower order terms, while the border rank upper bound is: $\binom{4 k^{2}+6 k+1}{2 k+1} /(2 k+1)^{2} \simeq(2 k)^{4 k}$ plus lower order terms. The even case is similar.

For $S^{3} \mathbb{C}^{n}$, we have the more precise result:
Theorem 10.2.2.3. $[\mathbf{B R} 13] \mathfrak{k}_{2 n}\left(v_{3}\left(\mathbb{P}^{n-1}\right)\right)=\mathbb{P} S^{3} \mathbb{C}^{n}$.
Compare this with the secant variety $\sigma_{r}\left(v_{3}\left(\mathbb{P}^{n-1}\right)\right)$ which equals $\mathbb{P} S^{3} \mathbb{C}^{n}$ when $r=r_{n, 3}=\left\lceil\frac{1}{n}\binom{n+2}{3}\right\rceil \sim \frac{n^{2}}{6}$.

For tensors, the same methods show $\mathfrak{l k}_{2(\mathbf{a}+\mathbf{b}+\mathbf{c}-2)}(\operatorname{Seg}(\mathbb{P} A \times \mathbb{P} B \times \mathbb{P} C),[a \otimes b \otimes c])=$ $\mathbb{P}(A \otimes B \otimes C)$ (J. Buczynski, personal communication). The precise filling $r$ for the cactus variety is not known. However, since for $\mathbb{C}^{\mathbf{m}} \otimes \mathbb{C}^{\mathbf{m}} \otimes \mathbb{C}^{\mathbf{m}}$ it is at most $6 \mathbf{m}-4$, one will never prove super-linear border rank bounds for tensors with determinantal equations.

### 10.3. Non-normality of $\mathcal{D e} t_{n}$

We have already seen that $C h_{n}(V)$ is not normal, but that the coordinate ring of the normalization, the coordinate ring of the orbit $G L(V) \cdot\left[x_{1} \cdots x_{n}\right]$, and the coordinate ring of $C h_{n}(V)$ are all closely related. For the determinant we know far less. Not surprisingly, $\mathcal{D e t}_{n}$ is not normal. I explain the proof in this section, which unfortunately is indirect, so it does not indicate relations between the three rings. By Remark 9.7.2.3 and the boundary component of §6.7.1, we know that the coordinate ring of the normalization of $\mathcal{D e} t_{n}$ and the polynomial part of the coordinate ring of the orbit do not agree, but we know little about their difference. Here I at least show that the normalization is distinct from $\mathcal{D e t}_{n}$.

I begin with generalities on $G L(W)$-orbits of points $P \in V$ with closed $S L(W)$-orbit. By a theorem of Kempf [Kem78, Cor. 5.1], the $S L(W)$-orbit of $P$ is closed if the $S L(W)$-isotropy group is not contained in any proper parabolic subgroup of $S L(W)$, which is the case for the $S L(W)$-stabilizers of $\operatorname{perm}_{n}$ and $\operatorname{det}_{n}$.

I follow [Lan15a] in this section, which gives an exposition of the results of [Kum13].
10.3.1. Generalities on $G L(W)$-orbit closures. Throughout this section I make the following assumptions and adopt the following notation:

Set up:

- $V$ is a $G L(W)$-module, and $P \in V$.
- $\mathcal{P}^{0}:=G L(W) \cdot P$ and $\mathcal{P}:=\overline{G L(W) \cdot P}$ respectively denote the $G L(W)$-orbit and $G L(W)$-orbit closure of $P$, and $\partial \mathcal{P}=\mathcal{P} \backslash \mathcal{P}^{0}$ denotes the boundary, which is assumed to be more than zero (otherwise $[\mathcal{P}]$ is homogeneous).

Assumptions:
(1) $P \in V$ is such that the $S L(W)$-orbit of $P$ is closed.
(2) The stabilizer $G_{P} \subset G L(W)$ is reductive, which is equivalent (by a theorem of Matsushima [Mat60]) to requiring that $\mathcal{P}^{0}$ is an affine variety.

This situation holds when $V=S^{n} W, \operatorname{dim} W=n^{2}$ and $P=\operatorname{det}_{n}$ or $\operatorname{perm}_{n}$ as well as when $\operatorname{dim} W=r n$ and $P=S_{n}^{r}:=\sum_{j=1}^{r} x_{1}^{j} \cdots x_{n}^{j}$, the sum-product polynomial, in which case $\mathcal{P}=\hat{\sigma}_{r}\left(C h_{n}(W)\right)$.
Lemma 10.3.1.1. [Kum13] Assumptions as in (10.3.1). Let $M \subset \mathbb{C}[\mathcal{P}]$ be a nonzero $G L(W)$-module, and let $\operatorname{Zeros}(M)=\{y \in \mathcal{P} \mid f(y)=0 \forall f \in M\}$ denote its zero set. Then $0 \subseteq \operatorname{Zeros}(M) \subseteq \partial \mathcal{P}$.

If moreover $M \subset I(\partial \mathcal{P})$, then as sets, $\operatorname{Zeros}(M)=\partial \mathcal{P}$.
Proof. Since $\operatorname{Zeros}(M)$ is a $G L(W)$-stable subset, if it contains a point of $\mathcal{P}^{0}$ it must contain all of $\mathcal{P}^{0}$ and thus $M$ vanishes identically on $\mathcal{P}$, which cannot happen as $M$ is nonzero. Thus $\operatorname{Zeros}(M) \subseteq \partial \mathcal{P}$. For the second assertion, since $M \subset I(\partial \mathcal{P})$, we also have $\operatorname{Zeros}(M) \supseteq \partial \mathcal{P}$.

Proposition 10.3.1.2. [Kum13] Assumptions as in (10.3.1). The space of $S L(W)$-invariants of positive degree in the coordinate ring of $\mathcal{P}, \mathbb{C}[\mathcal{P}]_{>0}^{S L(W)}$, is non-empty and contained in $I(\partial \mathcal{P})$. Moreover,
(1) any element of $\mathbb{C}[\mathcal{P}]_{>0}^{S L(W)}$ cuts out $\partial \mathcal{P}$ set-theoretically, and
(2) the components of $\partial \mathcal{P}$ all have codimension one in $\mathcal{P}$.

Proof. To study $\mathbb{C}[\mathcal{P}]^{S L(W)}$, consider the GIT quotient $\mathcal{P} / / S L(W)$ whose coordinate ring, by definition, is $\mathbb{C}[\mathcal{P}]^{S L(W)}$. It parametrizes the closed $S L(W)$-orbits in $\mathcal{P}$, so it is non-empty. Thus $\mathbb{C}[\mathcal{P}]^{S L(W)}$ is nontrivial.

Claim: every $S L(W)$-orbit in $\partial P$ contains $\{0\}$ in its closure, i.e., $\partial \mathcal{P}$ maps to zero in the GIT quotient. This will imply any $S L(W)$-invariant of positive degree is in $I(\partial \mathcal{P})$ because any non-constant function on the GIT quotient vanishes on the inverse image of [0]. Thus (1) follows from Lemma 10.3.1.1. The zero set of a single polynomial, if it is not empty, has codimension one, which implies the components of $\partial \mathcal{P}$ are all of codimension one, proving (2).

Let $\rho: G L(W) \rightarrow G L(V)$ denote the representation. It remains to show $\partial \mathcal{P}$ maps to zero in $\mathcal{P} / / S L(W)$. This GIT quotient inherits a $\mathbb{C}^{*}$ action via $\rho(\lambda \mathrm{Id})$, for $\lambda \in \mathbb{C}^{*}$. The normalization of $\mathcal{P} / / S L(W)$ is just the affine line $\mathbb{A}^{1}=\mathbb{C}$. To see this, consider the $\mathbb{C}^{*}$-equivariant map $\sigma: \mathbb{C} \rightarrow \mathcal{P}$ given by $z \mapsto \rho(z \mathrm{Id}) \cdot P$, which descends to a map $\bar{\sigma}: \mathbb{C} \rightarrow \mathcal{P} / / S L(W)$. Since the $S L(W)$-orbit of $P$ is closed, for any $\lambda \in \mathbb{C}^{*}, \rho(\lambda \mathrm{Id}) P$ does not map to zero in the GIT quotient, so $\bar{\sigma}^{-1}([0])=\{0\}$ as a set. Lemma 9.7.1.1 applies so $\bar{\sigma}$ is finite and gives the normalization. Finally, were there a closed nonzero orbit in $\partial \mathcal{P}$, it would have to equal $S L(W) \cdot \sigma(\lambda)$ for some $\lambda \in \mathbb{C}^{*}$ since $\bar{\sigma}$ is surjective. But $S L(W) \cdot \sigma(\lambda)$ is contained in the image of $\mathcal{P}^{0}$ in $\mathcal{P} / / S L(W)$.

Remark 10.3.1.3. That each irreducible component of $\partial \mathcal{P}$ is of codimension one in $\mathcal{P}$ is due to Matsushima [Mat60]. It is a consequence of his result that $\mathcal{P}^{0}$ is an affine variety if and only if the stabilizer is reductive.

The key to proving non-normality of $\hat{\mathcal{D e}} t_{n}$ and $\hat{\mathcal{P e r}} \hat{n} n_{n}$ is to find an $S L(W)$-invariant in the coordinate ring of the normalization (which has a $G L(W)$-grading), which does not occur in the corresponding graded component of the coordinate ring of $S^{n} W$, so it cannot occur in the coordinate ring of any $G L(W)$-subvariety.
Lemma 10.3.1.4. Assumptions as in (10.3.1). Let $V=S^{n} W$ and let $d$ be the smallest positive $G L(W)$-degree such that $\mathbb{C}\left[\mathcal{P}^{0}\right]_{d}^{S L(W)} \neq 0$. If $n$ is even (resp. odd) and $d<n \mathbf{w}$ (resp. $d<2 n \mathbf{w}$ ) then $\mathcal{P}$ is not normal.

Proof. Since $\mathcal{P}^{0} \subset \mathcal{P}$ is a Zariski open subset, we have the equality of $G L(W)$-modules $\mathbb{C}(\mathcal{P})=\mathbb{C}\left(\mathcal{P}^{0}\right)$. By restriction of functions $\mathbb{C}[\mathcal{P}] \subset \mathbb{C}\left[\mathcal{P}^{0}\right]$ and thus $\mathbb{C}[\mathcal{P}]^{S L(W)} \subset \mathbb{C}\left[\mathcal{P}^{0}\right]^{S L(W)}$. Now $\mathcal{P}^{0} / / S L(W)=\mathcal{P}^{0} / S L(W) \simeq \mathbb{C}^{*}$ because $S L(W) \cdot P$ is closed, so $\mathbb{C}\left[\mathcal{P}^{0}\right]^{S L(W)} \simeq \bigoplus_{k \in \mathbb{Z}} \mathbb{C}\left\{z^{k}\right\}$. Under this identification, $z$ has $G L(W)$-degree $d$. By Proposition 10.3.1.2, $\mathbb{C}[\mathcal{P}]^{S L(W)} \neq$ 0 . Let $h \in \mathbb{C}[\mathcal{P}]^{S L(W)}$ be the smallest element in positive degree. Then $h=z^{k}$ for some $k$. Were $\mathcal{P}$ normal, we would have $k=1$.

But now we also have a surjection $\mathbb{C}\left[S^{n} W\right] \rightarrow \mathbb{C}[\mathcal{P}]$, and by [How87, Prop. 4.3a], the smallest possible $G L(W)$-degree of an $S L(W)$-invariant in $\mathbb{C}\left[S^{n} W\right]$ when $n$ is even (resp. odd) is $\mathbf{w} n$ (resp. $2 \mathbf{w} n$ ) which would occur in $S^{\mathbf{w}}\left(S^{n} W\right)$ (resp. $S^{2 \mathbf{w}}\left(S^{n} W\right)$ ). We obtain a contradiction.
10.3.2. Case of $P=\operatorname{det}_{n}$ and $P=\operatorname{perm}_{n}$.

Theorem 10.3.2.1 (Kumar [Kum13]). For all $n \geq 3, \operatorname{Det}_{n}=\overline{G L_{n^{2}} \cdot\left[\operatorname{det}_{n}\right]}$ and $\mathcal{P e r m} m_{n}^{n}=\overline{G L_{n^{2}} \cdot\left[\operatorname{perm}_{n}\right]}$ are not normal. For all $n \geq 2 m$ (the range of interest), $\mathcal{P e r m}_{n}^{m}=\overline{G L_{n^{2}} \cdot\left[\ell^{n-m} \text { perm }_{m}\right]}$ is not normal.

I give the proof for $\mathcal{D e} t_{n}$, the case of $\mathcal{P e r m} n_{n}^{n}$ is an easy exercise. Despite the variety $\operatorname{Zeros}\left(\ell^{n-m} \operatorname{perm}_{m}\right)$ being much more singular than $\operatorname{Zeros}\left(\operatorname{perm}_{n}\right)$, the proof for $\mathcal{P e r m} n_{n}^{m}$ is more difficult, see [Kum13].

Proof. Let $\mathbb{C}\left[\mathcal{D e t}_{n}^{0}\right]_{k-G L}^{S L(W)}$ denote the degree $k G L$-degree component of $\mathbb{C}\left[\mathcal{D} e t_{n}^{0}\right]^{S L(W)}$ as defined in the proof of Lemma 10.3.1.4. I will show that when $n$ is congruent to 0 or $1 \bmod 4, \mathbb{C}\left[\mathcal{D e} t_{n}^{0}\right]_{n-G L}^{S L(W)} \neq 0$ and when $n$ is congruent to 2 or $3 \bmod 4, \mathbb{C}\left[\mathcal{D e} t_{n}^{0}\right]_{2 n-G L}^{S L(W)} \neq 0$. Since $n, 2 n<\left(n^{2}\right) n$ Lemma 10.3.1.4 applies.

The $S L(W)$-trivial modules are $\left(\Lambda^{n^{2}} W\right)^{\otimes s}=S_{\left(s^{n^{2}}\right)} W$. Write $W=$ $E \otimes F$. We want to determine the lowest degree trivial $S L(W)$-module that has a $G_{\operatorname{det}_{n}}=\left(S L(E) \times S L(F) / \mu_{n}\right) \rtimes \mathbb{Z}_{2}$ invariant. We have the
decomposition $\left(\Lambda^{n^{2}} W\right)^{\otimes s}=\left(\bigoplus_{|\pi|=n^{2}} S_{\pi} E \otimes S_{\pi^{\prime}} F\right)^{\otimes s}$, where $\pi^{\prime}$ is the conjugate partition to $\pi$. Thus $\left(\Lambda^{n^{2}} W\right)^{\otimes s}$ is the trivial $S L(E) \times S L(F)$ module $\left(S_{\left(n^{n}\right)} E \otimes S_{\left(n^{n}\right)} F\right)^{\otimes s}=S_{\left((s n)^{n}\right)} E \otimes S_{\left((s n)^{n}\right)} F$. Now consider the effect of the $\mathbb{Z}_{2} \subset G_{\operatorname{det}_{n}}$ with generator $\tau \in G L(W)$. It sends $e_{i} \otimes f_{j}$ to $e_{j} \otimes f_{i}$, so acting on $W$ it has +1 eigenspace $\left\{e_{i} \otimes f_{j}+e_{j} \otimes f_{i} \mid i \leq j\right\}$ and -1 eigenspace $\left\{e_{i} \otimes f_{j}-e_{j} \otimes f_{i} \mid 1 \leq i<j \leq n\right\}$. Thus it acts on the one-dimensional vector space $\left(\Lambda^{n^{2}} W\right)^{\otimes s}$ by $\left((-1)^{\binom{n}{2}}\right)^{s}$, i.e., by -1 if $n \equiv 2,3 \bmod 4$ and $s$ is odd and by 1 otherwise. We conclude that there is an invariant as needed for Lemma 10.3.1.4.
Remark 10.3.2.2. In the language of $\S 8.8 .2$, in the proof above we saw $k_{s^{n^{2}},(s n)^{n},(s n)^{n}}=1, s k_{(s n)^{n},(s n)^{n}}^{s^{n^{2}}}=1$ for all $s$ when $\binom{n}{2}$ is even, and $s k_{(s n)^{n},(s n)^{n}}^{s^{n^{2}}}=$ 1 for even $s$ when $\binom{n}{2}$ is odd and is zero for odd $s$.
Exercise 10.3.2.3: (2) Write out the proof of the non-normality of $\mathcal{P e r m} n_{n}^{n}$.
Exercise 10.3.2.4: (2) Show the same method gives another proof that $C h_{n}(W)$ is not normal.

Exercise 10.3.2.5: (2) Show a variant of the above holds for any reductive group with a nontrivial center (one gets a $\mathbb{Z}^{k}$-grading of modules if the center is $k$-dimensional), in particular it holds for $G=G L(A) \times G L(B) \times G L(C)$. Use this to show that $\sigma_{r}(\operatorname{Seg}(\mathbb{P} A \times \mathbb{P} B \times \mathbb{P} C))$ is not normal when $\operatorname{dim} A=$ $\operatorname{dim} B=\operatorname{dim} C=r>2$.

### 10.4. The minimal free resolution of the ideal generated by minors of size $r+1$

I give an exposition of Lascoux's computation of the minimal free resolution of the ideals of the varieties of matrices of rank at most $r$ from [Las78]. I expect it will be useful for the study of Valiant's hypothesis, as from it one can extract numerous algebraic properties of the determinant polynomial.

I follow the exposition in [ELSW15], which is based on the presentation in [Wey03].
10.4.1. Statement of the result. Let $E, F=\mathbb{C}^{n}$, give $E \otimes F$ coordinates $\left(x_{j}^{i}\right)$, with $1 \leq i, j \leq n$. Let $\hat{\sigma}_{r}=\hat{\sigma}_{r}\left(\operatorname{Seg}\left(\mathbb{P}^{n-1} \times \mathbb{P}^{n-1}\right)\right) \subset \mathbb{C}^{n} \otimes \mathbb{C}^{n}=E^{*} \otimes F^{*}$ denote the variety of $n \times n$ matrices of rank at most $r$. By "degree $S_{\pi} E^{\prime \prime}$, I mean $|\pi|=p_{1}+\cdots+p_{n}$, where $\pi=\left(p_{1}, \ldots, p_{n}\right)$. Write $\pi+\tilde{\pi}=\left(p_{1}+\right.$ $\left.\tilde{p}_{1}, \ldots, p_{n}+\tilde{p}_{n}\right)$.

Recall from $\S 8.1 .2$ that the weight (under $G L(E) \times G L(F)$ ) of a monomial $x_{j_{1}}^{i_{1}} \cdots x_{j_{q}}^{i_{q}} \in S^{q}(E \otimes F)$ is given by a pair of $n$-tuples $\left(\left(w_{1}^{E}, \ldots, w_{n}^{E}\right),\left(w_{1}^{F}, \ldots, w_{n}^{F}\right)\right)$ where $w_{s}^{E}$ is the number of $i_{\alpha}$ 's equal to $s$ and $w_{t}^{F}$ is the number of $j_{\alpha}$ 's equal to $t$.

Theorem 10.4.1.1. [Las78] Let $0 \rightarrow F_{N} \rightarrow \cdots \rightarrow F_{1} \rightarrow \operatorname{Sym}(E \otimes F)=$ $F_{0} \rightarrow \mathbb{C}\left[\hat{\sigma}_{r}\right] \rightarrow 0$ denote the minimal free resolution of $\mathbb{C}\left[\hat{\sigma}_{r}\right]$. Then
(1) $N=(n-r)^{2}$, i.e., $\hat{\sigma}_{r}$ is arithmetically Cohen-Macaulay.
(2) $\hat{\sigma}_{r}$ is Gorenstein, i.e., $F_{N} \simeq \operatorname{Sym}(E \otimes F)$, generated by $S_{(n-r)^{n}} E \otimes S_{(n-r)^{n}} F$. In particular $F_{N-j} \simeq F_{j}$ as $S L(E) \times S L(F)$ - modules, although they are not isomorphic as $G L(E) \times G L(F)$-modules.
(3) For $1 \leq j \leq N-1$, the space $F_{j}$ has generating modules of degree $s r+j$ where $1 \leq s \leq\lfloor\sqrt{j}\rfloor$. The modules of degree $r+j$ form the generators of the linear strand of the minimal free resolution.
(4) The generating module of $F_{j}$ is multiplicity free.
(5) Let $\alpha, \beta$ be (possibly zero) partitions such that $l(\alpha), l(\beta) \leq s$. Independent of the lengths (even if they are zero), write $\alpha=$ $\left(\alpha_{1}, \ldots, \alpha_{s}\right), \beta=\left(\beta_{1}, \ldots, \beta_{s}\right)$. The space of degree sr $+j$ generators of $F_{j}$, for $1 \leq j \leq N$ is the module

$$
\begin{equation*}
M_{j, r s+j}=\bigoplus_{\substack{|\alpha|| | \beta==-s^{2} \\ l(\alpha), l(\beta) \leq s}} S_{\left(s^{r+s}\right)+\left(\alpha, 0^{r}, \beta^{\prime}\right)} E \otimes S_{\left(s^{r+s}\right)+\left(\beta, 0^{r}, \alpha^{\prime}\right)} F . \tag{10.4.1}
\end{equation*}
$$

The Young diagrams of the modules are depicted in Figure 1 below.


Figure 10.4.1. Partition $\pi$ and pairs of partitions $\left(s^{r+s}\right)+$ $\left(\alpha, 0^{r}, \beta^{\prime}\right)=w \cdot \pi$ and $\left(s^{r+s}\right)+\left(\beta, 0^{r}, \alpha^{\prime}\right)=\pi^{\prime}$ it gives rise to in the resolution (see $\S 10.4 .5$ for explanations).
(6) In particular the generator of the linear component of $F_{j}$ is

$$
\begin{equation*}
M_{j, j+r}=\bigoplus_{a+b=j-1}=S_{a+1,1^{r+b}} E \otimes S_{b+1,1^{r+a}} F \tag{10.4.2}
\end{equation*}
$$

Remark 10.4.1.2. The module $M_{j, j+r}$ admits a basis as follows: form a size $r+j$ submatrix using $r+b+1$ distinct rows, repeating a subset of $a$ rows to have the correct number of rows and $r+a+1$ distinct columns, repeating a subset of $b$ columns, and then performing a "tensor Laplace expansion" as described below.
10.4.2. The Koszul resolution. The minimal free resolution of $\operatorname{Sym}(V)_{>0}$ is given by the exact complex

$$
\begin{equation*}
\cdots \rightarrow S^{q-1} V \otimes \Lambda^{p+2} V \rightarrow S^{q} V \otimes \Lambda^{p+1} V \rightarrow S^{q+1} V \otimes \Lambda^{p} V \rightarrow \cdots \tag{10.4.3}
\end{equation*}
$$

The maps are given by the transpose of exterior derivative (Koszul) map $d_{p, q}: S^{q} V^{*} \otimes \Lambda^{p+1} V^{*} \rightarrow S^{q-1} V^{*} \otimes \Lambda^{p+2} V^{*}$. Write $d_{p, q}^{T}: S^{q-1} V \otimes \Lambda^{p+2} V \rightarrow$ $S^{q} V \otimes \Lambda^{p+1} V$. The Pieri rule (§8.1.3) implies the $G L(V)$-decomposition $S^{q} V \otimes \Lambda^{p+1} V=S_{\left(q, 1^{p+1}\right)} V \oplus S_{\left(q+1,1^{p}\right)} V$, so the kernel of $d_{p, q}^{T}$ is the first module, which also is the image of $d_{p+1, q-1}^{T}$.

Explicitly, $d_{p, q}^{T}$ is the composition of polarization $\left(\Lambda^{p+2} V \rightarrow \Lambda^{p+1} V \otimes V\right)$ and multiplication:

$$
S^{q-1} V \otimes \Lambda^{p+2} V \rightarrow S^{q-1} V \otimes \Lambda^{p+1} V \otimes V \rightarrow S^{q} V \otimes \Lambda^{p+1} V
$$

For the minimal free resolution of any ideal, the linear strand will embed inside (10.4.3).
10.4.3. Geometry of the terms in the linear strand. For $T \in S^{\kappa} V \otimes V^{\otimes j}$, and $P \in S^{\ell} V$, introduce notation for multiplication on the first factor, $T \cdot P \in S^{\kappa+\ell} V \otimes V^{\otimes j}$. Write $F_{j}=M_{j} \cdot \operatorname{Sym}(V)$. As always, $M_{0}=\mathbb{C}$. Note that $F_{1}=M_{1} \cdot \operatorname{Sym}(E \otimes F)$, where $M_{1}=M_{1, r+1}=\Lambda^{r+1} E \otimes \Lambda^{r+1} F$, the size $r+1$ minors which generate the ideal. The syzygies among the elements of $F_{1}$ are generated by

$$
M_{2}=M_{2, r+2}:=S_{1^{r+2}} E \otimes S_{21^{r}} F \oplus S_{21^{r}} E \otimes S_{1^{r+2}} F \subset I\left(\sigma_{r}\right)_{r+2} \otimes V
$$

(i.e., $F_{2}=M_{2} \cdot \operatorname{Sym}(E \otimes F)$ ) where elements in the first module may be obtained by choosing $r+1$ rows and $r+2$ columns, forming a size $r+2$ square matrix by repeating one of the rows, then doing a "tensor Laplace expansion" as follows:

In the case $r=1$ the highest weight vector of $M_{2}$ is

$$
\begin{align*}
S_{123}^{1 \mid 12}: & =\left(x_{2}^{1} x_{3}^{2}-x_{2}^{2} x_{3}^{1}\right) \otimes x_{1}^{1}-\left(x_{1}^{1} x_{3}^{2}-x_{1}^{2} x_{3}^{1}\right) \otimes x_{2}^{1}+\left(x_{1}^{1} x_{2}^{2}-x_{2}^{1} x_{1}^{2}\right) \otimes x_{3}^{1}  \tag{10.4.4}\\
& =M_{23}^{12} \otimes x_{1}^{1}-M_{13}^{12} \otimes x_{2}^{1}+M_{12}^{12} \otimes x_{3}^{1}
\end{align*}
$$

a tensor Laplace expansion of a $3 \times 3$ matrix with first row repeated. In general $M_{J}^{I}$ will denote the minor obtained from the submatrix with indices $I, J$. To see (10.4.4) is indeed a highest weight vector, first observe that it has the correct weights in both $E$ and $F$, and that in the $F$-indices $\{1,2,3\}$ it is skew and that in the first two $E$ indices it is also skew. Finally to see it is a highest weight vector note that any raising operator sends it to zero. Also note that under the multiplication map $S^{2} V \otimes V \rightarrow S^{3} V$ the element maps to zero, because the map corresponds to converting a tensor Laplace expansion to an actual one, but the determinant of a matrix with a repeated row is zero.

In general, a basis of $S_{\pi} E \otimes S_{\mu} F$ is indexed by pairs of semi-standard Young tableau in $\pi$ and $\mu$. In the linear strand, all partitions appearing are hooks, a basis of $S_{a, 1^{b}} E$ is given by two sequences of integers taken from [ $n$ ], one weakly increasing of length $a$ and one strictly increasing of length $b$, where the first integer in the first sequence is at least the first integer in the second sequence.

A highest weight vector in $S_{21^{r}} E \otimes S_{1^{r+2}} F$ is

$$
S_{1, \ldots, r+2}^{1 \mid 1, \ldots, r+1}=M_{2, \ldots, r+2}^{1, \ldots, r+1} \otimes x_{1}^{1}-M_{1,3, \ldots, r+1}^{1, \ldots, r+1} \otimes x_{2}^{1}+\cdots+(-1)^{r} M_{1, \ldots, r+1}^{1, \ldots, r+1} \otimes x_{r+2}^{1}
$$

a tensor Laplace expansion of a size $r+2$ matrix with repeated first row. The same argument as above shows it has the desired properties. Other basis vectors are obtained by applying lowering operators to the highest weight vector, so their expressions will be more complicated.

Remark 10.4.3.1. If we chose a size $r+2$ submatrix, and perform a tensor Laplace expansion of its determinant about two different rows, the difference of the two expressions corresponds to a linear syzygy, but such linear syzygies are in the span of $M_{2}$. These expressions are important for comparison with the permanent, as they are the only linear syzygies for the ideal generated by the size $r+1$ sub-permanents, where one takes the permanental Laplace expansion.

Continuing, $F_{3}$ is generated by the module

$$
M_{3, r+3}=S_{1^{r+3}} E \otimes S_{3,1^{r}} F \oplus S_{2,1^{r+1}} E \otimes S_{2,1^{r+1}} F \oplus S_{3,1^{r}} E \otimes S_{1^{r+3}} F \subset M_{2} \otimes V
$$

These modules admit bases of double tensor Laplace type expansions of a square submatrix of size $r+3$. In the first case, the highest weight vector is obtained from the submatrix whose rows are the first $r+3$ rows of the original matrix, and whose columns are the first $r$-columns with the first column occuring three times. For the second module, the highest weight vector is obtained from the submatrix whose rows and columns are the first $r+2$ such, with the first row and column occuring twice. A highest weight
vector for $S_{3,1^{r}} E \otimes S_{1^{r+3}} F$ is

$$
\begin{aligned}
S_{1, \ldots, r+3}^{11 \mid 1, \ldots, r+1} & =\sum_{1 \leq \beta_{1}<\beta_{2} \leq r+3}(-1)^{\beta_{1}+\beta_{2}} M_{1, \ldots, \hat{\beta_{1}}, \ldots, \hat{\beta}_{2}, \ldots, r+3}^{1, \ldots, r+1} \otimes\left(x_{\beta_{1}}^{1} \wedge x_{\beta_{2}}^{1}\right) \\
& =\sum_{\beta=1}^{r+3}(-1)^{\beta+1} S_{1, \ldots, \ldots, \ldots,,_{r+3}^{1 \mid 1, \ldots, i_{r+1}} \otimes x_{\beta}^{1} .} .
\end{aligned}
$$

Here $S_{1, \ldots, \hat{\beta}, \ldots, r+3}^{1 \mid 1, \ldots, i_{r+1}}$ is defined in the same way as the highest weight vector.
A highest weight vector for $S_{2,1^{r+1}} E \otimes S_{2,1^{r+1}} F$ is

$$
\begin{aligned}
S_{1 \mid 1, \ldots, r+2}^{1 \mid 1, \ldots, r+3} & =\sum_{\alpha, \beta=1}^{r+3}(-1)^{\alpha+\beta} M_{1, \ldots, \hat{\beta}, \ldots, i+2}^{1, \ldots, \hat{\alpha}, \ldots, r+2} \otimes\left(x_{1}^{\alpha} \wedge x_{\beta}^{1}\right) \\
& =\sum_{\beta=1}^{r+3}(-1)^{\beta+1} S_{1 \mid 1, \ldots, \hat{\beta}, \ldots, r+2}^{1, \ldots, r+2} \otimes x_{\beta}^{1}-\sum_{\alpha=1}^{r+3}(-1)^{\alpha+1} S_{1, \ldots, r+2}^{1 \mid 1, \ldots, \hat{\alpha}, \ldots, r+3} \otimes x_{1}^{\alpha}
\end{aligned}
$$

Here $S_{1 \mid 1, \ldots, \hat{\beta}, \ldots, r+2}^{1, \ldots, r+2}, S_{1, \ldots, r+2}^{1 \mid 1, \ldots, \hat{\alpha}, \ldots, r+3}$ are defined in the same way as the corresponding highest weight vectors.
Proposition 10.4.3.2. The highest weight vector of $S_{p+1,1^{r+q}} E \otimes S_{q+1,1^{r+p}} F \subset$ $M_{p+q+1, r+p+q+1}$ is

$$
\begin{aligned}
& S_{1 q \mid 1, \ldots, r+p+1}^{1^{p} \mid 1, \ldots, r+q+1}= \\
& \sum_{\substack{I \subset[r+q+1]| || |=q, J \subset[r+p+1],|J|=p}}(-1)^{|I|+|J|} M_{1, \ldots, \ldots, \hat{j}_{1}, \ldots, \hat{j}_{p}, \ldots,(r+p+1)}^{1, \ldots \hat{i}_{1}, \ldots, \hat{i}_{q}, \ldots,(r+q+1)} \otimes\left(x_{j_{1}}^{1} \wedge \cdots \wedge x_{j_{p}}^{1} \wedge x_{1}^{i_{1}} \wedge \cdots \wedge x_{1}^{i_{q}}\right) .
\end{aligned}
$$

A hatted index is one that is omitted from the summation.
Proof. It is clear the expression has the correct weight and is a highest weight vector, and that it lies in $S^{r+1} V \otimes \Lambda^{p+q} V$. I now show it maps to zero under the differential.

Under the map $d^{T}: S^{r+1} V \otimes \Lambda^{p+q} V \rightarrow S^{r} V \otimes \Lambda^{p+q+1} V$, the element $S_{1 q \mid 1, \ldots, r+p+1}^{1^{p} \mid 1, \ldots, r+q+1}$ maps to:

$$
\begin{aligned}
& \sum_{\substack{I \subset[r+q+1]|I|=q, J \subset[r+p+1]| | J| |=p}}(-1)^{|I|+|J|} \\
& {\left[\sum_{\alpha \in I}(-1)^{p+\alpha} M_{1, \ldots, \hat{y}_{1}, \ldots, \hat{i}_{p}, \ldots,(r+p+1)}^{1, \ldots, \hat{i}_{1}, \ldots,(r+q+1)} x_{1}^{i_{\alpha}} \otimes\left(x_{j_{1}}^{1} \wedge \cdots \wedge x_{j_{p}}^{1} \wedge x_{1}^{i_{1}} \wedge \cdots \wedge \hat{x}_{1}^{i_{\alpha}} \wedge \cdots \wedge x_{1}^{i_{q}}\right)\right.} \\
& \left.+\sum_{\beta \in J}(-1)^{\beta} M_{1, \ldots, \hat{y}_{1}, \ldots, \hat{y}_{p}, \ldots,(r+p+1)}^{1, \ldots, \hat{i}_{1}, \ldots, \hat{i}_{q}, \ldots,(r+q+1)} x_{j_{\beta}}^{1} \otimes\left(x_{j_{1}}^{1} \wedge \cdots \wedge \hat{x}_{j_{\beta}}^{1} \wedge \cdots \wedge x_{j_{p}}^{1} \wedge x_{1}^{i_{1}} \wedge \cdots \wedge x_{1}^{i_{q}}\right)\right]
\end{aligned}
$$

Fix $I$ and all indices in $J$ but one, call the resulting index set $J^{\prime}$, and consider the resulting term
$\sum_{\beta \in[r+p+1] \backslash J^{\prime}}(-1)^{f\left(\beta, J^{\prime}\right)} M_{1, \ldots, \hat{j}_{1}^{\prime}, \ldots, \hat{j}_{p-1}^{\prime}, \ldots,(r+p+1)}^{1, \ldots, \hat{i}_{1}, \ldots, \hat{i}_{q}, \ldots(r+q+1)} x_{\beta}^{1} \otimes\left(x_{j_{1}^{\prime}}^{1} \wedge \cdots \wedge x_{j_{p-1}^{\prime}}^{1} \wedge x_{1}^{i_{1}} \wedge \cdots \wedge x_{1}^{i_{q}}\right)$
where $f\left(\beta, J^{\prime}\right)$ equals the number of $j^{\prime} \in J$ less than $\beta$. This term is the Laplace expansion of the determinant of a matrix of size $r+1$ which has its first row appearing twice, and is thus zero.

Notice that if $q, p>0$, then $S_{1^{q} \mid 1, \ldots, r+p+1}^{1^{p} \mid 1, \ldots, r+q+1}$ is the sum of terms including $S_{1^{q-1} \mid 1, \ldots, r+p+1}^{1^{p} \mid 1, \ldots, r+q} \otimes x_{1}^{r+q+1}$ and $S_{1^{q} \mid 1, \ldots, r+p}^{1^{p-1} \mid 1, \ldots, r+q+1} \otimes x_{r+p+1}^{1}$. This implies the following corollary:
Corollary 10.4.3.3 (Roberts [Rob17]). Each module $S_{a, 1^{r+b}} E \otimes S_{b, 1^{r+a}} F$, where $a+b=j$ that appears with multiplicity one in $F_{j, j+r}$, appears with multiplicity two in $F_{j-1, j+r}$ if $a, b>0$, and multiplicity one if $a$ or $b$ is zero. The map $F_{j, j+r+1} \rightarrow F_{j-1, j+r+1}$ restricted to $S_{a, 1^{r+b}} E \otimes S_{b, 1^{r+a}} F$, maps nonzero to both $\left(S_{a-1,1^{r+b}} E \otimes S_{b, 1^{r+a-1}} F\right) \cdot E \otimes F$ and $\left(S_{a, 1^{r+b-1}} E \otimes S_{b-1,1^{r+a}} F\right)$. $E \otimes F$.

Proof. The multiplicities and realizations come from applying the Pieri rule. (Note that if $a$ is zero the first module does not exist and if $b$ is zero the second module does not exist.) That the maps to each of these is non-zero follows from the observation above.

Remark 10.4.3.4. In [Rob17] it is proven more generally that all the natural realizations of the irreducible modules in $M_{j}$ have non-zero maps onto every natural realization of the module in $F_{j-1}$. Moreover, the constants in all the maps are determined explicitly.
10.4.4. Comparison with the ideal generated by sub-permanents. Let $E, F=\mathbb{C}^{n}, V=E \otimes F$, let $I_{\kappa}^{\text {perm }_{n}, \kappa} \subset S^{\kappa}(E \otimes F)$ denote the span of the sub-permanents of size $\kappa$ and let $I^{\text {perm }_{n}, \kappa} \subset \operatorname{Sym}(E \otimes F)$ denote the ideal it generates. Note that $\operatorname{dim}\left(I_{\kappa}^{\text {perm }}, \kappa\right)=\binom{n}{\kappa}^{2}$. Fix complete flags $0 \subset E_{1} \subset \cdots \subset E_{n}=E$ and $0 \subset F_{1} \subset \cdots \subset F_{n}=F$. Write $\mathfrak{S}_{E_{j}}$ for the copy of $\mathfrak{S}_{j}$ acting on $E_{j}$ and similarly for $F$.

Write $T_{E} \subset S L(E)$ for the maximal torus (diagonal matrices). Recall from Theorem 6.6.2.2 that $G_{\text {perm }_{n}}$ is $\left[\left(T_{E} \times \mathfrak{S}_{E}\right) \times\left(T_{F} \times \mathfrak{S}_{F}\right)\right] \rtimes \mathbb{Z}_{2}$, divided by the $n$-th roots of unity.

Introduce the notation $\tilde{\mathfrak{S}}_{\kappa}=\mathfrak{S}_{\kappa} \times \mathfrak{S}_{n-\kappa} \subset \mathfrak{S}_{n}$, and if $\pi$ is a partition of $\kappa$, write $\widetilde{[\pi]}=[\pi] \times[n-\kappa]$ for the $\tilde{\mathfrak{S}}_{\kappa}$-module that is $[\pi]$ as an $\mathfrak{S}_{\kappa}$-module and trivial as an $\mathfrak{S}_{n-\kappa}$-module. For finite groups $H \subset G$, and an $H$-module $W, \operatorname{Ind}_{H}^{G} W=\mathbb{C}[G] \otimes_{\mathbb{C}[H]} W$ is the induced $G$-module, which has dimension
equal to $(\operatorname{dim} W)|G| /|H|$ (see, e.g, $[\mathbf{F H 9 1}, \S 3.4]$ ), and that $\operatorname{dim}[\pi]$ is given by the hook-length formula (see, e.g., [FH91, p. 50]). These two facts give the dimensions asserted below.

As an $\mathfrak{S}_{E_{n}} \times \mathfrak{S}_{F_{n}}$-module the space $I_{\kappa}^{\text {perm }_{n}, \kappa}$ decomposes as (10.4.5)
$\left.\operatorname{Ind} d_{\tilde{\mathfrak{S}}_{E_{\kappa}} \times \tilde{\mathfrak{S}}_{F_{\kappa}}}^{\mathcal{S}_{F_{n}}}{\widetilde{[\kappa]_{E_{\kappa}}}}_{E_{\kappa}}^{\otimes[\kappa]_{F_{\kappa}}}\right)=\left([n]_{E} \oplus[n-1,1]_{E} \oplus \cdots \oplus[n-\kappa, \kappa]_{E}\right) \otimes\left([n]_{F} \oplus[n-1,1]_{F} \oplus \cdots \oplus[n-\kappa, \kappa]_{F}\right)$.
The space of linear syzygies $M_{2, \kappa+1}:=\operatorname{ker}\left(I_{\kappa}^{\text {perm }_{n}, \kappa} \otimes V \rightarrow S^{\kappa+1} V\right)$ is the $\mathfrak{S}_{E_{n}} \times \mathfrak{S}_{F_{n}}$-module

This module has dimension $2 \kappa\binom{n}{\kappa+1}^{2}$. A spanning set for it may be obtained geometrically as follows: for each size $\kappa+1$ sub-matrix, perform the permanental tensor Laplace expansion along a row or column, then perform a second tensor Laplace expansion about a row or column and take the difference. An independent set of such for a given size $\kappa+1$ sub-matrix may be obtained from the expansions along the first row minus the expansion along the $j$-th for $j=2, \ldots, \kappa+1$, and then from the expansion along the first column minus the expansion along the $j$-th, for $j=2, \ldots, \kappa+1$.

Remark 10.4.4.1. Compare this with the space of linear syzygies for the determinant, which has dimension $\frac{2 \kappa(n+1)}{n-\kappa}\binom{n}{\kappa+1}^{2}$. The ratio of their sizes is $\frac{n+1}{n-\kappa}$, so, e.g., when $\kappa \sim \frac{n}{2}$, the determinant has about twice as many linear syzygies, and if $\kappa$ is close to $n$, one gets nearly $n$ times as many.
Theorem 10.4.4.2. [ELSW15] $\operatorname{dim} M_{j+1, \kappa+j}=\binom{n}{\kappa+j}^{2}\left({ }^{2(\kappa+j-1)}{ }_{j}\right)$. As an $\mathfrak{S}_{E} \times \mathfrak{S}_{F}$-module,

$$
\begin{equation*}
M_{j+1, \kappa+j}=\operatorname{Ind} d_{\tilde{\mathfrak{S}}_{E_{\kappa+j}} \times \tilde{\mathfrak{S}}_{F_{\kappa+j}} \times \mathfrak{S}_{F_{n}}}\left(\bigoplus_{a+b=j}\left[\widetilde{\kappa+b, 1^{a}}\right]_{E_{\kappa+j}} \otimes\left[\widetilde{\kappa+a, 1^{b}}\right]_{F_{\kappa+j}}\right) . \tag{10.4.6}
\end{equation*}
$$

The $\binom{n}{\kappa+j}^{2}$ is just the choice of a size $\kappa+j$ submatrix, the $\binom{2(\kappa+j-1)}{j}$ comes from choosing a set of $j$ elements from the set of rows union columns. Naïvely there are $\binom{2(\kappa+j)}{j}$ choices but there is redundancy as with the choices in the description of $M_{2}$.

The proof proceeds in two steps. As described below, one first gets "for free" the minimal free resolution of the ideal generated by $S^{\kappa} E \otimes S^{\kappa} F$. Write the generating modules of this resolution as $\tilde{M}_{j}$. We then locate the generators of the linear strand of the minimal free resolution of our ideal, whose generators we denote $M_{j+1, \kappa+j}$, inside $\tilde{M}_{j+1, \kappa+j}$ and prove the assertion.

To obtain $\tilde{M}_{j+1}$, use the Howe-Young endofunctor mentioned in $\S 8.11 .1 .1$ that comes from the involution on symmetric functions that takes the Schur function $s_{\pi}$ to $s_{\pi^{\prime}}$ (see, e.g. [Mac95, §I.2]). This involution extends to an endofunctor of $G L(V)$-modules and hence of $G L(E) \times G L(F)$-modules, taking $S_{\lambda} E \otimes S_{\mu} F$ to $S_{\lambda^{\prime}} E \otimes S_{\mu^{\prime}} F$ (see [AW07, §2.4]). This is only true as long as the dimensions of the vector spaces are sufficiently large, so to properly define it one passes to countably infinite dimensional vector spaces.

Applying this functor to the resolution (10.4.1), one obtains the resolution of the ideal generated by $S^{\kappa} E \otimes S^{\kappa} F \subset S^{\kappa}(E \otimes F)$. The $G L(E) \times G L(F)-$ modules generating the linear component of the $j$-th term in this resolution are:

$$
\begin{align*}
\tilde{M}_{j, j+\kappa-1} & =\bigoplus_{a+b=j-1} S_{\left(a, 1^{\kappa+b}\right)^{\prime},} E \otimes S_{\left(b, 1^{\kappa+a}\right)^{\prime}} F  \tag{10.4.7}\\
& =\bigoplus_{a+b=j-1} S_{\left(\kappa+b+1,1^{a-1}\right)} E \otimes S_{\left(\kappa+a+1,1^{b-1}\right)} F .
\end{align*}
$$

Moreover, by Corollary 10.4.3.3 and functoriality, the map from $S_{\left(\kappa+b+1,1^{a-1}\right)} E \otimes S_{\left(\kappa+a+1,1^{b-1}\right)} F$ into $\tilde{M}_{j-1, j+\kappa-1}$ is non-zero to the copies of $S_{\left(\kappa+b+1,1^{a-1}\right)} E \otimes S_{\left(\kappa+a+1,1^{b-1}\right)} F$ in
$\left(S_{\kappa+b, 1^{a-1}} E \otimes S_{\kappa+a+1,1^{b-2} F}\right) \cdot(E \otimes F)$ and $\left(S_{\kappa+b+1,1^{a-2}} E \otimes S_{\kappa+a, 1^{b-1} F}\right) \cdot(E \otimes F)$, when $a, b>0$.

Inside $S^{\kappa} E \otimes S^{\kappa} F$ is the ideal generated by the sub-permanents (10.4.5) which consists of the regular weight spaces $\left(p_{1}, \ldots, p_{n}\right) \times\left(q_{1}, \ldots, q_{n}\right)$, where all $p_{i}, q_{j}$ are either zero or one. (Each sub-permanent has such a weight, and, given such a weight, there is a unique sub-permanent to which it corresponds.) The set of regular vectors in any $E^{\otimes m} \otimes F^{\otimes m}$ spans a $\mathfrak{S}_{E} \times \mathfrak{S}_{F^{-}}$ submodule.

The linear strand of the $j$-the term in the minimal free resolution of the ideal generated by (10.4.5) is thus a $\mathfrak{S}_{E} \times \mathfrak{S}_{F}$-submodule of $\tilde{M}_{j, j+\kappa-1}$. We claim this sub-module is the span of the regular vectors. In other words:
Lemma 10.4.4.3. [ELSW15] $M_{j+1, \kappa+j}=\left(\tilde{M}_{j+1, \kappa+j}\right)_{r e g}$.
For the proof, see [ELSW15]. Theorem 10.4.4.2 follows because if $\pi$ is a partition of $\kappa+j$ then the weight $(1, \ldots, 1)$ subspace of $S_{\pi} E_{\kappa+j}$, considered as an $\mathfrak{S}_{E_{\kappa+j}}$-module, is [ $\pi$ ] by Theorem 9.2.2.1, and the space of regular vectors in $S_{\pi} E \otimes S_{\mu} F$ is $\left.\operatorname{In} d_{\tilde{\mathfrak{G}}_{E_{\kappa+j}} \times \mathfrak{G}_{F} \times \tilde{\mathfrak{G}}_{F_{\kappa+j}}}(\widetilde{[\pi]}]_{E} \otimes \widetilde{[\mu]}\right)$.
10.4.5. Proof of Theorem 10.4.1.1. The variety $\hat{\sigma}_{r}$ admits a desingularization by the geometric method of [Wey03], namely consider the Grassmannian $G\left(r, E^{*}\right)$ and the vector bundle $p: \mathcal{S} \otimes F \rightarrow G\left(r, E^{*}\right)$ whose fiber over $x \in G\left(r, E^{*}\right)$ is $x \otimes F$. (Although we are breaking symmetry here, it
will be restored in the end.) The total space admits the interpretation as the incidence variety

$$
\left\{(x, \phi) \in G\left(r, E^{*}\right) \times \operatorname{Hom}\left(F, E^{*}\right) \mid \phi(F) \subseteq x\right\}
$$

and the projection to $\operatorname{Hom}\left(F, E^{*}\right)=E^{*} \otimes F^{*}$ has image $\hat{\sigma}_{r}$. One also has the exact sequence

$$
0 \rightarrow \mathcal{S} \otimes F^{*} \rightarrow \underline{E^{*} \otimes F^{*}} \rightarrow \mathcal{Q} \otimes F^{*} \rightarrow 0
$$

where $E^{*} \otimes F^{*}$ denotes the trivial bundle with fiber $E^{*} \otimes F^{*}$ and $\mathcal{Q}=\underline{E^{*}} / \mathcal{S}$ is the quotient bundle. As explained in [Wey03], letting $q: \mathcal{S} \otimes F^{*} \rightarrow E^{*} \otimes F^{*}$ denote the projection, $q$ is a desingularization of $\hat{\sigma}_{r}$, the higher direct images $\mathcal{R}_{i} q^{*}\left(\mathcal{O}_{\mathcal{S} \otimes F^{*}}\right)$ are zero for $i>0$, and so by [Wey03, Thms. 5.12,5.13] one concludes $F_{i}=M_{i} \cdot \operatorname{Sym}(E \otimes F)$ where

$$
\begin{aligned}
M_{i} & =\bigoplus_{j \geq 0} H^{j}\left(G\left(r, E^{*}\right), \Lambda^{i+j}\left(\mathcal{Q}^{*} \otimes F\right)\right) \\
& =\bigoplus_{j \geq 0} \bigoplus_{|\pi|=i+j} H^{j}\left(G\left(r, E^{*}\right), S_{\pi} \mathcal{Q}^{*}\right) \otimes S_{\pi^{\prime}} F
\end{aligned}
$$

One now uses the Bott-Borel-Weil theorem to compute these cohomology groups. Homogeneous vector bundles on the Grassmannian $G(r, n)$ are indexed by sequences $\left(k_{1}, \ldots, k_{n}\right)$ where $k_{1} \geq k_{2} \cdots \geq k_{r}$ and $k_{r+1} \geq k_{r+2} \geq$ $\cdots \geq k_{n}$. An algorithm for implementing the Bott-Borel-Weil theorem is given in [Wey03, Rem. 4.1.5]: If $\pi=\left(p_{1}, \ldots, p_{q}\right)$ (where we must have $p_{1} \leq n$ to have $S_{\pi^{\prime}} F$ non-zero, and $q \leq n-r$ as $\left.\operatorname{rank} \mathcal{Q}=n-r\right)$, then $S_{\pi} \mathcal{Q}^{*}$ is the vector bundle corresponding to the sequence

$$
\begin{equation*}
\left(0^{r}, p_{1}, \ldots, p_{n-r}\right) \tag{10.4.8}
\end{equation*}
$$

The dotted Weyl action by $\sigma_{i}=(i, i+1) \in \mathfrak{S}_{n}$ is

$$
\sigma_{i} \cdot\left(\alpha_{1}, \ldots, \alpha_{n}\right)=\left(\alpha_{1}, \ldots, \alpha_{i-1}, \alpha_{i+1}-1, \alpha_{i}+1, \alpha_{i+2}, \ldots, \alpha_{n}\right)
$$

and one applies simple reflections to try to transform $\alpha$ to a partition until one either gets a partition after $u$ simple reflections, in which case the cohomology group $H^{u}\left(G\left(r, E^{*}\right), S_{\pi} \mathcal{Q}^{*}\right)$ is equal to the module associated to the partition one ends up with and all other cohomology groups are zero, or one ends up on a wall of the Weyl chamber, i.e., at one step one has $\left(\beta_{1}, \ldots, \beta_{n}\right)$ with some $\beta_{i+1}=\beta_{i}+1$, in which case there is no cohomology.

In our case, we need to move $p_{1}$ over to the first position in order to obtain a partition, which means we need $p_{1} \geq r+1$, and then if $p_{2}<2$ we are done, otherwise we need to move it etc... The upshot is we can get cohomology only if there is an $s$ such that $p_{s} \geq r+s$ and $p_{s+1}<s+1$, in which case we get

$$
S_{\left(p_{1}-r, \ldots, p_{s}-r, s^{r}, p_{s+1}, \ldots, p_{n-r}\right)} E \otimes S_{\pi^{\prime}} F
$$

contributing to $H^{r s}\left(G\left(r, E^{*}\right), S_{\pi} \mathcal{Q}^{*}\right)$. Say we are in this situation, then write $\left(p_{1}-r-s, \ldots, p_{s}-r-s\right)=\alpha,\left(p_{s+1}, \ldots, p_{n-r}\right)=\beta^{\prime}$, so

$$
\left(p_{1}-r, \ldots, p_{s}-r, s^{r}, p_{s+1}, \ldots, p_{n-r}\right)=\left(s^{r+s}\right)+\left(\alpha, 0^{r}, \beta^{\prime}\right)
$$

and moreover we may write

$$
\pi^{\prime}=\left(s^{r+s}\right)+\left(\beta, 0^{r}, \alpha^{\prime}\right)
$$

proving Theorem 10.4.1.1. The case $s=1$ gives the linear strand of the resolution.

## Hints and Answers to Selected Exercises

## Chapter 1.

1.1.15.1 In general, the trilinear map associated to a bilinear form is $(u, v, \gamma) \mapsto$ $\gamma(T(u, v))$. Let $z_{v}^{* u}$ denote the linear form that eats a matrix and returns its $(u, v)$-th entry. Since $(X Y)_{k}^{i}=\sum_{j} x_{j}^{i} y_{k}^{j}$, the associated trilinear map is $\left(X, Y, z_{v}^{* u}\right) \mapsto \sum_{j} x_{j}^{u} y_{v}^{j}$. On the other hand, $\operatorname{trace}(X Y Z)=\sum_{i, j, k} x_{j}^{i} y_{k}^{j} z_{i}^{k}$. Now observe that both these agree, e.g., on basis vectors.

## Chapter 2.

2.1.1.4 For the second assertion, a generic matrix will have nonzero determinant. For the last assertion, first say $\operatorname{rank}(f)=r^{\prime} \leq r$ and let $v_{1}, \ldots, v_{\mathbf{v}}$ be a basis of $V$ such that the kernel is spanned by the last $\mathbf{v}-r^{\prime}$ vectors. Then the matrix representing $f$ will be nonzero only in the upper $r^{\prime} \times r^{\prime}$ block and thus all minors of size greater than $r^{\prime}$ will be zero. Next say $\operatorname{rank}(f)=s>r$. In the same manner, we see the upper right size $s$ submatrix will have a nonzero determinant. Taking a Laplace expansion, we see at least one size $r+1$ minor of it is nonzero. In any other choice of basis minors expressed in the new basis are linear combinations of minors expressed in the old, so we conclude. For the last assertion, since polynomials are continuous, if $f_{t}$ is in the zero set of all size $(r+1)$ minors, so will its continous limit $f$.
2.1.1.5 $v \in V$ goes to the map $\beta \mapsto \beta(v)$.
2.1.2.1 A multi-linear map is determined by its action on bases of $A_{1}^{*}, \ldots, A_{n}^{*}$.
2.1.2.4 See (4.1.1).
2.1.5.2 Write an arbitrary rank two tensor as $\left(\alpha_{1} a_{1}+\alpha_{2} a_{2}\right) \otimes\left(\beta_{1} b_{1}+\beta_{2} b_{2}\right) \otimes\left(\gamma_{1} c_{1}+\right.$ $\left.\gamma_{1} c_{2}\right)+\left(\alpha_{1}^{\prime} a_{1}+\alpha_{2}^{\prime} a_{2}\right) \otimes\left(\beta_{1}^{\prime} b_{1}+\beta_{2}^{\prime} b_{2}\right) \otimes\left(\gamma_{1}^{\prime} c_{1}+\gamma_{1}^{\prime} c_{2}\right)$ where the Greek letters are arbitrary constants and show they cannot be chosen to make the tensor equal to $a_{1} \otimes b_{1} \otimes c_{2}+a_{1} \otimes b_{2} \otimes c_{1}+a_{2} \otimes b_{1} \otimes c_{1}$.
2.1.5.4 See §3.1.6.
2.1.6.1 For example, take $a_{1} \otimes b_{1} \otimes c_{2}+a_{1} \otimes b_{2} \otimes c_{1}+a_{2} \otimes b_{1} \otimes c_{1}+\sum_{j=3}^{r} a_{j} \otimes b_{j} \otimes c_{j}$.
2.1.6.2 If $T=\sum_{i=1}^{r} a_{i} \otimes b_{i} \otimes c_{i}$, then, letting $\pi_{A}: A \rightarrow A /\left(A^{\prime}\right)^{\perp}$ be the projection, and similarly for $B, C$, then $T_{A^{\prime} \otimes B^{\prime} \otimes C^{\prime}}=\sum_{i=1}^{r} \pi_{A}\left(a_{i}\right) \otimes \pi_{B}\left(b_{i}\right) \otimes \pi\left(c_{i}\right)$.
2.1.7.2 First assume $\underline{\mathbf{R}}(T)=\mathbf{R}(T)$ and write $T=a_{1} \otimes b_{1} \otimes c_{1}+\cdots+$ $a_{r} \otimes b_{r} \otimes c_{r}$. Then $T\left(A^{*}\right)=\operatorname{span}\left\{b_{1} \otimes c_{1}, \ldots, b_{r} \otimes c_{r}\right\}$ so $\underline{\mathbf{R}}(T) \geq \operatorname{rank} T_{A}$. Now use that ranks of linear maps are determined by polynomials (the minors of the entries) to conclude.
2.2.1.2 Say $T=\sum_{j=1}^{\mathrm{b}} a_{j} \otimes b_{j} \otimes c_{j}$ and this is an optimal expression. Since $T_{A}$ is injective, the $a_{j}$ must be a basis. Let $\alpha^{j}$ be the dual basis, so $T\left(\alpha^{j}\right)=$ $b_{j} \otimes c_{j}$ has rank one. These span. In the other direction, say the image is $\operatorname{span}\left\{b_{1} \otimes c_{1}, \ldots, b_{\mathbf{b}} \otimes c_{\mathbf{b}}\right\}$. then for each $j$ there must be some $\alpha^{j} \in A^{*}$ with $T\left(\alpha^{j}\right)=b_{j} \otimes c_{j}$. Since $T_{A}$ is injective, these form a basis of $A$, so we must have $T=\sum_{j=1}^{\mathbf{b}} a_{j} \otimes b_{j} \otimes c_{j}$ with $a_{j}$ the dual basis vectors.
2.2.2.2 Use Exercise 2.1.7.4, taking three matrices in $A^{*}$, e.g. Id, a matrix with all 1's just below the diagonal and zero elsewhere and a matrix with 1's just above the diagonal and zeros elsewhere.
2.3.3.2 First assume $T=e_{I}=e_{i_{1}} \wedge \cdots \wedge e_{i_{k}}$ and take $\mu=e^{L}$ and $\zeta=e^{J}$. Then

$$
\begin{aligned}
& \mu\lrcorner T=\left\{\begin{array}{c}
e_{I \backslash L} \\
0 \text { if } L \not \subset I
\end{array} \quad \text { if } \quad L \subset I\right. \\
& \zeta\lrcorner T=\left\{\begin{array}{c}
e^{J \backslash I} \\
0 \text { if } I \not \subset J
\end{array} \quad \text { if } \quad I \subset J\right.
\end{aligned}
$$

and $\left\langle e^{J \backslash I}, e_{I \backslash L}\right\rangle=0$, because they have no indices in common. By linearity we get zero for any linear combination of such $e^{J}, e_{L}$ 's so we see that $G(k, V)$ is in the zero set of the equations. (Any element of $G(k, V)$ is equivalent to [ $e_{I}$ ] after a change of basis and our equations are independent of the choice of basis.)
Now for simplicity assume $T=e_{I_{1}}+e_{I_{2}}$ where $I_{1}, I_{2}$ have at least two indices different. Take $\zeta=e^{I_{1} \cup F}$ where $F \subset I_{2}, F \not \subset I_{1}$ and $I_{2} \not \subset I_{1} \cup F$. Then $\zeta\lrcorner T=e^{F}$. Take $\mu=e^{I_{2} \backslash F}$ so $\left.\mu\right\lrcorner T=e_{F}$. We conclude.
Any element of $\Lambda^{k} V$ not in $\hat{G}(k, V)$ can be degenerated to be of the form $T$, so we conclude in general.
2.4.2.2 Show that for $X \in \Lambda^{p-1} A \otimes B^{*}, T_{A}^{\wedge p}(a \wedge X)=-a \wedge T_{A}^{\wedge p-1}(X)$.

It is sufficient to consider the case $q=p-1$. Say $X \in \operatorname{ker}\left(T_{A}^{\wedge p-1}\right)$. Then $a \wedge X \in \operatorname{ker}\left(T_{A}^{\wedge p}\right)$ so $a \wedge X=0$ for all $a \in A$. But this is not possible.
2.5.1.4 trace $(f)$.
2.5.1.5 Use Exercise 2.5.1.2.
2.5.2.2 Extend the $a_{j}$ to a basis of $A$ and consider the induced basis of $\Lambda^{q+1} A$. Write out $X_{j} \wedge a_{j}$ with respect to the induced basis and compare coefficients.
2.5.2.3 Use a variant of Lemma 2.5.2.1.
2.5.3.2 Apply the proof of Theorem 2.5.2.6 to $M_{\langle p, p, 2\rangle}$.

## Chapter 3.

3.1.4.2 By the first part of the exercise, every point on the Chow variety is a projection of a point of the form $v_{1} \otimes \cdots \otimes v_{d}$, for some $v_{j} \in V$, but the projection of $v_{1} \otimes \cdots \otimes v_{d}$ is $v_{1} \cdots v_{d}$.
3.1.4.3 The ideal is generated by $p_{3}^{2}-p_{2} p_{4}, p_{2}^{2}-p_{0} p_{4}$. Note that we simply are throwing away the polynomials with $p_{1}$. The point $p_{3}$, corresponding to the polynomial $x^{3} y$ is on a tangent line to $v_{4}\left(\mathbb{P}^{1}\right)$, while the point $p_{22}$, corresponding to the polynomial $x^{2} y^{2}$ is not.
3.1.4.5 The ideal is generated by $p_{2}^{2}-p_{1} p_{3}, p_{1} p_{2}-p_{0} p_{3}, p_{1}^{2}-p_{0} p_{2}$.
3.1.4.8 Say $f(X)=Z_{1} \cup Z_{2}$ and note that $X=f^{-1}\left(Z_{1}\right) \cup f^{-1}\left(Z_{2}\right)$.
3.2.1.4 Recall from Exercise 2.5.1.9 that $\otimes_{j} M_{\left\langle\mathbf{1}_{j}, \mathbf{m}_{j}, \mathbf{n}_{j}\right\rangle}=M_{\left\langle\Pi_{j} \mathbf{l}_{j}, \Pi_{k} \mathbf{m}_{k}, \Pi_{l} \mathbf{n}_{l}\right\rangle}$. Set $N=\mathbf{n m l}$ and consider $M_{\langle N\rangle}=M_{\langle\mathbf{m}, \mathbf{n}, \mathbf{l}\rangle} \otimes M_{\langle\mathbf{n}, \mathbf{l}, \mathbf{m}\rangle} \otimes M_{\langle\mathbf{l}, \mathbf{m}, \mathbf{n}\rangle}$.
3.2.2.1 Consider

$$
\left(\begin{array}{ll}
0 & 0 \\
0 & \infty \\
0 & \infty
\end{array}\right)
$$

3.3.1.3 Since the border rank of points in $G L(A) \times G L(B) \times G L(C) \cdot T$ equals the border rank of $T$, the border rank of points in the closure cannot increase.
3.4.9.3 Instead of the curve $a_{0}+t a_{1}$ use $a_{0}+t a_{1}+t^{2} a_{q+1}$ and similarly for $b, c$.
3.4.6.3 Use Proposition 3.2.1.7.
3.5.3.3 If $\mathbf{R}_{h}(T)=\underline{\mathbf{R}}(T)$, then when writing $T=\lim _{t \rightarrow 0} T(t)$, we may take $t \in \mathbb{Z}_{h+1}$.
3.5.3.4 If we are multiplying polynomials of degrees $d_{1}$ and $d_{2}$, then their product has degree $d_{1} d_{2}$, so the answer is the same as if we were working over $\mathbb{Z}_{d_{1} d_{2}}$.

## Chapter 4.

4.1.1.1 If one uses the images of the standard basis vectors, one gets:
$M_{\langle 2\rangle}=\left(\begin{array}{ll}0 & -1 \\ 1 & -1\end{array}\right)^{\otimes 3}+\left(\begin{array}{ll}1 & 0 \\ 0 & 0\end{array}\right)^{\otimes 3}+\left(\begin{array}{ll}0 & 1 \\ 0 & 1\end{array}\right)^{\otimes 3}+\left(\begin{array}{cc}0 & 0 \\ -1 & 1\end{array}\right)^{\otimes 3}+\left\langle\left(\begin{array}{ll}0 & 1 \\ 0 & 0\end{array}\right) \otimes\left(\begin{array}{ll}0 & 0 \\ 1 & 0\end{array}\right) \otimes\left(\begin{array}{ll}1 & -1 \\ 1 & -1\end{array}\right)\right\rangle_{\mathbb{Z}_{3}}$.
4.2.2.2 $\left(1,-\frac{1}{2}\right),\left(-\frac{1}{2}, 1\right),\left(-\frac{1}{2},-\frac{1}{2}\right)$.
4.5.2.1 $\left(\begin{array}{lll}* & 0 & 0 \\ * & * & * \\ 0 & 0 & *\end{array}\right)$.
4.7.5.1 If a line goes through $[a \otimes b \otimes c]$, then it must be contained in $\mathbb{P} \hat{T}_{[a \otimes b \otimes c]} S e g(\mathbb{P} A \times$ $\mathbb{P} B \times \mathbb{P} C)$.

## Chapter 5.

5.1.4.5 First note that if $x$ is generic, it is diagonalizable with distinct eigenvalues so if $x$ is generic, then $\operatorname{dim} C(x)=\mathbf{b}$. Then observe that $\operatorname{dim}(C(x))$ is semi-continuous as the set $\{y \mid \operatorname{dim} C(y) \leq p\}$ is an algebraic variety. Alternatetively, and more painfully, compute the centralizer of elements in Jordan canonical form.
5.2.1.1 See the proof of Proposition 5.2.1.2 in the case $k=1$.
5.3.1.4 For the lower bound use Koszul flattenings, for the upper, write $T$ as the sum of the first AFT tensor and the remainder and bound the border rank of each.
5.3.1.8 For the lower bound, use the substitution method. For the upper, consider the rank decomposition of the structure tensor of $\mathbb{C}\left[\mathbb{Z}_{2 \mathbf{m}-1}\right]$, which, using the DFT, has rank and border rank $\mathbf{m}$. Show that this tensor degenerates to the tensor corresponding to the centralizer of a regular nilpotent element.
5.4.3.5 $\overline{G \cdot x}$ is a union of orbits, so the boundary is a union of orbits all of dimension strictly less than $\operatorname{dim}(G \cdot x)$.
5.4.5.2 $\binom{\mathbf{n}+j-2}{j-1}=\operatorname{dim} S^{j-1} \mathbb{C}^{\mathbf{n}-1}$ so the sum may be thought of as computing the dimension of $S^{m-1} \mathbb{C}^{\mathbf{n}}$ where each summand represents basis vectors (monomials) where e.g., $x_{1}$ appears to the power $m-j$.
5.4.5.3 Without loss of generality assume $2 \leq i \leq j$. For $j=2,3$ the inequality is straightforward to check, so assume $j \geq 4$. Prove the inequality 5.4.3 by induction on $\mathbf{n}$. For $\mathbf{n}=i j$ the inequality follows from the combinatorial interpretation of binomial coefficients and the fact that the middle one is the largest.

We have $\binom{\mathbf{n}+1-1+i j-1}{i j-1}=\binom{\mathbf{n}-1+i j-1}{i j-1} \frac{\mathbf{n}-1+i j}{\mathbf{n}},\binom{\mathbf{n}+1-j+i-1}{i-1}=\binom{\mathbf{n}-j+i-1}{i-1} \frac{\mathbf{n}-j+i}{n-j+1}$ and $\binom{\mathbf{n}+1-i+j-1}{j-1}=\binom{\mathbf{n}-i+j-1}{j-1} \frac{\mathbf{n}-i+j}{\mathbf{n}-i+1}$. By induction it is enough to prove that:

$$
\begin{equation*}
\frac{\mathbf{n}-1+i j}{\mathbf{n}} \geq \frac{\mathbf{n}-j+i}{\mathbf{n}-j+1} \frac{\mathbf{n}-i+j}{\mathbf{n}-i+1} . \tag{10.4.9}
\end{equation*}
$$

This is equivalent to:

$$
i j-1 \geq \frac{\mathbf{n}(i-1)}{\mathbf{n}-j+1}+\frac{\mathbf{n}(j-1)}{\mathbf{n}-i+1}+\frac{\mathbf{n}(i-1)(j-1)}{(\mathbf{n}-j+1)(\mathbf{n}-i+1)} .
$$

As the left hand side is independent from $\mathbf{n}$ and each fraction on the right hand side decreases with growing $\mathbf{n}$, we may set $\mathbf{n}=i j$ in inequality 10.4.9. Thus it is enough to prove:

$$
2-\frac{1}{i j} \geq\left(1+\frac{i-1}{i j-j+1}\right)\left(1+\frac{j-1}{i j-i+1}\right) .
$$

Then the inequality is straightforward to check for $i=2$, so assume $i \geq 3$.
5.4.4.3 For any $z \in v_{n}(\operatorname{Seg}(\mathbb{P} E \times \mathbb{P} F)), G_{\operatorname{det}_{n}, z}$, the group preserving both $\operatorname{det}_{n}$ and $z$, is isomorphic to $P_{E} \times P_{F}$, where $P_{E}, P_{F}$ are the parabolic subgroups of matrices with zero in the first column except the $(1,1)$-slot, and $z$ is in the $G_{\operatorname{det}_{n}, z}$-orbit closure of any $q \in v_{n}(\mathbb{P} W)$.
5.4.4.4 Notice that fixing $k=[(\mu \otimes v) \otimes(\nu \otimes w) \otimes(\omega \otimes u)]$ is equivalent to fixing a partial flag in each $U, V$ and $W$ consisting of a line and a hyperplane containing it. Let $[a \otimes b \otimes c] \in \operatorname{Seg}(\mathbb{P} A \times \mathbb{P} B \times \mathbb{P} C)$. If $[a] \notin \operatorname{Seg}\left(\mathbb{P} U^{*} \times \mathbb{P} V\right)$ then the orbit is not closed, even under the torus action on $V$ or $U^{*}$ that is compatible with the flag. So without loss of generality, we may assume $[a \otimes b \otimes c] \in S e g\left(\mathbb{P} U^{*} \times \mathbb{P} V \times \mathbb{P} V^{*} \times \mathbb{P} W \times \mathbb{P} W^{*} \times \mathbb{P} U\right)$. Write $a \otimes b \otimes c=$ $\left(\mu^{\prime} \otimes v^{\prime}\right) \otimes\left(\nu^{\prime} \otimes w^{\prime}\right) \otimes\left(\omega^{\prime} \otimes u^{\prime}\right)$. If, for example $v^{\prime} \neq v$, we may act with an element of $G L(V)$ that preserves the partial flag and sends $v^{\prime}$ to $v+\epsilon v^{\prime}$. Hence $v$ is in the closure of the orbit of $v^{\prime}$. As $G_{M_{\langle U, V, W\rangle}, k}$ preserves $v$ we may continue, reaching $k$ in the closure.
5.6.2.6 $\mathcal{A}$ has basis $x_{J}:=x_{1}^{j_{1}} \cdots x_{n}^{j_{n}}$ with $0 \leq j_{s}<a_{s}$. Let $e_{J}$ be the dual basis. Then $T_{\mathcal{A}}=\sum_{i_{s}+j_{s}<a_{s}} e_{I} \otimes e_{J} \otimes x_{I+J}$. Write $x_{K}^{*}=x_{1}^{a_{1}-k_{1}-1} \cdots x_{n}^{a_{n}-k_{n}-1}$. Then $T_{\mathcal{A}}=\sum_{i_{s}+j_{s}+k_{s}<a_{s}} e_{I} \otimes e_{J} \otimes e_{K}$.
5.6.3.1 Show that if $n \in \operatorname{Rad}(\mathcal{A})$ is not nilpotent, then there is some prime ideal of $\mathcal{A}$ not containing $n$.

## Chapter 6.

6.1.4.2 Use that $\frac{1}{1-\lambda t}=\sum_{j} \lambda^{j} t^{j}$.
6.2.2.4 Consider $\lim _{\epsilon \rightarrow 0} \frac{1}{\epsilon}\left((x+\epsilon y)^{n}-x^{n}\right)$.
6.2.2.7 Respectively, taking $k=\left\lfloor\frac{n}{2}\right\rfloor$ one gets the ranks are $\binom{n}{\left\lfloor\frac{n}{2}\right\rfloor},\binom{n}{\left\lfloor\frac{n}{2}\right\rfloor}^{2}$, and $\binom{n}{\left\lfloor\frac{n}{2}\right\rfloor}^{2}$.
6.2.3.1 $\quad N_{M}^{*} \sigma_{r}^{0}=\operatorname{ker} M \otimes(\operatorname{Image} M)^{\perp}=\operatorname{ker} M \otimes \operatorname{ker} M^{T} \subset U \otimes V^{*}$. The second equality holds because for a linear map $f: V \rightarrow W$, Image $(f)^{\perp}=$ $\operatorname{ker}\left(f^{T}\right)$.
6.3.3.3 The space of matrices with last two columns equal to zero is contained in $Z\left(\operatorname{perm}_{m}\right)_{\text {sing }}$.
6.3.4.5 Let $\hat{Q} \in S^{2} V$ be the corresponding quadratic form (defined up to scale). Take a basis $e_{1}, \ldots, e_{\mathbf{v}}$ of $V$ such that $e_{1}, \ldots, e_{k}$ correspond to a linear space on $Q$, so $Q\left(e_{s}, e_{t}\right)=0$ for $0 \leq s, t \leq k$. But $Q$ being smooth says $\hat{Q}$ is non-degenerate, so for each $e_{s}$, there must be some $e_{f(s)}$ with $Q\left(e_{s}, e_{f(s)}\right) \neq 0$.
6.4.2.3 Parametrize $C$ by a parameter $s$ and $\tau(C)$ by $s$ and a parameter $t$ for the line, then differentiate.
6.4.3.1 Consider a curve $([\bar{x}(t)],[\bar{H}(t)]) \in \mathcal{I}$. Note that $\langle\bar{x}(t), \bar{H}(t)\rangle \equiv 0$ where $\langle\cdot, \cdot\rangle: V \times V^{*} \rightarrow \mathbb{C}$ is the pairing. Now consider $\left.\frac{d}{d t}\right|_{t=0}\langle\bar{x}(t), \bar{H}(t)\rangle$.
6.4.6.1 First note that perm $_{m}$ evaluated on a matrix whose entries are all one is $m$ !. Then perform a permanental Laplace expansion about the first row.
6.5.2.2 Note that $\frac{\partial R}{\partial x_{i}}=\sum_{j} \frac{\partial^{2} R}{\partial x_{i} \partial x_{j}}$ and now consider the last nonzero column.
6.6.1.3 In this case the determinant is a smooth quadric.
6.6.1.5 Linear spaces on a variety $X$ through a point $x \in X$ must be contained in $\mathbb{P} \hat{T}_{x} X$.
6.6.2.1 $\left\{\operatorname{perm}_{2}=0\right\}$ is a smooth quadric.

## Chapter 7.

7.1.1.3 Consider (where blank entries are zero)

$$
\operatorname{det}\left(\begin{array}{ccccccc}
0 & & x_{1} & & x_{2} & & x_{3} \\
x_{1} & \ell & & & & & \\
& x_{1} & \ell & & & & \\
x_{2} & & & \ell & & & \\
& & & x_{2} & \ell & & \\
x_{3} & & & & & \ell & \\
& & & & & x_{3} & \ell
\end{array}\right)=\ell^{7-3}\left(x_{1}^{3}+x_{2}^{3}+x_{3}^{3}\right)
$$

7.4.2.7 Take $x_{m}=x_{m+1}=\cdots=x_{N}=0$.
7.4.2.9 Let $\omega$ be a primitive $q$-th root of unity. Let $x_{1}, \ldots, x_{N}$ denote the standard basis of $\mathbb{C}^{N}$. Consider the vector $\left(1, \omega, \omega^{2}, \ldots, \omega^{q-1}, 0, \ldots, 0\right)$ and its shifts by zeros.
7.6.1.3 In degree $d+\tau$, this ideal consists of all polynomials of the form $\ell_{1}^{d} Q_{1}+\ell_{2}^{d} Q_{2}$ with $Q_{1}, Q_{2} \in S^{\tau} \mathbb{C}^{n^{2}}$, which has dimension $2 \operatorname{dim} S^{\tau} \mathbb{C}^{n^{2}}-$
$\operatorname{dim} S^{\tau-(d)} \mathbb{C}^{n^{2}}$ because the polynomials of the form $\ell_{1}^{d} \ell_{2}^{d} Q_{3}$ with $Q_{3} \in S^{\tau-(d)} \mathbb{C}^{n^{2}}$ appear in both terms.

## Chapter 8.

8.1.2.2 Say we have a weight vector $z \in V^{\otimes d}$ weight $\left(j_{1}, \ldots, j_{\mathbf{v}}\right)$ with $j_{i}<$ $j_{i+1}$. Consider the matrix $g$ that is the identity plus a vector with one non-zero entry in the $(i, i+1)$ slot. Then $g z$ is a non-zero vector of weight $\left(j_{1}, \ldots, j_{i-1}, j_{i}+1, j_{i+1}-1, \ldots, j_{\mathbf{v}}\right)$.
8.1.4.1 The weight of the one-dimensional representation $\operatorname{det}^{-1}$ is $(-1, \ldots,-1)$.
8.1.4.2 Consider the linear form $v \mapsto \operatorname{det}_{\mathbf{v}}\left(v_{1}, \ldots, v_{\mathbf{v}-1}, v\right)$.
8.1.5.2 $g \cdot e_{1} \wedge \cdots \wedge e_{\mathbf{v}}=\operatorname{det}(g) e_{1} \wedge \cdots \wedge e_{\mathbf{v}}$
8.2.1.2 By linearity, for any $P_{1}, P_{2}$, the rank of the linear map $U^{*} \rightarrow W$ associated to $P_{1}+P_{2}$ is at most the sum of the ranks of the maps associated to $P_{1}$ and $P_{2}$.
8.4.1.2 A highest weight vector of any copy of $S_{\pi} V^{*}$ is constructed skewsymmetrizing over $l(\pi)$ vectors. For the other direction, the zero set of any $P \in S^{\delta}\left(S^{d} \mathbb{C}^{k}\right)$ is a proper subvariety of $S^{d} \mathbb{C}^{k}$.
8.5.3.4 Under the action of a basis vector in $\mathfrak{g l}(E \otimes F)$, since it is by Leibnitz rule, at most one variable in each monomial can be changed. So whatever highest weight vectors appear in the tangent space, their weight can differ by at most one in each of $E, F$ from $\left(\left(1^{n}\right),\left(1^{n}\right)\right)$. But there is only one partition pair with this property that occurs in $S^{n}(E \otimes F)$, namely $\left.\left(2,1^{n-1}\right),\left(2,1^{n-1}\right)\right)$.
8.6.8.2 We need $\operatorname{Hom}_{\mathfrak{S}_{d}}\left([\pi]^{*},[\mu]\right) \neq 0$. But $[\pi]^{*} \simeq[\pi]$. By Schur's lemma $\operatorname{Hom}_{\mathfrak{S}_{d}}([\pi],[\mu]) \neq 0$ if and only if $[\pi]=[\mu]$.
8.6.8.3 If the multiplicity were greater than one, $\pi$ would not be irreducible by Schur's lemma

### 8.6.1.1 Prove an algebra version of Schur's lemma.

8.6.4.2 If $V$ is an irreducible $G$-module, then $V^{*} \otimes V$ is an irreducible $G \times G$ module.)
8.7.1.3

$$
\begin{aligned}
S^{d}(E \otimes F) & =\left[(E \otimes F)^{\otimes d}\right]_{d} \\
& =\left(E^{\otimes d} \otimes F^{\otimes d}\right)^{\mathfrak{S}_{d}} \\
& \left.\left.=\left[\left(\oplus_{|\pi|=d} S_{\pi} E \otimes[\pi]\right)\right) \otimes\left(\oplus_{|\mu|=d} S_{\mu} F \otimes[\mu]\right)\right)\right]^{\mathfrak{S}_{d}} \\
& =\oplus_{|\mu|,|\pi|=d} S_{\pi} E \otimes S_{\mu} F \otimes([\pi] \otimes[\mu])^{\mathfrak{S}_{d}}
\end{aligned}
$$

Now use Exercise 8.6.8.2.
8.7.2.2 $c_{\pi^{\prime}}=\sum_{\sigma \in \mathfrak{S}_{\pi^{\prime}}} \delta_{\sigma} \sum_{\sigma \in \mathfrak{S}_{\pi}} \operatorname{sgn}(\sigma) \delta_{\sigma}$. Now show $c_{\left(1^{d}\right)} c_{\pi}=c_{\pi^{\prime}}$.
8.11.1.2 $\sum_{j=0}^{m}\binom{m}{j}^{2}=\binom{2 m}{m}$ as $\Lambda^{m}(E \oplus F)=\sum_{j} \Lambda^{j} E \otimes \Lambda^{m-j} F$.

## Chapter 9.

9.1.2.2 Highest weight vectors here correspond to partitions with at most $d$ parts.
9.5.2.1 The map $M a t_{n} \rightarrow M a t_{n} / / G L_{n}$ sends a matrix to the coefficients of its characteristic polynomial, i.e., the elementary symmetric functions of its eigenvalues.
9.5.2.2 Say $P \in \mathbb{C}(\operatorname{Nor}(Z / / \Gamma))$, satifies a monic polynomial with coefficients in $\mathbb{C}[\operatorname{Nor}(Z / / \Gamma)]$. Note that $\mathbb{C}(\operatorname{Nor}(Z / / \Gamma)) \subset \mathbb{C}(Z)^{\Gamma}$, and of course $\mathbb{C}[\operatorname{Nor}(Z / / \Gamma)] \subset \mathbb{C}[Z]$, so by the normality of $Z, P \in \mathbb{C}[Z]$, but $\mathbb{C}[\operatorname{Nor}(Z / / \Gamma)]=\mathbb{C}[Z] \cap \mathbb{C}(Z)^{\Gamma}$.
9.6.3.2 By the Pieri formula, one can have at most three parts. On the other hand, $C h_{d}\left(\mathbb{C}^{2}\right)=\mathbb{P} S^{d} \mathbb{C}^{2}$.

## Chapter 10.

10.1.1.2 For $P \in S^{e} V^{*}, P(x)=0$ is equivalent to $\left\langle\bar{P}, x^{e}\right\rangle=0$.
10.1.4.3 Consider $R$ with $I_{R}=\left(x_{1}^{d_{1}+1}, \ldots, x_{n-1}^{d_{n-1}+1}\right)$.

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[^0]:    ${ }^{1}$ To this day, it is not known if there is an even more efficient algorithm than the FFT. See [Val77, Lok08, KLPSMN09, GHIL].

[^1]:    ${ }^{2}$ This phrase is due to Howard Karloff

