# NOTES FOR MATH 847 (REPRESENTATION STABILITY) 

STEVEN V SAM

## Contents

1. Introduction ..... 1
2. Representation theory ..... 3
3. Representations of combinatorial categories ..... 26
4. Homological stability for symmetric groups ..... 33
5. Representation stability for configuration spaces ..... 38
6. Algebraic geometry from tensors ..... 41
7. More on FI-modules ..... 49
8. $\Delta$-modules ..... 60
Appendix A. Review of homological algebra ..... 64
References ..... 68

## 1. Introduction

1.1. Prerequisites. Prior familiarity with homological algebra will be useful but not essential. I'll review the key points as we go along, but will not give the proofs of basic facts. For details, any textbook on modern algebra should suffice.

This is an example-driven theory, and we will go through some examples from group theory, representation theory, topology, and algebraic geometry. I also will not assume these as prerequisites, and again will just state the facts that we'll need as they come up.
1.2. Motivation. We will study two aspects of representation stability in this course:

- A generalization of homological stability: one is given a sequence of abelian groups $/$ vector spaces $V_{1}, V_{2}, \ldots$ together with homomorphisms $V_{1} \xrightarrow{f_{1}} V_{2} \xrightarrow{f_{2}} \cdots$. These groups may arise as the (co)homology of some family of objects, and the phenomena of interest is when $f_{i}$ are isomorphisms for $i \gg 0$.

As an example, one might consider the homology of symmetric groups $\Sigma_{n}$, and more specifically the first homology group. In general, for any group $G, \mathrm{H}_{1}(G)=$ $G /[G, G]$, the abelianization of $G([G, G]$ is the subgroup of $G$ generated by elements of the form $\left.x y x^{-1} y^{-1}\right)$. The standard inclusions $\Sigma_{1} \rightarrow \Sigma_{2} \rightarrow \cdots$ induces maps $\mathrm{H}_{1}\left(\Sigma_{1}\right) \rightarrow \mathrm{H}_{1}\left(\Sigma_{2}\right) \rightarrow \cdots$. While $\mathrm{H}_{1}\left(\Sigma_{1}\right)$ is trivial, $\mathrm{H}_{1}\left(\Sigma_{i}\right) \cong \mathbf{Z} / 2$ for $i \geq 2$, and one can check $f_{i}$ is an isomorphism for $i \geq 2$.

For another example, one might take $X_{n}$ to be the space of unordered $n$-tuples of distinct points in the plane and consider its homology groups.

- A method for exploiting symmetry to prove uniformity statements across a family of examples. We'll develop this more later, but for an example in algebraic geometry, consider a tensor product of vector spaces $V_{1} \otimes \cdots \otimes V_{n}$. Every element is a sum of simple tensors (those of the form $v_{1} \otimes \cdots \otimes v_{n}$ for $v_{i} \in V_{i}$ ); define the rank of an element to be $r$ if $r$ is the least number of simple tensor summands needed. Define the border rank to be $\leq r$ if an element can be arbitrarily approximated by rank $r$ tensors. When $n=2$, this all reduces to the usual notion of rank for matrices.

In general, it is hard to determine the (border) rank of a tensor, but a theorem of Draisma-Kuttler that we will discuss says that one can test this by the vanishing of polynomials of degree at most $C(r)$, where $C(r)$ depends only on $r$, but not on $n$ or $\operatorname{dim} V_{i}$.

Remark 1.2.1. Why is tensor rank interesting? Here's one example. Multiplication of $n \times n$ matrices can be represented by a bilinear map $\mathbf{k}^{n^{2}} \otimes \mathbf{k}^{n^{2}} \rightarrow \mathbf{k}^{n^{2}}$, which is the same as a tensor in $\left(\mathbf{k}^{n^{2}}\right)^{*} \otimes\left(\mathbf{k}^{n^{2}}\right)^{*} \otimes \mathbf{k}^{n^{2}}$. The rank of this tensor is the minimum number of multiplications (in $\mathbf{k}$ ) needed to perform matrix multiplication (see [La, §1.1, 1.2]). The naive algorithm gives an upper bound of $n^{3}$, but one can do better. In particular, for $n=2$, the rank is actually 7 , while the precise value is already unknown for $n=3$ (it's at least 19 and at most 23).

The formula is a bit complicated, but the idea can be illustrated with multiplication of complex numbers. The standard way gives four real multiplications:

$$
(a+b i)(c+d i)=(a c-b d)+(b c+a d) i
$$

Alternatively, we have

$$
\begin{array}{r}
(a+b i)(c+d i)=\left(k_{1}-k_{3}\right)+\left(k_{1}+k_{2}\right) i \\
k_{1}=c(a+b), \quad k_{2}=a(d-c), \quad k_{3}=b(c+d),
\end{array}
$$

so you can get away with only 3 real multiplications at the cost of more additions. However, in some scenarios, minimizing the number of multiplications at the cost of increasing additions is desirable (for example, doing things by hand or with floating point approximation).

To elaborate on the first point, assume the $V_{i}$ are finite-dimensional vector spaces over a field $\mathbf{k}$ for simplicity. An obvious consequence of stability for the sequence $V_{1} \rightarrow V_{2} \rightarrow \cdots$ is that the sequence $\left(\operatorname{dim} V_{i}\right)_{i \geq 0}$ is eventually constant.

One of the first classes of examples of representation stable sequences is to consider a sequence $V_{i}$ as before, where now $V_{i}$ is a representation of the symmetric group $\Sigma_{i}$. The transition maps $f_{i}: V_{i} \rightarrow V_{i+1}$ should be $\Sigma_{i}$-equivariant, where we embed $\Sigma_{i} \subset \Sigma_{i+1}$ by having it fix $i+1$, and they should also obey another condition which we omit for now.

The right generalization of stability is finite generation. For $i \leq j$, set $f_{j, i}=f_{j} \circ f_{j-1} \circ \cdots \circ f_{i}$. Say that $x_{1}, \ldots, x_{d}$ with $x_{i} \in V_{n_{i}}$ generate the sequence $V$ if every $x \in V_{n}$ can be written as a $\mathbf{k}\left[\Sigma_{n}\right]$-linear combination of the $f_{n, n_{i}}\left(x_{i}\right)$.

Example 1.2.2. (1) Let $V_{n}=\mathbf{k}^{n}$ with $\Sigma_{n}$ acting by permutations on the coordinates. The map $f_{n}: \mathbf{k}^{n} \rightarrow \mathbf{k}^{n+1}$ is the inclusion $\left(a_{1}, \ldots, a_{n}\right) \mapsto\left(a_{1}, \ldots, a_{n}, 0\right)$. The element $1 \in V_{1}$ generates the sequence.
(2) Let $V_{n}=\mathbf{k}^{\binom{n}{2}}$ with a basis consisting of 2-element subsets of $\{1, \ldots, n\}$ with the usual permutation action. The maps $f_{n}$ are the obvious ones. Then $e_{\{1,2\}} \in V_{2}$ generates. More generally, we can do the same with $k$-element subsets of $\{1, \ldots, n\}$.
(3) Let $V_{n}=\mathbf{k}^{n^{2}}$ with a basis consisting of ordered pairs of elements in $\{1, \ldots, n\}$ with the usual permutation action. Then $e_{(1,1)} \in V_{1}$ and $e_{(1,2)} \in V_{2}$ gives a generating set.

Finite generation implies a number of things about the representations $V_{i}$, but one consequence avoiding representation theory is that the sequence $\left(\operatorname{dim} V_{i}\right)_{i \geq 0}$ grows like a polynomial function for $i \gg 0$. This definition of sequences of $\Sigma_{n}$-representations is a bit clumsy, and we'll see later that the language of FI-modules will give a more natural notion to work with.

There are many sources of examples of these kinds of sequences. Related to the above example, we can take $X_{n}$ to be the space of ordered $n$-tuples of distinct points in the plane and consider its homology groups. Each of these carries a natural action of the symmetric group $\Sigma_{n}$.

The representation theory of symmetric groups is a well-developed theory, especially over fields of characteristic 0 . This will be the focus for the first part of the course.

## 2. Representation theory

We won't go through most of the proofs for this part, just enough to get a feeling for how to work these objects. Some of this material is found in more detail in my symmetric functions notes:
http://www.math.wisc.edu/~svs/740/notes.pdf
2.1. Schur-Weyl duality. We will need to know that the irreducible complex representations of $\Sigma_{n}$ are indexed by partitions of size $n$. There are several ways to get to this, but we will use Schur-Weyl duality to motivate the indexing.

Let $A$ be a finite-dimensional algebra over a field $\mathbf{k}$. Recall that an $A$-module $M$ is simple if its only submodules are 0 and $M$ itself. An $A$-module $M$ is semisimple if, for every submodule $N \subset M$, there exists another submodule $N^{\prime} \subset M$ such that $N \cap N^{\prime}=0$ and $N+N^{\prime}=M$, i.e., the smallest submodule containing both $N$ and $N^{\prime}$ is $M$. We denote this relation by $M=N \oplus N^{\prime}$. We say that $A$ is semisimple if all of its finite-dimensional modules are semisimple.

Note that if $M$ is semisimple and finite-dimensional, then we can decompose $M$ as a direct sum of simple modules.

Let $\mathrm{M}_{n}(\mathbf{k})$ be the algebra of $n \times n$ matrices with entries in $\mathbf{k}$. Given a module $M$ in general, we let $\operatorname{End}(M)$ denote the ring of linear endomorphisms (self-maps) of $M$. We won't prove the next statement, but it can be found in many abstract algebra textbooks.

Lemma 2.1.1. For any $n_{1}, \ldots, n_{r}$, the product $\prod_{i=1}^{r} \mathrm{M}_{n_{i}}(\mathbf{k})$ is a semisimple algebra. If $\mathbf{k}$ is algebraically closed, then every semisimple finite-dimensional algebra is of this form where $r$ is the number of isomorphism classes of irreducible representations, and the $n_{i}$ are their dimensions.

Given a group $G$, the group algebra $\mathbf{k}[G]$ is the vector space with basis $\left\{e_{g} \mid g \in G\right\}$ with multiplication $e_{g} e_{h}=e_{g h}$. A $\mathbf{k}[G]$-module is the same thing as a linear representation of $G$, so we will use these perspectives interchangeably.

Lemma 2.1.2. Let $\mathbf{k}$ be a field, $G$ be a finite group, and suppose that $|G|$ is invertible in $\mathbf{k}$. Then $\mathbf{k}[G]$ is semisimple.

Proof. Let $M$ be a representation of $G$ and let $N \subset M$ be a submodule. Choose any linear projection $\pi: M \rightarrow N$, i.e., $\pi(x)=x$ for all $x \in N$. Define a new map $\psi: M \rightarrow N$ by

$$
\psi(m)=\frac{1}{|G|} \sum_{g \in G} g \psi\left(g^{-1} m\right)
$$

Then

- $\psi(x)=x$ for all $x \in N$ : given $x \in N, g^{-1} x \in N$ since $N$ is a submodule, and $\psi\left(g^{-1} x\right)=g^{-1} x$, so the sum simplifies to $x$,
- $\psi$ is $G$-equivariant, i.e., $h \psi(m)=\psi(h m)$ for all $h \in G$ and $m \in M$ :

$$
h \psi(m)=\frac{1}{|G|} \sum_{g \in G} h g \psi\left(g^{-1} m\right)=\frac{1}{|G|} \sum_{g^{\prime} \in G} g^{\prime} \psi\left(\left(g^{\prime}\right)^{-1} h m\right)=\psi(h m)
$$

where in the second equality, we do the change of indexing $g^{\prime}=h g$.
Let $N^{\prime}=\operatorname{ker} \psi$. Then $M=N \oplus N^{\prime}$.
We will be concerned with $G=\Sigma_{n}$ the symmetric group, and $\mathbf{k}=\mathbf{C}$, the field of complex numbers.

Given an algebra $A$ acting on a vector space $E$, let

$$
B=\operatorname{End}_{A}(E)=\{\varphi: E \rightarrow E \mid \varphi(a e)=a \varphi(e) \text { for all } a \in A \text { and } e \in E\}
$$

Proposition 2.1.3 (Double centralizer theorem). Let $E$ be a finite-dimensional vector space. Suppose $A \subset \operatorname{End}(E)$ is a semisimple subalgebra.
(1) $A=\operatorname{End}_{B}(E)$.
(2) $B$ is semisimple.
(3) As a representation of $A \times B, E$ decomposes as

$$
E=\bigoplus_{i \in I} V_{i} \otimes W_{i}
$$

where the $V_{i}$ are all of the irreducible $A$-modules and the $W_{i}$ are all of the irreducible $B$-modules. In particular, the correspondence $V_{i} \leftrightarrow W_{i}$ is a bijection between the irreducible modules of $A$ and $B$.

Proof. Since $A$ is semisimple, we can decompose $E$ into a direct sum of simple $A$-modules:

$$
E=\bigoplus_{i \in I} V_{i}^{\oplus c_{i}}=\bigoplus_{i \in I} V_{i} \otimes W_{i}
$$

where $W_{i}=\operatorname{Hom}_{A}\left(V_{i}, E\right)$ is the multiplicity space of $V_{i}$. Then $B=\bigoplus_{i \in I} \operatorname{End}\left(W_{i}\right)$, so it is semisimple, and the simple $B$-modules are the $W_{i}$. It is also clear from this decomposition that $A=\bigoplus_{i \in I} \operatorname{End}\left(V_{i}\right)=\operatorname{End}_{B}(E)$, and so $A$ is semisimple by Lemma 2.1.1.

Now we come to our example of interest. Let $V=\mathbf{C}^{n}$ and let $E=V^{\otimes d}$. Then $\Sigma_{d}$ acts on $E$ by permuting the tensor factors, i.e.,

$$
\sigma \cdot \sum v_{1} \otimes \cdots \otimes v_{d}=\sum v_{\sigma^{-1}(1)} \otimes \cdots \otimes v_{\sigma^{-1}(d)}
$$

In particular, this defines a homomorphism $\mathbf{C}\left[\Sigma_{d}\right] \rightarrow \operatorname{End}(E)$, and we let $A$ be the image. By definition, quotients of semisimple algebras are still semisimple, so $A$ is semisimple by Lemma 2.1.2.

There is also an action of $\mathbf{G L}(V)=\mathbf{G} \mathbf{L}_{n}(\mathbf{C})$ (the group of invertible linear transformations of $V$ ) by the natural action of linear change of coordinates, i.e.,

$$
g \cdot \sum v_{1} \otimes \cdots \otimes v_{d}=\sum\left(g v_{1}\right) \otimes \cdots \otimes\left(g v_{d}\right)
$$

This gives a group homomorphism $\mathbf{G L}(V) \rightarrow \mathbf{G L}(E)$, and we let $B$ be the algebra generated by the image. Evidently, we have $B \subseteq \operatorname{End}_{A}(E)$ and $A \subseteq \operatorname{End}_{B}(E)$.

Proposition 2.1.4 (Schur-Weyl duality). $B=\operatorname{End}_{A}(E)$. In particular, $A=\operatorname{End}_{B}(E)$ and we have a decomposition

$$
E=\bigoplus_{i \in I} V_{i} \otimes W_{i}
$$

where the $V_{i}$ are all of the irreducible modules for $A$ (and hence a subset of the irreducible representations of $\Sigma_{d}$ ) and the $W_{i}$ are all of the irreducible modules for $B$ (and hence a subset of the irreducible representations of $\mathbf{G L}(V)$ ).

Proof. We have an identification $\operatorname{End}\left(V^{\otimes d}\right)=\operatorname{End}(V)^{\otimes d}$, and $\operatorname{End}_{A}(E)$ is identified with the symmetric elements in $\operatorname{End}(V)^{\otimes d}$. By the next lemma, $\operatorname{End}_{A}(E)$ is spanned by the elements $\varphi^{\otimes d}=\varphi \otimes \cdots \otimes \varphi$ where $\varphi \in \operatorname{End}(V)$. We claim that $\varphi^{\otimes d} \in B$. Note that tid $+\varphi$ is invertible for all but finitely many values of $t$. In particular, $(\operatorname{tid}+\varphi)^{\otimes d} \in B$. Linear spaces are closed (in the Euclidean topology, or Zariski topology if you prefer), so we can take the limit $t \rightarrow 0$ to conclude that $\varphi^{\otimes d} \in B$.

The remainder of the proposition follows from the double centralizer theorem.
Lemma 2.1.5. Let $W$ be a vector space over a field of characteristic 0 . Then the $\Sigma_{d^{-}}$ invariant elements in $W^{\otimes d}$ are linearly spanned by elements of the form $w^{\otimes d}=w \otimes \cdots \otimes w$ for $w \in W$.

Proof. Pick a basis $x_{1}, \ldots, x_{n}$ for $W$. The subspace spanned by the $w^{\otimes d}$ is a $\mathbf{G L}(W)$ invariant subspace of the symmetric invariants of $W^{\otimes d}$, which we can identify with the space $\operatorname{Sym}^{d}(W)$ of homogeneous polynomials in variables $x_{1}, \ldots, x_{n}$, so it suffices to show that any nonzero $\mathbf{G L}(W)$-invariant subspace of $\operatorname{Sym}^{d}(W)$ is the whole space. This has a standard basis consisting of monomials.

So let $U$ be any $\mathbf{G L}(W)$-invariant subspace of $\operatorname{Sym}^{d}(W)$. First, consider a diagonal matrix in $\mathbf{G L}(W)$ with generic entries. Using a Vandermonde-type argument, one can show that if $f \in U$, then all of the monomials with nonzero coefficient in $f$ are also in $U$. Now, given any monomial $x_{1}^{d_{1}} \cdots x_{n}^{d_{n}} \in U$, we can conclude that any other monomial is also in $U$. We illustrate this with an example. Let $g \in \mathbf{G L}(W)$ be the matrix which fixes $x_{2}, \ldots, x_{n}$ and sends $x_{1}$ to $x_{1}+x_{2}$. Then $\left(x_{1}+x_{2}\right)^{d_{1}} x_{2}^{d_{2}} \cdots x_{n}^{d_{n}} \in U$ and $x_{1}^{d_{1}-1} x_{2}^{d_{2}+1} \cdots x_{n}^{d_{n}}$ has a nonzero coefficient (this is where characteristic 0 is important) and hence also belongs to $U$. By lowering and raising exponents appropriately, we can get to any other monomial.

Remark 2.1.6. The statement about $\operatorname{Sym}^{d}(W)$ is not correct in positive characteristic. For example, let $W=\mathbf{k}^{n}$ for a field of characteristic $p$. Then the space spanned by $x_{1}^{p}, \ldots, x_{n}^{p}$ in $\operatorname{Sym}^{p}(W)$ is closed under $\mathbf{G L}(W)$. However, it cannot be identified with the subspace of $\Sigma_{p}$-invariants in $W^{\otimes p}$ in a $\mathbf{G L}(W)$-equivariant way. This subspace of $\Sigma_{p}$-invariants is also known as the divided power of $W$. With some more care, the statement about divided powers can be proven in more generality, but we won't need it.

The next question is to identify the indexing set $I$ in the decomposition given by SchurWeyl duality. We'll use some information about $\mathbf{G L}_{n}(\mathbf{C})$ to do this, and go via symmetric functions.
2.2. Symmetric functions. Let $x_{1}, \ldots, x_{n}$ be a finite set of indeterminates. The symmetric group on $n$ letters is $\Sigma_{n}$. It acts on $\mathbf{Z}\left[x_{1}, \ldots, x_{n}\right]$, the ring of polynomials in $n$ variables and integer coefficients, by substitution of variables:

$$
\sigma \cdot f\left(x_{1}, \ldots, x_{n}\right)=f\left(x_{\sigma(1)}, \ldots, x_{\sigma(n)}\right) .
$$

The ring of symmetric polynomials is the set of fixed polynomials:

$$
\Lambda(n):=\left\{f \in \mathbf{Z}\left[x_{1}, \ldots, x_{n}\right] \mid \sigma \cdot f=f \text { for all } \sigma \in \Sigma_{n}\right\} .
$$

This is a subring of $\mathbf{Z}\left[x_{1}, \ldots, x_{n}\right]$.
We will also treat the case $n=\infty$. Let $x_{1}, x_{2}, \ldots$ be a countably infinite set of indeterminates. Let $\Sigma_{\infty}$ be the group of all permutations of $\{1,2, \ldots\}$. Let $R$ be the ring of power series in $x_{1}, x_{2}, \ldots$ of bounded degree. Hence, elements of $R$ can be infinite sums, but only in a finite number of degrees.

Write $\pi_{n}: \Lambda \rightarrow \Lambda(n)$ for the homomorphism which sets $x_{n+1}=x_{n+2}=\cdots=0$.
Remark 2.2.1. (For those familiar with inverse limits.) There is a ring homomorphism $\pi_{n+1, n}: \Lambda(n+1) \rightarrow \Lambda(n)$ obtained by setting $x_{n+1}=0$. Furthermore, $\Lambda(n)=\bigoplus_{d \geq 0} \Lambda(n)_{d}$ where $\Lambda(n)_{d}$ is the subgroup of homogeneous symmetric polynomials of degree $d$. The map $\pi_{n+1, n}$ restricts to a map $\Lambda(n+1)_{d} \rightarrow \Lambda(n)_{d}$; set

$$
\Lambda_{d}={\underset{\gtrless}{\underset{n}{2}}}_{\lim _{n}} \Lambda(n)_{d} .
$$

Then

$$
\Lambda=\bigoplus_{d \geq 0} \Lambda_{d}
$$

Note that we aren't saying that $\Lambda$ is the inverse limit of the $\Lambda(n)$; the latter object includes infinite sums of unbounded degree.

Then $\Sigma_{\infty}$ acts on $R$, and we define the ring of symmetric functions

$$
\Lambda:=\left\{f \in R \mid \sigma \cdot f=f \text { for all } \sigma \in \Sigma_{\infty}\right\} .
$$

Again, this is a subring of $R$.
Example 2.2.2. Here are some basic examples of elements in $\Lambda$ (we will study them more soon):

$$
\begin{aligned}
p_{k} & :=\sum_{i \geq 1} x_{i}^{k} \\
e_{k} & :=\sum_{i_{1}<i_{2}<\cdots<i_{k}} x_{i_{1}} x_{i_{2}} \cdots x_{i_{k}} \\
h_{k} & :=\sum_{i_{1} \leq i_{2} \leq \cdots \leq i_{k}} x_{i_{1}} x_{i_{2}} \cdots x_{i_{k}} .
\end{aligned}
$$

Sometimes, we want to work with rational coefficients instead of integer coefficients. In that case, we'll write $\Lambda_{\mathbf{Q}}$ or $\Lambda(n)_{\mathbf{Q}}$ to denote the appropriate rings.
2.3. Polynomial representations of general linear groups. Let $\mathbf{G L}_{n}(\mathbf{C})$ denote the group of invertible $n \times n$ complex matrices.

A polynomial representation of $\mathbf{G L}_{n}(\mathbf{C})$ is a homomorphism $\rho: \mathbf{G L}_{n}(\mathbf{C}) \rightarrow \mathbf{G L}(V)$ where $V$ is a $\mathbf{C}$-vector space, and the entries of $\rho$ can be expressed in terms of polynomials (as soon as we pick a basis for $V$ ).

A simple example is the identity map $\rho: \mathbf{G L}_{n}(\mathbf{C}) \rightarrow \mathbf{G L}_{n}(\mathbf{C})$. Slightly more sophisticated is $\rho: \mathbf{G L}_{2}(\mathbf{C}) \rightarrow \mathbf{G L}\left(\operatorname{Sym}^{2}\left(\mathbf{C}^{2}\right)\right)$ where $\operatorname{Sym}^{2}\left(\mathbf{C}^{2}\right)$ is the space of degree 2 polynomials in $x, y$ (which is a basis for $\mathbf{C}^{2}$ ). The homomorphism can be defined by linear change of coordinates, i.e.,

$$
\rho(g)\left(a x^{2}+b x y+c y^{2}\right)=a(g x)^{2}+b(g x)(g y)+c(g y)^{2} .
$$

If we pick the basis $x^{2}, x y, y^{2}$ for $\operatorname{Sym}^{2}\left(\mathbf{C}^{2}\right)$, this can be written in coordinates as

$$
\mathbf{G L}_{2}(\mathbf{C}) \rightarrow \mathbf{G L}_{3}(\mathbf{C})
$$

$$
\left(\begin{array}{cc}
g_{1,1} & g_{1,2}  \tag{2.3.1}\\
g_{2,1} & g_{2,2}
\end{array}\right) \mapsto\left(\begin{array}{ccc}
g_{1,1}^{2} & g_{1,1} g_{1,2} & g_{1,2}^{2} \\
2 g_{1,1} g_{2,1} & g_{1,1} g_{2,2}+g_{1,2} g_{2,1} & 2 g_{1,2} g_{2,2} \\
g_{2,1}^{2} & g_{2,1} g_{2,2} & g_{2,2}^{2}
\end{array}\right) .
$$

More generally, we can define $\rho: \mathbf{G L}_{n}(\mathbf{C}) \rightarrow \mathbf{G L}\left(\operatorname{Sym}^{d}\left(\mathbf{C}^{n}\right)\right)$ for any $n, d$. Another important example uses exterior powers instead of symmetric powers, so we have $\rho: \mathbf{G L}_{n}(\mathbf{C}) \rightarrow$ $\mathbf{G L}\left(\wedge^{d}\left(\mathbf{C}^{n}\right)\right)$.

An important invariant of a polynomial representation $\rho$ is its character: define

$$
\operatorname{char}(\rho)\left(x_{1}, \ldots, x_{n}\right):=\operatorname{Tr}\left(\rho\left(\operatorname{diag}\left(x_{1}, \ldots, x_{n}\right)\right)\right)
$$

where $\operatorname{diag}\left(x_{1}, \ldots, x_{n}\right)$ is the diagonal matrix with entries $x_{1}, \ldots, x_{n}$ and $\operatorname{Tr}$ denotes trace.
Lemma 2.3.2. $\operatorname{char}(\rho)\left(x_{1}, \ldots, x_{n}\right) \in \Lambda(n)$.
Proof. Each permutation $\sigma \in \Sigma_{n}$ corresponds to a permutation matrix $M(\sigma)$ : this is the matrix with a 1 in row $\sigma(i)$ and column $i$ for $i=1, \ldots, n$ and 0 's everywhere else. Then

$$
M(\sigma)^{-1} \operatorname{diag}\left(x_{1}, \ldots, x_{n}\right) M(\sigma)=\operatorname{diag}\left(x_{\sigma(1)}, \ldots, x_{\sigma(n)}\right) .
$$

Now use that the trace of a matrix is invariant under conjugation:

$$
\begin{aligned}
\operatorname{char}(\rho)\left(x_{1}, \ldots, x_{n}\right) & =\operatorname{Tr}\left(\rho\left(\operatorname{diag}\left(x_{1}, \ldots, x_{n}\right)\right)\right) \\
& =\operatorname{Tr}\left(\rho(M(\sigma))^{-1} \rho\left(\operatorname{diag}\left(x_{1}, \ldots, x_{n}\right)\right) \rho(M(\sigma))\right) \\
& =\operatorname{Tr}\left(\rho\left(M(\sigma)^{-1} \operatorname{diag}\left(x_{1}, \ldots, x_{n}\right) M(\sigma)\right)\right) \\
& =\operatorname{Tr}\left(\rho\left(\operatorname{diag}\left(x_{\sigma(1)}, \ldots, x_{\sigma(n)}\right)\right)\right) \\
& =\operatorname{char}(\rho)\left(x_{\sigma(1)}, \ldots, x_{\sigma(n)}\right) .
\end{aligned}
$$

Example 2.3.3. - The character of the identity representation is $x_{1}+x_{2}+\cdots+x_{n}$.

- The character of the representation $\rho: \mathbf{G L}_{n}(\mathbf{C}) \rightarrow \mathbf{G L}\left(\operatorname{Sym}^{d}\left(\mathbf{C}^{n}\right)\right)$ is

$$
h_{n}\left(x_{1}, \ldots, x_{n}\right)=\sum_{1 \leq i_{1} \leq \cdots \leq i_{d} \leq n} x_{i_{1}} \cdots x_{i_{d}} .
$$

- The character of the representation $\rho: \mathbf{G L}_{n}(\mathbf{C}) \rightarrow \mathbf{G L}\left(\wedge^{d}\left(\mathbf{C}^{n}\right)\right)$ is

$$
e_{n}\left(x_{1}, \ldots, x_{n}\right)=\sum_{1 \leq i_{1}<\cdots<i_{d} \leq n} x_{i_{1}} \cdots x_{i_{d}}
$$

A few remarks that aren't easy to see right now (though we may revisit):

- The set of characters in $\Lambda(n)$ generates all of $\Lambda(n)$ as an abelian group.
- If we are more careful about how to define characters in infinite-dimensional settings, we get that characters of polynomial representations of $\mathbf{G} \mathbf{L}_{\infty}(\mathbf{C})$ are elements of $\Lambda$.
- The character determines the representation up to isomorphism: if $\operatorname{char}(\rho)=\operatorname{char}\left(\rho^{\prime}\right)$, then $\rho$ and $\rho^{\prime}$ define isomorphic representations: one can be obtained from the other by a suitable isomorphism of the underlying vector spaces $V$ and $V^{\prime}$.

If we take these remarks as fact for now, this gives one motivation for studying $\Lambda$. This representation-theoretic interpretation of $\Lambda$ will clarify various definitions and constructions we will encounter. A few basic ones that we can see now:

- If $\rho_{i}: \mathbf{G L}_{n}(\mathbf{C}) \rightarrow \mathbf{G L}\left(V_{i}\right)$ are polynomial representations for $i=1,2$, we can form the direct sum representation $\rho_{1} \oplus \rho_{2}: \mathbf{G L}_{n}(\mathbf{C}) \rightarrow \mathbf{G L}\left(V_{1} \oplus V_{2}\right)$ via

$$
\left(\rho_{1} \oplus \rho_{2}\right)(g)=\left(\begin{array}{cc}
\rho_{1}(g) & 0 \\
0 & \rho_{2}(g)
\end{array}\right)
$$

and

$$
\operatorname{char}\left(\rho_{1} \oplus \rho_{2}\right)\left(x_{1}, \ldots, x_{n}\right)=\operatorname{char}\left(\rho_{1}\right)\left(x_{1}, \ldots, x_{n}\right)+\operatorname{char}\left(\rho_{2}\right)\left(x_{1}, \ldots, x_{n}\right)
$$

- There's also a multiplicative version using tensor product. If $\rho_{i}: \mathbf{G} \mathbf{L}_{n}(\mathbf{C}) \rightarrow \mathbf{G L}\left(V_{i}\right)$ are polynomial representations for $i=1,2$, we can form the tensor product representation $\rho_{1} \otimes \rho_{2}: \mathbf{G L}_{n}(\mathbf{C}) \rightarrow \mathbf{G L}\left(V_{1} \otimes V_{2}\right)$ via (assuming $\rho_{1}(g)$ is $\left.N \times N\right)$ :

$$
\left(\rho_{1} \otimes \rho_{2}\right)(g)=\left(\begin{array}{cccc}
\rho_{1}(g)_{1,1} \rho_{2}(g) & \rho_{1}(g)_{1,2} \rho_{2}(g) & \cdots & \rho_{1}(g)_{1, N} \rho_{2}(g) \\
\rho_{1}(g)_{2,1} \rho_{2}(g) & \rho_{1}(g)_{2,2} \rho_{2}(g) & \cdots & \rho_{1}(g)_{2, N} \rho_{2}(g) \\
\vdots & & & \\
\rho_{1}(g)_{N, 1} \rho_{2}(g) & \rho_{1}(g)_{N, 2} \rho_{2}(g) & \cdots & \rho_{1}(g)_{N, N} \rho_{2}(g)
\end{array}\right)
$$

(here we are multiplying $\rho_{2}(g)$ by each entry of $\rho_{1}(g)$ and creating a giant block matrix) and

$$
\operatorname{char}\left(\rho_{1} \otimes \rho_{2}\right)\left(x_{1}, \ldots, x_{n}\right)=\operatorname{char}\left(\rho_{1}\right)\left(x_{1}, \ldots, x_{n}\right) \cdot \operatorname{char}\left(\rho_{2}\right)\left(x_{1}, \ldots, x_{n}\right)
$$

Note that subtraction will not have any natural interpretation, and in general, the difference of two characters need not be a character. In general, the elements of $\Lambda(n)$ or $\Lambda$ can be thought of as virtual characters since every element is the difference of two characters.
2.4. Partitions. A partition of a nonnegative integer $n$ is a sequence $\lambda=\left(\lambda_{1}, \ldots, \lambda_{k}\right)$ such that $\lambda_{1} \geq \lambda_{2} \geq \cdots \geq \lambda_{k} \geq 0$ and $\lambda_{1}+\cdots+\lambda_{k}=n$. We will consider two partitions the same if their nonzero entries are the same. And for shorthand, we may omit the commas, so the partition $(1,1,1,1)$ of 4 can be written as 1111 . As a further shorthand, the exponential notation is used for repetition, so for example, $1^{4}$ is the partition $(1,1,1,1)$. We let $\operatorname{Par}(n)$ be the set of partitions of $n$, and denote the size by $p(n)=|\operatorname{Par}(n)|$. By convention, $\operatorname{Par}(0)$ consists of exactly one partition, the empty one.

## Example 2.4.1.

$$
\begin{aligned}
& \operatorname{Par}(1)=\{1\}, \\
& \operatorname{Par}(2)=\left\{2,1^{2}\right\}, \\
& \operatorname{Par}(3)=\left\{3,21,1^{3}\right\}, \\
& \operatorname{Par}(4)=\left\{4,31,22,21^{2}, 1^{4}\right\}, \\
& \operatorname{Par}(5)=\left\{5,41,32,31^{2}, 2^{2} 1,21^{3}, 1^{5}\right\} .
\end{aligned}
$$

If $\lambda$ is a partition of $n$, we write $|\lambda|=n$ (size). Also, $\ell(\lambda)$ is the number of nonzero entries of $\lambda$ (length). For each $i, m_{i}(\lambda)$ is the number of entries of $\lambda$ that are equal to $i$.

It will often be convenient to represent partitions graphically. This is done via Young diagrams, which is a collection of left-justified boxes with $\lambda_{i}$ boxes in row $i .{ }^{1}$ For example, the Young diagram

corresponds to the partition (5, 3, 2). Flipping across the main diagonal gives another partition $\lambda^{\dagger}$, called the transpose. In our example, flipping gives


So $(5,3,2)^{\dagger}=(3,3,2,1,1)$. In other words, the role of columns and rows has been interchanged. This is an important involution of $\operatorname{Par}(n)$ which we will use later.

We will use several different partial orderings of partitions:

- $\lambda \subseteq \mu$ if $\lambda_{i} \leq \mu_{i}$ for all $i$.
- The dominance order: $\lambda \leq \mu$ if $\lambda_{1}+\cdots+\lambda_{i} \leq \mu_{1}+\cdots+\mu_{i}$ for all $i$. Note that if $|\lambda|=|\mu|$, then $\lambda \leq \mu$ if and only if $\lambda^{\dagger} \geq \mu^{\dagger}$. So transpose is an order-reversing involution on the set of partitions of a fixed size.
- The lexicographic order: for partitions of the same size, $\lambda \leq^{R} \mu$ if $\lambda=\mu$, or otherwise, there exists $i$ such that $\lambda_{1}=\mu_{1}, \ldots, \lambda_{i-1}=\mu_{i-1}$, and $\lambda_{i}<\mu_{i}$. This is a total ordering.


### 2.5. Bases for $\Lambda$.

2.5.1. Monomial symmetric functions. Given a partition $\lambda=\left(\lambda_{1}, \lambda_{2}, \ldots\right)$, define the monomial symmetric function by

$$
m_{\lambda}=\sum_{\alpha} x^{\alpha}
$$

where the sum is over all distinct permutations $\alpha$ of $\lambda$. This is symmetric by definition. So for example, $m_{1}=\sum_{i \geq 1} x_{i}$ since all of the distinct permutations of $(1,0,0, \ldots)$ are integer

[^0]sequences with a single 1 somewhere and 0 elsewhere. By convention, $m_{0}=1$. Some other examples:
\[

$$
\begin{aligned}
m_{1,1} & =\sum_{i<j} x_{i} x_{j} \\
m_{3,2,1} & =\sum_{\substack{i, j, k \\
i \neq j, j \neq k, i \neq k}} x_{i} x_{j} x_{k} .
\end{aligned}
$$
\]

In general, $m_{1^{k}}=e_{k}$ and $m_{k}=p_{k}$.
Theorem 2.5.1. As we range over all partitions, the $m_{\lambda}$ form a basis for $\Lambda$.
Proof. They are linearly independent since no two $m_{\lambda}$ have any monomials in common. Clearly they also span: given $f \in \Lambda$, we can write $f=\sum_{\alpha} c_{\alpha} x^{\alpha}$ and $c_{\alpha}=c_{\beta}$ if both are permutations of each other, so this can be rewritten as $f=\sum_{\lambda} c_{\lambda} m_{\lambda}$ where the sum is now over just the partitions.

Corollary 2.5.2. $\Lambda_{d}$ has a basis given by $\left\{m_{\lambda}| | \lambda \mid=d\right\}$, and hence is a free abelian group of rank $p(d)=|\operatorname{Par}(d)|$.
Theorem 2.5.3. $\Lambda(n)_{d}$ has a basis given by $\left\{m_{\lambda}\left(x_{1}, \ldots, x_{n}\right)| | \lambda \mid=d, \ell(\lambda) \leq n\right\}$.
2.5.2. Elementary symmetric functions. Recall that we defined

$$
e_{k}=\sum_{i_{1}<i_{2}<\cdots<i_{k}} x_{i_{1}} x_{i_{2}} \cdots x_{i_{k}}
$$

For a partition $\lambda=\left(\lambda_{1}, \ldots, \lambda_{k}\right)$, define the elementary symmetric function by

$$
e_{\lambda}=e_{\lambda_{1}} e_{\lambda_{2}} \cdots e_{\lambda_{k}} .
$$

Note $e_{\lambda} \in \Lambda_{|\lambda|}$.
Theorem 2.5.4. The $e_{\lambda}$ form a basis of $\Lambda$.
Theorem 2.5.5. The set $\left\{e_{\lambda}\left(x_{1}, \ldots, x_{n}\right)\left|\lambda_{1} \leq n,|\lambda|=d\right\}\right.$ is a basis of $\Lambda(n)_{d}$.
Proof. If $\lambda_{1}>n$, then $e_{\lambda_{1}}\left(x_{1}, \ldots, x_{n}\right)=0$, so $e_{\lambda}\left(x_{1}, \ldots, x_{n}\right)=0$. Hence under the map $\pi_{n}: \Lambda \rightarrow \Lambda(n)$, the proposed $e_{\lambda}$ span the image. The number of such $e_{\lambda}$ in degree $d$ is $\left|\left\{\lambda\left|\lambda_{1} \leq n,|\lambda|=d\right\} \mid\right.\right.$, which is the same as $|\{\lambda|\ell(\lambda) \leq n,|\lambda|=d\} \mid$ via the transpose $\dagger$, and this is the rank of $\Lambda(n)_{d}$, so the $e_{\lambda}$ form a basis.

Remark 2.5.6. The previous two theorems say that the elements $e_{1}, e_{2}, e_{3}, \ldots$ are algebraically independent in $\Lambda$, and that the elements $e_{1}, \ldots, e_{n}$ are algebraically independent in $\Lambda(n)$. This is also known as the "fundamental theorem of symmetric functions".
2.5.3. Complete homogeneous symmetric functions. For a partition $\lambda=\left(\lambda_{1}, \ldots, \lambda_{k}\right)$, define the complete homogeneous symmetric functions by

$$
h_{\lambda}=h_{\lambda_{1}} \cdots h_{\lambda_{k}} .
$$

Theorem 2.5.7. The $h_{\lambda}$ form a basis for $\Lambda$.
Theorem 2.5.8. $h_{1}, \ldots, h_{n}$ are algebraically independent generators of $\Lambda(n)$, and the set $\left\{h_{\lambda}\left(x_{1}, \ldots, x_{n}\right)\left|\lambda_{1} \leq n,|\lambda|=d\right\}\right.$ is a basis of $\Lambda(n)_{d}$.
2.5.4. Power sum symmetric functions. Recall we defined

$$
p_{k}=\sum_{n \geq 1} x_{n}^{k}
$$

For a partition $\lambda=\left(\lambda_{1}, \ldots, \lambda_{k}\right)$, the power sum symmetric functions are defined by

$$
p_{\lambda}=p_{\lambda_{1}} \cdots p_{\lambda_{k}}
$$

Theorem 2.5.9. The $p_{\lambda}$ are linearly independent.
Remark 2.5.10. The $p_{\lambda}$ do not form a basis for $\Lambda$. For example, in degree 2, we have

$$
\begin{aligned}
p_{2} & =m_{2} \\
p_{1,1} & =m_{2}+2 m_{1,1}
\end{aligned}
$$

and the change of basis matrix has determinant 2 , so is not invertible over $\mathbf{Z}$. However, they do form a basis for $\Lambda_{\mathbf{Q}}$.
Theorem 2.5.11. $p_{1}, \ldots, p_{n}$ are algebraically independent generators of $\Lambda(n)_{\mathbf{Q}}$ and the set $\left\{p_{\lambda}\left(x_{1}, \ldots, x_{n}\right)\left|\lambda_{1} \leq n,|\lambda|=d\right\}\right.$ is a basis for $\Lambda(n)_{\mathbf{Q}, d}$.
2.5.5. The involution $\omega$. Since the $e_{i}$ are algebraically independent, we can define a ring homomorphism $f: \Lambda \rightarrow \Lambda$ by specifying $f\left(e_{i}\right)$ arbitrarily. ${ }^{2}$ Define

$$
\omega: \Lambda \rightarrow \Lambda
$$

by $\omega\left(e_{i}\right)=h_{i}$, where recall that $h_{k}=\sum_{i_{1} \leq \cdots \leq i_{k}} x_{i_{1}} \cdots x_{i_{k}}$.
Theorem 2.5.12. $\omega$ is an involution, i.e., $\omega^{2}=1$. Equivalently, $\omega\left(h_{i}\right)=e_{i}$.
Furthermore, we can define a finite analogue of $\omega$, the ring homomorphism $\omega_{n}: \Lambda(n) \rightarrow$ $\Lambda(n)$, given by $\omega_{n}\left(e_{i}\right)=h_{i}$ for $i=1, \ldots, n$.

Theorem 2.5.13. $\omega_{n}^{2}=1$, and $\omega_{n}$ is invertible. Equivalently, $\omega_{n}\left(h_{i}\right)=e_{i}$ for $i=1, \ldots, n$.
2.5.6. A scalar product. Define a bilinear form $\langle\rangle:, \Lambda \otimes \Lambda \rightarrow \mathbf{Z}$ by setting

$$
\left\langle m_{\lambda}, h_{\mu}\right\rangle=\delta_{\lambda, \mu}
$$

where $\delta$ is the Kronecker delta (1 if $\lambda=\mu$ and 0 otherwise). In other words, if $f=\sum_{\lambda} a_{\lambda} m_{\lambda}$ and $g=\sum_{\mu} b_{\mu} h_{\mu}$, then $\langle f, g\rangle=\sum_{\lambda} a_{\lambda} b_{\lambda}$ (well-defined since both $m$ and $h$ are bases). At this point, the definition looks completely unmotivated. However, this inner product is natural from the representation-theoretic perspective, which we'll mention in the next section (without proof).

In our setup, $m$ and $h$ are dual bases with respect to the pairing. We will want a general criteria for two bases to be dual to each other. To state this criterion, we need to work in two sets of variables $x$ and $y$ and in the ring $\Lambda \otimes \Lambda$ where the $x$ 's and $y$ 's are separately symmetric.
Lemma 2.5.14. Let $u_{\lambda}$ and $v_{\mu}$ be bases of $\Lambda$ (or $\Lambda_{\mathbf{Q}}$ ). Then $\left\langle u_{\lambda}, v_{\mu}\right\rangle=\delta_{\lambda, \mu}$ if and only if

$$
\sum_{\lambda} u_{\lambda}(x) v_{\lambda}(y)=\prod_{i, j}\left(1-x_{i} y_{j}\right)^{-1}
$$

[^1]Corollary 2.5.15. The pairing is symmetric, i.e., $\langle f, g\rangle=\langle g, f\rangle$.
Proof. The condition above is the same if we interchange $x$ and $y$, so $\left\langle m_{\lambda}, h_{\mu}\right\rangle=\left\langle h_{\mu}, m_{\lambda}\right\rangle$. Now use bilinearity.
Proposition 2.5.16.

$$
\begin{aligned}
\prod_{i, j}\left(1-x_{i} y_{j}\right)^{-1} & =\sum_{\lambda} m_{\lambda}(x) h_{\lambda}(y) \\
& =\sum_{\lambda} z_{\lambda}^{-1} p_{\lambda}(x) p_{\lambda}(y) .
\end{aligned}
$$

Proposition 2.5.17. $\omega$ is an isometry, i.e., $\langle f, g\rangle=\langle\omega(f), \omega(g)\rangle$.
The bilinear form $\langle$,$\rangle is positive definite, i.e., \langle f, f\rangle>0$ for $f \neq 0$.
2.5.7. Schur functions. Let $\lambda$ be a partition. A semistandard Young tableaux (SSYT) $T$ is an assignment of natural numbers to the Young diagram of $\lambda$ so that the numbers are weakly increasing going left to right in each row, and the numbers are strictly increasing going top to bottom in each column.

Example 2.5.18. If $\lambda=(4,3,1)$, and we have the assignment

$$
\begin{array}{|l|l|l|l}
\hline a & b & c & d \\
\hline e & f & g \\
\hline h & & \\
\hline
\end{array}
$$

then, in order for this to be a SSYT, we need to have

- $a \leq b \leq c \leq d$,
- $e \leq f \leq g$,
- $a<e<h$,
- $b<f$, and
- $c<g$.

An example of a SSYT is | 1 | 1 | 3 | 5 |
| :--- | :--- | :--- | :--- |
| 2 | 3 | 4 |  |
| 7 |  |  |  |

The type of a SSYT $T$ is the sequence type $(T)=\left(\alpha_{1}, \alpha_{2}, \ldots\right)$ where $\alpha_{i}$ is the number of times that $i$ appears in $T$. We set

$$
x^{T}=x_{1}^{\alpha_{1}} x_{2}^{\alpha_{2}} \cdots
$$

Given a pair of partitions $\mu \subseteq \lambda$, the Young diagram of $\lambda / \mu$ is the Young diagram of $\lambda$ with the Young diagram of $\mu$ removed. We define a SSYT of shape $\lambda / \mu$ to be an assignment of natural numbers of this Young diagram which is weakly increasing in rows and strictly increasing in columns.

Example 2.5.19. If $\lambda=(5,3,1)$ and $\mu=(2,1)$, then

is a SSYT if

- $a \leq b \leq c$,
- $d \leq e$, and
- $a<e$.

We define the type of $T$ and $x^{T}$ in the same way.
Given a partition $\lambda$, the Schur function $s_{\lambda}$ is defined by

$$
s_{\lambda}=\sum_{T} x^{T}
$$

where the sum is over all SSYT of shape $\lambda$. Similarly, given $\mu \subseteq \lambda$, the skew Schur function $s_{\lambda / \mu}$ is defined by

$$
s_{\lambda / \mu}=\sum_{T} x^{T}
$$

where the sum is over all SSYT of shape $\lambda / \mu$. Note that this is a strict generalization of the first definition since we can take $\mu=\varnothing$, the unique partition of 0 .

We can make the same definitions in finitely many variables $x_{1}, \ldots, x_{n}$ if we restrict the sums to be only over SSYT that only use the numbers $1, \ldots, n$.
Example 2.5.20. $s_{1,1}\left(x_{1}, x_{2}, \ldots, x_{n}\right)$ is the sum over SSYT of shape $(1,1)$. This is the same as a choice of $1 \leq i<j \leq n$, so $s_{1,1}\left(x_{1}, \ldots, x_{n}\right)=\sum_{1 \leq i<j \leq n} x_{i} x_{j}=e_{2}\left(x_{1}, \ldots, x_{n}\right)$, and by the same reasoning, $s_{1,1}=e_{2}$ in infinitely many variables. More generally, $s_{1^{k}}=e_{k}$ for any $k$.

Also, $s_{k}=h_{k}$ since a SSYT of shape $(k)$ is a choice of $i_{1} \leq i_{2} \leq \cdots \leq i_{k}$.
For something different, consider $s_{2,1}\left(x_{1}, x_{2}, x_{3}\right)$. There are 8 SSYT that of shape $(2,1)$ that only use $1,2,3$ :


From this, we can read off that $s_{2,1}\left(x_{1}, x_{2}, x_{3}\right)$ is a symmetric polynomial. Furthermore, it is $m_{2,1}\left(x_{1}, x_{2}, x_{3}\right)+2 m_{1,1,1}\left(x_{1}, x_{2}, x_{3}\right)$.
Theorem 2.5.21. For any $\mu \subseteq \lambda$, the skew Schur function $s_{\lambda / \mu}$ is a symmetric function.
We now focus on Schur functions. Suppose $\lambda$ is a partition of $n$. Let $K_{\lambda, \alpha}$ be the number of SSYT of shape $\lambda$ and type $\alpha$, this is called a Kostka number. The previous theorem says $K_{\lambda, \alpha}=K_{\lambda, \sigma(\alpha)}$ for any permutation $\sigma$, so it's enough to study the case when $\alpha$ is a partition. By the definition of Schur function, we have

$$
s_{\lambda}=\sum_{\mu \vdash n} K_{\lambda, \mu} m_{\mu} .
$$

An important special case is when $\mu=1^{n}$. Then $K_{\lambda, 1^{n}}$ is the number of SSYT that use each of the numbers $1, \ldots, n$ exactly once. Such a SSYT is called a standard Young tableau, and $K_{\lambda, 1^{n}}$ is denoted $f^{\lambda}$.

Theorem 2.5.22. If $K_{\lambda, \mu} \neq 0$, then $\mu \leq \lambda$ (dominance order). Also, $K_{\lambda, \lambda}=1$. In particular, the $s_{\lambda}$ form a basis for $\Lambda$.
Corollary 2.5.23. $\left\{s_{\lambda}| | \lambda \mid=d\right\}$ is a basis for $\Lambda_{d}$.
Corollary 2.5.24. $\left\{s_{\lambda}\left(x_{1}, \ldots, x_{n}\right)| | \lambda \mid=d, \ell(\lambda) \leq n\right\}$ is a basis for $\Lambda(n)_{d}$.
Proof. Note that if $\ell(\lambda)>n$, there are no SSYT only using $1, \ldots, n$, so $s_{\lambda}\left(x_{1}, \ldots, x_{n}\right)=0$. Hence the set in question spans $\Lambda(n)_{d}$. Since $\Lambda(n)_{d}$ is free of rank equal to the size of this set, it must also be a basis.
2.6. Schur functors. For the material in this section, see [Wey, Chapter 2]. What we call $\mathbf{S}_{\lambda}$ is denoted by $\mathbf{L}_{\lambda^{\prime}}$ there.

Definition 2.6.1. Let $R$ be a commutative ring and $E$ a free $R$-module. Let $\lambda$ be a partition with $n$ parts and write $m=\lambda_{1}$. We use $S^{n} E$ to denote the $n$th symmetric power of $E$. The Schur functor $\mathbf{S}_{\lambda}(E)$ is the image of the map

$$
\bigwedge_{1}^{\lambda_{1}^{\dagger}} E \otimes \cdots \otimes \bigwedge^{\lambda_{m}^{\dagger}} E \xrightarrow{\Delta} E^{\otimes \lambda_{1}^{\dagger}} \otimes \cdots \otimes E^{\otimes \lambda_{m}^{\dagger}}=E^{\otimes \lambda_{1}} \otimes \cdots \otimes E^{\otimes \lambda_{n}} \xrightarrow{\mu} S^{\lambda_{1}} E \otimes \cdots \otimes S^{\lambda_{n}} E
$$

where the maps are defined as follows. First, $\Delta$ is the product of the comultiplication maps $\bigwedge^{i} E \rightarrow E^{\otimes i}$ given by

$$
e_{1} \wedge \cdots \wedge e_{i} \mapsto \sum_{w \in \Sigma_{i}} \operatorname{sgn}(w) e_{w(1)} \otimes \cdots \otimes e_{w(i)}
$$

The equals sign is interpreted as follows: pure tensors in $E^{\otimes \lambda_{1}^{\dagger}} \otimes \cdots \otimes E^{\otimes \lambda_{m}^{\dagger}}$ can be interpreted as filling the Young diagram of $\lambda$ with vectors along the columns, which can be thought of as pure tensors in $E^{\otimes \lambda_{1}} \otimes \cdots \otimes E^{\otimes \lambda_{n}}$ by reading via rows. Finally, $\mu$ is the multiplication $\operatorname{map} E^{\otimes i} \rightarrow S^{i} E$ given by $e_{1} \otimes \cdots \otimes e_{i} \mapsto e_{1} \cdots e_{i}$.

In particular, note that $\mathbf{S}_{\lambda} E=0$ if the number of parts of $\lambda$ exceeds rank $E$.
Example 2.6.2. Take $\lambda=(3,2)$. Then the map is given by

$$
\begin{aligned}
\left(e_{1} \wedge e_{2}\right) \otimes\left(e_{3} \wedge e_{4}\right) \otimes e_{5} & \mapsto \frac{e_{1} \mid e_{3} e_{5}}{e_{2} \mid e_{4}}-\frac{e_{2}\left|e_{3}\right| e_{5}}{e_{1} e_{4}}-\frac{e_{1}\left|e_{4}\right| e_{5}}{e_{2} \mid e_{3}}+\frac{e_{2} e_{4} \mid e_{5}}{e_{1} e_{3}} \\
& \mapsto\left(e_{1} e_{3} e_{5} \otimes e_{2} e_{4}\right)-\left(e_{2} e_{3} e_{5} \otimes e_{1} e_{4}\right) \\
& -\left(e_{1} e_{4} e_{5} \otimes e_{2} e_{3}\right)+\left(e_{2} e_{4} e_{5} \otimes e_{1} e_{3}\right)
\end{aligned}
$$

The construction of $\mathbf{S}_{\lambda} E$ is functorial with respect to $E$ : given an $R$-linear map $f: E \rightarrow F$, we get an $R$-linear map $\mathbf{S}_{\lambda}(f): \mathbf{S}_{\lambda} E \rightarrow \mathbf{S}_{\lambda} F$ such that $\mathbf{S}_{\lambda}(f \circ g)=\mathbf{S}_{\lambda}(f) \circ \mathbf{S}_{\lambda}(g)$. This has two consequences: $\mathbf{S}_{\lambda} E$ is naturally a representation of $\mathbf{G L}(E)$ (in fact, it is more special because the action is defined for all linear operators of $E$, not just the invertible ones), and we can also construct $S_{\lambda} \mathcal{E}$ when $\mathcal{E}$ is a vector bundle.

Fix a basis $e_{1}, \ldots, e_{n}$ for $E$. Given a tableau $T$, we get an element $e_{T}$ in $\bigwedge^{\lambda_{1}^{\dagger}} E \otimes \cdots \otimes \bigwedge^{\lambda_{m}^{\dagger}} E$ by taking the tensor product of the wedge products of the standard basis vectors coming from the entries in each column. For example, if $T=$| 1 | 3 |
| :--- | :--- |
|  | 5 |
|  | 4 | then $e_{T}=\left(e_{1} \wedge e_{2}\right) \otimes\left(e_{3} \wedge e_{4}\right) \otimes e_{5}$.

Definition 2.6.3. Given a box $b=(i, j) \in \lambda$, its content is $c(b)=j-i$ and its hook length is $h(b)=\lambda_{i}-i+\lambda_{j}^{\dagger}-j+1$.
Example 2.6.4. Let $\lambda=(4,3,1)$. Then $\lambda^{\dagger}=(3,2,2,1)$. The contents and hook lengths are given as follows:

$$
c: \begin{array}{|l|l|l|l}
\hline 0 & 1 & 2 & 3 \\
\hline-1 & 0 & 1 \\
\hline-2 & & \\
\hline-2
\end{array} \quad h: \begin{array}{|l|l|l|l|}
\hline 6 & 4 & 3 & 1 \\
\hline 4 & 2 & 1 \\
\hline 1 & & \\
\hline
\end{array}
$$

Theorem 2.6.5. The Schur functor $\mathbf{S}_{\lambda} E$ is a free $R$-module. If rank $E=n$, then

$$
\operatorname{rank} \mathbf{S}_{\lambda} E=\prod_{b \in \lambda} \frac{n+c(b)}{h(b)}
$$

A basis for $\mathbf{S}_{\lambda} E$ is given by the images of the $e_{T}$ as $T$ ranges over all SSYT of shape $\lambda$ using the numbers $1, \ldots, n$.

If $R$ is a field, then the weight of $e_{T}$ is $x^{T}$, and hence the character of $\mathbf{S}_{\lambda} E$ is the Schur function $s_{\lambda}\left(x_{1}, \ldots, x_{n}\right)$.

The point is that to calculate things like tensor products of Schur functors, it suffices to understand how to multiply Schur functions. This is most useful in characteristic 0, where we have mentioned that the character completely encodes all information about the representation. We will discuss this in more detail later.
Remark 2.6.6. The construction we gave is unmotivated, but we can at least explain the indexing using generalities on irreducible representations of complex reductive groups, such as $\mathbf{G L}_{n}(\mathbf{C})$. A Borel subgroup $B$ is a maximal (Zariski) closed, connected, solvable (algebraic) subgroup. For $\mathbf{G L}_{n}(\mathbf{C})$, one can take the subgroup of upper-triangular matrices. Fixing $B$ and an irreducible representation $V$, there is a unique, up to scalar, nonzero vector $v \in V$ (highest weight vector) and an (algebraic) group homomorphism $\lambda: B \rightarrow \mathbf{C}^{\times}$(highest weight) such that $b \cdot v=\lambda(b) v$ for all $b \in B$.

For the upper-triangular matrices, an algebraic group homomorphism $\lambda: B \rightarrow \mathbf{C}^{\times}$takes the form $\lambda(b)=b_{1}^{\lambda_{1}} \cdots b_{n}^{\lambda_{n}}$ where $b_{1}, \ldots, b_{n}$ are the diagonal entries and $\lambda_{i} \in \mathbf{Z}$, so we can identify them with elements of $\mathbf{Z}^{n}$. It turns out that highest weights satisfy an additional constraint $\lambda_{1} \geq \lambda_{2} \geq \cdots \geq \lambda_{n}$. The indexing for the Schur functors $\mathbf{S}_{\lambda} E$ is chosen so that it has highest weight $\lambda$.
2.7. Pieri's rule. Since the $s_{\lambda}$ are a basis, we have unique expressions

$$
\begin{equation*}
s_{\mu} s_{\nu}=\sum_{\lambda} c_{\mu, \nu}^{\lambda} s_{\lambda}, \tag{2.7.1}
\end{equation*}
$$

and the $c_{\mu, \nu}^{\lambda}$ are called Littlewood-Richardson coefficients. We will see some special cases soon and study this in more depth later. From the definition, we have

$$
c_{\mu, \nu}^{\lambda}=c_{\nu, \mu}^{\lambda} .
$$

Applying $\omega$ to (2.7.1), we get

$$
\begin{equation*}
c_{\mu, \nu}^{\lambda}=c_{\nu^{\dagger}, \mu^{\dagger}}^{\lambda^{\dagger}} \tag{2.7.2}
\end{equation*}
$$

We can give an interpretation for the Littlewood-Richardson coefficients in the special case where $\mu$ (or $\nu$ ) has a single part or all parts equal to 1 . Say that $\lambda / \nu$ is a horizontal strip if no column in the skew Young diagram of $\lambda / \nu$ contains 2 or more boxes. Similarly, say that $\lambda / \nu$ is a vertical strip if no row in the skew Young diagram of $\lambda / \nu$ contains 2 or more boxes.

Theorem 2.7.3 (Pieri rule). - If $\mu=\left(1^{k}\right)$, then

$$
c_{\left(1^{k}\right), \nu}^{\lambda}= \begin{cases}1 & \text { if }|\lambda|=|\nu|+k \text { and } \lambda / \nu \text { is a vertical strip } \\ 0 & \text { otherwise }\end{cases}
$$

In other words,

$$
s_{\nu} s_{1^{k}}=\sum_{\lambda} s_{\lambda}
$$

where the sum is over all $\lambda$ such that $\lambda / \nu$ is a vertical strip of size $k$.

- If $\mu=(k)$, then

$$
c_{(k), \nu}^{\lambda}= \begin{cases}1 & \text { if }|\lambda|=|\nu|+k \text { and } \lambda / \nu \text { is a horizontal strip } \\ 0 & \text { otherwise }\end{cases}
$$

In other words,

$$
s_{\nu} s_{k}=\sum_{\lambda} s_{\lambda}
$$

where the sum is over all $\lambda$ such that $\lambda / \nu$ is a horizontal strip of size $k$.
Example 2.7.4. To multiply $s_{\lambda}$ by $s_{k}$, it suffices to enumerate all partitions that we can get by adding $k$ boxes to the Young diagram of $\lambda$, no two of which are in the same column. For example, here we have drawn all of the ways to add 2 boxes to $(4,2)$ :


So

$$
s_{4,2} s_{2}=s_{6,2}+s_{5,3}+s_{5,2,1}+s_{4,4}+s_{4,3,1}+s_{4,2,2} .
$$

Corollary 2.7.5. $s_{\nu} e_{\mu}=\sum_{\lambda} K_{\lambda^{\dagger} / \nu^{\dagger}, \mu} s_{\lambda}$.
Corollary 2.7.6. $s_{\nu} h_{\mu}=\sum_{\lambda} K_{\lambda / \nu, \mu} s_{\lambda}$.
Recall that for $|\lambda|=n, f^{\lambda}=K_{\lambda, 1^{n}}$ is the number of standard Young tableaux of shape $\lambda$.
Corollary 2.7.7. $s_{1}^{n}=\sum_{\lambda \vdash n} f^{\lambda} s_{\lambda}$.
Proof. The Pieri rule says that to multiply $s_{1}^{n}$, we first enumerate all sequences $\lambda^{(1)} \subset \lambda^{(2)} \subset$ $\cdots \subset \lambda^{(n)}$ where $\left|\lambda^{(i)}\right|=i$. Then the result is the sum of $s_{\lambda}$ with multiplicity given by the number of sequences with $\lambda^{(n)}=\lambda$. But such sequences are in bijection with standard Young tableaux: label the unique box in $\lambda^{(i)} / \lambda^{(i-1)}$ with $i$.

Remark 2.7.8. From the interpretation of $s_{\lambda}$ as the character of an irreducible representation $\mathbf{S}_{\lambda}$, and the fact that polynomial representations are direct sums of irreducible ones, we can reinterpret the Littlewood-Richardson coefficient as the multiplicity of $\mathbf{S}_{\lambda}$ in the decomposition of the tensor product of $\mathbf{S}_{\mu} \otimes \mathbf{S}_{\nu}$. From this, it is immediate that $c_{\mu, \nu}^{\lambda} \geq 0$.

The Pieri rule describes the decomposition of the tensor product of $\mathbf{S}_{\lambda}$ with an exterior power $\bigwedge^{k}$, respectively, a symmetric power $\mathrm{Sym}^{k}$.

### 2.8. Tensor categories.

2.8.1. Definitions. We will make use of tensor categories. We won't really need the precise definition, but we'll go over the main points. Let $\mathcal{A}$ be an abelian category. Roughly, this means that the objects can be thought of as the category of modules over some ring. In particular, the category of modules over a fixed ring is an example. Practically, it means that notions such as kernels, cokernels, exactness, etc. can be defined, and analogues of the isomorphism theorems for modules hold. Also, the morphisms between two fixed objects form an abelian group.

Even though we haven't made the definition precise, here is one easy construction that we will make a lot of use of. First, given two categories $\mathcal{C}$ and $\mathcal{A}$, the functor category $\operatorname{Fun}(\mathcal{C}, \mathcal{A})$ is the category whose objects are functors $\mathcal{C} \rightarrow \mathcal{A}$ and whose morphisms are natural
transformations. If $\mathcal{A}$ is abelian, then so is $\operatorname{Fun}(\mathcal{C}, \mathcal{A})$ : things such as kernels, cokernels, etc. are computed pointwise, meaning that $\operatorname{ker}(F \rightarrow G)(X)=\operatorname{ker}(F(X) \rightarrow G(X))$, etc. We'll usually apply this when $\mathcal{A}$ is the category of modules over a ring, or usually just the category of vector spaces over a field.

A monoidal structure, or tensor product, on $\mathcal{A}$ is a functor $\otimes: \mathcal{A} \times \mathcal{A} \rightarrow \mathcal{A}$, written $X \otimes Y$ rather than $\otimes(X, Y)$, which satisfies associativity in the sense that we are given an isomorphism $(X \otimes Y) \otimes Z \cong X \otimes(Y \otimes Z)$ which satisfies some axioms, together with a unit object $1_{\mathcal{A}}$ (or just 1 ) which satisfies $1 \otimes X \cong X \otimes X \otimes 1$. Again, we won't make so much use of the fine details of the definition, rather we'll work with a few concrete examples. Given a tensor product, a symmetry is a natural isomorphism $\tau_{X, Y}: X \otimes Y \cong Y \otimes X$ given for all objects $X, Y$ such that $\tau_{X, Y}=\tau_{Y, X}^{-1}$ that satisfy the braid relations. To expand on the last point: we can apply the $\tau_{X, X}$ to a 3 -fold tensor product $X^{\otimes 3}$ by either working on first two factors (call it $\tau_{1,2}$ ) or the last two factors (call it $\tau_{2,3}$ ). The braid relation is:

$$
\tau_{1,2} \tau_{2,3} \tau_{1,2}=\tau_{2,3} \tau_{1,2} \tau_{2,3}
$$

This is the relation that the transpositions $(1,2)$ and $(2,3)$ in $\Sigma_{3}$ satisfy. In particular, this implies that the symmetric group $\Sigma_{n}$ naturally acts on any $n$-fold tensor product $X^{\otimes n}$.

Example 2.8.1. The usual tensor product gives a monoidal structure on the category of modules over a commutative ring.

Given two tensor categories $\mathcal{A}, \mathcal{B}$, a tensor functor is a functor $F: \mathcal{A} \rightarrow \mathcal{B}$ together with natural isomorphisms $F(X \otimes Y) \cong F(X) \otimes F(Y)$ and $F\left(1_{\mathcal{A}}\right) \cong 1_{\mathcal{B}}$ satisfying some axioms. Again, we will not go into the details.

An important feature is that we can redo a lot of multilinear algebra in the setting of a symmetric tensor category. For example, the symmetric square of an object $S^{2}(X)$ is the cokernel of the map

$$
\tau_{X, X}-1_{X \otimes X}: X \otimes X \rightarrow X \otimes X
$$

and similarly one can define higher symmetric powers by taking the quotient of these maps for all adjacent positions. We define the exterior power $\bigwedge^{n}(X)$ as the image of the map

$$
\sum_{\sigma \in \Sigma_{n}} \operatorname{sgn}(\sigma) \sigma
$$

The multiplication maps $S^{n}(X) \otimes S^{m}(X) \rightarrow S^{n+m}(X)$ are defined as well. We can also define the comultiplication maps $\bigwedge^{n}(X) \rightarrow \bigwedge^{i}(X) \otimes \bigwedge^{n-i}(X)$ as follows. Pick coset representatives $\sigma_{1}, \ldots, \sigma_{N}$ for $\Sigma_{n} / \Sigma_{i} \times \Sigma_{n-i}$ and consider the map

$$
\sum_{i=1}^{N} \operatorname{sgn}\left(\sigma_{i}\right) \sigma_{i}
$$

on $X^{\otimes n}$. This maps $\bigwedge^{n}(X)$ into $\bigwedge^{i}(X) \otimes \bigwedge^{n-i}(X)$ and is the desired map. Once we have these maps, we can also define arbitrary Schur functors $\mathbf{S}_{\lambda}(X)$ using the definition from the previous section.
2.8.2. Example: chain complexes. Given a commutative ring $R$, a chain complex $V_{\bullet}$ is a sequence of $R$-modules and $R$-linear maps (differentials):

$$
\cdots \xrightarrow{d_{i}} V_{i} \xrightarrow{d_{i-1}} V_{i-1} \xrightarrow{d_{i-2}} \cdots
$$

such that $d_{i-1} d_{i}=0$ for all $i$. A map of chain complexes $f: V_{\bullet} \rightarrow W_{\bullet}$ is a sequence of linear maps $f_{i}: V_{i} \rightarrow W_{i}$ such that $f_{i} d_{i}^{V}=d_{i}^{W} f_{i+1}$ for all $i$. The notions of image complex, kernel, cokernel, etc. can be computed term by term. Given chain complexes $V_{\bullet}$ and $W_{\bullet}$, their tensor product is the chain complex with spaces

$$
(V \otimes W)_{n}=\bigoplus_{i+j=n} V_{i} \otimes W_{j}
$$

and the differentials are the sums of the maps

$$
V_{i} \otimes W_{j} \rightarrow\left(V_{i-1} \otimes W_{j}\right) \oplus\left(V_{i} \otimes W_{j-1}\right)
$$

where the map is $d_{i-1}^{V} \otimes 1+(-1)^{i} 1 \otimes d_{j-1}^{W}$.
The symmetry isomorphism $V_{\bullet} \otimes W_{\bullet} \cong W_{\bullet} \otimes V_{\bullet}$ is defined as the sum of the isomorphisms

$$
\begin{aligned}
V_{i} \otimes W_{j} & \rightarrow W_{j} \otimes V_{i} \\
\sum x \otimes y & \mapsto(-1)^{i j} \sum y \otimes x
\end{aligned}
$$

Remark 2.8.2. We can think of the category of $R$-modules as a subcategory of the category of chain complexes $V_{\bullet}$ where $V_{i}=0$ for $i \neq 0$. Then the symmetry just defined restricts to the usual one. However, we could also think of $R$-modules as complexes where $V_{i}=0$ for $i \neq 1$. Then the symmetry restricts to a twisted version of the usual symmetry which will come up later.
2.8.3. Example: symmetric sequences. Let $\mathbf{F B}$ be the category of finite sets and bijections.

Let $\mathbf{k}$ be a field and let $\mathbf{V e c}_{\mathbf{k}}$ the category of $\mathbf{k}$-vector spaces.
Set $\mathcal{V}_{\mathbf{k}}:=\operatorname{Fun}\left(\mathbf{F B}, \mathbf{V e c}_{\mathbf{k}}\right)$. Given an object $F \in \mathcal{V}_{\mathbf{k}}$ and a finite set $S$, we get a vector space with an action of $\operatorname{Aut}(S)$, which is isomorphic to a symmetric group. This information is already encoded by just considering the objects of the form [n], but it will be convenient for various constructions to consider all finite sets. So the data of an object of $\mathcal{V}_{\mathbf{k}}$ is the same as a sequence of representations $V_{n}$ of symmetric groups $\Sigma_{n}$, one for each $n$.

Given a vector space $M$, define $M[i] \in \mathcal{V}_{\mathbf{k}}$ by

$$
S \mapsto \begin{cases}M & \text { if }|S|=i \\ 0 & \text { if }|S| \neq i\end{cases}
$$

and all morphisms act by the identity.
Put a monoidal structure on $\mathcal{V}_{\mathbf{k}}$ as follows. Given $V, W \in \mathcal{V}_{\mathbf{k}}$, we define their tensor product to be

$$
(V \otimes W)(S)=\bigoplus_{T \subseteq S} V(T) \otimes_{\mathbf{k}} W(S \backslash T)
$$

This has a symmetry $\tau$ given by interchanging factors. The unit is $\mathbf{k}[0]$.
Remark 2.8.3. This is one instance where it can be convenient to use all finite sets, not just those of the form $[n]$. However, we can also phrase things from this perspective: given sequences $\left(V_{n}\right)$ and $\left(W_{n}\right)$, their tensor product is the sequence defined by

$$
(V \otimes W)_{n}=\bigoplus_{i=0}^{n} \operatorname{Ind}_{\Sigma_{i} \times \Sigma_{n-i}}^{\Sigma_{n}}\left(V_{i} \otimes W_{n-i}\right)
$$

where for finite groups $H \subset G$, and an $H$-representation $V$, the induction is defined by $\operatorname{Ind}_{H}^{G} V=\mathbf{k}[G] \otimes_{\mathbf{k}[H]} V$ and $\Sigma_{i} \times \Sigma_{n-i}$ is a subgroup of $\Sigma_{n}$ that only permutes $1, \ldots, i$ and $i+1, \ldots, n$ separately.

So we can define tensor powers $V^{\otimes n}$, and $\tau$ allows one to define symmetric powers $\operatorname{Sym}^{n}(V)$ and exterior powers $\bigwedge^{n}(V)$ as quotients of $V^{\otimes n}$.

Example 2.8.4.

$$
\begin{aligned}
U[1]^{\otimes n}(S) & = \begin{cases}\bigoplus_{\sigma \in \operatorname{Aut}(S)} U^{\otimes S} \cong\left(U^{\otimes n}\right)^{\oplus n!} & \text { if }|S|=n \\
0 & \text { else }\end{cases} \\
\operatorname{Sym}^{n}(U[1])(S) & = \begin{cases}U^{\otimes S} & \text { if }|S|=n \\
0 & \text { else }\end{cases}
\end{aligned}
$$

The action of $\operatorname{Aut}(S)$ on $U^{\otimes S}$ is by permuting tensor factors.
$\nu_{\mathbf{k}}$ has a subcategory $\mathcal{V}_{\mathbf{k}}^{\mathrm{f}}$ of "locally finite" objects: those functors such that $\operatorname{dim} V(S)<\infty$ for all $S$.
2.8.4. Example: polynomial functors. Now assume that $\mathbf{k}$ is an infinite field. Write $\mathbf{V e c}_{\mathbf{k}}$ for the category of finite dimensional vector spaces over $\mathbf{k}$.

Definition 2.8.5. A functor $F: \mathbf{V e c}_{\mathbf{k}} \rightarrow \mathbf{V e c}_{\mathbf{k}}$ is polynomial if for all $V, V^{\prime}$, the map

$$
\operatorname{Hom}_{V^{\prime} \mathbf{c}_{\mathbf{k}}}\left(V, V^{\prime}\right) \rightarrow \operatorname{Hom}_{\mathbf{V e c}_{\mathbf{k}}}\left(F(V), F\left(V^{\prime}\right)\right)
$$

is defined by polynomial functions. If these polynomials are homogeneous of degree $d$, we say that $F$ is a degree $d$ polynomial functor. Let $\mathbf{P o l}_{\mathbf{k}}$ be the category of polynomial functors, and let $\mathbf{P o l}_{\mathbf{k}, d}$ be the subcategory of degree $d$ polynomial functors.

Note that a polynomial functor is naturally a direct sum of its homogeneous parts, so we have $\mathbf{P o l}_{\mathbf{k}}=\bigoplus_{d \geq 0} \mathbf{P o l}_{\mathbf{k}, d}$.

The Schur functors $\mathbf{S}_{\lambda}$ are examples of polynomial functors. They are homogeneous of degree $|\lambda|$.

Theorem 2.8.6. If $\mathbf{k}$ is a field of characteristic 0 , then every irreducible polynomial functor is a Schur functor.
$\mathrm{Pol}_{\mathrm{k}}$ has a tensor product:

$$
\left(F \otimes F^{\prime}\right)(V)=F(V) \otimes_{\mathbf{k}} F^{\prime}(V)
$$

and a symmetry $\tau$ which interchanges the factors.
$\mathbf{P o l}_{\mathbf{k}}$ has a subcategory $\mathbf{P o l}_{\mathbf{k}}^{\mathrm{f}}$ of locally finite objects: those functors such that the degree $d$ piece is a finite length functor for all $d$, i.e., there are only finitely many subfunctors.
2.8.5. Algebras and modules. Let $(\mathcal{A}, \otimes)$ be a monoidal category with unit 1 . An algebra is an object $A$ together with maps

$$
\mu: A \otimes A \rightarrow A, \quad e: 1 \rightarrow A
$$

which should satisfy the axioms of being an associative, unital algebra (thinking of $\mu$ as multiplication and $e$ as inclusion of the unit). More precisely, it means that the following
diagram should commute:

and $\mu \circ(e \otimes \mathrm{id}): 1 \otimes A \rightarrow A \otimes A \rightarrow A$ is the same as the natural isomorphism $1 \otimes A \cong A$ and similarly with $\mu \circ(\mathrm{id} \otimes e)$.

Given an algebra $A$, a (left) $A$-module is an object $M$ together with a map $\mu: A \otimes M \rightarrow M$ such that the following diagram commutes:

and axioms similar to the unit axioms for $A$ are satisfied.
If $\mathcal{A}$ is also equipped with a symmetry $\tau$, then we can also define the notion of a commutative algebra. This is an algebra as above such that $\mu=\mu \circ \tau$ as maps $A \otimes A \rightarrow A$. The definition of a module over a commutative algebra remains unchanged.
2.9. Categorical version of Schur-Weyl duality. The decomposition of $s_{1}^{n}$ can be interpreted as a decomposition of the tensor power of a vector space

$$
\left(\mathbf{C}^{d}\right)^{\otimes n}=\bigoplus_{\substack{\lambda \vdash n \\ \ell(\lambda) \leq d}} \mathbf{S}_{\lambda}\left(\mathbf{C}^{d}\right)^{\oplus f^{\lambda}} .
$$

Hence the multiplicity space of $\mathbf{S}_{\lambda}\left(\mathbf{C}^{d}\right)$ has dimension $f^{\lambda}$. By Schur-Weyl duality, we see that $f^{\lambda}$ is the dimension of an irreducible representation of $\Sigma_{n}$, which we will call $\mathbf{M}_{\lambda}$. In the next section, we will discuss a construction for $\mathbf{M}_{\lambda}$. For now, one can take the multiplicity space of $\mathbf{S}_{\lambda} \mathbf{C}^{d}$ in $\left(\mathbf{C}^{d}\right)^{\otimes n}$.

Hence, the indexing set in Schur-Weyl duality can be taken to be the set of partitions of size $n$ with at most $d$ parts. If $d \geq n$, this latter condition is unnecessary.

Define a functor $\Phi: \operatorname{Fun}\left(\mathbf{F B}, \mathbf{V e c}_{\mathbf{k}}\right) \rightarrow \mathbf{P o l}_{\mathbf{k}}$ by $V \mapsto \Phi_{V}$ where $\Phi_{V}$ is defined by:

$$
\Phi_{V}(W)=\bigoplus_{n \geq 0}\left(V_{n} \otimes_{\mathbf{k}} W^{\otimes n}\right)^{\Sigma_{n}}
$$

and we write $V_{i}$ for $V([i])$. Here $\Sigma_{n}$ acts on $W^{\otimes n}$ by permuting factors, and the superscript denotes invariants.

Proposition 2.9.1. $\Phi$ is a tensor functor between $\mathcal{V}_{\mathbf{k}}^{\mathrm{f}}$ and $\mathbf{P o l}_{\mathbf{k}}^{\mathrm{f}}$.
If $\mathbf{k}$ has characteristic 0 , then $\Phi$ is an equivalence between $\mathcal{V}_{\mathbf{k}}^{\mathrm{f}}$ and $\mathbf{P o l}_{\mathbf{k}}^{\mathrm{f}}$.

Proof. Pick $V, V^{\prime} \in \operatorname{Fun}\left(\mathbf{F B}, \mathbf{V e c}_{\mathbf{k}}\right)$. Then we have

$$
\begin{aligned}
\left(\Phi_{V} \otimes \Phi_{V^{\prime}}\right)(W) & =\Phi_{V}(W) \otimes \Phi_{V^{\prime}}(W) \\
& =\left(\bigoplus_{n \geq 0}\left(V_{n} \otimes W^{\otimes n}\right)^{\Sigma_{n}}\right) \otimes\left(\bigoplus_{n \geq 0}\left(V_{n}^{\prime} \otimes W^{\otimes n}\right)^{\Sigma_{n}}\right) \\
& =\bigoplus_{n \geq 0} \bigoplus_{i=0}^{n}\left(V_{i} \otimes W^{\otimes i}\right)^{\Sigma_{i}} \otimes\left(V_{n-i}^{\prime} \otimes W^{\otimes(n-i)}\right)^{\Sigma_{n-i}} \\
& =\bigoplus_{n \geq 0} \bigoplus_{i=0}^{n}\left(\operatorname{Ind}_{\Sigma_{i} \times \Sigma_{n-i}}^{\Sigma_{n}}\left(V_{i} \otimes V_{n-i}^{\prime}\right) \otimes W^{\otimes n}\right)^{\Sigma_{n}} \\
& =\Phi_{V \otimes V^{\prime}}(W) .
\end{aligned}
$$

The fourth equality comes from hom-tensor adjunction (Frobenius reciprocity). To elaborate: since everything is finite-dimensional, we can write (below Res denotes the restriction of a representation to a subgroup)

$$
\begin{aligned}
\left(\operatorname{Ind}_{\Sigma_{i} \times \Sigma_{n-i}}^{\Sigma_{n}}\left(V_{i} \otimes V_{n-i}^{\prime}\right) \otimes W^{\otimes n}\right)^{\Sigma_{n}} & =\operatorname{Hom}_{\Sigma_{n}}\left(\left(W^{\otimes n}\right)^{*}, \operatorname{Ind}_{\Sigma_{i} \times \Sigma_{n-i}}^{\Sigma_{n}}\left(V_{i} \otimes V_{n-i}^{\prime}\right)\right) \\
& =\operatorname{Hom}_{\Sigma_{i} \times \Sigma_{n-i}}\left(\operatorname{Res}_{\Sigma_{i} \times \Sigma_{n-i}}^{\Sigma_{n}}\left(W^{\otimes n}\right)^{*}, V_{i} \otimes V_{n-i}^{\prime}\right) \\
& =\operatorname{Hom}_{\Sigma_{i}}\left(\left(W^{\otimes i}\right)^{*}, V_{i}\right) \otimes \operatorname{Hom}_{\Sigma_{n-i}}\left(\left(W^{\otimes(n-i)}\right)^{*}, V_{n-i}^{\prime}\right) \\
& =\left(V_{i} \otimes W^{\otimes i}\right)^{\Sigma_{i}} \otimes\left(V_{n-i}^{\prime} \otimes W^{\otimes(n-i)}\right)^{\Sigma_{n-i}} .
\end{aligned}
$$

For the last statement, first note that every locally finite object of $\mathcal{V}_{\mathbf{k}}$ is a direct sum of $\mathbf{M}_{\lambda}$, where we think of $\mathbf{M}_{\lambda}$ as a functor that has nonzero values only for sets of size $|\lambda|$ and only finitely many copies are used for each $\lambda$. By Schur-Weyl duality, the image of $\mathbf{M}_{\lambda}$ is $\mathbf{S}_{\lambda}$. In particular, $\Phi$ induces a bijection between the sets of irreducible objects of both categories, and every object is a direct sum of irreducible objects (note that $\Phi$ preserves direct sums), so we are done.

Remark 2.9.2. For a direct construction of $\mathrm{M}_{\lambda}$, see $[\mathrm{F}, \S 7.2]$.
Corollary 2.9.3. Given partitions $\lambda, \mu$ of size $n$ and $m$, we have

$$
\operatorname{Ind}_{\Sigma_{n} \times \Sigma_{m}}^{\Sigma_{n+m}} \mathbf{M}_{\lambda} \boxtimes \mathbf{M}_{\mu}=\bigoplus_{\nu} \mathbf{M}_{\nu}^{\oplus c_{\lambda, \mu}^{\nu}} .
$$

Furthermore, for any partition $\nu$ of size $p$, have

$$
\operatorname{Res}_{\Sigma_{i} \times \Sigma_{p-i}}^{\Sigma_{p}} \mathbf{M}_{\nu}=\bigoplus_{\lambda, \mu}\left(\mathbf{M}_{\lambda} \boxtimes \mathbf{M}_{\mu}\right)^{\oplus c_{\lambda, \mu}^{\nu}}
$$

Proof. The first part comes from the fact that the Littlewood-Richardson coefficients describe multiplication for GL-representations. The second part follows from the first by Frobenius reciprocity.
2.10. Infinite number of variables. There is a third symmetric monoidal category which will be convenient to use. We have inclusions

$$
\begin{aligned}
\mathbf{G} \mathbf{L}_{n}(\mathbf{C}) & \rightarrow \mathbf{G L}_{n+1}(\mathbf{C}) \\
X & \mapsto\left(\begin{array}{cc}
X & 0 \\
0 & 1
\end{array}\right),
\end{aligned}
$$

and we define

$$
\mathbf{G} \mathbf{L}_{\infty}(\mathbf{C})=\bigcup_{n \geq 1} \mathbf{G L}_{n}(\mathbf{C})
$$

This group consists of infinite invertible matrices which differ from the identity matrix only in finitely many entries. Let $T$ denote the subgroup of diagonal matrices. Say that a polynomial representation of $\mathbf{G} \mathbf{L}_{\infty}(\mathbf{C})$ is $T$-semisimple if its restriction to $T$ decomposes as a direct sum of 1-dimensional representations. These will be polynomial representations of $T$, which are classified by non-negative integer sequences $\left(\alpha_{1}, \alpha_{2}, \ldots\right)$. Furthermore, since it comes a polynomial representation of $\mathbf{G L}_{\infty}(\mathbf{C})$, these sequences can only have a finite number of nonzero entries. The representation $\mathbf{C}_{\alpha}=\mathbf{C}$ indexed by $\alpha$ is described by

$$
\operatorname{diag}\left(x_{1}, x_{2}, \ldots\right) \cdot 1=\prod_{i \geq 1} x_{i}^{\alpha_{i}}
$$

We let $\operatorname{Rep}(\mathbf{G L})$ be the category of polynomial $T$-semisimple representations of $\mathbf{G L}_{\infty}(\mathbf{C})$ such that each representation $\mathbf{C}_{\alpha}$ appears with finite multiplicity. Given $V \in \operatorname{Rep}(\mathbf{G L})$, let $V_{n}$ be the sum of the $\mathbf{C}_{\alpha}$ such that $\sum_{i} \alpha_{i}=n$. Then $V=\bigoplus_{n \geq 0} V_{n}$ is a decomposition of $V$ into a direct sum of subrepresentations. Say that $V$ is locally finite if each $V_{n}$ has a finite composition series. We let $\operatorname{Rep}^{\mathrm{f}}(\mathbf{G L})$ be the full subcategory of $\operatorname{Rep}(\mathbf{G L})$ consisting of the locally finite representations. Let $V_{\alpha}$ be the direct sum of the $\mathbf{C}_{\alpha}$ that appear in $V$.

Given an object $V$ of $\operatorname{Rep}(\mathbf{G L})$, let $V(n)$ denote the direct sum of those $V_{\alpha}$ such that $\alpha_{i}=0$ for $i>n$. Then $V(n)$ can also be thought of as the invariant subspace of diagonal matrices whose first $n$ entries are 1. Since the subgroup $\mathbf{G L}_{n}(\mathbf{C})$ commutes with these diagonal matrices, it acts on $V(n)$. Furthermore, we have $V=\bigcup_{n \geq 1} V(n)$. Note that there is a natural tensor product on objects of $\operatorname{Rep}(\mathbf{G L})$ (the conditions we imposed are closed under tensor products).

Lemma 2.10.1. The assignment $\mathbf{C}^{n} \mapsto V(n)$ is a polynomial functor.
Proof. Since $V$ is a polynomial representation, the function $\rho: \mathbf{G L}_{\infty}(\mathbf{C}) \rightarrow \mathbf{G L}(V)$ can be extended to $\operatorname{End}\left(\mathbf{C}^{\infty}\right) \rightarrow \operatorname{End}(V)$ (the entries of $\rho$ are polynomials and don't require invertibility of the input matrix). Given any linear map $f: \mathbf{C}^{n} \rightarrow \mathbf{C}^{m}$, we can extend this to an element $f \in \operatorname{End}\left(\mathbf{C}^{\infty}\right)$ and this gives a linear map $V(n) \rightarrow V(m)$ given by polynomial functions. Since each $\mathbf{C}_{\alpha}$ appears with finite multiplicity, $V(n)$ is finite-dimensional.

Call the polynomial functor $\mathbf{C}^{n} \mapsto V(n)$ by $F_{V}$.
Theorem 2.10.2. The assignment $V \mapsto F_{V}$ defines an equivalence between $\operatorname{Rep}(\mathbf{G L})$ and Pol of symmetric monoidal categories.
Proof. To define the inverse, let $F$ be a polynomial functor. The inclusion maps $\mathbf{C}^{n} \rightarrow \mathbf{C}^{n+1}$ give maps $F\left(\mathbf{C}^{n}\right) \rightarrow F\left(\mathbf{C}^{n+1}\right)$ (also inclusions because the projection $\mathbf{C}^{n+1} \rightarrow \mathbf{C}^{n}$ induce a left inverse after applying $F$ ). Then we get a polynomial representation of $\mathbf{G L}_{\infty}(\mathbf{C})$ by taking $\bigcup_{n \geq 1} F\left(\mathbf{C}^{n}\right)$. The fact that the symmetric monoidal structures is preserved just comes from the fact that it is defined by a tensor product of vector spaces in both cases.

Given $V \in \operatorname{Rep}(\mathbf{G L})$, we can define its character by

$$
\operatorname{char}(V)=\sum_{\alpha} \operatorname{dim}\left(V_{\alpha}\right) \prod_{i \geq 1} x_{i}^{\alpha_{i}}
$$

Lemma 2.10.3. $\operatorname{char}(V) \in \Lambda$.

Proof. Proof is similar to the case of polynomial representations of $\mathbf{G L}_{n}(\mathbf{C})$.
This is one natural context for symmetric functions, and the assignment $V \mapsto V(n)$ on the level of characters corresponds to setting $x_{i}=0$ for $i>n$.

Given a polynomial functor $F$ which is a direct sum of Schur functors (i.e., any element in $\mathbf{P o l}^{\mathrm{f}}$ ), we define $\ell(F)$ to the supremum of $\ell(\lambda)$ such that $\mathbf{S}_{\lambda}$ appears in $F$. Note that $\mathbf{S}_{\lambda}\left(\mathbf{C}^{n}\right) \neq 0$ if and only if $n \geq \ell(\lambda)$, so $\ell(F)$ tells us how big $n$ must be so that $F\left(\mathbf{C}^{n}\right)$ doesn't lose any information (i.e., kill any submodules). An advantage to working with $\operatorname{Rep}(\mathbf{G L})$ is that we are essentially working with the case $n=\infty$, and the equivalence above is stating that we never lose information in this case.
2.11. Littlewood-Richardson rule. We now want to understand the decomposition

$$
\operatorname{Ind}_{\Sigma_{n} \times \Sigma_{m}}^{\Sigma_{n+m}}\left(\mathbf{M}_{\lambda} \otimes \mathbf{M}_{\mu}\right)=\bigoplus_{\nu} \mathbf{M}_{\nu}^{\oplus c_{\lambda, \mu}^{\nu}}
$$

or equivalently, the decomposition

$$
\mathbf{S}_{\lambda} \otimes \mathbf{S}_{\mu}=\bigoplus_{\nu} \mathbf{S}_{\nu}^{\oplus c_{\lambda, \mu}^{\nu}}
$$

where the sum is over all partitions $\nu$ of $n+m$ and $c_{\lambda, \mu}^{\nu}$ is a non-negative integer. The numbers $c_{\lambda, \mu}^{\nu}$ are called the Littlewood-Richardson coefficients, and the eponymous rule gives a combinatorial description of them.

There are many descriptions for this rule, and we will formulate it via lattice words. First consider the skew-diagram $\nu / \lambda$. We fill the boxes with positive integers so that $i$ appears exactly $\mu_{i}$ times. Then $c_{\lambda, \mu}^{\nu}$ counts the number of such fillings which satisfy the properties (Littlewood-Richardson tableaux):

- semistandard: the entries are weakly increasing from left to right in each row, and the entries are strictly increasing from top to bottom in each column
- lattice word: Read the entries right to left in each row, starting with the top row to get a sequence of positive integers (reading word). Then each initial segment of this sequence has the property that for each $i, i$ occurs at least as many times as $i+1$.
See [F, §5, §7.3] or [Mac, §I.9]. See also [Sta, Appendix 7.A.1.3] for some other formulations of the rule. Here are some simple consequences of the Littlewood-Richardson rule:
- If $c_{\lambda, \mu}^{\nu} \neq 0$, then $\lambda \subseteq \nu$ and $\mu \subseteq \nu$.
- For all partitions $\lambda, \mu, c_{\lambda, \mu}^{\lambda+\mu}=1$, and $c_{\lambda, \mu}^{\lambda \cup \mu}=1$ where $\lambda \cup \mu$ denotes the partition obtained by sorting the sequence $(\lambda, \mu)$. To prove these, fill the Young diagram of $\mu$ with the number $i$ in each box in the $i$ th row. Append the $i$ th row to the $i$ th row of $\lambda$ to see $c_{\lambda, \mu}^{\lambda+\mu} \geq 1$. Append the $i$ th column to the $i$ th column of $\lambda$ to see $c_{\lambda, \mu}^{\lambda \cup \mu} \geq 1$. The reverse inequalities follow by the extremality of these shapes.
- For all integers $N>0$, we have $c_{N \lambda, N \mu}^{N \nu} \geq c_{\lambda, \mu}^{\nu}$, which can be seen by "stretching" the Littlewood-Richardson tableau. As a consequence, if $c_{\lambda, \mu}^{\nu}>0$, then $c_{N \lambda, N \mu}^{N \nu}>0$ for any $N>0$. The converse of this statement is also true, i.e., if $c_{N \lambda, N \mu}^{N \nu}>0$ for some $N>0$, then $c_{\lambda, \mu}^{\nu}>0$. This is a highly non-trivial fact known as the saturation theorem, see [KT1, DW1, KM] for different proofs of it. Furthermore, the function $N \mapsto C_{N \lambda, N \mu}^{N \nu}$ is a polynomial in $N \geq 0$ for any fixed choice of $\lambda, \mu, \nu$ [DW2, Corollary 3].

And here are some properties which are not obvious from the Littlewood-Richardson rule, but follow easily from the representation-theoretic interpretation:

- Symmetry: $c_{\lambda, \mu}^{\nu}=c_{\mu, \lambda}^{\nu}$. One way to give a symmetric combinatorial rule for $c_{\lambda, \mu}^{\nu}$ is to use the plactic monoid and jeu de taquin [F, §2, §5.1].
- Transpose symmetry: $c_{\lambda^{\dagger}, \mu^{\dagger}}^{\nu^{\dagger}}=c_{\lambda, \mu}^{\nu}$.

Example 2.11.1. We calculate $c_{(3,1),(4,2,1)}^{(5,3,2,1)}=3$. The Littlewood-Richardson tableaux are

|  |  |  | 1 | 1 |
| :---: | :---: | :---: | :---: | :---: |
|  | 1 | 1 |  |  |
| 2 | 2 |  |  |  |
| 3 |  |  |  |  |



The reading words are 1111223,1121213 , and 1121312 , respectively. It is easier to calculate this number after swapping the roles of $(3,1)$ and $(4,2,1)$ :

2.12. A few more formulas. By the equivalences above, the induction product on symmetric group representations corresponds to the tensor product of polynomial functors, or tensor product of representations in $\operatorname{Rep}(\mathbf{G L})$.

Proposition 2.12.1. We have a $\mathbf{G L}_{n} \times \mathbf{G L}_{m}$-equivariant decomposition

$$
\mathbf{S}_{\nu}\left(\mathbf{C}^{n} \oplus \mathbf{C}^{m}\right)=\bigoplus_{\lambda, \mu}\left(\mathbf{S}_{\lambda}\left(\mathbf{C}^{n}\right) \boxtimes \mathbf{S}_{\mu}\left(\mathbf{C}^{m}\right)\right)^{\oplus c_{\lambda, \mu}^{\nu}}
$$

Proof. Set $k=|\nu|$. We have

$$
\mathbf{S}_{\nu}\left(\mathbf{C}^{n} \oplus \mathbf{C}^{m}\right)=\operatorname{Hom}_{\Sigma_{k}}\left(\mathbf{M}_{\nu},\left(\mathbf{C}^{n+m}\right)^{\otimes k}\right) .
$$

Now $\left(\mathbf{C}^{n+m}\right)^{\otimes k}$ can be written as $\bigoplus_{I} \bigotimes_{j=1}^{k} A_{j, I}$ where the sum is over all subsets of $\{1, \ldots, k\}$ and $A_{j, I}=\mathbf{C}^{n}$ if $j \in I$ and $A_{j, I}=\mathbf{C}^{m}$ if $j \notin I$. For every $N$, the symmetric group $\Sigma_{k}$ preserves the sum $\bigoplus_{|I|=N} \bigotimes_{j} A_{j, I}$, and this representation is the induced representation

$$
\operatorname{Ind}_{\Sigma_{N} \times \Sigma_{k-N}}^{\Sigma_{k}}\left(\left(\mathbf{C}^{n}\right)^{\otimes N} \otimes\left(\mathbf{C}^{m}\right)^{\otimes(k-n)}\right)
$$

Hence by Frobenius reciprocity, we can write

$$
\begin{aligned}
\mathbf{S}_{\nu}\left(\mathbf{C}^{n} \oplus \mathbf{C}^{m}\right) & =\bigoplus_{N=0}^{k} \operatorname{Hom}_{\Sigma_{N} \times \Sigma_{k-N}}\left(\mathbf{M}_{\nu}| |_{\Sigma_{N} \times \Sigma_{k-N}}^{\Sigma_{k}},\left(\mathbf{C}^{n}\right)^{\otimes N} \otimes\left(\mathbf{C}^{m}\right)^{\otimes(k-N)}\right) \\
& =\bigoplus_{N=0}^{k} \bigoplus_{\lambda, \mu} \operatorname{Hom}_{\Sigma_{N} \times \Sigma_{k-N}}\left(\mathbf{M}_{\lambda} \otimes \mathbf{M}_{\mu},\left(\mathbf{C}^{n}\right)^{\otimes N} \otimes\left(\mathbf{C}^{m}\right)^{\otimes(k-N)}\right)^{\oplus c_{\lambda, \mu}^{\nu}} \\
& =\bigoplus_{\lambda, \mu}\left(\mathbf{S}_{\lambda}\left(\mathbf{C}^{n}\right) \boxtimes \mathbf{S}_{\mu}\left(\mathbf{C}^{m}\right)\right)^{\oplus c_{\lambda, \mu}^{\nu}} .
\end{aligned}
$$

This gives a $\mathbf{G L} L_{n}(\mathbf{C}) \times \mathbf{G} \mathbf{L}_{m}(\mathbf{C})$-equivariant decomposition of $\mathbf{S}_{\nu}\left(\mathbf{C}^{n} \oplus \mathbf{C}^{m}\right)$, and we see that the multiplicities are described by Littlewood-Richardson coefficients.

The induction product of symmetric group representations is natural from the perspective of our equivalence with polynomial functors, but one can also consider tensor products of symmetric group representations. Again, we will have certain decompositions

$$
\mathbf{M}_{\lambda} \otimes \mathbf{M}_{\mu}=\bigoplus_{\nu} \mathbf{M}_{\nu}^{g_{\lambda, \mu}^{\nu}}
$$

In fact, the $\mathbf{M}_{\lambda}$ are isomorphic to their own duals since they can be defined over the rational numbers (it is enough to be defined over the real numbers), so we have

$$
g_{\lambda, \mu}^{\nu}=\operatorname{dim}_{\mathbf{C}}\left(\mathbf{M}_{\lambda} \otimes \mathbf{M}_{\mu} \otimes \mathbf{M}_{\nu}\right)^{\Sigma_{n}}
$$

from which it is clear that the coefficient is symmetric in $\lambda, \mu, \nu$, so we denote them by $g_{\lambda, \mu, \nu}$. These are the Kronecker coefficients and it is a notoriously difficult problem to give subtraction-free combinatorial rules for them. A special case is when $\lambda=(n)$, since then $\mathbf{M}_{\lambda}$ is the trivial representation, in which case we get

$$
g_{(n), \mu, \nu}=\delta_{\mu, \nu}
$$

Another easy case is when $\lambda=\left(1^{n}\right)$ since the $\mathbf{M}_{\left(1^{n}\right)}$ is the sign representation. It turns out that $\mathbf{M}_{\lambda} \otimes \operatorname{sgn}=\mathbf{M}_{\lambda^{\dagger}}$, and hence

$$
g_{\left(1^{n}\right), \mu, \nu}=\delta_{\mu, \nu^{\dagger}} .
$$

Proposition 2.12.2. We have a $\mathbf{G L}_{n} \times \mathbf{G L}_{m}$-equivariant decomposition

$$
\mathbf{S}_{\nu}\left(\mathbf{C}^{n} \otimes \mathbf{C}^{m}\right)=\bigoplus_{\lambda, \mu}\left(\mathbf{S}_{\lambda}\left(\mathbf{C}^{n}\right) \boxtimes \mathbf{S}_{\mu}\left(\mathbf{C}^{m}\right)\right)^{\oplus g_{\lambda, \mu, \nu}} .
$$

Proof. First, we have

$$
\left(\mathbf{C}^{n} \otimes \mathbf{C}^{m}\right)^{\otimes k}=\bigoplus_{|\nu|=k} \mathbf{S}_{\nu}\left(\mathbf{C}^{n} \otimes \mathbf{C}^{m}\right) \boxtimes \mathbf{M}_{\nu}
$$

as representations of $\mathbf{G} \mathbf{L}_{n m}(\mathbf{C}) \times \Sigma_{k}$. Alternatively, as $\mathbf{G L}(n) \times \mathbf{G L}(m) \times \Sigma_{k}$ representations, we have

$$
\begin{aligned}
\left(\mathbf{C}^{n}\right)^{\otimes k} \otimes\left(\mathbf{C}^{m}\right)^{\otimes k} & =\left(\bigoplus_{|\lambda|=k} \mathbf{S}_{\lambda}\left(\mathbf{C}^{n}\right) \boxtimes \mathbf{M}_{\lambda}\right) \otimes\left(\bigoplus_{|\mu|=k} \mathbf{S}_{\mu}\left(\mathbf{C}^{m}\right) \boxtimes \mathbf{M}_{\mu}\right) \\
& =\bigoplus_{|\lambda|=|\mu|=|\nu|=k}\left(\mathbf{S}_{\lambda}\left(\mathbf{C}^{n}\right) \boxtimes \mathbf{S}_{\mu}\left(\mathbf{C}^{m}\right) \boxtimes \mathbf{M}_{\nu}\right)^{\oplus g_{\lambda, \mu, \nu}} .
\end{aligned}
$$

The result follows by taking the $\mathbf{M}_{\nu}$-isotypic component of both expressions.
Corollary 2.12.3 (Cauchy identities). For each $k$, we have $\mathbf{G L}_{n}(\mathbf{C}) \times \mathbf{G L}_{m}(\mathbf{C})$-equivariant decompositions

$$
\begin{aligned}
\operatorname{Sym}^{k}\left(\mathbf{C}^{n} \otimes \mathbf{C}^{m}\right) & =\bigoplus_{|\lambda|=k} \mathbf{S}_{\lambda}\left(\mathbf{C}^{n}\right) \otimes \mathbf{S}_{\lambda}\left(\mathbf{C}^{m}\right), \\
\bigwedge^{k}\left(\mathbf{C}^{n} \otimes \mathbf{C}^{m}\right) & =\bigoplus_{|\lambda|=k} \mathbf{S}_{\lambda}\left(\mathbf{C}^{n}\right) \otimes \mathbf{S}_{\lambda^{+}}\left(\mathbf{C}^{m}\right) .
\end{aligned}
$$

Of course, when $n, m$ are small relative to $k$, we don't need to sum over all partitions: recall that $\mathbf{S}_{\lambda}\left(\mathbf{C}^{n}\right)=0$ if $n<\ell(\lambda)$.

## 3. REpRESENTATIONS OF COMBINATORIAL CATEGORIES

3.1. Twisted commutative algebras. Let $\mathbf{k}$ be a commutative ring. A twisted commutative algebra (tca) is a commutative algebra in the tensor category of symmetric sequences over $\mathbf{k}$. Let's unpack this definition.

Let $A=\bigoplus_{n \geq 0} A_{n}$ be a $\mathbf{Z}_{\geq 0}$-graded, associative, unital $\mathbf{k}$-algebra such that each $A_{n}$ has a linear action of the symmetric group $\Sigma_{n}$. Embed $\Sigma_{n} \times \Sigma_{m}$ as a subgroup of $\Sigma_{n+m}$ by having $\Sigma_{n}$ act on $1, \ldots, n$ and $\Sigma_{m}$ act on $n+1, \ldots, n+m$ in the natural way. Then $A$ is a twisted commutative algebra (tca) if, for all $n$ and $m$,

- the multiplication map

$$
A_{n} \otimes A_{m} \rightarrow A_{n+m}
$$

is $\Sigma_{n} \times \Sigma_{m}$-equivariant with respect to the embedding just specified, and

- $\tau(x y)=y x$ where $x \in A_{n}, y \in A_{m}$, and $\tau \in \Sigma_{n+m}$ is the permutation that swaps $\{1, \ldots, n\}$ and $\{n+1, \ldots, n+m\}$ in order, i.e., $\tau(i)=m+i$ if $1 \leq i \leq n$ and $\tau(n+j)=j$ for $1 \leq j \leq m$.
Example 3.1.1. Let $E$ be a free $\mathbf{k}$-module, and $A_{n}=E^{\otimes n}$ (by convention, $A_{0}=\mathbf{k}$ ). The action of $\Sigma_{n}$ is by permuting tensor factors, and multiplication is concatenation of tensors. This is the free tca (i.e., a symmetric algebra in the category of symmetric sequences) generated in degree 1 by $E$, in the sense that $A_{n}=\operatorname{Sym}^{n}\left(A_{1}\right)$ for all $n$ and $A_{0}=\mathbf{k}$. We will denote it by $\operatorname{Sym}(E\langle 1\rangle)$.

When using all sets, we have $A(S)=E^{\otimes S}$, where we can think of $E^{\otimes S}$ as the space of tensors indexed by the set $S$ (more precisely, this should be the colimit of $E^{\otimes n}$ over all bijections $S \cong[n]$, where $n=|S|$ ).

We say that $A$ is finitely generated if it is a quotient of $\operatorname{Sym}(V)$ for some finite length symmetric sequence $V$.

We can also define modules over tca's. In fact, we have a general definition using the monoidal language. We can unpack it as follows. Let $A$ be a tca. An $A$-module $M$ is a graded abelian group $M=\bigoplus_{n \geq 0} M_{n}$ such that each $M_{n}$ is a linear $\Sigma_{n}$-representation, $M$ is a graded $A$-module in the usual sense, and such that the multiplication map

$$
A_{n} \otimes M_{m} \rightarrow M_{n+m}
$$

is $\Sigma_{n} \times \Sigma_{m}$-equivariant. A module is finitely generated if it can be generated by finitely many elements $m_{1}, \ldots, m_{n}$ under the action of $A$ and the symmetric groups. We are interested in the category of $A$-modules, and specifically, in formal properties of finitely generated $A$-modules.

In the language of symmetric sequences, a module $M$ over a tca $A$ is finitely generated if it is a quotient of $A \otimes V$ for some finite length symmetric sequence $V$ (i.e., $\sum \operatorname{dim} V_{n}<\infty$ ).

An important notion for us is noetherianity. A module is noetherian if all of its submodules are finitely generated. A tca is noetherian if all of its finitely generated modules are noetherian. The next facts follow from some standard algebra, and we leave it to the reader:

Proposition 3.1.2. Let $\mathbf{k}\left[\Sigma_{n}\right]$ be the symmetric sequence which is the regular representation of $\Sigma_{n}$ (i.e., $\Sigma_{n}$ acting on its group algebra by multiplication on the left) in degree $n$ and 0 elsewhere. $A$ tca $A$ is noetherian if and only if $A \otimes \mathbf{k}\left[\Sigma_{n}\right]$ is noetherian for all $n$.

Proposition 3.1.3. Let $M$ be a module over a tca. The following are equivalent:
(1) $M$ is noetherian.
(2) Submodules of $M$ satisfy the ascending chain condition, i.e., given a chain of submodules $N_{1} \subseteq N_{2} \subseteq \cdots$ of $M$, we must have $M_{i}=M_{i+1}=\cdots$ for some $i$.
(3) Every collection of submodules of $M$ has a maximal element with respect to inclusion.
3.2. Alternative models for tca's generated in degree 1. The tca $\operatorname{Sym}(E\langle 1\rangle)$ and its modules are an important example, especially when $\operatorname{dim} E=1$. We can give a different model for $\operatorname{Sym}(E\langle 1\rangle)$-modules in terms of functor categories. To do this, pick a basis $e_{1}, \ldots, e_{d}$ for E.

Let $\mathbf{F I}_{d}$ be the category such that:

- the objects are finite sets $S$,
- a morphism $S \rightarrow T$ is a pair $(f, g)$ where $f: S \rightarrow T$ is an injective function and $g: T \backslash f(S) \rightarrow[d]$ is an arbitrary function ("coloring").
Given two morphisms $(f, g): S \rightarrow T$ and $\left(f^{\prime}, g^{\prime}\right): T \rightarrow U$, we define the composition to be $\left(f^{\prime \prime}, g^{\prime \prime}\right): S \rightarrow U$ where $f^{\prime \prime}=f^{\prime} \circ f$, and $g^{\prime \prime}: U \backslash f^{\prime \prime}(S) \rightarrow[d]$ is given by $g^{\prime \prime}(u)=g^{\prime}(u)$ if $u$ is not in the image of $f^{\prime}$, and $g^{\prime \prime}(u)=g^{\prime}(v)$ if $u=f^{\prime}(v)$ and $v$ is not in the image of $f$.
Proposition 3.2.1. The category $\operatorname{Fun}\left(\mathbf{F I}_{d}, \operatorname{Mod}_{\mathbf{k}}\right)$ is equivalent to the category of $\operatorname{Sym}(E\langle 1\rangle)$ modules.
Proof. Set $A=\operatorname{Sym}(E\langle 1\rangle)$. Let $M$ be an $A$-module. Define a functor $F_{M}: \mathbf{F I}_{d} \rightarrow \operatorname{Mod}_{\mathbf{k}}$ on objects by setting $F_{M}(S)=M(S)$ for each finite set $S$. Consider the multiplication map $A \otimes M \rightarrow M$ evaluated on a set $T$ (indexed slightly differently):

$$
\bigoplus_{U \subseteq T} E^{\otimes(T \backslash U)} \otimes M(U) \rightarrow M(T)
$$

Given an $\mathbf{F I}_{d}$-morphism $(f, g): S \rightarrow T$, we consider the map $E^{\otimes(T \backslash U)} \otimes M(U) \rightarrow M(T)$ above with $U=f(S)$. The coloring $g: T \backslash U \rightarrow[d]$ gives a tensor product $x$ of basis elements in $E^{\otimes(T \backslash U)}$, and hence we get a map $M(S) \rightarrow M(T)$ by considering the restriction to $x \otimes M(U) \rightarrow M(T)$ and identifying $U$ with $S$ via $f$.

Conversely, given a functor $F: \mathbf{F I}_{d} \rightarrow \operatorname{Mod}_{\mathbf{k}}$, we define an $A$-module $M_{F}$ on objects by $M_{F}(S)=F(S)$. To define the multiplication map, we need to define the maps $E^{\otimes(T \backslash U)} \otimes$ $M(U) \rightarrow M(T)$ as above for every $U \subseteq T$ and every tensor product $x$ of basis elements. Let $f: U \rightarrow T$ be the inclusion map and let $g: T \backslash U \rightarrow[d]$ be the map corresponding to $x$. Then we get a map $F(U) \rightarrow F(T)$ which we use to define the desired map.

The two constructions are inverse to each other, so we get the desired equivalence.
The special case $d=1$ will be studied in more depth later. In that case, we write FI instead of $\mathbf{F I}_{1}$. This is an abbreviation for $\mathbf{F}$ inite $\mathbf{I n j e c t i o n s . ~}$

It will also be useful to have descriptions of $\operatorname{Sym}(E\langle 1\rangle)$ and its modules in the polynomial functor and $\mathbf{G} \mathbf{L}_{\infty}$-models. In that case, we assume that our field is of characteristic 0 . Chasing through the equivalence $\Phi$ from $\S 2.9$, the polynomial functor corresponding to $\operatorname{Sym}(E\langle 1\rangle)$ in degree $n$ is

$$
W \mapsto\left(E^{\otimes n} \otimes W^{\otimes n}\right)^{\Sigma_{n}} \cong \operatorname{Sym}^{n}(E \otimes W),
$$

so we get the functor which assigns $W$ to the symmetric algebra on $E \otimes W$. Given a linear map $W \rightarrow W^{\prime}$, we get an induced linear map $E \otimes W \rightarrow E \otimes W^{\prime}$ and hence a ring homomorphism $\operatorname{Sym}(E \otimes W) \rightarrow \operatorname{Sym}\left(E \otimes W^{\prime}\right)$. Recall the Cauchy identity (Corollary 2.12.3) that

$$
\operatorname{Sym}(E \otimes W)=\bigoplus_{\lambda} \mathbf{S}_{\lambda} E \otimes \mathbf{S}_{\lambda} W
$$

This is a $\mathbf{G L}(E) \times \mathbf{G} \mathbf{L}(W)$-equivariant decomposition, and hence gives a decomposition into polynomial functors if we ignore the $\mathbf{G L}(E)$-action. Note that the action of $\mathbf{G L}(E)$ is not present in any of the definitions for $\operatorname{Sym}(E\langle 1\rangle)$-modules. In this language, Proposition 3.1.2 translates to saying that the module $W \mapsto \operatorname{Sym}(E \otimes W) \otimes W^{\otimes n}$ is noetherian for all $n$.

In the model of representations of $\mathbf{G} \mathbf{L}_{\infty}(\mathbf{C})($ see $\S 2.10), \operatorname{Sym}(E\langle 1\rangle)$ is the algebra

$$
\operatorname{Sym}\left(E \otimes \mathbf{C}^{\infty}\right)
$$

with the action of $\mathbf{G L}_{\infty}(\mathbf{C})$ which acts on $\mathbf{C}^{\infty}$ in the usual way. The latter is a very useful model to keep in mind. A module over this tca is then a $\mathbf{G} \mathbf{L}_{\infty}(\mathbf{C})$ representation (satisfying the conditions in $\S 2.10$ with a compatible action of $\operatorname{Sym}\left(E \otimes \mathbf{C}^{\infty}\right)$. In this case, Proposition 3.1.2 translates to saying that the module $\operatorname{Sym}\left(E \otimes \mathbf{C}^{\infty}\right) \otimes\left(\mathbf{C}^{\infty}\right)^{n}$ is noetherian for all $n$.
3.3. Bounded tca's. In this section, we continue to work over a field of characteristic 0 .

Given a sum of Schur functors $F=\bigoplus_{\lambda} \mathbf{S}_{\lambda}^{\oplus c_{\lambda}}$, we defined $\ell(F)$ to be the supremum of $\ell(\lambda)$ such that $c_{\lambda} \neq 0$. We say that $F$ is bounded if $\ell(F)<\infty$. Given a tca $A$ in the polynomial functor model, we say that $A$ is bounded if it is bounded as a polynomial functor.

Given any vector space $\mathbf{C}^{n}$, we can evaluate a tca $A$ (and an $A$-module $M$ ) to get an algebra $A\left(\mathbf{C}^{n}\right)$ with a $\mathbf{G L} L_{n}(\mathbf{C})$-action together with a module $M\left(\mathbf{C}^{n}\right)$ with a compatible action. This gives a function from the lattice of submodules of $M$ to the lattice of submodules of $M\left(\mathbf{C}^{n}\right)$ which sends $N \subset M$ to $N\left(\mathbf{C}^{n}\right) \subset M\left(\mathbf{C}^{n}\right)$.

Proposition 3.3.1. Let $A$ be a finitely generated tca. If $M$ is a bounded $A$-module and $n \geq \ell(M)$, then the map of lattices above is injective. In particular, $M$ is noetherian.

Proof. Recall that $N\left(\mathbf{C}^{n}\right)=0$ if and only if $n<\ell(N)$. But $\ell(N) \leq \ell(M)$, so $N\left(\mathbf{C}^{n}\right) \neq 0$ as long as $N \neq 0$. This proves the injectivity statement. For the noetherianity statement, write $A$ as a quotient of $\operatorname{Sym}(V)$ with $V$ a finite length symmetric sequence. Then $A\left(\mathbf{C}^{n}\right)$ is generated as an algebra (in the usual sense) by $V\left(\mathbf{C}^{n}\right)$, a finite dimensional vector space, and $M\left(\mathbf{C}^{n}\right)$ is a finitely generated $A\left(\mathbf{C}^{n}\right)$ in the usual sense. Hence $A\left(\mathbf{C}^{n}\right)$ is a noetherian algebra, and the submodules of $M\left(\mathbf{C}^{n}\right)$ satisfy the ascending chain condition.

Proposition 3.3.2. If $A$ is a bounded tca, then every finitely generated $A$-module is bounded. In particular, $A$ is noetherian.

Proof. Let $M$ be a finitely generated $A$-module and write $M$ as a quotient of $A \otimes V$ with $V$ a finite length symmetric sequence. By $\S 2.11, \ell(A \otimes V)=\ell(A)+\ell(V)<\infty$. Furthermore, $\ell(M) \leq \ell(A \otimes V)$, so $M$ is bounded.

Proposition 3.3.3. The tca $\operatorname{Sym}(E\langle 1\rangle)$ is bounded. In particular, it is noetherian.
Proof. It follows immediately from the Cauchy identity that $\ell(\operatorname{Sym}(E\langle 1\rangle))=\operatorname{dim} E$.
3.4. Noetherianity in general. Via our equivalence of definitions, we see that $\operatorname{Fun}\left(\mathbf{F I}_{d}, \operatorname{Mod}_{\mathbf{k}}\right)$ satisfies some noetherian property when $\mathbf{k}$ is a field of characteristic 0 . One can ask whether this holds for general fields. The answer is yes, and we discuss this now and consider when $\mathbf{F I}_{d}$ is replaced by a more general category. This follows the strategy in [SS1].

We now shift our language slightly. Let $\mathcal{C}$ be a category. ${ }^{3}$ A representation of $\mathcal{C}$ (or a $\mathcal{C}$-module) over $\mathbf{k}$ is a functor $\mathcal{C} \rightarrow \operatorname{Mod}_{\mathbf{k}}$. A map of $\mathcal{C}$-modules is a natural transformation. We write $\operatorname{Rep}_{\mathbf{k}}(\mathcal{C})$ for the category of representations, which is abelian. Let $M$ be a representation of $\mathcal{C}$. A subrepresentation is a functor $N$ such that $N(x)$ is a subspace of $M(x)$ for all $x$ and which is closed under all operations $M(f)$ for all morphisms $f$. By an element of $M$ we mean an element of $M(x)$ for some object $x$ of $\mathcal{C}$. Given any set $S$ of elements of $M$, there is a smallest subrepresentation of $M$ containing $S$; we call this the subrepresentation generated by $S$. We say that $M$ is finitely generated if it is generated by a finite set of elements. For a morphism $f: x \rightarrow y$ in $\mathcal{C}$, we typically write $f_{*}$ for the given map of $\mathbf{k}$-modules $M(x) \rightarrow M(y)$.

Let $x$ be an object of $\mathcal{C}$. Define a representation $P_{x}$ of $\mathcal{C}$ by $P_{x}(y)=\mathbf{k}\left[\operatorname{Hom}_{\mathcal{C}}(x, y)\right]$, i.e., $P_{x}(y)$ is the free left $\mathbf{k}$-module with basis $\operatorname{Hom}_{\mathcal{C}}(x, y)$, the set of morphisms $x \rightarrow y$. For a morphism $f: x \rightarrow y$, we write $e_{f}$ for the corresponding element of $P_{x}(y)$. If $M$ is another representation then $\operatorname{Hom}\left(P_{x}, M\right)=M(x)$. This shows that $\operatorname{Hom}\left(P_{x},-\right)$ is an exact functor, and so $P_{x}$ is a projective object of $\operatorname{Rep}_{\mathbf{k}}(\mathcal{C})$. We call it the principal projective at $x$. A $\mathcal{C}$-module is finitely generated if and only if it is a quotient of a finite direct sum of principal projective objects.

An object of $\operatorname{Rep}_{\mathbf{k}}(\mathcal{C})$ is noetherian if every ascending chain of subobjects stabilizes; this is equivalent to every subrepresentation being finitely generated. The category $\operatorname{Rep}_{\mathbf{k}}(\mathcal{C})$ is locally noetherian if every finitely generated object in it is.
Proposition 3.4.1. The category $\operatorname{Rep}_{\mathbf{k}}(\mathcal{C})$ is locally noetherian if and only if every principal projective is noetherian.
Proof. By definition, if $\operatorname{Rep}_{\mathbf{k}}(\mathrm{C})$ is locally noetherian then so is every principal projective. Conversely, suppose every principal projective is noetherian. Let $M$ be a finitely generated object. Then $M$ is a quotient of a finite direct sum $P$ of principal projectives. Since noetherianity is preserved under finite direct sums, $P$ is noetherian. And since noetherianity descends to quotients, $M$ is noetherian. This completes the proof.

Let $\Phi: \mathcal{C} \rightarrow \mathcal{C}^{\prime}$ be a functor. There is then a pullback functor on representations $\Phi^{*}: \operatorname{Rep}_{\mathbf{k}}\left(\mathcal{C}^{\prime}\right) \rightarrow$ $\operatorname{Rep}_{\mathbf{k}}(\mathcal{C})$ given by $M \mapsto M \circ \Phi$. We study how $\Phi^{*}$ interacts with finiteness conditions. The following definition is of central importance:
Definition 3.4.2. We say that $\Phi$ satisfies property (F) (F for finite) if the following condition holds: given any object $x$ of $\mathcal{C}^{\prime}$ there exist finitely many objects $y_{1}, \ldots, y_{n}$ of $\mathcal{C}$ and morphisms $f_{i}: x \rightarrow \Phi\left(y_{i}\right)$ in $\mathcal{C}^{\prime}$ such that for any object $y$ of $\mathcal{C}$ and any morphism $f: x \rightarrow \Phi(y)$ in $\mathcal{C}^{\prime}$, there exists a morphism $g: y_{i} \rightarrow y$ in $\mathcal{C}$ such that $f=\Phi(g) \circ f_{i}$.

The following proposition is the motivation for introducing property (F).
Proposition 3.4.3. A functor $\Phi: \mathcal{C} \rightarrow \mathcal{C}^{\prime}$ satisfies property $(\mathrm{F})$ if and only if $\Phi^{*}$ takes finitely generated objects of $\operatorname{Rep}_{\mathbf{k}}\left(\mathcal{C}^{\prime}\right)$ to finitely generated objects of $\operatorname{Rep}_{\mathbf{k}}(\mathcal{C})$.
Proof. Assume that $\Phi$ satisfies property (F). It suffices to show that $\Phi^{*}$ takes principal projectives to finitely generated representations. Thus let $P_{x}$ be the principal projective of $\operatorname{Rep}_{\mathbf{k}}\left(\mathcal{C}^{\prime}\right)$ at an object $x$. Note that $\Phi^{*}\left(P_{x}\right)(y)$ has for a basis the elements $e_{f}$ for $f \in$ $\operatorname{Hom}_{\mathcal{C}^{\prime}}(x, \Phi(y))$. Let $f_{i}: x \rightarrow \Phi\left(y_{i}\right)$ be as in the definition of property (F). Then the $e_{f_{i}}$ generate $\Phi^{*}\left(P_{x}\right)$. The converse is left to the reader (and not used in this paper).

[^2]Proposition 3.4.4. Suppose that $\Phi: \mathcal{C} \rightarrow \mathfrak{C}^{\prime}$ is an essentially surjective functor. Let $M$ be an object of $\operatorname{Rep}_{\mathbf{k}}\left(\mathrm{C}^{\prime}\right)$ such that $\Phi^{*}(M)$ is finitely generated (resp. noetherian). Then $M$ is finitely generated (resp. noetherian).

Proof. Let $S$ be a set of elements of $\Phi^{*}(M)$. Let $S^{\prime}$ be the corresponding set of elements of $M$. (Thus if $S$ contains $m \in \Phi^{*}(M)(y)$ then $S^{\prime}$ contains $m \in M(\Phi(y))$.) If $N$ is a subrepresentation of $M$ containing $S^{\prime}$ then $\Phi^{*}(N)$ is a subrepresentation of $\Phi^{*}(M)$ containing $S$. It follows that if $N$ (resp. $N^{\prime}$ ) is the subrepresentation of $M$ (resp. $\Phi^{*}(M)$ ) generated by $S^{\prime}($ resp. $S)$, then $N^{\prime} \subset \Phi^{*}(N)$. Thus if $S$ generates $\Phi^{*}(M)$ then $\Phi^{*}(N)=\Phi^{*}(M)$, which implies $N=M$ since $\Phi$ is essentially surjective, i.e., $S$ generates $M$. In particular, if $\Phi^{*}(M)$ is finitely generated then so is $M$.

Now suppose that $\Phi^{*}(M)$ is noetherian. Given a subrepresentation $N$ of $M$, we obtain a subrepresentation $\Phi^{*}(N)$ of $\Phi^{*}(M)$. Since $\Phi^{*}(M)$ is noetherian, it follows that $\Phi^{*}(N)$ is finitely generated. Thus $N$ is finitely generated, and so $M$ is noetherian.

Corollary 3.4.5. Let $\Phi: \mathcal{C} \rightarrow \mathcal{C}^{\prime}$ be an essentially surjective functor satisfying property $(\mathrm{F})$ and suppose $\operatorname{Rep}_{\mathbf{k}}(\mathcal{C})$ is noetherian. Then $\operatorname{Rep}_{\mathbf{k}}\left(\mathrm{C}^{\prime}\right)$ is noetherian.

Proof. Let $M$ be a finitely generated $\complement^{\prime}$-module. Then $\Phi^{*}(M)$ is finitely generated by Proposition 3.4.3, and therefore noetherian, and so $M$ is noetherian by Proposition 3.4.4.
3.5. Noetherian posets. Let $X$ be a partially ordered set (poset). Then $X$ satisfies the ascending chain condition (ACC) if every ascending chain in $X$ stabilizes, i.e., given $x_{1} \leq$ $x_{2} \leq \cdots$ in $X$ we have $x_{i}=x_{i+1}$ for $i \gg 0$. The descending chain condition (DCC) is defined similarly. An anti-chain in $X$ is a sequence $x_{1}, x_{2}, \ldots$ such that $x_{i} \not \leq x_{j}$ for all $i \neq j$. An ideal in $X$ is a subset $I$ of $X$ such that $x \in I$ and $x \leq y$ implies $y \in I$. We write $\mathcal{J}(X)$ for the poset of ideals of $X$, ordered by inclusion. For $x \in X$, the principal ideal generated by $x$ is $\{y \mid y \geq x\}$. An ideal is finitely generated if it is a finite union of principal ideals. The following result is left to the reader as an exercise.

Proposition 3.5.1. The following conditions on $X$ are equivalent:
(a) The poset $X$ satisfies $D C C$ and has no infinite anti-chains.
(b) Given a sequence $x_{1}, x_{2}, \ldots$ in $X$, there exists $i<j$ such that $x_{i} \leq x_{j}$.
(c) The poset $\mathcal{J}(X)$ satisfies $A C C$.
(d) Every ideal of $X$ is finitely generated.

The poset $X$ is noetherian if the above conditions are satisfied. Where we say " $X$ is noetherian," one often sees " $\leq$ is a well-quasi-order" in the literature. Similarly, where we say " $X$ satisfies DCC" one sees " $\leq$ is well-founded."

Proposition 3.5.2. Let $X$ be a noetherian poset and let $x_{1}, x_{2}, \ldots$ be a sequence in $X$. Then there exists an infinite sequence of indices $i_{1}<i_{2}<\cdots$ such that $x_{i_{1}} \leq x_{i_{2}} \leq \cdots$.

Proof. Let $I$ be the set of indices such that $i \in I$ and $j>i$ implies that $x_{i} \nleftarrow x_{j}$. If $I$ is infinite, then there is $i<i^{\prime}$ with $i, i^{\prime} \in I$ such that $x_{i} \leq x_{i^{\prime}}$ by definition of noetherian and hence contradicts the definition of $I$. So $I$ is finite; let $i_{1}$ be any number larger than all elements of $I$. Then by definition of $I$, we can find $x_{i_{1}} \leq x_{i_{2}} \leq \cdots$.

Proposition 3.5.3. Let $X$ and $Y$ be noetherian posets. Then $X \times Y$ is noetherian.

Proof. Let $\left(x_{1}, y_{1}\right),\left(x_{2}, y_{2}\right), \ldots$ be an infinite sequence in $X \times Y$. Since $X$ is noetherian, there exists $i_{1}<i_{2}<\cdots$ such that $x_{i_{1}} \leq x_{i_{2}} \leq \cdots$ (Proposition 3.5.2). Since $Y$ is noetherian, there exists $i_{j}<i_{j^{\prime}}$ such that $y_{i_{j}} \leq y_{i_{j^{\prime}}}$, and hence $\left(x_{i_{j}}, y_{i_{j}}\right) \leq\left(x_{i_{j^{\prime}}}, y_{i_{j^{\prime}}}\right)$.

Given a poset $X$, let $X^{\star}$ be the set of finite words $x_{1} \cdots x_{n}$ with $x_{i} \in X$. We define $x_{1} \cdots x_{n} \leq x_{1}^{\prime} \cdots x_{m}^{\prime}$ if there exist $1 \leq i_{1}<\cdots<i_{n} \leq m$ such that $x_{j} \leq x_{i_{j}}^{\prime}$ for $j=1, \ldots, n$.
Theorem 3.5.4 (Higman's lemma [Hi]). If $X$ is a noetherian poset, then so is $X^{\star}$.
Proof. Suppose that $X^{\star}$ is not noetherian. We use Nash-Williams' theory of minimal bad sequences [NW] to get a contradiction. A sequence $w_{1}, w_{2}, \ldots$ of elements in $X^{\star}$ is bad if $w_{i} \ddagger w_{j}$ for all $i<j$. We pick a bad sequence minimal in the following sense: for all $i \geq 1$, among all bad sequences beginning with $w_{1}, \ldots, w_{i-1}$ (this is the empty sequence for $i=1), \ell\left(w_{i}\right)$ is as small as possible. Let $x_{i} \in X$ be the first element of $w_{i}$ and let $v_{i}$ be the subword of $w_{i}$ obtained by removing $x_{i}$. By Proposition 3.5.2, there is an infinite sequence $i_{1}<i_{2}<\cdots$ such that $x_{i_{1}} \leq x_{i_{2}} \leq \cdots$. Then $w_{1}, w_{2}, \ldots, w_{i_{1}-1}, v_{i_{1}}, v_{i_{2}}, \ldots$ is a bad sequence because $v_{i_{j}} \leq w_{i_{j}}$ for all $j$, and $v_{i_{j}} \leq v_{i_{j^{\prime}}}$ would imply that $w_{i_{j}} \leq w_{i_{j^{\prime}}}$. It is smaller than our minimal bad sequence, so we have reached a contradiction. Hence $X^{\star}$ is noetherian.
3.6. Monomial representations and Gröbner bases. Let $\mathcal{C}$ be an essentially small category and let Set be the category of sets. Fix a functor $S: \mathcal{C} \rightarrow \mathbf{S e t}$, and let $P=\mathbf{k}[S]$, i.e., $P(x)$ is the free k-module on the set $S(x)$.

Given $f \in S(x)$, we write $e_{f}$ for the corresponding element of $P(x)$. An element of $P$ is a monomial if it is of the form $\lambda e_{f}$ for some $\lambda \in \mathbf{k}$ and $f \in S(x)$. A subrepresentation $M$ of $P$ is monomial if $M(x)$ is spanned by the monomials it contains, for all objects $x$.

To connect arbitrary subrepresentations of $P$ to monomial subrepresentations, we need a theory of monomial orders. Let WO be the category of well-ordered sets and strictly orderpreserving functions. There is a forgetful functor $\mathbf{W O} \rightarrow$ Set. An ordering on $S$ is a lifting of $S$ to WO. More concretely, an ordering on $S$ is a choice of well-order on $S(x)$, for each $x \in \mathcal{C}$, such that for every morphism $x \rightarrow y$ in $\mathcal{C}$ the induced map $S(x) \rightarrow S(y)$ is strictly order-preserving. We write $\preceq$ for an ordering; $S$ is orderable if it admits an ordering.

Suppose $\preceq$ is an ordering on $S$. Given non-zero $\alpha=\sum_{f \in S(x)} \lambda_{f} e_{f}$ in $P(x)$, we define the initial term of $\alpha$, denoted init $(\alpha)$, to be $\lambda_{g} e_{g}$, where $g=\max _{\preceq}\left\{f \mid \lambda_{f} \neq 0\right\}$. The initial variable of $\alpha$, denoted $\operatorname{init}_{0}(\alpha)$, is $g$. Now let $M$ be a subrepresentation of $P$. We define the initial representation of $M$, denoted $\operatorname{init}(M)$, as follows: $\operatorname{init}(M)(x)$ is the $\mathbf{k}$-span of the elements init $(\alpha)$ for non-zero $\alpha \in M(x)$. The name is justified by the following result.
Proposition 3.6.1. Notation as above, $\operatorname{init}(M)$ is a monomial subrepresentation of $P$.
Proof. Let $\alpha=\sum_{i=1}^{n} \lambda_{i} e_{f_{i}}$ be an element of $M(x)$ with each $\lambda_{i}$ non-zero, ordered so that $f_{i} \prec f_{1}$ for all $i>1$. Thus $\operatorname{init}(\alpha)=\lambda_{1} e_{f_{1}}$. Let $g: x \rightarrow y$ be a morphism. Then $g_{*}(\alpha)=$ $\sum_{i=1}^{n} \lambda_{i} e_{g_{*}\left(f_{i}\right)}$. Since $g_{*}: S(x) \rightarrow S(y)$ is strictly order-preserving, we have $g_{*}\left(f_{i}\right) \prec g_{*}\left(f_{1}\right)$ for all $i>1$. Thus $\operatorname{init}\left(g_{*}(\alpha)\right)=\lambda_{1} e_{g f_{1}}$, or, in other words, $\operatorname{init}\left(g_{*}(\alpha)\right)=g_{*}(\operatorname{init}(\alpha))$. This shows that $g_{*}$ maps $\operatorname{init}(M)(x)$ into $\operatorname{init}(M)(y)$, and so $\operatorname{init}(M)$ is a subrepresentation of $P$. That it is monomial follows immediately from its definition.
Proposition 3.6.2. If $N \subseteq M$ are subrepresentations of $P$ and $\operatorname{init}(N)=\operatorname{init}(M)$, then $M=N$.

Proof. Assume that $M(x) \neq N(x)$ for $x \in \mathcal{C}$. Let $K \subset S(x)$ be the set of all elements which appear as the initial variable of some element of $M(x) \backslash N(x)$. Then $K \neq \varnothing$, so has
a minimal element $f$ with respect to $\preceq$. Pick $\alpha \in M(x) \backslash N(x)$ with $\operatorname{init}_{0}(\alpha)=f$. By assumption, there exists $\beta \in N(x)$ with $\operatorname{init}(\alpha)=\operatorname{init}(\beta)$. But then $\alpha-\beta \in M(x) \backslash N(x)$, and $\operatorname{init}_{0}(\alpha-\beta) \prec \operatorname{init}_{0}(\alpha)$, a contradiction. Thus $M=N$.

Let $M$ be a subrepresentation of $P$. A set $\mathfrak{G}$ of elements of $M$ is a Gröbner basis of $M$ if $\{\operatorname{init}(\alpha) \mid \alpha \in \mathfrak{G}\}$ generates $\operatorname{init}(M)$. Note that $M$ has a finite Gröbner basis if and only if $\operatorname{init}(M)$ is finitely generated. As usual, we have:

Proposition 3.6.3. Let $\mathfrak{G}$ be a Gröbner basis of $M$. Then $\mathfrak{G}$ generates $M$.
Proof. Let $N \subseteq M$ be the subrepresentation generated by $\mathfrak{G}$. Then init $(N)$ contains init $(\alpha)$ for all $\alpha \in \mathfrak{G}$, and so $\operatorname{init}(N)=\operatorname{init}(M)$. Thus $M=N$ by Proposition 3.6.2.

We now come to our main result.
Theorem 3.6.4. Suppose $\mathbf{k}$ is noetherian, $S$ is orderable, and $|S|$ is noetherian. Then every subrepresentation of $P$ has a finite Gröbner basis. In particular, $P$ is a noetherian object of $\operatorname{Rep}_{\mathbf{k}}(\mathcal{C})$.

Proof. It suffices to show that every monomial subrepresentation is finitely generated. Let $Q$ be the direct product of the poset of ideals in $\mathbf{k}$ with $|S|$. Suppose we have a strictly increasing chain of monomial subrepresentations $M_{1} \subset M_{2} \subset \cdots \subset P$. For each $i$, pick a monomial $\lambda_{i} e_{f_{i}}=m_{i} \in M_{i} \backslash M_{i-1}$ and let $I_{i}$ be the ideal in $\mathbf{k}$ generated by all $\lambda_{i}$ which are coefficients of $e_{f_{i}}$ appearing in any monomial in $M_{i}$. Then we get a sequence $\left(I_{1}, f_{1}\right),\left(I_{2}, f_{2}\right), \ldots$ in $Q$ which are incomparable, contradicting the fact that $|S|$, and hence $Q$ (by Proposition 3.5.3), is noetherian.
3.7. Gröbner categories. Let $\mathcal{C}$ be an essentially small category. For an object $x$, let $S_{x}: \mathcal{C} \rightarrow$ Set be the functor given by $S_{x}(y)=\operatorname{Hom}_{\mathcal{C}}(x, y)$. Note that $P_{x}=\mathbf{k}\left[S_{x}\right]$.

Definition 3.7.1. We say that $\mathcal{C}$ is Gröbner if, for all objects $x$, the functor $S_{x}$ is orderable and the poset $\left|S_{x}\right|$ is noetherian. We say that $\mathcal{C}$ is quasi-Gröbner if there exists a Gröbner category $\mathcal{C}^{\prime}$ and an essentially surjective functor $\mathcal{C}^{\prime} \rightarrow \mathcal{C}$ satisfying property ( F ).

Theorem 3.7.2. Let $\mathcal{C}$ be a quasi-Gröbner category. Then for any noetherian ring $\mathbf{k}$, the category $\operatorname{Rep}_{\mathbf{k}}(\mathcal{C})$ is noetherian.

Proof. First suppose that $\mathcal{C}$ is a Gröbner category. Then every principal projective of $\operatorname{Rep}_{\mathbf{k}}(\mathcal{C})$ is noetherian, by Theorem 3.6.4, and so $\operatorname{Rep}_{\mathbf{k}}(\mathcal{C})$ is noetherian by Proposition 3.4.1.

Now suppose that $\mathcal{C}$ is quasi-Gröbner, and let $\Phi: \mathcal{C}^{\prime} \rightarrow \mathcal{C}$ be an essentially surjective functor satisfying property (F), with $\mathcal{C}^{\prime}$ Gröbner. Then $\operatorname{Rep}_{\mathbf{k}}\left(\mathrm{C}^{\prime}\right)$ is noetherian, by the previous paragraph, and so $\operatorname{Rep}_{\mathbf{k}}(\mathcal{C})$ is noetherian by Corollary 3.4.5.

Remark 3.7.3. If the functor $S_{x}$ is orderable, then the group $\operatorname{Aut}(x)$ admits a well-order compatible with the group operation, and is therefore trivial. Thus, in a Gröbner category, there are no non-trivial automorphisms.

Proposition 3.7.4. The cartesian product of finitely many (quasi-) Gröbner categories is (quasi-) Gröbner.

We leave the details to the reader.
3.8. Example: $\mathbf{F I}_{d}$-modules. Let $d$ be a positive integer. Define $\mathbf{O I}_{d}$ to be the ordered version of $\mathbf{F I}_{d}$ : its objects are totally ordered finite sets and a morphisms between $S$ and $T$ is a pair $(f, g)$ with $f$ an order-preserving injection, and $g$ is an arbitrary function $g: T \backslash f(S) \rightarrow$ $[d]$. When $d=1$, we will write $\mathbf{F I}$ and $\mathbf{O I}$ instead of $\mathbf{F I}_{1}$ and $\mathbf{O I}_{1}$.

Let $\Sigma=\{0, \ldots, d\}$, and let $\mathcal{L}$ be the subset of $\Sigma^{\star}$ consisting of words $w_{1} \cdots w_{r}$ in which exactly $n$ of the $w_{i}$ are equal to 0 . Partially order $\mathcal{L}$ using the subsequence order, i.e., if $s:[i] \rightarrow \Sigma$ and $t:[j] \rightarrow \Sigma$ are words then $s \leq t$ if there exists $I \subseteq[j]$ such that $s=\left.t\right|_{I}$.
Lemma 3.8.1. The poset $\mathcal{L}$ is noetherian.
Proof. Noetherianity is an immediate consequence of Higman's lemma (Theorem 3.5.4).
Our main result about $\mathbf{O I}_{d}$ is the following theorem.
Theorem 3.8.2. The category $\mathbf{O I}_{d}$ is Gröbner.
Proof. Let $n$ be a non-negative integer, and regard $x=[n]$ as an object of $\mathcal{C}$. Pick $(f, g) \in$ $\operatorname{Hom}_{\mathcal{e}}([n],[m])$. Define $h:[m] \rightarrow \Sigma$ to be the function which is 0 on the image of $f$, and equal to $g$ on the complement of the image of $f$. One can recover $(f, g)$ from $h$ since $f$ is required to be order-preserving and injective. This construction therefore defines an isomorphism of posets $i: S_{x} \rightarrow \mathcal{L}$. It follows that $\left|\mathcal{C}_{x}\right|$ is noetherian. Furthermore, putting the lexicographic order on $\mathcal{L}$ (using the standard order on $\Sigma$ ) gives a lift of the functor $S_{x}$ to a well-ordering. Thus $\mathcal{C}$ is Gröbner.

Remark 3.8.3. The results about OI can be made more transparent with the following observation: the set of order-increasing injections $f:[n] \rightarrow[m]$ is naturally in bijection with monomials in $x_{0}, \ldots, x_{n}$ of degree $m-n$ by assigning the monomial $\mathbf{m}_{f}=\prod_{i=0}^{n} x_{i}^{f(i+1)-f(i)-1}$ using the convention $f(0)=0$ and $f(n+1)=m+1$. Given $g:[n] \rightarrow\left[m^{\prime}\right]$, there is a morphism $h:[m] \rightarrow\left[m^{\prime}\right]$ with $g=h f$ if and only if $\mathbf{m}_{f}$ divides $\mathbf{m}_{g}$. Thus the monomial subrepresentations of $P_{n}$ are in bijection with monomial ideals in the polynomial ring $\mathbf{k}\left[x_{0}, \ldots, x_{n}\right]$.

In fact, this also shows that finitely generated OI-modules (and hence finitely generated FI-modules) have eventually polynomial growth when $\mathbf{k}$ is a field, i.e., the function $d \mapsto$ $\operatorname{dim}_{\mathbf{k}} M([d])$ is a polynomial for $d \gg 0$. This can also be put in the general framework of "lingual categories" as in [SS1], but we will omit that discussion.

Theorem 3.8.4. The forgetful functor $\Phi: \mathbf{O I}_{d} \rightarrow \mathbf{F I}_{d}$ satisfies property $(\mathrm{F})$. In particular, $\mathbf{F I}_{d}$ is quasi-Gröbner.
Proof. Let $x=[n]$ be a given object of $\mathbf{F I}_{d}$. If $y$ is any totally ordered set, then any morphism $f: x \rightarrow y$ can be factored as $x \xrightarrow{\sigma} x \xrightarrow{f^{\prime}} y$, where $\sigma$ is a permutation and $f^{\prime}$ is order-preserving. It follows that we can take $y_{1}, \ldots, y_{n!}$ to all be $[n]$, and $f_{i}: x \rightarrow \Phi\left(y_{i}\right)$ to be the $i$ th element of the symmetric group $S_{n}$ (under any enumeration). This establishes the claim. Since $\mathbf{O I}_{d}$ is Gröbner, this shows that $\mathbf{F I}_{d}$ is quasi-Gröbner.

Corollary 3.8.5. If $\mathbf{k}$ is left-noetherian then $\operatorname{Rep}_{\mathbf{k}}\left(\mathbf{F I}_{d}\right)$ is noetherian.

## 4. Homological stability for symmetric groups

Some of the applications of representation stability and the theory we've been developing was motivated by generalizations of homological stability. So we'll pause to prove homological
stability for the symmetric groups to get a feeling for what sort of result this is and to illustrate the use of spectral sequences in this subject.

See $\S A .3$ for basic definitions and properties of group homology. We have inclusions of symmetric groups $\Sigma_{1} \rightarrow \Sigma_{2} \rightarrow \cdots$ which induce maps on homology

$$
\mathrm{H}_{i}\left(\Sigma_{1}\right) \rightarrow \mathrm{H}_{i}\left(\Sigma_{2}\right) \rightarrow \cdots
$$

for each $i$, and Nakaoka's stability theorem says that the map $\mathrm{H}_{i}\left(\Sigma_{n-1}\right) \rightarrow \mathrm{H}_{i}\left(\Sigma_{n}\right)$ is an isomorphism for $n>2 i$. We will follow Kerz's proof [Ke] of this fact.
4.1. The complex of injective words. Let $m$ be a fixed positive integer. An injective word in the alphabet $[m]$ is a sequence $\left(i_{1}, \ldots, i_{n}\right)$ of elements in $[m]$, such that $i_{j} \neq i_{k}$ for $j \neq k$. Let $C_{n}(m)$ be the free abelian group with basis indexed by the injective words of length $n$ in the alphabet $[m]$. By convention, $C_{0}(m)=\mathbf{Z}$ with 1 basis element given by the empty sequence. We will abuse notation and use $\left(i_{1}, \ldots, i_{n}\right)$ to denote the basis vector in $C_{n}(m)$. Define a map

$$
\begin{aligned}
d: C_{n}(m) & \rightarrow C_{n-1}(m) \\
\left(i_{1}, \ldots, i_{n}\right) & \mapsto \sum_{j=1}^{n}(-1)^{j+1}\left(i_{1}, \ldots, \hat{i}_{j}, \ldots, i_{n}\right)
\end{aligned}
$$

where the hat denotes that we are omitting that term. A quick check shows that $d^{2}=0$, so we have a chain complex $C_{*}(m)$.

A scalar multiple of a sequence is called a term. If $x \in C_{n}(m)$ is a sum of sequences, then a number $i \in[m]$ is said to appear in $x$ if it is an element in one of the sequences. Finally, we define a partially defined product on $C_{n}(m)$ as follows. If $x \in C_{n}(m)$ and $x^{\prime} \in C_{\ell}(m)$ are terms with $x=N\left(i_{1}, \ldots, i_{n}\right)$ and $x^{\prime}=M\left(i_{1}^{\prime}, \ldots, i_{\ell}^{\prime}\right)$, then we define the product $x x^{\prime}=N M\left(i_{1}, \ldots, i_{n}, i_{1}^{\prime}, \ldots, i_{\ell}^{\prime}\right)$ assuming that this is still an injective word. We extend this to sums whenever all of the corresponding products make sense. When the product is defined for $x \in C_{n}(m)$ and $x^{\prime} \in C_{\ell}(m)$, it satisfies the Leibniz rule in the sense that

$$
d\left(x x^{\prime}\right)=d(x) x^{\prime}+(-1)^{n} x d\left(x^{\prime}\right)
$$

Lemma 4.1.1. If $x \in C_{n}(m)$ is a cycle and there is a number $i$ not appearing in $x$, then $x$ is a boundary.

Proof. By the Leibniz rule, we have $d((i) x)=x-(i) d(x)$. The second quantity is $x$ since $d(x)=0$.
Theorem 4.1.2 (Farmer). $\mathrm{H}_{i}\left(C_{*}(m)\right)=0$ for $i<m$.
Proof. This is clear for $i=0$. To prove the rest of the result, we do induction on $m$, starting with $m=2$. In that case, we have

$$
C_{2}(2) \rightarrow C_{1}(2) \rightarrow C_{0}(2) .
$$

The kernel of the second map is generated by $(1)-(2)=d(1,2)$, so $\mathrm{H}_{1}\left(C_{*}(2)\right)=0$.
Now let $c \in C_{n}(m)$ be a cycle with $n<m$. If there is some letter not appearing in $c$, then we can apply Lemma 4.1.1 and conclude $c$ is a boundary. It will suffice to show that there is some letter not appearing in $c-d(w)$ for some $w \in C_{n+1}(m)$, so we now show how to construct such a $w$.

Suppose there is a letter $\alpha$ that appears as the first entry in some basis element in $c$. Then write

$$
c=\sum_{j}(\alpha) c_{j}+c^{\prime}
$$

where the $c_{j}$ are terms and $\alpha$ does not appear as the first entry in $c^{\prime}$. Since $n<m$, for each $j$ there is a letter $\alpha_{j} \neq \alpha$ that does not appear in $c_{j}$. Then we have

$$
\begin{aligned}
c-d\left(\sum_{j}\left(\alpha_{j}, \alpha\right) c_{j}\right) & =c-\sum_{j}\left((\alpha) c_{j}-\left(\alpha_{j}\right) d\left((\alpha) c_{j}\right)\right) \\
& =c^{\prime}+\sum_{j}\left(\left(\alpha_{j}\right) c_{j}-\left(\alpha_{j}, \alpha\right) d\left(c_{j}\right)\right)
\end{aligned}
$$

and $\alpha$ does not appear as the first entry in the last expression. By induction on $i$, we will now show that we can add a boundary to the above expression so that $\alpha$ does not appear in the first $i$ entries.

So suppose $\alpha$ does not appear in the first $i$ entries of the cycle $c$. Write

$$
c=\sum_{j} s_{j}(\alpha) c_{j}+c^{\prime}
$$

where $\ell\left(s_{j}\right)=i, s_{j}$ and $c_{j}$ do not contain $\alpha$, the $c_{j}$ are distinct for different $j$, and $c^{\prime}$ does not contain $\alpha$ in the first $i+1$ entries. Then we have

$$
\begin{aligned}
0=d(c) & =\sum_{j}\left(d\left(s_{j}\right)(\alpha) c_{j}+(-1)^{i} s_{j} d\left((\alpha) c_{j}\right)\right)+d\left(c^{\prime}\right) \\
& =\sum_{j}\left(d\left(s_{j}\right)(\alpha) c_{j}+(-1)^{i} s_{j} c_{j}+(-1)^{i+1} s_{j}(\alpha) d\left(c_{j}\right)\right)+d\left(c^{\prime}\right) .
\end{aligned}
$$

By our assumption, the only summands where $\alpha$ appears in the $i$ th entry are $d\left(s_{j}\right)(\alpha) c_{j}$, so $\sum_{j} d\left(s_{j}\right)(\alpha) c_{j}=0$. Furthermore, the $c_{j}$ are distinct, so we conclude that $d\left(s_{j}\right)=0$ for all $j$. Now, $\ell\left(s_{j}\right)=i<m-1-\ell\left(c_{j}\right)$ and $s_{j}$ uses only at most $m-1-\ell\left(c_{j}\right)$ letters, so by our induction hypothesis, we conclude there exists $s_{j}^{\prime}$ using these same $m-1-\ell\left(c_{j}\right)$ letters such that $d\left(s_{j}^{\prime}\right)=s_{j}$. In particular, the element

$$
z=\sum_{j} s_{j}^{\prime}(\alpha) c_{j}
$$

is well-defined, and we have

$$
\begin{aligned}
c-d(z) & =c-\sum_{j}\left(s_{j}(\alpha) c_{j}+(-1)^{i+1} s_{j}^{\prime} c_{j}+(-1)^{i+2} s_{j}^{\prime}(\alpha) d\left(c_{j}\right)\right) \\
& =c^{\prime}-\sum_{j}\left((-1)^{i+1} s_{j}^{\prime} c_{j}+(-1)^{i} s_{j}^{\prime}(\alpha) d\left(c_{j}\right)\right)
\end{aligned}
$$

Again, $\alpha$ does not appear in the first $i+1$ entries of the last expression. Finally, once we have $i=n$, we have found the desired element $c-d(w)$.

### 4.2. Nakaoka's theorem.

Theorem 4.2.1 (Nakaoka). The map $\mathrm{H}_{i}\left(\Sigma_{n-1}\right) \rightarrow \mathrm{H}_{i}\left(\Sigma_{n}\right)$ is an isomorphism for $n>2 i$.
Proof. We go by induction on $n$. For $n=3$, we only consider $i=1$, in which case we have

$$
\mathrm{H}_{1}\left(\Sigma_{2}\right)=\mathrm{H}_{1}\left(\Sigma_{3}\right)=\mathbf{Z} / 2
$$

is the abelianization of the group, and the induced map is the identity.
For $n \geq 4$, define a new complex $C_{*}^{\prime}(n)$ by $C_{i}^{\prime}(n)=C_{i+1}(n)$ for $i \geq 0$. Then we have removed the term $C_{0}(n)$, so $\mathrm{H}_{0}\left(C_{*}^{\prime}(n)\right)=\mathbf{Z}$ and $\mathrm{H}_{i}\left(C_{*}^{\prime}(n)\right)=0$ for $0<i<n-1$.

Note that $\Sigma_{n}$ acts on the $C_{i}^{\prime}(n)$ by permuting the letters and that the differential $d$ is equivariant for these actions. We can define the homology of a group with coefficients in a chain complex (this is a special case of hyperhomology), but we won't need the precise definition. We just need to know that there are two spectral sequences which both converge to this value [Br, §VII.5]:

$$
\begin{array}{r}
\mathrm{E}_{p, q}^{2}=\mathrm{H}_{p}\left(\Sigma_{n} ; \mathrm{H}_{q}\left(C_{*}^{\prime}(n)\right)\right) \Longrightarrow \mathrm{H}_{p+q}\left(\Sigma_{n} ; C_{*}^{\prime}(n)\right) \\
\mathrm{E}_{p, q}^{1}=\mathrm{H}_{q}\left(\Sigma_{n} ; C_{p}^{\prime}(n)\right) \Longrightarrow \mathrm{H}_{p+q}\left(\Sigma_{n} ; C_{*}^{\prime}(n)\right) .
\end{array}
$$

Consider the first spectral sequence. As mentioned before, $\mathrm{H}_{q}\left(C_{*}^{\prime}(n)\right)=0$ for $0<q<n-1$. Hence the only terms of degree $d$ with $0 \leq d<n-1$ is $\mathrm{H}_{d}\left(\Sigma_{n} ; \mathrm{H}_{0}\left(C_{*}^{\prime}(n)\right)\right)=\mathrm{H}_{d}\left(\Sigma_{n}\right)$. Furthermore, on each page, there is no nonzero differential to these terms, so

$$
\mathrm{H}_{d}\left(\Sigma_{n} ; C_{*}^{\prime}(n)\right)=\mathrm{E}_{d, 0}^{\infty}=\mathrm{H}_{d}\left(\Sigma_{n}\right) \quad(0 \leq d<n-1)
$$

Now consider the second spectral sequence. The $\Sigma_{n}$-representation $C_{p}^{\prime}(n)$ is the permutation representation on the set of injective words of length $p+1$. This has a transitive action of $\Sigma_{n}$ with stabilizer subgroup $\Sigma_{n-p-1}$, so it is isomorphic to the induced module $\operatorname{Ind} \Sigma_{\Sigma_{n-p-1}}^{\Sigma_{n}} \mathbf{Z}$, so

$$
\mathrm{H}_{q}\left(\Sigma_{n} ; C_{p}^{\prime}(n)\right)=\mathrm{H}_{q}\left(\Sigma_{n-p-1}\right)
$$

by Shapiro's lemma. The map $C_{p}^{\prime}(n) \rightarrow C_{p-1}^{\prime}(n)$ corresponding to forgetting the $j$ th entry becomes the map

$$
\mathrm{H}_{q}\left(\Sigma_{n-p-1}\right) \rightarrow \mathrm{H}_{q}\left(\Sigma_{n-p}\right)
$$

induced by the inclusion of $\Sigma_{n-p-1}$ into $\Sigma_{n-p}$ as the permutations that fix $j$. These all differ by conjugation by an element of $\Sigma_{n-p}$, so are the same for all $j$. Hence the differential induces a map on homology which is an alternating sum of the same map, and we see that it is 0 if $p$ is odd and the inclusion map if $p$ is even. Hence the spectral sequence looks like

$$
\begin{aligned}
& \mathrm{H}_{2}\left(\Sigma_{n-1}\right)<{ }^{0} \mathrm{H}_{2}\left(\Sigma_{n-2}\right) \longleftarrow \mathrm{H}_{2}\left(\Sigma_{n-3}\right) \leftarrow_{\longleftarrow}^{0} \cdots \quad \mathrm{H}_{2}\left(\Sigma_{2}\right) \quad \mathrm{H}_{2}\left(\Sigma_{1}\right) \quad 0 \\
& \mathrm{H}_{1}\left(\Sigma_{n-1}\right) \stackrel{0}{\leftarrow} \mathrm{H}_{1}\left(\Sigma_{n-2}\right) \longleftarrow \mathrm{H}_{1}\left(\Sigma_{n-3}\right) \leftarrow 0<\cdots \quad \mathrm{H}_{1}\left(\Sigma_{2}\right) \quad \mathrm{H}_{1}\left(\Sigma_{1}\right) \quad 0 \\
& \mathrm{H}_{0}\left(\Sigma_{n-1}\right) \stackrel{0}{\longleftarrow} \mathrm{H}_{0}\left(\Sigma_{n-2}\right) \longleftarrow \mathrm{H}_{0}\left(\Sigma_{n-3}\right) \stackrel{0}{\longleftarrow} \cdots \quad \mathrm{H}_{0}\left(\Sigma_{2}\right) \quad \mathrm{H}_{0}\left(\Sigma_{1}\right) \quad \mathbf{Z}
\end{aligned}
$$

where the maps alternate between being 0 and coming from the standard inclusion.
In the second page, the leftmost column stays the same, i.e., $\mathrm{E}_{0, i}^{2}=\mathrm{H}_{i}\left(\Sigma_{n-1}\right)$ for all $i \geq 0$. Now suppose that $n>2 i$. The terms with the same degree $i$ are of the form $\mathrm{H}_{i-k}\left(\Sigma_{n-k-1}\right)$ and so by induction, the maps from the previous symmetric group or to the next symmetric
group are isomorphisms on homology. This implies that the only term of degree $i$ on the $\mathrm{E}^{2}$ page is $\mathrm{E}_{i, 0}^{2}$. Finally, we claim that $\mathrm{E}_{0, i}^{\infty}=\mathrm{E}_{0, i}^{2}$ since the only term that can map to it on the $\mathrm{E}^{r}$ page is $\mathrm{E}_{r, i-r+1}^{r}$, and this is also 0 by the induction hypothesis. In particular, this implies that we have an identification $\mathrm{H}_{i}\left(\Sigma_{n-1}\right)=\mathrm{H}_{i}\left(\Sigma_{n}\right)$.

It remains to show that this is induced by the inclusion map, but we will omit this detail since it involves a more detailed analysis of the spectral sequence.

Remark 4.2.2. Nakaoka's theorem remains true if we replace $\mathbf{Z}$ by any coefficient ring $\mathbf{k}$ (as long as $\Sigma_{n}$ acts by the identity on $\mathbf{k}$ ). For example, the proof above is not sensitive to the coefficients we have chosen except in the calculation of $\mathrm{H}_{1}\left(\Sigma_{n}\right)$, but the general statement is that it is $\mathbf{Z} / 2 \otimes_{\mathbf{z}} \mathbf{k}$ for $n \geq 2$.
4.3. Twisted homological stability. Let $M$ be an FI-module over some coefficient ring $\mathbf{k}$. Let $M(n)$ denote the value of $M$ on the set $[n]$. This carries an action of the symmetric group $\Sigma_{n}$. Let $\iota:[n] \rightarrow[n+1]$ be the standard inclusion, i.e., $\iota(i)=i$ for all $i$ and let $f: \Sigma_{n} \rightarrow \Sigma_{n+1}$ also be the standard inclusion. Then for all $\sigma \in \Sigma_{n}$ and $m \in M(n)$, we have $\iota_{*}(\sigma m)=f(\sigma) \iota_{*}(m)$ since $\iota \circ \sigma=f(\sigma)$ as morphisms in FI. So by functoriality of group homology, we have induced maps for all $i$ :

$$
\mathrm{H}_{i}\left(\Sigma_{n} ; M(n)\right) \rightarrow \mathrm{H}_{i}\left(\Sigma_{n+1} ; M(n+1)\right) .
$$

When $M$ is the constant functor, i.e., $M_{n}=\mathbf{k}$ and all morphisms map to the identity, then we considered these maps above and Nakaoka's theorem tells us that they are isomorphisms for $n \gg i$ (more precisely, $n>2 i$ ). We now extend this to any finitely generated FI-module over a noetherian ring $\mathbf{k}$.

We say that the modules $M(n)$ with the maps $\iota_{*}: M(n) \rightarrow M(n+1)$ satisfy twisted homological stability if the maps on homology are isomorphisms for $n \gg i$.

So now assume that $\mathbf{k}$ is noetherian.
Lemma 4.3.1. The principal projective module $P_{d}$ satisfies twisted homological stability.
Proof. For each $n, P_{d}(n)$ is the permutation representation on the set of injective words of length $d$ in the alphabet $[n]$. As discussed in the proof of Nakaoka's theorem, this is the induced representation $\operatorname{Ind}_{\Sigma_{n-d}}^{\Sigma_{n}} \mathbf{k}$ and, via Shapiro's lemma, the map

$$
\mathrm{H}_{i}\left(\Sigma_{n} ; P_{d}(n)\right) \rightarrow \mathrm{H}_{i}\left(\Sigma_{n+1} ; P_{d}(n+1)\right)
$$

can be rewritten as

$$
\mathrm{H}_{i}\left(\Sigma_{n-d} ; \mathbf{k}\right) \rightarrow \mathrm{H}_{i}\left(\Sigma_{n+1-d} ; \mathbf{k}\right)
$$

which is the map coming from the standard inclusion. Hence twisted homological stability for $P_{d}$ reduces to Nakaoka's stability theorem.

Corollary 4.3.2. Any finitely generated FI-module $M$ satisfies twisted homological stability.
Proof. Consider the following two statements:
$\left(\mathrm{A}_{i}\right)$ : For any finitely generated FI-module $M$, the map

$$
\mathrm{H}_{i}\left(\Sigma_{n} ; M(n)\right) \rightarrow \mathrm{H}_{i}\left(\Sigma_{n+1} ; M(n+1)\right)
$$

is surjective for $n \gg i$.
$\left(\mathrm{B}_{i}\right)$ : For any finitely generated FI-module $M$, the map

$$
\mathrm{H}_{i}\left(\Sigma_{n} ; M(n)\right) \rightarrow \mathrm{H}_{i}\left(\Sigma_{n+1} ; M(n+1)\right)
$$

is an isomorphism for $n \gg i$.
We will show that $\left(\mathrm{B}_{i-1}\right)$ implies $\left(\mathrm{A}_{i}\right)$ for all $i$, and that $\left(\mathrm{A}_{i}\right)$ and $\left(\mathrm{B}_{i-1}\right)$ together imply $\left(\mathrm{B}_{i}\right)$ for all $i$. With the convention that $\mathrm{H}_{-1}\left(\Sigma_{n} ; M(n)\right)=0$ for all $n$, the statements $\left(\mathrm{A}_{-1}\right)$ and $\left(B_{-1}\right)$ are vacuously true.

Let $M$ be a finitely generated FI-module generated by homogeneous elements $x_{1}, \ldots, x_{r}$ and of degrees $d_{1}, \ldots, d_{r}$. Let $P=P_{d_{1}} \oplus \cdots \oplus P_{d_{r}}$ so that we have a canonical surjection $P \rightarrow M$, and let $K$ be the kernel. Then $K$ is finitely generated by the noetherian property for FI-modules. Fix $i$ and consider the long exact sequence on homology:


By Lemma 4.3.1, $\beta$ and $\varepsilon$ are isomorphisms for $n \gg i$.
Suppose $\left(\mathrm{B}_{i-1}\right)$ is true. Then $\delta$ is an isomorphism for $n \gg i$. Then the first part of the four lemma implies that $\gamma$ is surjective for $n \gg i$. Since $M$ was arbitrary, $\left(\mathrm{A}_{i}\right)$ is true.

Now suppose that $\left(\mathrm{A}_{i}\right)$ and $\left(\mathrm{B}_{i-1}\right)$ are both true, then $\alpha$ and $\gamma$ are surjective, and $\delta$ is an isomorphism for $n \gg i$. The second part of the four lemma then says that $\gamma$ is injective. Since $M$ was arbitrary, $\left(\mathrm{B}_{i}\right)$ is true.

## 5. Representation stability for configuration spaces

5.1. Definitions. Let $X$ be a topological space. For an integer $n \geq 0$, we set

$$
\operatorname{Conf}_{n}(X)=\left\{\left(x_{1}, \ldots, x_{n}\right) \in X^{n} \mid x_{i} \neq x_{j} \text { for } i \neq j\right\}
$$

Alternatively, this is the space of injective functions $[n] \rightarrow X$. It will be convenient to define $\operatorname{Conf}_{S}(X)$ to be the space of injective functions $S \rightarrow X$ for any finite set $S$.

Given an injection $S \rightarrow T$, we get a continuous function $\operatorname{Conf}_{T}(X) \rightarrow \operatorname{Conf}_{S}(X)$ by precomposition. So we have define a functor from the opposite category of FI to the category of topological spaces. To get something linear out of it, we just need to choose some (contravariant) functor from the category of topological spaces to the category of modules over some ring. We will consider singular cohomology with coefficients in a commutative ring $\mathbf{k}$, denoted $\mathrm{H}^{i}(X ; \mathbf{k})$ where $i \geq 0$ is some integer. We won't need to know the details of how this is constructed, but here are two key properties:
(1) $\mathrm{H}^{*}(X ; \mathbf{k})=\bigoplus_{i \geq 0} \mathrm{H}^{i}(X ; \mathbf{k})$ has the structure of a graded ring: given homogeneous elements in degrees $i$ and $j$, their product has degree $i+j$. Furthermore, it is gradedcommutative: if $x \in \mathrm{H}^{i}(X ; \mathbf{k})$ and $y \in \mathrm{H}^{j}(X ; \mathbf{k})$, then $x y=(-1)^{i j} y x$.
(2) Given a continuous function $f: X \rightarrow Y$, we get a corresponding $\mathbf{k}$-linear map $f^{*}: \mathrm{H}^{i}(Y ; \mathbf{k}) \rightarrow \mathrm{H}^{i}(X ; \mathbf{k})$, and these compose correctly: $(f \circ g)^{*}=g^{*} \circ f^{*}$. Furthermore, taking the sum of these maps over all $i, f^{*}$ is a ring homomorphism.
(3) (Künneth formula) If $\mathbf{k}$ is a field, then we have an isomorphism of graded rings

$$
\mathrm{H}^{*}(X \times Y ; \mathbf{k}) \cong \mathrm{H}^{*}(X ; \mathbf{k}) \otimes_{\mathbf{k}} \mathrm{H}^{*}(Y ; \mathbf{k}) .
$$

In particular,

$$
\mathrm{H}^{n}(X \times Y ; \mathbf{k}) \cong \bigoplus_{i=0}^{n} \mathrm{H}^{i}(X ; \mathbf{k}) \otimes_{\mathbf{k}} \mathrm{H}^{n-i}(Y ; \mathbf{k})
$$

In particular, every choice of topological space $X$ and nonnegative integer $i$, we get an FI-module with coefficients in $\mathbf{k}$ which sends a set $S$ to $\mathrm{H}^{i}\left(\operatorname{Conf}_{S}(X) ; \mathbf{k}\right)$ and an injection $f: S \rightarrow T$ to the map $f^{*}: \mathrm{H}^{i}\left(\operatorname{Conf}_{S}(X) ; \mathbf{k}\right) \rightarrow \mathrm{H}^{i}\left(\operatorname{Conf}_{T}(X) ; \mathbf{k}\right)$ which is induced by the continuous function $\operatorname{Conf}_{T}(X) \rightarrow \operatorname{Conf}_{S}(X)$ obtained by precomposing with $f$. Denote this FI-module by $\mathrm{H}^{i}\left(\operatorname{Conf}_{*}(X) ; \mathbf{k}\right)$.

From what we've discussed so far, a natural question to ask is: for which $X$ and $i$ is this FI-module finitely generated? We won't discuss the most general results known, but instead content ourselves with an easy to analyze case.
Example 5.1.1. Let's consider the case $X=\mathbf{R}^{d}$. The cohomology ring of $\operatorname{Conf}_{n}\left(\mathbf{R}^{d}\right)$ is the graded-commutative $\mathbf{k}$-algebra generated by variables $w_{a, b}$ of degree $d-1$ with $1 \leq a \neq b \leq n$ and modulo the relations:

- $w_{a, b}=(-1)^{d} w_{b, a}$,
- $w_{a, b}^{2}=0$,
- $w_{a, b} w_{a, c}+w_{b, c} w_{b, a}+w_{c, a} w_{c, b}=0$ for $a, b, c$ distinct.

The $w_{a, b}$ are elements of $\mathrm{H}^{d-1}\left(\operatorname{Conf}_{n}\left(\mathbf{R}^{d}\right) ; \mathbf{k}\right)$.
Given an injection $f:\{1, \ldots, n\} \rightarrow\{1, \ldots, m\}$, the action of $f^{*}$ is to replace each $w_{a, b}$ with $w_{f(a), f(b)}$. Note that this is well-defined: $f^{*}$ sends each of the relations to another relation for the cohomology ring of $\operatorname{Conf}_{m}\left(\mathbf{R}^{d}\right)$.

Assume now that $d>1$. We can see directly that $\mathrm{H}^{d-1}\left(\operatorname{Conf}_{*}\left(\mathbf{R}^{d}\right) ; \mathbf{k}\right)$ is finitely generated: it is either the permutation representation on the set of 2-element subsets of $[n]$ or a twist of this by the sign character, so is generated by $w_{1,2} \in H^{d-1}\left(\operatorname{Conf}_{2}\left(\mathbf{R}^{d}\right) ; \mathbf{k}\right)$. In general, $H^{i(d-1)}\left(\operatorname{Conf}_{*}\left(\mathbf{R}^{d}\right) ; \mathbf{k}\right)$ is a quotient of $\left(\mathrm{H}^{d-1}\left(\operatorname{Conf}_{*}\left(\mathbf{R}^{d}\right) ; \mathbf{k}\right)\right)^{\boxtimes i}$ and hence is finitely generated (see Exercise 5.1.2).

If $m=1$ and $\mathbf{k}$ is a field, one can show that $\mathrm{H}^{0}\left(\operatorname{Conf}_{n}(\mathbf{R}) ; \mathbf{k}\right)$ has dimension $n!$, and hence cannot be a finitely generated FI-module. This can either be done using the relations above, or noting that $\operatorname{Conf}_{n}(\mathbf{R})$ has $n$ ! connected components, each one indexed by the order in which the points $\left(x_{1}, \ldots, x_{n}\right)$ appear on the line, and the dimension of $\mathrm{H}^{0}$ is the number of connected components.

Exercise 5.1.2. If $M$ and $N$ are FI-modules over $\mathbf{k}$, then define an FI-module $M \boxtimes N$, the pointwise (or Segre) tensor product as follows. On objects, we set $(M \boxtimes N)(S)=$ $M(S) \otimes_{\mathbf{k}} N(S)$, and given an injection $f: S \rightarrow T$, we set $M(S) \otimes_{\mathbf{k}} N(S) \rightarrow M(T) \otimes_{\mathbf{k}} N(T)$ to be the tensor product of the maps $M(f) \otimes N(f)$.

If $M$ and $N$ are finitely generated, show that $M \boxtimes N$ is also finitely generated.
Hint: Reduce to the case when $M$ and $N$ are principal projectives.
5.2. A spectral sequence. For this part, see [To, $\S 2]$. Let $X$ be a real orientable manifold of dimension $d$ and let $\mathbf{k}$ be a field. Fix a positive integer $n$. Given $1 \leq a \leq n$, let $p_{a}: X^{n} \rightarrow X$ be the projection onto the $a$ th factor, i.e., $p_{a}\left(x_{1}, \ldots, x_{n}\right)=x_{a}$. Similarly, given $a \neq b$, let $p_{a, b}: X^{n} \rightarrow X^{2}$ be the map $p_{a, b}\left(x_{1}, \ldots, x_{n}\right)=\left(x_{a}, x_{b}\right)$.

Define a bigraded algebra $A(n)$ over $\mathrm{H}^{*}\left(X^{n} ; \mathbf{k}\right)$ with variables $G_{a, b}(1 \leq a \neq b \leq n)$ modulo the relations

- $G_{a, b}=(-1)^{d} G_{b, a}$,
- $G_{a, b}^{2}=0$,
- $G_{a, b} G_{a, c}+G_{b, c} G_{b, a}+G_{c, a} G_{c, b}=0$ for $a, b, c$ distinct,
- $p_{a}^{*}(x) G_{a, b}=p_{b}^{*}(x) G_{a, b}$ for $a \neq b$ and $x \in \mathrm{H}^{*}(X ; \mathbf{k})$.

Here the bidegree of $\mathrm{H}^{i}\left(X^{n} ; \mathbf{k}\right)$ is $(i, 0)$ and the bidegree of $G_{a, b}$ is $(0, d-1)$.
Let $A(n)_{i, j}$ be the space of bidegree $(i, j)$ elements in $A(n)$. As in [To, Theorem 1], there is a cohomology spectral sequence

$$
\mathrm{E}_{d}^{p, q}=A(n)_{p, q} \Longrightarrow \mathrm{H}^{p+q}\left(\operatorname{Conf}_{n}(X) ; \mathbf{k}\right)
$$

We won't discuss the differentials in this spectral sequence. All we will need to know is that the terms $A(n)_{p, q}$, as $n$ varies, form an FI-module, and that this spectral sequence can be upgraded to a spectral sequence of FI-modules.

First, given an injection $f:[n] \rightarrow[m]$, we get a continuous function $X^{m} \rightarrow X^{n}$ given by $\left(x_{1}, \ldots, x_{m}\right) \mapsto\left(x_{f(1)}, \ldots, x_{f(n)}\right)$, and hence a corresponding map on cohomology $f^{*}: \mathrm{H}^{*}\left(X^{n} ; \mathbf{k}\right) \rightarrow$ $\mathrm{H}^{*}\left(X^{m} ; \mathbf{k}\right)$. Using the Künneth formula $\mathrm{H}^{*}\left(X^{n} ; \mathbf{k}\right) \cong \mathrm{H}^{*}(X ; \mathbf{k})^{\otimes n}$, this can be described as follows. Let $e \in \mathrm{H}^{0}(X ; \mathbf{k})$ be the unit for multiplication. Then $\sum v_{1} \otimes \cdots \otimes v_{n}$ gets sent to $\sum w_{1} \otimes \cdots \otimes w_{m}$ where $w_{i}=e$ if $i$ is not in the image of $f$, and $w_{f(j)}=v_{j}$ otherwise.

Say that $X$ is of finite type if $\mathrm{H}^{*}(X ; \mathbf{k})$ is a finite-dimensional vector space. This holds if $X$ is compact, or Euclidean space, for instance.

Lemma 5.2.1. If $X$ is a finite type and connected manifold, then for each $i$, the FI-module $[n] \mapsto \mathrm{H}^{i}\left(X^{m} ; \mathbf{k}\right)$ is finitely generated.

Proof. Since $X$ is connected, $\mathrm{H}^{0}(X ; \mathbf{k})$ is 1-dimensional with basis $e$. By the Künneth formula,

$$
\mathrm{H}^{i}\left(X^{n} ; \mathbf{k}\right)=\bigoplus_{\substack{\left(j_{1}, \ldots, j_{n}\right) \in \mathbf{Z}_{\geq 0}^{n} \\ i=j_{1}+\cdots+j_{n}}} \mathrm{H}^{j_{1}}(X ; \mathbf{k}) \otimes \cdots \otimes \mathrm{H}^{j_{n}}(X ; \mathbf{k})
$$

At most $i$ of the $j_{k}$ can be positive, and we claim that the FI-module in question is generated in degrees $\leq i$. The point is that any element is a sum of simple tensors, and if there $n>i$, then at least one of the factors must be a multiple of $e$. So it comes from applying an injection from strictly smaller degree. By the finite type assumption, $\mathrm{H}^{i}\left(X^{n} ; \mathbf{k}\right)$ is finite-dimensional for all $n$, so $\sum_{n \leq i} \operatorname{dim} \mathrm{H}^{i}\left(X^{n} ; \mathbf{k}\right)<\infty$.

Theorem 5.2.2 (Church). If $X$ is a finite type, orientable, connected manifold of dimension $d \geq 2$, then for each $i$, the FI-module $\mathrm{H}^{i}\left(\operatorname{Conf}_{*}(X) ; \mathbf{k}\right)$ is finitely generated.

Proof. Combining the previous lemma with Example 5.1.1, we see that $[n] \mapsto A(n)_{p, q}$ is a finitely generated FI-module for all $p, q$. The spectral sequence above is compatible with the FI-structure, and hence the FI-module $\mathrm{H}^{i}\left(\operatorname{Conf}_{*}(X) ; \mathbf{k}\right)$ has a finite filtration whose quotients are computed by the spectral sequence. Namely, they are obtained from the $A(*)_{p, q}$ by taking homology. By the noetherian property for FI-modules, each time we take homology (which is the quotient of a submodule), we get another finitely generated FI-module, and so $\mathrm{H}^{i}\left(\operatorname{Conf}_{*}(X) ; \mathbf{k}\right)$ itself is a finitely generated FI-module.

Next, the symmetric group $\Sigma_{n}$ acts on $\operatorname{Conf}_{n}(X)$ by permuting the points, and this action is free (no point has a nontrivial stabilizer). Let

$$
\operatorname{UConf}_{n}(X)=\operatorname{Conf}_{n}(X) / \Sigma_{n}
$$

be the quotient space. Its points are unordered $n$-tuples of points in $X$. We have a CartanLeray spectral sequence (dual of $\left[\mathrm{Mc}, \S 8^{\text {bis }} .2\right]$ ):

$$
\begin{equation*}
\mathrm{E}_{2}^{p, q}=\mathrm{H}^{p}\left(\Sigma_{n} ; \mathrm{H}^{q}\left(\operatorname{Conf}_{n}(X) ; \mathbf{k}\right)\right) \Longrightarrow \mathrm{H}^{p+q}\left(\operatorname{UConf}_{n}(X) ; \mathbf{k}\right) . \tag{5.2.3}
\end{equation*}
$$

If $n$ ! is invertible in $\mathbf{k}$, then $\mathrm{H}^{p}\left(\Sigma_{n} ; \mathrm{H}^{q}\left(\operatorname{Conf}_{n}(X) ; \mathbf{k}\right)\right)=0$ whenever $p>0$, and so there are no nontrivial differentials in the spectral sequence, so we get the following result.
Proposition 5.2.4. If $n!$ is invertible in $\mathbf{k}$, then $\mathrm{H}^{i}\left(\operatorname{Conf}_{n}(X) ; \mathbf{k}\right)^{\Sigma_{n}}=\mathrm{H}^{i}\left(\operatorname{UConf}_{n}(X) ; \mathbf{k}\right)$.
Proposition 5.2.5. Given any FI-module $M$, the direct sum $\bigoplus_{n \geq 0}\left(M_{n}\right)_{\Sigma_{n}}$ has the structure of a $\mathbf{k}[t]$-module where multiplication by $t$ on $\left(M_{n}\right)_{\Sigma_{n}}$ is the action of any injection $[n] \rightarrow$ $[n+1]$. If $M$ is finitely generated, then so is this $\mathbf{k}[t]$-module.

Corollary 5.2.6 (Church). If $X$ is a finite type, orientable, connected manifold of dimension $\geq 2$, then for each $i$, the function $n \mapsto \operatorname{dim}_{\mathbf{Q}} \mathrm{H}^{i}\left(\operatorname{UConf}_{n}(X) ; \mathbf{Q}\right)$ is constant for $n \gg 0$.
Proof. We have $\mathrm{H}^{i}\left(\operatorname{Conf}_{n}(X) ; \mathbf{Q}\right)^{\Sigma_{n}}=\mathrm{H}^{i}\left(\operatorname{UConf}_{n}(X) ; \mathbf{Q}\right)$ by Proposition 5.2.4. By semisimplicity, we also have

$$
\operatorname{dim}_{\mathbf{Q}} \mathrm{H}^{i}\left(\operatorname{Conf}_{n}(X) ; \mathbf{Q}\right)^{\Sigma_{n}}=\operatorname{dim}_{\mathbf{Q}} \mathrm{H}^{i}\left(\operatorname{Conf}_{n}(X) ; \mathbf{Q}\right)_{\Sigma_{n}} .
$$

Now the result follows from Proposition 5.2 .5 and the structure theorem for finitely generated modules over $\mathbf{k}[t]$.

Naturally, one can ask what happens for fields of positive characteristic. Unfortunately, we have only proven a result about the $\Sigma_{n}$-homology of an FI-module, and not its cohomology, which behaves very differently. The situation can be analyzed, though it requires a finer understanding of the structure of FI-modules. This might be discussed later in the semester, time permitting. See [NS] for details.

## 6. Algebraic geometry from tensors

6.1. Review of Zariski topology. Let $R$ be a commutative ring. Let $\operatorname{Spec}(R)$ be the set of prime ideals in $R$. Given an ideal $I \subset R$, let $V(I)=\{\mathfrak{p} \in \operatorname{Spec}(R) \mid \mathfrak{p} \supseteq I\}$. $\operatorname{Spec}(R)$ is a topological space if we declare the subsets of the form $V(I)$ to be closed, as follows from this list of properties:

- $\bigcap_{s \in S} V\left(I_{s}\right)=V\left(\sum_{s \in S} I_{s}\right)$,
- $V(I J)=V(I) \cup V(J)$,
- $V(0)=\operatorname{Spec}(R)$,
- $V(R)=\varnothing$.

This is the Zariski topology on $\operatorname{Spec}(R)$.
Note that $V(I)=V(\sqrt{I})$ where $\sqrt{I}=\left\{f \in R \mid f^{n} \in I\right.$ for some $\left.n>0\right\}$ is the radical of $I$. Also, the only primes such that $\{\mathfrak{p}\}$ are closed sets are maximal ideals. We define a variant: $\operatorname{MSpec}(R)$ is the set of maximal ideals with a similarly defined topology. This is usually not well-behaved, though it will be useful in some things we do.

We will be primarily concerned with the case that $R=\mathbf{k}\left[x_{1}, \ldots, x_{n}\right]$ where $\mathbf{k}$ is an algebraically closed field. In this case, Hilbert's nullstellensatz tells us that the maximal ideals of $R$ are all of the form $\left(x_{1}-\alpha_{1}, \ldots, x_{n}-\alpha_{n}\right)$ where $\alpha_{i} \in \mathbf{k}$. So geometrically, we can think of $\operatorname{Spec}(R)$ as the vector space $\mathbf{k}^{n}$ with some additional points coming from non-maximal prime ideals. The picture of $\mathbf{k}^{n}$ usually gives good intuition, so if you haven't seen algebraic geometry before, it will usually be enough to rely on it. We call sets of the form $V(I)$ algebraic varieties, or just varieties.

Our main concern will be with recovering the ideal $I$ that defines a closed subset $V(I)$. However, by the above remark, $I$ cannot be recovered since $V(I)=V(\sqrt{I})$. Hence we will be concerned with finding $\sqrt{I}$, which is well-defined.

If $f_{1}, \ldots, f_{r}$ are polynomials generating $I,{ }^{4}$ then in the $\mathbf{k}^{n}$ model, the closed points of $V(I)$ are those points $\left(\alpha_{1}, \ldots, \alpha_{n}\right)$ such that $f_{i}\left(\alpha_{1}, \ldots, \alpha_{n}\right)=0$ for all $i$. Hence we will say that $f_{1}, \ldots, f_{r}$ define $V(I)$ set-theoretically (since they determine it as a subset), and we will say they generate the full ideal of $V(I)$ if $I=\sqrt{I}$ is radical. The meaning of this is that any function that vanishes on all closed points of $V(I)$ must be in $I$.

Remark 6.1.1. In the case $\mathbf{k}=\mathbf{C}$ is the field of complex numbers, we can also put the Euclidean topology on $\mathbf{C}^{n}$. Given polynomials $f_{1}, \ldots, f_{r}$, we get a continuous function $\mathbf{C}^{n} \rightarrow \mathbf{C}^{r}$ defined by $\left(\alpha_{1}, \ldots, \alpha_{n}\right) \mapsto\left(f_{1}(\alpha), \ldots, f_{r}(\alpha)\right)$, and so the preimage of 0 is closed in the Euclidean topology. This implies that the Euclidean topology is a refinement of the Zariski topology. In particular, the Zariski closure of a set always contains the Euclidean closure, and in particular, is closed under limits. One can make the notion of being closed under limits precise outside of the case $\mathbf{k}=\mathbf{C}$, though we won't make much use of it.
6.2. Border rank. Let $V_{1}, \ldots, V_{n}$ be vector spaces over an algebraically closed field $\mathbf{k}$ and $\mathbf{V}=V_{1} \otimes \cdots \otimes V_{n}$. An element of $\mathbf{V}$ of the form $v_{1} \otimes \cdots \otimes v_{n}$ with $v_{i} \in V_{i}$ is a simple tensor, and a general element has rank $\leq r$ if it can be expressed as a sum of $r$ simple tensors. Finally, an element in the Zariski closure of rank $\leq r$ tensors has border rank $\leq r$. They form an algebraic variety.

Example 6.2.1. The following example illustrates why the notion of border rank is needed, i.e., why tensor rank is not semicontinuous.

Let $n=3$ and $\operatorname{dim} V_{i}=2$, and for simplicity, take $\mathbf{k}=\mathbf{C}$. Pick bases $\left\{a_{1}, a_{2}\right\},\left\{b_{1}, b_{2}\right\},\left\{c_{1}, c_{2}\right\}$ of $V_{1}, V_{2}, V_{3}$, respectively. The element

$$
\begin{aligned}
v & =a_{1} \otimes b_{1} \otimes c_{1}+a_{1} \otimes b_{1} \otimes c_{2}+a_{1} \otimes b_{2} \otimes c_{1}+a_{2} \otimes b_{1} \otimes c_{1} \\
& =a_{1} \otimes b_{1} \otimes\left(c_{1}+c_{2}\right)+a_{1} \otimes b_{2} \otimes c_{1}+a_{2} \otimes b_{1} \otimes c_{1}
\end{aligned}
$$

has rank $\leq 3$, and it can be shown that it does not have rank $\leq 2$, so the rank is exactly 3 . However,

$$
v=\lim _{\varepsilon \rightarrow 0} \frac{1}{\varepsilon}\left((\varepsilon-1) a_{1} \otimes b_{1} \otimes c_{1}+\left(a_{1}+\varepsilon a_{2}\right) \otimes\left(b_{1}+\varepsilon b_{2}\right) \otimes\left(c_{1}+\varepsilon c_{2}\right)\right)
$$

shows that $v$ has border rank $\leq 2$.
Our first main goal is to prove the following theorem:
Theorem 6.2.2 (Draisma-Kuttler). For each $r$, there is a constant $C(r)$ such that the variety of border rank $\leq r$ tensors is cut out by polynomials of degree $\leq C(r)$. The constant is independent of $n$ and also the $V_{i}$.

Example 6.2.3. A familiar case is when $n=2$. Then $\mathbf{V}$ can be thought of as the space of $a \times b$ matrices where $a=\operatorname{dim} V_{1}$ and $b=\operatorname{dim} V_{2}$. In this case, a simple tensor is a rank 1 matrix, and rank $r$ in the tensor sense corresponds to rank $r$ in the matrix sense. Border rank $\leq r$ is equivalent to rank $\leq r$. The polynomials in this case are the determinants of all $(r+1) \times(r+1)$ submatrices, and so they are of degree $r+1$. When $n \geq 3$, the situation is more complicated.

Remark 6.2.4. For concreteness, consider the case $r=1$ and $n=3$ (the case of general $n$ follows in the same way, though the indexing is more cumbersome). Let $V_{1}, V_{2}, V_{3}$ be vector

[^3]spaces of dimensions $d_{1}, d_{2}, d_{3}$ and pick a basis $e_{1}, \ldots, e_{d_{i}}$ for $V_{i}$. Let $e_{i, j, k}=e_{i} \otimes e_{j} \otimes e_{k}$, which forms a basis as we vary $i, j, k$, and let $x_{i, j, k}$ be the dual basis. Suppose we're given a rank 1 tensor of the form
$$
v=\left(\sum_{i} \lambda_{i} e_{i}\right) \otimes\left(\sum_{j} \mu_{j} e_{j}\right) \otimes\left(\sum_{k} \nu_{k} e_{k}\right) .
$$

Then $x_{i, j, k}(v)=\lambda_{i} \mu_{j} \nu_{k}$. In particular, we have the relations

$$
x_{i, j, k} x_{i^{\prime}, j^{\prime}, k^{\prime}}=x_{a, b, c} x_{a^{\prime}, b^{\prime}, c^{\prime}}
$$

whenever $\left\{i, i^{\prime}\right\}=\left\{a, a^{\prime}\right\},\left\{j, j^{\prime}\right\}=\left\{b, b^{\prime}\right\}$, and $\left\{k, k^{\prime}\right\}=\left\{c, c^{\prime}\right\}$. We single out a particular class of these:

$$
\begin{equation*}
x_{i, j, k} x_{i^{\prime}, j^{\prime}, k^{\prime}}=x_{\min \left(i, i^{\prime}\right), \min \left(j, j^{\prime}\right), \min \left(k, k^{\prime}\right)} x_{\max \left(i, i^{\prime}\right), \max \left(j, j^{\prime}\right), \max \left(k, k^{\prime}\right)} . \tag{6.2.4a}
\end{equation*}
$$

By replacing any product of two variables with the product on the right side, we conclude that every polynomial of degree $m$ on the rank 1 locus is the span of monomials of the form

$$
x_{i_{1}, j_{1}, k_{1}} \cdots x_{i_{m}, j_{m}, k_{m}} \quad\left(i_{1} \leq \cdots \leq i_{m}, \quad j_{1} \leq \cdots \leq j_{m}, \quad k_{1} \leq \cdots \leq k_{m}\right)
$$

The conditions on the indices are independent of one another, so we see that this is the same as the dimension of $\operatorname{Sym}^{m}\left(V_{1}^{*}\right) \otimes \operatorname{Sym}^{m}\left(V_{2}^{*}\right) \otimes \operatorname{Sym}^{m}\left(V_{3}^{*}\right)$ (for a bijection, send the product above to $\left.\left(e_{i_{1}}^{*} \cdots e_{i_{m}}^{*}\right) \otimes\left(e_{j_{1}}^{*} \cdots e_{j_{m}}^{*}\right) \otimes\left(e_{k_{1}}^{*} \cdots e_{k_{m}}^{*}\right)\right)$. One can show that these monomials are linearly independent as functions on the rank 1 locus, so we see that the equations (6.2.4a) generate the full ideal vanishing on the rank 1 locus.

Unfortunately, this reasoning does not readily extend to the $r=2$ case, though the equations are known in that case via other methods, see Remark 6.4.3.

Remark 6.2.5. Part of the reason that the case $n>2$ is so hard to study is because the notions of border rank and rank behave so differently, and it is not even clear what the precise relation of the two are (does a function of one bound the other for all tensors?).

The general idea is to show that there is a way to take the limit as the parameters $n$ and $\operatorname{dim} V_{i}$ go to $\infty$ and then to study the resulting limit space. Border rank is preserved upon change of basis (with respect to each $V_{i}$ ) so there is a symmetry that can be used in this problem. The existence of the constant can be then deduced from an equivariant noetherianity property of the limit space.
6.3. Hillar-Sullivant theorem. We'll make use of a theorem of Hillar and Sullivant from [HS]. We're going to use it as a tool for proving Theorem 6.2.2, but see their paper for some other applications.

Let $\operatorname{Inc}\left(\mathbf{Z}_{>0}\right)$ be the set of strictly increasing functions $\mathbf{Z}_{>0} \rightarrow \mathbf{Z}_{>0}\left(\mathbf{Z}_{>0}\right.$ is the set of positive integers). It is a monoid under composition. Fix a positive integer $d$ and a commutative ring k. Define

$$
R=\mathbf{k}\left[x_{i, j} \mid 1 \leq i \leq d, j \in \mathbf{Z}_{>0}\right] .
$$

Remark 6.3.1. When $\mathbf{k}=\mathbf{C}$, we have studied this ring in the context of $\mathbf{F I}_{d}$-modules. Namely, $\mathbf{F I}_{d}$-modules are equivalent, via Schur-Weyl duality, to $R$-modules with a compatible action of $\mathbf{G L}_{\infty}(\mathbf{C})$.

We study the action of $\operatorname{Inc}\left(\mathbf{Z}_{>0}\right)$ on $R$ by $f\left(x_{i, j}\right)=x_{i, f(j)}$.

Theorem 6.3.2. Suppose $\mathbf{k}$ is a noetherian ring. Given an increasing chain of $\operatorname{Inc}\left(\mathbf{Z}_{>0}\right)$ invariant ideals $I_{1} \subseteq I_{2} \subseteq \cdots R$, we have $I_{j}=I_{j+1}$ for $j \gg 0$.

Equivalently, any ideal of $R$ that is closed under the action of $\operatorname{Inc}\left(\mathbf{Z}_{>0}\right)$ is finitely generated under multiplication by $R$ and the action of $\operatorname{Inc}\left(\mathbf{Z}_{>0}\right)$.

Proof. This can be proven using the techniques from $\S 3$, so we'll just outline the steps and leave the details to the reader.
(1) We define a term order on the variables by first setting

$$
x_{1,1}<x_{2,1}<\cdots<x_{d, 1}<x_{1,2}<x_{2,2}<\cdots<x_{d, 2}<x_{1,3}<\cdots .
$$

Now order the monomials using lexicographic order. Then given two monomials $m<m^{\prime}$, we have $n m<n m^{\prime}$ for any monomial $n$, and $f m<f m^{\prime}$ for any $f \in \operatorname{Inc}\left(\mathbf{Z}_{>0}\right)$. In particular, we may define an initial ideal for any ideal $I$ closed under $\operatorname{Inc}\left(\mathbf{Z}_{>0}\right)$ and it suffices to prove that it is finitely generated under the simultaneous action of $R$ and $\operatorname{Inc}\left(\mathbf{Z}_{>0}\right)$.
(2) Define a poset $P$ from the monomials in $R$ by $m \preceq m^{\prime}$ if the ideal generated by $m$ contains $m^{\prime}$, i.e., there exists a monomial $n$ and $f \in \operatorname{Inc}\left(\mathbf{Z}_{>0}\right)$ such that $m^{\prime}=n f\left(m^{\prime}\right)$. We can give a model for $P$ as follows. Given a monomial $m$, let $p$ be the largest index such that $x_{i, p}$ has nonzero exponent for some $i$, and define $w(m)=\left(w_{1}, \ldots, w_{p}\right)$ where $w_{i} \in \mathbf{Z}_{\geq 0}^{d}$ records the exponents of $x_{1, i}, \ldots, w_{d, i}$. If we partially order $\mathbf{Z}_{\geq 0}^{d}$ by $\left(\alpha_{1}, \ldots, \alpha_{d}\right) \leq\left(\beta_{1}, \ldots, \beta_{d}\right)$ if $\alpha_{i} \leq \beta_{i}$ for all $i$, then the order on $\left(\mathbf{Z}_{\geq 0}^{d}\right)^{\star}$ that we get is exactly the Higman order, so it is a noetherian poset by Theorem 3.5.4.
(3) Combining the above two steps, we get the theorem for the case when $\mathbf{k}$ is a field. For the general case, we can use the argument in the proof of Theorem 3.6.4.
6.4. Proof of Theorem $\mathbf{6 . 2 . 2}$. The outline of the proof is as follows:
(1) Reduce to the case that $\operatorname{dim} V_{i}=r+1$ for all $i$.
(2) Construct an infinite limit space $V^{\otimes \infty}$ together with an action of a group $G_{\infty}$. Also construct a limit of the border rank $\leq r$ elements $X^{\leq r}$.
(3) Show that $G$-equivariant closed subsets of an auxiliary space $Y^{\leq r}$, which contains $X^{\leq r}$, satisfy the descending chain condition.
(4) Translate the above properties into the desired theorem.
6.4.1. Reduction to bounded number of variables. First, we need the notion of a flattening: given a subset $S \subseteq[n]$, set $U=\bigotimes_{i \in S} V_{i}$ and $U^{\prime}=\bigotimes_{i \notin S} V_{i}$. Then $\mathbf{V}=U \otimes U^{\prime}$, but now we can identify it with the space of matrices. Given $\omega \in \mathbf{V}$, the corresponding matrix is called a flattening of $\omega$ (it depends on $S$ ). If $\omega$ has rank $r$, then its flattening is a sum of $r$ rank 1 matrices and hence has rank $\leq r$.

Lemma 6.4.1. It suffices to prove Theorem 6.2.2 when $\operatorname{dim} V_{i}=r+1$.
Proof. We claim that for any $\omega \in \mathbf{V}$, its rank is at most $r$ if and only if for each set of linear maps $\varphi_{i}: V_{i} \rightarrow \mathbf{k}^{r+1}$, the rank of its image is also at most $r$. The "if" direction is clear: if $\omega$ has rank $\leq r$, then so are all of its images.

To prove the other direction, suppose that $\omega$ has rank $\geq r+1$. We show that for each $i$, there exists a linear map $\varphi: V_{i} \rightarrow \mathbf{k}^{r+1}$ such that the image of $\omega$ has rank $\geq r+1$. Without loss of generality, we may assume $i=1$. We can reinterpret $\omega$ as a linear map

$$
V_{2}^{*} \otimes \cdots \otimes V_{n}^{*} \rightarrow V_{1}
$$

and we let $W$ be the image. If $\operatorname{dim} W \leq r+1$, let $\mathbf{k}^{r+1}$ be any subspace of $V_{1}$ containing $W$ and let $\varphi: V_{1} \rightarrow \mathbf{k}^{r+1}$ be a projection. Then the image of $\omega$ has the same rank as $\omega$. Otherwise, if $\operatorname{dim} W>r+1$, let $\varphi: V \rightarrow \mathbf{k}^{r+1}$ be any linear map such that $\varphi$ maps $W$ surjectively onto $\mathbf{k}^{r+1}$. The flattening of the image of $\omega$ is then $V_{2}^{*} \otimes \cdots \otimes V_{n}^{*} \rightarrow \mathbf{k}^{r+1}$ and has rank $r+1$, so the rank of the image of $\omega$ is $\geq r+1$.

Given the claim, if there is a constant that works for Theorem 6.2.2, then it will work in general: for any tensor product: we can pullback the equations that cut out the rank $\leq r$ locus in $\left(\mathbf{k}^{r+1}\right)^{\otimes n}$ along all possible choices of linear maps to get equations for the rank $\leq r$ locus in $V_{1} \otimes \cdots \otimes V_{n}$.

Exercise 6.4.2. Consider the rank 1 case of flattenings. Show that the equations (6.2.4a) are in the linear span of the $2 \times 2$ minors of the flattenings. Deduce that the flattening equations generate the full ideal of the locus of rank 1 tensors.

Remark 6.4.3. In fact for the rank 2 case, the $3 \times 3$ minors of the flattenings generate the full ideal of the locus of border rank $\leq 2$ tensors, at least when the field is of characteristic 0 . This is proven in [Ra].
6.4.2. Spaces of infinite tensors. In particular, we may as well fix a single space $V$ of dimension $r+1$ to replace all $V_{i}$. Let $X_{p}^{\leq r} \subset V^{\otimes p}$ be the variety of border rank $\leq r$ tensors, and let $Y_{p}^{\leq r} \subset V^{\otimes p}$ be the variety of tensors with flattening rank $\leq r$ under all possible flattenings. Then $X_{p}^{\leq r} \subseteq Y_{p}^{\leq r}$.

We also fix a linear functional $x_{0} \in V^{*}$ and use this to define maps

$$
\begin{aligned}
V^{\otimes(p+1)} & \rightarrow V^{\otimes p} \\
v_{1} \otimes \cdots \otimes v_{p} \otimes v_{p+1} & \mapsto x_{0}\left(v_{p+1}\right) v_{1} \otimes \cdots \otimes v_{p}
\end{aligned}
$$

If $\omega \in V^{\otimes(p+1)}$ has rank $\leq r$, then so does its image in $V^{\otimes p}$. A similar statement holds for flattening rank. Now set

$$
\begin{aligned}
& X^{\leq r}:={\underset{p}{\lim }}_{{\underset{p}{x}}^{\leq}}^{X_{p}^{\leq r}} \\
& Y^{\leq r}:=\underset{{\underset{p}{p}}^{\lim }}{{ }_{p}} Y_{p}^{\leq r} .
\end{aligned}
$$

Call $Y^{\leq r}$ the flattening variety. If we let $T_{p}$ denote the ring of polynomial functions on $V^{\otimes p}$, then the maps $V^{\otimes(p+1)} \rightarrow V^{\otimes p}$ induce maps $T_{p} \rightarrow T_{p+1}$ via pullback of functions, and we set

$$
T_{\infty}=\bigcup_{p} T_{p}
$$

We will think of this as the ring of polynomial functions on $V^{\otimes \infty}$, or its coordinate ring. We can describe this more concretely: $T_{p}$ is the symmetric algebra on $\left(V^{*}\right)^{\otimes p}$, so pick a basis $x_{0}, \ldots, x_{r}$ for $V^{*}$ (same $x_{0}$ as above). Then this naturally gives a basis $x_{I}=x_{i_{1}} \otimes \cdots \otimes x_{i_{p}}$ for $\left(V^{*}\right)^{\otimes p}$ ranging over all $I=\left(i_{1}, \ldots, i_{p}\right) \in\{0, \ldots, r\}^{p}$. We identify this with an infinite word in $\{0, \ldots, r\}$ by appending infinitely many 0 's at the end. The support of this infinite word is the set of indices with nonzero value. Then $T_{\infty}$ is the polynomial ring in the $x_{I}$ as $I$ ranges over all infinite words with finite support.

The space $V^{\otimes p}$ carries an action of $\Sigma_{p}$ by permuting tensor factors:

$$
\sigma \cdot \sum v_{1} \otimes \cdots \otimes v_{p}=\sum v_{\sigma^{-1}(1)} \otimes \cdots \otimes v_{\sigma^{-1}(p)}
$$

It also has an action of $\mathbf{G L}(V)^{p}$ by linear change of coordinates:

$$
\left(g_{1}, \ldots, g_{p}\right) \cdot \sum v_{1} \otimes \cdots \otimes v_{p}=\sum\left(g_{1} v_{1}\right) \otimes \cdots \otimes\left(g_{p} v_{p}\right)
$$

So the semidirect product

$$
G_{p}:=\Sigma_{p} \ltimes \mathbf{G} \mathbf{L}(V)^{p}
$$

acts on $V^{\otimes p}$. We have inclusions $G_{p} \subset G_{p+1}$ compatible with the projection with respect to $x_{0}$, so the union

$$
G_{\infty}=\bigcup_{p} G_{p}
$$

acts on $V^{\otimes \infty}$. Furthermore, $X^{\leq r}$ is closed under the action of $G_{\infty}$.
Both $Y^{\leq r}$ and $X^{\leq r}$ are closed under the action of $G_{\infty}$. Furthermore, $Y^{\leq r}$ is (set-theoretically) defined by taking the determinants of all $(r+1) \times(r+1)$ submatrices of all flattenings, which are in particular polynomials of degree $r+1$.

We will need a more precise version of this statement. Let $\mathbf{w}=\left(w_{1}, \ldots, w_{\ell}\right)$ be an $\ell$-tuple of distinct infinite words in $\{0, \ldots, p\}$ with finite support, and let $\mathbf{w}^{\prime}=\left(w_{1}^{\prime}, \ldots, w_{m}^{\prime}\right)$ be another one. Assume that the support of $w_{i}$ is disjoint from the support of $w_{j}^{\prime}$ for all $i, j$. Let $M\left[\mathbf{w} ; \mathbf{w}^{\prime}\right]$ be the $\ell \times m$ matrix whose $(i, j)$-entry is $x_{w_{i}+w_{j}^{\prime}}$. This is a submatrix of some flattening along subsets $I, J$ where $I$ contains the supports of all $w_{i}$ and $J$ contains the supports of all $w_{j}^{\prime}$.

Remark 6.4.4. In Remark 6.2.4, we saw that the ideal of the locus of rank 1 tensors is generated by polynomials of degree 2. Suppose we're dealing with the tensor product $\mathbf{V}=V_{1} \otimes \cdots \otimes V_{n}$. By iterating Cauchy's identity (Corollary 2.12.3), we see that, in characteristic 0 , the space of quadratic polynomials on $\mathbf{V}$ has a $\mathbf{G L}\left(V_{1}\right) \times \cdots \times \mathbf{G L}\left(V_{n}\right)$ equivariant decomposition

$$
\operatorname{Sym}^{2}\left(V_{1}^{*} \otimes \cdots \otimes V_{n}^{*}\right)=\bigoplus F_{1}\left(V_{1}^{*}\right) \otimes \cdots \otimes F_{n}\left(V_{n}^{*}\right)
$$

where each $F_{i}$ is either $\operatorname{Sym}^{2}$ or $\Lambda^{2}$, and the sum is over all possible ways for an even number of the $F_{i}$ to be $\Lambda^{2}$. The equations themselves span the subspace where a positive (and even) number of the $F_{i}$ are $\Lambda^{2}$.

Hence, as $n$ increases, the number of irreducible representations goes to infinity. Note that the statement above is not saying anything about all equations forming finitely many $G_{\infty}$-orbits, just that one can do this up to taking radical. In fact, one can show that it is not possible to realize all of the equations using just finitely many $G_{\infty}$-orbits, see [Dr1, §7] and also for a discussion of a possible fix.

Theorem 6.4.5. For each $r$, there is a finite list of pairs $\left(\mathbf{w}, \mathbf{w}^{\prime}\right)$ of pairs of words as above of length $r+1$ such that the flattening variety $Y^{\leq r}$ is set-theoretically defined by the $G_{\infty}$-orbits of $\operatorname{det} M\left[w ; w^{\prime}\right]$.

We omit the proof since it is involved and uses some algebraic geometry which we will not discuss. See [DrK, §4] for the details.
6.4.3. Equivariant noetherianity of $Y^{\leq r}$. A topological space $X$ is noetherian if for every descending chain of closed subsets $Y_{1} \supseteq Y_{2} \supseteq \cdots$, we have $Y_{i}=Y_{i+1}$ for $i \gg 0$.

Example 6.4.6. Noetherianity is not a property of "typical" topological spaces, for example, take $Y_{i}=\{n \in \mathbf{Z} \mid n \geq i\}$, considered as subsets of $\mathbf{R}$ with the Euclidean topology.

Noetherianity is ubiquitous in algebraic geometry. For example, $\mathbf{k}^{n}$ is noetherian under the Zariski topology since a descending chain of closed subsets is of the form

$$
V\left(I_{1}\right) \supseteq V\left(I_{2}\right) \supseteq \cdots
$$

for ideals $I_{1} \subseteq I_{2} \subseteq \cdots \subseteq \mathbf{k}\left[x_{1}, \ldots, x_{n}\right]$. Hilbert's basis theorem implies that every ideal is finitely generated. In particular, $I=\bigcup_{j} I_{j}$ is finitely generated, and these generators must live in some $I_{j}$, which implies that $I_{j}=I_{j+k}$ for any $k \geq 0$. Note, however, that noetherianity of the space $\mathbf{k}\left[x_{1}, \ldots, x_{n}\right]$ is strictly weaker than the statement of Hilbert's basis theorem, since (when $\mathbf{k}$ is algebraically closed) it only implies that $\sqrt{I_{j}}=\sqrt{I_{j+1}}$ for $j \gg 0$.

If we have a monoid $G$ acting continuously on $X$, then we say that $X$ is $G$-noetherian if for every descending chain of $G$-invariant closed subsets $Y_{1} \supseteq Y_{2} \supseteq \cdots$, we have $Y_{i}=Y_{i+1}$ for $i \gg 0$.

Exercise 6.4.7. Let $G$ be a monoid acting on a topological space $X$. Prove the following statements:
(1) If $X$ is $G$-noetherian, and $Y \subset X$ is a $G$-invariant subspace with the induced topology, then $Y$ is also $G$-noetherian.
(2) If $Y, Y^{\prime} \subset X$ are $G$-invariant subspaces and are both $G$-noetherian, then so is $Y \cup Y^{\prime}$.
(3) If $X^{\prime}$ is another space with a $G$-action and $f: X \rightarrow X^{\prime}$ is $G$-equivariant and continuous, then the image of $f$ is $G$-noetherian.
(4) Suppose $G$ is a group and let $G^{\prime} \subseteq G$ be a subgroup. Let $G^{\prime}$ act on $G \times X$ by $h \cdot(g, x)=\left(g h^{-1}, h x\right)$ for $h \in G^{\prime}$. If $X$ is $G^{\prime}$-noetherian, then the space $(G \times X) / G^{\prime}$ is $G$-noetherian.

Our goal is the following:
Theorem 6.4.8 (Draisma-Kuttler). The space $Y^{\leq r}$ is $G_{\infty}$-noetherian.
Proof. This is proven by induction on $r$. The case $r=0$ is trivial since $Y^{\leq 0}$ is a single point.
Now assume that $r>0$ and that $Y^{\leq(r-1)}$ is $G_{\infty}$-noetherian. By Theorem 6.4.5, we can find finitely many pairs $\mathbf{w}_{1}, \mathbf{w}_{1}^{\prime}, \ldots, \mathbf{w}_{N}, \mathbf{w}_{N}^{\prime}$ such that
(1) Each $\mathbf{w}_{i}$ has length $r$. If we write it as $\left(v_{1}, \ldots, v_{r}\right)$, then each $v_{i}$ has finite support.
(2) For each $i$, the components of $\mathbf{w}_{i}$ and $\mathbf{w}_{i}^{\prime}$ have disjoint support.
(3) $Y^{\leq(r-1)}$ is set-theoretically defined by the $G_{\infty}$-orbits of $\operatorname{det}\left(M\left[\mathbf{w}_{i} ; \mathbf{w}_{i}^{\prime}\right]\right)$.

For each $i$, define

$$
Z_{i}=Y^{\leq r} \backslash V\left(G_{\infty} \operatorname{det} M\left[\mathbf{w}_{i} ; \mathbf{w}_{i}^{\prime}\right]\right) .
$$

Then $Z_{i}$ is an open subset of $Y^{\leq r}$ which is closed under $G_{\infty}$, and we have

$$
Y^{\leq r}=Y^{\leq(r-1)} \cup Z_{1} \cup Z_{2} \cup \cdots \cup Z_{N} .
$$

By Exercise 6.4.7(2), it will suffice to show that each $Z_{i}$ is $G_{\infty}$-noetherian. So fix one such $i$ and set $Z=Z_{i}, \mathbf{w}_{i}=\left(w_{1}, \ldots, w_{r}\right), \mathbf{w}_{i}^{\prime}=\left(w_{1}^{\prime}, \ldots, w_{r}^{\prime}\right)$. Let $p$ be the maximum of the union of the supports of all $w_{i}$ and $w_{i}^{\prime}$.

Recall that we have a copy of $\Sigma_{\infty} \subset G_{\infty}$ which permutes tensor factors. Let $\Sigma \subset \Sigma_{\infty}$ be the subgroup of $\sigma$ such that $\sigma(i)=i$ for $i=1, \ldots, p$. Define

$$
Z^{\prime}=Y^{\leq r} \backslash V\left(\operatorname{det} M\left[\mathbf{w} ; \mathbf{w}^{\prime}\right]\right)
$$

Then $Z^{\prime}$ is invariant under $\Sigma$. By the next lemma, $Z^{\prime}$ is $\Sigma$-noetherian. Now define a map

$$
\begin{aligned}
\left(G_{\infty} \times Z^{\prime}\right) / \Sigma & \rightarrow Z \\
(g, z) & \mapsto g \cdot z
\end{aligned}
$$

This is both $G_{\infty}$-equivariant and surjective, so by Exercise 6.4.7, $Z$ is $G_{\infty}$-noetherian, and we are done.

Lemma 6.4.9. $Z^{\prime}$ is $\Sigma$-noetherian.
Proof. We continue to use the same notation from the previous proof.
Let $\Gamma \subset \operatorname{Inc}\left(\mathbf{Z}_{>0}\right)$ be the submonoid of increasing functions which are the identity on $1, \ldots, p$. Note that $\Gamma$ acts on $Y^{\leq r}$ : given an increasing function $f$ and a vector $v \in Y^{\leq r}$, there are only finitely many basis vectors $e_{I}$ needed to write $v$, and hence the union of the $I$ have finite support, let $q$ be the maximal value. Then for any two permutations $\sigma, \sigma^{\prime}$ whose inverses agree on $1, \ldots, q$, we have $\sigma(v)=\sigma^{\prime}(v)$. In particular, pick any $\sigma \in \Sigma$ whose inverse agrees with $f$ on $1, \ldots, q$ and define $f(v)=\sigma(v)$. It will suffice to show that $Z^{\prime}$ is $\Gamma$-noetherian, since the $\Gamma$-orbit of an equation is contained in its $\Sigma$-orbit.

Let $J$ be the set of infinite words $w$ in $\{0, \ldots, r\}$ such that the support of $w$ contains at most 1 element in positions beyond $p$. Let $\mathbf{k}^{J}$ be the vector space with basis indexed by $J$. Setting $R=\mathbf{k}\left[x_{w} \mid w \in J\right], \mathbf{k}^{J}$ is a subset of $\operatorname{Spec}(R)$. By setting the noetherian ring in Theorem 6.3.2 to be the polynomial ring over $\mathbf{k}$ in the variables $x_{w}$ where $w$ ranges over the (finitely many) words with support contained in $[p]$, we deduce that $R$ is $\Gamma$-noetherian, and hence the same is true for $\mathbf{k}^{J}$.

Then $\operatorname{det} M\left[\mathbf{w} ; \mathbf{w}^{\prime}\right]$ can be interpreted as a function on $\mathbf{k}^{J}$; let $Q$ be the open subset where this function is nonzero. We define a function $\pi: Z^{\prime} \rightarrow Q$ which sends a point to the values of its coordinates $x_{w}$ where $w \in J$. This is continuous and $\Gamma$-equivariant. In fact, it is an embedding. We omit the details and refer the reader to [Dr2, §7], but the point is that on the subspace $Z^{\prime}$, one can rewrite any coordinate $x_{w^{\prime}}$ in terms of the coordinates $x_{w}$ with $w \in J$. Assuming this embedding property, we see that $Z^{\prime}$ inherits $\Gamma$-noetherianity from $Q$, so we are done.

Remark 6.4.10. By Exercise 6.4.2, $Y^{\leq 1}=X^{\leq 1}$. So $Y^{\leq 1}$ consists of two $G_{\infty}$-orbits: the points $\{0\}$ and the set of nonzero rank 1 tensors. So the theorem is also trivial in this case.

Lemma 6.4.11. If $Z \subset Y^{\leq r}$ is a closed $G_{\infty}$-subset, then it is the common solution set of finitely many $G_{\infty}$-orbits of polynomials. In particular, it is defined by bounded degree polynomials.

Proof. If not, we can find an infinite sequence of polynomials $f_{1}, f_{2}, \ldots$ that all vanish on $Z$, but such that $Z_{i}$, the common solution set of $G_{\infty} f_{1}, \ldots, G_{\infty} f_{i}$, satisfies $Z_{1} \supsetneqq Z_{2} \supsetneqq \cdots$ in direct contradiction to (2).

Taking $Z=X^{\leq r}$, we get that it is defined by bounded degree polynomials.
6.4.4. Finishing the proof. To finish, pick $e_{0} \in V$ such that $x_{0}\left(e_{0}\right)=1$. Define

$$
\begin{aligned}
\tau: V^{\otimes p} & \rightarrow V^{\otimes(p+1)} \\
\omega & \mapsto \omega \otimes e_{0} .
\end{aligned}
$$

Given $\omega \in V^{\otimes p}$, we can define $\omega_{\infty}=\omega \otimes e_{0}^{\otimes \infty} \in V^{\otimes \infty}$, and their ranks are the same. So pulling back all of the equations vanishing on $X^{\leq r}$ gives equations that define the border rank $\leq r$ locus in $V^{\otimes p}$.
6.5. Variants. A similar situation can be considered with both exterior and symmetric powers of vector spaces instead of tensor powers. Given $V$, an element of $\bigwedge^{n} V$ of the form $v_{1} \wedge \cdots \wedge v_{n}$ is said to have rank 1. Rank and border rank are defined similarly. For $\mathrm{Sym}^{n} V$, elements of the form $v^{n}$ are said to have rank 1 (technically we should be working with the divided power instead of the symmetric power). We note some theorems:

Theorem 6.5.1 (Draisma-Eggermont [DE]). For each $r$, there is a constant $C(r)$ such that the variety of border rank $\leq r$ anti-symmetric tensors is cut out by polynomials of degree $\leq C(r)$. The constant is independent of $n$ and also $V$.

Theorem 6.5.2 (Sam [Sa]). Fix a field of characteristic 0. For each r, there is a constant $C(r)$ such that the full ideal of polynomials vanishing on the variety of border rank $\leq r$ symmetric tensors is generated by polynomials of degree $\leq C(r)$. The constant is independent of $n$ and also $V$.

The last theorem gives an ideal-theoretic statement which improves the set-theoretic statement (though is limited by its restriction on characteristic).

## 7. More on FI-modules

7.1. Asymptotic combinatorial properties. In a homework problem, you were asked to show that given a finitely generated FI-module $M$ over a field $\mathbf{k}$, the function $n \mapsto \operatorname{dim}_{\mathbf{k}} M_{n}$ is a polynomial for $n \gg 0$. Each $M_{n}$ is a representation of the symmetric group $\Sigma_{n}$. If $\mathbf{k}$ is a field of characteristic 0 , then $M_{n}$ decomposes into a direct sum of irreducible representations, which are indexed by partitions $\lambda$ of $n$, say that $\mathbf{M}_{\lambda}$ appears with multiplicity $m_{M, \lambda}$. We'd like to understand how these multiplicities behave.

Recall that the principal projective module $P_{n}$ satisfies $P_{n}(m)=\operatorname{Ind}_{\Sigma_{n} \times \Sigma_{m-n}}^{\Sigma_{m}} \mathbf{k}\left[\Sigma_{n}\right]$ where $\Sigma_{m-n}$ acts trivially on $\mathbf{k}\left[\Sigma_{n}\right]$. We have an action of $\Sigma_{n}$ on $P_{n}$ by pre-composition, which commutes with post-composition by injections (this is just the associativity of composing functions). In that sense, we can decompose $P_{n}$ into a direct sum

$$
P_{n}=\bigoplus_{\lambda \vdash n} P_{\lambda}^{\oplus \operatorname{dim} \mathbf{M}_{\lambda}}
$$

which follows from the decomposition of the group algebra

$$
\mathbf{k}\left[\Sigma_{n}\right]=\bigoplus_{\lambda \vdash n} \mathbf{M}_{\lambda}^{\oplus \operatorname{dim} \mathbf{M}_{\lambda}} .
$$

In particular, we have

$$
P_{\lambda}(m)=\operatorname{Ind}_{\Sigma_{n} \times \Sigma_{m-n}}^{\Sigma_{m}} \mathbf{M}_{\lambda}=\bigoplus_{\nu \vdash m} \mathbf{M}_{\nu}
$$

where, by Pieri's rule (Theorem 2.7.3), the last sum is over all $\nu$ containing $\lambda$ such that $\nu / \lambda$ is a horizontal strip of size $m-n$. The submodule structure of $P_{\lambda}$ can be determined completely from the following result whose proof we will omit.

Proposition 7.1.1. Given $\nu$ such that $\nu / \lambda$ is a horizontal strip, the submodule of $P_{\lambda}$ generated by $\mathbf{M}_{\nu}$ contains all $\mathbf{M}_{\nu^{\prime}}$ such that $\nu \subseteq \nu^{\prime}$.

Recall that $P_{n}$ satisfies the universal property that a map of FI-modules $P_{n} \rightarrow M$ is the same as an $\Sigma_{n}$-equivariant map $\mathbf{k}\left[\Sigma_{n}\right] \rightarrow M(n)$, which is again equivalent to a choice of arbitrary element in $M(n)$. Similarly, a map of FI-modules $P_{\lambda} \rightarrow M$ is the same as an $\Sigma_{n}$-equivariant map $\mathbf{M}_{\lambda} \rightarrow M(n)$.

To add a horizontal strip to $\lambda$, we add some subset of boxes to the first $\lambda_{1}$ columns, and the rest must go in the first row. Let $\mu^{1}, \ldots, \mu^{c}$ be all possible partitions that can be obtained by adding a horizontal strip to $\lambda$ by only adding boxes in the first $\lambda_{1}$ columns. Then we see that every partition in the decomposition of $P_{\lambda}(m)$ must be of the form $\mu^{i}+\left(m-\left|\mu^{i}\right|\right)$ for some $i$ and where addition for partitions is defined componentwise.

Example 7.1.2. If $\lambda=(3,1)$, then $c=6$ with $\mu^{1}=\lambda=(3,1), \mu^{2}=(3,1,1), \mu^{3}=(3,2)$, $\mu^{4}=(3,2,1), \mu^{5}=(3,3)$, and $\mu^{6}=(3,3,1)$. If we add a horizontal strip to $\lambda$ to get a partition of size $m$, then it must be one of
$(m-1,1), \quad(m-2,1,1), \quad(m-2,2), \quad(m-3,2,1), \quad(m-3,3), \quad(m-4,3,1)$.

Let $\mu$ be the largest of the $\mu^{i}$, i.e., it is obtained by adding all possible boxes in the first $\lambda_{1}$ columns to $\lambda$ (in the previous example, it is $(3,2,1)$ ). Then we have a map $P_{\mu} \rightarrow P_{\lambda}$ by the above discussion by taking the identity map $\mathbf{M}_{\mu} \rightarrow \mathbf{M}_{\mu} \subset P_{\lambda}(|\mu|-|\lambda|)$. If $\mathbf{M}_{\nu}$ is in the image of this map, then both $\nu / \lambda$ and $\nu / \mu$ must be horizontal strips, and this is only possible for partitions of the form $\mu+(k)$. In fact, by Proposition 7.1.1, $\mathbf{M}_{\mu+(k)}$ is in the image for all $k \geq 0$. In any case, one can argue that the sum of the subspaces $\mathbf{M}_{\mu+(k)}$ is an FI-submodule of $P_{\lambda}$, which we will temporarily denote by $V_{\mu}^{+}$.

Proposition 7.1.3. $V_{\mu}^{+}$is a finitely generated FI-module.
Proof. $V_{\mu}^{+}$is a submodule of $P_{\lambda}$, which is finitely generated, so the result follows from noetherianity of FI-modules.

Proposition 7.1.4. Let $M$ be a finitely generated FI-module over a field of characteristic 0 . For a partition $\lambda$, let $m_{M, \lambda}$ be the multiplicity of $\mathbf{M}_{\lambda}$ in $M(|\lambda|)$. Then, for any $\lambda$, the function $i \mapsto m_{M, \lambda+(i)}$ is constant for $i \gg 0$.

Proof. By Exercise 5.1.2, the Segre product of two finitely generated FI-modules is again finitely generated. In particular, we have a finitely generated FI-module of the form $[n] \mapsto$ $M(n) \otimes_{\mathbf{k}} V_{\lambda}^{+}(n)=M(n) \otimes_{\mathbf{k}} \mathbf{M}_{\lambda+(n-|\lambda|)}$. Note that

$$
\operatorname{dim}_{\mathbf{k}}\left(M(n) \otimes_{\mathbf{k}} \mathbf{M}_{\lambda+(n-|\lambda|)}\right)^{\Sigma_{n}}=\operatorname{dim}_{\mathbf{k}}\left(M(n) \otimes_{\mathbf{k}} \mathbf{M}_{\lambda+(n-|\lambda|)}\right)_{\Sigma_{n}}=m_{M, \lambda+(n-|\lambda|)}
$$

where the first equality follows from semisimplicity, and the second follows from self-duality of representations of $\Sigma_{n}$. In particular, Proposition 5.2 .5 implies that there is a finitely generated graded $\mathbf{k}[t]$-module whose degree $i$ piece has dimension $m_{M, \lambda+(i)}$, so this value is constant for $i \gg 0$.

Recall that in $\S 2.12$, we defined the Kronecker coefficients by

$$
g_{\lambda, \mu, \nu}=\operatorname{dim}_{\mathbf{C}}\left(\mathbf{M}_{\lambda} \otimes \mathbf{M}_{\mu} \otimes \mathbf{M}_{\nu}\right)^{\Sigma_{n}}
$$

where $\lambda, \mu, \nu$ are partitions of $n$. Alternatively, since $\Sigma_{n}$-representations are self-dual, this is the multiplicity of $\mathbf{M}_{\lambda}$ in $\mathbf{M}_{\mu} \otimes \mathbf{M}_{\nu}$.
Corollary 7.1.5 (Murnaghan). Let $\lambda, \mu, \nu$ be partitions of the same size. Then the function

$$
i \mapsto g_{\lambda+(i), \mu+(i), \nu+(i)}
$$

is constant for $i \gg 0$.
Proof. In the above proposition, take $M$ to be the Segre product of $V_{\mu}^{+}$and $V_{\nu}^{+}$.
The value

$$
\bar{g}_{\lambda, \mu, \nu}=\lim _{i \rightarrow \infty} g_{\lambda+(i), \mu+(i), \nu+(i)}
$$

is the stable Kronecker coefficient, and are in some sense, supposed to simpler than the usual Kronecker coefficients. Our indexing is different from the standard indexing in the literature: there one often replaces $\lambda, \mu, \nu$ by the partitions obtained by removing the first part.

Example 7.1.6. - When $\lambda=\mu=\nu=\varnothing$, then $g_{i, i, i}=1$ for all $i$ since $\mathbf{M}_{i}$ is the trivial representation of $\Sigma_{i}$.

- When $\lambda=\mu=\nu=(1,1)$, then $g_{(i+1,1),(i+1,1),(i+1,1)}$ is 1 if $i>0$ and is 0 for $i=0$. The case $i=0$ follows from the fact that $M_{(1,1)}$ is the sign representation of $\Sigma_{2}$. In general, for partitions of size $n, g_{\alpha, \beta,(n-1,1)}+1$ is the number of ways to remove a box from $\alpha$ and then add a box to get $\beta$.

More generally, one can add other sequences to $\lambda, \mu, \nu$. For example, fix 3 other partitions $\alpha, \beta, \gamma$ of the same size and consider the sequence

$$
i \mapsto g_{\lambda+i \alpha, \mu+i \beta, \nu+i \gamma}
$$

What kind of behavior does this exhibit? In fact, one has the following result (see [Ste, SS2]):
Theorem 7.1.7 (Stembridge, Sam-Snowden). If $g_{i \alpha, i \beta, i \gamma}=1$ for all $i$, then the sequence $i \mapsto g_{\lambda+i \alpha, \mu+i \beta, \nu+i \gamma}$ is constant for $i \gg 0$. Conversely, if this sequence is eventually constant for all $\lambda, \mu, \nu$, then $g_{i \alpha, i \beta, i \gamma}=1$ for all $i$.

Murnaghan's theorem is the special case when $\alpha=\beta=\gamma=(1)$ (cf. Example 7.1.6).
The techniques used here will not suffice to prove this generalization, which uses a little bit of algebraic geometry, see [SS2] for more details.
7.2. Serre quotient categories. To understand the role of the $V_{\lambda}^{+}$in the asymptotic behavior of finitely generated FI-modules, we use the notion of a Serre quotient of an abelian category.

Definition 7.2.1. Let $\mathcal{C}$ be an abelian category. A Serre subcategory $\mathcal{C}^{\prime}$ of $\mathcal{C}$ is a full subcategory (i.e., $\operatorname{Hom}_{\mathcal{P}^{\prime}}(x, y)=\operatorname{Hom}_{\mathfrak{C}}(x, y)$ for all $x, y \in \mathcal{C}^{\prime}$ ) such that given a short exact sequence of objects in $\mathcal{C}$

$$
0 \rightarrow A \rightarrow B \rightarrow C \rightarrow 0
$$

we have $B \in \mathcal{C}^{\prime}$ if and only if $A, C \in \mathcal{C}^{\prime}$.

This comes up as follows. Suppose we are given an exact functor $F: \mathcal{C} \rightarrow \mathcal{D}$ between abelian categories. Then the full subcategory on the objects $\{x \in \mathcal{C} \mid F(x)=0\}$ is a Serre subcategory. In this case, we can think of this subcategory as being the kernel of the functor $F$. Then, mimicking more familiar constructions in algebra, we will also want to define the quotient by a Serre subcategory.

Definition 7.2.2. Let $\mathcal{C}^{\prime} \subset \mathcal{C}$ be a Serre subcategory and let $x, y \in \mathcal{C}$ be objects. Consider the category $\mathcal{D}(x, y)$ whose objects are pairs $\left(x^{\prime}, y^{\prime}\right)$ where $x^{\prime} \subseteq x$ is a subobject such that $x / x^{\prime} \in \mathcal{C}^{\prime}$ and $y^{\prime} \subseteq y$ is a subobject such that $y^{\prime} \in \mathcal{C}^{\prime}$. We have a morphism $\left(x^{\prime}, y^{\prime}\right) \rightarrow\left(x^{\prime \prime}, y^{\prime \prime}\right)$ if $x^{\prime \prime} \subseteq x^{\prime}$ and $y^{\prime} \subseteq y^{\prime \prime}$. The assignment $\left(x^{\prime}, y^{\prime}\right) \mapsto \operatorname{Hom}_{\mathcal{P}^{\prime}}\left(x^{\prime}, y / y^{\prime}\right)$ is a functor on the category $\mathcal{D}(x, y)$.

The quotient category $\mathcal{C} / \mathcal{C}^{\prime}$ has the same objects as $\mathcal{C}$, and the hom set between $x$ and $y$ is defined by

$$
\operatorname{Hom}_{\mathcal{C} / \mathbb{E}^{\prime}}(x, y)=\underset{\left(x^{\prime}, y^{\prime}\right) \in \mathcal{D}(x, y)}{\lim } \operatorname{Hom}_{\mathbb{C}^{\prime}}\left(x^{\prime}, y / y^{\prime}\right)
$$

Recall that the colimit can be defined as the direct sum $\bigoplus_{\left(x^{\prime}, y^{\prime}\right) \in \mathcal{D}(x, y)} \operatorname{Hom}_{\mathcal{P}^{\prime}}\left(x^{\prime}, y / y^{\prime}\right)$ modulo the equivalence relation that identifies a morphism $x^{\prime} \rightarrow y / y^{\prime}$ with the corresponding morphism $x^{\prime \prime} \rightarrow y / y^{\prime \prime}$ if we have a morphism $\left(x^{\prime}, y^{\prime}\right) \rightarrow\left(x^{\prime \prime}, y^{\prime \prime}\right)$ in $\mathcal{D}(x, y)$.

To compose maps $x \rightarrow y$ and $y \rightarrow z$ in $\mathcal{C} / \mathcal{C}^{\prime}$, first pick representatives $f: x^{\prime} \rightarrow y / y^{\prime}$ and $g: y^{\prime \prime} \rightarrow z / z^{\prime}$. Let $x^{\prime \prime} \subseteq x^{\prime}$ be the preimage of $y^{\prime \prime} / y^{\prime}$ under $f$ and let $z^{\prime \prime}$ be the sum of $z^{\prime}$ and the image of $g$. Then we can restrict $f$ and $g$ to get representatives $f^{\prime}: x^{\prime \prime} \rightarrow y / y^{\prime}$ and $g^{\prime}: y^{\prime \prime} \rightarrow z / z^{\prime \prime}$, which are actually composable since the image of the first map is $y^{\prime \prime}$. The composition $x \rightarrow y \rightarrow z$ is then represented by $g^{\prime} \circ f^{\prime}$.

In more detail, a morphism $x \rightarrow y$ in the quotient category can be represented by a morphism $x^{\prime} \rightarrow y / y^{\prime}$ where $x / x^{\prime}$ and $y^{\prime}$ are both in $\mathcal{C}$. Given a subobject $y^{\prime \prime} \subseteq y^{\prime}$, we get a natural map $\operatorname{Hom}_{\mathcal{C}}\left(x^{\prime}, y / y^{\prime \prime}\right) \rightarrow \operatorname{Hom}_{\mathcal{C}}\left(x^{\prime}, y / y^{\prime}\right)$, and any morphism in the source represents the same morphism as its image. Similarly, given $x^{\prime} \subseteq x^{\prime \prime} \subseteq x, x / x^{\prime \prime}$ is a quotient of $x / x^{\prime}$ and hence belongs to $\mathcal{C}$, and we get a natural map $\operatorname{Hom}_{\mathcal{C}}\left(x^{\prime \prime}, y / y^{\prime}\right) \rightarrow \operatorname{Hom}_{\mathfrak{C}}\left(x^{\prime}, y / y^{\prime}\right)$ and any morphism in the source represents the same morphism as its image.

In particular, $x$ and $y$ are isomorphic in $\mathcal{C} / \mathcal{C}^{\prime}$ if and only if there are maps $x \stackrel{f}{\leftarrow} z \xrightarrow{g} y$ such that the kernel and cokernel of both $f, g$ are in $\mathcal{C}^{\prime}$.

In fact, $\mathcal{C} / \mathcal{C}^{\prime}$ is an abelian category, and the quotient functor $T: \mathcal{C} \rightarrow \mathcal{C} / \mathcal{C}^{\prime}$ is exact, and we call it the localization functor. It satisfies the following universal property: any exact functor $F: \mathcal{C} \rightarrow \mathcal{D}$ such that $F(x)=0$ for all $x \in \mathcal{C}^{\prime}$ factors through $T$, i.e., there is a functor $F^{\prime}: \mathcal{C} / \mathcal{C}^{\prime} \rightarrow \mathcal{D}$ such that $F=F^{\prime} \circ T$.
Example 7.2.3. Let $\mathcal{C}$ be the category of finitely generated $R$-modules for a PID $R$. Let $\mathcal{C}^{\prime}$ be the full subcategory on the torsion $R$-modules. Then $\mathcal{C} / \mathcal{C}^{\prime}$ is equivalent to the category of finite-dimensional $K$-vector spaces, where $K$ is the fraction field of $R$.

Example 7.2.4 (For those who know some algebraic geometry). Let $\mathbf{k}$ be a field and let $\mathcal{C}$ be the category of finitely generated graded $\mathbf{k}\left[x_{1}, \ldots, x_{n}\right]$-modules (here $\operatorname{deg} x_{i}=1$ for all $i$ ) and let $\mathfrak{C}^{\prime}$ be the full subcategory on the torsion modules. In this case, $\mathcal{C} / \mathrm{C}^{\prime}$ is equivalent to the category of coherent sheaves on projective space $\mathbf{P}_{\mathbf{k}}^{n-1}$.

Inspired by these examples, an important case to consider is when $\mathcal{C}$ is the category of modules of some kind, and $\mathcal{C}^{\prime}$ is the subcategory of "torsion" modules. Then $\mathcal{C} / \mathcal{C}^{\prime}$ can be thought of as the category of modules over the "fraction field", even if such an object does
not exist. Borrowing terminology from algebraic geometry, where the fraction field of an integral domain corresponds to the "generic point" of the corresponding affine scheme, we call $\mathcal{C} / \mathcal{C}^{\prime}$ the generic category.
7.3. Generic FI-modules. An element $x \in M(n)$ of an FI-module $M$ is torsion if there exists an integer $m$ and a non-invertible morphism $f:[n] \rightarrow[m]$ such that $f_{*}(x)=0$. Since $\Sigma_{m}$ acts transitively on all morphisms $[n] \rightarrow[m]$, this implies that $f_{*}(x)=0$ for all non-invertible morphisms $f:[n] \rightarrow[m]$. In particular, the set of torsion elements of an FImodule is a submodule, called the torsion submodule, and denoted $M^{\text {tors }}$. An FI-module $M$ is torsion if $M=M^{\text {tors }}$.

Lemma 7.3.1. If $M$ is finitely generated and torsion, then $M(m)=0$ for $m \gg 0$.
Proof. Say that $M$ is generated by $x_{1}, \ldots, x_{r}$, where $x_{i} \in M\left(n_{i}\right)$. Since each is torsion, there exists $m_{i}$ such that $f_{*}\left(x_{i}\right)=0$ for all $f:\left[n_{i}\right] \rightarrow[m]$ with $m \geq m_{i}$. In particular, $M(m)=0$ for $m \geq \max \left(m_{1}, \ldots, m_{r}\right)$.

We consider the above situation with $\mathcal{C}$ being the category of finitely generated FImodules, and $\mathcal{C}^{\prime}$ being the full subcategory of torsion modules (it is easy to verify that $\mathcal{C}^{\prime}$ is a Serre subcategory).

Given an integer $n$ and an FI-module $M$, let $M^{\geq n}$ be the FI-submodule of $M$ generated by all elements of degrees $\geq n$. Call it the degree $n$ truncation of $M$. In particular,

$$
M^{\geq n}(m)= \begin{cases}M(m) & \text { if } m \geq n \\ 0 & \text { else }\end{cases}
$$

Proposition 7.3.2. Two finitely generated FI-modules $M$ and $N$ are isomorphic in the generic category $\mathcal{C} / \mathfrak{C}^{\prime}$ if and only if $M^{\geq n} \cong N^{\geq n}$ for some $n$.

Proof. Note that $M / M^{\geq n}$ and $N / N^{\geq n}$ are both objects in $\mathcal{C}^{\prime}$, so if $M^{\geq n} \cong N^{\geq n}$, pick an isomorphism $\varphi$. Then $\varphi$ is an element in the colimit defining $\operatorname{Hom}_{\mathcal{C} / \mathrm{e}^{\prime}}(M, N)$, and its inverse $\varphi^{-1}$ is also an element in the colimit defining $\operatorname{Hom}_{\mathcal{C} / \mathrm{e}^{\prime}}(N, M)$, and their compositions in both directions represent the identity.

Conversely, if $M$ and $N$ are isomorphic in $\mathcal{C} / \mathcal{C}^{\prime}$, then pick an isomorphism $\varphi$ with inverse $\varphi^{-1}$. Then there exists $M^{\prime} \subseteq M$ and $N^{\prime} \subseteq N$ such that $M / M^{\prime}$ and $N^{\prime}$ are torsion, and an FImodule morphism $f: M^{\prime} \rightarrow N / N^{\prime}$ which represents $\varphi$ in the colimit. Similarly, there exists $N^{\prime \prime} \subseteq N$ and $M^{\prime \prime} \subseteq M$ such that $N / N^{\prime \prime}$ and $M^{\prime \prime}$ are torsion and an FI-module morphism $g: N^{\prime \prime} \rightarrow M / M^{\prime \prime}$ which represents $\varphi^{-1}$ in the colimit. Up to the equivalence relation on the colimit, their composition is defined on some pair $\left(M^{\prime \prime \prime}, N^{\prime \prime \prime}\right)$ such that $M / M^{\prime \prime \prime}$ and $N^{\prime \prime \prime}$ are torsion. By Lemma 7.3.1, these are 0 in large enough degrees, say larger than $n$. The restrictions of $f$ and $g$ to the degree $n$ truncations are then inverses of one another.

A particular case of truncations: if $|\lambda|=k$, then $\left(V_{\lambda}^{+}\right)^{\geq n}=V_{\lambda+(n-k)}^{+}$, as follows from Proposition 7.1.1. In particular, these truncations are all isomorphic to each other in the generic category $\mathcal{C} / \mathcal{C}^{\prime}$. In fact, since the size of the first part is irrelevant, let $\mu=\left(\lambda_{2}, \lambda_{3}, \ldots\right)$, and call the resulting object $L_{\mu}$ in $\mathcal{C} / \mathcal{C}^{\prime}$.

In fact, $L_{\mu}$ is a simple object: any nonzero subobject would be represented by a subobject $V_{\lambda+(n+k)}^{+}$in $V_{\lambda}$, and the quotient is torsion, so the only nonzero subobject of $L_{\mu}$ is $L_{\mu}$ itself. We can be more precise about the structure of the generic category (proof omitted):

Theorem 7.3.3. (1) Every simple object of $\mathcal{C} / \mathcal{C}^{\prime}$ is isomorphic to $L_{\mu}$ for some $\mu$, and they are all mutually non-isomorphic.
(2) Every object of $\mathcal{C} / \mathfrak{C}^{\prime \prime}$ has a finite composition series.

This structural result implies that many asymptotic properties about finitely generated FI-modules can be deduced from the corresponding properties of the $V_{\lambda}^{+}$, which involve "growing the first row to infinity".

Remark 7.3.4. Given an abelian category $\mathcal{C}$, say that an object is dimension 0 if every finitely generated subobject has a finite composition series. Let $\mathcal{C}^{\prime}$ be the full subcategory of dimension 0 objects. By induction, say that $M$ has dimension $d$ if $T(M) \in \mathcal{C} / \mathcal{C}^{\prime}$ has dimension $d-1$. For any object where this definition is not applicable, i.e., it does not become dimension 0 after any finite number of iterations of quotienting by dimension 0 objects, then define its dimension to be $\infty$.

The Gabriel-Krull dimension of $\mathcal{C}$ is defined to be the supremum of the dimension of all objects. Our discussion above says that the category of FI-modules has dimension 1. This is some measure of complexity of the category: when $R$ is a noetherian commutative ring, the Gabriel-Krull dimension of its module category agrees with its Krull dimension [ $\star$ Steven: double-check assumptions $\star$ ].
7.4. Semi-induced FI-modules. The projective FI-modules $P_{\lambda}$ in characteristic 0 are "induced modules" in the sense that we have

$$
P_{\lambda}(m) \cong \operatorname{Ind}_{\Sigma_{n} \times \Sigma_{m-n}}^{\Sigma_{m}} \mathbf{M}_{\lambda}
$$

for all $m \geq n$, and where $n=|\lambda|$.
Now consider FI-modules over a general field $\mathbf{k}$. Let $V$ be a (finite-dimensional) representation of $\Sigma_{n}$ for some $n$. If $x_{1}, \ldots, x_{r}$ generate $V$ as a $\Sigma_{n}$-module, then we have a surjection $\mathbf{k}\left[\Sigma_{n}\right]^{\oplus r} \rightarrow V$ via $\left(\alpha_{1}, \ldots, \alpha_{r}\right) \mapsto \sum_{i} \alpha_{i} x_{i}$. Let $K$ be the kernel. Then for each $m \geq n$, we have a short exact sequence

$$
0 \rightarrow \operatorname{Ind}_{\Sigma_{n} \times \Sigma_{m-n}}^{\Sigma_{m}} K \rightarrow \operatorname{Ind}_{\Sigma_{n} \times \Sigma_{m-n}}^{\Sigma_{m}} \mathbf{k}\left[\Sigma_{n}\right]^{\oplus r} \rightarrow \operatorname{Ind}_{\Sigma_{n} \times \Sigma_{m-n}}^{\Sigma_{m}} V \rightarrow 0 .
$$

The middle term is $P_{n}^{\oplus r}(m)$, and since the $\Sigma_{n}$ action on $P_{n}$ commutes with the FI-module structure, we have a submodule $\mathcal{J}(K) \subset P_{n}^{\oplus r}$ given by

$$
\mathcal{J}(K)(m)=\operatorname{Ind}_{\Sigma_{n} \times \Sigma_{m-n}}^{\Sigma_{m}} K
$$

and a quotient module $P_{n}^{\oplus r} \rightarrow \mathcal{J}(V)$ given by

$$
\mathcal{J}(V)(m)=\operatorname{Ind}_{\Sigma_{n} \times \Sigma_{m-n}}^{\Sigma_{m}} V .
$$

Since $V$ was arbitrary, the construction of $\mathcal{J}(V)$ makes sense in general, though it may not be clear that it is independent of the choice of generators. Alternatively, one can use the twisted commutative algebra perspective and define $V$ as $A \otimes V$ where $A$ is the free twisted commutative algebra generated by a single element of degree 1 (see $\S \S 3.1$ and 3.2).

In any case, we call $\mathcal{J}(V)$ an induced FI-module.
Exercise 7.4.1. Show that $\mathcal{J}(V)$ is a projective FI-module if and only if $V$ is a projective $\mathbf{k}\left[\Sigma_{n}\right]$-module (which is automatic if $\mathbf{k}\left[\Sigma_{n}\right]$ is semisimple, for example in characteristic 0 ).
Remark 7.4.2. In the language of relative homological algebra [ Ho ], $\mathcal{J}(V)$ is a relatively projective module with respect to the pair of "rings" $\left(\mathbf{k}[\mathbf{F I}], \mathbf{k}\left[\Sigma_{n}\right]\right)$. To define this in general, consider a ring $R$ with a subring $S$. Then an exact sequence of $R$-modules is $(R, S)$-exact if
the kernel of each map is a direct summand of the source as an $S$-module. Then an $R$-module $P$ is $(R, S)$-projective if, for every $(R, S)$-exact sequence

$$
0 \rightarrow M_{1} \xrightarrow{f} M_{2} \xrightarrow{g} M_{3} \rightarrow 0
$$

and every $R$-linear map $h: P \rightarrow M_{3}$, there is a lifting of this map $h^{\prime}: P \rightarrow M_{2}$ meaning that $h=g \circ h^{\prime}$.

An FI-module $M$ is semi-induced if it has a finite filtration $M=M_{r} \supset M_{r-1} \supset \cdots \supset$ $M_{1} \supset M_{0}=0$ such that $M_{i} / M_{i-1}$ is an induced FI-module for $i=1, \ldots, r$. Consider the functor FI $\rightarrow$ FI which sends a finite set $S$ to the disjoint union with a new element $S \amalg\{*\}$ and an injection $f: S \rightarrow T$ to the corresponding injection $f^{\prime}: S \amalg\{*\} \rightarrow T \amalg\{*\}$ which is $f$ on $S$ and sends $*$ to $*$. Given an FI-module $M$, we let $\Sigma M$ be the pullback of $M$ along this functor. Explicitly, this means that $(\Sigma M)(S)=M(S \amalg\{*\})$ and that the action of $f$ comes from the action of $f^{\prime}$. Then $\Sigma M$ is called the shift of $M$ since $(\Sigma M)(n)=M(n+1)$.

There is a canonical map $M \rightarrow \Sigma M$ where the map $M(S) \rightarrow(\Sigma M)(S)$ comes from the inclusion $S \rightarrow S \amalg\{*\}$. By composing, we also get a canonical map $M \rightarrow \Sigma^{n} M$ for each $n$.

Exercise 7.4.3. Show that the kernel of $M \rightarrow \Sigma^{n} M$ is a torsion module.
Let $\Delta M$ be the cokernel of the canonical map $M \rightarrow \Sigma M$ and let $h_{M}(n)=\operatorname{dim}_{\mathbf{k}} M(n)$. By the previous exercise and Lemma 7.3.1, the kernel of $M \rightarrow \Sigma M$ is 0 in sufficiently large degrees. In particular, for $n \gg 0$, we have

$$
h_{\Delta M}(n)=h_{\Sigma M}(n)-h_{M}(n)=h_{M}(n+1)-h_{M}(n) .
$$

Exercise 7.4.4. If $M$ is generated in degrees $\leq d$, show that $\Delta M$ is generated in degrees $\leq d-1$. By induction on $d$, and using the above equation, show that $h_{M}(n)$ agrees with a polynomial of degree $\leq d$ for $n \gg 0$.

The polynomial in the previous exercise will be denoted $p_{M}(n)$ and is called the Hilbert polynomial of $M$. Define the degree of $M$ to be the degree of $p_{M}(n)$, and denote it $\delta(M)$. The degree of the 0 polynomial is -1 , by convention. So assuming $\delta(M) \geq 0$, we have

$$
\begin{equation*}
\delta(\Delta M) \leq \delta(M)-1 \tag{7.4.5}
\end{equation*}
$$

In fact, for the same reason, the same holds for the cokernel of the composition $M \rightarrow \Sigma M \rightarrow$ $\Sigma^{n} M$ for any $n$.

We have the following important theorem about shifts and semi-induced modules [Na]:
Theorem 7.4.6 (Nagpal). If $M$ is a finitely generated FI-module, then $\Sigma^{n} M$ is semi-induced for $n \gg 0$.

We omit the proof, but note the following corollary.
Corollary 7.4.7. Let $M$ be a finitely generated FI-module and set $k=\delta(M)$. There exists a complex of FI-modules

$$
0 \rightarrow M \rightarrow I^{0} \rightarrow I^{1} \rightarrow \cdots \rightarrow I^{k} \rightarrow 0
$$

such that $I^{j}$ is semi-induced for all $j$, and the homology of the complex is torsion.
Proof. We prove this by induction on $k$. If $k=-1$, then there is nothing to show since $M$ is torsion and the complex $0 \rightarrow M \rightarrow 0$ satisfies the above conditions. In general for $k \geq 0$,
take $n$ large enough so that $\Sigma^{n} M$ is semi-induced and set $I^{0}=\Sigma^{n} M$. Then the cokernel $N$ of the map $M \rightarrow \Sigma^{n} M$ satisfies $\delta(N) \leq \delta(M)-1$ and so there exists a complex

$$
0 \rightarrow N \rightarrow I^{1} \rightarrow \cdots \rightarrow I^{k} \rightarrow 0
$$

such that each $I^{j}$ is semi-induced and the homology is torsion. Define $I^{0} \rightarrow I^{1}$ to be the composition $I^{0} \rightarrow N \rightarrow I^{1}$. Then the image of $I^{0} \rightarrow I^{1}$ is the same as the image of $N \rightarrow I^{1}$ and the kernel is the preimage of the kernel of $N \rightarrow I^{1}$ under the surjection $I^{0} \rightarrow N$. So the homology remains torsion and we have the desired complex.
7.5. Cohomology of FI-modules. In $\S 4.3$, we studied group homology with coefficients in a finitely generated FI-module $M$. In particular, in each degree, the direct sum $\bigoplus_{n} \mathrm{H}_{i}\left(\Sigma_{n} ; M(n)\right)$ has the structure of a finitely generated $\mathbf{k}[t]$-module. Now we investigate the situation with group cohomology, which will utilize Nagpal's theorem above.

Example 7.5.1. Let $\mathbf{k}$ be a field of positive characteristic $p>0$. Let $\mathbf{k}^{n}$ be the standard permutation representation of $\Sigma_{n}$, and let $V_{n}$ be the subspace of $\mathbf{k}^{n}$ consisting of vectors whose coordinates sum to 0 . If $p$ divides $n$, then the line spanned by the all 1 's vector is in $V_{n}$ and hence $V_{n}^{\Sigma_{n}}$ is 1-dimensional. Otherwise, it is 0-dimensional. In other words,

$$
\operatorname{dim}_{\mathbf{k}} \mathrm{H}^{0}\left(\Sigma_{n} ; V_{n}\right)=\left\{\begin{array}{ll}
1 & \text { if } p \text { divides } n \\
0 & \text { else }
\end{array},\right.
$$

so the sequence of dimensions exhibits periodic behavior. Note that $[n] \mapsto \mathbf{k}^{n}$ is a finitely generated FI-module and that $[n] \mapsto V_{n}$ is a submodule of it. We will see more generally that periodic behavior occurs in the cohomology of finitely generated FI-modules.

Let $M$ be an FI-module and, for each $i \geq 0$, define

$$
\mathcal{H}^{i}(M)=\bigoplus_{n \geq 0} \mathrm{H}^{i}\left(\Sigma_{n} ; M(n)\right), \quad \mathcal{H}^{*}(M)=\bigoplus_{i \geq 0} \mathcal{H}^{i}(M) .
$$

Then $\mathcal{H}^{*}(M)$ is bigraded where $\mathrm{H}^{i}\left(\Sigma_{n} ; M(n)\right)$ has bidegree $(i, n)$. Recall that in $\S 3$, we discussed the equivalence between the category of FI-modules and the category of modules over the twisted commutative algebra $A$ which satisfies $A_{n}=\mathbf{k}$ for all $n$. To do this, we defined a monoidal structure on the category of symmetric sequences, which in one form becomes:

$$
(V \otimes W)_{n}=\bigoplus_{i=0}^{n} \operatorname{Ind}_{\Sigma_{i} \times \Sigma_{n-i}}^{\Sigma_{n}}\left(V_{i} \otimes W_{n-i}\right)
$$

Given non-negatively bigraded vector spaces $A$ and $B$ (with finite-dimensional components), their tensor product is a bigraded vector space satisfying

$$
(A \otimes B)_{i, j}=\bigoplus_{\substack{i^{\prime}+i^{\prime \prime}=i \\ j^{\prime}+j^{\prime \prime}=j}} A_{i^{\prime}, j^{\prime}} \otimes B_{i^{\prime \prime}, j^{\prime \prime}}
$$

So we have a category of non-negatively bigraded vector spaces with this tensor product as its monoidal structure, and $\mathcal{H}^{*}(M)$ is an object of it.

Proposition 7.5.2. The assignment $M \mapsto \mathcal{H}^{*}(M)$ is a monoidal functor, i.e., $\mathcal{H}^{*}(M \otimes N)$ is naturally isomorphic to $\mathcal{H}^{*}(M) \otimes \mathcal{H}^{*}(N)$.

Proof. Fix $(i, n)$. Then we have

$$
\begin{aligned}
\mathcal{H}^{*}(M \otimes N)_{i, n} & =\mathrm{H}^{i}\left(\Sigma_{n} ;(M \otimes N)(n)\right) \\
& =\mathrm{H}^{i}\left(\Sigma_{n} ; \bigoplus_{k=0}^{n} \operatorname{Ind}_{\Sigma_{k} \times \Sigma_{n-k}}^{\Sigma_{n}}(M(k) \otimes N(n-k))\right) \\
& =\bigoplus_{k=0}^{n} \mathrm{H}^{i}\left(\Sigma_{k} \times \Sigma_{n-k} ; M(k) \otimes N(n-k)\right) \\
& =\bigoplus_{k=0}^{n} \bigoplus_{i^{\prime}+i^{\prime \prime}=i} \mathrm{H}^{i^{\prime}}\left(\Sigma_{k} ; M(k)\right) \otimes \mathrm{H}^{i^{\prime \prime}}\left(\Sigma_{n-k} ; N(n-k)\right) \\
& =\bigoplus_{\substack{i^{\prime}+i^{\prime \prime}=i \\
n^{\prime}+n^{\prime \prime}=n}} \mathcal{H}^{*}(M)_{i^{\prime}, n^{\prime}} \otimes \mathcal{H}^{*}(N)_{i^{\prime \prime}, n^{\prime \prime}},
\end{aligned}
$$

where the first two equalities are by definition, the third equality is Shapiro's lemma, and the fourth equality is the Künneth isomorphism.

In particular, since $A$ is an algebra, applying $\mathcal{H}$ to the multiplication map $A \otimes A \rightarrow A$ yields a map

$$
\mathcal{H}(A) \otimes \mathcal{H}(A) \rightarrow \mathcal{H}(A)
$$

By degree considerations, this restricts to a map $\mathcal{H}^{0}(A) \otimes \mathcal{H}^{0}(A) \rightarrow \mathcal{H}^{0}(A)$, and this is a commutative, associative algebra satisfying $\mathcal{H}^{0}(A)=\bigoplus_{n} \mathrm{H}^{0}\left(\Sigma_{n} ; \mathbf{k}\right)$. Since we have a canonical identification $\mathrm{H}^{0}\left(\Sigma_{n} ; \mathbf{k}\right)=\mathbf{k}$, we let $x^{[n]}$ be the element corresponding to $1 \in \mathbf{k}$.
Lemma 7.5.3. For all $n, m$, we have

$$
x^{[n]} x^{[m]}=\binom{n+m}{n} x^{[n+m]} .
$$

Proof. Note that the map $\mathrm{H}^{0}\left(\Sigma_{n} ; \mathbf{k}\right) \otimes \mathrm{H}^{0}\left(\Sigma_{m} ; \mathbf{k}\right) \rightarrow \mathrm{H}^{0}\left(\Sigma_{n+m} ; \mathbf{k}\right)$ is multiplication by some scalar, and we just have to show that it is $\binom{n+m}{m}$. In fact, this follows from the fact that the isomorphism coming from Shapiro's lemma can be described as a transfer map [Wei, Exercise 6.7.7]: in general, for a subgroup $H$ of $G$ and a representation $M$ of $G$, the transfer map is induced by the map $M^{H} \rightarrow M^{G}$ given by $m \mapsto \sum g m$ where the sum is over a set of representatives of $G / H$. In our case, $M=\mathbf{k}$ is trivial and so this is multiplication by the index of $\Sigma_{n} \times \Sigma_{m}$ in $\Sigma_{n+m}$, which is $\binom{n+m}{n}$.

This algebra is the divided power algebra in 1 variable, and we will denote it $\mathbf{D}$. Note that if $\mathbf{k}$ is a field of characteristic 0 , then $\mathbf{D} \cong \mathbf{k}[t]$ under the identification $t^{n} / n!=x^{[n]}$. Over a field of positive characteristic, $\mathbf{D}$ has many zerodivisors, and is not even finitely generated or noetherian.
Exercise 7.5.4. Let k be a field of characteristic $p>0$. Define a ring homomorphism $\mathbf{k}\left[y_{0}, y_{1}, \ldots\right] \rightarrow \mathbf{D}$ by $y_{i} \mapsto x^{\left[p^{i}\right]}$. Show that this is surjective and that the kernel is the ideal $\left(y_{0}^{p}, y_{1}^{p}, y_{2}^{p}, \ldots\right)$.

A commutative algebra $R$ is coherent if every finitely generated ideal is a finitely presented $R$-module.
Exercise 7.5.5. Let $R$ be a coherent ring and let $M$ and $N$ be finitely presented $R$-modules. Given a homomorphism $f: M \rightarrow N$, show that $\operatorname{ker} f$ and coker $f$ are finitely presented.

Proposition 7.5.6. D is a coherent ring.
Proof. If $\mathbf{k}$ is a field of characteristic 0 , then $\mathbf{D} \cong \mathbf{k}[t]$, so it is noetherian, and hence coherent.
Otherwise, suppose $\mathbf{k}$ is a field of positive characteristic $p>0$. As in the previous exercise, identify $\mathbf{D}$ with $\mathbf{k}\left[y_{0}, y_{1}, \ldots\right] /\left(y_{0}^{p}, y_{1}^{p}, \ldots\right)$. Any finite list of generators for an ideal only uses finitely many of the $y_{i}$, say the first $r$ of them. Note that $\mathbf{D}$ is a free module over the subring $R=\mathbf{k}\left[y_{0}, \ldots, y_{r}\right] /\left(y_{0}^{p}, \ldots, y_{r}^{p}\right)$, and so its presentation over $\mathbf{D}$ can be obtained by taking a free presentation over $R$ and then applying $-\otimes_{R} \mathbf{D}$.

Proposition 7.5.7. Suppose $\mathbf{k}$ is a field of positive characteristic $p>0$. Let $M$ be a finitely presented graded $\mathbf{D}$-module. Then $n \mapsto \operatorname{dim}_{\mathbf{k}} M_{n}$ is a periodic function, whose period is a power of $p$, for $n \gg 0$.
Proof. As in the previous proof, identify $\mathbf{D}$ with $\mathbf{k}\left[y_{0}, y_{1}, \ldots\right] /\left(y_{0}^{p}, y_{1}^{p}, \ldots\right)$. Pick a finite presentation of $M$, i.e., write it as the cokernel of a finite matrix with homogeneous entries from D. It involves only finitely many of the $y_{i}$, say $y_{1}, \ldots, y_{r}$, and let $R=\mathbf{k}\left[y_{0}, \ldots, y_{r}\right] /\left(y_{0}^{p}, \ldots, y_{r}^{p}\right)$. Let $M^{\prime}$ be the cokernel of the same matrix thought of now as a map of free modules over $R$.

Since $R$ is finite-dimensional over $\mathbf{k}$, the same is true for $M^{\prime}$. Furthermore, $\mathbf{D}$ is a free module over $R$ with basis given by the monomials in $y_{r+1}, y_{r+2}, \ldots$ with exponents between 0 and $p-1$. In particular, $M=M^{\prime} \otimes_{R} \mathbf{D}$, and as a vector space, $M$ is the tensor product of $M^{\prime}$ with the vector space whose basis is the set of monomials just listed. This implies that

$$
\sum_{n \geq 0}\left(\operatorname{dim}_{\mathbf{k}} M_{n}\right) t^{n}=\frac{\sum_{n \geq 0}\left(\operatorname{dim}_{\mathbf{k}} M_{n}^{\prime}\right) t^{n}}{1-t^{p^{r+1}}}
$$

As just discussed, the numerator on the right hand side is a polynomial and hence the function $n \mapsto \operatorname{dim}_{\mathbf{k}} M_{n}$ is periodic of period $p^{r+1}$ beyond the degree of the numerator.

The k-linear map $d: \mathbf{D} \rightarrow \mathbf{D}$ defined by $d\left(x^{[n]}\right)=x^{[n-1]}$ is a derivation: for all $f, g \in \mathbf{D}$, we have $d(f g)=d(f) g+f d(g)$. Let $M$ be a graded $\mathbf{D}$-module. A connection on $M$ is a linear map $\nabla: M \rightarrow M$ that sends a homogeneous element of degree $n$ to a homogeneous element of degree $n-1$, and satisfies $\nabla(f m)=d(f) m+f \nabla(m)$ for all $f \in \mathbf{D}$ and $m \in M$.

Proposition 7.5.8. Suppose $M$ has a connection $\nabla$. Then the map

$$
\begin{aligned}
& \psi: \mathbf{D} \otimes_{\mathbf{k}} \operatorname{ker} \nabla \rightarrow M \\
& \sum_{i} f_{i} \otimes m_{i} \mapsto \sum_{i} f_{i} m_{i}
\end{aligned}
$$

is an isomorphism. In particular, $M$ is a free $\mathbf{D}$-module.
Proof. First, we show that $\psi$ is injective. If not, then we have $\sum_{i} f_{i} m_{i}=0$ where the $m_{i}$ are linearly independent elements of ker $\nabla$. Furthermore, we may choose this so that the elements are homogeneous and so that $\operatorname{deg}\left(f_{i}\right)+\operatorname{deg}\left(m_{i}\right)$ is constant, and minimal amongst all relations that exist. In particular, $\operatorname{deg}\left(f_{i}\right)>0$ for some $i$. Applying $\nabla$ to this linear combination, we get

$$
0=\sum_{i}\left(d\left(f_{i}\right) m_{i}+f_{i} \nabla\left(m_{i}\right)\right)=\sum_{i} d\left(f_{i}\right) m_{i}
$$

This gives a linear combination of smaller degree, which contradicts our choice from before, so $\psi$ is injective as claimed.

Now we show that $\psi$ is surjective. Let $m \in M$ be an arbitrary element. Let $r$ be the minimal integer such that $\nabla^{r} m=0$ ( $r$ exists since $\nabla$ is degree decreasing). By induction on $r$, we show that $m$ is in the image of $\psi$. If $r=0$, there is nothing to show, so suppose it holds for all values strictly smaller than $r$. Note that

$$
m+\sum_{i=1}^{r-1}(-1)^{i} x^{[i]} \nabla^{i}(m) \in \operatorname{ker} \nabla
$$

By induction, each $\nabla^{i}(m)$ for $i>0$ is in the image of $\psi$. Since every element of ker $\nabla$ is in the image of $\psi$, we conclude that $m$ is in the image of $\psi$.

By Proposition 7.5.2 and Lemma 7.5.3, for each $i \geq 0, \mathcal{H}^{i}(M)=\bigoplus_{n \geq 0} H^{i}\left(\Sigma_{n} ; M(n)\right)$ has the structure of a $\mathbf{D}$-module.

Proposition 7.5.9. Let $\mathbf{k}$ be a field of positive characteristic. The restriction map $\mathrm{H}^{i}\left(\Sigma_{n} ; \mathbf{k}\right) \rightarrow$ $\mathrm{H}^{i}\left(\Sigma_{n-1} ; \mathbf{k}\right)$ gives a connection on $\mathcal{H}^{i}\left(P_{0}\right)$. Furthermore, ker $\nabla$ is 0 in degrees $>2 i$ and hence is finite-dimensional.

Proof. We omit the proof of the first part, see [NS, §4.2] for the calculation. The second part is the dual version of Nakaoka's stability theorem (Theorem 4.2.1) and the general fact that each of the (co)homology groups of a finite group with trivial coefficients is finite-dimensional (as follows from example, from the existence of the bar resolution [Wei, §6.5]).

Given a graded $\mathbf{D}$-module $N$ and an integer $n$, we let $N^{\geq n}$ be the submodule of $N$ generated by all elements of degree $\geq n$.

Theorem 7.5.10 (Nagpal-Snowden). Let $\mathbf{k}$ be a field of positive characteristic $p>0$. Let $M$ be a finitely generated FI-module over $\mathbf{k}$.
(1) There exists a finitely presented $\mathbf{D}$-module $N$ and an integer $m$ such that $N^{\geq m} \cong$ $\mathcal{H}^{i}(M) \geq{ }^{\geq m}$.
(2) The function $n \mapsto \operatorname{dim}_{\mathbf{k}} \mathrm{H}^{i}\left(\Sigma_{n} ; M(n)\right)$ is a periodic function, whose period is a power of $p$, for $n \gg 0$.

Proof. We reduce (1) to the case of semi-induced modules using Corollary 7.4.7, Exercise 7.5.5, and Proposition 7.5.6. Furthermore, using the long exact sequence on cohomology and doing induction on the length of the filtration for a semi-induced module whose quotients are induced, we can reduce to the case of an induced module $\mathcal{J}(V)$ with $V$ a representation of $\Sigma_{n}$.

Recall that in the tca perspective, $\mathcal{J}(V)$ is the tensor product $A \otimes V$ where $A$ is the tca freely generated by a single element of degree 1 . The action of $A$ is a map $A \otimes A \otimes V \rightarrow A \otimes V$, which is in fact the multiplication map on $A$ tensored with the identity on $V$. By Proposition 7.5.2, after applying $\mathcal{H}^{*}$, the corresponding map $\mathcal{H}^{*}(A) \otimes \mathcal{H}^{*}(A) \otimes \mathcal{H}^{*}(V) \rightarrow \mathcal{H}^{*}(A) \otimes \mathcal{H}^{*}(V)$ is the tensor product of the multiplication map on $\mathcal{H}^{*}(A)$ with the identity on $\mathscr{H}^{*}(V)$. $\mathrm{So}^{i}(\mathcal{J}(V))$ is the $\mathbf{D}$-module $\bigoplus_{j=0}^{i} \mathcal{H} \mathcal{H}^{i-j}(A) \otimes \mathrm{H}^{j}\left(\Sigma_{n} ; V\right)$, and in particular, is a finitely generated free D-module by Propositions 7.5.8 and 7.5.9.
(2) follows from (1) and Proposition 7.5.7.

Finally, we return to the discussion of the cohomology of unordered configuration spaces of a topological space $X$ in $\S 5$. Recall the definitions:

$$
\begin{aligned}
\operatorname{Conf}_{n}(X) & =\left\{\left(x_{1}, \ldots, x_{n}\right) \in X^{n} \mid x_{i} \neq x_{j} \text { for } i \neq j\right\}, \\
\operatorname{UConf}_{n}(X) & =\operatorname{Conf}_{n}(X) / \Sigma_{n}
\end{aligned}
$$

Theorem 7.5.11. Let $\mathbf{k}$ be a field of positive characteristic $p$. Let $X$ be a topological space with the property that the FI-module $[n] \mapsto \mathrm{H}^{i}\left(\operatorname{Conf}_{n}(X) ; \mathbf{k}\right)$ is finitely generated for all $i$. Then for each $i$, the function $n \mapsto \operatorname{dim}_{\mathbf{k}} \mathrm{H}^{i}\left(\operatorname{UConf}_{n}(X) ; \mathbf{k}\right)$ is periodic, with period a power of $p$, for $n \gg 0$.

Proof. Recall from our discussion on cohomology of configuration spaces that we have a Cartan-Leray spectral sequence (5.2.3):

$$
\mathrm{E}_{2}^{p, q}=\mathrm{H}^{p}\left(\Sigma_{n} ; \mathrm{H}^{q}\left(\operatorname{Conf}_{n}(X) ; \mathbf{k}\right)\right) \Longrightarrow \mathrm{H}^{p+q}\left(\operatorname{UConf}_{n}(X) ; \mathbf{k}\right) .
$$

We take the direct sum of these spectral sequences over $n$. We focus on the case $p+q=i$. To compute the $\mathrm{E}_{\infty}$ page along this degree, we have to do computations with finitely many other degrees. By Theorem 7.5.10, each of the terms agrees with a finitely presented D-module after some truncation. Since there are finitely many terms, we can find a single truncation that works for all of them. By Exercise 7.5.5 and Proposition 7.5.6, this implies that the $\mathrm{E}_{\infty^{-}}{ }^{-}$ terms in degree $i$ agree with a finitely presented D-module after some truncation. Finally, the term we want to compute, $\bigoplus_{n} \mathrm{H}^{i}\left(\operatorname{UConf}_{n}(X) ; \mathbf{k}\right)$, is then built out of these terms using a finite number of extensions, so is itself a finitely presented $\mathbf{D}$-module after some truncation. Now we use Proposition 7.5.7.

## 8. $\Delta$-MODULES

In $\S 6$, we studied the notion of tensor rank. Here we will come back to the rank 1 case from a different perspective and study its higher order syzygies.
8.1. Segre embeddings. First, we rephrase the setting. Given a vector space $V$, let $\mathbf{P}(V)$ denote the set of lines in $V$. Then $\mathbf{P}(V)$ is called the projectivization of $V$, and also projective space. Given a homogeneous polynomial $f \in \operatorname{Sym}\left(V^{*}\right)$, it does not make sense to ask for its value on a line $\ell \in \mathbf{P}(V)$, but it is well-defined whether $f(\ell)=0$ or $f(\ell) \neq 0$. In particular, to any collection $I$ of homogeneous polynomials, we can associate its vanishing locus $V(I)$. If we declare these to be closed subsets, then we get a topology on $\mathbf{P}(V)$, which is a graded variant of the Zariski topology on $\operatorname{Spec} R$ for a commutative ring $R$.

We focus on the Segre embedding:

$$
\begin{aligned}
\mathbf{P}\left(V_{1}\right) \times \cdots \times \mathbf{P}\left(V_{n}\right) & \rightarrow \mathbf{P}\left(V_{1} \otimes \cdots \otimes V_{n}\right) \\
\left(\ell_{1}, \ldots, \ell_{n}\right) & \mapsto \ell_{1} \otimes \cdots \otimes \ell_{n} .
\end{aligned}
$$

Exercise 8.1.1. Check that the Segre embedding is in fact injective.
Note that if we take the union of the set of lines in the image of the Segre embedding, we get exactly the rank 1 tensors (and 0 ). So the image is a closed subset cut out by the same equations described in Remark 6.2.4.

Remark 8.1.2. The Segre embedding is not continuous if we put the product topology on the left hand side. Instead, we give the product of projective spaces the subspace topology coming from this embedding.

The Segre embeddings have a certain "factorization property": for example, when $n=3$, we can break it up into two steps:

$$
\mathbf{P}\left(V_{1}\right) \times \mathbf{P}\left(V_{2}\right) \times \mathbf{P}\left(V_{3}\right) \rightarrow \mathbf{P}\left(V_{1} \otimes V_{2}\right) \times \mathbf{P}\left(V_{3}\right) \rightarrow \mathbf{P}\left(V_{1} \otimes V_{2} \otimes V_{3}\right)
$$

Using this, and a certain linearity property, all of the equations in Remark 6.2.4 can be built out of the simplest instance of a Segre embedding when $n=2$ and $\operatorname{dim} V_{1}=\operatorname{dim} V_{2}=2$.

Let's consider that case in depth, but we'll use dual spaces to make the functions not have a dual. Pick bases $e_{0}, e_{1}$ and $f_{0}, f_{1}$ for $V_{1}^{*}$ and $V_{2}^{*}$ and pick dual bases $x_{0}, x_{1}$ and $y_{0}, y_{1}$ so that $\operatorname{Sym}\left(V_{1}\right)=\mathbf{k}\left[x_{0}, x_{1}\right]$ and $\operatorname{Sym}\left(V_{2}\right)=\mathbf{k}\left[y_{0}, y_{1}\right]$. Let $z_{i, j}=x_{i} \otimes y_{j}$ so that $\operatorname{Sym}\left(V_{1} \otimes V_{2}\right)=\mathbf{k}\left[z_{0,0}, z_{0,1}, z_{1,0}, z_{1,1}\right]$. The image of the Segre embedding consists of all lines spanned by rank 1 matrices, so it is the solution set of the equation

$$
\operatorname{det}\left(\begin{array}{cc}
z_{0,0} & z_{0,1} \\
z_{1,0} & z_{1,1}
\end{array}\right)=z_{0,0} z_{1,1}-z_{1,0} z_{0,1}
$$

The image in this case is also known as the quadric surface. As a representation of $\mathbf{G L}\left(V_{1}\right) \times$ $\mathbf{G L}\left(V_{2}\right)$, the subspace generated by this equation in the space of degree 2 polynomials is

$$
\bigwedge^{2}\left(V_{1}\right) \otimes \bigwedge^{2}\left(V_{2}\right) \subset \operatorname{Sym}^{2}\left(V_{1} \otimes V_{2}\right)
$$

In fact, this implies that for arbitrary $V_{1}, V_{2}$, the subspace $\bigwedge^{2}\left(V_{1}\right) \otimes \bigwedge^{2}\left(V_{2}\right)$ consists of equations that vanish on the image of the Segre embedding since it is true for every choice of 2-dimensional summands $V_{1}^{\prime} \subset V_{1}$ and $V_{2}^{\prime} \subset V_{2}$ and every element of $\bigwedge^{2}\left(V_{1}\right) \otimes \bigwedge^{2}\left(V_{2}\right)$ is a linear combination of elements living in summands of the form $\Lambda^{2}\left(V_{1}^{\prime}\right) \otimes \Lambda^{2}\left(V_{2}^{\prime}\right)$. Alternatively, we can argue that the space of degree 2-equations must be a polynomial subfunctor of $\left(V_{1}, V_{2}\right) \mapsto \operatorname{Sym}^{2}\left(V_{1} \otimes V_{2}\right)$ in the following sense.

We recall the notion of polynomial functor from $\S 2.8 .4$ and generalize it at the same time. Let $\mathbf{k}$ be an infinite field and let $\mathbf{V e c}=\mathbf{V e c}_{\mathbf{k}}$ be the category of finite-dimensional vector spaces. Let $\mathbf{V e c}_{\mathbf{k}}^{n}=\mathbf{V e c}^{n}=\mathbf{V e c}_{\mathbf{k}} \times \cdots \times \mathbf{V e c}_{\mathbf{k}}$, the $n$-fold product of this category, so objects are $n$-tuples of vector spaces and morphisms are $n$-tuples of linear maps. A functor $F: \mathbf{V e c}^{n} \rightarrow$ Vec is a polynomial functor if, for every pair of $n$-tuples $\left(V_{1}, \ldots, V_{n}\right)$ and $\left(W_{1}, \ldots, W_{n}\right)$, the map

$$
\operatorname{Hom}_{\mathbf{k}}\left(V_{1}, W_{1}\right) \times \cdots \times \operatorname{Hom}_{\mathbf{k}}\left(V_{n}, W_{n}\right) \rightarrow \operatorname{Hom}_{\mathbf{k}}\left(F\left(V_{1}, \ldots, V_{n}\right), F\left(W_{1}, \ldots, W_{n}\right)\right)
$$

is given by polynomial functions.
In particular, if we pick bases as above with the same convention, then all equations of the form $z_{i, j} z_{i^{\prime}, j^{\prime}}-z_{i, j^{\prime}} z_{i^{\prime}, j}$ vanish on the image.

Now let's consider the factorization property. For simplicity, assume $n=3$ and $\operatorname{dim} V_{i}=2$ for all $i$. We let $z_{i, j}$ be the coordinate functions on $\mathbf{P}\left(V_{1}^{*} \otimes V_{2}^{*}\right)$ and let $z_{i, j, k}$ denote the coordinate functions on $\mathbf{P}\left(V_{1}^{*} \otimes V_{2}^{*} \otimes V_{3}^{*}\right)$. By repeating the above with the two vector spaces $V_{1}^{*} \otimes V_{2}^{*}$ and $V_{3}^{*}$, we conclude that the equations

$$
z_{i, j, k} z_{i^{\prime}, j^{\prime}, k^{\prime}}-z_{i, j, k^{\prime}} z_{i^{\prime}, j^{\prime}, k^{\prime}}
$$

vanish on the image of $\mathbf{P}\left(V_{1}^{*} \otimes V_{2}^{*}\right) \times \mathbf{P}\left(V_{3}^{*}\right)$. Since the triple product $\mathbf{P}\left(V_{1}^{*}\right) \times \mathbf{P}\left(V_{2}^{*}\right) \times \mathbf{P}\left(V_{3}^{*}\right)$ is in there, these equations also vanish there. However, we could have grouped together two of the vector spaces in two other ways $\left(V_{1}^{*}\right.$ with $V_{3}^{*}$, but also $V_{2}^{*}$ with $\left.V_{3}^{*}\right)$. In particular, the
following two types of equations also vanish on the triple product:

$$
\begin{gathered}
z_{i, j, k} z_{i^{\prime}, j^{\prime}, k^{\prime}}-z_{i, j^{\prime}, k} z_{i^{\prime}, j, k^{\prime}} \\
z_{i, j, k} z_{i^{\prime}, j^{\prime}, k^{\prime}}-z_{i^{\prime}, j, k} z_{i, j^{\prime}, k^{\prime}}
\end{gathered}
$$

By taking sums of these 3 types of equations, we can get all of the equations described in Remark 6.2.4. This argument extends to $n$-fold products of projective spaces for arbitrary $n$.

In particular, we see that the equation for the quadric surface "generates" all other quadratic equations for all Segre embeddings using this factorization property together with linearity. As we discussed, the ideals of the Segre embeddings are generated in degree 2, so this provides a complete description. Now we would like to axiomatize this structure and make precise in what sense this generation is happening.
8.2. $\Delta$-modules. Define a category $\mathbf{V e c}^{\Delta}$ as follows. Its objects are finite collections of vector spaces $\left\{V_{i}\right\}_{i \in I}$ and a morphism $\left\{U_{i}\right\}_{i \in I} \rightarrow\left\{V_{j}\right\}_{j \in J}$ is a surjection $f: J \rightarrow I$ together with, for each $i \in I$, a linear map $\varepsilon_{i}: U_{i} \rightarrow \bigotimes_{j \in f^{-1}(i)} V_{j}$. Given another morphism $\left\{V_{j}\right\}_{j \in J} \rightarrow$ $\left\{W_{k}\right\}_{k \in K}$ with surjection $g: K \rightarrow J$ and linear maps $\zeta_{j}: V_{j} \rightarrow \bigotimes_{k \in g^{-1}(j)} W_{k}$, the composition is given by the data $h: K \rightarrow I$ where $h=f \circ g$ and linear maps $\eta_{i}: U_{i} \rightarrow \bigotimes_{k \in h^{-1}(i)} W_{k}$ given by the composition of $\varepsilon_{i}: U_{i} \rightarrow \bigoplus_{j \in f^{-1}(i)} V_{j}$ with the tensor product of the maps $\zeta_{j}: V_{j} \rightarrow \bigotimes_{k \in g^{-1}(j)} W_{k}$ over all $j \in f^{-1}(i)$.

For each $n$, note that $\mathbf{V e c}^{n}$ is a subcategory (not full) of $\mathbf{V e c}^{\Delta}$ consisting of the objects $\left(V_{1}, \ldots, V_{n}\right)$ and where the morphisms are the ones where $f:[n] \rightarrow[n]$ is the identity. Hence, given a functor $F: \operatorname{Vec}^{\Delta} \rightarrow \Delta$, we get a sequence of functors $F_{n}: \mathrm{Vec}^{n} \rightarrow \mathrm{Vec}$ by restricting to this subcategory $\mathbf{V e c}^{n}$. We will say that a functor $\mathbf{V e c}^{\Delta} \rightarrow \mathbf{V e c}$ is polynomial if each of the $F_{n}$ is a polynomial functor in the sense described above. For short, polynomial functors $\mathrm{Vec}^{\Delta} \rightarrow$ Vec are called $\Delta$-modules.

Unfolding this definition, we see that a $\Delta$-module is a sequence of polynomial functors $F_{n}: \mathbf{V e c}^{n} \rightarrow$ Vec together with some extra data that relates them. Since they are polynomial functors, we get an action of $\prod_{i=1}^{n} \mathbf{G L}\left(V_{i}\right)$ on $F_{n}\left(V_{1}, \ldots, V_{n}\right)$. Also, on the objects of the form $\left(V_{1}, \ldots, V_{n}\right)$ with all $V_{i}$ equal to a single vector space $V$, there is an action of the symmetric group $\Sigma_{n}$ on $F_{n}(V, \ldots, V)$. The other key structure is that we have transition maps between the various $F_{n}$.
Example 8.2.1. The most basic example of a $\Delta$-module is the functor

$$
T:\left(V_{i}\right)_{i \in I} \mapsto \bigotimes_{i \in I} V_{i}
$$

To get other examples, we can compose this with a polynomial functor Vec $\rightarrow \mathbf{V e c}$, such as symmetric and exterior powers (or even tensor products of such functors). So other examples of $\Delta$-modules are,

$$
\begin{aligned}
& \left(V_{i}\right)_{i \in I} \mapsto \operatorname{Sym}^{r}\left(\bigotimes_{i \in I} V_{i}\right), \\
& \left(V_{i}\right)_{i \in I} \mapsto \bigwedge^{r}\left(\bigotimes_{i \in I} V_{i}\right)
\end{aligned}
$$

where $r$ is some fixed nonnegative integer. We call them $\operatorname{Sym}^{r}(T)$ and $\bigwedge^{r}(T)$. Let $\operatorname{Sym}(T)=$ $\bigoplus_{r \geq 0} \operatorname{Sym}^{r}(T)$ and $\bigwedge(T)=\bigoplus_{r \geq 0} \Lambda^{r}(T)$.

An element of a $\Delta$-module $F$ is an element of some $F\left(\left\{V_{i}\right\}_{i \in I}\right)$. Given any collection of elements, there is a smallest $\Delta$-submodule containing it, which we call the submodule generated by these elements. We say that $F$ is finitely generated if it can be generated by a finite collection of elements.

The basic functor $T$ is finitely generated by a nonzero element in $F(\mathbf{k})$ where $\mathbf{k}$ is the 1-dimensional vector space indexed by a singleton set.

Exercise 8.2.2. Show that the tensor power $T^{r}$ defined by

$$
T^{r}\left(\left\{V_{i}\right\}_{i \in I}\right)=\bigotimes_{i \in I} V_{i}^{\otimes r}
$$

is a finitely generated $\Delta$-module. In particular, finite tensor powers of symmetric and exterior powers (being quotients of $T^{r}$ ) are also finitely generated.

A key property is that the finite generation property is inherited by submodules (see [SS1, Theorem 9.2.3]):

Theorem 8.2.3. Let $F$ be a finitely generated $\Delta$-module. Then $F$ is noetherian, i.e., all $\Delta$-submodules of $F$ are again finitely generated.

A key example of a submodule of $\operatorname{Sym}^{r}(T)$ is to take the space of degree $r$ equations which vanish on the Segre embedding of $\prod_{i \in I} \mathbf{P}\left(V_{i}^{*}\right)$. The previous section explains that for $r=2$, this is generated by a single element. The general theorem says that this is finitely generated for any $r$.

As a final example, we note that one can define a Koszul complex:

$$
\cdots \rightarrow \operatorname{Sym}(T) \otimes \bigwedge^{r}(T) \rightarrow \cdots \rightarrow \operatorname{Sym}(T) \otimes \bigwedge^{2}(T) \rightarrow \operatorname{Sym}(T) \otimes T \rightarrow \operatorname{Sym}(T)
$$

This is the standard Koszul complex for $\operatorname{Sym}\left(\otimes_{i \in I} V_{i}\right)$ when evaluated on $\left\{V_{i}\right\}_{i \in I}$, and one just has to check that it is compatible with all transition maps (this follows from naturality of the Koszul complex). Taking the quotient of $\operatorname{Sym}(T)$ by the submodule of equations vanishing on the Segre embeddings, we get a $\Delta$-module that sends $\left\{V_{i}\right\}_{i \in I}$ to the homogeneous coordinate ring of the Segre embedding of $\prod_{i \in I} \mathbf{P}\left(V_{i}^{*}\right)$, call it Seg, and its degree $d$ piece $\mathrm{Seg}_{d}$.

The $r$ th homology of the tensor product of the Koszul complex with Seg (this computes Tor with the residue field $\mathbf{k}$ ) gives the space of $r$-syzygies of the Segre embedding. If we restrict to a single degree, then each of the terms involved is finitely generated, so the noetherianity result above implies:

Proposition 8.2.4. For each $r, d$, the $\Delta$-module given by

$$
\left\{V_{i}\right\}_{i \in I} \mapsto \operatorname{Tor}_{r}^{\operatorname{Sym}\left(\otimes_{i \in I} V_{i}\right)}\left(\operatorname{Seg}\left(\bigotimes_{i \in I} V_{i}\right), \mathbf{k}\right)_{d}
$$

is finitely generated. Here the subscript d denotes the space of r-syzygies of degree $d$.
For $r=1$, this is computing the space of minimal degree $d$ equations of the Segre embedding, so this space is only nonzero for $d=2$ by our previous discussions. In general, one
can show that the space of $r$-syzygies is 0 in degrees $>2 r,{ }^{5}$ so it is actually superfluous to include the degree $d$ in the above result.

Our original heuristic is that the equations of the Segre embedding are all generated, in some sense, by a single one (the $2 \times 2$ determinant). This result says that a similar thing is true if one considers the higher-order syzygies that these equations satisfy: they all come from finitely many basic syzygies. In general, one does not know how to write them all down.

## Appendix A. Review of homological algebra

For a thorough treatment of homological algebra, see Weibel's text [Wei].
A.1. Exact sequences and exact functors. Let $\mathcal{A}$ be an abelian category (the category of $R$-modules if you like). Suppose we are given a chain complex of objects in $\mathcal{A}$

$$
\mathbf{V}: \cdots \rightarrow V_{i+1} \xrightarrow{f_{i+1}} V_{i} \xrightarrow{f_{i}} V_{i-1} \rightarrow \cdots
$$

The $i$ th homology of $\mathbf{V}$, denoted $\mathrm{H}_{i}(\mathbf{V})$, is

$$
\mathrm{H}_{i}(\mathbf{V})=\operatorname{ker}\left(f_{i}\right) / \operatorname{image}\left(f_{i+1}\right) .
$$

The complex $\mathbf{V}$ is exact if $H_{i}(\mathbf{V})=0$ for all $i$. If our complex has finitely many terms, then we only require exactness everywhere except the end points. In general, elements of $\operatorname{ker}\left(f_{i}\right)$ for some $i$ are cycles and elements of image $\left(f_{i}\right)$ for some $i$ are boundaries.

Sequences of the form

$$
0 \rightarrow A \xrightarrow{f} B \xrightarrow{g} C \rightarrow 0
$$

are called short exact sequences. If $\mathcal{B}$ is another abelian category and $F: \mathcal{A} \rightarrow \mathcal{B}$ is a functor (with $F(0)=0$ ), we get a corresponding sequence $0 \rightarrow F(A) \rightarrow F(B) \rightarrow F(C) \rightarrow 0$. Then $F$ is

- exact if $0 \rightarrow F(A) \rightarrow F(B) \rightarrow F(C) \rightarrow 0$ is an exact sequence for all exact sequences $0 \rightarrow A \rightarrow B \rightarrow C \rightarrow 0$,
- left exact if $0 \rightarrow F(A) \rightarrow F(B) \rightarrow F(C)$ is an exact sequence for all exact sequences $0 \rightarrow A \rightarrow B \rightarrow C \rightarrow 0$,
- right exact if $F(A) \rightarrow F(B) \rightarrow F(C) \rightarrow 0$ is an exact sequence for all exact sequences $0 \rightarrow A \rightarrow B \rightarrow C \rightarrow 0$.
Let $\mathbf{k}$ be a commutative ring and let $R$ be a k-algebra. Let $\mathcal{A}$ be the category of left $R$-modules and let $\mathcal{B}$ be the category of $\mathbf{k}$-modules. Given a left $R$-module $M$, we have a functor $\operatorname{Hom}_{R}(M,-)$ given by $N \mapsto \operatorname{Hom}_{R}(M, N)$. Given a right $R$-module $M$, we also have a functor $M \otimes_{R}-$ given by $N \mapsto M \otimes_{R} N$.

Proposition A.1.1. - The functor $\operatorname{Hom}_{R}(M,-)$ is left-exact.

- The functor $M \otimes_{R}$ - is right-exact.

Finally, we state the "four lemmas" (for simplicity, just when we're dealing with the category of $R$-modules).

[^4]Lemma A.1.2 (Four lemmas). Consider the following diagram of $R$-modules:


Assume that rows are both exact.
(a) If $\beta$ and $\delta$ surjective and $\varepsilon$ is injective, then $\gamma$ is surjective.
(b) If $\beta$ and $\delta$ injective and $\alpha$ is surjective, then $\gamma$ is injective.
A.2. Derived functors. Let $F: \mathcal{A} \rightarrow \mathcal{B}$ be a left-exact functor. When $\mathcal{A}$ has "enough injectives" (satisfied if $\mathcal{A}$ is the category of left $R$-modules for some ring $R$ ) then there is a sequence of functors $\mathrm{R}^{i} F$ for $i \geq 0$, called the right derived functors of $F$ such that

- $\mathrm{R}^{i} F(M)=0$ whenever $M$ is an injective object,
- $\mathrm{R}^{0} F=F$,
- For every short exact sequence $0 \rightarrow A \rightarrow B \rightarrow C \rightarrow 0$ in $\mathcal{A}$, we get a long exact sequence

$$
0 \rightarrow F(A) \rightarrow F(B) \rightarrow F(C) \rightarrow \mathrm{R}^{1} F(A) \rightarrow \mathrm{R}^{1} F(B) \rightarrow \mathrm{R}^{1} F(C) \rightarrow \mathrm{R}^{2} F(A) \rightarrow \cdots .
$$

When $F=\operatorname{Hom}_{R}(M,-)$, we use the notation $\mathrm{R}^{i} F=\operatorname{Ext}_{R}^{i}(M,-)$ and $\mathrm{R}^{i} F(N)=$ $\operatorname{Ext}_{R}^{i}(M, N)$. These functors commute with taking finite direct sums, i.e., $\operatorname{Ext}_{R}^{i}\left(M, N \oplus N^{\prime}\right)=$ $\operatorname{Ext}_{R}^{i}(M, N) \oplus \operatorname{Ext}_{R}^{i}\left(M, N^{\prime}\right)$.

There is also a dual situation. Let $G: \mathcal{A} \rightarrow \mathcal{B}$ be a right-exact functor. When $\mathcal{A}$ has "enough projectives" (satisfied if $\mathcal{A}$ is the category of left $R$-modules for some ring $R$ ) then there is a sequence of functors $\mathrm{L}_{i} G$ for $i \geq 0$, called the left derived functors of $G$ such that

- $\mathrm{L}_{i} G(M)=0$ whenever $M$ is a projective object,
- $\mathrm{L}_{0} G=G$,
- For every short exact sequence $0 \rightarrow A \rightarrow B \rightarrow C \rightarrow 0$ in $\mathcal{A}$, we get a long exact sequence

$$
\cdots \rightarrow \mathrm{L}_{2} G(C) \rightarrow \mathrm{L}_{1} G(A) \rightarrow \mathrm{L}_{1} G(B) \rightarrow \mathrm{L}_{1} G(C) \rightarrow G(A) \rightarrow G(B) \rightarrow G(C) \rightarrow 0
$$

When $G=M \otimes_{R}-$, we use the notation $\mathrm{L}_{i} G=\operatorname{Tor}_{i}^{R}(M,-)$ and $\mathrm{L}_{i} G(N)=\operatorname{Tor}_{i}^{R}(M, N)$. These functors commute with arbitrary direct sums, i.e., $\operatorname{Tor}_{i}^{R}\left(M, \bigoplus_{\alpha} N_{\alpha}\right)=\bigoplus_{\alpha} \operatorname{Tor}_{i}^{R}\left(M, N_{\alpha}\right)$.
A.3. (Co)homology of groups. Let $\Gamma$ be a group. As discussed before, representations of $\Gamma$ over $\mathbf{k}$ are the same as left $\mathbf{k}[\Gamma]$-modules. We have the trivial module $\mathbf{k}$ with all elements of $\Gamma$ acting as the identity, which is both a left and right $\mathbf{k}[\Gamma]$-module. Given a representation $N$, we define group homology and cohomology with coefficients in $N$ by

$$
\begin{aligned}
& \mathrm{H}_{i}(\Gamma ; N)=\operatorname{Tor}_{i}^{\mathbf{k}[\Gamma]}(\mathbf{k}, N), \\
& \mathrm{H}^{i}(\Gamma ; N)=\operatorname{Ext}_{\mathbf{k}[\Gamma]}^{i}(\mathbf{k}, N) .
\end{aligned}
$$

Consider the special case $i=0$. Then

$$
\mathrm{H}_{0}(\Gamma ; N)=\mathbf{k} \otimes_{\mathbf{k}[\Gamma]} N \cong N_{\Gamma}=N /\{n-\gamma n \mid n \in N, \gamma \in \Gamma\}
$$

is the coinvariants of $N$ with respect to $\Gamma$, and

$$
\mathrm{H}^{0}(\Gamma ; N)=\operatorname{Hom}_{\mathbf{k}[\Gamma]}(\mathbf{k}, N)=N^{\Gamma}=\{n \in N \mid \gamma n=n \text { for all } \gamma \in \Gamma\}
$$

is the invariants. If $N=\mathbf{Z}$ with the trivial action, then we usually omit it from the notation. By what we said above, group homology commutes with taking arbitrary direct sums while group cohomology commutes with taking finite direct sums.

If $\mathbf{k}[\Gamma]$ is semisimple, then $N$ is both injective and projective for any representation $N$, so $\mathrm{H}_{i}(\Gamma ; N)=0$ and $\mathrm{H}^{i}(\Gamma ; N)=0$ for $i>0$. So this is interesting in the non-semisimple case, e.g., $\Gamma$ is finite and $\mathbf{k}$ is a field of positive characteristic $p$ where $p$ divides the order of $\Gamma$.

As mentioned, $\mathrm{H}^{i}(\Gamma ;-)$ and $\mathrm{H}_{i}(\Gamma ;-)$ are functors. Suppose we have a group homomorphism $f: \Gamma^{\prime} \rightarrow \Gamma$. If $M$ is a $\Gamma$-module, then it is also a $\Gamma^{\prime}$-module by the action $\gamma^{\prime} m=f\left(\gamma^{\prime}\right) m$ and we have induced maps

$$
\mathrm{H}_{i}\left(\Gamma^{\prime} ; M\right) \rightarrow \mathrm{H}_{i}(\Gamma ; M), \quad \mathrm{H}^{i}(\Gamma ; M) \rightarrow \mathrm{H}^{i}\left(\Gamma^{\prime} ; M\right),
$$

which are compatible with composition of group homomorphisms (but note that cohomology swaps the order of the groups). Given a $\Gamma$-equivariant map $\varphi: M \rightarrow N$, we also have a map $\mathrm{H}_{i}(\Gamma ; M) \rightarrow \mathrm{H}_{i}(\Gamma ; N)$, so we can compose with the above to get

$$
\mathrm{H}_{i}\left(\Gamma^{\prime} ; M\right) \rightarrow \mathrm{H}_{i}(\Gamma ; N) .
$$

In particular, this comes from a pair of maps $f: \Gamma^{\prime} \rightarrow \Gamma$ and $\varphi: M \rightarrow N$ such that $\varphi\left(\gamma^{\prime} m\right)=$ $f\left(\gamma^{\prime}\right) \varphi(m)$ for all $m \in M$ and $\gamma^{\prime} \in \Gamma^{\prime}$. We can do something similar with cohomology.

In the case when the action of $\Gamma=\Gamma^{\prime}$ and $f$ is conjugation by an element, and $M=\mathbf{k}$ is trivial, the induced map is the identity.

Let $\Omega \subset \Gamma$ be a subgroup. Given an $\Omega$-module $M$, we have the induction $\operatorname{Ind}_{\Omega}^{\Gamma} M=$ $\mathbf{k}[\Gamma] \otimes_{\mathbf{k}[\Omega]} M$ as well as the coinduction $\operatorname{Coind}_{\Omega}^{\Gamma} M=\operatorname{Hom}_{\mathbf{k}[\Omega]}(\mathbf{k}[\Gamma], M)$. If $\Omega$ is a finite-index subgroup in $\Gamma$, then these two are isomorphic.

Lemma A.3.1 (Shapiro's lemma). With the notation above, we have

$$
\mathrm{H}_{i}\left(\Gamma ; \operatorname{Ind}_{\Omega}^{\Gamma} M\right)=\mathrm{H}_{i}(\Omega ; M), \quad \mathrm{H}^{i}\left(\Gamma ; \operatorname{Coind}_{\Omega}^{\Gamma} M\right)=\mathrm{H}^{i}(\Omega ; M)
$$

A.4. Spectral sequences. While spectral sequences are rather involved structures, we will only emphasize enough to know how they can be applied to our examples of interest. Let $\mathcal{A}$ be an abelian category.

A homology spectral sequence (starting with $\mathrm{E}^{a}$ ) is the following data:

- An object $\mathrm{E}_{p, q}^{r}$ of $\mathcal{A}$ for every integer $r \geq a$ and pair of integers $p, q$.
- Morphisms $d_{p, q}^{r}: \mathrm{E}_{p, q}^{r} \rightarrow \mathrm{E}_{p-r, q+r-1}^{r}$ which are differentials in the sense that $\left(d^{r}\right)^{2}=0$.
- Isomorphisms $\mathrm{E}_{p, q}^{r+1} \cong \operatorname{ker}\left(d_{p, q}^{r}\right) / \operatorname{image}\left(d_{p+r, q-r+1}^{r}\right)$.

The degree of an object $\mathrm{E}_{p, q}^{r}$ is the quantity $p+q$, note that $d^{r}$ always lowers degree by 1 . The spectral sequence is bounded if for each $n$, there are only finitely many terms of degree $n$ in $\mathrm{E}^{a}$.

The sequence of objects with fixed superscript $r$ is the " $r$ th page" of the spectral sequence, and the last item says that taking homology with respect to the differentials $d^{r}$ gives the next page. The process of taking homology is called "turning the page". In our cases of interest, the spectral sequences will be "first quadrant", meaning that $\mathrm{E}_{p, q}^{r}=0$ if $p<0$ or $q<0$.

For given $p, q$, if we have $\mathrm{E}_{p, q}^{r}=\mathrm{E}_{p, q}^{r+1}$ for $r \gg 0$ (which happens in the bounded case), then we write $\mathrm{E}_{p, q}^{\infty}$ for this limiting value. A bounded spectral sequence converges to a collection $\left\{\mathrm{H}_{n}\right\}_{n}$ of objects of $\mathcal{A}$ if each $\mathrm{H}_{n}$ has a finite filtration

$$
0=F_{s} \mathrm{H}_{n} \subseteq F_{s+1} \mathrm{H}_{n} \subseteq \cdots \subseteq F_{t-1} \mathrm{H}_{n} \subseteq F_{t} \mathrm{H}_{n}=\mathrm{H}_{n}
$$

such that

$$
\mathrm{E}_{p, q}^{\infty} \cong F_{p} \mathrm{H}_{p+q} / F_{p-1} \mathrm{H}_{p+q} .
$$

The notation to indicate a convergent spectral sequence is

$$
\mathrm{E}_{p, q}^{a} \Longrightarrow \mathrm{H}_{p+q} .
$$

A cohomology spectral sequence is similar except we swap superscripts and subscripts and the indices increase instead of decrease and vice versa.

Morally, the purpose of a spectral sequence is to compute the desired quantity $\mathrm{H}_{n}$ from the initial data $\mathrm{E}_{p, q}^{a}$. In practice this is often difficult to carry out, but can be used to gain approximate information. For example, if $\mathcal{A}$ is the category of vector spaces, then in a convergent spectral sequence, we always have an inequality

$$
\sum_{p+q=n} \operatorname{dim}_{\mathbf{k}} \mathrm{E}_{p, q}^{a} \geq \operatorname{dim}_{\mathbf{k}} \mathrm{H}_{n}
$$

However, there are other interesting ways to use the existence of a spectral sequence to deduce information. For example, if $\mathrm{H}_{n}=0$, then it forces maps between various $\mathrm{E}_{p, q}^{r}$ to be isomorphisms.

A general and powerful spectral sequence is the Grothendieck spectral sequence. Suppose we have abelian categories $\mathcal{A}, \mathcal{B}, \mathcal{C}$ and right-exact functors $F: \mathcal{A} \rightarrow \mathcal{B}$ and $G: \mathcal{B} \rightarrow \mathcal{C}$. Then $G \circ F$ is also a right-exact functor. Suppose that $\mathcal{A}, \mathcal{B}$ have enough projectives (automatic if they are categories of modules) and suppose that $\left(\mathrm{L}_{i} G\right)(F(M))=0$ whenever $M$ is projective. Then we have a spectral sequence for any object $M$ :

$$
\mathrm{E}_{p, q}^{2}=\left(\mathrm{L}_{p} G\right)\left(\mathrm{L}_{q} F\right)(M) \Longrightarrow\left(\mathrm{L}_{p+q}(G F)\right)(M)
$$

## References

[Br] Kenneth S. Brown, Cohomology of groups, Corrected reprint of the 1982 original, Graduate Texts in Math. 87, Springer-Verlag, New York, 1994.
[Ch] Thomas Church, Homological stability for configuration spaces of manifolds, Invent. Math. 188 (2012), no. 2, 465-504, arXiv:1103.2441v3.
[CEF] Thomas Church, Jordan Ellenberg, Benson Farb, FI-modules and stability for representations of symmetric groups, Duke Math. J. 164 (2015), no. 9, 1833-1910, arXiv:1204.4533v3.
[CEFN] Thomas Church, Jordan S. Ellenberg, Benson Farb, Rohit Nagpal, FI-modules over Noetherian rings, Geom. Topol. 18 (2014), no. 5, 2951-2984, arXiv:1210.1854v2.
[DW1] Harm Derksen, Jerzy Weyman, Semi-invariants of quivers and saturation for Littlewood-Richardson coefficients, J. Amer. Math. Soc. 13 (2000), no. 3, 467-479.
[DW2] Harm Derksen, Jerzy Weyman, On the Littlewood-Richardson polynomials, J. Algebra 255 (2002), no. 2, 247-257.
[Dr1] Jan Draisma, Topological Noetherianity of polynomial functors, arXiv:1705.01419v1.
[Dr2] Jan Draisma, Noetherianity up to symmetry, Combinatorial algebraic geometry, Lecture Notes in Math. 2108, Springer, 2014, arXiv:1310.1705v2.
[DE] Jan Draisma, Rob H. Eggermont, Plücker varieties and higher secants of Sato's Grassmannian, J. Reine Angew. Math., to appear, arXiv:1402.1667v2.
[DrK] Jan Draisma, Jochen Kuttler, Bounded-rank tensors are defined in bounded degree, Duke Math. J. 163 (2014), no. 1, 35-63, arXiv:1103.5336v2.
[F] William Fulton, Young Tableaux. With applications to representation theory and geometry, London Math. Society Student Texts 35, Cambridge University Press, 1997.
[Hi] Graham Higman, Ordering by divisibility in abstract algebras, Proc. London Math. Soc. (3) 2 (1952), 326-336.
[HS] Christopher J. Hillar, Seth Sullivant, Finite Gröbner bases in infinite dimensional polynomial rings and applications, Adv. Math. 229 (2012), 1-25.
[Ho] G. Hochschild, Relative homological algebra, Trans. Amer. Math. Soc. 82 (1956), 246-269.
[KM] Michael Kapovich, John J. Millson, A path model for geodesics in Euclidean buildings and its applications to representation theory, Groups Geom. Dyn. 2 (2008), no. 3, 405-480, arXiv:math/0411182v2.
[KT1] Allen Knutson, Terence Tao, The honeycomb model of $G L_{n}(\mathbb{C})$ tensor products I: Proof of the saturation conjecture, J. Amer. Math. Soc. 12 (1999), no. 4, 1055-1090, arXiv:math/9807160v4.
[La] J. M. Landsberg, Tensors: Geometry and Applications, Graduate Studies in Math 128 (2012), AMS.
[Mac] I. G. Macdonald, Symmetric Functions and Hall Polynomials, second edition, Oxford Mathematical Monographs, Oxford, 1995.
[NW] C. St. J. A. Nash-Williams, On well-quasi-ordering finite trees, Proc. Cambridge Philos. Soc. 59 (1963), 833-835.
[JR] Rita Jimenez Rolland, Representation stability for the cohomology of the moduli space $\mathscr{M}_{g}^{n}$, Algebr. Geom. Topol. 11 (2011), no. 5, 3011-3041, arXiv:1106.0947v2.
[Ke] Moritz C. Kerz, The complex of words and Nakaoka stability, Homology Homotopy Appl. 7 (2005), no. 1, 77-85.
[Mc] John McCleary, A user's guide to spectral sequences, Second edition, Cambridge Studies in Advanced Mathematics 58, Cambridge University Press, Cambridge, 2001.
[Na] Rohit Nagpal, FI-modules and the cohomology of modular representations of symmetric groups, arXiv:1505.04294v1.
[NS] Rohit Nagpal, Andrew Snowden, Periodicity in the cohomology of symmetric groups via divided powers, arXiv:1705.10028v1.
[PY] Nicholas Proudfoot, Ben Young, Configuration spaces, FS ${ }^{\text {op }}$-modules, and Kazhdan-Lusztig polynomials of braid matroids, arXiv:1704.04510v2.
[Pu] Andrew Putman, Stability in the homology of congruence subgroups, Invent. Math. 202 (2015), no. 3, 987-1027, arXiv:1201.4876v5.
[PS] Andrew Putman, Steven V Sam, Representation stability and finite linear groups, Duke Math. J., to appear, arXiv:1408.3694v3.
[Ra] Claudiu Raicu, Secant varieties of Segre-Veronese varieties, Algebra Number Theory 6 (2012), no. 8, 1817-1868, arXiv:1011.5867v2.
[SS1] Steven V Sam, Andrew Snowden, Gröbner methods for representations of combinatorial categories, J. Amer. Math. Soc. 30 (2017), no. 1, 159-203, arXiv:1409.1670v3.
[SS2] Steven V Sam, Andrew Snowden, Proof of Stembridge's conjecture on stability of Kronecker coefficients, J. Algebraic Combin. 43 (2016), no. 1, 1-10, arXiv:1501.00333v2.
[Sa] Steven V Sam, Ideals of bounded rank symmetric tensors are generated in bounded degree, Invent. Math. 207 (2017), no. 1, 1-21, arXiv:1510.04904v2.
[Sta] Richard P. Stanley, Enumerative Combinatorics, Vol. 2, with a foreword by Gian-Carlo Rota and appendix 1 by Sergey Fomin, Cambridge Studies in Advanced Mathematics 62, Cambridge University Press, Cambridge, 1999.
[Ste] John R. Stembridge, Generalized stability of Kronecker coefficients, available from http://www. math.lsa.umich.edu/~jrs/.
[To] Burt Totaro, Configuration spaces of algebraic varieties, Topology 35 (1996), no. 4, 1057-1067.
[Wei] Charles A. Weibel, An Introduction to Homological Algebra, Cambridge University Press, Cambridge, 1994.
[Wey] Jerzy Weyman, Cohomology of Vector Bundles and Syzygies, Cambridge University Press, Cambridge, 2003.

Department of Mathematics, University of Wisconsin, Madison
URL: http://math.wisc.edu/~svs/
E-mail address: svs@math.wisc.edu


[^0]:    ${ }^{1}$ In the English convention, row $i$ sits above row $i+1$, in the French convention, it is reversed. There is also the Russian convention, which is obtained from the French convention by rotating by 45 degrees counter-clockwise.

[^1]:    ${ }^{2}$ Every element is uniquely of the form $\sum_{\lambda} c_{\lambda} e_{\lambda}$; since $f$ is a ring homomorphism, it sends this to $\sum_{\lambda} c_{\lambda} f\left(e_{\lambda_{1}}\right) f\left(e_{\lambda_{2}}\right) \cdots f\left(e_{\lambda_{\ell(\lambda)}}\right)$.

[^2]:    ${ }^{3}$ We assume $\mathcal{C}$ is essentially small, that is, the isomorphism classes of objects in $\mathcal{C}$ form a set. In almost all cases, this set is the natural numbers or some slight variant so it won't be an issue for us.

[^3]:    ${ }^{4} \mathbf{k}\left[x_{1}, \ldots, x_{n}\right]$ is noetherian by the Hilbert basis theorem, so we can always take $r$ finite

[^4]:    ${ }^{5}$ This uses the existence of the Taylor complex of a monomial ideal, a semicontinuity argument with Gröbner bases, and the fact that the degree 2 equations for the Segre embedding form a Gröbner basis with respect to some term ordering of the variables.

