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# Using character varieties: Presentations, invariants, divisibility and determinants 

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# Using character varieties: presentations, invariants, divisibility and determinants 

## BY

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# Submitted to the University of New Hampshire in partial fulfillment of the requirements for the degree of 

Doctor of Philosophy
in
Mathematics

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## UMİ

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This dissertation has been examined and approved.


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This thesis is dedicated to my grandmother, Margaret Hall, and to the memory of my grandmother, Ethel Smith.

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# ABSTRACT <br> Using character varieties: presentations, invariants, divisibility and determinants <br> by 

Jeffrey Allan Hall<br>University of New Hampshire, December, 1999

If $G$ is a finitely generated group, then the set of all characters from $G$ into a linear algebraic group is a useful (but not complete) invariant of $G$. In this thesis, we present some new methods for computing with the variety of $\mathrm{SL}_{2} \mathbb{C}$-characters of a finitely presented group. We review the theory of Fricke characters, and introduce a notion of presentation simplicity which uses these results. With this definition, we give a set of GAP routines which facilitate the simplification of group presentations. We provide an explicit canonical basis for an invariant ring associated with a symmetrically presented group's character variety. Then, turning to the divisibility properties of trace polynomials, we examine a sequence of polynomials $\boldsymbol{r}_{n}(a)$ governing the weak divisibility of a family of shifted linear recurrence sequences. We prove a discriminant/determinant identity about certain factors of $r_{n}(a)$ in an intriguing manner. Finally, we indicate how ordinary generating functions may be used to discover linear factors of sequences of discriminants.

Other novelties include an unusual binomial identity, which we use to prove a well-known formula for traces; the use of a generating function to find the inverse of a map $x^{n} \mapsto f_{n}(x)$; and a brief exploration of the relationship between finding the
determinants of a parametrized family of matrices and the Smith Normal Forms of the sequence.

| $\chi$ p | acter afforded by the representation $\rho$ |
| :---: | :---: |
| $\mathbb{C}(\mathbb{Q}, \mathbb{Z})$ | the complex numbers (rationals, integers) |
| $\mathbb{N}$ | the natural numbers $\mathbb{N}=\{1,2,3 . \cdots\}$ |
| $\mathrm{F}_{\mathrm{n}}$ | the free group on $n$ letters |
| X(G) | the character variety of a group G |
| R(G) | the coordinate ring of the variety $\mathrm{X}(\mathrm{G})$ |
| $\mathrm{L}=(\mathrm{P}) \mathrm{SL}_{2} \mathbb{C}$ | the (projective) special linear group |
| $\mathrm{H}<\mathrm{G}, \mathrm{H} \triangleleft \mathrm{G}$ | $H$ is a subgroup of $G, H$ is a normal subgroup of $G$ |
| [G:H] | the index of $\mathrm{H}<\mathrm{G}$ in G |
| $[\alpha, \beta]$ | for elements $\alpha \beta$ of a poset, the interval between $\alpha$ and $\beta$, including $\alpha$ and $\beta$ |
| $<\mathrm{S}>$ | the group or algebra generated by the set $S$, or freely freely generated by $S$ |
| [a.b] | $a^{-1} b^{-1} a b$ |
| $\mathrm{G}^{\prime}$ | the first derived subgroup of $\mathrm{G}, \mathrm{G}^{\prime}=<[\mathrm{G}, \mathrm{G}\}>$ |
| $A^{\text {B }}$ | $B^{-1} A B$ |
| $\mathrm{I}<\mathrm{R}, \mathrm{I} \triangleleft \mathrm{R}$ | $I$ is a subring of $R, I$ is an ideal of $R$ |
| (S) | the ideal generated by a set $S$ |
| $<\mathrm{S} \mid \mathrm{R}>$ | the group freely presented by generators $S$, with relations $R$ |
| $\mathbb{V}(\mathrm{I})$ | the variety which is defined by an ideal I |
| $\sqrt{\text { I }}$ | the radical of an ideal I |
| $\mathbb{C}[X]^{\mathbf{G}}$ | the algebra of invariants of matrix group $G$ acting on the vector space $\mathbb{C}^{\|X\|}$ |
| $x_{i}^{M}$ | action of matrix $M$ on variable $x_{i}: x_{i}^{\left[a_{i j}\right]}=\sum_{j}{ }^{2}{ }_{j i} x_{j}$ |
| $\prec$ | an admissible term order |
| $\mathrm{in}_{\prec} \mathrm{S}$ | the set of initial terms with respect to $\prec$ in a set $S$ |
| LTく ${ }_{\text {¢ }}$ | the monic leading term of $f$ with respect to $\prec$ |
| $L^{\text {L }}$ ¢ $f$ | the leading coefficient of $f$ with respect to $\prec$ |
| $\sigma_{n}$ | the $n$-th elementary symmetric function |
| * | the Reynolds operator: $* f=\frac{1}{n!} \sum_{g \in G} \mathrm{fg}^{\boldsymbol{g}}$ |
| $S_{n}$ | the permutation group on $n$ letters $\{1,2, \ldots, n\}$. We use the same notation for the group of $n \times n$ permutation matrices |
| $S_{r}^{(m)}$ | the group $S_{n}$, acting on m-element sets |
| distp, q | the distance between $p$ and $q$ |
| $S_{n}^{\{m\}}$ | the group $S_{n}$, acting on sets with $m$ or fewer elements |
| $a \vdash n$ | $a$ is a partition of integer $n$ (e.g., $(3,2,2) \vdash 7$ ) |
| $r^{2}$ | an edge in a graph (an edge between 1 and 2) |
| $\therefore 123$ | a hyperedge in a hypergraph (a hyperedge on vertices $1,2,3$ ) |

$\lfloor x\rfloor \quad$ the floor of $x$ (which is the largest integer smaller than $x$ )
$\operatorname{Res}(f, g), \Delta(f)$ the resultant of polynomials $f$ and $g$; the discriminant of $f$
$B_{t}(z) \quad$ the generalized binomial series with parameter $t$ (see [79, chapter 5], where $\mathrm{B}_{\mathrm{t}}(z)$ was introduced)
$P_{n}(j) \quad$ the number of $j$-tableaux shape with $n$ or fewer rows the "Whittemore variety" of a set of words $\diamond$ the minor of matrix $M$ obtained by deleting the $i-$ th row and the $j-$ th column the end of a proof; q.e.d.

## Chapter 1

# An introduction to the ring of Fricke 

## characters

Boone's construction of a group without solvable word problem in the 1950's was a development which would have delighted Bishop George Berkeley: a natural question about an algebraically important construction turned out to be algorithmically undecidable. There were two natural ways to explore this new territory: to construct new groups and semigroups with unsolvable word problem (respectively conjugacy problem, etc.), and to show that interesting groups and families of groups had solvable word problem (respectively conjugacy problem, etc.) The first approach has yielded a host of interesting (and often discouraging) undecidability results, the second has at least helped delineate the undecidable from the decidable in combinatorial group theory. Neither approach interests us here directly. Instead, our goal is to introduce and examine some algorithms which exploit properties of finitely generated groups, and of objects which are related to the study of representations of finitely presented groups, which are semi-decidable (i.e., if the property is true, then there is a procedure to verify it, whose running time may not be bounded.) These provide us with properties of elements in groups, weaker than the conjugacy problem, which are decidable for
all finitely-presented groups.

In this thesis, a representation of a group $G$ into a complex linear algebraic group $L$ is a homomorphism $\rho: G \rightarrow L$. If $G$ is a topological group, then we require that $\rho$ be continuous. The character of $G$ afforded by $\rho$ is the function $\chi_{\rho}: G \rightarrow \mathbb{C}$ defined by $\chi_{\rho}(\mathrm{g})=\operatorname{tr} \rho(\mathrm{g})$. If L is one-dimensional, then we call the character afforded by $\rho$ linear. Our focus will be largely on the representations of finitely presented groups $G$ into $L=\mathrm{SL}_{2} \mathbb{C}, \mathrm{PSL}_{2} \mathbb{C}$, or $\mathbb{C}$. In this chapter, all representations and characters will be into $\mathrm{SL}_{2} \mathbb{C}$.

G is said to be a finitely generated group if it is the homomorphic image of a finiterank free group under some homomorphism $\xi: F_{n} \rightarrow G$. (If $G$ is finitely generated, then there are in general many such homomorphisms.) Fix $\xi: F_{n} \rightarrow G$, a surjective homomorphism from a finite-rank free group to $G$. Now, any representation $\rho: G \rightarrow L$ determines a representation $\xi \circ \rho$ of $F_{n}$, and since $F_{n}$ is free, this representation is defined by the images of the free generators of $F_{n}$. Thus, the representation $\rho$ is uniquely determined by the images of a set of generators of G. The condition that $\rho$
is a homomorphism is the same as the condition that the diagram

$\rho$
commutes for any group $\mathcal{G}^{\prime}$ of rank $n$ and any choice of $\tau$; thus, the statement that a set of $n(d \times d)$-matrices are the images $\rho\left(g_{i}\right)$ of a finite set $\left\{g_{i}\right\}$ of generators of $G$ is may be restated as the statement that a set of polynomial relations on the $d^{2} n$ coordinates of the matrices is satisfied. By the Hilbert basis theorem, we may take this set of polynomials to be finite. The variety defined by these polynomials is called the representation variety of $G$. (Here, and in the future, we refer to the zero set of any set of polynomials over a field to be a "variety." We will not require that the set be irreducible: this is typical in the literature. Some authors might use the term "algebraic set.") We denote by $X(G)$ the quotient of the representation variety, obtained by identifying representations with equal traces. We embed $X(G)$ into the image of the representation variety under the trace map $\operatorname{tr}: \mathrm{SL}_{2} \mathbb{C} \rightarrow \mathbb{C}$. This manifold is actually a variety itself, called the character variety of G. It is an important invariant of the finitely generated group G.
$X(G)$, our main object of study, has a lesser role in the study of finitely presented
groups than its distant cousin, the character table, plays in the study of finite groups. A finite group's irreducible characters have bourded degree; and, by Cayley's theorem, every finite group of order $m$ has a faithful representation of degree $m$. Thus, knowledge of the representations up to degree $m$ completely determines the structure of the finite group.

But finitely presented groups do not, in general, have solvable word problem. On the other hand, given any set of matrices, and any finite word in this set, we may decide whether the product is the identity matrix by multiplying them together. In other words, the word problem for finite-dimensional matrix groups is algorithmically solvable [93]. The correctness of Buchberger's algorithm now implies:

Proposition 1.1 The following questions are algorithmically undecidable:
a) What is the smallest degree of a finite-dimensional faithful representation of $G$ ?
b) What is the smallest index [G:ker $\rho$ ] among finite-dimensional representations $\rho$ ?
c) Does G have a finite-dimensional faithful representation?
(For an explicit proof, see [93].) The finite-dimensional representation theory of a finitely generated group $G$ is a compromise between our desire to know the structure of a finitely generated group and the reality of the word problem. Despite the unsolvability of the word problem, we may inquire into the behavior of homomorphic
images of $G$ in algebraic groups. In fact, although $G$ may not have solvable word problem, its representations necessarily do.

We review the classical theory of $X(G)$.

Theorem 1.2 Let $\mathrm{G}=<\mathrm{g}_{1}, \ldots, \mathrm{~g}_{\boldsymbol{n}} \mid \mathrm{R}>$ be a finitely presented group. Then:
a) (Vogt) Let $\chi$ be an $S L_{2} \mathbb{C}$-character of $G$. Let $g \in G$. Then $\chi(g)$ is determined by the numbers

$$
\left\{t_{i}=\chi g_{i}\right\} \cup\left\{t_{i j}=\chi g_{i} g_{j} \mid 1 \leq i<j \leq n\right\} \cup \cdots \cup\left\{t_{123 \ldots n}=\chi g_{1} g_{2} \ldots g_{n}\right\}
$$

Indeed, more is true. The values of a character on the singletons, pairs, and triples in the above list suffice to determine the character on all of G . Vogt's theorem lets us index the characters with points

$$
\left(t_{1}, \ldots, t_{n}, t_{12}, \ldots, t_{(n-1) n}, t_{123}, \ldots, t_{(n-2)(n-1) n}\right)
$$

(A representation, of course, is determined by its value on generators. Part (a) is a finiteness theorem: characters are uniquely determined by their values on generators, and on ordered products of two and three generators.)
b) (Horowitz) The set of all possible characters is a variety, which is an invariant of the group $G$; i.e., any relation between elements of $X=\left\{t_{i}, t_{i j}, t_{i j k}\right\}$ is a polynomial. For example, if $n=2$, then $\mathrm{X}(\mathrm{G})$ is all of $\mathbb{C}^{3}$.

Proof: (a) We remind the reader that we are working with $2 \times 2$ matrices only. Each of the following identities is a trivial calculation. (1.2 is also a disguised version of the Cayley-Hamilton theorem, but that doesn't concern us here.)

$$
\begin{align*}
\operatorname{tr}(A B C)= & \operatorname{tr} A \operatorname{tr}(B C)+\operatorname{tr} B \operatorname{tr}(A C)+\operatorname{tr} C \operatorname{tr}(A B)-\operatorname{tr} A \operatorname{tr} B \operatorname{tr} C \\
& \quad-\operatorname{tr}(A C B)  \tag{1.1}\\
\operatorname{tr}(A B)= & (\operatorname{tr} A)(\operatorname{tr} B)-\operatorname{tr} A / B  \tag{1.2}\\
\operatorname{tr}(A)= & \operatorname{tr}\left(A^{-1}\right) \\
= & \operatorname{tr}\left(A^{B}\right) \\
\operatorname{tr}(A B)= & \operatorname{tr}(B A) \tag{1.3}
\end{align*}
$$

$$
\begin{align*}
\operatorname{tr}(A B C D)= & \frac{1}{2}(\operatorname{tr} A \operatorname{tr}(B C D)+\operatorname{tr} B \operatorname{tr}(C D A)+\operatorname{tr} C \operatorname{tr}(D A B)  \tag{1.4}\\
& +\operatorname{tr} D \operatorname{tr}(C A B)+\operatorname{tr}(A B) \operatorname{tr}(C D)-\operatorname{tr}(A C) \operatorname{tr}(B D) \\
& +\operatorname{tr}(A D) \operatorname{tr}(B C)-\operatorname{tr} A \operatorname{tr} B \operatorname{tr} C D-\operatorname{tr} C \operatorname{tr} D \operatorname{tr} A B \\
& -\operatorname{tr} D \operatorname{tr} A \operatorname{tr} B C-\operatorname{tr} B \operatorname{tr} C \operatorname{tr} D A+\operatorname{tr} A \operatorname{tr} B \operatorname{tr} C \operatorname{tr} D)
\end{align*}
$$

(1.4 is known as "Vogt's relation," the relation 1.1 is known as "Fricke's lemma," and 1.2 is called the "fundamental trace relation.")

Given any word $w$ with more than three generators, we may use 1.4 to express
tr $w$ using shorter words. Using 1.2, we can express traces of words of length 2 or 3 in terms of traces of length 2 or 3 without using inverses of generators (e.g.,

$$
\begin{aligned}
\operatorname{tr} a b^{-1} & =(\operatorname{tr} a)\left(\operatorname{tr} b^{-1}\right)-\operatorname{tr} a b \\
& =(\operatorname{tr} a)(\operatorname{tr} b)-\operatorname{tr} a b \\
\operatorname{tr} a b c^{-1} & =(\operatorname{tr} a b)(\operatorname{tr} c)-\operatorname{tr} a b c
\end{aligned}
$$

etc.) By conjugation, traces of words that are mis-ordered cyclically can be re-ordered (e.g., $\operatorname{tr} b c a=\operatorname{tr} a b c$.) Finally, 1.1 lets us express traces of length 3 words which are mis-ordered non-cyclically.
(b) See [57] or [33].

We will now review in some detail the computational mechanisms with which one may manipulate trace polynomials. First, however, we introduce a somewhat unusual binomial identity.

Theorem 1.3 For any positive integers $\mathrm{o}, \mathrm{l}$,

$$
\begin{equation*}
\sum_{m}\binom{20}{2 m}\binom{m}{l}=2^{20-2 l} \frac{o}{l!} \prod_{k=1}^{l-1}(2 o-l-k) \tag{1.5}
\end{equation*}
$$

Proof: We proceed via the Wilf-Zeilberger method. [99] The claim of the theorem is readily verified for $o \in\{1,2,3\}$ and $l \in\{1,2,3,4,5,6\}$. We exhibit a recurrence
relation which is satisfied by the sequence

$$
\begin{equation*}
f_{o}=\frac{\binom{20}{2 m}\binom{m}{l}}{2^{20-2 l o} \prod_{k=1}^{1 . .1-1}(20-l-k)} \tag{1.6}
\end{equation*}
$$

namely the first-order recurrence

$$
\begin{aligned}
& 2(2 l-2 o-1)(l-o-1)(k-2 o+1)(k+20) S_{o} \\
& -(l-2 o)(l-2 o-1)(k+l-2 o-1)+(k+l-2 o-1)=1
\end{aligned}
$$

where $S_{n}$ is the backwards shift operator $S_{n} f(n)=f(n-1)$. (This recurrence was found with Zeilberger's implementation of his "creative telescoping" algorithm [99], although it is readily verified by hand.) A similar recurrence relation exists for $l$. Since we have verified the recurrence for $l=1, o=1,2,3$, by induction the theorem is true for all positive integral $m$ when $l=1$. Likewise the theorem is true for all integral l, m. $\diamond$

Creative telescoping proofs, like the one above, are straightforward but unenlightening. We note that a proof of this identity by comparing the quotients of successive derivatives of Chebyshev polynomials is possible, but quite tedious, and not terribly insightful. I do not know a purely combinatorial proof of this result.

We use the identity to prove a useful formula, a form of which seems to be wellknown among some applied mathematicians [81]:

Theorem 1.4 Let K be a field. Consider the recurrence relation in $\mathrm{K}[\mathrm{x}]$

$$
\begin{align*}
T_{n+2} & =x T_{n+1}(x)-T_{n}(x)  \tag{1.7}\\
T_{1} & =x \\
T_{0} & =2
\end{align*}
$$

Then, for $\mathrm{n} \geq 1, \mathrm{~T}_{\mathrm{n}}(\mathrm{x})$ is a degree n , monic polynomial, which is either either even or odd, and the coefficient of $z^{n-2 l}$ is

$$
\begin{equation*}
\frac{(-1)^{l}}{l!} n \prod_{k=1}^{l-1}(n-l-k) \tag{1.8}
\end{equation*}
$$

for each $0<l \leq\left\lfloor\frac{n}{2}\right\rfloor$.

Proof: Write $c_{l}$ for the coefficient of $z^{n-2 l}$. Solving the recurrence relation for $T_{n}$ by standard techniques, we see that

$$
\begin{align*}
c_{l} & =\frac{1}{2^{n-1}}(-4)^{l} \sum_{m=-\infty}^{\infty}\binom{n}{2 m}\binom{m}{l}  \tag{1.9}\\
& =\frac{1}{2^{n-1}}(-4)^{l} \sum_{m=l}^{\left\lfloor\frac{n}{2}\right\rfloor}\binom{n}{2 m}\binom{m}{l} \tag{1.10}
\end{align*}
$$

so, considering odd and even $n$ separately, we have by the identity of the previous
theorem

$$
c_{l}=\frac{(-1)^{l}}{l!} n \prod_{k=1}^{l-1}(n-l-k) .
$$

## $\diamond$

Corollary 1.5 Let t be the trace of a $2 \times 2$ matrix $\mathcal{A} \in S L_{2}$. Then the trace of $A^{n}$ is given by

$$
\begin{equation*}
T_{n}(t)=t^{n}+\sum_{l=1}^{\lfloor n / 2\rfloor} c_{1} t^{n-2 l} \tag{1.11}
\end{equation*}
$$

Proof: By the fundamental trace identity, equation 1.2 , the trace of $A^{n}$ is given by 1.7. $\diamond$

Corollary 1.6 The character variety of a finite group, or more generally of a finitely generated torsion group of bounded exponent, has zero dimension.

Proof: If $S$ is the set of all exponents of elements of the group, then any group element's $\mathrm{SL}_{2}$ trace must satisfy one of $\left\{\mathrm{T}_{n}(x) \mid n \in S\right\}$. Each of these polynomials has a finite solution set, and so each of the coordinates takes on a discrete number of values. $\diamond$

For the reader's convenience, the first few $T_{\pi}$ 's are collected in table one.
The following formula, due to Jorgensen [75], extends the usefulness of 1.9 to arbitrary words on two letters.

2
x
$x^{2}-2$
$x^{3}-3 x$
$x^{4}-4 x^{2}+2$
$x^{5}-5 x^{3}+5 x$
$x^{6}-6 x^{4}+9 x^{2}-2$
$x^{7}-7 x^{5}+14 x^{3}-7 x$
$x^{8}-8 x^{6}+20 x^{4}-16 x^{2}+2$
$x^{9}-9 x^{7}+27 x^{5}-30 x^{3}+9 x$
$x^{10}-10 x^{8}+35 x^{6}-50 x^{4}+25 x^{2}-2$
$x^{11}-11 x^{9}+44 x^{7}-77 x^{5}+55 x^{3}-11 x$
$x^{12}-12 x^{10}+54 x^{8}-112 x^{6}+105 x^{4}-36 x^{2}+2$
$x^{13}-13 x^{11}+65 x^{9}-156 x^{7}+182 x^{5}-91 x^{3}+13 x$
$x^{14}-14 x^{12}+77 x^{10}-210 x^{8}+294 x^{6}-196 x^{4}+49 x^{2}-2$
Table 1.1 Trace polynomials of powers of an element

Proposition 1.7 Let

$$
w=x^{a_{1}} y^{b_{1}} x^{a_{2}} y^{b_{2}} \cdots x^{a_{k}} y^{b_{k}}
$$

be a cyclically reduced word on two letters. Let $\chi$ be a character of $F_{2}$, the free group on $\{x, y\}$. Write $\alpha+\alpha^{-1}$ for the number $\chi x$, and $\beta+\beta^{-1}$ for the number $\chi y$. Let $p=\chi(w)$, which by Vogt's identity is a polynomial in variables $\alpha+\alpha^{-1}, \beta+\beta^{-1}$, and $\chi(x y)$. Then the degree of $z=\chi(x y)$ in $p$ is the number $k$ of $\left(a_{i}, b_{i}\right) \neq(0,0)$ and the coefficient of $z^{k}$ is

$$
\begin{equation*}
\prod_{i}\left(\frac{\alpha^{a_{i}}-\alpha^{-a_{i}}}{\alpha-\alpha^{-1}}\right) \prod_{i}\left(\frac{\beta^{b_{i}}-\beta^{-b_{i}}}{\beta-\beta^{-1}}\right) \tag{1.12}
\end{equation*}
$$

We have generally found it more convenient to use Jorgensen's formula in a different
form:

$$
\begin{aligned}
& \prod_{i}\left(\frac{\alpha^{a_{i}}-\alpha^{-a_{i}}}{\alpha-\alpha^{-1}}\right) \prod_{i}\left(\frac{\beta^{b_{i}}-\beta^{-b_{i}}}{\beta-\beta^{-1}}\right)= \\
& \prod_{i}\left(\sum_{j}\left(\alpha^{a_{i}-1}+\alpha^{1-a_{i}}\right) \prod_{i} \sum_{J}\left(\beta^{b_{i}-1}+\beta^{1-b_{i}}\right)\right. \\
&= \prod_{i=1}^{k} T_{i} T_{j} .
\end{aligned}
$$

We introduce some notation. If $w$ is a word in a free group $F_{n}$, and if the character values

$$
S=\left\{t_{w_{1}}=\operatorname{tr} w_{1}, t_{w_{2}}=\operatorname{tr} w_{2}, \ldots, t_{w_{m}}=\operatorname{tr} w_{m} \mid w_{i} \in F_{n}\right\}
$$

generate the images of each word in $F_{\eta}$ under any $\left(\mathrm{SL}_{2}\right)$ character $\chi$, then the polynomial in $\mathbb{Z}[S]$ which defines the character image of $w$ is denoted either $\operatorname{tr} w$, or $t_{w}$. (Both are standard notations.) Usually, it is convenient order the letters of $F_{n}$ once and for all, and denote the character images of $x_{1}, x_{2}, x_{3}, \ldots, x_{1} x_{2}, \ldots$ etc. by $t_{1}, t_{2}, t_{3}, \ldots, t_{12}, \ldots$ etc. We will furthermore sometimes overload this notation, by considering inverses of generators, i.e., $t_{1-1}=\operatorname{tr} x_{1}^{-1}\left(=\operatorname{tr} x_{1}\right), t_{12^{-1}}=\operatorname{tr} x_{1} x_{2^{-1}}(=$ $\left.t_{1} t_{2}-t_{12}\right)$; and other powers of generators, e.g. $t_{12}=\operatorname{tr}\left(x_{1}^{2}\right)\left(=t_{1}^{2}-2\right)$.

Considering again a representation $\rho$ of a finitely presented group $G$, we have maps:

$$
\begin{array}{ccccc} 
& & & \psi \\
\operatorname{ker} \psi & \hookrightarrow & F_{n} & \rightarrow & G \\
\perp & & \downarrow & \searrow & \downarrow \rho \\
& & & & \\
& & & & \\
\left\{0_{\mathrm{SL}_{2}}\right\} & \hookrightarrow & \mathrm{SL}_{2} \mathbb{C} & \rightarrow & \mathrm{SL}_{2} \mathbb{C} \\
\downarrow \operatorname{tr} & & \downarrow \operatorname{tr} & & \downarrow \operatorname{tr} \\
\{2\} & \subseteq & \mathbb{C} & \supseteq & \mathbb{C}
\end{array}
$$

Let $I_{n}$ be the radical ideal defining the variety $X(G)$. The ring

$$
R\left(F_{n}\right)=\mathbb{C}\left[t_{1}, t_{2}, \ldots, t_{12}, \ldots, t_{123}, \ldots\right] / I_{n}=\mathbb{C}[X] / I_{n}
$$

is called the ring of Fricke characters of the free group $F_{n}=\left\langle g_{1}, \ldots, g_{n}\right\rangle$. (More generally, the ring of Fricke characters $R(G)$ of a finitely generated group $G$ is the coordinate ring of the character variety of $G$.) Computing inside $R\left(F_{n}\right)$ can be hard for at least four reasons:

1. $R\left(F_{n}\right)$ has $n+\binom{n}{2}+\binom{n}{3}=\frac{n\left(n^{2}+5\right)}{6}$ variables. The average and worst-case running times for most algorithms in computational commutative algebra depend exquisitely on the number of variables in ring presentations; for some illustrative examples, see [116] or [17].
2. $R\left(F_{n}\right)$ has high regularity.
3. The polynomials generating $I_{n}$, the irreducible ideal which defines the ideal $X\left(F_{n}\right)$ are the Gonzalez-Montesinos relations, whose definition we recall below. They are very symmetrical. Indeed, they admit an intransitive group of symmetries of order n!. Symmetry among the generators of an ideal cause algorithms such as Buchberger's algorithm to do work which does not move it towards its termination condition (for examples, see [115].)
4. The Gonzalez-Montesinos relations have $2\binom{n}{3}+\binom{n-2}{2}+\binom{n-3}{2}$ polynomials comprised of

$$
15\binom{n}{3}+12\binom{n-2}{2}+24\binom{n}{3}+10\binom{n-3}{2}
$$

(not all distinct) monomials. This imposes a real "book-keeping" cost as $n$ grows large.

For $n<4$ these are not serious objections to computing inside $R(G)$. For $n=1,2$, $I_{n}=\{0\}$, and so two characters $\bar{f}, \bar{g}$ are equal if and only if $f=g$ in $\mathbb{C}[X]$. For $\pi=3$, the results of [68] suffice $-I_{3}$ is principal over $\mathbb{Z}$,

$$
I_{3}=\left(t_{123}^{2}-P t_{123}+Q\right)
$$

where

$$
\begin{aligned}
P= & t_{1} t_{23}+t_{2} t_{13}+t_{3} t_{12}-t_{1} t_{2} t_{3} \\
Q= & t_{1}^{2}+t_{2}^{2}+t_{3}^{2}+t_{12}^{2}+t_{13}^{2}+t_{23}^{2}+t_{12} t_{13} t_{23} \\
& -t_{1} t_{2} t_{12}-t_{1} t_{3} t_{13}-t_{2} t_{3} t_{23}-4 .
\end{aligned}
$$

$t_{123}^{2}-P t_{123}+Q$ is thus evidently a Grobner basis for $I_{3}$ with respect to any term ordering, and so we have many normal-form algorithms for $R\left(F_{3}\right)$. Likewise, since Gonzalez-Acuna and Montesinos-Amilibia have provided an explicit set of polynomials generating $I_{n},[57]$, we may in principle find a Grobner basis for $R\left(F_{n}\right)$, and thus a normal-form procedure for $R\left(I_{n}\right)$. But when $n=4$, this is already a non-trivial calculation, and for larger $n$, our objections 1-4 above seem formidable.

Our primary object of study is $R\left(F_{n}\right)$. Its defining ideal is generated by the Gonzalez-Montesinos polynomials:

$$
\begin{equation*}
t_{a b c}^{2}-P_{a b c} t_{a b c}+Q_{a b c} \quad a<b<c \tag{1.13}
\end{equation*}
$$

for

$$
\begin{aligned}
P_{a b c}= & t_{a} t_{b c}+t_{b} t_{a c}+t_{c} t_{a b}-t_{a} t_{b} t_{c} \\
Q_{a b c}= & t_{a}^{2}+t_{b}^{2}+t_{c}^{2}+t_{a b}^{2}+t_{a c}^{2}+t_{b c}^{2}+t_{a b} t_{a c} t_{b c} \\
& -t_{a} t_{b} t_{a b}-t_{a} t_{c} t_{a c}-t_{b} t_{c} t_{b c}-4,
\end{aligned}
$$

$$
\left|\begin{array}{cccc}
t_{1}^{2}-4 & 2 t_{12}-t_{1} t_{2} & 2 t_{1 a}-t_{1} t_{a} & 2 t_{1 b}-t_{1} t_{b}  \tag{1.14}\\
2 t_{12}-t_{1} t_{2} & t_{2}^{2}-4 & 2 t_{2 a}-t_{2} t_{a} & 2 t_{2 b}-t_{2} t_{b} \\
2 t_{1 a}-t_{1} t_{a} & 2 t_{2 a}-t_{2} t_{a} & t_{a}^{2}-4 & 2 t_{a b}-t_{a} t_{b} \\
2 t_{1 b}-t_{1} t_{b} & 2 t_{2 b}-t_{2} t_{b} & 2 t_{a b}-t_{a} t_{b} & t_{b}^{2}-4
\end{array}\right| \quad 3 \leq a<b \leq n
$$

$$
\left|\begin{array}{cccc}
t_{1}^{2}-4 & 2 t_{12}-t_{1} t_{2} & 2 t_{13}-t_{1} t_{3} & 2 t_{1 a}-t_{1} t_{a}  \tag{1.15}\\
2 t_{12}-t_{1} t_{2} & t_{2}^{2}-4 & 2 t_{23}-t_{2} t_{3} & 2 t_{2 a}-t_{2} t_{a} \\
2 t_{13}-t_{1} t_{3} & 2 t_{23}-t_{2} t_{3} & t_{3}^{2}-4 & 2 t_{3 a}-t_{3} t_{a} \\
2 t_{1 b}-t_{1} t_{b} & 2 t_{2 b}-t_{2} t_{b} & 2 t_{b 3}-t_{b} t_{3} & 2 t_{a b}-t_{a} t_{b}
\end{array}\right| \quad 3<a<b \leq n
$$

$$
\begin{align*}
& \left(t_{123}-t_{132}\right)\left(2 t_{a b c}-t_{a} t_{b} t_{c}-t_{a} t_{b c}-t_{b} t_{a c}-t_{c} t_{a b}\right)  \tag{1.16}\\
& -\left|\begin{array}{cccc}
t_{1} & t_{1 a} & t_{1 b} & t_{1 c} \\
t_{2} & t_{2 a} & t_{2 b} & t_{2 c} \\
t_{3} & t_{3 a} & t_{3 b} & t_{3 c} \\
2 & t_{a} & t_{b} & t_{c}
\end{array}\right| \quad a<b<c \tag{1.17}
\end{align*}
$$

We introduce some useful notation. Let $G$ be presented by $<g_{1}, \ldots, g_{n}|R\rangle$; call this presentation $P$. By [57, Section three], the character variety $X(G)$ is the
intersection of two other varieties: $X\left(F_{n}\right)$ and

$$
\mathrm{V}\left(\left(\operatorname{tr} \mathrm{R}^{\prime}\right)-2\right)=\mathrm{V}\left(\mathrm{I}_{w}\right),
$$

where $\operatorname{tr} R^{\prime}=\left\{t_{r} \mid r \in\left(\{e\} \cup\left\{g_{1}, g_{2}, \ldots, g_{n}\right\}\right) R\right\}$ for some arbitrary choice $\left\{t_{r}\right\}$ of trace polynomials in $X$, one $t_{r}$ representing each $w \in R^{\prime}$. We will call the variety $W(P)=$ $V((\operatorname{tr} R))$ a Whittemore variety of the presentation $\left\langle g_{1}, \ldots, g_{n} \mid R\right\rangle$, which we write $V_{w}(G)=W(P)$. This is a double abuse of notation: the group G's Whittemore variety is really associated with the presentation $P:\left\langle g_{1}, \ldots, g_{n} \mid R\right\rangle$, and we have a choice of trace polynomials for each $r \in R$ whenever $n \geq 3$.

Example 1.8 We will defer choosing a canonical way of generating $I_{w}$ for a general finitely presented group until the next chapter, but we give here an ad hoc example for the group

$$
<f_{1}, f_{2}, f_{3}\left|f_{1} f_{2} f_{3}^{-1}, f_{2} f_{3} f_{1}^{-1}, f_{3} f_{1} f_{2}^{-1}\right\rangle
$$

which is the finite Fibonacci group $\mathrm{F}(2,3)$. The character variety $\mathrm{X}\left(\mathrm{F}_{3}\right)$ is the variety of the principal ideal

$$
t_{123}^{2}-P_{123} t_{123}+Q_{123}
$$

where

$$
\begin{aligned}
P= & t_{1} t_{23}+t_{2} t_{13}+t_{3} t_{12}-t_{1} t_{2} t_{3} \\
Q= & t_{1}^{2}+t_{2}^{2}+t_{3}^{2}+t_{12}^{2}+t_{13}^{2}+t_{23}^{2}+t_{12} t_{13} t_{23} \\
& -t_{1} t_{2} t_{12}-t_{1} t_{2} t_{13}-t_{2} t_{3} t_{23}-4
\end{aligned}
$$

as in 1.13 above; while $\mathrm{V}_{w}(\mathrm{~F}(2,3))$ may be chosen to be the variety defined by

$$
\begin{align*}
\left(\operatorname{tr} f_{1} f_{2} f_{3}^{-1}, \operatorname{tr} f_{2} f_{3} f_{1}^{-1}, \operatorname{tr} f_{3} f_{1} f_{2}^{-1}\right) &  \tag{1.18}\\
& =\left(\underline{t_{3} t_{12}}-t_{123}, \underline{t_{1} t_{23}}-t_{123}, \underline{t_{2} t_{13}}-t_{123}\right)
\end{align*}
$$

Note that each polynomial has its degree-lexicographic leading monomial underlined. Since these terms are relatively prime, the set (1.18) forms a Grobner basis. (A Grobner basis like this one, where the leading terms are relatively prime, is called a structural Grobner basis. See [117].)

As the term $\mathrm{t}_{123}$ occurs in none of the leading terms of this graded Grobner basis, the coordinate ring of $\mathrm{V}_{w}(\mathrm{~F}(2,3))$ has codimension of at least one. A coset enumeration (using the system [40], for example) shows that $|\mathcal{F}(2,3)|=8<\infty$, so that $\operatorname{dim} X(F(2,3))=\operatorname{dim} X\left(F_{n}\right) \cap V_{w}=0$, and in so particular $X(F(2,3)) \neq V_{w}$, by corollary 1.6. $\diamond$

When we are calculating $\mathrm{V}_{\mathrm{w}}$ for a group, the nicest sorts of relators which we might encounter are perfect powers, since we have an easy way to write down $\operatorname{tr} w^{n}-2=$
$T_{n}\left(t_{w}\right)-2$ using corollary 1.5. We can also go backwards; the next theorem tells us how to write $\mathrm{x}^{\mathfrak{n}}$ in terms of the polynomials $\mathrm{T}_{\mathrm{n}}(\mathrm{x})$.

Theorem 1.9 Define $\left\langle\binom{ n}{k}\right\rangle$ by $\left\langle\binom{ n}{k}\right\rangle=\binom{n}{k}$ if $n$ and $k$ are integers; $\left\langle\binom{ n}{k}\right\rangle=0$ otherwise. Define the polynomials $\widetilde{T}_{n}(x)$ by

$$
\tilde{T}_{n}(x)=\left\{\begin{array}{cc}
1 & \text { if } n=0 \\
T_{n}(x) & \text { if } n>0
\end{array}\right.
$$

Then

$$
x^{n}=\sum_{k}\left\langle\binom{ n}{\frac{k}{2}}\right\rangle\left(\widetilde{T}_{(n-k)}(x)\right)
$$

This expression for $x^{n}$ as a linear combination of the polynomials $T_{n}(x)$ is unique.

Proof: The statement of the theorem is equivalent to the statement that the matrix illustrated in figure 1-1 is the inverse of the matrix in figure 1-2. Let $N_{1}$ be the $n \times n$ upper-left submatrix of $M$. Let $N_{2}$ be the $n \times \pi$ upper-left submatrix of of the matrix in figure 1-2. Let $\alpha$ be the $a-$ th row of $N_{2}, \beta$ the $b-$ th column of $N_{1}$. Index the entries of $\alpha$, starting at the right, with $\alpha_{0}$ the rightmost entry of the row. (Continue this sequence to the left, so that the sequence $\left\{\alpha_{i}\right\}$ runs through all of the binomial coefficients.) Likewise index the entries of $\beta$, starting from the top, with $\beta_{0}$ the top
entry of the column. The generating function of the sequence $\left\{\alpha_{i}\right\}$ is clearly

$$
A_{a}(z)=\left(1+z^{2}\right)^{a} z^{n-a-1}
$$

Likewise the generating function of the sequence $\left\{\beta_{i}\right\}$ is

$$
\mathrm{B}_{b}(z)=z^{\mathrm{b}}\left(1-z^{2}\right)\left(1+z^{2}\right)^{-b-1}
$$

(Proof: The numbers $N_{1}(x, y)$ satisfy the recurrence

$$
\begin{equation*}
N_{1}(x, y)=N_{1}(x-1, y-1)-N_{1}(x, y-2) . \tag{1.19}
\end{equation*}
$$

This is immediate from the defining relation $T_{n+2}=x T_{n+1}(x)-T_{n}(x)$. Multiply each side of (1.19) by $z^{n}$, and sum over all $n \geq 0$. We have

$$
\sum_{n} N_{1}(x, y) z^{n}=\sum_{n} N_{1}(x-1, y-1) z^{n}-\sum_{n} N_{1}(x, y-2) z^{n}
$$

$$
\mathrm{B}_{x}(z)=z \mathrm{~B}_{x-1}(z)-z^{2} \mathrm{~B}_{x}(z)
$$

$$
B_{x}(z)=\frac{z}{1+z^{2}} B_{x-1}(z)
$$

Since $B_{0}(z)=\frac{2}{1+z^{2}}-1$, we have

$$
\mathrm{B}_{x}(z)=\left(\frac{z}{1+z^{2}}\right)^{x}\left(\frac{1-z^{2}}{1+z^{2}}\right)
$$

as claimed.)
Now consider the generating function of the convolution $\alpha * \beta$, which is the product $A_{a}(z) B_{b}(z):$

$$
\left(A_{a} B_{b}\right)(z)=z^{n-1+(b-a)}\left(1+z^{2}\right)^{a-b-1}\left(1-z^{2}\right)
$$

We examine the coefficient of $z^{n-1}$ in this power series. There are three cases:
$a<b$ Then $\left(A_{a} B_{b}\right)(z)=z^{n-1} F(z)$, where $F(z)$ is analytic at $z=0$, so the coefficient of $z^{n-1}$ in $\left(A_{a} B_{b}\right)(z)$ is 0.
$\mathrm{a}=\mathrm{b}$ Then the coefficient of $z^{\mathrm{n}-1}$ in $\left(A_{a} B_{b}\right)(z)$ is obviously 1 .
$a>b$ Then $\left(A_{a} B_{b}\right)(z)$ is a polynomial. The coefficient of $z^{n-1}$ in $\left(A_{a} B_{b}\right)(z)$ is the coefficient of $z^{a-b}$ in

$$
\left(1+z^{2}\right)^{a-b-1}\left(1-z^{2}\right)
$$

This coefficient is 0 if $(a-b)$ is odd, and is

$$
\binom{a-b-1}{\frac{a-b}{2}}-\binom{a-b-1}{\frac{a-b}{2}-1}
$$

$$
M=\left(\begin{array}{cccccccc}
1 & & & & & & & \\
0 & 1 & & & & & & \\
-2 & 0 & 1 & & & & & \\
0 & -3 & 0 & 1 & & & & \\
2 & 0 & -4 & 0 & 1 & & & \\
0 & 5 & 0 & -5 & 0 & 1 & & \\
-2 & 0 & 9 & 0 & -6 & 0 & 1 & \\
0 & -7 & 0 & 14 & 0 & -7 & & \ddots
\end{array}\right)
$$

Figure 1-1 The matrix $M$

$$
M^{-1}=\left(\begin{array}{cccccccc}
1 & & & & & & & \\
0 & 1 & & & & & \\
1 & 0 & 1 & & & & & \\
0 & 3 & 0 & 1 & & & & \\
6 & 0 & 4 & 0 & 1 & & & \\
0 & 10 & 0 & 5 & 0 & 1 & & \\
20 & 0 & 15 & 0 & 6 & 0 & 1 & \\
0 & 35 & 0 & 21 & 0 & 7 & & \ddots
\end{array}\right)
$$

Figure 1-2 The inverse of matrix $M$
if ( $a-b$ ) is even. By the symmetry property of Pascal's triangle,

$$
\binom{a-b-1}{\frac{a-b}{2}}-\binom{a-b-1}{\frac{a-b}{2}-1}=0
$$

when $(a-b-1)$ is odd.

So $(\alpha * \beta)_{n-1}$ is 0 if $a \neq b$, and 1 if $a=b$. But $(\alpha * \beta)_{n-1}$ is precisely the dot product of the $a$ - th row of $N_{2}$ with the $b$ - th column of $N_{1}$, so the matrices $N_{1}$ and $N_{2}$ are inverses of each other. $\diamond$

## Chapter 2

## Simplification of group presentations:

## dimension theory aids string matching

In this brief chapter, we explore definitions of group presentation simplicity that use geometric information taken from the group presentation. Our approach is firmly experimental.

Suppose that $G=\langle Y \mid R\rangle$ and $H=\langle Z \mid S\rangle$ are two finite presentations of groups. Then it is a well-known theorem of Tietze that that G is isomorphic to H if and only if the presentation $\langle Y \mid R\rangle$ may be obtained from $\langle Z \mid S\rangle$ by a finite sequence of "Tietze transformations:"

1. Adding a new generator $g$, and a new relation $g w$, where $w$ is any word in $Z$
2. Deleting a generator $g$, and a relation $g w$, where $w$ is a word in $Z-\{g\}$, and no other relation in $S$ uses the generator $g$
3. Adding a relation that is a consequence of other relations in $S$
4. Deleting a relation that is a consequence of the other relations in $S$.
(If two presentations $\langle Y \mid R\rangle$ and $\langle Z \mid S\rangle$ of the same group are not finite, then in general we may not find a finite sequence of Tietze transformations transforming one into the
other, even if $|Y|,|Z|<\infty$. For example, $\left\langle a \mid a, a^{2}, a^{3}, \ldots\right\rangle$ is clearly a presentation of the trivial group, and no Tietze transformation can ever yield a presentation where there are not infinitely many relators, which uses only one generator. But

$$
\left\langle a, b \mid a, b, w_{2}(a, b), w_{2}(a, b), \ldots\right\rangle
$$

is such a presentation of the trivial group, so it can't be obtained from $\left\langle a \mid a, a^{2}, a^{3}, \ldots\right\rangle$ in only finitely many steps.)

The search-space of Tietze transforms of a finitely presented group is an important object of study in computational group theory. A typical application of searching through the space of Tietze transforms of a finitely presented group is the problem of simplifying a presentation. What does it mean for one presentation of a finitely presented group to be "simpler" than another? Some typical definitions are that a presentation is simpler if it has fewer generators, or fewer relations, or smaller total relator length. Given a presentation $\langle\gamma \mid R\rangle$, there is a (finite) path of Tietze transformations that transforms $<\mathrm{Y} \mid \mathrm{R}>$ into a "simplest" presentation. But, for general finitely presented groups, the isomorphism problem is unsolvable; and thus, if our definition of simplicity admits a unique simplest presentation then the problem of finding this path is algorithmically undecidable.

In practice, then, before attempting to simplify a finite presentation of group $G$ via a sequence of Tietze transformations, one chooses a binary relation $\leq$ on the set of all group presentations, defining $<Y \mid R>$ to be simpler than $<Z \mid S>$ if
$<\mathrm{Y}|\mathrm{R}>\leq<\mathrm{Z}| \mathrm{S}>$. Our search space is an oriented graph, where the vertices are finite presentations of groups isomorphic to $G$, and where edges correspond to Tietze transforms. We search through this graph, until we find a locally minimal (or acceptable minimal) element. The complexity of this graph gives the whole theory a very computational flavor. The hardest part of simplifying a presentation by searching the space of Tietze transformations is determining the search-space itself: transformations of type (2) and (4) require a common-substring search, which is much harder than merely transforming a presentation, or comparing the number of generators or relation lengths [113, Section 6.4] [64]. This is a distinctive aspect of this problem, which is often absent in other applications of combinatorial search.

In this chapter, we will explore a new definition of simplicity for a group presentation. Our definition has the advantage that it allows us to avoid common-substring searching when deciding whether to apply Tietze transformations of type (1) and (2). Briefly, we will consider one presentation $P_{1}=<x_{1}, \ldots, x_{n} \mid R_{1}>$ of a finitely presented group $G$ to be simpler than another presentation $P_{2}=\left\langle x_{1}, \ldots, x_{m} \mid R_{2}\right\rangle$, if the Gonzalez-Montesinos presentation of the character variety of $F_{n}$ "contributes" less, via presentation $P_{1}$, to the Hilbert polynomial of the ideal $I(X(G))$ than the corresponding presentation of $F_{m}$ via presentation $P_{2}$. (We will clarify what we mean by "contributes less" below.)

Let us first consider a simple example. The triangle group ( $2,3, \infty$ ) may be presented as:

$$
\Gamma=<g_{1}, g_{2}, g_{3} \mid g_{1}^{2}=\left(g_{1} g_{2}\right)^{3}=g_{3} g_{1}^{-1} g_{2}^{2}=e>
$$

Its character variety's Gonzalez-Montesinos presentation is

$$
T_{1}=\left\{x_{1}^{2}-2, x_{12}^{3}-3 x_{12}, x_{1} x_{2}, x_{23}-x_{2} x_{123}-x_{1} x_{3}+x_{13}, x_{123}^{2}-P x_{123}+Q\right\}
$$

where $P$ and $Q$ are as defined in Chapter 1. Buchberger's criterion states that a set $G$ of polynomials is a Grobner basis if and only if, for each $f_{1}, f_{2} \in \Gamma$,

$$
S\left(f_{1}, f_{2}\right)=\frac{\operatorname{lcm}\left(\operatorname{in}\left(f_{1}\right), \operatorname{in}\left(f_{2}\right)\right)}{L C\left(f_{1}\right) \operatorname{in}\left(f_{1}\right)} f_{1}-\frac{\operatorname{lcm}\left(\operatorname{in}\left(f_{1}\right), \operatorname{in}\left(f_{2}\right)\right)}{L C\left(f_{2}\right) L C\left(f_{2}\right)} f_{2}
$$

reduces to 0 on division by $G$. If we order, say, by total degree, then $T_{1}$ is not a Grobner basis, since

$$
S\left(x_{1}^{2}-2, x_{1} x_{2} x_{23}-x_{2} x_{123}-x_{1} x_{3}+x_{13}\right)=x_{1}^{2} x_{3}+x_{1} x_{2} x_{123}-x_{1} x_{13}-2 x_{2} x_{23}
$$

has remainder $x_{1} x_{2} x_{123}-x_{1} x_{13}-2 x_{2} x_{23}-2 x_{3}$ on division by $T_{1}$. A Tietze transformation of type (2), however, yields a presentation

$$
\Gamma=<g_{1}, g_{2} \mid g_{1}^{2}=\left(g_{1} g_{2}\right)^{3}=e>
$$

whose Gonzalez-Montesinos presentation is

$$
\begin{gathered}
T_{2}=\left\{x_{1}^{2}-4, x_{12}^{3}-3 x_{12}, x_{1}^{2} x_{2}-x_{1} x_{12}-4 x_{2}, x_{12}^{2} x_{1}-x_{12} x_{2}-x_{1}-x_{2}\right. \\
\left.x_{12}^{2} x_{2}-x_{12} x_{1}-x_{2}-x_{1}\right\}
\end{gathered}
$$

which is a Grobner basis.
If we wished to use a search algorithm to simplify a group presentation so that its Gonzalez-Montesinos polynomials are closer to being a Grobner basis, then we would need a way to compare two polynomials and determine which is more "Grobnerlike." Some ways to measure how far $T_{1}$ is from being a Grobner basis include the number of terms in the remainders of the S-polynomials, or the maximum degree of the remainders. But these measures would be computationally difficult to compute, and this would dominate the cost of the search. Moreover, most of the polynomials in the set of Gonzalez-Montesinos polynomials for a presentation will be determinantal or Fibonacci-type polynomials, and there is no obvious way to use our knowledge of the special properties of such polynomials to hasten the computation of these the remainders of their S-polynomials. We will instead take a slightly roundabout approach.

Recall from Chapter 1, corollary 1.5, that if $u=\operatorname{tr} w$, then the coefficient of $u^{n-2 l}$ in $\operatorname{tr} w^{n}$ is

$$
\begin{array}{cc}
1, & l=0 \\
\frac{(-1)^{l}}{l!} n \prod_{k=1}^{l-1}(n-l-k), & l>0 .
\end{array}
$$

This allows us to calculate the Fricke polynomials for relations of the form $w^{n}$ in $\Gamma$ efficiently, since the coefficient of $u^{n-2 l-2}$ may be found iteratively from the coefficient of $u^{n-2 l}$.

But calculating the Fricke polynomial of each potential new relation would easily dominate a search for a simpler group presentation, and probably be impractical for most interesting examples. Instead, one reasonable approach might be to evaluate these polynomials at some integral point. Following the advice of [81], an efficient way to do perform this calculation is to use the relation for Chebyshev polynomials of the first kind

$$
2 T_{n}(x) T_{m}(x)=T_{n+m}(x)+T_{n-m}(x)
$$

which implies the relation

$$
\left(\operatorname{tr} x^{n}\right)\left(\operatorname{tr} x^{m}\right)=\operatorname{tr} x^{n+m}+\operatorname{tr} x^{n-m}
$$

for Fricke polynomials. But the difficulty of guessing a good point (or a good set of points) at which to evaluate the trace polynomials in this manner makes this approach
seem unpromising.
In our example above, we saw that $\Gamma$ had a presentation in which its ideal of character relations' "natural" presentation was a structural Grobner basis. It would be nice if this were the norm, rather than the exception. In general, sets of polynomials whose leading terms with respect to a given term order are relatively prime are rather rare; and likewise, two randomly selected words in a free group will rarely have trace polynomials with relatively prime leading terms. (Of course, relators which arise in practical computations are not at all uniformly selected from the set of all words in a free group: they are, for purely physical reasons, restricted to words of some reasonable length.)

Since isomorphic finitely generated groups have isomorphic character varieties, the dimension of the character variety is invariant under Tietze transformations. For $\mathrm{SL}_{2} \mathbb{C}$ characters, the character variety of $G=<Y|R>,|Y|<\infty$ is the intersection of the character variety of the free group on $|\mathrm{Y}|$ letters, $X\left(\mathrm{~F}_{\mathrm{Y} \mid}\right)$, and the Whittemore variety arising from the relators $R, W(R)$. The automorphisms of $G$ induce morphisms of the character variety $X(G)$. (Since inner automorphisms of $G$ fix characters, we have an injection Out $G \hookrightarrow$ Aut $X(G)$.)

Neither of $X\left(F_{|Y|}\right)$ or $W(R)$ are invariant under automorphisms of $G$. However, since there exists a smallest $n$ such that $G$ is generated by $n$ elements, the dimension of the smallest such $X\left(F_{n}\right)$ is an invariant of the group $G$. Likewise, smallest dimension of a Whittemore variety corresponding to a set of relators defining $G$ is a well-defined
invariant of $G$. Let's denote the rank of $G$ by rank $G$, and the smallest dimension of a Whittemore variety corresponding to a set of relators defining $G$ (which may be on a larger set of generators) by $\operatorname{dim}_{w} G$. If $G=<Y|R\rangle$, with $|Y|=\operatorname{rank} G$, and if $<Z \mid S>$ is another finite presentation of $G$, then it is not necessarily true that $W(R)=W(S)$ implies that $|Z|=\operatorname{rank} G$ (as may easily be seen, for example, by the Klein four-group

$$
\left.<x, y\left|x^{2}, y^{2},(x y)^{2}\right\rangle=<x, y, z \mid x^{2}, y^{2},(x y)^{2}, x y z^{-1}>.\right)
$$

But it is true that, if $<Z \mid S>$ is obtained from $<Y \mid R>$ by a Tietze transformation of type (1), then $\operatorname{dim} W(R) \leq \operatorname{dim} W(S)$. For, adjoining a new generator to $<Y|R\rangle$ adds $1+|\mathrm{Y}|+\binom{|Y|}{2}$ new variables to the ideal of character relations, and adding new relation $t w$ for some $w \in Y$ adds the polynomial $p=(\operatorname{tr} t w)-2 . p$ cannot depend only on traces of words in $\langle\mathrm{Y}\rangle$. Furthermore, we have

Proposition 2.1 Suppose that $<\mathrm{Z} \mid \mathrm{S}>$ is obtained from $<\mathrm{Y} \mid \mathrm{R}>$ by a Tietze transformation of type (1), adjoining variable $z$ to $Y$, where $z=w$. Then $\operatorname{dim} W(R) \leq \operatorname{dim} W(S)$, and furthermore $\operatorname{dim} W(R)<\operatorname{dim} W(S)$ if $\operatorname{tr} z^{-1} w$ does not involve any of the new variable $\left\{\mathrm{t}_{z}, \mathrm{t}_{\mathrm{zy}}, \mathrm{t}_{\mathrm{zxy}} \mid \mathrm{x}, \mathrm{y} \in \mathrm{Y}\right\}$. In particular, a Tietze transformation of type (2) will always reduce the dimension of a Whittemore variety of a presentation, if the generator $z$ deleted only appears in a word of form $z^{-1} w$, where $\operatorname{tr} z^{-1} W$ does not involve each variable $\left\{t_{z}, t_{z y}, t_{z x y} \mid x, y \neq z\right\}$.

We are now in a position to define (and to try to justify) our definition of "simplicity" of a group presentation. We will say that a presentation of a group is simpler than another if its Whittemore variety has smaller dimension. Why is this a reasonable definition of simplicity? First, we should note that, given any presentation $P$, it is possible to find a sequence of Tietze transformations which reduce $P$ to "simplest" (really "locally simplest") form $P^{\prime}$; i.e., any transformation either makes $P^{\prime}$ less simple, or keeps it equally simple. Secondly, by the above proposition, a Tietze transformation of type (1) which makes a presentation less simple increases the number of generators. What about transformations of type (3) and (4)?

Theorem 2.2 Let $<Y=\left\{g_{1}, g_{2}, \ldots, g_{n}\right\} \mid R>$ be a finite presentation. Then
a A Tietze transformation of type (1) either increases the dimension of the Whittemore variety, or keeps it the same.
b A Tietze transformation of type (2) either reduces the dimension of the Whittemore variety, or keeps it the same.
c A Tietze transformation of type (3) or type (4) does not change the dimension of the Whittemore variety.

Proof: Parts $a$ and $b$ were proved in the last proposition. We will show $\mathbf{c}$ for Tietze transformation of type (3); this implies the statement for transformations of type (4), since type (4) transformations are inverses of transformations of type (3). Suppose
that a new relator T is added, which is the consequence of two relators $w, x$ in R; i.e., up to a cyclic permutation, $r=x w$. Adjoin variables $\operatorname{tr} x, \operatorname{tr} w$ to our set of indeterminates, and denote by I the ideal of character relations arising from the original relators. In particular,

$$
I^{\prime}=\left(\operatorname{tr} x-2, \operatorname{tr} w-2, \operatorname{tr} x g_{i}-t_{i}, \operatorname{tr} w g_{i}-t_{i}, \ldots\right) \subset I .
$$

We proceed by induction on the word-length of $x$. Suppose $|x|=1$, say $x=g_{1}$. Then we have immediately that $\operatorname{tr} r=\operatorname{tr} g_{1} w \equiv 2 \bmod I^{\prime}$. Furthermore, it is easy to see that $\operatorname{tr} \mathrm{g}_{\mathrm{i}} \mathrm{T} \equiv \mathrm{t}_{\mathrm{i}} \bmod \mathrm{I}^{\prime}$, since

$$
\begin{aligned}
\operatorname{tr} g_{1} r & =\operatorname{tr} g_{1}^{2} w \\
& =t_{1} \operatorname{tr} g_{1} w-\operatorname{tr} w \\
& \equiv t_{1}^{2}-\operatorname{tr} w \bmod I^{\prime} \\
& \equiv 4-2=2 \equiv t_{1} \bmod I^{\prime}
\end{aligned}
$$

and for $\mathrm{j} \neq 1$,

$$
\begin{aligned}
\operatorname{tr} g_{j} g_{1} w & =t_{1} \operatorname{tr} t_{j} w-\operatorname{tr} w g_{j} g_{1}^{-1} \\
& \equiv t_{1} t_{j}-\operatorname{tr} w g_{j} g_{1}^{-1} \bmod I^{\prime} \\
& \equiv t_{1} t_{j}-t_{j} \operatorname{tr} g_{1}^{-1} w+\operatorname{tr} g_{1}^{-1} g_{j}^{-1} w \bmod I^{\prime} \\
& \equiv t_{1 j} \operatorname{tr} w-\operatorname{tr} g_{j} g_{1} w \bmod I^{\prime} \\
2 \operatorname{tr} g_{j} g_{1} w & \equiv 2 t_{j} \cdot \bmod I^{\prime}
\end{aligned}
$$

Now suppose that part c is true when one relator has length $<\mathrm{n}$. Let $|x|=n$. By repeated application of the induction hypothesis, and the fundamental trace identity,

$$
\begin{aligned}
\operatorname{tr} g_{j} x w & =\operatorname{tr} x \operatorname{tr} t_{j} w-\operatorname{tr} w g_{j} x^{-1} \\
& \equiv \operatorname{tr} x t_{j}-\operatorname{tr} w g_{j} x^{-1} \bmod I^{\prime \prime} \\
& \equiv \operatorname{tr} x t_{j}-t_{j} \operatorname{tr} x^{-1} w+\operatorname{tr} x^{-1} g_{j}^{-1} w \bmod I^{\prime \prime} \\
& \equiv \operatorname{tr} x g_{j} \operatorname{tr} w-\operatorname{tr} g_{j} x w \bmod I^{\prime \prime} \\
2 \operatorname{tr} g_{j} x w & \equiv 2 t_{j} \bmod I^{\prime \prime} \\
\operatorname{tr} g_{j} x w & \equiv t_{j} \bmod I^{\prime \prime}
\end{aligned}
$$

where $j \in\{1,2, \ldots, n\}$ and $I^{\prime \prime}$ is generated by $I^{\prime}$, together with all elements of the form $\operatorname{tr} g_{j} y v-t_{j}$, for relators $y, v$ with $|y|<n$. Finally, we note that $\operatorname{tr} g_{j} x w-2 \equiv 0$, by substituting " $g_{j}=e$ " in the above calculation.

The simplest example of Theorem 2.2 is the free group of rank two $P_{1}=\langle x, y \mid\rangle$. Here, the Whittemore variety is unique: $W\left(P_{1}\right)=\mathbb{C}^{3}$, and so $\operatorname{dim} W\left(P_{1}\right)=3$. (It is clear that in fact $\operatorname{dim}_{W} P_{1}=3$.) A Tietze transformation of type (1) gives us the isomorphic group presentation $P_{2}=\langle x, y, z \mid z\rangle$. Now we have

$$
\begin{aligned}
W\left(P_{2}\right) & =\mathbb{V}\left(t_{z}-2, t_{x z}-t_{x}, t_{y z}-t_{y}, t_{z z}-t_{z}\right) \\
& =\mathbb{V}\left(t_{z}-2, t_{x z}-t_{x}, t_{y z}-t_{y}, 4-2-2\right) \\
& =\mathbb{V}\left(t_{z}-2, t_{x z}-t_{x}, t_{y z}-t_{y}\right)
\end{aligned}
$$

and so $\operatorname{dim} W\left(P_{2}\right)=4$.
If instead we apply a different Tietze transformation of type (1) to $P_{1}$, perhaps $P_{3}=\left\langle x, y, z \mid z y^{-1} x^{-1}\right\rangle$, then

$$
\begin{aligned}
W\left(P_{3}\right)= & \mathbb{V}\left(\operatorname{tr} z y^{-1} x^{-1}-2, \operatorname{tr} x z y^{-1} x^{-1}-t_{x}, \operatorname{tr} y z y^{-1} x^{-1}-t_{y}, \operatorname{tr} z^{2} y^{-1} x^{-1}-t_{z}\right) \\
= & \mathbb{V}\left(t_{z} t_{x y}-t_{x y z}-2, t_{z} t_{y}-t_{y z}-t_{x}, t_{y z} t_{x y}-t_{y} t_{x y z}+t_{x z}-t_{y},\right. \\
& \left.t_{z}^{2} t_{x y}-t_{z} t_{x y z}-t_{x y}-t_{z}\right)
\end{aligned}
$$

A Grobner basis for the polynomials defining $W\left(P_{3}\right)$, with respect to degreelexicographic order with $t_{x}>t_{y}>t_{y z}>t_{x y z}>t_{x y}>t_{z}>t_{x z}$ is:

$$
\begin{gathered}
\left\{t_{z} t_{y}-t_{y z}-t_{x}, \underline{t_{y}} t_{x y z}-t_{z} t_{y z}+t_{y}-t_{z}, \underline{t_{x y}}-t_{z}, \underline{t_{z}^{2}}-t_{x y z}-2,\right. \\
\left.\underline{t_{x} t_{x y z}}+t_{x}-t_{y z}-t_{x y z}-2, \underline{t_{z} t_{x}}-t_{y}-t_{z} \underline{t_{\underline{x}}^{2}}+t_{x} t_{y z}-t_{y}^{2}-t_{x}-t_{y z}\right\} .
\end{gathered}
$$

Inspection of the leading terms, which generate the initial ideal, shows that $\operatorname{dim} W\left(P_{3}\right)=3$.

This example, as simple as it is, exhibits two omissions in our definition of presentation simplicity. Firstly, we have not decided on a canonical way to choose the Whittemore variety of a presentation. Secondly, and more disturbingly, it seems that the dimension of $W(P)$ is much too coarse a measure of simplicity: after all, $P_{3}$ "looks" quite a bit different than $P_{1}$, yet they both have the same $\operatorname{dim} W(P)$.

Definition 2.3 (The canonical Whittemore variety of a presentation) Let $\mathrm{P}=<$ $\mathrm{Y} \mid \mathrm{R}>$ be a presentation. We will choose $\mathrm{W}(\mathrm{P})$ as follows: For each $w \in(\mathrm{Y} \cup\{\mathrm{e}\}) \mathrm{R}$, assign variable $\mathrm{p}(w)$ the value tr $w$.

1. First, we apply the fundamental trace relation repeatedly to make $p$ squarefree (i.e., we transform $p$ into an equivalent Fricke polynomial consisting of the traces of square-free words.) We always choose to eliminate the left-most syllable of length greater than one, and we always cyclically permute the word so that the syllable being permuted is on the right. (Note that this means that, in general, cyclic permutations of a relator may give different trace polynomials to the Whittemore variety. We choose this definition for its ease of computation, as opposed, say, to choosing the longest syllable, or choosing alphabetically.)

Example 2.4 We would transform $t_{12^{23-14}} \mapsto t_{3^{-1412^{2}}} \mapsto t_{3^{-1412}} t_{2}-t_{3^{-141}}$
2. We eliminate inverses in words in the same manner:

Example $2.5 t_{3^{-1} 41} \mapsto t_{413^{-1}} \mapsto t_{41} t_{3}-t_{413}\left(=t_{14} t_{3}-t_{134}\right)$
3. Each word in the trace expression now is a square-free, inverse-free word. For each word of length $\mathrm{N}>3$, apply Vogt's identity, with

$$
A=\text { the first } N-3 \text { letters of the word. }
$$

Repeat until all the words have length 3.
4. Finally, apply Fricke's lemma in the unique way, so that each trace is a trace of an ordered word on 3 or less letters.
(Fractional dimension of an ideal) Let $\prec$ be a graded term order on $\mathbb{C}[X]$. Let $\Gamma$ be the reduced lattice basis of $M$, a monomial ideal in $\mathbb{C}[X]$. The fractional dimension of $M$ is the real number

$$
\sum_{l \in \mathcal{A}} \min _{p \in \Gamma} \frac{\operatorname{dist}(l, p)}{\operatorname{dist}(l, p)+1}
$$

Let $\prec$ be a graded term order on $\mathbb{C}[X]$, and $I$ be an ideal in $\mathbb{C}[X]$. then the fractional dimension of I with respect to graded term order $\prec$ is the fractional dimension of in $\mathrm{in}_{\prec} \mathrm{I}$.

In the appendix (page 94) we present a GAP package that searches the space of Tietze transforms of a finitely presented group, attempting to minimize the "fractional" dimension of the Whittemore variety. We illustrate this technique with a GAP
session.
gap> Read( "FPFricke.g" ); Read( "fibs.g" );
fpfricke, Version . 899927
("fibs.g" defines the generalized Fibonacci groups [73]

$$
F(r, n)=\left\langle x_{1}, x_{2}, \ldots, x_{n} \mid\left\{x_{i} x_{i+1} \cdots x_{i+r-1} x_{i+r}^{-1} \mid i \in\{1,2, \ldots, n\}\right\}\right\rangle
$$

where the subscripts in the relators are taken mod $n$.)

```
gap> G := FibonacciGroup( 7,5 );
Group( f.1, f.2, f.3, f.4, f.5 )
gap> P := PresentationFpGroup( G );
<< presentation with 5 gens and 5 rels of total length 40 >>
gap> T := Copy( P );
<< presentation with 5 gens and 5 rels of total length 40 >>
gap> TzWhit( P );
gap> P;
<< presentation with 5 gens and 5 rels of total length 24 >>
```

We have reduced the total size of the relators from 40 characters to 24 , by using TzWhit, which tries to reduce the "partial" dimension of the Whittemore variety's defining polynomials coming directly from pairs of generators. Could we do better by considering triples of generators?

```
gap> TzWhitTriples( T );
#I new generator is _x6
gap> T;
<< presentation with 5 gens and 6 rels of total length 32 >>
gap> U := Copy( P );
<< presentation with 5 gens and 5 rels of total length 24 >>
gap> TzWhitTriples( U );
gap> U;
```

Apparently not! Let's try simplifying $F(7,5)$ with GAP's own presentation simplification function, TzGoGo.

```
V := Copy( P );
gap> TzGoGo( V );
#I there are 2 generators and 2 relators of total length 54
gap> G;
Group( f.1, f.2, f.3, f.4, f.5 )
gap> W := PresentationFpGroup( G );;
gap> TzGoGo( W );
#I there are 2 generators and 2 relators of total length 46
```

Can we do better than this, by applying TzWhit?

```
gap> TzWhit( V );
```

\#I new generator is _x7
gap> V;
<< presentation with 2 gens and 3 rels of total length 79 >>
gap> TzGoGo( V );

Apparently not.
gap> TzWhitTriples( V );
gap> V;
<< presentation with 2 gens and 2 rels of total length 79 >>
gap> TzGoGo( V );
gap> V;
<< presentation with 2 gens and 2 rels of total length 79 >>

It gets worse:

```
gap> TzWhit( V );
#I new generator is _x8
gap> V;
    << presentation with 2 gens and 3 rels of total length 104 >>
    gap> TzWhitTriples( V );
    gap> V;
    << presentation with 2 gens and 2 rels of total length 104 >>
    gap> TzGoGo( V );
```

```
#I there are 2 generators and 2 relators of total length 104
gap> TzWhit( V );
#I new generator is _x9
gap> V;
<< presentation with 2 gens and 3 rels of total length 129 >>
gap> TzGoGo( V );
#I there are 2 generators and 2 relators of total length }12
gap> TzWhit( V );
#I new generator is _x10
gap> V;
<< presentation with 2 gens and 3 rels of total length 154 >>
```

The actual relators here are:

```
gap> G6 := FpGroupPresentation( V );
Group( f.4, _x10 )
gap> G6.relators;
[f.4~2*_x10--1*f.4-4*_x10^-1*f.4~5*_x10--1*f.4~4*_x10~-1*f.4~3*
_x10*f.4*_x10^-1*f.4~4*_x10^-1*f.4~5*_x10^-1*f.4^4*_x1\
0^-1*f.4-6*_x10^-1*f.4-4*_x10^-1*f.4~5*_x10^-1*f.4~4*_x10^-1*f.4^4,
f.4^-3*_x10*f.4^-6*_x10*f.4^-4*_x10*f.4^-5*_x10*f.4^-4*_x10*f.4^-5*_x10*f.
~2*_x10^-1*f.4~4*_x10^-1*f.4~5*_x10^-1*f.\
4-4*_x10^-1*f.4-3*_x10*f.4^-4*_x10*f.4^-6*_x10*f.4^-4*_x10
```

```
*f.4^-5*_x10*f.4^-4*_x10*f.4^-2 ]
```

It would seem from this example (and from numerous other examples) that this definition of group presentation simplicity gives, at best, a minimal advantage over other standard presentation simplification techniques. On the other hand, the computational simplicity of our definition of presentation simplicity is quite appealing.

## Chapter 3

# Some invariant theory of the symmetric 

## group

In this chapter, we present an algorithm that can be used as part of a normal-form algorithm for $R\left(F_{n}\right)$. Our motivation, roughly, is as follows. We observe that $S_{n}$ acts on the Gonzalez-Montesinos relations ( 1.13 through 1.17 on page 16 ) in the natural way, and that a complete set of orbit representatives has 4 polynomials, for each $n>4$. Following [115] we may describe the points in $\mathbb{V}\left(I_{n}\right)=X\left(F_{n}\right)$ by examining separately a Grobner basis for these polynomials, and a Grobner basis for the ideal of the orbit variety of $\mathbb{C}^{\left(n^{3}+5 n\right) / 6}$ modulo this representation of $S_{n}$. Although this approach works for small $n$, it quickly bogs down in the complexity of finding and manipulating the invariant ring of this $\frac{1}{6}\left(n^{3}+5 n\right)$-dimensional permutation representation of $S_{n}$. But, as we shall see, some careful use of "constructive" Polya theory lets us extend this method considerably. More precisely, we are able to give a practical normal-form algorithm for the invariant ring of these representations of $S_{n}$. We do not however claim that this would be a practical method to find the normal form of an an element of $R(G)$, for a general finitely-presented group $G$.

An $n \times n$ complex matrix $\left[a_{i j}\right]$ acts on $\mathbb{C}\left[x_{1}, x_{2}, \ldots, x_{n}\right]=\mathbb{C}[X]$ by transforming
generators as follows: $x_{i}^{\left[a_{i j}\right]}=\sum_{j} a_{j i} x_{j}$. the action being extended to all of $\mathbb{C}[X]$. The fixed points in $\mathbb{C}[X]$ under the action of a complex matrix group $G$ form a $\mathbb{C}$-algebra, denoted by $\mathbb{C}[X]^{G}$, the "ring of invariants" of $G$. Rings of invariants of finite groups received much detailed study during this century. We mention particularly that $\mathbb{C}[X]^{G}$ is finitely generated (the "Hilbert finiteness theorem," see [39, section 1.4.1]) and that $\mathbb{C}[\mathrm{X}]^{\mathrm{G}}$ is Cohen-Macaulay (the Hochster-Eagon theorem [66]). When encountering a finitely generated algebra, we naturally ask whether we can list or describe a set of its generators; whether there exists a set of generators with nice properties; and even whether we may write down such a set.

The question of the existence of a procedure for constructing a generating set for $\mathbb{C}[X]^{G}$ was solved by Noether [97] who found a degree bound for a certain generating set of $\mathbb{C}[X]^{G}$. She then showed that this implies that there is a finite set of polynomials, whose image under the Reynolds operator generates all of $\mathbb{C}[X]^{G}$. (We recall the definition of the Reynolds operator below.) In this chapter, we will give a canonical basis (in the sense of Robbiano and Sweedler, whose work we briefly survey) for the ring of invariants of a particular representation of the symmetric group. We will also present a very explicit normal-form algorithm that uses this canonical basis.

A "canonical basis" (also called a "sagbi basis") B of an algebra $A \subseteq \mathbb{C}[X]$, with respect to a term order $\prec$, is a generating set for $A$ such that

$$
\left\langle\operatorname{in}_{\prec} B>=<\operatorname{in}_{\prec} A>.\right.
$$

[110]. For example, consider $\mathbb{C}[X]^{\mathbf{S}_{n}}$, the ring of symmetric functions, which are generated by

$$
\begin{array}{rlc}
\sigma_{1} & = & x_{1}+x_{2}+\cdots+x_{n} \\
\sigma_{2} & = & x_{1} x_{2}+x_{1} x_{3}+\cdots+x_{n-1} x_{n} \\
& \vdots & \vdots \\
\sigma_{n} & = & x_{1} x_{2} \ldots x_{n}
\end{array}
$$

the "elementary symmetric functions." (Here, we view $S_{n}$ as the group of $\pi \times \pi$ permutation matrices; i.e. $\mathbb{C}[X]^{S_{n}}$ is the ring of invariants of the symmetric group with the obvious action.) The elementary symmetric functions form a canonical basis for $\mathbb{C}[X]^{S_{n}}$ with respect to degree-lexicographic order (indeed, with respect to any term order [110, an observation attributed to Sturmfels].)

A finite canonical basis for an algebra allows us to write normal forms modulo $\mathcal{A}$, much as a Grobner basis allows us to write normal forms modulo an ideal. (When we write "modulo an algebra $A$," we mean of course modulo $A$ as a vector space. Also, we will consistently write $<B>$ for the subalgebra generated by $B$; ( $B$ ) for the ideal generated by B.) Instead of the familiar division algorithm of Grobner basis theory, in this algebra case we must use the "subduction" (subalgebra reduction) algorithm. The subduction algorithm proceeds as follows: Suppose that the algebra $B \subset \mathbb{C}[X]$ is generated by the set $\left\{f_{1}, f_{2}, \ldots, f_{n}\right\}$. To reduce $f \in \mathbb{C}[X]$ modulo $\langle B\rangle$, for $B=\left\{f_{1}, f_{2}, \ldots, f_{n}\right\}:[116]$

Algorithm 3.1 Subduction Algorithm. Given $f$ and $B=\left\{f_{1}, f_{2}, \ldots, f_{n}\right\}$. while $f \notin \mathbb{C}$ and there are non-negative integers $\left\{i_{1}, \ldots, i_{r}\right\}$ so that

$$
\begin{equation*}
\operatorname{in}_{\prec} f=c \cdot \prod_{j=1}^{r} \operatorname{in}_{\prec}\left(f_{j}\right)^{i_{j}}, c \neq 0 d o \tag{3.1}
\end{equation*}
$$

output $c f_{1}^{i_{1}} f_{2}^{i_{2}} \cdots f_{r}^{i_{r}}$
replace $f$ by $f-c f_{1}^{i_{1}} f_{2}^{i_{2}} \cdots f_{\tau}^{i_{r}}$
od
output $f$ as the remainder

When $\left\{f_{1}, f_{2}, \ldots, f_{n}\right\}$ is a canonical basis, then this algorithm will always return a remainder of zero for any $f \in B$. Conversely, if a set $\left\{f_{1}, f_{2}, \ldots, f_{n}\right\}$ subduces each $f \in B$ to zero then the set $\left\{f_{1}, f_{2}, \ldots, f_{n}\right\}$ is a canonical basis [98]. The expressions $c f_{1}^{i_{1}} f_{2}^{i_{2}} \cdots f_{r}^{i_{r}}$ found in this procedure, which we think of as monomials in $\mathbb{C}\left[" f_{i} "\right]$, are called superpositions, and are analogous to the more familiar S-polynomials of Grobner basis theory.

The permutation group $S_{n}$ has a natural representations as a group of $2^{n} \times 2^{n}$ matrices, and also as a group of $\binom{n}{m} \times\binom{ n}{m}$ matrices. We present an explicit canonical basis for the invariant ring of these representations of the symmetric group $S_{n}$ acting on the power set of $\{1,2, \ldots, n\}$. In other words, we consider the representation by
permutation matrices arising from this action of $S_{n}$, and let this matrix group act on

$$
X_{\binom{n}{m}}:=\left\{x_{12 \ldots m}, x_{12 \ldots(m-1)(m+1)}, \ldots, x_{(n-m+1) \ldots n}\right\}
$$

by permuting the indices of the elements of $X_{\binom{n}{m}}$ :

$$
\sigma\left(x_{j_{1}, \ldots, j_{m}}\right)=x_{\left\{\sigma_{i}, \ldots, \ldots j_{m} i\right.}
$$

for $\sigma \in S_{n}, 1 \leq j_{1}<\ldots<j_{m} \leq n$.
First, let's fix some notation. As is customary, we denote by $S_{n}^{(m)}$ this representation of $S_{n}$ acting pointwise on m-element subsets of $\{1 \ldots n\}$. Analogously, by $S_{n}^{\{m\}}$ we will denote the symmetric group $S_{n}$ acting on subsets with $m$ or fewer elements. (So, for example,

$$
\begin{gathered}
S_{n}^{(1)}=S_{n}^{\{1\}}=S_{n} \\
S_{n}^{(1)}+S_{n}^{(2)}+\cdots+S_{n}^{(m)}=S_{n}^{\{m\}}
\end{gathered}
$$

etc. )
We will henceforth write $\mathbb{C}[X]^{S_{n}^{(m)}}$ for $\mathbb{C}\left[X_{\left(\frac{n}{m}\right)}\right]^{S_{n}^{(m)}}$, and $\mathbb{C}[X]^{S_{n}^{(m)}}$ for $\mathbb{C}\left[X_{\left(\frac{n}{m}\right)} \cup \ldots \cup X_{\left(\frac{n}{2}\right)} \cup X_{\left(\frac{n}{1}\right)}\right]^{S_{n}^{\prime m i}}$. In [114] , Stanley introduced an ingenious scheme for
the construction of a set of generators for

$$
\mathbb{C}[X]^{s_{m}^{(m)}}
$$

see also [90]. We call a hypergraph with exactly $m$ vertices on each hyperedge an " $m$ hypergraph". We take the convention that graphs (resp. hypergraphs) are without loops, but may have multiple edges (resp. hyperedges.) Encode each monomial with coefficient 1,

$$
X_{\binom{n}{m}}^{\alpha}:=\prod_{i \leq i_{1}<\ldots<i_{m} \leq n} x_{i j}^{\alpha\left(i_{1}, \ldots k\right.},
$$

with the m-hypergraph on vertices $\{1,2, \ldots, n\}$, where there are exactly $\alpha\left(i_{1}, i_{2}, \ldots, i_{m}\right)$ hyperedges on the $m$ vertices $1 \leq i_{1}<\cdots<i_{m} \leq n$. Let $*=*_{s_{n}^{(m)}}: \mathbb{C}[X] \rightarrow \mathbb{C}[X]$ be the so-called "Reynolds operator" of the matrix group $S_{\mathfrak{n}}^{(\mathrm{m})}$, which averages the action of a group on a polynomial:

$$
* f=\frac{1}{n!} \sum_{g \in S_{n}^{(m)}} f^{g} .
$$

Then $* f \in \mathbb{C}_{n}^{(m)}$, and furthermore any $f$ in the ring of invariants may be written uniquely as a $\mathbb{C}$-linear combination of Reynolds operators of monomials. By extending linearly the identification of hypergraphs with vertices from $\{1,2, \ldots, n\}$ with monomials in $\mathbb{C}[X]$ to the free vector space generated by these graphs, Merris and

## Watkins observed:

Proposition 3.2 The Reynolds operator, restricted to the free vector space $M$ generated by a complete set of representatives of non-isomorphic hypergraphs with $k$ hyperedges on vertex set $\{1,2, \ldots, n\}$, is a vector space isomorphism from $M$ to the space $\left(\mathbb{C}[X]^{S_{n}^{(m)}}\right)_{k}=\left\{f \in \mathbb{C}_{n}^{S_{n}^{(m)}} \mid \operatorname{deg} f=k\right\}$.

Proof: see [90].
(By the way, Merris and Watkins were motivated by some computational problems in Polya theory. For example, by calculating the permutation character of $G=S_{n}^{(2)}$, and invoking Molien's theorem, which expresses the Hilbert series of $\mathbb{C}[X]^{\mathbf{G}}$ in terms of the characteristic polynomials of the elements of $G$, they verified a generating function for $a_{k}^{n}$, the number of nonisomorphic graphs on $n$ vertices with $k$ edges:

$$
\begin{aligned}
\sum_{k=0}^{\infty} a_{k}^{n} z^{k} & =\frac{1}{n!} \sum_{\sigma \in S_{n}^{(2)}} \frac{1}{\operatorname{det}(I-z \sigma)} \\
& =\frac{1}{n!} \sum_{\sigma}\left(1-z^{a}\right)^{-1}\left(1-z^{b}\right)^{-1} \cdots
\end{aligned}
$$

where $a, b, \ldots$ are the lengths of disjoint cycles of permutation $\sigma$.)
For the convenience of the reader, we briefly recall some classical invariant theory. A recent survey of algorithmic invariant theory can be found in [115]. An important fact about invariant rings is that their degree doesn't depend much on the action of the group:

Theorem 3.3 Let $\mathrm{G}<G L\left(\mathbb{C}^{\mathrm{d}}\right)$ be a finite matrix group. Then any $\mathrm{d}+1$ elements
of $\mathbb{C}\left[x_{1}, \ldots x_{\mathrm{d}}\right]^{\mathrm{G}}$ are $\mathbb{C}$-algebraically dependent. There exist d independent elements of $\mathbb{C}\left[x_{1}, \ldots x_{d}\right]^{G}$, thus $\operatorname{dim} \mathbb{C}[X]^{G}=\mathrm{d}$.

Proof: This is essentially theorems I and II of chapter XVII of Burnside's treatise [23]. Our proof basically follows his. Another proof may be found in [115, Theorem 2.1]. By the Hilbert finiteness theorem, $\mathbb{C}\left[x_{1}, \ldots x_{d}\right]^{G}$ is a finitely generated integral domain. Thus, by the Noether normalization theorem, $\operatorname{dim} \mathbb{C}\left[x_{1}, \ldots x_{d}\right]^{\mathbf{G}}$, the maximum length of a chain of prime ideals, is the length of any maximal chain of prime ideals. Since $G$ is finite, $\mathbb{C}\left[x_{1}, \ldots x_{d}\right]^{G}$ is Cohen-Macaulay. We exhibit a chain of length $d$, which is the longest possible chain in a Cohen-Macaulay subring of $\mathbb{C}\left[x_{1}, \ldots x_{d}\right]^{G}$.

Recall that if $I, J \subset R$ are ideals, $I+J$ is the smallest ideal that contains both $I$ and J. Define

$$
\begin{gathered}
I_{0}=\left(* x_{1}\right) \\
I_{1}=\left(* x_{1} x_{2}\right)+I_{0} \\
\vdots \\
I_{d-1}=\left(* x_{1} \cdots x_{d}\right)+I_{d-2} \\
I_{d}=\mathbb{C}[X]^{G} .
\end{gathered}
$$

These ideals evidently are a filtration for $\mathbb{C}\left[x_{1}, \ldots x_{d}\right]^{G}$. Recall that the set of associated primes Ass(I) of an ideal $I \subset R$ is the set of prime ideals of $R$ which annihilate some element of the module $R / I$. For each $i$, choose minimal $J_{i} \in$ Ass $I_{i}$. For $1 \leq i<d, I_{i}$ contains $* x_{1} x_{2} \cdots x_{d} \notin I_{i-1}$. Also, $I_{d}$ contains $1 \notin I_{d-1}$. So we have
the maximal chain of prime ideals

$$
\mathrm{J}_{0} \leq \mathrm{J}_{1} \leq \cdots \leq \mathrm{J}_{\mathrm{d}-1} \leq \mathrm{I}_{\mathrm{d}}=\mathbb{C}[X]^{\mathrm{G}}
$$

and thus $\operatorname{dim} \mathbb{C}\left[x_{1}, \ldots, x_{d}\right]^{G}=\mathrm{d} . \diamond$
(An interesting discussion of this result may be found in [114].)
Given a permutation group, let us say the permutation group $S_{n}^{(2)}$, then clearly any symmetric polynomial lies in the invariant ring $\mathbb{C}[X]^{\mathbf{s}_{n}^{(2)}}$. The elementary symmetric polynomials

$$
\begin{aligned}
\sigma_{1} & =x_{12}+x_{13}+\cdots+x_{(n-1) n} \\
& =\sum x_{i j}=* x_{12} \\
\sigma_{2} & =x_{12} x_{13}+\cdots \\
& =\frac{1}{2} *\left(x_{12} x_{13}\right)+\frac{1}{2} *\left(x_{12} x_{34}\right)
\end{aligned}
$$

generate the symmetric polynomials (and indeed, as we have noted, are a canonical basis for them.) It is well-known that $\mathbb{C}[X]^{S_{n}^{(2)}}$ is a free module over the ring $\left.<\sigma_{1}, \ldots\right\rangle$ (the existence of such a "regular system of parameters" is a popular way to define Cohen-Macaulayness.) We call these homogeneous invariants the primary invariants of the group $S_{n}^{(2)}$, and seek a finite set of invariants, called the secondary invariants,
which together with the primary invariants generate the whole invariant ring. (Such a set of secondary invariants exists, since invariant rings of finite groups are CohenMacaulay.) By corollary 2.7.10 of [115] the set of polynomials

$$
S=\left\{* m \mid m \text { is a descent monomial of } S_{\binom{n}{2}}\right\}
$$

are a set of secondary invariants of $S_{n}^{(2)}$. A "descent monomial" is a monic monomial which is associated to a permutation. Since $S_{\binom{n}{2}}$ has $\binom{n}{2}!$ members, this set of secondary invariants is clearly less than optimal. To "pick out" a minimal set of secondary invariants, one may use the Hilbert series of $\mathbb{C}[X]^{s_{n}^{(2)}}$, which tells us the number of algebraically independent invariants of a given degree. (Hilbert series of invariant rings may be found by a fundamental theorem of Molien, and are usually called Molien series in his honor.) The Molien series of $G \mathbb{C}[X]^{s_{n}^{(2)}}$ is, by Molien's theorem,

$$
\begin{aligned}
\frac{1}{n!} \sum_{M \in S_{n}^{(2)}} \frac{1}{\operatorname{det}(I-z M)} & =\frac{1}{n!} \sum_{M \in S_{n}^{(2)}} \prod_{i=1}^{\binom{n}{2}}\left(1-z^{i}\right)^{-l_{i}(M)} \\
& =\frac{1}{n!} \sum_{\sigma \in S_{n}} \prod_{i=1}^{\binom{n}{2}}\left(1-z^{i}\right)^{(2-n)\left(l_{i}(\sigma)\right)}
\end{aligned}
$$

(where $l_{i}(M)$ is the number of cycles of length $i$ in $M$ )

$$
=\sum_{a+\pi} \frac{\prod_{i=1}^{\operatorname{length}(a)}\left(1-z^{a(i)}\right)^{(n-2)}}{\prod_{i=1}^{\operatorname{length}(a)}\left(a(i)!i^{a(i)}\right)}
$$

$a \vdash n$ denoting that vector $a$ is a partition of integer $n$.
We do not take this approach here. Instead, we present a canonical subalgebra basis for the ring $\mathbb{C}[X]^{(2)}$, and a version of the subduction algorithm which uses this (very large) basis implicitly.

Theorem 3.4 Recall that we have associated labelled hypergraphs with monomials.
a) Let S be a set of equivalence-class representatives of labelled hypergraphs partitioned by isomorphism on vertices $\{1,2, \ldots, \mathrm{n}\}$, with k m-hyperedges (resp., k hyperedges, each edge on $k$ or fewer vertices.) Then $* S$ is a canonical basis for $\left\langle\left(\mathbb{C}[X]^{s_{n}^{(m)}}\right)_{k}\right\rangle$ (the algebra generated by the set of degree $k$ elements of the ring $\left.\mathbb{C}[X]_{n}^{S_{n}^{(m)}}\right)$ (resp. $\left.\left\langle\left(\mathbb{C}[X]^{S_{n}^{(m!}}\right)_{k}\right\rangle\right)$ with respect to any term order. Furthermore, $S$ contains a regular system of parameters $f\left\langle\left(\mathbb{C}[X]^{\prime}[X]^{s_{n}^{(m)}}\right)_{k}\right\rangle\left(\right.$ resp. $\left.\left\langle\left(\mathbb{C}[X]^{S_{n}^{[m]}}\right)_{k}\right\rangle\right)$ under the filtration induced by degree grading.
b) The image under $*$ of a set of representatives of all labelled hypergraphs up to isomorphism on vertices $\{1,2, \ldots, n\}$ with $\binom{n}{m}$ or fewer hyperedges, each on $m$ vertices (resp. $m$ or fewer vertices) is a canonical basis for $\mathbb{C}[X]^{s_{n}^{(m)}}$ (resp. $\mathbb{C}[X]^{S_{n}^{(m)}}$ ) with respect to any term order.

Proof: Fix any term order $\prec$.
a) Let $f$ be in $\left(\mathbb{C}[X]^{S_{m}^{(m)}}\right)_{k} \backslash\{0\}$, for $k>0$. Then (in $\left.f f\right) / L C_{\prec}(f)$ (where $L C_{\prec}(f)$ denotes the leading coefficient of $f$ with respect to $\prec$ ) is a hypergraph with $k$ mhyperedges, since $\operatorname{deg}$ in $_{\downarrow} f=k$. Since $f$ is in $\mathbb{C}[X]^{!^{m!}!}$,

$$
f-L C(f) \cdot *\left(\operatorname{in}_{\prec} f / L C(f)\right)=f_{1}
$$

is in $\mathbb{C}[X]^{s_{n}^{m i}}$, and furthermore $f_{1}$ is again homogeneous of degree $k$, and has fewer terms than $f$ (since $f$ is not constant.) Continuing this process, we eventually reach a constant, which is zero since $f$ is homogeneous of positive degree. Thus, the set of monomials associated to the set of hypergraphs with $\mathfrak{n}$ vertices and $k$ m-hyperedges generate $<\left(\mathbb{C}[X]^{S_{n}^{\prime m} \mid}\right)_{k}>$ as an algebra. Also, each step in our reduction is a step in the subduction algorithm, thus these hypergraphs are a canonical basis for

$$
<\left(\mathbb{C}[X]^{\mathrm{s}^{\prime m i}}\right)_{k}>
$$

If $f \in\left(\mathbb{C}[X]^{S_{n}^{(m)}}\right)_{k} \backslash\{0\}, k>0$, then (in $\left.f\right) / L_{C_{\prec}}(f)$ is a hypergraph with $k$ hyperedges, each with $m$ or fewer vertices. The proof proceeds as above.
b) Since $\mathbb{C}[X]^{s_{n}^{\prime m}}=\bigoplus_{\mathbf{k}}\left(\mathbb{C}[X]^{s_{m}^{\{m i}}\right)_{k}$, we have by part (a) an (infinite) canonical basis for $\mathbb{C}[X]^{S_{n}^{m!}}$ : the image under $*$ of all hypergraphs (up to isomorphism) on $\{1 \ldots n\}$, with each hyperedge on $m$ vertices. We would like to express each $f \in$ $\left(\mathbb{C}[X]^{s_{n}^{(m)}}\right)_{D}, D>\binom{n}{m}$ as an algebraic combination of $\left\{* f_{l} \mid f_{l} \in\left(\mathbb{C}[X]^{s_{n}^{(m)}}\right)_{l}, l \leq\binom{ n}{m}\right\}$. Now products of Reynolds operators of simple hypergraphs (hypergraphs without
multiple hyperedges) must have multiple hyperedges, thus the set of simple hypergraphs, each of whose Reynolds operator is contained in some $\left(\mathbb{C}[X]^{S_{n}^{(m)}}\right)_{\mathfrak{l}}, l \leq\binom{ n}{m}$, contains $\binom{n}{m}$ algebraically independent elements. Statement (b) is now true by theorem 3.3. $\diamond$

Although we are most interested in using the theorem in the case when $m \leq 3$, it is nonetheless unfortunate that these canonical bases grow so quickly (there are many 3hypergraphs on 9 vertices.) For this reason, any computer algorithm which operates on an explicit canonical basis for $\mathbb{C}[X]^{S_{n}^{(3)}}$ will be confined, for practical reasons, to small $n$. But the subduction algorithm itself only requires that we be able to exhibit elements of the canonical basis satisfying condition 3.1.

Algorithm 3.5 Algorithm (Subduction for $\mathbb{C}[X]^{S_{n}^{|m|}}$.) Given: $f$ in $\mathbb{C}[X]$,
Output: a normal form for $f \bmod \mathbb{C}[X]^{\mathrm{s}_{\mathrm{n}}^{(m)}}$, together with a sequence of superpositions
$F:=f$.
done $:=$ false
repeat
if F is constant then
output $F$ as a term of the expression for $f$ done $:=$ true
elseif there are hypergraphs $\left\{\Gamma_{i}\right\}$ so that:
a) $\mathrm{in}_{\prec} * \Gamma_{\mathrm{i}}=\Gamma_{\mathrm{i}}$
b)each $\Gamma_{\mathfrak{i}}$ has $\binom{n}{m}$ edges or less, and
c) $\mathrm{in}_{\prec} \mathrm{F}=\prod_{i} \Gamma_{i}$
then
output $L C(f) \cdot \prod_{i} * \Gamma_{\mathfrak{i}}$ as a term of the superposition
$\mathrm{F}:=\mathrm{f}-L C(\mathrm{f}) \cdot \Pi_{\mathrm{i}} * \Gamma_{\mathrm{i}}$
else
done := true
$f$
until done
output F as a normal form for f .

We have considerable leeway in our choice of $\left\{\Gamma_{i}\right\}$ at each step. But part (d) of the theorem (together with the theory of canonical bases - see e.g. [98] or [92]) guarantee that the algorithm will terminate with a normal form for $f \bmod \mathbb{C}[X]^{s_{n}^{[m]}}$.

As an example, let's look at

$$
\begin{equation*}
f_{1}=x_{12}^{2}+x_{13}^{4} x_{123}+2 \in \mathbb{C}\left[X_{4}\right] \tag{3.2}
\end{equation*}
$$

with lex term order, $x_{1} \succ x_{2} \succ \ldots \succ x_{123} \succ \ldots \succ x_{234}$. We have

$$
\operatorname{LT}\left(f_{1}\right) / \operatorname{LC}\left(f_{1}\right)=\mathscr{V}^{2}
$$

so we write $f_{1}$ as

$$
f_{1}=*\left(x_{12}^{2}\right)+\left(f_{1}-*\left(x_{12}^{2}\right)\right)
$$

$$
=*\left(x_{12}^{2}\right)+\underbrace{-x_{13}^{2}-x_{14}^{2}-x_{34}^{2}+x_{13}^{4} x_{123}+2}_{f_{2}}
$$

Now $\operatorname{LT}\left(f_{2}\right) / L C\left(f_{2}\right)=3 \begin{array}{ll}1,2 & 3\end{array}$ on vertices $\{1,2,3\}$ ) cannot be written as hypergraphs such that $\operatorname{in}_{\prec} \Gamma=$ in $_{\prec} * \Gamma$, and thus $f_{2}$ is the normal form for $f_{1} \bmod \mathbb{C}\left[X_{\binom{n}{2}}, X_{\binom{n}{3}}\right]^{]_{n}^{s m i}}$.

How may we recognize when an appropriate choice $\left\{\Gamma_{i}\right\}$ exists? In other words, how may we recognize that an appropriate superposition 3.1 is available? (We restrict ourselves to hypergraphs with the same number of vertices on each hyperedge, since all of the hyperedges in the orbit of a hyperedge lie on the same number of vertices.) One reasonable idea is to greedily choose higher-weight hyperedges to form a maximal graph $\Gamma_{\mathrm{i}}$ so that in $\Gamma_{\mathrm{i}}=\mathrm{in} \mathrm{n}_{\prec} * \Gamma_{\mathrm{i}}$ (let's call such a $\Gamma_{\mathrm{i}}$ an "initial graph" (or "initial hypergraph.") If this can be done, we have reduced the problem to a smaller graph $\operatorname{LT}\left(f_{2}\right) /\left(L C\left(f_{2}\right) \cdot \Gamma_{i}\right)$. But for the monomial

$$
f=x_{12}^{2} x_{13}^{2} x_{14}
$$

with lex order $\prec$, this greedy algorithm would first choose $\Gamma_{1}=x_{12}^{2}$, and then fail for $f / x_{12}^{2}=x_{13}^{2} x_{14}$, which is worrisome since $f$ is already an initial graph. (Clearly, on the other hand, if $\prec$ were degree ordering, this greedy algorithm would work for $f$ ) . Let us examine this situation more closely, using the language of greedoids [83].

A "greedoid" is a sort of generalized matroid. Consider the following five conditions on the ordered pair $(E, \mathcal{G} \subseteq \operatorname{PowerSet}(E)):$

1. $\emptyset \in \mathcal{G}$
2. (Accessibility Axiom) $X \in \mathcal{G}, X \neq \varnothing$ implies that there is $x \in X$ such that $X \backslash\{x\} \in \mathcal{G}$
3. (Exchange Axiom) $X, Y \in \mathcal{G}|X|=|Y|+1$ implies that there is $x \in X \backslash Y$ such that $Y \cup\{x\} \in \mathcal{G}$.
4. If $A \subset B \in \mathcal{G}$, then $A \in E$.
5. $\mathcal{G}$ is closed under union.

If the pair $(E, \mathcal{G} \subseteq$ PowerSet $(E)$ ) satisfies conditions (1)-(3), it is called a (simple) greedoid. If the pair satisfies conditions (1) - (4), it is called a matroid. If the pair satisfies conditions (1) - (3) and (5), it is called an antimatroid. The elements of $\mathcal{G}$ are called the feasible sets of the greedoid. A maximal feasible set is called a basis. If we furthermore allow elements of $\mathcal{G}$ to be multisets, rather than sets, the greedoid is called non-simple.

See [22] or [83] for a survey of greedoids and their applications to combinatorial optimization.

The multiplicity-free initial graphs (in other words, the square-free monomials) form a greedoid $\mathcal{G}$ as (hyper-)edge sets of $\{1 . . n\}$ under $\subseteq$. We declare $\oslash$ to be the initial graph of 1 , so that $\oslash \in \mathcal{G}$. The inclusion poset of $\mathcal{G}$, when $m=2$, and $\prec$ is lex order, $x_{12} \succ x_{13} \succ \cdots$ is shown in figure 3-1.

We have drawn the graphs in 3-1 so that if $\Gamma_{1}$ is to the left of $\Gamma_{2}$, then $\Gamma_{1} \prec \Gamma_{2}$. Note that

1) $\mathcal{G}$ is not a matroid when $n>2 . \mathcal{G}$ is not closed under union (e.g.

2) $\quad \mathrm{K}_{n}$, the complete graph, is the unique basis (maximal feasible set) of $\mathcal{G}$. Thus $\operatorname{rank}(\mathcal{G})=\operatorname{rank}\left(K_{n}\right)=n!$.
3) By construction, any objective function which chooses $\prec$-greater graphs would be compatible with this greedoid structure, thus the greedy algorithm applied to $\mathcal{G}$ will find a basis which maximizes this function - in other words, $\mathrm{K}_{\mathrm{n}}$.

If $\Gamma \in \mathcal{G}$, then the interval $[\varnothing, \Gamma]$ is a greedoid where the greedy algorithm produces $\Gamma$, which is the unique basis. Let $\Gamma$ be a graph. We define a greedoid $\mathcal{G}(\Gamma)$ as follows: If $\Gamma^{\prime \prime}$ is a multiplicity-free initial graph subgraph of $\Gamma$, which occurs $\tau$ times in $\Gamma$, then


Figure 3-1 Inclusion poset for $\mathcal{G}$
the multiset (subgraph)

$$
\{\underbrace{}_{\tau \text { times }}, e, \ldots, e \mid e \text { an edge of } \Gamma\}
$$

is a feasible set of $\mathcal{G}(\Gamma)$, which we call the initial graph branching matroid of $\Gamma$. (If $\Gamma$ contains no initial graph, we define $\mathcal{G}(\Gamma)=\varnothing$.) Non-simple greedoids maximize compatible objective functions with the greedy algorithm in the same way that simple greedoids do - in this case to provide us with a $\prec$-maximal basis (included initial graph) of $\Gamma$. Indeed, this greedy algorithm is just a mutation of Prim's 'visit an unvisited node first' procedure to find a maximal-weight spanning tree in a graph.

The following theorem gives us a termination condition for the subduction algorithm 3.5.

Theorem 3.6 Let $\Gamma_{0}$ be the basis found by greedily searching for a maximal initial graph subgraph of $\Gamma$, in the greedoid $\mathcal{G}(\Gamma)$. If

$$
\Gamma \backslash \Gamma_{0}
$$

is non-empty and contains no initial graph, then condition 3.1 in the subduction algorithm for $\mathbb{C}[X]^{S_{n}}$ cannot be satisfied for $F=\Gamma$.

Proof: Suppose on the contrary that there exist initial graphs $\Gamma_{1}, \Gamma_{2}, \ldots \Gamma_{\mathfrak{l}}$ so that

$$
\Gamma=\Gamma_{1} \Gamma_{2} \ldots \Gamma_{l}
$$

Then, since $\prec$ is a term order, $\Gamma$ is an initial graph, so $\Gamma \in \mathcal{G}(\Gamma)$, and $\mathcal{G}(\Gamma)$ is just the interval $[\varnothing, \Gamma]$ in $\mathcal{G}\left(\mathrm{K}_{n}\right)$. But then $\Gamma \backslash \Gamma_{0}=\varnothing$, since $\Gamma$ is itself the unique basis in $[\oslash, \Gamma] \diamond$

If we encode a graph by the $\frac{\mathrm{n}^{2}-n}{2}$-tuple consisting of the entries in the upper triangular part of the graph's adjacency matrix, (e.g. 3 has adjacency matrix

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

so we write code $\left.\left(\int_{3}^{1_{2}^{2}}\right)=(1,1,0)\right)$, then the task of greedily searching $\mathcal{G}(\Gamma)$ with respect to lex-order $\prec$ is:

## repeat

new := largest code (considered as a binary integer) obtained by changing the first zero which has no ones to the right of it to one
until new is maximal in $\Gamma$ and is not isomorphic to any graph previously considered

The number of operations that must be performed in this loop is dominated by the difficulty of checking each of the graphs constructed for isomorphism with previous graphs. Consider the monomial associated to graph


If $m$ contains a spanning initial graph subgraph, then this graph has three vertices (1,2, and 4) of degree 2, one vertex (2) of degree 3, and one vertex (5) of degree 1. If greedily adding an edge would cause us to exceed one of these degree bounds, then we have found our basis of $\mathcal{G}(\Gamma)$. R. Read [107], [78] has considered this problem in the context of the problem of listing all simple graphs on $\eta$ vertices up to isomorphism. He showed that, for lex order, it suffices in the above loop to ensure that the code of new is canonical - is, maximal under all permutations of the vertices. (He coined the term orderly algorithm for this sort of graph-cataloging procedure.)

To greedily search $\mathcal{G}(\Gamma)$, we start with $\varnothing$, and add the single edge 1-2. We add edge 1-3, which does not exceed our vertex degree bounds, and which is canonical. We may not add edges 1-4 or 1-5, since degree(1) now equals 2 . We add edge $2-3$, which is canonical. We add edge 2-4, which is canonical. We may not add 2-5. Adding edge 4-5 completes our search.

Checking a graph for canonicality can take as many as $n!$ comparisons. But for
most graphs, when $\preceq$ is lex order, the situation is not that dire -the new graph is by definition canonical.

We summarize our discussion with an efficient normal form algorithm for $\mathbb{C}[X]^{\mathrm{s}_{n}^{(2)}}$.

Algorithm 3.7 Algorithm: Subduction for $\mathbb{C}[X]^{S_{n}^{(2)}}$ with respect to lexicographic order. Given: $f \in \mathbb{C}[X]$

Output: A normal form for $f \bmod \mathbb{C}[X]^{S_{n}^{(2)}}$.
$F:=f$
done := false
repeat

$$
\begin{aligned}
& \text { if } \mathrm{F} \in \mathbb{C} \text { then } \\
& \quad \text { done }:=\text { true } \\
& \text { else } \\
& \text { new }:=(0,0, \cdots, 0) \\
& \text { newer }:=\text { new } \\
& \text { success }:=\text { false "orderly algorithm") } \\
& \text { repeat }
\end{aligned}
$$

change the first zero in new, which is to the right of all ones, to one
if this new edge is in F then

$$
\text { success }:=\text { true }
$$

```
                    new:= newer
                    fi
            until success
            or
            (no edge that would be represented by a position to the right of
            rightmost one is in F)
            if success then
            m:= multiplicity of new in F
            F:=F-LC(F)\cdot(*new)}\mp@subsup{)}{}{m
            else
                        done:= true
            fi
    f
until done
output F as a normal form for f
```


## Chapter 4

## Divisibility properties of trace polynomials

### 4.4.1 The shifted trace polynomials; strong and weak divisibility

In this chapter, we give some divisibility properties of character relations and trace polynomials.

Theorem 4.1 (horowitz) Let $w_{1}=g_{1}^{k}, w_{2}=g_{1}^{2}, g_{1} \in F_{n}$. If $\operatorname{gcd}(k . l)=1$, then

$$
\operatorname{gcd}\left(t_{w_{1}}-2, t_{w_{2}}-2\right)=t_{g_{1}^{\operatorname{gcd}(k . t)}}-2
$$

In other words, the sequence of polynomials $\left\{\mathrm{t}_{\mathrm{g}_{\mathrm{i}}}-2\right\}_{\mathfrak{i}}$ is a strong divisibility sequence.

We begin our proof with a version of Theorem B of [94]

Proposition 4.2 For $z \in \mathbb{N}$, let $\left\{f_{n}\right\}$ be the sequence of integers determined by

$$
\begin{align*}
f_{n+2} & =a\left(f_{n+1}+c\right)+b\left(f_{n}+c\right)-c  \tag{4.1}\\
f_{0} & =0 \\
f_{1} & =z
\end{align*}
$$

Suppose that $\mathrm{a}, \mathrm{b}, \mathrm{c}$ are pairwise relatively prime.
Then, for sufficiently large prime d , there is an integer k so that

$$
\mathrm{k}|n \Rightarrow \mathrm{~d}| \mathrm{f}_{\mathrm{n}}
$$

Proof: Set $\mathfrak{I}=\left\{n|d| f_{n}\right\}$. Since $f_{0}=0, \mathfrak{I}$ is not empty. If $\mathfrak{I} \neq\{0\}$, then let $n$ be the element of $\mathfrak{I}$ with smallest absolute value. Then, following [94], there is an integer $\mathrm{d}^{\prime}$ so that $\mathrm{d}>\mathrm{d}^{\prime}$ implies

$$
f_{n+k}+(-b)^{k} f_{n-k} \equiv 0 \quad \bmod d
$$

Thus, $2 n, 3 n, \ldots \in \mathfrak{I}$, and likewise $-n,-2 n, \ldots \in \mathcal{I} . \diamond$

Corollary 4.3 (to proposition) $f_{\operatorname{gcd}(m, n)}= \pm C \operatorname{gcd}\left(f_{m}, f_{n}\right)$.

Proof: (of theorem) We apply the proposition with $c=2, z=t_{1}=\operatorname{tr} g_{1}$. By the recurrence 1.7, and corollary 1.5, $\left(\operatorname{tr} x^{n}\right)-2=f_{n}\left(t_{x}\right)$ for any word $x$. Also, $\operatorname{deg}_{t_{x}}\left(\operatorname{tr} x^{n}\right)=n$, so $\operatorname{deg}_{t_{1}}\left(t_{g_{1} \operatorname{gcd}(k .1)}\right) \leq \operatorname{gcd}(k . l)$. For an infinite number of primes $d$, we have by the above corollary

$$
\operatorname{gcd}\left(t_{w_{1}}-\left.2\right|_{t_{1}=z}, t_{w_{2}}-\left.2\right|_{t_{1}=z} \equiv \pm\left.\left(t_{1 \operatorname{scd}(k, l)}-2\right)\right|_{t_{1}=z} \bmod d\right.
$$

Since each of the polynomials $t_{w_{1}}-2, t_{w_{2}}-2$, and $t_{1 g e d(K, L)}-2$ are monic, the theorem follows. $\diamond$

Let $\left\{a_{n}\right\}$ be a sequence of elements of a unique factorization domain. $\left\{a_{n}\right\}$ is called a divisibility sequence [119] if $a_{n} \mid a_{m}$ implies that $n \mid m . \quad a_{n}$ is called a strong divisibility sequence if, in addition $a_{\operatorname{gcd}(m, n)}= \pm \operatorname{gcd}\left(a_{m}, a_{n}\right)$. Corollary 4.3 states that the shifted trace polynomials $f_{n}$ are a strong divisibility sequence. For each integer $n \geq 0$, let $h_{n}$ be the polynomial

$$
h_{n}=\frac{f_{n}-2}{\operatorname{lcm}\left\{f_{i}-2|i| n\right\}}
$$

The sequence of polynomials $h_{n}(x)$ was first considered, in a different context, by Horadam, Loh and Shannon [67], who noted that $f_{n}(x)$ is a divisibility sequence. For positive integer $a$, define a sequence of polynomials

$$
g_{n+2}=x g_{n+1}-g_{n}
$$

$$
\begin{aligned}
& g_{0}=a \\
& g_{1}=x
\end{aligned}
$$

with associated $H_{n}=g_{n}-a$,

$$
h_{n}=\frac{H_{n}}{\operatorname{lcm}\left\{H_{n}|i| n\right\}}
$$

|  | $\operatorname{tr} y^{n}-2$ |
| ---: | :--- |
| 1 | $x-2$ |
| 2 | $(x-2)(x+2)$ |
| 3 | $(x-2)(x+1)^{2}$ |
| 4 | $(x-2)(x+2) x^{2}$ |
| 5 | $(x-2)\left(x^{2}+x-1\right)^{2}$ |
| 6 | $(x-2)(x+2)(x-1)^{2}(x+1)^{2}$ |
| 7 | $(x-2)\left(x^{3}+x^{2}-2 x-1\right)^{2}$ |
| 8 | $(x-2)(x+2)\left(x^{2}-2\right)^{2} x^{2}$ |
| 9 | $(x-2)(x+1)^{2}\left(x^{3}-3 x+1\right)^{2}$ |
| 10 | $(x-2)(x+2)\left(x^{2}-x-1\right)^{2}\left(x^{2}+x-1\right)^{2}$ |
| 11 | $(x-2)\left(x^{5}+x^{4}-4 x^{3}-3 x^{2}+3 x+1\right)^{2}$ |
| 12 | $(x-2)(x+2)(x-1)^{2}(x+1)^{2}\left(x^{2}-3\right)^{2} x^{2}$ |
| 13 | $(x-2)\left(x^{6}+x^{5}-5 x^{4}-4 x^{3}+6 x^{2}+3 x-1\right)^{2}$ |
| 14 | $(x-2)(x+2)\left(x^{3}-x^{2}-2 x+1\right)^{2}\left(x^{3}+x^{2}-2 x-1\right)^{2}$ |
| 15 | $(x-2)(x+1)^{2}\left(x^{2}+x-1\right)^{2}\left(x^{4}-x^{3}-4 x^{2}+4 x+1\right)^{2}$ |

Table 4.1 The shifted trace polynomials

We aim to generalize the results of [67] to consider the notion of "weak divisibility." Let $k \in \mathbb{Z}$, and let $\left\{a_{n}\right\}$ be a sequence in a unique factorization domain. We call $\left\{a_{n}\right\}$ $a^{2}$ " $k$-weak divisibility sequence" if, for all $l>0$ such that $k \mid l, a_{k}$ divides $a_{l}$. Clearly, divisibility sequences are weak divisibility sequences for all $k$.

Since we have seen that $\alpha=2$ implies that $H_{n}$ is a divisibility sequence, we ask: for what $a \in \mathbb{C}$ is $H_{n}$ a $k$-weak divisibility sequence? In the polynomial ring $\mathbb{C}[x, a]$, let $<$ be degree-lexicographic order [31], with $a<x$. The ideal $I=(x-a)$ is principal, and so $\{x-a\}$ is a Grobner basis for $I$. Let $r_{n}(x, a) \in \mathbb{C}[x, a]$ be the normal form of $H_{n}(x, a) \bmod$. If a term of $r_{n}$ included the variable $x$ to a positive power, then the leading term of $T_{n}$ would also include $x$ to a positive power, and we could reduce $r_{n}$ by $x-a$. But this would contradict the assumption that $r_{n}$ is the normal form of a polynomial. Thus we have

|  | $\operatorname{tr} y^{n}$ |
| ---: | :--- |
| 1 | $x$ |
| 2 | $x^{2}-2$ |
| 3 | $\left(x^{2}-3\right) x$ |
| 4 | $x^{4}-4 x^{2}+2$ |
| 5 | $\left(x^{4}-5 x^{2}+5\right) x$ |
| 6 | $\left(x^{2}-2\right)\left(x^{4}-4 x^{2}+1\right)$ |
| 7 | $\left(x^{6}-7 x^{4}+14 x^{2}-7\right) x$ |
| 8 | $x^{8}-8 x^{6}+20 x^{4}-16 x^{2}+2$ |
| 9 | $\left(x^{2}-3\right)\left(x^{6}-6 x^{4}+9 x^{2}-3\right) x$ |
| 10 | $\left(x^{2}-2\right)\left(x^{8}-8 x^{6}+19 x^{4}-12 x^{2}+1\right)$ |
| 11 | $\left(x^{10}-11 x^{8}+44 x^{6}-77 x^{4}+55 x^{2}-11\right) x$ |
| 12 | $\left(x^{4}-4 x^{2}+2\right)\left(x^{8}-8 x^{6}+20 x^{4}-16 x^{2}+1\right)$ |
| 13 | $\left(x^{12}-13 x^{10}+65 x^{8}-156 x^{6}+182 x^{4}-91 x^{2}+13\right) x$ |
| 14 | $\left(x^{2}-2\right)\left(x^{12}-12 x^{10}+53 x^{8}-104 x^{6}+86 x^{4}-24 x^{2}+1\right)$ |
| 15 | $\left(x^{2}-3\right)\left(x^{4}-5 x^{2}+5\right)\left(x^{8}-7 x^{6}+14 x^{4}-8 x^{2}+1\right) x$ |

Table 4.2 The trace polynomials factored over the integers
$r_{n}(x, a) \in \mathbb{C}[a]$. The roots of the polynomials $r_{n}$ thus are the complex numbers $a$ such that $H_{n}$ is an $n$-weak divisibility sequence. The first few $r_{n}$ are listed in table 4.3.

### 4.4.2 A discriminant identity

We now fix some notation for the rest of the chapter.

Notation 4.4 We denote by $f_{n}(x)$ the $(n)^{\text {th }}$-degree polynomial

$$
\sum_{i}(-1)^{\lfloor i / 2\rfloor}\binom{n-\lfloor(i+1) / 2\rfloor}{\lfloor i / 2\rfloor} x^{i}
$$

We denote by $\alpha_{n, i}$ the coefficient of $x^{n-i+1}$ in $f_{n}$. When $n$ is clear, we write $\alpha_{i}$ instead

$$
\begin{aligned}
& 0 \\
& -2 a+a^{2} \\
& -2 a-a^{2}+a^{3} \\
& -2 a^{2}-a^{3}+a^{4} \\
& 2 a^{2}-3 a^{3}-a^{4}+a^{5} \\
& -2 a+3 a^{2}+3 a^{3}-4 a^{4}-a^{5}+a^{6} \\
& -2 a-3 a^{2}+6 a^{3}+4 a^{4}-5 a^{5}-a^{6}+a^{7} \\
& -4 a^{2}-6 a^{3}+10 a^{4}+5 a^{5}-6 a^{6}-a^{7}+a^{8} \\
& 4 a^{2}-10 a^{3}-10 a^{4}+15 a^{5}+6 a^{6}-7 a^{7}-a^{8}+a^{9} \\
& -2 a+5 a^{2}+10 a^{3}-20 a^{4}-15 a^{5}+21 a^{6}+7 a^{7}-8 a^{8}-a^{9}+a^{10}
\end{aligned}
$$

Table 4.3 The remainder polynomials $r_{n}(a)$

|  |  |
| :--- | :--- |
| $r_{1}$ | 0 |
| $r_{2}$ | $a(-2+a)$ |
| $r_{3}$ | $a(a+1)(-2+a)$ |
| $r_{4}$ | $a^{2}(a+1)(-2+a)$ |
| $r_{5}$ | $a^{2}(-2+a)\left(a^{2}+a-1\right)$ |
| $r_{6}$ | $a(a-1)(-2+a)(a+1)\left(a^{2}+a-1\right)$ |
| $r_{7}$ | $a(a-1)(-2+a)(a+1)\left(a^{3}+a^{2}-2 a-1\right)$ |
| $r_{8}$ | $a^{2}(-2+a)\left(a^{2}-2\right)\left(a^{3}+a^{2}-2 a-1\right)$ |
| $r_{9}$ | $a^{2}(a+1)(-2+a)\left(a^{2}-2\right)\left(a^{3}-3 a+1\right)$ |
| $r_{10}$ | $a(a+1)(-2+a)\left(a^{2}+a-1\right)\left(-a-1+a^{2}\right)\left(a^{3}-3 a+1\right)$ |
| $r_{11}$ | $a(-2+a)\left(a^{2}+a-1\right)\left(-a-1+a^{2}\right)\left(a^{5}+a^{4}-4 a^{3}-3 a^{2}+3 a+1\right)$ |
| $r_{12}$ | $a^{2}(a-1)(-2+a)(a+1)\left(a^{2}-3\right)\left(a^{5}+a^{4}-4 a^{3}-3 a^{2}+3 a+1\right)$ |
| $r_{13}$ | $a^{2}(a-1)(-2+a)(a+1)\left(a^{2}-3\right)\left(a^{6}+a^{5}-5 a^{4}-4 a^{3}+6 a^{2}+3 a-1\right)$ |
| $r_{14}$ | $a(-2+a)\left(a^{3}+a^{2}-2 a-1\right)\left(a^{3}-a^{2}-2 a+1\right)\left(a^{6}+a^{5}-5 a^{4}-4 a^{3}+6 a^{2}+3 a-1\right)$ |

Table 4.4 The remainder polynomials $r_{n}(a)$, factored over the integers
of $a_{n, i}$.
It is easy to show by induction that $f_{\lfloor(n-1) / 2\rfloor}(a)$ divides $r_{n}(a)$ for all integers $n \geq 0$. (See table 4.4.)

We have

$$
\alpha_{2 m, i}=\left\{\begin{aligned}
(-1)^{K}\binom{m+k}{2 k}, & i=2 K+1 \\
(-1)^{K}\binom{m+k-1}{2 k}, & i=2 K
\end{aligned}\right.
$$

and an analogous formula for $\alpha_{2 m+1, i}$.
We give a surprising formula for the discriminant of the polynomials $f_{n}(a)$ :

Theorem 4.5 If $n \geq 1$, the discriminant of $f_{n}, \Delta\left(f_{n}\right)$, is $(2 n+1)^{n-1}$.

Proof: To avoid an otherwise oppressive notation, we will first assume that $n$ is divisible by 4 , and write $n=2 \mathrm{~m}$. We indicate at the conclusion of the proof the minor changes needed when $n$ is not divisible by 4 . Our proof is in three parts. In part one, we show that $(2 n+1)^{n-1} \mid \Delta\left(f_{n}\right)$. In part two, we show that $\left|\Delta\left(f_{n}\right)\right| \leq(2 n+1)^{n-1}$. Finally, in part three, we show that $\Delta\left(f_{n}\right)$ has positive sign.

Part one: Our convention is that the Sylvester matrix of the resultant of two polynomials $f, g$ looks like

$$
\left(\begin{array}{c}
\text { Coefficients of } f \\
\text { Coefficients of } f \\
\vdots \\
\text { Coefficients of } g \\
\text { Coefficients of } g \\
\vdots
\end{array}\right)
$$

so, since $f_{2 m}(x)$ is a polynomial of degree $2 m$, the Sylvester matrix $S$ of $\Delta\left(f_{n}\right)$ is:

| $\alpha_{1}$ | $\alpha_{2}$ | $\alpha_{3}$ | $\cdots$ | $\alpha_{n-1}$ | $\alpha_{n}$ | $\alpha_{n+1}$ | 0 | - |  |
| :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: |
|  | $\alpha_{1}$ | $\alpha_{2}$ | $\cdots$ | $\alpha_{n-2}$ | $\alpha_{n-1}$ | $\alpha_{n}$ | $\alpha_{n+1}$ | 0 | $\ldots$ |
|  |  |  |  | $\alpha_{1}$ | $\alpha_{2}$ | $\alpha_{3}$ | ... | .. | $\alpha_{n+1}$ |
| $n \alpha_{1}$ | $(n-1) \alpha_{2}$$n \alpha_{1}$ | $(n-2) \alpha_{3}$ | $\cdots$ | $2 \alpha_{n-1}$ | $\alpha_{n}$ | $0 \alpha_{n+1}$ | $\cdots$ |  |  |
|  |  | $(n-1) \alpha_{2}$ | $(n-2) \alpha_{3}$ | $\ldots$ |  |  |  |  |  |
|  |  |  | : |  | $(n-1) \alpha_{2}$ | $\cdots$ | : |  |  |
|  |  |  |  | 0 | $n \alpha_{1}$ | $(n-1) \alpha_{2}$ | $\cdots$ |  | $\alpha_{n}$ |

Transpose the rows of $S$, so that each of the first $2 m-1$ rows coming from the coefficients of $f_{n}^{\prime}(x)$ are followed by the corresponding row from $f_{n}(x)$

$$
S^{\prime}=\left(\begin{array}{cccc}
2 m\binom{m}{m} & -(2 m-1)\binom{m}{m-1} & -(2 m-2)\binom{m-1}{m-1} & (2 m-3)\binom{m-1}{m-2}  \tag{4.2}\\
\binom{m}{m} & -\binom{m}{m-1} & -\binom{m-1}{m-1} & \binom{m-1}{m-2} \\
& 2 m\binom{m}{m} & -(2 m-1)\binom{m}{m-1} & -(2 m-2)\binom{m-1}{m-1} \\
& \binom{m}{m} & \binom{m}{m-1} & -\binom{m-1}{m-1} \\
& 2 m\binom{m}{m} & -(2 m-1)\binom{m}{m-1} & \cdots \\
& & \binom{m}{m} & -\binom{m}{m-1} \\
& & \vdots & \cdots
\end{array}\right)
$$

Since we have assumed that $n \equiv 0 \bmod 4, S^{\prime}$ is obtained from $S$ by an even number of row exchanges. So the determinant of this matrix is the determinant of $S$, which is a polynomial in $m$. An irreducible polynomial $p(m)$ divides the polynomial $(\operatorname{det} S)(m)$ if and only if each root $r$ of this polynomial $p(m)$, when substituted for $m$ in the above matrix, gives a matrix of determinant zero. If $S(r)$ is a matrix whose row rank is less than $2 n-1$, then there is a non-trivial linear combination of the rows which is the zero vector. Let us write down such a linear combination. There
are ( $n-1$ ) triples of consecutive rows which look like

$$
\begin{array}{ccc}
(2 m)\binom{m}{m} & -(2 m-1)\binom{m}{m-1} & \cdots \\
\binom{m}{m} & -\binom{m}{m-1} & \cdots  \tag{4.3}\\
(2 m+1)\binom{m-1}{m} & \binom{m}{m} & \cdots
\end{array}
$$

We write each of the binomial coefficients as a polynomial as follows: for $k \geq 0$, $\binom{m}{k}$ is identified with the polynomial $\frac{1}{k!} m^{k}$, and $\binom{m+k}{m-l}$ is rewritten as $\binom{m+k}{l-k}$. For $k<0,\binom{m}{k}$ is identified with the polynomial 0 . (The reader should be aware that in general these polynomials yield the appropriate values for the binomial coefficients only when $m$ is a positive integer.) In particular, when we write an expression like $\left.\binom{m+k}{m-l}\right|_{m=a}$ we intend that a be substituted for $m$ in the polynomial $\binom{m+k}{-l-k}$.

We claim that twice the first row in 4.3 plus the second equals the third, when $m=r=-\frac{1}{4}$. This is verified by direct calculation: since, for positive integers $a, b$ :

$$
A\binom{a}{b}-B\binom{a-1}{b-1}=0
$$

if and only if $\mathrm{Aa}-\mathrm{Bb}=0$, we have

$$
\left.\left(2(2 m-2 i)\binom{m+i}{m-i}+\binom{m+i}{m-i}\right)\right|_{m=-\frac{1}{4}}=\left.(2 m-2 i+1)\binom{m+i-1}{m-i-1}\right|_{m=-\frac{1}{4}}
$$

for each $i$ if and only if

$$
\left.(4 m-4 i+1)(m+i)\right|_{m=-\frac{1}{4}}=\left.(2 i)(2 m-2 i+1)\right|_{m=-\frac{1}{4}}
$$

which is clearly true. So

$$
2(2 m-k+1) \alpha_{k}+\alpha_{k}=(2 m-k+2) \alpha_{k-1}
$$

for even $k$; and similarly for odd $k$. Thus, $\operatorname{det}(S)(m)$ evaluated at $m=-\frac{1}{4}$ has determinant zero, since adding twice the first row and ( -1 ) times the third row to the second row yields a row of zeroes. But likewise, twice the third row plus the fourth equals the fifth, etc. There are ( $n-1$ ) rows of zeros in this new matrix; so, $m=-\frac{1}{4}$ is a root of multiplicity at least $(n-1)$ :

$$
\operatorname{det} S=C\left(m+\frac{1}{4}\right)^{K}
$$

for some $K \geq 2 m-1$. In the next part, we will show that $K \leq 2 m-1$. (In particular, this means that the determinant is nonzero.) Our method of proof will give us as a bonus that $C \mid 4^{k}$, so that $4^{k}=C$.

Part Two: Let $r_{1} \geq r_{2} \geq \ldots \geq r_{n}$ be the roots of $f_{n}$. By definition, $\Delta\left(f_{n}\right)=$ $\prod_{i \neq j}\left(r_{i}-r_{j}\right)$. Since the geometric mean of a set of positive numbers is less the arith-
metic mean when the numbers are not all equal, we have that

$$
\begin{aligned}
\sqrt[(n-1) n]{\prod_{i \neq j}\left|\left(r_{i}-r_{j}\right)\right|} & <\frac{1}{n(n-1)} \sum_{i \neq j}\left|r_{i}-r_{j}\right| \\
& \leq \frac{1}{n} \sum_{i}\left|r_{i}\right|
\end{aligned}
$$

i.e.

$$
\begin{equation*}
\left|\Delta\left(f_{n}\right)\right|<\left[\left(\frac{1}{n} \sum_{i}\left|r_{i}\right|\right)^{n}\right]^{n-1} \tag{4.4}
\end{equation*}
$$

By Graeffe's method from the theory of symmetric functions, $\left|T_{1}\right| \doteq\left|\frac{\alpha_{1}}{1}\right|,\left|\tau_{2}\right| \doteq$ $\left|\frac{\alpha_{2}}{\alpha_{1}}\right|$, etc. In particular,

$$
\begin{aligned}
\sum\left|r_{i}\right| & \leq 2 \sum\left|\frac{\alpha_{i}}{\alpha_{i-1}}\right| \\
& =2 \sum \frac{\binom{m+a}{m-a}}{\binom{m+a-1}{m-a}}+2 \sum \frac{\binom{m+a}{m-a+1}}{\binom{m+a}{m-a}} \\
& =2 \sum \frac{m+a}{2 a}+2 \sum \frac{2 a}{m-a+1} \\
& \leq 2 \int_{1}^{m} \frac{m+a}{2 a}+2 \int_{0}^{m} \frac{2 a}{m-a+1} \\
& =(1+\ln m) m-1+(4 \ln (m+1)-4) m+4 \ln (m+1) \\
& <\sqrt[n]{2} \sqrt{2 n+1} \quad(\text { for } n>0 .)
\end{aligned}
$$

Plugging this last inequality into 4.4 gives us that $\left|\Delta\left(f_{n}\right)\right|<2(2 n+1)^{n-1}$. Since the discriminant of $f_{n}$ is clearly an integer, we have $\left|\Delta\left(f_{n}\right)\right| \leq(2 n+1)^{n-1}$.

Part Three: We define a new matrix, $S^{\prime \prime}$, obtained by adding twice each odd row of $S^{\prime}$ to any row immediately below it, subtracting each odd row from any row immediately above it, and multiplying each even row by $\frac{1}{(4 m+1)} . S^{\prime \prime}$ looks like

$$
S^{\prime \prime}=\left(\begin{array}{ccc}
2 m\binom{m}{m} & -(2 m-1)\binom{m}{m-1} & \cdots \\
\frac{(4 m+1)\binom{m}{m}}{(4 m+1)} & \frac{\left(-(4 m-1)\binom{m-1}{m}-2 m\binom{m}{m}\right)}{(4 m+1)} & \cdots \\
0 & 2 m\binom{m}{m} & \cdots \\
0 & \vdots
\end{array}\right)
$$

Clearly $\operatorname{det} S=\operatorname{det} S^{\prime}=(4 m+1)^{2 m-1} \operatorname{det} S^{\prime \prime}$. It is our task to show that $\operatorname{det} S^{\prime \prime}=1$. In light of parts one and two above, we know that $\operatorname{det} S^{\prime \prime}=-1,0$, or 1 . Thus, it suffices to show that $\operatorname{det} S^{\prime \prime}>0$.

Let $M$ be a connected $k \times k$ minor of $S^{\prime \prime}$, which does not have a zero on the main diagonal. (As the name implies, a connected minor of a matrix $M$ is a minor which is obtained by deleting rows and columns only at the beginning and end of $M$.)

For a matrix $M$, the notation $\left.M\right|_{i} ^{j}$ denotes the minor obtained by deleting the $i-$ th row and the $j-$ th column of matrix $M$. There is a recurrence, originally popularized by C. Dodgson, which relates the determinants of a square matrix with the determinants of its connected minors. Dodgson's determinant identity, [6], [38], states
that

$$
\begin{equation*}
\operatorname{det}\left(\left.\left(\left.M\right|_{n} ^{n}\right)\right|_{1} ^{1}\right)(\operatorname{det} M)=\left(\left.\operatorname{det} M\right|_{I} ^{1}\right)\left(\left.\operatorname{det} M\right|_{n} ^{n}\right)-\left(\left.\operatorname{det} M\right|_{n} ^{1}\right)\left(\left.\operatorname{det} M\right|_{I} ^{\Omega}\right) \tag{4.5}
\end{equation*}
$$

The recurrence 4.5 is non-linear. Clearly, if we specify values for the determinants of connected one and two-dimensional minors of $M$, the recurrence 4.5 gives a unique value for $\operatorname{det} M$ if and only if there are no "interior" connected minors with zero determinant.

Let $z=\frac{k+1}{4}$. Let $M$ be any $l \times l$ connected minor of $M$, for $l \leq k$. We use 4.5 to show that the $k \times k$ minor $M$ satisfies the following three properties:

P1 If $k>2$, then $\operatorname{det} M$ is an integer-valued polynomial (a polynomial which sends integers to integers.)

P2 Suppose $l=2$; i.e. $M=\left(\begin{array}{ll}a(m) & b(m) \\ c(m) & d(m)\end{array}\right)$ is a $2 \times 2$ minor of $M$. Then

$$
|a(z) d(z)|>|b(z) c(z)| .
$$

P3 $|\operatorname{det} M|(m)$ is increasing for all $m \geq z$.

Property (P1) is an immediate consequence of part one of the proof. We will prove properties (P2) and (P3) using the identity 4.5. We will first note that they are true for 1-dimensional and 2-dimensional connected minors of the matrix $S^{\prime \prime}$.

$$
\left.\begin{array}{c}
\left(\begin{array}{lllllll}
+ & - & - & + & + & - & - \\
+ & - & - & + & + & - & - \\
+ \\
0 & + & - & - & + & + & - \\
0 & + & - & - & + & + & - \\
0 & - & + & - & - & + & + \\
0 & - \\
0 & 0 & + & - & - & + & + \\
\hline & 0 & 0 & + & - & - & + \\
\hline
\end{array}\right. \\
0
\end{array}\right)
$$

Figure 4-1 The "sign condensation" of S"

The $2 \times 2$ connected minors of $S^{\prime}$ look like one of

$$
\left(\begin{array}{cc}
\binom{m+i}{m-1} & \pm\binom{ m+i}{m-i-1} \\
\mp(2 m+C+1)\binom{m+i-1}{m-i} & (2 m+C)\binom{m+i}{m-i}
\end{array}\right)
$$

whose determinant is

$$
\frac{(m+i)\left(2 m^{2}+\text { lower order terms }\right) m!^{2}}{(2 i)!^{2}(m-i)!^{2}}
$$

$$
\left(\begin{array}{cc}
\binom{m+i}{m-i-1} & \pm\binom{ m+i+1}{m-i-1} \\
\mp(2 m+C+1)\binom{m+i}{m-i} & (2 m+C)\binom{m+i}{m-i-1}
\end{array}\right)
$$

whose determinant is

$$
\begin{aligned}
& \frac{2(i+1)\left(-8 m^{2}+\text { lower order terms }\right)(m+i)!^{2}}{(2 i+2)!^{2}(m-i-1)!^{2}(m-i)} \\
& \left(\begin{array}{cc}
(2 m+C+1)\binom{m+i-1}{m-i} & \pm(2 m+C)\binom{m+i}{m-i} \\
\binom{m+i-1}{m-i} & \pm\binom{ m+i}{m-i}
\end{array}\right)
\end{aligned}
$$

whose determinant is

$$
\pm\binom{ m+i-1}{m-i}\binom{m+i}{m-i}
$$

or

$$
\left(\begin{array}{cc}
(2 m+C)\binom{m+i}{m-i} & \pm(2 m+C-1)\binom{m+i}{m-i-1} \\
\binom{m+i}{m-i} & \pm\binom{ m+i}{m-i-1}
\end{array}\right)
$$

whose determinant is

$$
\pm\binom{ m+i}{m-i}\binom{m+i}{m-i-1}
$$

Examining the 1- and 2-dimensional connected minors, we see that (P2)-(P3) are satisfied when $l=1,2$.

Since the degree of $\left.\left.M\right|_{1} ^{1} M\right|_{n} ^{n}$ never equals the degree of $\left.\left.M\right|_{n} ^{1} M\right|_{\Omega} ^{n}$, except when the minor M in the interior of $\mathrm{S}^{\prime}$ has a zero on its diagonal, the right-hand side of 4.5 is never identically zero. Thus, the recurrence 4.5 has unique solution, with initial values for the $1 \times 1$ and $2 \times 2$ connected submatrices of $S$. In particular, the properties (P2) - (P3) follow by induction.

Since there are no interior zeros, 4.5 gives the value of $\operatorname{det} M(m)$ for any $m \geq z$. In particular, since $z$ lies to the right of all the zeros and poles of $\operatorname{det} M(m)$, by (P3), the sign of $\operatorname{det} \mathrm{M}(z)$ is determined by the signs of the entries of $\mathrm{M}(z)$. (For example,
the case when $k=8$ is shown in table 4-1.) Clearly, when $M$ lies in the north-east corner of $S^{\prime \prime}$, and $k$ is odd, we have that the sign of $M(z)$ is positive. But when $M=S^{\prime \prime}$, then $M(z)=\Delta\left(f_{n}\right)$ by construction. Thus $\Delta\left(f_{n}\right) \geq 0$. Combining parts one, two, and three, we now have that

$$
\Delta\left(f_{n}\right)=(2 n+1)^{n-1}
$$

when $n \equiv 0 \bmod 4$.
The case $n \equiv 2 \bmod 4$ is the same as the case $n \equiv 0 \bmod 4$, except that the transformation of $S$ to $S^{\prime}$ in 4.2 on page 73 changes the sign of the determinant of $S$. The determinants of $k \times k$ connected minors of $S^{\prime}$ now alternate in sign, giving the desired positive determinant $(2 n+1)^{n-1}$ for $S$. The case $n \equiv 1 \bmod 4$ is the same as the case $n \equiv 0 \bmod 4$, except that the first two rows of the matrix $S^{\prime}$ now looks like

$$
\left(\begin{array}{cccc}
(2 m+1)\binom{m}{m} & (2 m)\binom{m+1}{m} & -(2 m-1)\binom{m+1}{m-1} & \cdots \\
\binom{m}{m} & \binom{m}{m-1} & -\binom{m-1}{m-1} & \cdots
\end{array}\right)
$$

where $(2 m+1)=n$. The case $n \equiv 3 \bmod 4$ is similar to the case $n \equiv 2 \bmod 4 . \diamond$
Two invaluable resources for determinant identities are the surveys [99] and [7]. For the reader's convenience, we illustrate a "condensation" using the recurrence 4.5 from the second half of the proof, for the Sylvester matrix of polynomial $f_{4}(x)$ in
figure 4-2 on page 92. We note that

- Condition (P3) is important, because if there are interior zeros at any step, the difference equation 4.5 , together with the $(2 n-1)^{2}+(2 n-2)^{2}$ initial values at level $l=1$ and $l=2$, may not have unique solution.
- This is quite an unusual application of Dodgson's recurrence. More typically, the connected $n \times n$ minors of a matrix $M(m)$ are imbedded in a family of matrices $M_{n}(a, b)$ - and the recurrence 4.5 becomes an integral recurrence relation with variables $a, b, M_{n}$. Proving a determinant identity is now a matter of showing that the family of matrices $M_{n}(a, b)$ satisfies the recurrence and that an adequate set of initial conditions are satisfied. By using 4.5 instead essentially to bound the degree of $(2 \pi+1)$ in $(\operatorname{det} S)(m)$, we avoided the (usually very hard) problem of finding a useful parameterization of the minors of $S$.
- Of course, one might try to find a more "mechanical" method of identifying the factors of the determinant of a matrix (part one of the proof.) Here is a brief description of the process which led to the discovery of theorem 4.5. We may write the first two rows of the matrix $S^{\prime}$ as:

$$
\begin{array}{ccc}
-1\binom{m}{m} & 2\binom{m}{m-1} & \cdots \\
\binom{m}{m} & -\binom{m}{m-1} & \cdots
\end{array}
$$

without changing the determinant. Clearly, the first non-zero entry of each
even row must cancel the entry immediately above it; so we may assume without loss of generality that the linear combination of rows for which we are searching is obtained by multiplying the matrix $S(r)$ by the vector $\left(1, x_{1}, 1, x_{2}, \ldots, 1, x_{2 m-1}, X\right)^{T}$. Since the result is the zero vector, the sequence $\left[x_{i}\right]_{i}$ satisfies, for each $1 \leq N<2 m$

$$
\left(x_{N}-1\right)+\left(x_{N-1}-2\right) \alpha_{N-1}+\cdots+\left(x_{1}-N\right) \alpha_{N}=0
$$

$$
x_{N}=\sum_{l=1}^{N-1}\left(x_{l}-N+l-1\right) \alpha_{N-l+1}=\sum_{l=1}^{N-1} x_{l} \alpha_{N-l+l}+\sum_{l=1}^{N-1}(N-l+1) \alpha_{N-l+1}
$$

i.e.

$$
\begin{equation*}
\sum_{l=1}^{N} x_{l} \alpha_{N-l+1}=\sum_{k=1}^{N} k \alpha_{k} \tag{4.6}
\end{equation*}
$$

In [79] the "generalized binomial series" $\mathrm{B}_{\mathrm{t}}$ is introduced. It satisfies, and is well-defined by, the relation

$$
\frac{B_{t}(z)^{R}}{1-t+t B_{t}(z)^{-1}}=\sum_{k \geq 0}\binom{t k+R}{k} z^{k}
$$

When $t=1$, this is the binomial series; when $t=2$, we have the generating
function for the Catalan numbers. More generally, we have

$$
B_{t}(z)=\sum_{k \geq 0}\binom{t k+1}{k} \frac{1}{t k+1} z^{k}
$$

Consider the function

$$
\begin{aligned}
\mathrm{G}_{1}(z) & =\sum_{k \geq 0}\binom{m+k}{2 k} z^{k} \\
& =\frac{2 \mathrm{~B}_{1 / 2}(z)^{[m+1]}}{\mathrm{B}_{1 / 2}(z)+1} .
\end{aligned}
$$

(The second equality is the result of the identity $B_{1 / 2}(z)=\frac{1}{B_{1 / 2}(-z)}$.) The even terms of $\mathrm{G}_{\mathrm{l}}(z)$ may be extracted as

$$
E(z)=\frac{B_{1 / 2}(z)^{[m+1]}+B_{1 / 2}(z)^{[-m]}}{B_{1 / 2}(z)+1}
$$

Likewise the odd terms of the function

$$
\begin{aligned}
G_{2}(z) & =\sum_{k \geq 0}\binom{m+k-1}{2 k} z^{k} \\
& =\frac{2 \mathrm{~B}_{1 / 2}(z)^{\left[m+\frac{1}{2}\right]}}{\mathrm{B}_{1 / 2}(z)+1}
\end{aligned}
$$

are

$$
O(z)=\frac{B_{1 / 2}(z)^{\left[m+\frac{1}{2}\right]}-B_{1 / 2}(z)^{\left[\frac{1}{2}-m\right]}}{B_{1 / 2}(z)+1}
$$

Thus,

$$
\begin{aligned}
\sum_{k}\left|\alpha_{k+1}\right| z^{k}= & \sum_{k}\binom{m+\left\lfloor\frac{k+1}{2}\right\rfloor}{ m-\left\lfloor\frac{k+2}{2}\right\rfloor} z^{k} \\
= & E(z)+O(z) \\
= & \frac{1}{B_{1 / 2}(z)+1}\left(B_{1 / 2}(z)^{m+1}+B_{1 / 2}(z)^{-m}+B_{1 / 2}(z)^{m+\frac{1}{2}}\right. \\
& \left.\quad-B_{1 / 2}(z)^{\frac{1}{2}-m}\right)
\end{aligned}
$$

and

$$
\begin{aligned}
G(z)= & \sum_{k \geq 0} \alpha_{k} z^{k} \\
= & \frac{1}{z}(E(i z)+i O(i z)) \\
= & \frac{1}{z}\left(\left(\frac{i z+\sqrt{4-z^{2}}}{2}\right)^{2}+1\right)^{-1} \cdot \\
& \left(\left(\frac{i z+\sqrt{4-z^{2}}}{2}\right)^{2 m+2}+\left(\frac{i z+\sqrt{4-z^{2}}}{2}\right)^{-2 m}\right. \\
& \left.\left(\frac{i z+\sqrt{4-z^{2}}}{2}\right)^{2 m+1}+i\left(\frac{i z+\sqrt{4-z^{2}}}{2}\right)^{1-2 m}\right) .
\end{aligned}
$$

Now the left-hand side of 4.6 is the convolution of the sequences $\left[x_{k}\right]_{k}$ and $\left[\alpha_{k-1}\right]_{k}$; the right-hand side of 4.6 is the convolution of the sequence $[1]_{k}$ and the derivative sequence of $\left[\alpha_{k}\right]_{k}$. So, writing 4.6 in terms of generating functions,
we get

$$
H(z)(z G(z))=\frac{1}{1-z} G^{\prime}(z)
$$

$$
H(z)=\frac{G^{\prime}(z)}{G(z)} \frac{1}{z-z^{2}}
$$

where $H(z)$ is the generating function for $\left[x_{k}\right]_{k}$.
For $M-l<l$, we must have

$$
\begin{equation*}
-\sum\left(x_{k}+2 m-M\right) \alpha_{M-k}=X(2 m-M+1) \alpha_{M} \tag{4.7}
\end{equation*}
$$

so that $\frac{\sum\left(x_{k}+2 m-M\right) \alpha_{M-k}}{(2 m-M+1) \alpha_{M}}$ must be constant. Writing the cases $M=0, M=1$ in terms of generating functions in variable $z$, and substituting $z=0$, we see by inspecting special cases that this happens only when $m=-\frac{1}{4} . S$, an integral matrix, clearly has an integral determinant, so each of the equations 4.7 must be satisfied for some rational $m \in \mathbb{C}$. We have discovered that ( $m+\frac{1}{4}$ ) is a linear factor (and, apparently, the only linear factor) of $\operatorname{det} S$. The force of theorem 4.5 is that S has no other linear factors.

- At the risk of diverging even further from the topic of this thesis, it is reasonable to ask: if we have an $n \times n$ matrix $M$, whose ( $\boldsymbol{i}-\boldsymbol{j}$ )-th entry is some function
of $i, j$, and $n$, how could we hope to find a matrix $N$, with the same determinant as $M$, so that the recurrence 4.5 for $\operatorname{det} N$ has unique solution for a given set of initial values? I.e., if we have a determinant identity which we wish to prove, how could we guess a transformation of $M$ which yields a matrix for which Dodgson's determinant recurrence is useful?

One possibility, if you are a graduate student with a lot of free time on your hands, is to try every determinant-preserving matrix transformation that you can think of until you find something that works. We very briefly sketch a more systematic plan of attack here.

Suppose that we have an integral matrix $M$, and that we have successfully used the recurrence

$$
(\operatorname{det} M)=\frac{\left(\left.\operatorname{det} M\right|_{1} ^{1}\right)\left(\left.\operatorname{det} M\right|_{n} ^{n}\right)-\left(\left.\operatorname{det} M\right|_{n} ^{1}\right)\left(\left.\operatorname{det} M\right|_{1} ^{n}\right)}{\left(\left.\left.\operatorname{det} M\right|_{n} ^{n}\right|_{\mathfrak{1}} ^{1}\right)}
$$

to evaluate $\operatorname{det} M$ (so that there are no interior zeros.) Then each connected minor was integral, and if we kept track of the determinants of the connected minors, then we could write down the Smith normal form of $M$, and the Hermite normal for of $M$. (This is actually a practical algorithm for finding the Smith and Hermite normal forms of integral matrices. The advantages of this method are its easy scalability to parallel machines, and the fact that the integers in intermediate calculations don't grow very fast as $n \rightarrow \infty$. The disadvantage of this method is that, if an interior zero is found, then one must start over,
applying an elementary row or column operation to eliminate the zero. For a discussion of the technique, see [25] or [8].) On the other hand, a transformation of $M$ to a matrix $M^{\prime}$ which changes the sort of possible parameterizations $M(a, b)$ must not respect the Smith normal form of $M$ - otherwise, we could have gone from $M$ to $M^{\prime}$ with an $\mathrm{SL}_{2} \mathbb{Z}$ transformation.

If the determinant of an $n \times n$ matrix has prime factorization $p_{1}^{\mathbf{c}_{1}} p_{2}^{\mathbf{c}_{2}} \cdots p_{k}^{\mathbf{e}_{k}}$, then there are

$$
P(M)=P_{n}\left(e_{1}\right) P_{n}\left(e_{2}\right) \cdots P_{n}\left(e_{k}\right)
$$

possible choices for the invariant factors of $M$ (where $P_{n}(j)$ is the number of $j$-tableaux shape with $n$ or fewer rows - i.e. the number of partitions of $j$ into $n$ or fewer parts.) Recall that a Smith Matrix is a matrix in Smith normal form: a diagonal matrix s.t. $a_{i i} \mid a_{(i+1),(i+1)}$ for each diagonal entry $a_{(i+1),(i+1)}, i>0$. There are clearly $P(M)$ Smith matrices with the same determinant as $M$. Each such matrix $\bar{M}$ which is $\mathrm{GL}_{2} \mathbb{Q}$ similar to $M$, but only one them is $\mathrm{SL}_{2} \mathbb{Z}$ similar to $M$. Writing down each of these matrices, and the $\mathrm{GL}_{2} \mathbb{Q}$ matrices which turn it into $M$, is purely mechanical. If $M_{n}(m)$ is family of square matrices, each with parameter $m$ and dimension $n$, then we get a family of possible transforms of $M_{n}(m)$ amenable to analysis with Dodgson's recurrence 4.5 by choosing some $n^{\prime}, m^{\prime}$, and finding the second-order recurrence relation which is satisfied by each $\overline{M_{n^{\prime}}\left(m^{\prime}\right)}, \overline{M_{n^{\prime}+1}\left(m^{\prime}\right)}$, and $\overline{M_{n^{\prime}+2}\left(m^{\prime}\right)}$. Repeating the process with another
$m^{\prime \prime}$ gives another family of transforms of $M_{n}(m)$ - and any possible transform of $M_{n}(m)$ at these values of $n$ must work for both $m^{\prime}$ and $m^{\prime \prime}$.

Here is a contrived but illustrative example. Consider the $n \times n$ matrix $\mathcal{A}_{n}$ defined by

$$
a_{i j}=\left\{\begin{array}{cc}
0, & i \neq j \\
1, & i=j<n \\
\prod_{1 \leq p<q \leq n}(q-p), & i=j=n
\end{array}\right.
$$

The first few values of $\left|A_{n}\right|$ are $1,1,2,12=2^{2} 3,288=2^{5} 3^{2}, 34560=2^{8} 3^{3} 5$. $A_{n}$ is already in Smith normal form. For $i<j<n$, divide the term ( $j-i$ ) from the $n$-th row, and multiply the $j$-th row by ( $j-i$ ). We have $G L_{n}$-transformed $A_{n}$ into an integral matrix with the same determinant, but with the following Smith normal form:

By judiciously adding rows and columns, it is easy to $S L_{n}$-transform $V^{\prime}$ into
the matrix

$$
V_{n=}\left(\begin{array}{cccc}
1 & 1 & 1 & \ldots \\
1 & 2 & 3 & \\
1^{2} & 2^{2} & 3^{2} & \\
\vdots & & & \ddots
\end{array}\right)
$$

which is a Vandermonde matrix whose determinant is $\prod_{1 \leq p<q \leq n}(q-p)$. So we have discovered an unsurprising formula for the determinant of $A_{n}$. (We have also discovered an obscure fact about the Smith normal form of these Vandermonde matrices.)

- Returning from this long digression to the matter at hand, one might ask what the Smith normal forms of the Sylvester matrices of the $f_{n}^{\prime} s$ look like. Actually tracing the divisors of $\left(\operatorname{det} \tilde{S}^{\prime}\right)(m)$ for connected minors $\tilde{S^{\prime}}$ of $S^{\prime}$ yields:

Theorem 4.6 Let $S$ be the Smith normal form of the Sylvester matrix for $\Delta\left(f_{n}\right)=$ $\operatorname{Res}\left(f_{n}, f_{n}^{\prime}\right)$. Then $S$ is a diagonal matrix, consisting of " 1 "s in the first $n$ diagonal positions, and " $(2 n-1)$ "'s in the remaining $n-1$ positions.

$$
\begin{aligned}
& S=\left(\begin{array}{rrrrrrr}
1 & -2 & -3 & 1 & 1 & 0 & 0 \\
4 & -6 & -6 & 1 & 0 & 0 & 0 \\
0 & 0 & 1 & -2 & -3 & 1 & 1 \\
0 & 1 & -2 & -3 & 1 & 1 & 0 \\
0 & 4 & -6 & -6 & 1 & 0 & 0 \\
0 & 0 & 4 & -6 & -6 & 1 & 0 \\
0 & 0 & 0 & 4 & -6 & -6 & 1
\end{array}\right) \\
& \longrightarrow\left(\begin{array}{cccccc}
-(2) & (2)(3) & -(3) & 1 & 0 & 0 \\
(2)^{2} & (2)^{3}(3) & (2)^{3}(3) & (7) & 0 & 0 \\
0 & -(2) & (2)(3) & -(3) & 1 & 0 \\
0 & (2)^{2} & (2)^{3}(3) & (2)^{3}(3) & (7) & 0 \\
0 & 0 & -(2) & (2)(3) & -(3) & 1 \\
0 & 0 & (2)^{2} & (2)^{3}(3) & (2)^{3}(3) & (7)
\end{array}\right) \\
& \rightarrow\left(\begin{array}{ccccc}
(2)^{2}(3)^{2} & -(2)^{3}(3)^{2} & -(3)^{2}(5) & 0 & 0 \\
-(2) & -(2)^{5} & (19) & (7) & 0 \\
0 & (2)^{2}(3)^{2} & -(2)^{3}(3)^{2} & -(3)^{2}(5) & 0 \\
0 & -(2) & -(2)^{5} & (19) & (7) \\
0 & 0 & (2)^{2}(3)^{2} & -(2)^{3}(3)^{2} & -(3)^{2}(5)
\end{array}\right) \\
& \rightarrow\left(\begin{array}{cccc}
-(2)(3)^{3} & -(3)^{2}(13) & -(3)^{2}(5) & 0 \\
(2)^{2}(3)^{2} & (2)(3)^{3}(5) & (3)^{2}(13) & 0 \\
0 & -(2)(3)^{3} & -(3)^{2}(13) & -(3)^{2}(5) \\
0 & (2)^{2}(3)^{2} & (2)(3)^{3}(5) & (3)^{2}(13)
\end{array}\right) \\
& \rightarrow\left(\begin{array}{ccc}
(2)^{2}(3)^{4} & -(3)^{4} & 0 \\
-(2)(3)^{3} & (3)^{3}(13) & (3)^{2}(13) \\
0 & (2)^{2}(3)^{4} & -(3)^{4}
\end{array}\right) \rightarrow\left(\begin{array}{cc}
(3)^{4}(5) & -(3)^{4} \\
(2)^{2}(3)^{4} & (3)^{4}(7)
\end{array}\right) \\
& \rightarrow\left((3)^{6}\right)
\end{aligned}
$$

Figure 4-2 Dodgson condensation for the Sylvester matrix of $f_{4}(x)$.

## Appendix: Program listing

## Some routines for simplifying group presentations

## To accompany Chapter 4.

## 

## 

## 

## 

## 

## 

## 

NUMPAIRSCHKTzWHIT := 20;
NUMTRIPLESCHKTzWHIT := 20;
emptyVectors := [ [0], [0,0,0], [0,0,0,0,0,0,0] ];
RememberMonabc := [];

```
```

RememberMonabc[3] := []; RememberMonabc[3][1] := \square];
RememberMonabc[3] [1] [2] := []; RememberMonabc[3] [1] [2][3] := 7;

## 

## setRememberMonabc sets up the n'th row of the RememberMonabc table for

function Monabc
setRememberMonabc := function( n )
local count, i,j,k;
if not IsBound( RememberMonabc[n] ) then
count := n + Binomial( n, 2);
RememberMonabc[n] := [];
for i in [1..n] do
RememberMonabc[n][i] := [];
for j in [(i+1)..n] do
RememberMonabc[n][i][j] := [];
for k in [(j+1)..n] do
count := count + 1;
RememberMonabc[n][i][j][k] := count;
od;
od;
od;
fi;
end;

## 

## MonomialVector returns a vector of zeroes, whose length is

## n + C(n, 2) + C(n,3)

MonomialVector := function( n )
local l, i, c;
if IsBound( emptyVectors[n] ) then
return ShallowCopy( emptyVectors[n] );
else
c := (n + Binomial(n,2) +Binomial(n,3));
l := [];
I[c] := 0;

```
```

        for i in [1..c] do
            I[i] := 0;
        od;
        emptyVectors[n] := ShallowCopy( I );
        return l;
    fi;
    end;

## set up the remember tables

for i in [4..45] do
MonomialVector( i );
setRememberMonabc( i );
od;

## 

## Sorted returns the sorted version of its single argument, Without

## changing the argument.

Sorted := function( 1 )
local L;
L := ShallowCopy( I );
Sort( L );
return L;
end;

## 

## ourRelatorRepresentatives cyclically reduces a list of relators

ourRelatorRepresentatives := function( 1 )
if I = [IdWord] then
return [IdWord];
else
return RelatorRepresentatives( I );
fi;
end;

```
```


## 

## positionInSet returns the position of l or 1^-1 in L,

## where L is a sorted list without holes of length len.

## 

## No checking is done to see whether l<> [ and length(l) = 1.

positionInSet := function( I,L, len )
local left, right, middle, tmp;
left := 1; right := len;
while left + 1 < right do
middle := QuoInt( right-left, 2 ) + left;
tmp := I[middle];
if tmp < l then
left := middle;
elif tmp > l^-1 then \# By convention, l < l^-1
right := middle;
else
left := middle;
right := middle;
fi;
od;
if L[left] < l then
return right;
else
return left;
fi;
end;

## 

## positionInList returns the position of the first occurence of

## l or 1"-1 in L

## L need not be sorted, and may have holes.

positionInList := function( 1,L )

```
```

    return Minimum( Position( L, I ),
    Position( L, 1^-1 ) );
    end;

```
```


## 

```
##
## LettersInWord returns a list of the letters in the word, in order
## LettersInWord returns a list of the letters in the word, in order
## of first appearance.
## of first appearance.
LettersInWord := function( w, varList )
LettersInWord := function( w, varList )
local W, thisLetter, I;
local W, thisLetter, I;
    W := Copy( w ); l := \square;
    W := Copy( w ); l := \square;
    while (W <> IdWord) do
    while (W <> IdWord) do
        thisLetter := Subword( W, 1, 1);
        thisLetter := Subword( W, 1, 1);
        if not (thisLetter in varlist) then
        if not (thisLetter in varlist) then
            thisLetter := thisLetter^-1;
            thisLetter := thisLetter^-1;
        fi;
        fi;
        Add( l, thisLetter );
        Add( l, thisLetter );
        W := EliminatedWord( W, thisLetter, IdWord );
        W := EliminatedWord( W, thisLetter, IdWord );
    od;
    od;
    return l;
    return l;
end;
end;
##
##
## IsSortedWord returns true if w is sorted lexicographically
## IsSortedWord returns true if w is sorted lexicographically
IsSortedWord := function( w )
IsSortedWord := function( w )
local count, tmp, tmp1, tmp2, len;
local count, tmp, tmp1, tmp2, len;
    Ien := LengthWord( w );
    Ien := LengthWord( w );
    tmp1 := Subword( w, 1, 1);
    tmp1 := Subword( w, 1, 1);
    tmp := true;
    tmp := true;
    count := 2;
    count := 2;
    while (count <= len) and tmp do
    while (count <= len) and tmp do
        tmp2 := Subword( w, count, count);
        tmp2 := Subword( w, count, count);
        tmp := tmp and (tmp1 <= tmp2);
        tmp := tmp and (tmp1 <= tmp2);
        tmp1 := tmp2;
```

        tmp1 := tmp2;
    ```
```

        count := count+1;
        od;
    return tmp;
    end;

## 

## LettersInWordList returns a list of the letters in the words of a

## list, in order of first appearance.

LettersInWordList := function( 1, varlist )
local L, W, thisLetter, li, len, i, j;
L := Copy(I); Ii := []; Ien := Length( 1 );
for i in [1..len] do
W := L[i];
while (W <> IdWord) do
thisLetter := Subword( W, 1, 1);
if not (thisLetter in varlist) then
thisLetter := thisLetter^-1;
fi;
Add( li, thisLetter );
W := EliminatedWord( W, thisLetter, IdWord );
for j in [i..len] do
L[j] := EliminatedWord( L[j], thisLetter, IdWord );
od;
od;
od;
return li;
end;

## 

## Monab( a,b, VARLIST, n) returns the position of the coordinate of

tr_{ab}

## where a, b are variables in VARLIST, a<>B, where n = |VARLIST|

## No type- or bound-checking is performed.

```
```

Monab := function( a, b, VARLIST, n )
local A,B;
A := positionInSet( a, VARLIST, n );
B := positionInSet( b, VARLIST, n );
return A* (n+1) - (A~2 + A)/2 + B - A;
end;

## 

## Monabc( a,b,c, VARLIST, n) returns the position of the coordinate of

## tr_{abc},

## a<b<c

## where a, b are variables in VARLIST, a<>B, where n = |VARLIST|

## No type- or bound-checking is performed.

Monabc := function( a, b, c, VARLIST, n)
local A,B,C;
A := positionInSet( a, VARLIST, n );
B := positionInSet( b, VARLIST, n );
C := positionInSet( c, VARIIST, n );
setRememberMonabc( n );
return RememberMonabc[n] [A] [B] [C];
end;

## 

## PosFirstNonInverse returns the position of the first generator

## in w which is not raised to a negative power.

## W is assumed to be square-free

## 

## If w has only positive exponents, returns false

PosFirstNonInverse := function( w, varlist )
local count;
for count in [1..LengthWord( w )] do
if not (Subword( w, count, count) in varlist) then

```
```

                return count;
            fi;
        od;
        return false;
    end;

## 

## leqDLEX returns true if m1 <= m2 with degree-lex order

## 

## m1, m2 are monomials, represented as a list of integers

## No type-checking is performed.

leqDLEX := function( m1, m2 )
local s1, s2;
s1 := Sum(m1);
s2 := Sum(m2);
return ( (s1 < s2)
or
( (s1 = s2) and m1 <= m2) );
end;

## 

## MaxMonPair returns the largest of m1, m2, with respect

## to leqDLEX order

MaxMonPair := function(m1, m2 )
if leqDLEX(m1, m2) then
return m2;
else
return m1;
fi;
end;

## 

## MaxMonTrip returns the largest of m1, m2,m3

MaxMonTrip := function(m1, m2, m3)
return MaxMonPair( MaxMonPair(m1, m2), m3 );

```
end;
```


## 

## MaxMon returns the largest monomial in its argument, which

## must be a list.

## If the argument is empty, returns "false"

MaxMon := function( 1 )
local max, i;
max := false;
for i in l do
max := MaxMonPair( max, i);
od;
return max;
end;

## 

## Forward reference of function LTFrickeChar1

LTFrickeChar1 := function( w, \# a word
varList, \# the variables
m );
end;

## 

## LTFrickeCharSquareFree returns the lead monomial of a normal form of

## the Fricke character (trace polynomial) of square-free word w,

## multiplied by the monomial m.

## w is assumed to be cyclically reduced.

LTFrickeCharSquareFree := function( w, \# a word
VARLIST, \# the variables
m )
local len, numvars, M, letters, count, ES,
tmp, a, b, c, max, maxPos, W;
len := LengthWord( w );

```
```

M := ShallowCopy(m);
numvars := Length( VARLIST );
if len = 1 then
tmp := positionInSet( Subword( w, 1, 1), VARIIST, len );
M[tmp] := M[tmp] + 1;
return M;
elif len = 0 then
return M;
elif len = 2 then
\# two cases: something like "ab" or something like "ab--1"
\# the leading term of trace(ab) is just t_{ab}
\# the leading term of trace(ab~-1) is t_a * t_b
a := Subword(w,1,1);
b := Subword(}(\pi,2,2)
if (a in VARIIST) and (b in VARLIST) then
tmp := Sorted( [ a,b ] );
tmp := Monab( tmp[1], tmp[2], VARLIST, numvars );
M[tmp] := M[tmp] + 1;
return M;
else
tmp := positionInSet( a, VARLIST, numvars );
M[tmp] := M[tmp] + 1;
tmp := positionInSet( b, VARLIST, numvars );
M[tmp] := M[tmp] + 1;
return M;
fi;
elif len = 3 then
a := Subword(%,1,1);
b := Subword(w,2,2);
c := Subword(ఒ,3,3);
\#4 cases: "abc", or "acb", or "abc^-1", or "acb^-1"
ES := 0;

```
```

if a in VARLIST then
ES := ES + 1;
else
ES := ES -1;
fi;
if b in VARLIST then
ES := ES + 1;
else
ES := ES -1;
fi;
if c in VARLIST then
ES := ES + 1;
else
ES := ES -1;
fi;
if ES < O then
return LTFrickeCharSquareFree( w--1, VARLIST, m);
elif ES = 3 then
if IsSortedWord( w ) then \# "abc" case
tmp := Monabc( a, b, c, VARLIST, numvars);
M[tmp] := M[tmp] + 1;
return M;
else \#"acb" case
tmp := positionInSet( a, VARLIST, numvars );
M[tmp] := M[tmp] + 1;
tmp := positionInSet( b, VARLIST, numvars );
M[tmp] := M[tmp] + 1;
tmp := positionInSet( c, VARLIST, numvars );
M[tmp] := M[tmp] + 1;
return M;
fi;
else \# "abc"-1" ... one letter has exponent -1, the others +1
\#
\# so that IT( trace( abc^-1 ) ) = tr (ab)tr(c)
if ExponentSumWord( w, a) < 0 then
tmp := a; a := c; c := tmp;

```
```

            elif ExponentSumWord( w, b) < 0 then
                tmp := b; b := c; c := tmp;
            fi;
            tmp := Sorted( [a,b ] );
            tmp := Monab( tmp[1], tmp[2], VARLIST, numvars );
            M[tmp] := M[tmp] + 1;
            tIMP := positionInSet( c, VARLIST, numvars );
            M[tmp] := M[tmp] + 1;
            return M;
        fi;
    else \# len >= 4, so there is the possibility of repeated letters
\# get rid of repeated letters
for count in [1..len-1] do
tmp := Subword(w, count, count);
tmp := Minimum(PositionWord( w, tmp, count+1 ),
PositionWord( W, tmp^-1, count+1 ));
if tmp <> false then
a := Subword( w, 1, count);
b := ourRelatorRepresentatives( ReducedRrsWord( [Subword( w,
count+1,
tmp )]
)) [1];
if tmp = len then
c := IdWord;
else
c := Subword( w, tmp+1, len);
fi;
tmp := ourRelatorRepresentatives( ReducedRrsWord( [c*a] ))[1];
return LTFrickeChar1( tmp,
VARLIST,
m)

```
```

                LTFrickeCharSquareFree( b,
                                    VARLIST,
                                    m );
    fi;
    od;
\# if w survives the for loop, then w is square-free,
\# with no repeated letters
tmp := false;
count := 0;
max := Subword( w, 1, 1);
maxPos := 1;
while (tmp = faise) and count < len do
count := count + 1;
a := Subword( w, count, count);
if not (a in VARIIST) then
tmp := count;
fi;
if a < max then
maxPos := count;
max := a;
fi;
od;
if tmp = false then \# all of the exponents of w are 1
if maxPos = len then
W := max*Subword( w, 1, maxPos-1);
elif maxPos > 1 then
W := max*Subword( w, maxPos +1, len)*Subword( w, 1, maxPos-1);
else
W := w;
fi;
M := MonomialVector( numvars );
\# Now, if f = max, and b,c,d are the last letters of W, and a is
everything

```
```

    # between, then by Vogt's identity, the leading term of t_w = t_W is
    one of:
\#
\# t_{fac}*t_b*t_d
\# t_{fa}*t_b*t_{cd}
\# t_{fad}*t_b*t_c
\# since these are the terms of highest degree, when w=(fabcd) is a
word with
\# positive exponents and no repeated letters.
return (
m
+
MaxMonTrip( LTFrickeCharSquareFree( Subword(W,1,len-3)*Subword(W,len-1,len-1),
VARLIST, M)
+
LTFrickeCharSquareFree( Subword(W,Ien-2,len-2),
VARLIST, M)
+
LTFrickeCharSquareFree( Subword(W,Ien,len),
VARLIST, M),

# i.e., t_{fac}*t_b*t_d

                    LTFrickeCharSquareFree( Subword(W,1,1en-3),
                                    VARLIST, M)
                +
                            LTFrickeCharSquareFree( Subword(W,Ien-2,1en-2),
                                    VARLIST,M)
                    LTFrickeCharSquareFree( Subword(W,Ien-1,len-1)*Subword(W,len,len),
                            VARLIST,M),
    
# i.e. t_{fa}*t_b*t_{cd}

    LTFrickeCharSquareFree( Subword(W,1,len-3)*Subword(W,len,len),
                                    VARIIST, M)
        +
            LTFrickeCharSquareFree(Subrord(W,len-2,1en-2),
                                    VARLIST, M
                +
                    ITFrickeCharSquareFree( Subword(W,Ien-1,Ien-1),
                        VARLIST, M) )
    
# i.e. t_{fad}*t_b*t_c

```
    )
```

                    );
    else # finally, we have the case where w is square-free, without
    repeated
\# generators, but with at least one exponent -1 at position tmp
a := Subword(w,tmp,tmp);
return LTFrickeCharSquareFree( EliminatedWord( w,
a,
a^-1),
VARLIST,
m );
fi;
fi;
end; \# LTFrickeCharSquareFree

## 

## LTFrickeChar1 returns the lead monomial of a normal form of

## the Fricke character (trace polynomial) of word w, multiplied

## by the monomial m

## w is assumed to be cyclically reduced.

LTFrickeChar1 := function( w, \# a word
VARLIST, \# the variables
m )
local M, W, c, tmp, squareFreePart, len, numvars;
M := ShallowCopy( m );
squareFreePart := IdWord;
W := Copy( w );
len := LengthWord(W);
numvars := Length( VARLIST );
while (W <> IdWord) do
tmp := Subword(W, 1, 1);
if (len > 1)
and

```
```

                (tmp = Subword(W, 2, 2))
        then
            c := positionInSet( tmp, VARLIST, numvars );
            M[c] := M[c] + 1;
        else
        squareFreePart := squareFreePart*tmp;
        fi;
        W := tmp--1 * W;
        len := len -1;
    od;
    return LTFrickeCharSquareFree( squareFreePart, VARIIST, M );
    end; \# LTFrickeChar1

## 

## ITFrickeChar returns the lead monomial of a normal form of

## the Fricke character (trace polynomial) of word w

LTFrickeChar := function( w, \# a word
varList) \# the variables
return LTFrickeChar1( ourRelatorRepresentatives( ReducedRrsWord([m]) ) [1],
Set(varList),
MonomialVector( Length( varList ) ) );
end;

## 

## nudgeSet increases its argument S by 1, where S, a boolean list,

## is considered to be a binary integer written with least significant

## digit first, with false=0, true=1. In case of overflow, nudgeSet

## changes nothing, and returns "false"; otherwise "true" is returned.

nudgeSet := function( S )
local p, carry, lenS;
IenS := Length( S );
if SizeBlist( S ) = 0 then

```
```

        return false;
    else
        carry := true;
        p := 1;
        while carry and p< lenS do
            if S[p] then
                carry := true;
                S[p] := false;
                p := p + 1;
            else
                S[p] := true;
                carry := false;
                return true;
            fi;
        od;
    fi;
    end;

## 

## DimMonomialIdeal returns the degree of the Hilbert polynomial of

## the ideal generated by the monomials in list L, The ideal <L> is

## assumed to be over a field of characteristic zero.

## We do this by actually looking at each of the subsets of the

## variables, in order to use GAP 3.4 kernel functions in preference

## to library functions.

## As the number of variables grows, then the issue of avoiding library

## functions in favor of internal GAP functions becomes moot, of course.

DimMonomialIdeal := function( L )
local l, S, least, M, nummons, numvars, count, covers, p, truesies,
noOverflow;

```
```

I := Set( L ); \# get rid or duplicates - might not always be useful

```
I := Set( L ); # get rid or duplicates - might not always be useful
if I = [] then
if I = [] then
    return false;
    return false;
fi;
fi;
nummons := Length( 1 );
nummons := Length( 1 );
numvars := Length( 1[1] );
numvars := Length( 1[1] );
least := numvars;
least := numvars;
# for each monomial in l, make a boolean list, with entry "true" in
# for each monomial in l, make a boolean list, with entry "true" in
# the place of each variable used in the monomial
```


# the place of each variable used in the monomial

```
```

    M := [];
    for count in [1..nummons] do
        M[count] := List( l[count],
                        i -> i > 0);
        IsBlist( M[count] );
    od;
    # Search through the subsets of [1..numvars]
    truesies := List( [1..numvars], i -> true );
    IsBlist( truesies );
    S := ShallowCopy( truesies );
    S[1] := false;
    IsBlist( S );
    least := numvars;
    noOverflow := true;
    while noOverfiow do
        covers := true;
        count := 1;
        while covers and count <= nummons do
            # Oddly, it's quicker to use the kernel function SizeBlist, than
            # to find the OR of the elements of S AND B.
            covers := covers
                and
                        SizeBlist( ( IntersectionBlist( S, M[count] ) ) ) > 0;
            count := count + 1;
        od;
        if covers then # found a set of variables in each of the monomials
            least := Minimum( least, SizeBlist( S ) );
            noDverflow := nudgeSet(S );
        else # ignore large sets
            repeat
                noDverflow := mudgeSet( S );
            until (not noDverflow) and SizeBlist( S ) > least;
        fi;
    od;
    return numvars - least;
    end; \# DimMonomialIdeal

```
```


## 

## Calculates the part of the squared "fractional dimension"

## contributed by the coordinate axis i in monomial ideal M

## Returns 2-15-1 if the fractional dimension contributed is 1

FracDimSqrForAxis := function(M, l)
local least, m, tmp;
least := false;
for m in M do
tmp := m[1];
if tmp < least then
least := tmp;
fi;
od;
if least <> false then
return least;
else
return 0;
fi;
end;

## 

## Calculates the "fractional dimension" of a monomial ideal

FracDimSqrPerAxis := function( M )
Iocal i, tmp;
if M = [] or (not IsList( M )) then
return false;
fi;
tmp:= [];
for i in [1..Length(M[1])] do
tmp[i] := FracDimSqrForAxis(M, i);
od;

```
```

    return tmp;
    end;

## 

## IsStructGP tests M to see if it is pairwise relatively prime

IsStructGB := function( M )
return false = PositionProperty( Sum ( List( M,
m -> List( m,
i -> SignInt(i)
)
)
),
i -> i > 1
);
end;

## 

## TzWhit tries to simplify the presentation P, by searching the space

## of Tietze transforms of P, to minimize the dimension of the monomial

## ideal generated by the leading terms of the variety of the ideal of

## SL_2 character relations which come from the relations of P, as

## presented by the function LTFrickeChar.

## 

## TzWhit uses, and is modelled upon, the low- and high- level

## functions for Tietze transformations in GAP.

TzWhit := function( arg ) \# arg = [P,
\# numtries, (default is 10)
\# useMostFreqP (default is false) ];
local tietze, count, rels, relators, RealRelators, m, n, tmp, FD, pp,
A, B, LeadingMonomials, p, pairs, pairs1, couples, x, g, elims,
left, right, i, j, current, total,
P, numtries, useMostFreqP ;
if Length( arg ) = 3 then \# this flag is undocumented, and indeed
hardiy ever works.

```
```

            useMostFreqP := arg[3];
    else useMostFreqP := false;
    fi;
    if Length( arg ) = 2 then numtries := arg[2];
    else numtries := 10;
    fi;
    P := arg[1];
    #TzFindCyclicJoins( P ); # do some preprocessing, and run consistency
    checks
tietze := P.tietze;
n := tietze[TZ_NUMGENS];
if n < 2 then
return; \# nothing to do to P
fi;
rels := tietze[TZ_RELATORS];
m := tietze[TZ_NUMRELS];
RealRelators := List( [1..m],
j -> TzWord( tietze, rels[j] ) ); \# the presentation
relators
\# relators is a list of the group relators, plus the relators times left-
\# multiplied by each generator
relators := [];
for tmp in tietze[TZ_GENERATORS] do
Append( relators, tmp*RealRelators );
od;
LeadingMonomials := List( relators,
r -> LTFrickeChar( r,
Copy( tietze[TZ_GENERATORS] )
) );
FD := FracDimSqrPerAxis( LeadingMonomials );
\# TzMostFrequentPairs returns a list of lists for each of the
\# NUMPAIRSCHKTzWHIT most frequently occuring length 2 subwords (a`eb^f)
\# in the relators. The format returned is:
\# [ frequency, a,b, x]

```
```

    # where x=0 for e=1,f=1
    # x=1 for e=1, f=-1
    # x=2 for e=-1, f=1
    # x=3 for e=f=-1
    if useMostFreqP then
    ```
```

        pairs := TzMostFrequentPairs( P, NUMPAIRSCEKTzWHIT, useMostFreqP );
    ```
        pairs := TzMostFrequentPairs( P, NUMPAIRSCEKTzWHIT, useMostFreqP );
        # Pick out those squares which have the same exponent
        # Pick out those squares which have the same exponent
    # Order them by the increase in fractional codimension which each
    # Order them by the increase in fractional codimension which each
    # might contribute if they are replaced by a new generator
    # might contribute if they are replaced by a new generator
    couples := \square;
    couples := \square;
    IsSet( couples );
    IsSet( couples );
    for p in pairs do
    for p in pairs do
        if p[4] = 0 or p[4] = 3 then
        if p[4] = 0 or p[4] = 3 then
            pp := Sorted( p{[2,3]});
            pp := Sorted( p{[2,3]});
            A := pp[1];
            A := pp[1];
            B := pp[2];
            B := pp[2];
            AddSet( couples,
            AddSet( couples,
                        Concatenation( [FD[A*(n+1) - (A^2 + A)/2 + B - A ]],
                        Concatenation( [FD[A*(n+1) - (A^2 + A)/2 + B - A ]],
                                    pp
                                    pp
                                    ) );
                                    ) );
            fi;
            fi;
        od;
        od;
    else
    else
    # Choose the NUMPAIRSCHKTzWHIT smallest-frac-codim pairs, order them by
    # Choose the NUMPAIRSCHKTzWHIT smallest-frac-codim pairs, order them by
the
the
    # increase in frac. codim. each might contribute if they were replaced.
    # increase in frac. codim. each might contribute if they were replaced.
    # This avoids string matching. We use a modified quicksort selection
    # This avoids string matching. We use a modified quicksort selection
method
method
    # (see e.g. Sedgewick, "Algorithms", Addison--Wesley, 1983)
    # (see e.g. Sedgewick, "Algorithms", Addison--Wesley, 1983)
        pairs := [];
        pairs := [];
        pairs1 := [];
        pairs1 := [];
        total := 0;
        total := 0;
        for A in [1..n-1] do #chuck out the zeros
        for A in [1..n-1] do #chuck out the zeros
        for B in [(A+1)..n] do
        for B in [(A+1)..n] do
            tmp := FD[A*(n+1) -(A^2 + A)/2 + B - A];
            tmp := FD[A*(n+1) -(A^2 + A)/2 + B - A];
            if tmp <> 0 then
            if tmp <> 0 then
                if tmp = 1 then
                if tmp = 1 then
                    Add( pairs1,
                    Add( pairs1,
                            [tmp, A, B] );
```

                            [tmp, A, B] );
    ```
```

                else
                        Add( pairs,
                            [tmp, A, B] );
                    total := total + tmp;
                fi;
        fi;
    od;
    od;
couples := [];
right := Length( pairs );
if right < NUMPAIRSCHKTzWHIT then
couples := Concatenation( pairs,
pairs1{[1..Minimum( NOMPAIRSCHKTzWHIT,
Length( pairs1 )
)
]} );
elif right = 0 then
couples := [];
else
left := 1;
i := 0; j := right;
current := QuoInt( total*2*NOMPAIRSCHKTzWHIT, tmp-2 );
\# e.g. a weighted average of pairs
\# best partition strategy if the values of FD[pairs]
\# were uniformly distributed
while left < right do
repeat
repeat
i := i + 1;
until pairs[i][1] <= current;
repeat
j := j - 1;
until pairs[j][1] >= current;
tmp := pairs[i];
pairs[i] := pairs[j];
pairs[j] := tmp;
until j <= i;
pairs[j] := pairs[i];
pairs[i] := pairs[right];
pairs[right] := tmp;
if i >= NOMPAIRSCHKTzWHIT then
right := i -1;

```
```

            fi;
            if i <= NUNPAIRSCHKTzWHIT then
                    left := i + 1;
            fi;
            current := pairs[right][1];
            i := left - 1;
            j := right;
        od;
        couples := pairs{[1..NUMPAIRSCHKTzWHIT]};
    fi;
    fi;
elims := [];
for p in couples{[1..Minimum( numtries, Length( couples ) )]} do
\# Add a new generator
AddGenerator( P );
x := P.generators[ Length( P.generators ) ];
\# Add relation x^-1*a*b
AddRelator( P,
x - 1*TzWord( tietze, [p[2], p[3]] ) );
\# choose the generator in {a,b} contributing the least fract. codim.
elims := [];
if FD[p[2]] < FD[p[3]] then
AddSet( elims, p[3] );
else
AddSet( elims, p[2] );
fi;
od;

# replace each ab or (ab)--1 by the generators introduced

P.searchSimultaneous := Maximum( 20, Length( couples )+10);
TzCheckRecord( P );
TzSearch( P );

# for each (a,b) eliminate the generator with the largest FracDimSqr

for g}\mathrm{ in elims do
TzEliminateGen(P , g);
od;

```
```

end; \# TzWhit

## 

## TzWhitTriples is like TzWhit, but for triples rather than pairs.

TzWhitTriples := function( arg ) \# arg = [P,
\# numtries, (default is 10) ]
local tietze, count, rels, relators, m, n, tmp, FD, Pp,
A, B, C, LeadingMonomials, p, triples, triples1, triplets, x, g,
elims, left, right, i, j, current, total,
P, numtries;
if Length( arg ) = 2 then numtries := arg[2];
else numtries := 10;
fi;
P := arg[1];
TzFindCyclicJoins( P ); \# do some preprocessing, and run consistency
checks
tietze := P.tietze;
n := tietze[TZ_NUMGENS];
if n< 3 then
return; \# nothing to do to P
fi;
rels := tietze[TZ_RELATORS];
m] := tietze[TZ_NUMRELS];
relators := List( [1..m],
j -> TzWord( tietze, rels[j] ) );
LeadingMonomials := List( relators,
I -> LTFrickeChar( r,
Copy( tietze[TZ_GENERATORS] )
) );
FD := FracDimSqrPerAxis( LeadingMonomials );

```
```

    # Choose the NUMTRIPLESCHKTzWHIT smallest-frac-codim triples, order them
    by the
\# increase in frac. codim. each might contribute if they were replaced.
\# This avoids string matching. We use a modified quicksort selection
method
\# (see e.g. Sedgewick, "Algorithms", Addison--Wesley, 1983)
triples := [];
triples1 := [];
total := 0;
setRememberMonabc( n );
for A in [1..n-2] do \#chuck out the zeros
for B in [(A+1)..n] do
for C in [(B+1)..n] do
tmp := FD[ RememberMonabc[n] [A] [B][C] ];
if tmp <> 0 then
if tmp = 1 then
Add( triples1,
[tmp, A, B, C] );
else
Add( triples,
[tmp, A, B, C] );
total := total + tmp;
fi;
fi;
od;
od;
od;
triplets := [];
right := Length( triples );
if right < NUMTRIPLESCHKTzWHIT then
triplets := Concatenation( triples,
triples1{[1..Minimum( NUMTRIPLESCHKTzWHIT,
Length( triples1 )
)
]} );
elif right = 0 then
triplets := [];
else
left := 1;
i := 0; j := right;
current := QuoInt( total*2*NUMTRIPLESCHKTzWHIT, tmp^2 );

```
```

                # e.g. a weighted average of triples
                # best partition strategy if the values of FD[triples]
                # were uniformly distributed
    while left < right do
        repeat
            repeat
                i := i + 1;
            until triples[i][1] <= current;
            repeat
                    j := j - 1;
            until triples[j][1] >= current;
            tmp := triples[i];
            triples[i] := triples[j];
            triples[j] := tmp;
        until j <= i;
    triples[j] := triples[i];
    triples[i] := triples[right];
    triples[right] := tmp;
    if i >= NUMTRIPLESCHKTzWHIT then
            right := i -1;
        fi;
        if i <= NUMTRIPLESCHKTzWHIT then
            left := i + 1;
        fi;
        current := triples[right][1];
        i := left - 1;
        j := right;
    od;
    triplets := triples{[1..NUMTRIPLESCHKTzWHIT]};
    fi;
elims := [];
for p in triplets{[1..Minimum( namtries, Length( triplets ) )]} do
\# Add a new generator
AddGenerator( P );
x := P.generators[ Length( P.generators ) ];
\# Add relation x^-1*a*b*c
AddRelator( P,
x^-1*IzWord( tietze, [p[2], p[3], p[4]] ) );

```
```

    # choose the generator in {a,b,c} contributing the least fract. codim.
    elims := [];
    if FD[p[2]] <= FD[p[3]] then
        if FD[p[3]]<= FD[p[4]] then
            AddSet( elims, p[4] );
        else
                AddSet( elims, p[3] );
        fi;
    elif FD[p[2]] <= FD[p[4]] then
        AddSet( elims, p[4] );
    else
        AddSet( elims, p[2] );
    fi;
    od;
    # replace each abc or (abc) - -1 by the generators introduced
    P.searchSimultaneous := Maximum( 20, Length( triplets )+10 );
    TzCheckRecord( P );
    TzSearch( P );
    # for each (a,b,c) eliminate the generator with the largest FracDimSqr
    for g}\mathrm{ in elims do
    TzEliminateGen(P , g);
    od;
end; \# TzWhitTriples

```

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