# LINEARIZING QUADRATIC TRANSFORMATIONS IN GENETIC ALGEBRAS 

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## 1. Introduction

In the deterministic theory of randomly mating infinite populations, in which there is no differential selection or fertility, certain types of quadratic transformations connecting one generation with the succeeding one are susceptible of a complete mathematical treatment because of their inherently simple structure. A quadratic transformation of an algebra is one which involves quadratic functions of the coordinates. The technique used to study these transformations was first introduced by Haldane [9] in a genetical context in 1930 in connection with polyploids and is a method of linearizing the transformations by extending the original vector space sufficiently, using functions of the coordinates, until the transformation becomes linear. The process is described as linearization and was studied in more detail, by Bennett [2] in connection with linked loci, and in an algebraical context by Holgate [11]. Holgate's paper gives a useful brief introduction to genetic algebras.

In this paper we study some questions which arise directly from Holgate's paper and obtain some more explicit results.

We shall use the abbreviation GA for a genetic algebra as defined by Schafer [17], for which we shall use the canonical form given by Gonshor [8]. Let $A_{n}$ denote the general GA of dimension $n+1$ with canonical basis $c_{0}, c_{1}, \ldots, c_{n}$. We shall assume for simplicity that $c_{0}$ is an idempotent element. Algebras arising in practice usually possess idempotents; Gonshor has given conditions for their existence [7]. Then multiplication is defined by

$$
c_{i} c_{j}=c_{j} c_{i}=\sum_{k=0}^{n} \lambda_{i j k} c_{k},
$$

where

$$
\begin{gathered}
\lambda_{000}=1, \quad \lambda_{0 j k}=0 \text { for } k<j, \quad \lambda_{00 k}=0 \text { for } k>0, \\
\lambda_{i j k}=0 \text { for }(i, j)>0 \text { and } k \leqslant \max (i, j) .
\end{gathered}
$$

$\lambda_{0 j j}$ are called the train roots of $A_{n}$. Those of the $\lambda_{i j k}$ which are not zero by definition will be called the constants, or structure constants of the algebra. $A_{n}$ is a commutative but usually non-associative algebra. It is
a baric algebra [5]. The baric property implies that if $x$ is a general element representing a population, it can be written uniquely as

$$
x=c_{0}+x_{1} c_{1}+x_{2} c_{2}+\ldots+x_{n} c_{n},
$$

where the coefficient of $c_{0}$ is 1 . We say that $x$ is of unit weight, the weight of an element in general being the coefficient of $c_{0}$. A special train algebra [5,7] can be defined as a GA in which all powers of the ideal with basis $\left\{c_{1}, \ldots, c_{n}\right\}$ (the nilideal) are ideals.

The plenary powers representing random mating between discrete non-overlapping generations are $x, x^{[2]}, x^{[3]}, \ldots, x^{[n]}, \ldots$, where

$$
x^{[n]}=x^{[n-1]} x^{[n-1]}
$$

and derive from the successive application of the quadratic transformation $\varphi: x \rightarrow x^{2}$ which in general will be a quadratic function of $x_{1}, x_{2}, \ldots, x_{n}$.
$B_{n}$ will denote the corresponding linearized vector space (with respect to $\varphi$ and $x$ ) which will have a basis defined by coordinates which are monomial functions $x_{1}{ }^{k_{1}} x_{2}{ }^{k_{2}} \ldots x_{n}{ }^{k_{n}}$ (where the $k_{i}$ are non-negative integers). These will be called the linearizing (coordinate) functions. This technique is subsequently defined and discussed very fully. The present problem is two-fold:
(i) to find the dimension of $B_{n}$ exactly, recursively, or asymptotically; and
(ii) to describe precisely the monomials required or, equivalently, to generate them explicitly.
We shall work in terms of a fixed canonical basis. However, as is well known, such bases are not unique-although the train roots are invariant. We shall show that the dimension of $B_{n}$ is independent of the basis in $A_{n}$. Haldane, working in a non-algebraic context, used the natural genetic basis and it is intuitively evident that the dimension of $B_{n}$ is an invariant of $A_{n}$ and the particular quadratic transformation under consideration. However, the linearizing functions are not characterized uniquely. With respect to a canonical basis in $A_{n}$ they can be taken as monomials, but for a general basis in $A_{n}$ they will usually be homogeneous polynomial functions in $n$ variables.

The origins of this paper are interesting and have three contemporaneous sources. Initially it was a remark by Holgate on finding an asymptotic value for $\operatorname{dim} B_{n}$ in the general case. Independently I was looking at the genetic algebra of polyploids with several linked loci. It seemed to me that the techniques necessary were an extension of the case for the genetic algebra of diploids with several linked loci but treated in a rather different way than had hitherto been attempted. This was by extending a technique of Etherington [5] and elucidating more fully the explicit nature
of the plenary roots and in particular the question of multiplicities. This in turn depended on a thorough and explicit understanding of the linearization technique. We study in a subsequent paper the induced linear transformation in $B_{n}$ and the plenary roots of the algebra with respect to quadratic transformations representing random mating.

## 2. Linearization

To illustrate some features of the linearization technique which Holgate developed [11] let us consider the following example discussed by him.

Example 1. Consider the algebra of tetraploidy $A_{2}$.

|  | $c_{0}$ | $c_{1}$ | $c_{2}$ |
| :--- | :---: | :---: | :---: |
| $c_{0}$ | $c_{0}$ | $\frac{1}{2} c_{1}$ | $\frac{1}{6} c_{2}$ |
| $c_{1}$ |  | $\frac{1}{6} c_{2}$ | 0 |
| $c_{2}$ |  |  | 0 |
| $x=$ |  |  |  |
| $x$ | $=c_{0}+x_{1} c_{1}+x_{2} c_{2}$. |  |  |

The quadratic transformation $\varphi: x \rightarrow x^{2}=c_{0}+x_{1} c_{1}+\left(\frac{1}{3} x_{2}+\frac{1}{6} x_{1}{ }^{2}\right) c_{2}$ can be represented as acting on the coordinates to give

$$
1 \varphi=1, \quad x_{1} \varphi=x_{1}, \quad x_{2} \varphi=\frac{1}{3} x_{2}+\frac{1}{6} x_{1}^{2} .
$$

The linearization of $\varphi$ defines a set

$$
M_{2}=\left\{1, x_{1}, x_{1}^{2}, x_{2}\right\}
$$

called the linearizing set of monomial (coordinate) functions with cardinality 4 , and we write

$$
\operatorname{Card} M_{2}=4
$$

$M_{2}$ then defines uniquely (up to isomorphism) an induced vector space $B_{2}$ of vectors

$$
\left(u_{0}, u_{1}, u_{11}, u_{2}\right)
$$

where the $u_{i}$ 's are coordinates with respect to a basis which we may write as

$$
c_{0}, c_{1}, c_{1} \otimes c_{1}, c_{2}
$$

Clearly $\operatorname{dim} B_{2}=\operatorname{Card} M_{2} . B_{2}$ in this case is isomorphic to $\mathbf{R}^{4}$.
We shall call the map $R: A_{2} \rightarrow B_{2}$,

$$
x=c_{0}+x_{1} c_{1}+x_{2} c_{2} \rightarrow\left(1, x_{1}, x_{1}^{2}, x_{2}\right)
$$

which maps the plane of unit weight in $A_{2}$ onto a variety $V$ in $B_{2}$, the linearization map.
$\tilde{\varphi}$ may be defined by linearity on the whole vector space $B_{2}$ :

$$
u_{0} \tilde{\varphi}=u_{0}, \quad u_{1} \tilde{\varphi}=u_{1}, \quad u_{11} \tilde{\varphi}=u_{11}, \quad u_{2} \tilde{\varphi}=\frac{1}{3} u_{2}+\frac{1}{6} u_{11}
$$

By examining the structure of $\tilde{\varphi}$ in $B_{2}$ we can, in particular, deduce its action on the variety $V$, and by projecting back onto $A_{2}$ we can deduce the action of $\varphi$.

The representation of $\tilde{\varphi}$ as a matrix in this basis is given by the matrix $A$,

$$
A=\left[\begin{array}{llll}
1 & 0 & 0 & 0 \\
0 & 1 & 0 & 0 \\
0 & 0 & 1 & \frac{1}{6} \\
0 & 0 & 0 & \frac{1}{3}
\end{array}\right]
$$

It is sometimes convenient in studying $\tilde{\varphi}$ to look at its action on the basis vectors. This is given by the rows of $A$ :

$$
c_{0} \tilde{\varphi}=c_{0}, \quad c_{1} \tilde{\varphi}=c_{1}, \quad c_{1} \otimes c_{1} \tilde{\varphi}=c_{1} \otimes c_{1}+\frac{1}{6} c_{2}, \quad c_{2} \tilde{\varphi}=\frac{1}{3} c_{2}
$$

Remarks. 1. The linearizing functions are not unique. We could define, for instance, the function $t=x_{1}{ }^{2}-4 x_{2}$. Then

$$
\begin{aligned}
x_{2} \varphi & =\frac{1}{3} x_{2}+\frac{1}{6}\left(t+4 x_{2}\right)=x_{2}+\frac{1}{6} t, \\
t \tilde{\varphi} & =\left(x_{1} \varphi\right)^{2}-4\left(x_{2} \varphi\right) \\
& =x_{1}{ }^{2}-4\left(x_{2}+\frac{1}{6} t\right)=x_{1}{ }^{2}-4 x_{2}-\frac{2}{3} t=\frac{1}{3} t .
\end{aligned}
$$

In order to obtain homogeneous linearizing functions we can take $x=x_{0} c_{0}+x_{1} c_{1}+x_{2} c_{2}$ where $x_{0}=1$, and then $t=x_{1}{ }^{2}-4 x_{2} x_{0}$. This polynomial function was the one used by Haldane working in the usual genetic basis and taking

$$
x=x_{0} A A+x_{1} A a+x_{2} a a, \quad \sum_{i=0}^{2} x_{i}=1
$$

It is quite clear that in either basis we can always take our linearizing functions as homogeneous polynomials by multiplying suitably by either $x_{0}=1$ or $\sum x_{i}=1$.
2. Clearly we are concerned with choosing a minimum number of linearizing functions in order to obtain a linearization of the problem. By a linearization we shall mean a minimal linearization.
3. It can sometimes happen that we can only obtain a canonical basis over the complex numbers, Heuch [10], in which case we obtain a complex induced vector space.
4. Linearization is very easily effected by working in a canonical basis. The linearizing functions which arise are monomials

$$
x_{1}{ }_{1}^{\alpha_{1}} x_{2}^{\alpha_{2}} \ldots x_{n}{ }^{\alpha_{n}}
$$

We can call this a canonical linearization. Linearization in an arbitrary basis is usually difficult to obtain and will yield, in general, linearizing functions which are polynomials in $\mathbf{R}\left[x_{1}, x_{2}, \ldots, x_{n}\right]$ or $\mathbf{C}\left[x_{1}, x_{2}, \ldots, x_{n}\right]$, where these denote polynomial rings over the real and complex numbers respectively:
5. It is sometimes helpful to look at the transformation $\varphi$ as defining a system of difference equations.
Writing $x_{n}=x_{n 0} c_{0}+x_{n 1} c_{1}+x_{n 2} c_{2}$ to denote the population in the $n$th generation, we have

$$
\begin{aligned}
& x_{n+1,0}=x_{n 0} \varphi=x_{n 0}=x_{00}=1, \\
& x_{n+1,1}=x_{n 1} \varphi=x_{n 1}=x_{01}, \\
& x_{n+1,2}=x_{n 2} \varphi=\frac{1}{3} x_{n 2}+\frac{1}{8} x_{n 1}{ }^{2} .
\end{aligned}
$$

We can easily solve this explicitly by recursion. In this paper we are only concerned with obtaining a canonical linearization and describing the dimension of the induced vector space.

As Holgate [11] showed, we can, however, always take monomials to linearize quadratic transformations $x \rightarrow x^{2}$ in (Schafer) genetic algebras.
We shall now prove the main theorem of this section, namely, that the dimension of $B_{n}$ is independent of the basis of $A_{n}$ (not necessarily the canonical basis) and it is also independent of the linearizing functions and a fortiori of the particular construction used to obtain them. Thus we shall have shown the uniqueness of $B_{n}$ (up to isomorphism) with respect to the particular quadratic transformation $x \rightarrow x^{2}$, where $x$ is a general element of unit weight in $A_{n}$. However, the linearizing functions in the natural genetic basis will not have a simple structure and can always be taken as homogeneous polynomials in $n+1$ variables $x_{0}, x_{1}, x_{2}, \ldots, x_{n}$. In the proof we shall have to use the result that non-singular transformations of a basis in $A_{n}$ induce non-singular transformations on the space of homogeneous polynomial functions. Thus, for example, different quadratic, cubic, quartic, ... functions in $x_{0}, x_{1}, x_{2}, \ldots, x_{n}$ are mapped onto different quadratic, cubic, quartic,... functions respectively in $y_{0}, y_{1}, y_{2}, \ldots, y_{n}$ by a non-singular transformation. For this result we use a property of Schläfian matrices (also known as induced matrices). Finally, we shall see that although the linearizing functions are not unique, even in a particular basis, they are unique up to their degree. By that I mean that the same number of linearizing functions of the same degree are needed, whatever the basis in $A_{n}$, in order to linearize the quadratic transformation.

We remind the reader of some standard terminology using a simple example. Consider a two-dimensional vector space with coordinates ( $x_{1}, x_{2}$ ) and a non-singular linear transformation $x=y A$, giving new coordinates $\left(y_{1}, y_{2}\right)$, so that

$$
x_{1}=a_{11} y_{1}+a_{21} y_{2}, \quad x_{2}=a_{12} y_{1}+a_{22} y_{2}
$$

The homogeneous quadratic form $Q=a x_{1}{ }^{2}+2 b x_{1} x_{2}+c x_{2}{ }^{2}$ will be transformed into $Q=a^{\prime} y_{1}{ }^{2}+2 b^{\prime} y_{1} y_{2}+c^{\prime} y_{2}{ }^{2}$, where

$$
\left(a^{\prime}, b^{\prime}, c^{\prime}\right)=(a, b, c)\left[\begin{array}{ccc}
a_{11}^{2} & a_{11} a_{21} & a_{21}^{2} \\
2 a_{11} a_{12} & a_{11} a_{22}+a_{21} a_{12} & 2 a_{21} a_{22} \\
a_{12}^{2} & a_{12} a_{22} & a_{22}{ }^{2}
\end{array}\right]
$$

We say that $A$ induces a transformation $\tilde{A}:(a, b, c) \rightarrow\left(a^{\prime}, b^{\prime}, c^{\prime}\right)$ defined by the above matrix, which is a second-order Schläflian, Muir [15].

In general for a homogeneous form $f$ of degree $r$,

$$
f=\sum_{i} a_{i} x_{0}^{\alpha_{i}} x_{1}{ }^{\beta_{i}} x_{2}{ }^{\gamma_{4}} \ldots x_{n}^{\nu_{i}},
$$

where $\alpha_{i}+\beta_{i}+\ldots+\nu_{i}=r$ for all $i$, a linear transformation $\left(x_{i}\right) \rightarrow\left(y_{i}\right)$ induces a linear transformation of the coefficients of the form, given by $\left(a_{i}\right) \rightarrow\left(a_{i}^{\prime}\right)$, which is given by a Schläflian matrix of the $r$ th order.

By a result of Schläfli [18], quoted by Muir in [15], this matrix has a determinant equal to $|A|^{\alpha}$, where $\alpha=\binom{n+r-1}{n}$. Hence if $A$ is nonsingular, so is $\tilde{A}$.

This proves that different $r$ th degree homogeneous forms are mapped onto different $r$ th degree forms under the induced transformation of a change of basis in $A_{n}$. Hence the linearizing functions are unique up to their degree and the number of homogeneous linearizing functions of a given degree is an invariant of the algebra $A_{n}$ with respect to a given quadratic transformation. If one chooses a minimal set of linearizing functions with respect to any basis then this will define uniquely the dimension of $B_{n}$.

Thus we have the following result.
Theorem 1: the first fundamental theorem of genetic algebra. The induced vector space corresponding to a linearization of a quadratic transformation is unique to within an isomorphism.

Corollary. The number of linearizing functions of a given degree is independent of the basis in $A_{n}$.

## 3. The general Schafer genetic algebra

We consider an arbitrary (Schafer) genetic algebra (GA) $A_{n}$ with the multiplication table defined as in $\S 1$ and we shall assume that none of the constants of the algebra is zero. It is easily seen that in this case $A_{n}$ is a special train algebra. This condition gives the most general genetic algebra $A_{n}$ of a particular dimension in the sense that any other genetic algebra of the same dimension can only have at most the same number of non-zero structure constants as $A_{n}$. Letting some of these be zero decreases $\operatorname{dim} B_{n}$. Hence the general case gives us a maximal dimension for $B_{n}$ for a particular quadratic transformation and for a general element of unit weight in a GA of dimension $n$.

Holgate [11] has given examples illustrating the linearization technique for $A_{1}$ and $A_{2}$. We illustrate the case for $A_{3}$, and state the results for $A_{4}$.

The algebra $A_{3}$

|  | $c_{0}$ | $c_{1}$ | $c_{2}$ | $c_{3}$ |
| :---: | :---: | :---: | :---: | :---: |
| $c_{0}$ | $c_{0}$ | $\lambda_{011} c_{1}+\lambda_{012} c_{2}+\lambda_{013} c_{3}$ | $\lambda_{022} c_{2}+\lambda_{023} c_{3}$ | $\lambda_{033} c_{3}$ |
| $c_{1}$ |  | $\lambda_{112} c_{2}+\lambda_{113} c_{3}$ | $\lambda_{123} c_{3}$ | 0 |
| $c_{2}$ |  |  | $\lambda_{223} c_{3}$ | 0 |
| $c_{3}$ |  |  |  | 0 |

$$
x=c_{0}+x_{1} c_{1}+x_{2} c_{2}+x_{3} c_{3}
$$

$x \varphi=x^{2}$ represents the quadratic transformation.
The transformation $\varphi$ may be viewed as acting on the coordinates $1, x_{1}, x_{2}$, $x_{3}$ as follows:

$$
\begin{aligned}
1 \varphi & =1 \\
x_{1} \varphi & =2 \lambda_{011} x_{1} \\
x_{2} \varphi & =2 \lambda_{012} x_{1}+2 \lambda_{022} x_{2}+\lambda_{112} x_{1}{ }^{2} \\
x_{3} \varphi & =2 \lambda_{013} x_{1}+2 \lambda_{023} x_{2}+2 \lambda_{033} x_{3}+\lambda_{113} x_{1}^{2}+2 \lambda_{123} x_{1} x_{2}+\lambda_{223} x_{2}{ }^{2} .
\end{aligned}
$$

We define new variables $y_{1}=x_{1}{ }^{2}, y_{2}=x_{1} x_{2}, y_{3}=x_{2}{ }^{2}$.
In the transition from $x$ to $x \varphi, x_{i} x_{j}$ is replaced by $x_{i} \varphi \cdot x_{j} \varphi$, so for the induced transformation we define $\left(x_{i} x_{j}\right) \tilde{\varphi}=x_{i} \varphi \cdot x_{j} \varphi$. Hence we have

$$
\begin{aligned}
& y_{1} \tilde{\varphi}=4 \lambda_{011}{ }^{2} y_{1}, \\
& y_{2} \tilde{\varphi}=2 \lambda_{011} x_{1}\left(2 \lambda_{012} x_{1}+2 \lambda_{022} x_{2}+\lambda_{112} x_{1}^{2}\right), \\
& y_{3} \tilde{\varphi}=\left(2 \lambda_{012} x_{1}+2 \lambda_{022} x_{2}+\lambda_{112} x_{1}^{2}\right)^{2} .
\end{aligned}
$$

We now need to define further variables in order to linearize these equations:

Then

$$
y_{4}=x_{1}^{3}, \quad y_{5}=x_{1}^{4}, \quad y_{6}=x_{1}{ }^{2} x_{2}
$$

$$
\begin{aligned}
& y_{4} \tilde{\varphi}=8 \lambda_{011}{ }^{3} y_{4} \\
& y_{5} \tilde{\varphi}=16 \lambda_{011}{ }^{4} y_{5} \\
& y_{6} \tilde{\varphi}=4 \lambda_{011}{ }^{2} x_{1}{ }^{2}\left(2 \lambda_{012} x_{1}+2 \lambda_{022} x_{2}+\lambda_{112} x_{1}{ }^{2}\right)
\end{aligned}
$$

The original quadratic transformation is now completely linearized. Altogether we needed six monomials to linearize the transformation. Hence $B_{3}$ is a 10 -dimensional vector space induced by the set $M_{3}$ of coordinate functions,

$$
M_{3}=\left\{1, x_{1}, x_{2}, x_{3}, x_{1}^{2}, x_{1} x_{2}, x_{2}^{2}, x_{1}^{3}, x_{1}^{4}, x_{1}^{2} x_{2}\right\} .
$$

Similarly for $A_{1}$ and $A_{2}$ we have the corresponding spaces $B_{1}$ and $B_{2}$ induced by $M_{1}=\left\{1, x_{1}\right\}$ and $M_{2}=\left\{1, x_{1}, x_{1}{ }^{2}, x_{2}\right\}$ respectively. For $A_{4}, B_{4}$ is 36 -dimensional and is induced by

$$
\begin{aligned}
& M_{4}=\left\{1, x_{1}, x_{2}, x_{1}{ }^{2}, x_{1}{ }^{3}, x_{3}, x_{1} x_{2}, x_{4}, x_{2}{ }^{2}, x_{1} x_{3}, x_{1}{ }^{2} x_{2}, x_{1}{ }^{4}, x_{1}{ }^{5}, x_{2} x_{3}, x_{1} x_{2}{ }^{2},\right. \\
& x_{1}{ }^{3} x_{2}, x_{1}{ }^{2} x_{3}, x_{1}{ }^{6}, x_{3}{ }^{2}, x_{1} x_{2} x_{3}, x_{1}{ }^{2} x_{2}{ }^{2}, x_{1}{ }^{3} x_{3}, x_{1}{ }^{4} x_{2}, x_{1}{ }^{7}, x_{1}{ }^{5} x_{2} \text {, } \\
& \left.x_{1}{ }^{3} x_{2}{ }^{2}, x_{1}{ }^{2} x_{2} x_{3}, x_{2}{ }^{2} x_{3}, x_{1} x_{2}{ }^{3}, x_{1}{ }^{4} x_{3}, x_{1}{ }^{8}, x_{1}{ }^{6} x_{2}, x_{1}{ }^{4} x_{2}{ }^{2}, x_{1}{ }^{2} x_{2}{ }^{3}, x_{2}{ }^{4}, x_{2}{ }^{3}\right\} .
\end{aligned}
$$

It is convenient to define a polynomial function corresponding to each of these spaces:

$$
\begin{aligned}
P_{1}\left(1, x_{1}\right) & =1+x_{1} \\
P_{2}\left(1, x_{1}, x_{2}\right) & =\left(1+x_{1}\right)^{2}+x_{2}, \\
P_{3}\left(1, x_{1}, x_{2}, x_{3}\right) & =\left(\left(1+x_{1}\right)^{2}+x_{2}\right)^{2}+x_{3}, \\
P_{4}\left(1, x_{1}, x_{2}, x_{3}, x_{4}\right) & =\left(\left(\left(1+x_{1}\right)^{2}+x_{2}\right)^{2}+x_{3}\right)^{2}+x_{4}
\end{aligned}
$$

By expanding these polynomials (from the outer brackets) one can see that the constituent monomials are given by the sets $M_{1}, M_{2}, M_{3}, M_{4}$ respectively. Thus we have produced explicitly a generating polynomial for the coordinate functions. We may call such a polynomial $P$ a generator for $M$. By abuse of language we may say that $P$ generates the vector space $B$. This polynomial plays a key role in obtaining the properties of $B_{n}$.

We shall write $P_{n}$ for $P_{n}\left(1, x_{1}, \ldots, x_{n}\right)$.
Theorem 2. The linearizing set $M_{n}$ of coordinate functions for the quadratic transformation $\varphi: x \rightarrow x^{2}$, where $x$ is a general element of unit weight in $A_{n}$, is generated by the polynomial $P_{n}$ defined recursively by $P_{n}=P_{n-1}{ }^{2}+x_{n}$, for $n=1,2, \ldots, P_{0}=1$.

Proof. The theorem is clearly true for $A_{n}$, with $n \leqslant 4$.
Assume the theorem is true for $A_{n-1}$ and suppose that $B_{n-1}$ is induced by the set of monomials

$$
M_{n-1}=\left\{1, y_{1}, y_{2}, y_{3}, y_{4}, y_{5}, z_{1}, y_{6}, z_{2}, z_{3}, y_{7}, \ldots\right\}
$$

obtained by expanding $P_{n-1}$ and relabelling the monomials by the variables $y_{i}$ and $z_{i}$ listed in a special (not quite uniquely determined) order as follows. The $y_{i}$ 's are either the linear or quadratic variables occurring in the transformation equations of $A_{n-1}$; they will be called primary variables. They are listed in the order in which they occur in the transformation equations. Thus $y_{1}=x_{1}, y_{2}=x_{2}, y_{3}=x_{1}^{2}, y_{4}=x_{3}, y_{5}=x_{1} x_{2}, \ldots$, and we exclude any repetitions of such variables. The $z_{i}$ 's are further monomials generated by the quadratic primary variables in order to linearize $\varphi$. Thus $z_{1}=x_{1}{ }^{3}$ is generated by $\tilde{\varphi}$ acting on $y_{5}=x_{1} x_{2}$. By abuse of language we shall say that $x_{1} x_{2}$ generates $x_{1}{ }^{3}$. If $z_{i}$ does not occur amongst any of the variables to its left, namely amongst the previously listed $y$ 's or $z$ 's, it is included in the listing; otherwise it is excluded. This avoids any repetitions. Similarly each $z_{i}$ may generate further $z$ 's. These are listed in the ordering as they are generated, apart from repetitions. Thus any $z_{i}$ will be generated by some $y$ or $z$ to the left of it. We call the $z$ 's secondary variables.

Consider $A_{n} ; x_{n} \varphi$ contains a term in $x_{n-1}{ }^{2}$ since all the $\lambda$ 's are assumed to be non-zero. Now $x_{n-1} \varphi$ is a linear combination of all the primary variables of $M_{n-1}$ for the same reason, and

$$
x_{n-1}^{2} \tilde{\varphi}=\left(x_{n-1} \varphi\right)^{2}
$$

Hence the induced transformation on $x_{n-1}{ }^{2}$ will define further variables which are all the possible pairwise products (including squares) of the primary variables of $M_{n-1}$.

Suppose that the primary variables of $M_{n-1}$ are $\left\{1, y_{1}, y_{2}, \ldots, y_{p}\right\}$. Hence $y_{i} y_{j}$ are contained in $M_{n}$ for all $i, j=1, \ldots, p$.

We must show that all the other pairwise products (including squares) of the variables (primary and secondary) of $M_{n-1}$ are included in $M_{n}$. There are two cases to consider.

Case 1: the variables $y_{i} z_{j} \in M_{n}$ for $y_{i}, z_{j} \in M_{n-1}$ for all $i, j$. We use induction on $j$. Consider $y_{i} z_{1} . z_{1}$ is generated by $y_{5}$, that is, $y_{5} \tilde{\varphi}=\alpha_{1} z+\ldots$ and $y_{i} y_{5} \in M_{n}$,

$$
\begin{aligned}
\left(y_{i} y_{5}\right) \tilde{\varphi} & =\left(y_{i}\right) \tilde{\varphi}\left(y_{5}\right) \tilde{\varphi} \\
y_{i} \tilde{\varphi} & =\beta y_{i}+\ldots
\end{aligned}
$$

Hence the induced transformation on $y_{i} y_{5}$ defines a variable $y_{i} z_{1}$. This therefore is included in $M_{n}$.

Now we make the inductive hypothesis that $y_{i} z_{j} \in M_{n}$ for fixed $i$ and all $j<k$. Consider $y_{i} z_{k}$. By the ordering $z_{k}$ is generated by either a $y$ variable or a $z$ variable to the left of it in the listing. If the former case holds

$$
\begin{aligned}
& y_{j} \tilde{\varphi}=\gamma z_{k}+\ldots \quad \text { for some } j \\
& y_{i} \tilde{\varphi}=\beta y_{i}+\ldots
\end{aligned}
$$

The induced transformation on $y_{i} y_{j}$ will define $y_{i} z_{k}$. Hence $y_{i} z_{k}$ belongs to $M_{n}$. On the other hand, if $z_{k}$ is generated by $z_{j}$, for $j<k, y_{i} z_{j} \in M_{n}$ by the inductive hypothesis and

$$
\left(y_{i}\right) \tilde{\varphi}=\beta y_{i}+\ldots, \quad z_{j} \tilde{\varphi}=\delta z_{k}+\ldots
$$

Hence the induced transformation on $y_{i} z_{j}$ defines the variable $y_{i} z_{k}$. Therefore $y_{i} z_{k}$ belongs to $M_{n}$.

We conclude that $y_{i} z_{j} \in M_{n}$ for all $j$, fixed $i$. Since $y_{i}$ was arbitrary, this is true for all $i$ and $j$, including $y_{0}=1$.

Case 2: the variables $z_{i} z_{j} \in M_{n}$ for all $z_{i}, z_{j}$ in $M_{n-1}$. By the ordering $z_{i}$ is generated from $z_{i^{\prime}}$ or $y_{i^{\prime \prime}}$ and $z_{j}$ from $z_{j^{\prime}}$ or $y_{j^{\prime \prime}}$. This gives four possibilities:
(a) $z_{i}$ is generated from $y_{i^{\prime \prime}}$ and $z_{j}$ from $y_{j^{\prime \prime}}$;
(b) $z_{i}$ is generated from $z_{i^{\prime}}\left(i^{\prime}<i\right)$ and $z_{j}$ from $z_{j^{\prime}}\left(j^{\prime}<j\right)$;
(c) $z_{i}$ is generated from $y_{i^{\prime \prime}}$ and $z_{j}$ from $z_{j^{\prime}}\left(j^{\prime}<j\right)$;
(d) $z_{i}$ is generated from $z_{i^{\prime}}\left(i^{\prime}<i\right)$ and $z_{j}$ from $y_{j^{\prime \prime}}$.

We examine each sub-case in turn.
Sub-case (a). We must show that $z_{i} z_{j} \in M_{n}$. Since $y_{i^{\prime \prime}}$ and $y_{j^{\prime \prime}}$ are primary, $y_{i^{n}} y_{j^{\prime \prime}} \in M_{n}$,

$$
\left(y_{i^{n}} y_{j^{n}}\right) \tilde{\varphi}=\left(\alpha z_{i}+\ldots\right)\left(\beta z_{j}+\ldots\right)
$$

Hence $z_{i} z_{j} \in M_{n}$.
Sub-case (b). Here we use double induction on $i$ and $j$. For $i=j=1$, $z_{1}{ }^{2} \in M_{n}$ since

$$
y_{5}{ }^{2} \tilde{\varphi}=\left(\alpha z_{1}+\ldots\right)\left(\alpha z_{1}+\ldots\right)
$$

Assume that $z_{i^{\prime}} z_{j^{\prime}} \in M_{n}$ for all $i^{\prime}<i, j^{\prime}<j$. We shall show that this implies $z_{i} z_{j} \in M_{n}$ and hence by the second principle of induction (applied twice) $z_{i} z_{j} \in M_{n}$ for all $z_{i}, z_{j}$ in $M_{n-1}$ whenever Sub-case (b) is true. Since $z_{i} z_{j^{\prime}} \in M_{n}$ by the inductive hypothesis

$$
\left(z_{i^{\prime}} z_{j^{\prime}}\right) \tilde{\varphi}=\left(\alpha z_{i}+\ldots\right)\left(\beta z_{j}+\ldots\right)
$$

defines $z_{i} z_{j}$. Hence $z_{i} z_{j} \in M_{n}$.
Sub-case (c). This follows easily by using Case l, since $y_{i^{n} z_{j^{\prime}} \in M_{n}}$. Hence

$$
\left(y_{i^{\circ}} z_{j^{\prime}}\right) \tilde{\varphi}=\left(\alpha z_{i}+\ldots\right)\left(\beta z_{j}+\ldots\right)
$$

which implies $z_{i} z_{j} \in M_{n}$.

Sub-case (d). The proof is similar to (c).
We have therefore shown that all pairwise products (including squares) of the variables (primary and secondary) of $M_{n-1}$ are included in $M_{n}$.

Finally, we must show that the other primary variables in $x_{n} \varphi$ (excluding $x_{n-1}{ }^{2}$ ) will occur in $P_{n-1}{ }^{2}$. This is obvious since any such variable $x_{i} x_{j}$, with $i<n, j<n$, will occur as a product of linear primary variables in $M_{n-1}$. The secondary variables they generate will clearly be in $P_{n-1}{ }^{2}$.

Hence the linearized space of $A_{n}$ is given precisely by all the monomials in

$$
P_{n}=P_{n-1}^{2}+x_{n} .
$$

Another characterization of the monomial functions is given by a linear diophantine inequality. This can be converted to a linear diophantine equality. Such equalities are important in the theory of partitions. We can also represent the monomials using weight functions.
If $x_{i}$ has weight $w\left(x_{i}\right)=2^{i-1}$, for $i=1,2, \ldots$, we define the weight of $x_{1}{ }_{1}^{a_{1}} x_{2}^{a_{2}} \ldots x_{n}^{a_{n}}$ as $\sum_{i=1}^{n} a_{i} w\left(x_{i}\right)$.

The use of weight functions is simply a combinatorial device for investigating $\operatorname{dim} B_{n}$. There is no connection with the concept of baric weight.

Proposition 1. The monomials required to linearize the quadratic transformation $x \rightarrow x^{2}$, where $x$ is a general element of unit weight in $A_{n}$, are all those monomials of weight less than or equal to $2^{n-1}$. Equivalently $x_{1}{ }^{a_{1}} x_{2}{ }_{2} \ldots \ldots x_{n} a_{n}$ is such a monomial if $a_{1}, a_{2}, \ldots, a_{n}$ are integral non-negative solutions of the inequality

$$
a_{1}+2 a_{2}+2^{2} a_{3}+2^{3} a_{4}+\ldots+2^{n-1} a_{n} \leqslant 2^{n-1}, \text { for } n \geqslant 1 .
$$

Proof. Clearly the proposition is true for $A_{1}$. Consider

$$
\left\{x_{1}^{a_{1}} x_{2}^{a_{2}}: a_{1}+2 a_{2} \leqslant 2, a_{i} \geqslant 0, a_{i} \in \mathbf{N}\right\} .
$$

The solutions of the inequality are $\{(0,0),(0,1),(1,0),(2,0)\}$ which when substituted in $x_{1}{ }_{1} x_{2}{ }_{2}^{a_{2}}$ define the set $M_{2}$. Hence it is true for $A_{2}$.

Assume the truth of the proposition for $A_{n-1}$. Now $P_{n}=P_{n-1}{ }^{2}+x_{n}$ and by the inductive hypothesis $P_{n-1}$ is the set of monomials $m$ such that
 monomials $m_{i} m_{j}$ and $x_{n}$, where $m_{i}, m_{j}$ are in $M_{n-1}$. Now

$$
\begin{aligned}
w\left(m_{i} m_{j}\right) & =w\left(m_{i}\right)+w\left(m_{j}\right) \leqslant 2^{n-2} \cdot 2=2^{n-1} \\
w\left(x_{n}\right) & =2^{n-1}
\end{aligned}
$$

Hence $P_{n}$ consists of all the monomials $m$ in $x_{1}, \ldots, x_{n}$ such that $w(m) \leqslant 2^{n-1}$, which establishes the truth of the proposition.

If $m=x_{1}{ }^{a_{1}} x_{2}{ }^{a_{2}} \ldots x_{n}{ }^{a_{n}}$, then

$$
w(m)=a_{1}+2 a_{2}+2^{2} a_{3}+\ldots+2^{n-1} a_{n} \leqslant 2^{n-1}, \quad \text { for } n \geqslant 1
$$

Hence $a_{1}, \ldots, a_{n}$ are non-negative integral solutions of this diophantine inequality.

Proposition 2. The number of monomials in $M_{n}$ is precisely the number of integral, non-negative solutions of the linear diophantine equation

$$
a_{1}+2 a_{2}+2^{2} a_{3}+\ldots+2^{n} a_{n+1}=2^{n}
$$

Proof. See the next section.
We can also define generating functions for the set $M_{n}$ and for $u_{n}=\operatorname{dim} B_{n}=\operatorname{Card} M_{n}$. Define

$$
\begin{aligned}
\Phi(x ; t)= & (1+ \\
& \left.x_{1} t+x_{1}{ }^{2} t^{2}+x_{1}{ }^{3} t^{3}+\ldots\right)\left(1+x_{2} t^{2}+x_{2}{ }^{2} t^{4}+\ldots\right) \\
& \times\left(1+x_{3} t^{4}+x_{3} t^{8}+\ldots\right)\left(1+x_{4} t^{8}+\ldots\right) \ldots \\
= & \prod_{i=1}^{\infty}\left(1-x_{i} t^{2-1}\right)^{-1}
\end{aligned}
$$

where these are formal power series. Then clearly $\Phi(x ; t)$ is a form of generating function for $M_{n}$ in the sense that $M_{n}$ will consist of all the coefficients of powers of $t^{i}$ such that $i \leqslant 2^{n-1}$. We may obtain a proper generating function for $M_{n}$ by considering $(1-t)^{-1} \Phi(x ; t)$. Then the coefficient of $t^{2^{n-1}}$ in this expansion will give precisely the set $M_{n}$.

Similarly, by considering

$$
(1-t)^{-1} \Phi(1 ; t)=(1-t)^{-1} \prod_{i=1}^{\infty}\left(1-t^{2-1}\right)^{-1}
$$

we will have a generating function for $u_{n}$. Here $u_{n}$ is the coefficient of $t^{2^{n-1}}$.
We have thus given various alternative characterizations of $M_{n}$ and its cardinality, and hence of $B_{n}$ and its dimension.

## 4. The dimension of the linearized space

A general problem in the theory of partitions is to partition a given number $n$ with respect to some given quantities $b_{1}, b_{2}, \ldots, b_{m}$. This is equivalent to solving the linear diophantine equation

$$
\begin{equation*}
y_{1} b_{1}+y_{2} b_{2}+\ldots+y_{m} b_{m}=n \tag{1}
\end{equation*}
$$

for integral, non-negative values $y_{1}, y_{2}, \ldots, y_{m}$. The number of solutions of this equation is called the denumerant $D(m, n)$ (with respect to the base which is the fixed set of summands $b_{1}, b_{2}, \ldots, b_{m}$ ). Sylvester was the first to investigate this problem and gave some very deep results using complex variable techniques. Several other writers, mentioned by Dickson in
[4, Chapter 3], have given some remarkable explicit formulae for the number of solutions of (1). However, they are usually impossible to evaluate for large $m$ and $n$, and are of little practical value. E. T. Bell [1] has given a modern treatment of this problem. Formulae for $D(m, n)$ for special cases and small $m$ and $n$ can be obtained. A brief survey is given by Riordan in [16].
It is well known that if we define a generating function for $D(m, n)$ by

$$
\Phi_{m}(t)=\sum_{n} D(m, n) t^{n},
$$

then
We define

$$
\Phi_{m}(t)=\left\{\left(1-t^{b_{1}}\right)\left(1-t^{b_{2}}\right) \ldots\left(1-t^{b_{m}}\right)\right\}^{-1} .
$$

$$
\begin{aligned}
D(m, 0)=1 & (m=0,1,2, \ldots), \\
D(0, n)=0 & (n=1,2,3, \ldots) .
\end{aligned}
$$

Proof of Proposition 2. This is based on a recursive relation derived by Blom and Fröberg [3, p. 63].
Consider

$$
\begin{aligned}
\Phi_{n}(t) & =\left\{(1-t)\left(1-t^{2}\right) \ldots\left(1-t^{2 n-1}\right)\right\}^{-1}, \\
\Phi_{n+1}(t) & =\left\{(1-t)\left(1-t^{2}\right) \ldots\left(1-t^{2}\right)\right\}^{-1} .
\end{aligned}
$$

We have

$$
\left.\begin{array}{rl}
\Phi_{n+1}(t)=(1-t)^{-1} \Phi_{n}\left(t^{2}\right) \\
=\left(1+t+t^{2}+t^{3}+\ldots\right)(D(n, 0) & +t^{2} D(n, 1)+t^{2} D(n, 2) \\
& +\ldots+t^{2^{n}-4} D\left(n, 2^{n-1}-2\right)
\end{array}\right)+t^{2^{n}-2} D\left(n, 2^{n-1}-1\right) .
$$

Equating coefficients of $t^{2}$ in both sides, we obtain

$$
D\left(n+1,2^{n}\right)=D(n, 0)+D(n, 1)+D(n, 2)+\ldots+D\left(n, 2^{n-1}\right),
$$

where $D(m, n)$ is with respect to the set of summands $1,2,2^{2}, \ldots, 2^{m-1}$. The left-hand side is the number of solutions of the equality

$$
a_{1}+2 a_{2}+2^{2} a_{3}+\ldots+2^{n} a_{n+1}=2^{n},
$$

and the right-hand side is the number of solutions in the inequality of Proposition 1. This completes the proof.
Our interest centres on a very special case of this problem when the base consists of powers of 2 , so that

$$
b_{1}=1, \quad b_{2}=2, \quad b_{3}=2^{2}, \quad \ldots, \quad b_{m}=2^{m-1}
$$

and $n=2^{m-1}$.
Thus, from Proposition 2, $\operatorname{dim} B_{n}=D\left(n+1,2^{n}\right)$. In this case a more complete and easier recurrence solution can be given than in the general
problem. A complete analysis of this case is given by Blom and Fröberg [3]. We shall also use a result of Mahler's [13] to give an asymptotic value for $u_{n}=\operatorname{dim} B_{n}$.

A number of recursive procedures are possible for computing $u_{n}$, depending on which method is used for characterizing $B_{n}$.

Method 1: using the polynomials $P_{n}$. The number of different monomials in the expansion of $P_{n}$ will be denoted by $u_{n}$, and the number of different monomials in $P_{n}{ }^{m}(m=0,1,2, \ldots)$ will be denoted by $u_{n, m}$.

We can easily derive a difference equation for $u_{n, m}$ by expanding

$$
P_{n+1}^{m}=\left(\left(P_{n-1}^{2}+x_{n}\right)^{2}+x_{n+1}\right)^{m}
$$

and forming the corresponding equation for $u_{n+1, m}$. Collecting terms together we eventually obtain

$$
\begin{aligned}
u_{n+1, m}=\left[u_{n-1,4 m}\right. & \left.+u_{n-1,4 m-2}\right]+2\left[u_{n-1,4 m-4}+u_{n-1,4 m-6}\right] \\
& +3\left[u_{n-1,4 m-8}+u_{n-1,4 m-10}\right] \\
& +\ldots+m\left[u_{n-1,4}+u_{n-1,2}\right]+m+1 \quad \text { for } n \geqslant 2 .
\end{aligned}
$$

One easily finds

$$
u_{1, m}=m+1, \quad u_{2, m}=(m+1)^{2}, \quad u_{3, m}=1+\frac{1}{3}\left[4 m^{3}+12 m^{2}+11 m\right] .
$$

Further explicit solutions for small values of $n$ are possible in terms of Bernoulli polynomials. Using the recursion formula one can compute $u_{n}$.

For the general algebra $A_{n}$ (where the $\lambda$ 's are non-zero) we obtain the following table:

$$
\begin{array}{ccccccccc} 
& A_{0} & A_{1} & A_{2} & A_{3} & A_{4} & A_{5} & A_{6} & A_{7} \\
u_{n} & 1 & 2 & 4 & 10 & 36 & 202 & 1828 & 27337
\end{array}
$$

Method 2: using a diophantine equality. An alternative procedure is to use the characterization of $B_{n}$ by the diophantine equality in Proposition 2. Blom and Fröberg [3] have given explicit details of the recursion procedure based on Bell's formula [1]. We give their result. Our sequence is
where

$$
u_{m-1}=c_{m 1}+1, \text { for } m=1,2, \ldots
$$

$$
\begin{aligned}
c_{m+1, i} & =\sum_{j=i-1}^{2 i} a_{i j} c_{m j} \quad(i, m \geqslant 1) \\
c_{i j} & =0 \quad(j \geqslant i) \\
c_{m 0} & =1 \quad(\text { all } m)
\end{aligned}
$$

and $a_{i \jmath}$ is given by either

$$
a_{i+1, j+1}=2 a_{i j}+a_{i, j-1} \quad\left(i \geqslant 0, j \geqslant 0, a_{00}=1\right)
$$

or

$$
a_{i j}=\binom{i}{j-i} 2^{2 i-j}+\binom{i-1}{j-i+1} 2^{2 i-j-2}
$$

They calculate first a table of the coefficients $a_{i j}$ and then derive $c_{m i}$. Using this method we easily find $u_{7}=27337$, whereas the calculation is more protracted by the first method.

Mahler [13], in the course of a more general problem, investigated similar diophantine equations but with a base in terms of an arbitrary integer $a \geqslant 2$. Using his result, we obtain the asymptotic formula

$$
\log u_{n} \sim \frac{\left(\log 2^{n}\right)^{2}}{2 \log 2}=n^{2} \log \sqrt{ } 2 \quad(\text { as } n \rightarrow \infty)
$$

## 5. A special case

Here we examine a class of (Schafer) genetic algebras $A_{n}$, with the usual canonical basis $c_{0}, c_{1}, \ldots, c_{n}$, which require only quadratic functions of the coordinates $x_{1}, \ldots, x_{n}$ to linearize quadratic transformations denoting random mating. This is related to a question raised by Moran [14, Chapter 2].

Proposition 3. If $A_{n}$ is a Schafer genetic algebra with idempotent $c_{0}$ such that $c_{i} c_{j}=\lambda_{i j n} c_{n}$, for $1 \leqslant i, j<n, \lambda_{i j n} \neq 0$, that is, the result of multiplication in the nil-ideal is a multiple of $c_{n}$, then $A_{n}$ is an algebra which requires only quadratic functions to linearize a quadratic transformation $\varphi: x \rightarrow x^{2}$, where $x$ is a general element of unit weight in $A_{n}$.

Proof. The multiplication tables for $A_{2}$ and $A_{3}$ are:


By direct computation we find that the corresponding linear spaces $B_{n}$ are induced by the sets of coordinate functions

$$
\begin{aligned}
M_{2} & =\left\{1, x_{1}, x_{2}, x_{1}^{2}\right\} \\
M_{3} & =\left\{1, x_{1}, x_{2}, x_{3}, x_{1} x_{2}, x_{2}^{2}, x_{1}^{2}\right\}
\end{aligned}
$$

respectively, and in each case $\operatorname{dim} B_{n}(n=2,3)$ is the maximum dimension for such a transformation of $A_{n}$; if any of the $\lambda_{i j k}=0(1 \leqslant i, j<n) \operatorname{dim} B_{n}$ would decrease.

Consider $A_{n}$ with the above multiplication table. Now $x_{n} \varphi$ includes the maximum number of quadratic terms which can occur and the induced transformations on these will not introduce any other functions because of the linearity of $x_{i} \varphi$, for $i<n$.

Remark. Any number of train roots may be zero. This will not alter $B_{n}$. The reason is that we are restricted to quadratic functions only.

If we extend the concept of a Bernstein algebra [12] we can define an $n t h$-order Bernstein algebra as one in which equilibrium is reached after exactly $n$ generations of panmixia:

$$
x^{[n+2]}=x^{[n+1]} .
$$

It is seen that the algebras in this section are all second-order Bernstein algebras, $x^{[4]}=x^{[3]}$, when all the train roots apart from 1 are zero.

The Bernstein algebras discussed in [12] would correspond to first-order Bernstein algebras.

The space $B_{n}$ is generated by
and

$$
\left(1+x_{1}+x_{2}+\ldots+x_{n-1}\right)^{2}+x_{n}
$$

$$
\operatorname{dim} B_{n}=\frac{1}{2}\left(n^{2}+n+2\right) .
$$

The following table gives the dimensions of the induced vector spaces for the algebras $A_{n}$ :

$$
\begin{array}{ccccccccc} 
& A_{0} & A_{1} & A_{2} & A_{3} & A_{4} & A_{5} & A_{6} & A_{7} \\
\operatorname{dim} B_{n} & 1 & 2 & 4 & 7 & 11 & 16 & 22 & 29
\end{array}
$$

## 6. Conclusions

Holgate's Theorem 2 [11] is an important theorem as it represents a major step forward in constructing a definitive and fundamental theory of genetic algebras, as we show in a subsequent paper. However, it is incomplete since the eigenvalues of the induced linear transformation explicitly depend on the constructed monomials which are not unique and are basis dependent, and the minimality of $B_{n}$ is only shown with respect to a particular canonical basis. In this paper we prove that $\operatorname{dim} B_{n}$ is an invariant of $A_{n}$ with respect to $\varphi: x \rightarrow x^{2}$; it then follows that the eigenvalues of $\tilde{\varphi}$ are also invariant from the theory in a subsequent paper where linearization is shown to be unique (up to similarity).

By considering the most general case of a genetic algebra we easily obtain some elegant generating polynomials in closed form. This is not generally the case for simpler algebras where many of the structure constants are zero. Methods of defining generating functions for $M_{n}$ and $u_{n}$ are described. The use of diophantine equalities enables us to examine the asymptotic behaviour of $u_{n}$.

The (Schafer) genetic algebras requiring only quadratic functions to linearize them have very simple structures. Again it is more convenient to consider the most general of such algebras.

Although we have assumed the existence of an idempotent (for simplicity), it is easily seen that this assumption can be dispensed with. Our results are clearly valid over the real and complex number fields.

Further, we see that the concept of linearization is a fundamental idea as the class of algebras in which quadratic transformations may be linearized is strictly larger than the class of Schafer algebras. Further, it is a technique which is clearly applicable in a wider context than in genetics and to more general non-linear transformations.

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