# ON THE ASYMPTOTIC DIRECTIONS OF THE S-DIMENSIONAL OPTIMUM GRADIENT METHOD 

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# On the Asymptotic Directions of the s-dimensional Optimum Gradient Method* 

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## Abstract

The optimum s-gradient method for minimizing a positive definite quadratic function $f(x)$ on $E_{n}$ has long been known to converge for s > 1 . For these $\underline{s}$ the author studies the directions from which the iterates $x_{k}$ approach their limit, and extends to s > 1 a theory proved by Akaike for $s=1$. It is shown that $f\left(x_{k}\right)$ can never converge to its minimum value faster than linearly, except in degenerate cases where it attains the minimum in one step.

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## 1. Introduction and Summary.

To minimize a smooth real-valued function $f(x)$ of $n$ real variables, the optimum s-gradient method has been described by Birman [3], Faddeev and Faddeeva [5], Khabaza [8], and others. We here consider the model function $f(x)=\frac{\frac{1}{2}}{} x^{T} A x$, where $A$ is a positive definite matrix. Then each iterate $x_{k}$ is equal to its error. The convergence of the method was proved long ago--see (2.14)--and the question now under study is to find the asymptotic manner in which the iterates $\mathrm{x}_{\mathrm{k}} \rightarrow \theta$, the null vector.

For $s=1$ it was conjectured by Forsythe and Motzkin [7] and proved by Akaike [1]--see (4.12)--that the iterates $\mathrm{x}_{\mathrm{k}}$ converge to $\theta$ by asymptotically alternating between two directions--the "cage" of Stiefel [10]. Thus the convergence of $f\left(x_{k}\right)$ to 0 for $s=1$ is known to be linear, and no faster than linear, for any start $x_{0}$ that is not an eigenvector. Moreover, if coordinates are chosen so that $A$ is a diagonal matrix, then the two asymptotic directions have only two nonzero components. Finally, any direction with only two nonzero components is invariant under two steps of the optimum l-gradient method.

In the present paper the author has extended most of the known results to arbitrary $s>1$. The main result (3.8) shows that the directions of the even iterates $x_{2 k}$ have as a limit set a continuum $R$ (which might be a single direction). Moreover, each direction of $R$ is invariant under two steps of the optimum s-gradient method. Let A be a diagonal matrix. It is then shown in (3.10) that in the optimum s-gradient process $f\left(x_{k}\right)$ converges to 0 no faster than linearly for any initial vector $x_{0}$ with at least $s+1$ nonzero components.

Theorem (4.7) shows that all vectors of $R$ have between $s+1$ and $2 s$ nonzero coordinates, inclusive, Theorem (4.8) says that any direction with $s+1$ nonzero components is invariant under two steps of the method, for any $s$. Examples are shown in Sec. 4 of directions with this invariance and with as many as 2 s nonzero components,

Experimental evidence from computer runs for $s=2$ suggests strongly that $R$ is always a single point, just it has been proved to be for $s=1$. The author conjectures without proof that $R$ is a single point for all s, so that. $\mathrm{x}_{\mathrm{k}} \rightarrow \theta$ in an alternating manner completely analogous to the case with $s=1$ 。

The author is aware that for minimizing quadratic functions $f(x)$ in practice, the conjugate-gradient method of Hestenes and Stiefel (see [5]) may usually be expected to be superior to the optimum s-gradient methods, although Khabaza [8] claims superiority for the 3-gradient method in some cases. For nonquadratic functions $f(x)$ the relative merits of the methods are less clear, The purpose of the present investigation was the intellectual one of trying to understand the asymptotic behavior of the various gradient methods for quadratic functions, The author expects that this information may have some useful application -to the minimization of general smooth functions $f(x)$.

## 2. The Optimum s-gradient Method for Quadratic Functions.

Let $f(x)$ be real for all $x$ in real euclidean $n-s p a c e ~ E_{n}$. Let $f(x)$ assume a minimum value for a unique $x$, which can be taken as $\theta$, the origin of $E_{n}$, without loss of generality in the analysis. The advantage of using $\theta$ is that the iterate $\mathrm{X}_{\mathrm{k}}$ is then also its own error $\mathrm{X}_{\mathrm{k}}-\theta$ as a minimizing vector. We wish to analyze certain asymptotic properties of a class of optimum gradient methods for finding the minimum of $f(x)$.

The simplest $f$ to analyze is the quadratic function

$$
\begin{equation*}
\mathrm{f}(\mathrm{x})=\frac{1}{2} \mathrm{x}^{\mathrm{T}} \mathrm{Ax}, \tag{2.1}
\end{equation*}
$$

where $A$ is a symmetric, positive definite, nonderogatory matrix of order $n$. Moreover, (2.1) represents the local behavior at $\theta$ of $f(x)=f(0)$ for most sufficiently smooth functions f. The author conjectures that the theorems proved below for a quadratic function apply essentially also to any sufficiently smooth function $f$ which is locally like (2.1). In this paper only quadratic functions will be studied. See Daniel [4] for an investigation comparing gradient methods for quadratic and nonquadratic-functions in Hilbert space.,

In the various gradient methods one starts with an arbitrary vector $\mathrm{x}_{\mathrm{O}}$, and computes a sequence $\left\{\mathrm{x}_{\mathrm{k}}\right\}$ converging to $\theta$. We assume all arithmetic to be exact, and round-off error is not considered in this paper.

Let $Z_{k}=$ grad $f\left(x_{k}\right)=A x_{k}$ denote the gradient of $f$ at $x_{k}$. In the optimum l-gradient method [5], $x_{k+l}$ is taken to be the unique
point on the line $L_{1}=\left\{x_{k}+\alpha A x_{k}:-\infty<\alpha<\infty\right\}$ for which $F(\alpha)=$ $f\left(x_{k}+\alpha A x_{k}\right)$ is a minimum．（The existence and uniqueness of $x_{k+1}$ result from the easily proved fact that $F(\alpha)$ is a quadratic function of $\alpha$ with $\mathrm{F}^{\prime \prime}(\alpha)>0$ ．）The line $\mathrm{L}_{1}$ through $\mathrm{x}_{\mathrm{k}}$ is called the line of steepest descent of $f(x)$ at $x_{k}$ 。

For $x \in L_{1}$ ，grad $f(x)=A\left(x_{k}+a A x_{k}\right)=A x_{k}+a A^{2} x_{k}$ ．We there－ fore consider the 2－dimensional plane through $\mathrm{x}_{\mathrm{k}}$ ，

$$
L_{2}=\left\{x_{k}+\alpha_{1} A x_{k}+\alpha_{2} A^{2} x_{k}:-\infty<\alpha_{1}<\infty,-\infty<\alpha_{2}<\infty\right\}
$$

and call it the 2－plane of steepest descent of $f(x)$ at $x_{k}$ ． By extension，for any integer $s(1 \leq s \leq n)$ let

$$
L_{s}=\left\{x_{k}+\sum_{i=1}^{S} \alpha_{i} A^{i} x_{k}:-\infty<\alpha_{i}<\infty \quad \text { (all i) }\right\}
$$

be the s－dimensional plane of steepest descent of $f(x)$ at $x_{k}$ 。 Since $A$ is not derogatory，$A x_{k}, \ldots, A^{n} x_{k}$ are linearly independent，provided $\mathrm{x}_{0}$ is a vector whose minimum polynomial is of degree n ．In that case $\mathrm{L}_{\mathrm{n}}$ is the whole space $\underset{\mathrm{n}}{\mathrm{E}}$ 。

In the optimum s－gradient method［5］for minimizing the quadratic function $f$ of（2．1），the point $x_{k+1}$ is defined to be the unique point $y$ in $L_{s}$ for which $f(y)$ is a minimum（ $k=0,1, \ldots$ ）。（Again existence and uniqueness follow from the positive definiteness of $A$ ．） It is the optimum s－gradient methods that we shall analyze in this paper．

We now give two representations of the minimizing $\left\{\alpha_{i}\right\}$ which are useful in the analysis．Actual computing algorithms for the optimum
s-gradient method often proceed differently, and find $x_{k+1}$ by taking $s$ steps of the conjugate gradient method, starting from $x_{k}$. See [5]. We concentrate on the gradients $z_{k}=A x_{k}$.

## First representation

Let

$$
x_{k+1}=x_{k}+\gamma_{l} A x_{k}+\ldots .+\gamma_{s} A^{S} x_{k} .
$$

Then the gradient of $f(x)$ at $x_{k+1}$ is

$$
\begin{equation*}
z_{k+1}=z_{k}+\gamma_{l} A z_{k}+\ldots+\gamma_{S} A^{S} z_{k} \tag{2.2}
\end{equation*}
$$

Since $x_{k+1}$ minimizes $f(y)$ for $y \in L_{s}$, the vector $z_{k+1}$ must be orthogonal to $L_{S}$. For this it is necessary and sufficient that $z_{k+1}$ be orthogonal to $z_{k}, A z_{k}$, . . . $A^{s-1} z_{k}$. Then $\gamma_{1}$, . . . 7, are determined by the s conditions

$$
\left\{\begin{array}{l}
\left(z_{k}, z_{k+1}\right)=\left(z_{k}, z_{k}\right)+\gamma_{l}\left(z_{k}, A z_{k}\right)+\ldots+\gamma_{s}\left(z_{k}, A^{S} z_{k}\right)=0 \\
\quad \cdots \\
\left(A^{s-1} z_{k}, z_{k+1}\right)=\left(A^{S-1} z_{k}, z_{k}\right)+\gamma_{l}\left(A^{S-1} z_{k}, A z_{k}\right)+\ldots+\gamma_{S}\left(A^{S-1} z_{k}, A^{S} z_{k}\right)=0 .
\end{array}\right.
$$

Here $(u, v)=u^{T} v+v^{T} u$ denotes the inner product of two column vectors. Since $\left(A^{p}, A^{q} z\right)=\left(z, A^{p+q_{z}}\right)=z^{T} A^{p+q_{z}}$, we may write the above equations as

$$
\left\{\begin{array}{c}
z_{k}^{T} z_{k}+\gamma_{l} z_{k}{ }^{T} A z_{k}+\ldots+\gamma_{s} z_{k} T_{A} z_{k}=0 \\
z_{k}{ }^{T} A z_{k}+\gamma_{l} z_{k}{ }^{T} A^{2} z_{k}+\ldots+\gamma_{s} z_{k}{ }^{T} A^{s+1} z_{k}=0  \tag{2.3}\\
\cdots \cdots \\
z_{k} T_{A}{ }^{s-1} z_{k}+\gamma_{l} z_{k}{ }^{T} A^{s} z_{k}+\ldots+\gamma_{s} z_{k}{ }^{T} A^{2 s-1} z_{k}=0
\end{array}\right.
$$

As long as $z_{k}, A z_{k}$, . . . $A^{s-1} z_{k}$ are linearly independent, the equations (2.3) determine the minimizing $\gamma_{1}$, . . . $\gamma_{s}$ uniquely.

## Second representation

Let $q_{S}(t)=t^{s}+\beta_{s-1} t^{s-1}+\ldots+\beta_{0}$ denote any monic polynomial of degree s, with $\beta_{0} \neq 0$. Then

$$
q_{S}(A) z_{k}=A^{S} z_{k}+\beta_{S-1} A^{s-1} z_{k}+\ldots+\beta_{0} z_{k}
$$

and

$$
\begin{equation*}
\frac{q_{S}(A) z_{k}}{q_{S}(0)}=\frac{1}{\beta_{0}} A^{s} z_{k}+\frac{\beta_{S-1}}{\beta_{0}} A^{s-1} z_{k}+\ldots+z_{k} . \tag{2.4}
\end{equation*}
$$

Comparing (2.4) with (2.2), we see that we can write

$$
\begin{equation*}
z_{k+1}=\frac{p_{S}(A)}{p_{S}(0)} z_{k}, \tag{2.5}
\end{equation*}
$$

where

$$
\mathrm{p},(\mathrm{t}) \text { is the particular polynomial }
$$

$$
\begin{equation*}
\mathrm{p},(\mathrm{t})=\mathrm{t}^{\mathrm{s}}+\frac{\not \Varangle_{\mathrm{s}-1}}{\gamma_{\mathrm{s}}} t^{\mathrm{s}-1} \quad \ldots+\frac{y_{1}}{\gamma_{s}} t+\frac{1}{\gamma_{s}} . \tag{2.6}
\end{equation*}
$$

Now p,(t) is a certain orthogonal polynomial. Without loss of generality assume $A$ to be the diagonal matrix

where $0<\lambda_{\rceil_{\perp}}<\lambda_{\rho}<. \quad .<\lambda_{r_{\perp}}$ are its eigenvalues (distinct because $A$ is not derogatory).
(2.8) Definition. In the coordinate system corresponding to (2.7), let the nonzero vector $z$ be $\left(\zeta_{1}, \ldots, \zeta_{n}\right)$ ! Let orthogonality of two polynomials $p(t), \quad q(t)$ (relative to $z$ ) be defined by

$$
\langle p(t), q(t))=\sum_{i=1}^{n} p\left(\lambda_{i}\right) q\left(\lambda_{i}\right) \zeta_{i}^{2}=0 .
$$

(2.9) Definition. Let $P_{S}(t ; z)=t^{s}+\ldots$ be the unique monic polynomial of degree $s$ that, relative to $z$, is orthogonal in the sense of (2.8) to all polynomials of degree $\leq s-1$.

Note that $P_{S}(t ; z)$ depends only on the direction of $z$, and not its magnitude. I.e., $P_{S}(t ; z)=P_{S}(t ; \alpha z)$, for all real $\alpha \neq 0$.
(2.10) Theorem. The polynomial $p,(t)$ of (2.5), (2.6) is the_orthoponally nomial $P_{S}\left(t ; z_{k}\right)$ defined in (2.9).

We shall not prove (2.10). For a related proof see, for example, p. 349 of [5]. The basic reason for (2.10) is the isomorphism, well expounded by Stiefel [ll], between orthogonality of the polynomials
$p(t), \quad q(t)$ in the sense of (2.8) and geometric orthogonality of the vectors $p(A) z, q(A) z$ in $E_{n}$. That is,

$$
(p(t), q(t)\rangle=(p(A) z, q(A) z)
$$

Hence the conditions (2.3) asserting the orthogonality of the vector $z_{k+1}=P_{s}\left(A ; z_{k}\right) z_{k} / P_{s}\left(0 ; z_{k}\right)$ to $z_{k}, A z_{k}, A^{2} z_{k}, \ldots, A^{s-1} z_{k}$ are equivalently asserting the orthogonality of the polynomial $P_{S}\left(t ; z_{k}\right)$ to the polynomials $1, t, t^{2}, \ldots, t^{s-1}$.

In summary $z_{k+1}$ is uniquely determined from $z_{k}$ by the formula

$$
\begin{equation*}
z_{k+1}=\frac{P_{S}\left(A ; z_{k}\right)}{P_{S}\left(0 ; z_{k}\right)} z_{k} \tag{2.11}
\end{equation*}
$$

Moreover,

$$
\begin{equation*}
x_{k+1}=\frac{P_{s}\left(A ; z_{k}\right)}{P_{s}\left(0 ; z_{k}\right)} x_{k} \tag{2.12}
\end{equation*}
$$

Relation (2.12) is the basis for a proof by Birman[3] that in the optimum s-gradient method $f\left(x_{k}\right)$ converges to 0 linearly, or faster. To be precise, let $\sigma=\left(\lambda_{n}+\lambda_{1}\right)\left(\lambda_{n}-\lambda_{1}\right)^{-1}$. Let $T_{S}(t)$ denote the Chebyshev polynomial on $[-1,1]$, normalized so that $\max _{-1}<t<{ }_{I}\left|T_{S}(t)\right|=$ 1. Let

$$
Q_{s}(u)=T_{s}\left(\frac{\lambda_{n}+\lambda_{1}-2 u}{\lambda_{n}-\lambda_{l}}\right) .
$$

Then $Q_{S}(0)=T_{S}(\sigma)>1$, and $\left|Q_{S}(t)\right| \leq 1$, for $\lambda_{1} \leq t \leq \lambda_{n}$. It is

Birman's proof goes as follows:
(2.13)

$$
\begin{aligned}
f\left(x_{k+1}\right) & =f\left(\frac{P_{s}\left(A ; z_{k}\right)}{P_{s}\left(0 ; z_{k}\right)} x_{k}\right) \\
& \leq f\left(\frac{Q_{S}(A)}{Q_{s}(0)} x_{k}\right), \begin{array}{c}
\text { because } P_{s}\left(t ; z_{k}\right) \text { is the poly- } \\
\text { nomial that minimizes } f\left(x_{k+1}\right)
\end{array} \\
& =\frac{1}{\left[Q_{s}(0)\right]^{2}} x_{k}^{T} Q_{s}(A) A Q_{s}(A) x_{k} \\
& =\frac{1}{\left[Q_{s}(0)\right]^{2}} \sum_{i=1}^{n} \lambda_{i}\left[Q_{s}\left(\lambda_{i}\right)\right]^{2}\left[\xi_{i}(k)\right]^{2}
\end{aligned}
$$

$$
\leq \frac{1}{\left[Q_{s}(0)\right]^{2}} \sum_{i=1}^{n} \lambda_{i}\left[\xi_{i}^{(k)}\right]^{2}
$$

$$
=\frac{1}{\left[T_{S}(\sigma)\right]^{2}} f\left(x_{k}\right)
$$

Hence

$$
\begin{equation*}
\sqrt{f\left(x_{k}\right)} \leq \frac{1}{\left[T_{S}(\sigma)\right]^{k}} \sqrt{f\left(x_{0}\right)} \tag{2.14}
\end{equation*}
$$

proving the convergence of $f\left(x_{k}\right)$ to 0 to be linear or faster.
(2.15) Definition. For $\alpha=0, \pm 1, \pm 2, \ldots$, let the moments $\mu_{\alpha}$ of $z=\left(\zeta_{1}, \ldots, \zeta_{n}\right)^{T}$ be defined by

$$
\mu_{\alpha}=\sum_{i=1}^{n} \lambda_{i}^{\alpha} \zeta_{i}^{2}
$$

(2.16) Theorem. Fix $s>1$. Except for a constant factor, the orthogonal polynomial $P_{S}(t ; z)$ of (2.9) can be expressed by the determinant

$$
P_{s}(t ; z)=\left|\begin{array}{cccc}
\mu_{0} \mu_{1} & \ldots & \mu_{s-1} & 1  \tag{2.17}\\
\mu_{1} \mu_{2} & \ldots & \mu_{s} & t \\
\cdots & \ldots & . & \\
\mu_{s} \mu_{s+1} & & \mu_{2 s-1} & t^{s}
\end{array}\right|
$$

The proof is left to the reader.
In the next theorem we give an explicit representation for the ratio $f\left(x_{k+1}\right) /\left(f\left(x_{k}\right)\right.$ in terms of the moments of ${ }^{z} k$.
(2.18) Theorem. Fix $s>1$. Let $x_{k}$ be any vector in the optimum s-gradient method, and let $\mu_{\alpha}$ be the moments defined by (2.15) for the gradient vector $z_{k}=A x_{k}$. Then

$$
\frac{f\left(x_{k+1}\right)}{f\left(x_{k}\right)}=\frac{\left\lvert\, \begin{array}{cccc}
\mu_{-1} & \mu_{0} & \mu_{1} & \cdots \\
\mu_{s-1} & \mu_{s-1} \\
\mu_{0} & \mu_{1} & \mu_{2} & \cdots
\end{array}\right.}{\mu_{s}} \begin{aligned}
& \mu_{s} \\
& \mu_{s+1} \\
& \mu_{-1} \\
& \mu_{-1}
\end{aligned} \mu_{2 s-1} .
$$

where $M_{-1}$ i.s the mingr determinant of $\mu_{-1}$ in the above determinant:

$$
M_{-1}=\left|\begin{array}{ccc}
\mu_{1} \mu_{2} & \ldots \mu_{s} \\
\mu_{2} \mu_{3} & \ldots * \mu_{s+1} \\
\cdot & \ldots & \\
\mu_{s} \mu_{s+1} & \ldots & \mu_{2 s-1}
\end{array}\right|
$$

Proof. We have $2 f\left(x_{k}\right)=x_{k}{ }^{T} A x_{k}=z_{k}{ }^{T} A^{-1} z_{k}=\mu_{-1} . \quad$ To simplify the notation, let $z_{k}=\left(\zeta_{1}, \ldots \zeta_{n}\right)^{T}$ and $z_{k+1}=\left(\zeta_{1}, \text {, } \zeta_{n}^{\prime}\right)^{T}$. Then

$$
\begin{align*}
\zeta_{i}^{\prime} & =\frac{P_{s}\left(\lambda_{i} ; z_{k}\right)}{P_{s}\left(0 ; z_{k}\right)} \zeta_{i}  \tag{2.11}\\
& =\frac{P_{s}\left(\lambda_{i} ; z_{k}\right)}{(-1)^{S} M_{-1}} \zeta_{i}
\end{align*}
$$

where we use the representation (2.17) for $P_{S}\left(t ; z_{k}\right)$. Then
(2.19) $2 f\left(x_{k+1}\right)=z_{k+1} T_{A}^{-1} z_{k+1}=\frac{1}{M^{2}} \sum_{i=1}^{\sim}\left[P_{S}\left(\lambda_{i} ; z_{k}\right)\right]^{2} \frac{1}{\lambda_{i}} \zeta_{i}^{2}$

$$
\frac{1}{M_{-1}^{2}} \sum_{i=1}^{n} P_{s}\left(\lambda_{i} ; z_{k}\right) \frac{P_{s}\left(\lambda_{i} ; z_{k}\right)}{\lambda_{i}} \zeta_{i}^{2}
$$

-Now $P_{S}\left(t ; z_{k}\right)$ is orthogonal-in the sense of (2.8) to all polynomials of degree $\leq s-1$. Hence the only term of $P_{s}\left(\lambda_{i} ; z_{k}\right) / \lambda_{i}$ that contributes anything nonzero to the sum (2.19) is the term $(-1)^{S_{M}}{ }_{-1} / \lambda_{i}$. Hence

$$
2 f\left(x_{k+1}\right)=\frac{(-1)^{s}}{M_{-1}} \sum_{i=1}^{n} P_{s}\left(\lambda_{i} ; z_{k}\right) \zeta_{i}^{2} / \lambda_{i}
$$

$$
\begin{aligned}
& =\frac{(-1)^{s}}{M_{-1}}\left|\begin{array}{lllll}
\mu_{0} & \mu_{1} & \cdots & \mu_{S-1} & \mu_{-1} \\
\mu_{1} & \mu_{2} & \cdots & \mu_{S} & \mu_{0} \\
\mu_{S} & \mu_{S+1} & \cdots & \mu_{2 S-1} & \mu_{S-1}
\end{array}\right|
\end{aligned}
$$

Dividing $2 f\left(x_{k+1}\right)$ by $2 f\left(x_{k}\right)=\mu_{-1}$ and rearranging the columns of the last determinant proves theorem (2.18).
(2.20) Corollary. In the notation of theorem (2.18), for $s=1$,

$$
\begin{equation*}
\frac{f\left(x_{k+1}\right)}{f\left(x_{k}\right)}=\frac{\mu_{1} \mu_{-1}-\mu_{0}^{2}}{\mu_{1}^{-}-1} \tag{2.21}
\end{equation*}
$$

If $n=2$ and $s=1$, then
(2.22) $\frac{f\left(x_{k+1}\right)}{f\left(x_{k}\right)}=\frac{\zeta_{1}^{2} \zeta_{2}^{2}\left(\lambda_{2}-\lambda_{1}\right)^{2}}{\left(\lambda_{2} \zeta_{1}^{2}+\lambda_{1} \zeta_{2}^{2}\right)\left(\lambda_{1} \zeta_{1}^{2}+\lambda_{2} \zeta_{2}^{2}\right)}=c^{2}=c^{2}\left(x_{k}\right)$.

Proof. The second expression comes from the first by using (2.15) and (2.21), where $z_{k}=\left(\zeta_{1}, \zeta_{2}\right) \frac{T}{1}$ with some algebraic manipulation.
(2.23) Corollary. The expression (2.22) for $f\left(x_{k+1}\right) / f\left(x_{k}\right)$ is unchanged, if $\left(\zeta_{1}, \zeta_{2}\right)^{T}$ is changed to $\left(\zeta_{2},-\zeta_{1}\right)^{T}$.

The inequality (2.13) yields an upper bound for the expression in (2.18) . We may state this result in the form of the following inequality, valid for $s=1,2, \quad . \quad$.
(2.24) $\left|\begin{array}{ccccc}\mu_{-1} & \mu_{0} & \mu_{1} & \ldots & \mu_{s-1} \\ \mu_{0} & \mu_{1} & \mu_{2} & \ldots & \mu_{s} \\ & \ldots & \ldots & \\ \mu_{s-1} & \mu_{s} & \mu_{s+1} . & * * & \mu_{2 s-1-}\end{array}\right|<\left|\begin{array}{ccc}\mu_{1} & \ldots & \mu_{s} \\ \ldots & \ldots \\ \mu_{s} & \ldots & \mu_{2 s-1}\end{array}\right|<\frac{1}{\left[T_{s}\left(\frac{\lambda_{n}+\lambda_{1}}{\lambda_{n}-\lambda_{1}}\right)\right]^{2^{*}}}$

This is essentially the inequality of Meinardus [Ba], who derived it by the same argument for a slightly different iteration in which $\|x\|^{2}$ is minimized instead of $f(x)$.

The special case for $s=1$,

$$
\left|\begin{array}{cc}
\mu_{-1} & \mu_{0}  \tag{2.25}\\
\mu_{0} & \mu_{1}
\end{array}\right| / \mu_{-1} \mu_{1} \leq \frac{1}{\left[\mathbb{T}_{1}\left(\frac{\lambda_{n}+\lambda_{1}}{\lambda_{n}-\lambda_{1}}\right)\right]^{2}}=\left(\frac{\lambda_{n}-\lambda_{1}}{\lambda_{n}+\lambda_{1}}\right)^{2}
$$

is a well-known inequality of Kantorovich; see (8) on p. 410 of [5].
It was stated by. Birman [3] that the bound (2.14) is sharp, in the sense that for each $s$ and each given $\lambda_{1}, \lambda_{n}(s<n)$, one can find $A$ and $x_{0}$ so that (2.14) is an equality for all $k$. This is done by finding a set of $\lambda_{i}$ and $\xi_{i}{ }^{(0)}$ so that the shifted Chebyshev polynomial $Q,(t)$ is (up to a scalar factor) identical with $P_{s}\left(t ; z_{0}\right)$ and so that $\left|Q_{s}\left(\lambda_{i}\right)\right|=1$ for each eigenvalue $\lambda_{i}$. This is known to be possible because the Chebyshev polynomials, like cosines, are orthogonal with respect to summation over certain points.

However, Birman did not investigate the actual manner or rate of convergence of $f\left(x_{k}\right)$ to 0 in the optimum s-gradient method for a general given $A$ and $x_{0}$. He left open the question of whether the convergence of $f\left(x_{k}\right)$ to 0 might actually be faster than linear in certain nontrivial cases.

For $s=1$ Forsythe and Motzkin [7] conjectured that if ${ }_{S_{1}}^{(0)} n_{n}(0)$ $\neq 0$, then $\left(k_{\xi_{i}}=\right.$ ○ $\left(\left\|x_{k}\right\|\right)$, as $k \rightarrow \infty$, for all i with $1<i<n$. In words, $x_{k} \rightarrow \theta$ asymptotically in the 2 -space $\pi_{l, n}$ spanned by the eigenvectors belonging to $\lambda_{1}$ and $\lambda_{\mathrm{n}}$. The conjecture was proved by Forsythe and Motzkin (unpublished) only for $n=3$. Akaike [l] proved the conjecture for arbitrary $n$. In an unpublished manuscript Arms [2] had found a similar proof. We give a proof in (4.12) as a consequence of our result (3.8) for the s-gradient method.

Suppose the optimum l-gradient process is performed entirely in the two-dimensional space $\pi_{1, n}$. Then, if $x_{0} \in \pi_{1, n}$ and $x_{0}$ is not an eigenvector, it is easy to prove that:
(i) $x_{0}, x_{2}, x_{4}, \ldots$ are all collinear vectors, and that $x_{1}, x_{3}, x_{5}, \ldots$ are also collinear in another direction. Furthermore, $x_{2 k+2}=c^{2} x_{2 k}$ and $x_{2 k+1}=c^{2} x_{2 k-1}$, for all $k$. Here $c^{2}$ is given by (2.22). The basic reason why these vectors are collinear is that the gradients $z_{k+1}$ and $z_{k}$ must always be perpendicular in any optimum gradient method.
(ii) Moreover, for each $k=0,1, . . . \quad f\left(x_{k+1}\right)=c^{2} f\left(x_{k}\right)$. This is an immediate consequence of Corollary (2.23). Hence $f\left(x_{k}\right)>0$ in a strictly linear fashion, like the $k$-th term of a convergent geometric series, even though the vectors $\mathrm{x}_{\mathrm{k}}$ alternate between two fixed directions.

It is a consequence of the Forsythe-Motzkin-Arms-Akaike result on the manner of convergence of $x_{k}$ to $\theta$ in $E_{n}$ for $s=1$ that the iteration behaves asymptotically, as $k \rightarrow \infty$, as though it were entirely in the two-space $\pi_{1, n}$. The vectors $x_{k}$ behave ultimately as though they had resulted from an iteration started with some $x_{0}^{*}$ in $\pi_{1}, n$ 。 In particular, we find that $f\left(x_{k}\right) \searrow 0$ linearly, in the sense that

$$
\lim _{k \rightarrow \infty} \frac{f\left(x_{k+1}\right)}{f\left(x_{k}\right)}=c^{2}\left(x_{0}^{*}\right)
$$

However, the vector $\mathrm{x}_{0}^{*}$ is an extremely complex function of $\mathrm{x}_{\mathrm{O}}$ 。
Till now, the asymptotic nature of the optimum s-gradient method has not been described for $s>1$. This problem, posed on $\mathrm{p}, 314$ of Forsythe [6], is studied in the next section

## 3. Asymptotic Behavior of the s-gradient Method.

We are still assuming A to have distinct positive eigenvalues $\lambda_{1}<\lambda_{2}<\cdot . \quad .<\lambda_{n}$. Fix any $s$ with $1<s$. Motivated by (2.11) and by Akaike's approach [I] for $s=1$, we shall consider the transformation

$$
\begin{equation*}
W^{\prime}=P_{S}(A ; w) w \tag{3.1}
\end{equation*}
$$

Here $w \neq \theta$ and $P_{S}(t ; w)=t^{s}+\ldots$ is the orthogonal polynomial defined in (2.9). Let

$$
\begin{equation*}
\varphi(w)=\frac{\left\|w^{\prime}\right\|^{2}}{\|w\|^{2}}=\frac{\left\|P_{s}(A ; w)\right\|^{2}}{\|w\|^{2}} \tag{3.2}
\end{equation*}
$$

where $\|u\|$ denotes the euclidean length of $u$.

$$
\begin{gathered}
\text { Similarly, if } w^{\prime} \neq \theta, \text { let } w^{\prime \prime}=P_{S}\left(A ; w^{\prime}\right) w^{\prime}, \text { so that } \\
\varphi\left(w^{\prime}\right)=\frac{|M|^{2}}{|M N|^{2}} .
\end{gathered}
$$

The following theorem is of basic importance to our analysis of the asymptotic behavior of the s-gradient method.
(3.3) Theorem. Let $\psi$ be the angle between $w$ and $w^{\prime \prime}$ • For any w such that $w^{\prime \prime} \neq \theta$, we have

$$
\varphi(w)=\frac{\| \|^{2} \|^{2}}{\|w\|^{2}}=\cos ^{2} \psi \frac{\left\|w^{\prime \prime}\right\|^{2}}{\left\|w^{\prime}\right\|^{2}} \leq \frac{\left\|w^{\prime \prime}\right\|^{2}}{\left\|w^{\prime}\right\|^{2}}=\varphi\left(w^{\prime}\right)
$$

and there is equality if and only if $w^{\prime \prime}=\mathrm{cw}$ for some scalar $\mathrm{c}>0$.
Proof. By the Cauchy-Schwarz inequality and the definition of $\psi$,

$$
\begin{equation*}
\left(w^{T} w^{\prime}\right)^{2}=\cos ^{2} \psi\|w\|^{2}\left\|w^{\prime \prime}\right\|^{2} \leq\|w\|^{2}\left\|w^{\prime \prime}\right\|^{2} \tag{3.4}
\end{equation*}
$$

with equality if and only if $w=c w "$, for $c \neq 0$.
Now

$$
\begin{aligned}
\left.|M|\right|^{2}-w^{T} w^{\prime \prime} & =\left\|P_{S}(A ; w) w\right\|^{2}-w^{T} P_{S}\left(A ; w^{:}\right) P_{S}(A ; w) w \\
& =w^{T}\left[P_{S}(A ; w)\right]^{2} w-w^{T} P_{S}\left(A ; w^{\prime}\right) P_{S}(A ; w) w \\
& =w^{T} P_{S}(A ; w)\left\{P_{S}(A ; w)-P_{S}\left(A ; w^{\prime}\right)\right\} w \\
& =w^{T} P_{S}(A ; w) D(A) w \\
& =0,
\end{aligned}
$$

by (2.3), because $D(t)$ is a polynomial of degree at most $s-1$, since the leading terms $t^{s}$ cancel. Hence $\left\|^{\prime}\right\|^{2}=\mathcal{W}^{\prime} w^{\prime \prime}$, whence

$$
\left\|w^{\prime}\right\|^{4}=\left(w^{T} w^{\prime \prime}\right)^{2}
$$

Combining (3.4) with(3.5), we have

$$
\left|W_{w}\left\|^{4}=\cos ^{2} \psi\right\| w\left\|^{2}\right\| w^{\prime \prime}\left\|^{2} \leq\right\| w\left\|^{2} \mid M\right\|^{2},\right.
$$

with equality if and only if $w^{\prime \prime}=\mathrm{cw}$. That $\mathrm{c}>0$ follows from the fact that $w^{\mathbb{T}^{n}}=\left\|w^{\mathbf{r}}\right\|^{<\hat{>}} 0$. This proves theorem (3.3).
（3．6）Definition．Fix $s$ with $1 \leq s \leq n-1$ ．Fix a euclidean coordinate system in $\mathrm{E}_{\mathrm{n}}$ so A takes the form（2．7）．Let $\Sigma$ be the unit sphere in $E_{n}$ 。 Define $\Sigma^{*} \subset \Sigma$ to consist of all unit vectors $y$ with at least $s+l$ nonzero components．We define a transformation $\mathrm{T}: \quad \Sigma^{*} \rightarrow \Sigma^{*}$ ，as follows：For each y in $\Sigma^{*}$ ，let $\mathrm{y}^{\prime}=\mathrm{Ty}=\mathrm{w} /\|\mathrm{w}\|$ ， where $w=P_{S}(A ; y) y . \quad$（That $w \neq \theta$ and $y^{\prime} \in \Sigma^{*}$ are proved in theorem（5．1）．）
（3．7）Definition．By a continuum we mean a closed connected set in $\mathrm{E}_{\mathrm{n}}$ ， with the understanding that a single point is a continuum．
（3．8）Theorem，Fix $s$ with $1_{-} \leq s_{-} \leq n-1$ ．Let $y_{0}=\left(\eta_{1}^{(0)}, \ldots, \eta_{n}^{(0)}\right)^{T}$ be any vector in $\Sigma^{*} \xrightarrow{(0)} \underline{w \perp t} \neq \eta_{1}^{0}(i=1, \ldots, n) \quad$ For $k=0,1, \ldots$ ， define $y_{k+1}=T y_{k}$ ，where was defined in（3．6）．Then the set of limit points of the sequence $\left\{y_{2 k}: k=0,1,2, \ldots\right\}$ of normalized gradients is a continuum $R \subset \Sigma^{*}$ 。 Moreover，for any point $r$ in $R$ ， we have $r=T^{2} r=T(T r)$ ．

Proof，Let $w_{0}=y_{0}$ ．For $k=0,1, \ldots$, let $w_{k+1}=P_{s}\left(A ; y_{k}\right) w_{k}$ ， －where $P_{s}(t ; y)$ was defined in（2．9）．It is easily shown that $y_{k}=$ $\mathrm{w}_{\mathrm{k}} /\left\|\mathrm{w}_{\mathrm{k}}\right\|$ ，for all k ．Since $\mathrm{n} \geq \mathrm{s}+1$ components of $\mathrm{w}_{0}$ are nonzero， it follows from theorem（5．1）that at least $s+l$ components of $w_{k}$ are nonzero for $k=1,2, \ldots$ ．Hence no $w_{k}=\theta$ 。

Let $w_{k}=\left(\omega_{1}^{(k)}, \ldots, \omega_{n}^{(k)}\right)^{T}$ ．By theorem（3．3），

$$
\varphi\left(w_{0}\right) \leq \varphi\left(w_{\perp}\right) \leq \ldots \leq \varphi\left(w_{k}\right) \leq \ldots 。
$$

But for each $k$ the $s$ zeros of $P_{s}\left(t ; w_{k}\right)$ lie in the interval $\left(\lambda_{1}, \lambda_{n}\right)$. Hence $\left|P_{s}\left(t ; w_{k}\right)\right| \leq\left(\lambda_{n}-\lambda_{1}\right)^{s}$, for $\lambda_{1} \leq t \leq \lambda_{n}$, and so

$$
\begin{aligned}
\varphi\left(w_{k}\right) & =\frac{\left\|w_{k+1}\right\|^{2} \| P_{s}}{\|\left. w_{k}^{\prime}\right|^{2}}=\frac{\left(A ; w_{k}\right) \|^{2}}{\|\left. w_{k}^{\prime}\right|^{2}} \\
& =\frac{\sum_{i=1}^{n}\left[P_{s}\left(\lambda_{i} ; w_{k}\right)\right]^{2}\left[\omega_{i}^{(k)}\right]^{2}}{\sum_{i=1}^{n}\left[\omega_{i}^{(k)}\right]^{2}} \\
& \leq\left(\lambda_{n}-\lambda_{-1}\right)^{s}, \text { for all } k
\end{aligned}
$$

As a monotone bounded sequence, $\left\{\varphi\left(w_{k}\right)\right\}$ has a limit $L$. Hence

$$
\begin{equation*}
\varphi\left(w_{k+1}\right)-\varphi\left(w_{k}\right) \rightarrow 0 \quad(\text { as } k \rightarrow \infty) \tag{3.9}
\end{equation*}
$$

But, by theorem (3.3),

$$
\varphi\left(w_{k+1}\right)-\varphi\left(w_{k}\right)=\frac{\left\|w_{k+2}\right\|^{2}}{\left\|w_{k+1}\right\|^{2}}-\frac{\left\|w_{k+1}\right\|^{2}}{\left|w_{k}\right|^{2}}
$$

(3.10)

$$
=\frac{\left\|w_{k+2}\right\|^{2}}{\left\|w_{k+1}\right\|^{2}}\left[1-\cos ^{2} \psi_{k}\right]
$$

where $\psi_{k}$ is the angle between $w_{k}$ and $w_{k+2}$. Then, by (3.9), $\cos ^{2} \psi_{k} \rightarrow 1$, and $\psi_{k} \rightarrow 0$, as $k \rightarrow \infty$. (Since $c>0$ in (3.3), $\psi_{k} \neq \pi 0$ (1)

Now consider the set $Y$ of unit vectors $\left\{y_{2 k}: k=0,1,2, \ldots\right\}$. As an infinite subset of the compact unit sphere $\Sigma,\left\{y_{2 k}\right\}$ has limit points; let $R$ be the set of all limit points of $Y$. Since $\psi_{k} \rightarrow 0$,
as $k \rightarrow \infty$, we have $\left\|y_{2 k+2}-y_{2 k}\right\| \rightarrow 0$, as $k \rightarrow \infty$. Then, as Ostrowski shows on $p .203$ of [9], the set $R$ must be a continuum in the sense of (3.7).

Let $r$ be any point of $R$. Then there is a subsequence $\left(y_{2 k .1}\right\}$ converging to r . Since $\left\|\mathrm{y}_{2 k_{1}+2}-\mathrm{y}_{2 k_{i}}\right\| \rightarrow 0$, we have also that $\mathrm{y}_{2 \mathrm{k}_{i}+2}=\mathrm{T}^{2} \mathrm{y}_{2 k_{i}} \rightarrow \mathrm{r}$. But T is a continuous transformation. Hence $T^{2}{ }^{2}{ }_{2 k_{i}} \rightarrow T^{2} r$, and $T^{2} r=r$. Since $T^{2} r=r$, we see from theorem (5.1) that $r \in \Sigma^{*}$. Hence $R \subset \Sigma^{*}$. This completes the proof of theorem (3.8).

The author has programmed a number of test cases with $s=2$, to investigate the nature of the set $R$. In every case, $R$ appeared to be a single point. The author conjectures that $R$ is always a single point in theorem (3.8). So far, this has been proved only for $s=1$, and we give the proof in (4.12).

The following theorem shows one way in which one might be able to prove that $R$ consists always of a single point.
(3.11) Theorem, Suppose in the proof of theorem (3.8) that $\varphi\left(w_{k}\right)$ were to converge to $L$ so rapidly that, for some $\alpha<1$.

$$
\begin{equation*}
0 \leq \varphi\left(w_{k+1}\right)-\varphi\left(w_{k}\right) \leq \alpha\left[\varphi\left(w_{k}\right)-\varphi\left(w_{k l}\right)\right], \quad \text { for all } k . \tag{3.12}
\end{equation*}
$$

Then $R$ would consist of a-single point.

Proof. If (3.12) held, then the following infinite series would be convergent:

$$
\begin{equation*}
\sum_{1}^{\infty}\left[\varphi\left(w_{k+1}\right)-\varphi\left(w_{k}\right)\right]^{\frac{1}{2}}<\infty \tag{3.13}
\end{equation*}
$$

as is seen from (3.12), by the ratio test. It is shown in (3.10) that

$$
\begin{equation*}
\left[\varphi\left(w_{k+1}\right)-\varphi\left(w_{k}\right)\right]^{\frac{1}{2}} \sim \sin \left|\psi_{k}\right|, \quad \text { as } k \rightarrow \infty \tag{3.14}
\end{equation*}
$$

where $\psi_{k}$ is the angle between the vectors $w_{k}$ and $w_{k+2}$. Then, from (3.13) and (3.14), we would have

$$
\begin{equation*}
\sum_{4}^{\infty}\left|\psi_{k}\right|<\infty \tag{3.15}
\end{equation*}
$$

Now, let $y_{k}=w_{k} /\left\|w_{k}\right\|$ be the unit vector in the direction of $w_{k}$. It would follow from (3.15) that

$$
\sum_{k=0}^{\infty}\left\|y_{2 k+2}-y_{2 k}\right\|<\infty,
$$

whence

$$
\begin{equation*}
\sum_{k=0}^{\infty}\left(y_{2 k+2}-y_{2 k}\right) \tag{3.16}
\end{equation*}
$$

would be an absolutely convergent series of vectors. Since

$$
y_{2 k}=\sum_{h=0}^{k-1}\left(y_{2 h+2}-y_{2 h}\right)+y_{0}
$$

we see that the sequence $\left\{\mathrm{y}_{2 \mathrm{k}}\right\}$ would then have one limit point. This proves the theorem (3.11).

However, the author sees no way to prove (3.12) nor the conjecture. The following theorem proves that, whether $R$ has one point or an infinite number, $f\left(x_{k}\right) \rightarrow 0$ no faster than linearly.
(3.17) Theorem. Fix $s$ with $1 \leq s \leq n-1$. Given any $A$ in the form (2.7). Let $x_{0}=\left(\xi_{1}^{(0)}, \ldots, \xi_{n}^{(0)}\right)^{T}$ be any vector in $E_{n}$ with $m$ nonzero components. Then in the $f\left(x_{k}\right)$ mum $s$-gradient method converges to 0 in the following ways:
(i) If $m \leq s$, then $x_{1}=\theta, f\left(x_{1}\right)=0$, and the iteration terminates in one step.
(ii) If $s+l \leq m$, then the convergence of $f\left(x_{k}\right)$ to 0 is asymptotically linear, in the sense that there exist constants $c_{1}, c_{2}$ depending on $x_{0}$, with

$$
\begin{equation*}
0<c_{1} \leq \frac{f\left(x_{2 k+2^{\prime}}\right)}{f\left(x_{2 k}\right)} \leq c_{2}<1^{\prime} \quad \text { for all } k \tag{3.18}
\end{equation*}
$$

Proof. We may ignore any zero components of $x_{0}$, as they remain zero throughout the iteration. We are thus minimizing $f(x)$ in $E_{m}$. Proof of (i): If $m \leq s$, then the subspace $L_{S}$ defined in Sec. 2 is $E_{m}$. Hence $x_{1}=\theta$ and $f\left(x_{1}\right)=0$, the minimum of $f(x)$ in $E_{m}$. Proof of (ii): That

$$
\frac{f\left(x_{2 k+2}\right)}{f\left(x_{2 k}\right)} \leq c_{2}<1
$$

follows from the chain of inequalities (2.13). We have to prove the inequalities involving $c_{1}$.

Given $x_{0}$ with at least $s+1$ nonzero components. By theorem (5 . 1 ) all other vectors $x_{k}$ have at least $s+1$ nonzero components, so that no $x_{k}=\theta$. By theorem (3.8), the normalized gradient vectors $y_{2 k}$ have as a limit set a continuum $R$. For each point $r$ in $R$, we have $T^{2} r=r$. Suppose a position vector $x$ were such that $r=A x /\|A x\| \in R$. That is, $x$ would be in the direction of $A^{-1} r$. Let $x^{\prime \prime}$ be the result of two steps of the optimum s-gradient method applied to $x$. Since $T^{2} r=r$, we see that $x^{\prime \prime}$ would be in the same direction as $x$. Hence

$$
\begin{equation*}
x^{\prime \prime}=y x \text { and so } f\left(x^{\prime \prime}\right)=y^{2} f(x) \tag{3.19}
\end{equation*}
$$

for some Y with $0<\gamma=\gamma(r)<1$.
I.e., for each point $r$ of $R$ there is a positive real number $\gamma(r)$ such that whenever the gradient of a vector x lies in the direction of $r$, then (3.19) holds.

Let $C$ be the minimum of $y(r)$ for $r \in R$. Since $R$ is compact, the minimum is assumed and $C>0$. Hence

$$
\begin{equation*}
0<c^{2} \leq \frac{f\left(x^{\prime \prime}\right)}{f(x)} \tag{3.20}
\end{equation*}
$$

for all $x$ such that $A x /\|A x\| \in R$.
Now the ratio $f\left(x^{\prime \prime}\right) / f(x)$ is a continuous function of $x$. Let $N(R) \subset \Sigma$ be such a neighborhood of $R$ that

$$
\begin{equation*}
\frac{1}{2} C^{2} \leq \frac{f\left(x^{\prime \prime}\right.}{f x f} \tag{3.21}
\end{equation*}
$$

for all $x$ with $A x /\|A x\|$ in $N(R)$. Consider the sequence $\left\{x_{2 k}\right\}$. Let $z_{2 k}=A x_{2 k}$, and let $y_{2 k}=z_{2 k} /\left\|_{2 k}\right\|$. By theorem (3.8), the $\left\{\mathrm{y}_{2 k}\right\}$ have R as a limit set. Hence there is a $K$ such that for $k \geq k$, all $y_{2 k}$ lie in $N(R)$. By (3.21) then for $k \geq k$ we have

$$
\frac{1}{2} C^{2} \leq \frac{f\left(x_{2 k+2}\right)}{f\left(x_{2 k}\right)}
$$

Letting $c_{1}=\frac{1}{2} c^{2}$ completes the proof of the theorem.
Actually we could have taken $c_{1}=c^{2}-\epsilon$, for any $\epsilon>0$.
(3.22) Corollary. With the hypotheses of theorem (3.17), there exist constants $\mathrm{d}_{1}, \mathrm{~d}_{2}$ with

$$
0<d_{1} \leq \frac{f\left(x_{k+1}\right)}{f\left(x_{k}\right)} \leq d_{2}<1, \quad \text { for all } k
$$

Proof. The corollary follows from theorem (3.17), the inequalities (2.13), and the fact that $f\left(x_{k}\right) \searrow 0$, as $k \wedge \infty$.
(3.23) Theorem . Fix $s \gg_{-}$. Let $x_{0}$ be any vector such that $x_{2}$ -is parallel to $x_{0}$ in the optimum s-gradient method. In other words, $z_{0} /\left\|z_{0}\right\|$ is in the set $F(A)$ of (4.5), where $z_{0}=A x_{0}$. Then (3.24) $\frac{f\left(x_{k+1}\right)}{f\left(x_{k}\right)}=c^{2} \quad(k=0,1,2, .$.$) ,$
where $c^{2}$ depends on $A$ and on $x_{0}$.
Remark. The import of this theorem is that, although the $x_{k}$
alternate between two fixed directions, as $k \rightarrow \infty$, the ratio (3.24) is constant for all $k$, and does not alternate.

Proof of (3.23). We first note from Corollary (2.23) that the theorem is true for $s=1$, and that (2.22) gives a formula for $c^{2}$ in terms of the two nonzero components $\zeta_{1}$, $\zeta_{2}$ of $z_{0}$.

For any fixed $s>1$, let $\pi$ be the 2 -space spanned by $x_{0}$ and $x_{1}$. Let $f_{\pi}(x)$ be the restriction of $f(x)=\frac{1}{2} x^{T} A x$ to the subspace $\pi$. Then the vectors $x_{0}, x_{1}, x_{2}, \ldots$ can be shown by a geometrical argument to be the successive iterates of the optimum l-gradient method for finding the minimum of $f,(x)$ in $\pi$, starting with $x_{0}$. Then (3.24) for $s=1$ states that

$$
\frac{f_{\pi}\left(x_{k+1}\right)}{f_{\pi}\left(x_{k}\right)}=c^{2}
$$

for some constant $c^{2}$ depending on the eigenvalues of $f_{\pi}$, Since $f_{\pi}(x)=f(x)$ in $\pi$, this proves the theorem for $s$.

Presumably theorem (3.23) could somehow be proved from theorem (2.18), just as the case $s=1$ follows from (2.22). Corollary (3.22) could also be proved from theorem (3.23).

## 4. Nature of the Asvmptotic Directions.

We should like to characterize as well as we can the possible limiting vectors $r \in R$ of the (normalized) gradient vectors $y_{2 k}$ of theorem (3.8). Since $T^{2} r=r$, for $r$ in $R$, we have

$$
\begin{align*}
c r & =P_{s}(A ; \operatorname{Tr}) P_{S}(A ; r)  \tag{4.1}\\
& =Q_{2 s}(A) r
\end{align*}
$$

where $c>0$ is a constant and $Q_{r},(t)$ is the product of the two polynomials $P_{s}(t ; T r)$ and $P_{s}(t ; r)$ Letting $r=\left(\rho_{1}, \cdots . \rho_{n}\right)^{T}$, we have

$$
\begin{equation*}
\mathrm{cp}_{1}=Q_{2 s}\left(\lambda_{i}\right) \rho_{i} \quad(i=1, \ldots n) \tag{4.2}
\end{equation*}
$$

Recall from p. 44 of [12] that $P_{S}(t ; T r)=t^{s}+.$. and $P_{S}(t ; r)$ $=t^{s}+\ldots$ are polynomials of degree $s$, each with $s$ distinct real zeros in the open interval $\left(\lambda_{1}, \lambda_{n}\right)$. Hence $Q_{2 s}(t)=t^{2 s}+.$. is a polynomial of degree 2 s with 2 s real zeros in the interval $\left(\lambda_{1}, \lambda_{\mathrm{n}}\right)$, counting double zeros twice, if any. Now $c>0$ in (4.2), which implies that for each i

$$
(4.3)
$$

$$
\begin{equation*}
Q_{2 s}(\text { hi })=c>0 \text { or } \rho_{i}=0 \tag{4.3}
\end{equation*}
$$

$$
(i=1, \ldots \ldots n)
$$

Since $Q_{1}$, ( $t$ ) vanishes for some $t$ in $\left(\lambda_{1}, \lambda_{n}\right)$, the equation $Q_{2 s}(t)$ = c > 0 can have 2'3'4' $\boldsymbol{m}^{\prime}$ or 2 s distinct real roots, which we call $\mu_{䒑_{j}}(j=1$, . . . m), and number so that

$$
\mu_{1}<\mu_{2}<\ldots<\mu_{m}
$$

(Here we count a multiple root of $Q,(t)=c$ only once.) Thus

$$
Q_{2 s}\left(\mu_{j}\right)=c \quad(j=1, \ldots m)
$$

By (4.3) each $\lambda_{i}$ for which $\rho_{i} \neq 0$ is one of the $\mu_{\mathcal{J}}$.
(4.4) Definition. Given $x_{0}$. Let $R$ be the set of limiting points of the normalized gradients $\left\{y_{2 k}: k=0^{\prime} 1^{\prime} \ldots\right\}$ of the optimum s-gradient method starting from $x_{0}$. For any vector $r=\left(\rho_{1}, ~ . ~ . ~ . ~ \rho_{n}\right)^{T}$ in $R$, let $S$ be the set of $\lambda_{i}$ for which $\rho_{i} \neq 0$. Any such set is called an asymptotic spectrum of the optimum s-gradient method for the given $x_{0}$. Any $r$ in $R$ is called an asymptotic gradient vector of the same iteration.

Note that $R$ depends on $A$ and $x_{0}$, and we occasionally write $R\left(x_{0}, A\right)$ to make the dependence explicit. Note that $S$ is a property of $r$ only, and only indirectly of $x_{0}$.
-(4.5) Definition. For a given $A$, we define the invariant set $F(A)$ of the optimum s-gradient method to be the set of unit vectors $r$ such that $T^{2} r=r$.

We have shown in theorem (3.8) that, for any $x_{0}, R\left(x_{0}, A\right) \subset F(A)$. It is never true that $R\left(x_{0}, A\right)=F(A)$. However, it is true that

$$
F(A)=\bigcup_{x_{0} \in E_{n}} R\left(x_{0}, A\right) .
$$

For, if $r \in F(A)$, then $T^{2} r=r$, so that $r=R(r, A)$.
(4.6) Theorem. Given $x_{0}=\left(\xi_{l}^{(0)} \text {, . . . . } \xi_{n}^{(0)}\right)^{T}$ with $\xi_{i}^{(0)} \neq 0 \quad(i=1$, $\ldots, \mathrm{n})$. Assume $\mathrm{s}<\mathrm{n}$. Then both eigenvalues $\lambda_{1}$ and $\lambda_{\mathrm{n}}$ belong to all asymptotic spectra $S$ of the optimum s-gradient method starting with x

Proof. Assume that $\lambda_{q}(q<n)$ is the largest eigenvalue in the asymptotic spectrum $S$ corresponding to an asymptotic vector $r$ of $R\left(x_{0}, A\right)$. The zeros of each $P_{S}\left(t ; z_{k}\right)(k=0,1$, . . . ) lie in the open interval $\left(\lambda_{1}, \lambda_{n}\right)$. Hence $P_{S}\left(\lambda_{n} ; z_{k}\right) \neq 0$ for all $k$. Hence $\eta_{n}^{(2 k)} \neq 0$ for all $k$, where the $\eta_{i}^{(2 k)}$ are the components of $y_{2 k}=z_{2 k} /\left\|z_{2 k}\right\|$.

Let $\tau$ be the largest zero of $P_{S}(t ; T r) P_{S}(t ; r)$. Since the zeros of both $P_{S}(t ; T r)$ and $P_{S}(t ; r)$ lie in the open interval $\left(\lambda_{1}, \lambda_{q}\right)$, we see that $P_{S}(t ; T r) P_{S}(t ; r) \lambda$, as $t \lambda$, for $t>\tau$. Hence

$$
c^{2}=P_{s}\left(\lambda_{q} ; \operatorname{Tr}\right) P_{s}\left(\lambda_{q} ; r\right)<P_{s}\left(\lambda_{n} ; \operatorname{Tr}\right) P_{s}\left(\lambda_{n} ; r\right) .
$$

-But then, by continuity,

$$
P_{s}\left(\lambda_{q} ; z_{2 k+1}\right) P_{s}\left(\lambda_{q} ; z_{2 k}\right) \leq \sigma P_{s}\left(\lambda_{n} \cdot z_{2 k+1}\right) P_{s}\left(\lambda_{n} ; z_{2 k}\right)
$$

for some $\sigma<1$ and all $k>k$. Since all $\prod_{n}^{(r)} \neq 0$, and since $\eta_{q}^{\left(2 k_{j}\right)} \rightarrow \rho_{q} \neq 0$, for a certain subsequence $k_{j^{\prime}}$ this means that $\left.\mid \eta_{n}^{\left(2 k_{j}\right.}\right) \mid \rightarrow \infty$, as $j \rightarrow \infty$. This is impossible, since all $y^{(2 k)}$ lie
on the unit sphere．Hence $q=n$ ，and $\lambda_{n}$ is in the asymptotic spectrum S．

The proof that $\lambda_{1}$ is in $S$ is analogous．
（4．7）Theorem．Given $x_{0}$ with $\xi_{\dot{i}}^{(0)} \neq 0 \quad(i=1, \ldots, n)$ ；assume that $s>n$ ．Let $m$ be the number of eigenvalues in any asymptotic spectrum $S$ of the optimum s－gradient method．Then

$$
s+1 \leq m \leq 2 s
$$

Proof．Let $r \in R$ be an asymptotic gradient vector corresponding to a given S ．As shown after（4．3），the asymptotic spectrum $S$ is a subset of the set of $t$ for which $P_{S}(t ; T r) P_{S}(t ; r)=c$ ，and the number of such $t$ is between 2 and 2 s ．

However，if $m \leq s$ ，one step of the optimum gradient method would carry $r$ into $\theta$ ，and so $r$ could not belong to $R$ ．Hence $s+1 \leq m \leq 2 s$.
（4．8）Theorem．Suppose $s<n$. Let $x_{0}=\left(\xi_{1}^{(0)} \text { ，．．．．} \xi_{n}^{(0)}\right)^{T}$ be any －vector in $E_{n} \underline{\text { with exactly }} \underline{+1 \text { nonzero components }} \xi_{i}^{(0)}$－The⿵⺆⿻二丨冂刂 $x_{0}, x_{2}, x_{4}$ ，are all collinear vectors．That is，the normalized gradient vector $y_{0}=A x_{0} /\left\|A x_{0}\right\|$ is in the invariant set $F(A)$ of（4．5）．

Proof．Let $z_{0}=A x_{0}$ ．It will suffice to prove that $z_{2}=c_{0} z_{0}$ ， for some positive constant $c_{0}$ ．Without loss of generality we may assume that $n=s+1$ ，since the components for which $\xi_{i}^{(0)}=0$ remain zero．

By (2.2)

$$
\begin{equation*}
z_{1}=z_{0}^{+} \gamma_{1} A z_{0}^{+} \cdot \cdots+\gamma_{S} A_{Z_{0}}^{S} \tag{4.9}
\end{equation*}
$$

and $\gamma_{7}$, - on $\gamma_{s}$ are so chosen that $z_{I_{+}}$is orthogonal to $z_{0}, A z z_{0}$, . $A^{s-1} z_{0}$. Because $s+1$ components of $z_{0}$ are nonzero, the $s$ vectors $z_{O}, A z_{O}, \ldots . . A^{s-l_{z}} z_{0}$ are linearly independent. Hence the set $\left\{z_{0}, A z_{0}\right.$, $\left.\ldots, A^{s-1} z_{0}\right\}$ forms a basis for the subspace of $E_{S+1}$ orthogonal to $Z_{I}$. Next, $z_{2}$ is formed as a linear combination of $z_{1}, A z_{1}, \ldots$, $A^{S} z_{1}$ which is orthogonal to $z_{1}, A z_{1}, \ldots, A^{S-1} z_{1}$. Since $z_{2}$ is orthogonal to $z_{1}$, it is expressible in terms of the basis $z_{0}$, • $A^{s-1} z_{0}$ :

$$
\begin{equation*}
z_{2}={ }^{c_{0} z_{0}}+c_{1} A z_{0}+\ldots+c_{s-1} A_{z_{0}-1} \tag{4.10}
\end{equation*}
$$

We shall prove that $c_{1}=c_{2}=. .=c_{s 1}=0$.
Take the inner product of (4.10) with $\mathrm{Az}_{1}$ :

$$
\begin{align*}
& +c_{S-1} Z_{1}{ }^{T} A_{Z_{0}} . \tag{4.11}
\end{align*}
$$

But $z_{1} \mathrm{~T}_{\mathrm{Az}}^{2}{ }_{2}=\mathrm{z}_{2}^{\mathrm{T}} \mathrm{Az}_{1}=0$ because $\mathrm{z}_{2}$ is orthogonal to $\mathrm{Az}{ }_{1}$. And $z_{1}^{T} A z_{0}=z_{1} A^{T} z_{0}=\ldots .=z_{1}^{T} A^{s-1} z_{O}=0^{\prime}$ because $z_{1}$ is orthogonal to $A z_{0}, A^{2} z_{0}, \ldots A^{s-1} z_{0}$. And $z_{1} T^{T} S_{0} \neq 0$, since otherwise by (4.11)
$z_{l}$ would be $\theta$. It then follows from (4.11) that $c_{s l}=0$.

Next, taking the inner product of (4.10) with $A^{2} Z_{1}$ and using the same argument and the fact that $c_{s_{1}}=0$, we show that $c_{s 2}=0$. After taking the inner product of (4.10) with $A z_{1}, A^{2} z_{1}, \cdots, A^{s-1} Z_{1}$, we will have proved that $c_{s-1}=. . c_{2}=c_{1}=0$. Then, from (4.10), $z_{2} . c_{0} z_{0}$. That $c_{0}>0$ follows from the proof of (3.3). This completes the proof of theorem (4.8).

Theorem (4.8) implies that any $s+1$ eigenvalues of $A$ can be in the asymptotic spectrum for some start $x_{0}$. Moreover, any vector $r$ with exactly $s+1$ nonzero components can be an asymptotic gradient vector of an iteration. This extends to $s \geq 2$ the known fact for the ordinary optimum l-gradient method in 2 dimensions that any initial gradient direction is repeated at every other step of the iteration. See the end of Sec. 2 above, or p. 214 of Ostrowski [9].

That for all s the period of the iteration in theorems (3.8) and (4.8) is 2, and not higher than 2, was a surprising fact to the author. However, the experiments of Khabaza [8] for $s=3$ suggest the period 2.

For $s=1$ we have $s+1=2 s=2$, and then by theorem (4.7) all the vectors invariant under two steps of the optimum l-gradient method are of the type covered in theorem (4.8). From this we can now show for $s=1$ that the limiting set $R$ of theorem (3.8) is actually a single point. The following is a modification of Akaike's proof in [1] of the Forsythe-Motzkin conjecture [7].
(4.12) Theorem (Akaike). Let $s=1$. Let $y_{0}=\left(\eta_{1}^{(0)}, \ldots, \eta_{n}^{(0)}\right)^{T}$ be any vector in $\Sigma^{*}$ with $\eta_{i}^{(0)} \neq 0 \quad(i=1, \ldots, n)$. Then the sequence $\left\{y_{2 k}: k=0,1, \ldots\right\}$ of normalized gradients converges to a single point $r$ whose spectrum is $\left\{\lambda_{l}, \lambda_{n}\right\}$. Moreover, $T^{2} r=r$.

Proof. By theorem (3.8) the set of unit vectors $\left\{y_{2 k}: k=0,1, \ldots\right\}$ has a continuum $R$ as a limit set. By theorem (4.7), for any $r \in R$ the corresponding spectrum $S$ of $r$ has only 2 eigenvalues in it (for $s+1$ $=2 \mathrm{~s}=2$ ). Now by theorem (4.6) the two eigenvalues in $S$ must be $\lambda_{1}$ and $\lambda_{n}$. Let $r$ be any point of $R$; let $r=\left(\rho_{1}, 0, \ldots, \rho_{n}\right)^{T}$, with $\rho_{1}^{2}+\rho_{n}^{2}=1$. Then $P_{1}(t ; r)=t-\mu$, where $\mu=\lambda_{1} \rho_{1}^{2}+\lambda_{n} \rho_{n}^{2}$. Hence $P_{1}(A ; r) r=\left(\left(\lambda_{1}-\mu\right) \rho_{1}, 0, \ldots, 0,\left(\lambda_{n}-\mu\right) \rho_{n}\right)^{T}$.

By the proof of theorem (3.8),

$$
\begin{array}{rlr}
L & =\lim _{k \rightarrow \infty} \varphi\left(w_{k}\right)=\varphi(r)=\left\|P_{1}(A ; r) r\right\|^{2}, \quad \text { since }\left\|\|^{2}=1\right. \\
& =\left(\lambda_{1}-\mu\right)^{2} \rho_{1}^{2}+\left(\lambda_{n}-\mu\right)^{2} \rho_{n}^{2}, &
\end{array}
$$

or

$$
\begin{equation*}
L=\left(\lambda_{\mathrm{n}}-\lambda_{\mathrm{l}}\right)^{2} \rho_{\mathrm{l}}^{2} \rho_{\mathrm{n}}^{2} . \tag{4.13}
\end{equation*}
$$

Now $L$ is a number determined by the iteration, $\lambda_{1}$ and $\lambda_{n}$ are given eigenvalues, and $\rho_{1}^{2}+\rho_{n}^{2}=1$. Hence the pair ${ }^{2} \rho_{1}, \rho_{n}^{2}$ are determined by (4.13), up to an interchange at most. Hence the set $R$ can have at most eight vectors in it, if all permutations of signs are considered. But then, since $R$ is a continuum, it must consist of a single point, which we call r . Then $\mathrm{y}_{2 \mathrm{k}} \rightarrow \mathrm{r}$, as $\mathrm{k} \rightarrow \infty$. This proves theorem (4.12).

$$
\text { Actually, if } \mathrm{r}=\left(\rho_{1}, \ldots, \rho_{\mathrm{n}}\right)^{\mathrm{T}} \text {, then } \operatorname{Tr}=\left(\rho_{\mathrm{n}}, \ldots . \rho_{1}\right)^{\mathrm{T}}
$$

where all components $\underline{p}^{2}=0$ for $l<i<n$. Then $r=\lim y_{2 k}$ and $\operatorname{Tr}=\lim \mathrm{y}_{2 k+1}$, as $k \rightarrow \infty$. So, the directions of the gradient vectors $z_{k}$ alternately approach the directions of $r$ and $\operatorname{Tr}$, as $k \rightarrow \infty$.

The reason we cannot extend our proof of theorem (4.12) to $\mathrm{s}>1$ is that the equation analogous to (4.13) involves between $s+1$ and 2 s unknown components of $r$, and we do not see how to limit $r$ to a finite number of vectors. Even for $s=2$, theorem (4.8) shows that all vectors $r$ with 3 nonzero components are invariant under $T^{2}$. Prescribing the vector $r$ to have unit length and prescribing the value of $L$ reduce the number of free parameters in $r$ to 1 . But, so far as the author can see, there remain $\infty^{l}$ possible limiting vectors $r$ in $R$.

Moreover, for an even number $s>1$, there are asymptotic spectra containing more than $s+1$ eigenvalues, as will now be demonstrated. We shall consider only spectra with symmetry about a midpoint. We do not know whether there are asymptotic spectra with more than $s+1$ eigenvalues without such a symmetry.

We shall first examine possible asymptotic spectra with an even number 2 q of eigenvalues. Let the eigenvalues in $S$ be $a-\mu_{q}$, $a-\mu_{q-1}, \quad a-\mu_{1}, a+\mu_{1}, 1 . \cdot . a+\mu_{q-1}, a+\mu_{q}$, where $0<a-\mu_{q}$ and $0<\mu_{1}<. . \quad .<\mu_{\mathrm{q}}$. Let us consider unit vectors $r$ with symmetric components $\rho_{q} . ., \rho_{1}, \rho_{1}, \ldots, \rho_{q}, \quad$ corresponding to the respective -points of the spectrum.

Because of the symmetry about the point $t=a$, the orthogonal polynomials $P_{2 k}(t ; r), P_{2 k+1}(t ; r)$ associated with $S$ and the $\left\{\rho_{i}^{2}\right\}$ satisfy the conditions

$$
\begin{equation*}
P_{2 k}(t ; r)=g_{k}\left((t-a)^{2}\right) \tag{4.14}
\end{equation*}
$$

where $\mathrm{g}_{\mathrm{k}}$ is a monic polynomial of degree k ;

$$
\begin{equation*}
P_{2 k+1}(t ; r)=(t-a) h_{k}\left((t-a)^{2}\right) \tag{4.15}
\end{equation*}
$$

where $h_{k}$ is a monic polynomial of degree $k$.
By symmetry, the even and odd polynomials $P_{k}(t ; r)$ are automatically orthogonal. By (4.14) orthogonality of the $\mathrm{P}_{2 \mathrm{k}}(\mathrm{t} ; \mathrm{r})$ among themselves can be expressed in the form

$$
\begin{equation*}
\sum_{i=1}^{q} g_{j}\left(\mu_{i}^{2}\right) g_{k}\left(\mu_{i}^{2}\right) \rho_{i}^{2}=0 \quad\left(j, k=0^{\prime} 1^{\prime} \ldots, \operatorname{l}\right. \tag{4.16}
\end{equation*}
$$

Thus the $g$, (t) are themselves orthogonal polynomials over the set $\mu_{1}^{2}$, $\mu_{q}^{2}$ with the weight factors $\rho_{l}^{2}$, $\quad \rho_{q}^{2}$. Moreover, $\hat{g}_{\mathrm{k}}(\mathrm{t})=$ $(-1)^{k} g_{k}\left(a^{2}-t\right)$ are monic orthogonal polynomials over the transformed set $\hat{S}=\left\{a^{2}-\mu_{q}^{2}, \ldots, a^{2}-\mu_{1}^{2}\right\}$ with the same weights ${ }^{2} \rho_{1}, \ldots, \rho_{n}^{2}$ 。 Note that $\left|\hat{g}_{k}(0)\right|=\left|P_{2 k}(0 ; r)\right|$ and that $\left|\hat{g}_{k}\left(a^{2}-\mu_{i}^{2}\right)\right|=\left|P_{2 k}\left(a \pm \mu_{i} ; r\right)\right|$ for $i=1, \ldots, q$. Hence $\left|\hat{g}_{k}(t) / \hat{g}_{k}(0)\right|$ has the same constant value over the set $\hat{S}$ that $\left|P_{2 k}(t ; r) / P_{2 k}(0 ; r)\right|$ has over the set $S$.

By (4.15) the orthogonality of the $P_{2 k+1}$ among themselves can be -expressed as

$$
\begin{equation*}
\sum_{i=1}^{q} h_{j}\left(\mu_{i}^{2}\right) h_{k}\left(\mu_{i}^{2}\right) \mu_{i}^{2} \rho_{i}^{2}=0 \quad(j, k=0,1, \ldots j \neq k) \tag{4.17}
\end{equation*}
$$

Thus the $\hat{h}_{k}(t)=(-1)^{k} h_{k}\left(a^{2}-t\right)$ are monic orthogonal polynomials over the set $\hat{S}=\left\{a^{2}-\mu_{\underline{q}}^{2}, \ldots, a^{2}-\mu_{1}^{2}\right\}$ with the different weights $\mu_{1}^{2} \rho_{1}^{2}, \ldots$. $\mu_{q}^{2} \rho_{q}^{2}$. Note that $\left|\hat{h}_{k}(0)\right|=\left|h_{k}\left(a^{2}\right) \|=P_{2 k+1}(0 ; r)\right| / a$, and that

$$
\left|\hat{h}_{k}\left(a^{2}-\mu_{i}^{2}\right)\right|=\left|h_{k}\left(\mu_{i}^{2}\right)\right|=\left|P_{2 k+1}\left(a \pm \mu_{i} ; r\right)\right| / \mu_{i}
$$

Thus constancy of $\left|P_{2 k+1}(t ; r)\right|$ over $S$ does not imply constancy of $\left|\hat{h}_{k}(t)\right|$ over 5 . The even and odd polynomials transform differently. By means of these orthogonal polynomials $\hat{\mathrm{g}}_{\mathrm{k}}$ we can reduce the problem of the invariance of the $r$ under two steps of the optimum 2s-gradient method over $S$ to the problem of the invariance of an optimum s-gradient method over $\hat{S}$ in a space of half the dimension. To be precise, the above relations imply the following result, which we do not prove.
(4.18) Theorem. If $s$ is even and $s+1<2 q \leq 2 s$, then the vector $r=\left(\rho_{q}, \rho_{1}, \rho_{1}, \cdot \rho_{q}\right)^{T}$ (with no $\rho_{i}=0$ ) is in $\eta_{m}$ the invariant set (4.5) for the optimum s-gradient method for the diagonal matrix of 2q nonzero elements

$$
\operatorname{diag}\left(a-\mu_{q}, \quad a-\mu_{1}, a+\mu_{1}, \ldots, a+\mu_{q}\right)
$$

if and only in the vector $\hat{r}=\left(\rho_{1}, \ldots . \rho_{q}\right)^{T} \quad$ (with no $\rho_{i}=0$ ) is in the invariant set for the optimum ( $s / 2$ )-gradient method for the diagonal matrix of $q$ nonzero elements

$$
\operatorname{diag}\left(a^{2}+\mu_{1}^{2}, \ldots, a^{2}+\mu_{q}^{2}\right)
$$

Moreover, when iterations exist with these invariance properties, if $z_{0}=r /\|r\|$ and $\hat{z}_{o}=\hat{r} /\|\hat{r}\|$, then $\left\|z_{k}\right\|=\left\|\hat{z}_{k}\right\|$ for $k=0^{\prime} 1^{\prime} 2$, . . . . where $z_{k}$ and $\hat{z}_{k}$ are the gradient vectors of the respective iterations.

We do not know a comparable theorem for odd integers s.
As an application of theorem (4.18), we can show that for any $s$ of the form $s=2^{p}(p=0,1,2, \ldots)$, there exist vectors with $2 s$ nonzero components that are in the invariant set of some optimum s-gradient method. For $p=0$ this is theorem (4.12), and is true for any diagonal matrix of two positive elements $\operatorname{diag}\left(a^{2}+\mu_{1}^{2}, a^{2}+\mu_{2}^{2}\right)$ and any vector $r=\left(\rho_{1}, \rho_{2}\right)^{T}$, Application of the first sentence of (4.18) leads to $s=2$ with any matrix of form $\operatorname{diag}\left(b^{2}+v_{1}^{2}, b^{2}+v_{2}^{2}, b^{2}+v_{3}^{2}\right.$, $b_{2}+\nu_{4}^{2}$ ) where $b^{2}+\nu_{1}^{2}=a-\mu_{2}, b^{2}+\nu_{2}^{2}=a-\mu_{1 \prime} \quad b^{2}+\nu_{3}^{2}=a+\mu_{1}$, $b^{2}+\nu_{4}^{2}=a+\mu_{2}$, and corresponding vector $c\left(\rho_{2}, \rho_{1}, \rho_{1}, \rho_{2}\right)^{T}$. Another application of ( 4.18 ) leads to $s=4$ with the matrix

$$
\operatorname{diag}\left(b-\nu_{4}, \cdot \cdot \cdot \cdot b-\nu_{1}, b+\nu_{1}, \cdot \cdot . b+\nu_{4}\right)
$$

and corresponding vector $c^{\prime}\left(\rho_{2}, \rho_{1}, \rho_{1}, \rho_{2}, \rho_{2}, \rho_{1}, \rho_{1}, \rho_{2}^{\top}\right)$ 。 It is clear that the process may be continued to $s=2^{p}$ for any $p$.

Note from theorem (4.7) that 2 s is the maximal number of nonzero components in any vector in the invariant set for an optimum s-gradient method. Our above example illustrates the maximal case.

We next consider 'symmetric asymptotic spectra with an odd number $2 q+1$ of eigenvalues $a-\mu_{q}, \ldots, a-\mu_{1}, a, a+\mu_{1}, \ldots, a+\mu_{q}$ and a corresponding symmetric vector

$$
\left(\rho_{q}, \bullet \bullet_{v} \rho_{l}, \rho_{o}, \rho_{1}, \ldots \circ, \rho_{q}\right)^{T}
$$

invariant under $\mathrm{T}^{2}$. Then again the orthogonal polynomials take the
forms (4.14) and (4.15). The odd polynomials are still defined by the condition (4.17), but the condition (4.16) must be replaced by

$$
\begin{gather*}
2 \sum_{i=1}^{n} g_{j}\left(\mu_{i}^{2}\right) g_{k}\left(\mu_{i}^{2}\right) \rho_{i}^{2}+g_{j}(0) g_{k}(0) \rho_{o}^{2}=0  \tag{4.19}\\
(j, k=0,1, \cdots \cdots \neq k)
\end{gather*}
$$

The analog of theorem (4.18) is now stated, but not proved:
(4.20) Theorem. If $s$ is even and $s+1 \leq 2 q+1<2 s$, then the vector $r=\left(\rho_{q}, \ldots, \rho_{1}, \rho_{o}, \rho_{1}, \quad, \rho_{q}\right)^{T}$ (with no $\rho_{i}=0$ ) is in the invariant set (4.5) for the optimum s-gradient method for the diagonal matrix of $2 q+1$ nonzero elements

$$
\operatorname{diag}\left(a-\mu_{q}, \ldots, a-\mu_{1}, a, a+\mu_{1}, \ldots, a+\mu_{q}\right)
$$

if and only if the vector $\hat{r}=\left(\rho_{0} / \sqrt{2}, \rho_{1} \text {, e Am } \rho_{q}\right)^{T} \quad\left(\right.$ with no $\left.\rho_{i}=0\right)$ is in the invariant set for the optimum (s/2) -gradient method for the diagonal matrix of $q+1$ elements

$$
\operatorname{diag}\left(a^{2}, a^{2} \mu_{1}^{2}, \cdots a^{2}+\mu_{q}^{2}\right)
$$

Moreover, when iterations exist with these invariance properties, if $z_{0}=r /\left\|r_{0}\right\|$ and $\hat{z}_{0}=\hat{r} /\|\hat{r}\|$, then $\left\|z_{k}\right\|=\left\|\hat{z}_{k}\right\|$ for $k=0,1,2, . \quad . \quad$. where $z_{k}$ anal $\hat{z}_{k}$ are tine gradient vectors of the respective iterations.

If $s$ is odd, then the set of $2 q+1$ eigenvalues $\left(a-\mu_{q}, \cdots\right.$, $a-\mu_{1}, a, a+\mu_{1}$, $\left.\boxminus a+\mu_{q}\right\}$ can never be the asymptotic spectrum of an optimum s-gradient iteration.

The first two sentences are strict analogs of theorem (4.18). The third is true because $P_{2 k+1}(a ; r)=0$ for all $k$.

The signs of the $\rho_{i}$ are of no importance in theorems (4.18) and (4.20), and any $\rho_{i}$ could be left alone or replaced by $-\rho_{i}$ independently at any place it is mentioned.
5. Singular and Derogatory Quadratic Forms; Zero Components.

Two restrictions placed on $A$ above are really irrelevant--that A be regular and nonderogatory. If $A$ is singular, then for some $p>1$, we have $\lambda_{1}=\ldots=\lambda_{\mathrm{p}}=0<\lambda_{\mathrm{p}+1}<\ldots<\lambda_{\mathrm{n}}$. Then it follows from (2.12) that

$$
\xi_{i}^{(k+l)}=\xi_{i}^{(k)}, \quad \text { for } l \leq i \leq p ; k=0,1,2, \ldots,
$$

while all components $\xi_{i}^{(k)} \rightarrow 0$, as $k \rightarrow \infty$, for $p+1 \leq i \leq n$. On the other hand $f(x)=\frac{1}{2} x^{T} A x=\sum_{i=1}^{n} \lambda_{i} \xi_{i}^{2}=\sum_{i=p+1}^{n} \lambda_{i} \xi_{i}^{2}$. Thus $f(x)$ is minimized for all vectors in the subspace $N$ where $\xi_{1}=\ldots .=\xi_{p}=0$, and the gradient methods proceed from $x_{0}$ to the closest point $x_{\infty}$ of $N$, with $\operatorname{allx}_{k}-x_{\infty}$ and all gradients $z_{k}$ located in the orthogonal complement of N .

If A is derogatory, it has multiple eigenvalues but a complete set of eigenvectors (because $A$ is symmetric). Suppose, for example, that $0<\lambda_{1}=\lambda_{2}=\ldots .=\lambda_{r}<\lambda_{r+1}<\ldots<\lambda_{n}$, and suppose that

$$
x_{0}=\left(\xi_{1}^{(0)}, \ldots, \xi_{r}^{(0)}, \xi_{r+1}^{(0)}, \ldots, \xi_{\mathrm{n}}^{(0)}\right)^{\mathrm{T}}
$$

Now the orthogonal basis of eigenvectors belonging to $\lambda_{1}$, . . . . $\lambda_{r}$ is not uniquely defined. Our preceding analysis required at various places (egg., in the proof of (4.8) that the $\lambda_{i}$ be distinct for each nonzero component $\xi_{-1}^{(0)}$, but zero components $\xi_{i}^{(0)}$ were ignored. If any of $\xi_{2}^{(0)}$, $\xi_{r}^{(0)}$ are nonzero, make an orthogonal transformation of the eigenvector basis so that $x_{0}$ takes the form

$$
x_{0}=\left(\left(\xi_{1}^{(0) 2}+\ldots+\xi_{r}^{(0) 2}\right)^{\frac{1}{2}}, 0, \ldots, 0, \xi_{r+1}^{(\Omega)}, \ldots, \xi_{n}^{(0)}\right)^{T}
$$

Then drop the new zero components $\xi_{2}$. . . . $\xi_{r}$ entirely, and effectively reduce $A$ to a nonderogatory matrix $A$ of order $n-r+1$.

Thus, in effect, only the set and number of distinct nonzero eigenvalues of $A$ have a real relevance to the gradient methods for quadratic functions $\frac{1}{2} x^{T} A x$.

Moreover, zero components of any $\mathrm{x}_{\mathrm{k}}$ should be ignored, and the order of A reduced by unity for each zero component $\xi_{i}^{(k)}$ that occurs at any stage of the iteration.

If fewer than $s+1$ components of any $x_{k}$ are nonzero, then $x_{k+1}=\theta$ and the iteration terminates at once. Hence we have always insisted that at least $s+1$ components of $x_{0}$ be nonzero. Even so, one may ask, might not enough $P_{s}\left(\lambda_{i} ; z_{h}\right)$ be "accidentally" zero, so that for some later $\mathrm{x}_{\mathrm{k}}$ fewer than $\mathrm{s}+1$ components are nonzero? The answer is negative, as the following theorem shows:
(5.1) Theorem. Assume $s+1<n$. Assume $\xi_{i}^{(k)} \neq 0$ for $i=1$, . . . n. Then at least $s+1$ components $\xi_{i}^{(k+1)} \neq 0$.

Proof. By (2.12), $\xi_{i}^{(k+l)}=P_{s}\left(\lambda_{i} ; z_{k}\right) \xi_{i}^{(k)}$, up to a multiplicative constant that does not matter, where $P_{S}\left(t ; z_{k}\right)$ is the orthogonal polynomial of degree $s$ over the set $\left\{\lambda_{l}, \cdot \lambda_{n}\right\}$ with weights $\left[\zeta_{i}^{(k)}\right]^{2}$. We shall prove that there exist $s+1$ eigenvalues out of the $\lambda_{i}$ :

$$
\begin{equation*}
\lambda_{i}^{\prime}<\lambda_{2}^{\prime}<\ldots<\lambda_{s+1}^{\prime}, \tag{5.2}
\end{equation*}
$$

such that $P_{s}\left(\lambda_{i}^{\prime} ; z_{k}\right) P_{s}\left(\lambda_{i+1}^{\prime} ; z_{k}\right)<0$ for $i=1, \ldots, s$. A fortiori, $P_{s}\left(\lambda_{-1} ; z_{k}\right) \neq 0$ for $i=1,2, . . \quad . s+1$, and the theorem will have been proved.

If the above sign-alternation property is false, then let $q<s$ be the largest integer such that we can find $\left\{\lambda_{i}^{\prime}\right\}$ with

$$
\begin{equation*}
P_{S}\left(\lambda_{i}^{\prime} ; z_{k}\right) P_{S}\left(\lambda_{i+l}^{\prime} ; z_{k}\right)<0 \quad \text { for } i=1, \cdot . \cdot q-1 \tag{5.3}
\end{equation*}
$$

(Clearly some $q \geq 2$ exists, or else $P_{S}\left(\lambda_{i} ; z_{k}\right)$ would always be of one sign and hence $P_{s}$ could not be orthogonal to $P_{0}=1$. Then pick $\mu_{1}, \ldots, \mu_{q-1}$ with

$$
\lambda_{1}^{\prime}<\mu_{1}<\lambda_{2}^{\prime}<\mu_{2}<\ldots<\lambda_{q-1}^{\prime}<\mu_{q-1}<\lambda_{q}^{\prime}
$$

so that, if $Q(t)=\left(t-\mu_{1}\right) .0\left(t-\mu_{q-1}\right)$, then $P_{S}\left(\lambda_{i} ; z_{k}\right) Q\left(\lambda_{i}\right)>0$ for all $i=1$, . . . . n . (We omit details of the construction.) Then

$$
\left\langle P_{s}\left(t ; z_{k}\right), Q(t)\right\rangle=\sum_{i=1}^{n} P_{S}\left(\lambda_{i} ; z_{k}\right) Q\left(\lambda_{i}\right)\left[\zeta_{i}^{(k)}\right]^{2}>0
$$

-so that $\underset{S}{P}$ and $Q$ are not orthogonal. But, since $Q$ is of degree q-1 $\leq s-1, P_{s}$ must be orthogonal to $Q$. This contradiction completes the proof of theorem (5.1).

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