# SOME PROPERTIES OF THE METHOD 

 OF STEEPEST ASCENT;By<br>CHARLES HENRY JOHNSON<br>!<br>Bachelor of Science<br>Bradley University Peoria, Illinois 1949<br>Master of Science<br>Bradley University Peoria, Illinois<br>1950

Submitted to the Faculty of the Graduate School of the Oklahoma State University
in partial fulfillment of the requirements
for the degree of DOCTOR OF PHILOSOPHY

August, 1963

# SOME PROPERTIES OF THE METHOD OF STEEPEST ASCENT 

Thesis Approved:


Navel H. Muts
Albeit \& Torsion


Ancentuacluier
Dean of the Graduate School

## AGKNOWLEDGMENTS

I am indebted to Dr, Leroy Folks for suggesting the problem and for his valuable aid and direction.

I am also indebted to the National Science Foundation and the Ely Lilly Foundation for their financial support,

Especially, I am grateful ta Dr. Harrison S, Mendenhall for suggesting Oklahoma State University to me and to Dr, Carl E. Marshall for his interest and encouragement.

My thanks are extended to the members of my advisory committee, Professors John E. Hoffman, Robert D. Morrison, George W. Newell, Herbert Scholz Jr., and David L. Weeks.

I wish to acknowledge the fine support of my family throughout my studies and during the preparation of this thesis.

My thanks are extended to Mrs. Beverly Richardson for her careful typing of this thesis.

## TABLE OF CONTENTS

Chapter Page
I. INTRODUCTION ..... 1
II . ESTIMATION OF FIRST ORDER PARTIAL DERIVATIVES AND DESIGN SIZE ..... 4
Definition of Basic Terms ..... 4
Estimation of the Fifst Order Partials ..... 7
Responses at the Points of a: Square Grid ..... 14
Choice of Design Size ..... 18
III . PROPERTIES OF THE METHOD OF
STEEPEST ASCENT ..... 24
An Invariance Property of the
Steepest Ascent Path ..... 25
The Predicted Response ..... 44
Upper and Lower Bounds for $\beta_{11} / \beta_{22}$ ..... 48
Gradient Properties ..... 53
IV . THE CHOIĆE OF SCALE ..... 66
The Lack of Scale Invariance of the Steepest Ascent Path. ..... 67
Experimentation Along the Transformed Path of Steepest Ascent ..... 72
The Choice Between the Lines $\mathrm{L}_{\mathrm{d}}$ and $\mathrm{L}_{\mathrm{p}}$ ..... 78
V . SUMMARY ..... 85
Conclusions ..... 85
Areas for Future Research ..... 87
SELECTED BIBLIOGRAPHY ..... 89

## LIST OF TABLES

Table Page
I. Design Size Necessary to Minimize M.S.E, $\hat{\beta}_{1}$. ..... 23
II. Upper and Lower Bounds for $\beta_{11} / \beta_{22}$ ..... 52

## LIST OF FIGURES

Figure ..... Page

1. Relative Magnitude of Observed and Predicted Responses ..... 47
2. Intersection of the Gradient With the Axes of the ContourSystem , . . . . . . . . . . . . . . . . . . 60
3. Relationship Between $\mathrm{m}_{1}, \mathrm{~m}_{2}, \mathrm{~m}_{3}$, and the Contour System Axes . . . . . . . . . . . . . . . . 65
4. The Effect of a Change of Scale on the Steepest Ascent Path . . . . . . . . . . . . . . . . . . . . 73
5. An Appropriate Grid for Horizontal Elongated Contours , 80
6. The Slope of $L_{p}$ is Positive; There is no Basis for a Decision . . . . . . . . . . . . . .. . . . . 81
7. The Slope of $L_{p}$ is Positive; the Line $L_{p}$ Should be.$~ 82$
8. The Slope of $L_{p}$ is Negative; the Line $L_{t}$ Should be Chosen . . . . . . . . . . . . . . . . . . 83
9. The Slope of $L_{p}$ is Negative; the Line $L_{p}$ Should be Chosen . . . . . . . . . . . . . . . . . .84

## CHAPTER I

## INTRODUCTION

A fairly common objective in industry is to find that level of the ingredients and operating charcteristics which maximizes or minimizes some characteristic of the end product. For example, a lens maker might be interested in minimizing the coefficient of expansion of the finished lens. He might feel that the decisive factors were the concentration of a particular ingredient, the rate at which the molten glass was allowed to cool, and the length of time the glass was retained in a liquid state before pouring.

Frequently, it is not practical to operate at the particular factor Levels which optimizes the end product characteristic of interest. It may be that the cost or time involved is prohibitive. The practical limits on the factors involved define the factor space or the experimental region.

Let the controlled continuous variables, $x_{1}, x_{2}, \ldots, x_{p}$ be the independent factors and let $\eta$ be the true measure of the characteristic of the end product. $\eta$ is called the response and is dependent on the variables $x_{1}, x_{2}, \cdots, x_{p}$. Then the statistical model may be written as

$$
\begin{equation*}
y=\eta+e \tag{1.1}
\end{equation*}
$$

where $y$ is the observed response and $e$ is a measure of the failure of the mathematical model to conform to the real world. It will be assumed that $e \sim\left(0, \sigma^{2}\right)$.

In many situations where time and expense are prime considerations, it is imperative to adopt a strategy for approaching an optimal or stationary point in the experimental region using as few experiments as possible. One proposed strategy is the method of steepest ascent put forward as a statistical technique by G. E. P. Box and K. B. Wilson (1) in 1951. Further descriptions of this technique are given by O. L. Davies (2) and W. C. Cochran and G. M. Cox (3).

Box and Wilson (1) have shown that the maximum gain in response in proceeding a distance $r$ from a point 0 to a point $P$ in a $k$-demensional space is achieved approximately by varying the factors in proportion to their first order partial derivatives at 0. The direction thus determined is called the direction of steepest ascent.

In practice, the controlled variables are coded and the direction of steepest ascent estimated in terms of the coded variables. This is done by using the method of least squares to fit a plane to the response surface and using the coefficients of the fitted plane as estimates of the partial derivatives at 0 . It will be shown in Chapter II that these estimates are unbiased. The steepest ascent path is uncoded and further experiments are performed along this path until it is felt that a maximum has been attained. Another set of experiments

1 Note: ( ) refersto Selected Bibliography.
is run using this point as the design center and the foregoing procedure is repeated. The process is iterated until the experimenter feels that the process is being run in the proximity of a stationary point.

In chapter II some theprems regarding the estimation of the coefficients of the fitted plane will be presented. The content of a number of these theorems is no doubt known to experimenters but do not seem to have been formalized. Also an examination will be made of the relationship between variance, bias, mean square error, and spread of the points in the design. A design size is suggested so as to minimize the mean square error.

In chapter III a number of theorems regarding the path of steepest ascent will be presented. An interesting property of invariance of the steepest ascent line will be examined and a number of theorems involving the gradient line will also be given. Since the gradient to the factor space contours does not involve error, the study obviously suffers to the extent that the mathematical model fails to represent the experimental situation. In spite of this drawback, it is felt that some of the results will provide the experimenter with additional insight.

When the direction of steepest ascent in the coded variables is transformed to the uncoded variables, it is found that, in general, the resulting direction is no longer normal to the contours. This lack of invariance under a change of scale was noted by Box and Wilson (1) and more recently by O. Kempthorne (4). Chapter IV deals with the effect of the choice of units used in the coding and will examine in some detail the lack of invariance under a change of scale.

## CHAPTER II

## ESTIMATION OF FIRST ORDER PARTIAL DERIVATIVES AND DESIGN SIZE

Certain terms occur frequently in the ensuing discussion and it is necessary to establish their meaning at the outset.

## Definition of Basic Terms

Definition 2.1. The space in which the controlled variables are allowed to vary will be called the factor space.

A point in the factor space will be denoted by $x=\left(x_{1}, x_{2}, \ldots\right.$, $x_{p}$ ). It will be assumed throughout that the experimenter can attain the point x without error.

Definition 2.2. The true response at a point $x$ will be denoted by $\eta\left(x_{1}, x_{2}, \ldots, x_{p}\right)$ or more commonly by $\eta$.

The symbol $\eta$ will always be taken as a polynomial in the variables $x_{1}, x_{2}, \ldots, x_{p}$. For example, the true response might be assumed to be quadratic, in which case $\eta$ would have the form,

$$
\begin{equation*}
\eta=\beta_{0}+\beta_{1} x_{1}+\beta_{2} x_{2}+\beta_{11} x_{1}^{2}+\beta_{12} x_{1} x_{2}+\beta_{22} x_{2}^{2} . \tag{2.1}
\end{equation*}
$$

Definition 2.3. The observed response at the point $x$ will simply be called the "response" and will be denoted by $y$.

The response $y$ is a measurable quantity the magnitude of which depends on the particular value of $x$. For example, a bacteriologist might be concerned with the effect of temperature and humidity on the bacteria count in a culture In this case the factors are temperature and humidity while the response is the number of bacteria expressed in appropriate units.

Definition 2.4. The random error, denoted by $e_{\text {, }}$ is the difference between the true response and the observed response.

The random errors will be assumed to be independent of each other, to have a mean of 0 , and a variance of $\sigma^{2}$.

In conformance with definitions 2,2,2.3, and 2.4, y, $\eta$, and e are related by the equation

$$
\begin{equation*}
y=\eta+e . \tag{2.2}
\end{equation*}
$$

When it is necessary to distinguish between responses at different points or to fix attention at a particular point in the factor space, the subscript $i$ will be used. In this event, (2.2) is written as $y_{i}=\beta_{0}+\beta_{1} x_{l i}+\beta_{2} x_{2 i}+\beta_{11} x_{l i}{ }^{2}+\beta_{12} x_{1 i} x_{2 i}+\beta_{22^{x}}{ }_{2 i}{ }^{2}+e_{i} .(2.3)$

In matrix notation this may be written as

$$
y=\left(1, x_{1}, x_{2}, x_{1}^{2}, x_{1} x_{2}, x_{2}^{2}\right) \quad\left|\begin{array}{l}
\beta_{0} \\
\beta_{1} \\
\beta_{2} \\
\beta_{11} \\
\beta_{12} \\
\beta_{22}
\end{array}\right|+\quad \text { e }
$$

or more compactly by

$$
\begin{equation*}
y=X \beta+e . \tag{25}
\end{equation*}
$$

Many times it is desirable to consider concurrently a number of points in the factor space together with the associated responses and random errors. For this purpose the vector concept is well suited. N responses corresponding to N points in the factor space together with the $N$ random errors may be summarized by the matrix equation

$$
\left(\begin{array} { l } 
{ y _ { 1 } } \\
{ y _ { 2 } } \\
{ \cdot } \\
{ \cdot } \\
{ \cdot } \\
{ y _ { N } }
\end{array} \left|=\left|\begin{array}{llll}
1 & x_{11} & \cdots & x_{p l} \\
1 & x_{12} & \cdots & x_{p 2} \\
\cdot & \cdot & & \cdot \\
\cdot & \cdot & & \cdot \\
\cdot & \cdot & & \cdot \\
1 & x_{1 N} & \cdots & x_{p N}
\end{array}\right|\left(\begin{array}{l}
\beta_{0} \\
\beta_{1} \\
\cdot \\
\cdot \\
\cdot \\
\beta_{p N}
\end{array} \left\lvert\,+\left(\begin{array}{l}
e_{1} \\
e_{2} \\
\cdot \\
\cdot \\
\cdot \\
e_{N}
\end{array}\right)(2.6)\right.\right.\right.\right.
$$

or more compactly by

$$
\begin{equation*}
Y=X \beta+E . \tag{2.7}
\end{equation*}
$$

The matrix $X$ as displayed in (2.6) shows at a glance the configuration of the points in the factor space. When the dimension of the matrix $X$ is not apparent from the context it will be indicated.

Definition 2.5. The matrix. $X$ as displayed in (2.6) will be known as the design matrix for the assumed model.

## Estimation of the First Order Partials

The method of steepest ascent calls for experimentation in a direction which is determined by the estimates of the first order partial derivatives of $\eta$ evaluated at the origin. From a geometrical standpoint, this is equivalent to fitting a layperplane to the response surface over some small initial region in the factor space and proceeding in a direction determined by the coefficients of the fitted plane. This strategy has an intuitive appeal for the type of responses likely to be encountered in practice.

In the case of quadratic type responses, the type about which this thesis is chiefly concerned, the design matrix may be chosen so that the first order partials are estimated readily without bias.

Definition 2.6. The response predicted at a point $x$ by means of a hyperplane fitted from a design by the method of least squares will be denoted by $\quad y$.

The response predicted by means of a plane at a point $x=\left(x_{1}, x_{2}\right)$ is

$$
\begin{equation*}
\hat{y}=\hat{\beta}_{0}+\hat{\beta}_{1} x_{1}+\hat{\beta}_{2} x_{2} \tag{2.8}
\end{equation*}
$$

The question of what type of design to use in order to estimate the partial derivatives is an important one. It is especially desirable to have estimates which are unbiased and the ensuing discussion sheds
light upon why factorial designs are a common choice for this purpose.
The $X$ and $\beta$ matrices of (2.7) may be partitioned into their constant, linear, and quadratic parts. $Y$ may be written

$$
\begin{equation*}
Y=X_{0} B_{0}+X_{1} B_{1}+X_{2} B_{2}+E \tag{2.9}
\end{equation*}
$$

where

$$
\begin{aligned}
& B_{0}=1, \quad B_{1}=\left|\begin{array}{l}
\beta_{1} \\
\beta_{2}
\end{array}\right|, \quad B_{2}=\left|\begin{array}{c}
\beta_{11} \\
\beta_{12} \\
\beta_{22}
\end{array}\right|,
\end{aligned}
$$

In matrix notation, the fitted plane is

$$
\begin{equation*}
\hat{y}=\hat{B}_{0}+\hat{X}_{1} \hat{B}_{1} \tag{2.10}
\end{equation*}
$$

where

$$
\begin{equation*}
\hat{B}_{1}=\left(X_{1}^{\prime} X_{1}\right)^{-1} X_{1}^{\prime} Y \tag{2.11}
\end{equation*}
$$

Now

$$
\begin{equation*}
E\left(\hat{B}_{1}\right)=\left(X_{1}^{\prime} X_{1}\right)^{-1} X_{1}^{\prime}\left(X_{0} B_{0}+X_{1} B_{1}+X_{2} B_{2}+E\right) \tag{2.12}
\end{equation*}
$$

or

$$
\begin{equation*}
E\left(\stackrel{\rightharpoonup}{B}_{1}\right)=\left(X_{1}^{\prime} X_{1}\right)^{-1} X_{1}^{\prime} X_{0} B_{0}+B_{1}+\left(X_{1}^{\prime} X_{1}\right)^{-1} X_{1}^{\prime} X_{2} B_{2} \tag{2.13}
\end{equation*}
$$

Hence if $X_{1} X_{0}=\phi$ and $X_{1} X_{2}=\phi$, then $E\left(\hat{B}_{1}\right)=B_{1}$ and $\hat{B}_{1}$
will be an unbiased estimate of $\mathrm{B}_{1}$. Factorial designs enjoy the property that $X_{1}{ }^{\prime} X_{0}=\phi \quad$ and $X_{1}{ }^{\prime} X_{2}=\phi$ if the variables are properly coded. Also it is well known from least squares the ory that the variance-covariance matrix of $\hat{\mathrm{B}}_{1}$ is the same whether the quadratic terms are present or not.

Definition 2.7. A factorial design whose points form a 2-dimensional factor space will be called a rectangular grid, a square grid, or a unit grid according as the factor points form a rectangle, a square, or a unit square. If the length of the side of the square grid is 2 h , the grid will be said to be of size $h$.

Theorem 2.1. Let $y=\eta+e$ where
$\eta=\beta_{0}+\beta_{1} x_{1}+\beta_{2} x_{2}+\beta_{11} x_{1}{ }^{2}+\beta_{12} x_{1} x_{2}+\beta_{22} x_{2}{ }^{2}$
and let

$$
\begin{equation*}
\hat{y}=\hat{\beta}_{0}+\hat{\beta}_{1} x_{1}+\hat{\beta}_{2} x_{2} \tag{2,15}
\end{equation*}
$$

be fitted from a rectangular grid whose center is at ( 0,0 ).
Then if $e=0$,

$$
\left.\frac{\partial \eta}{\partial x_{i}}\right|_{(0,0)} \quad=\hat{\beta}_{i}, \quad: i=1,2
$$

Proof: Let the design matrix be

$$
X=\left|\begin{array}{rrr}
1 & -h & -k  \tag{2,16}\\
1 & -h & k \\
1 & h & -k \\
1 & h & k
\end{array}\right|
$$

Then by least squares,

$$
\left|\begin{array}{l}
\hat{\beta}_{0}  \tag{2.17.}\\
\hat{\beta}_{1} \\
\hat{\beta}_{2}
\end{array}\right|=\operatorname{diag} \cdot\left(\frac{1}{4}, \frac{1}{4 h^{2}}, \left.\frac{1}{4 k^{2}}| | \begin{array}{cccc}
1 & 1 & 1 & 1 \\
-h & -h & h & h \\
-k & k & -k & k
\end{array}| | \begin{array}{c}
y_{1} \\
y_{2} \\
y_{3} \\
y_{4}
\end{array} \right\rvert\,\right.
$$

so that

$$
\hat{\beta}_{1}=\frac{1}{4 h^{2}} \cdot h\left(4 h \beta_{1}\right)=\beta_{1}
$$

and

$$
\begin{equation*}
\hat{\beta}_{2}=\frac{1}{4 k^{2}} \cdot k\left(4 k \beta_{2}\right)=\beta_{2} . \tag{2,18}
\end{equation*}
$$

Now

$$
\frac{\partial \eta}{\partial x_{1}}=\beta_{1}+2 \beta_{11} x_{1}+\beta_{12} x_{2}
$$

and

$$
\begin{equation*}
\frac{\partial \eta}{\partial x_{2}}=\beta_{2}+\beta_{12} x_{1}+2 \beta_{22} x_{2} \tag{2.19}
\end{equation*}
$$

which at (0, 0) reduce to

$$
\begin{equation*}
\frac{\partial \eta}{\partial x_{1}}=\beta_{1} \text { and } \frac{\partial \eta}{\partial x_{2}}=\beta_{2} \tag{2.20}
\end{equation*}
$$

The question naturally arises as to whether the coefficients of the plane fitted without error are equal to the first order partials at an arbitrary point other than $(0,0)$. That the answer is in the
affirmative is the content of theorem 2.2.

Theorem 2.2. Let $y$ and $\hat{y}$ be as in (2.14) and (2.15) respectively and let $y$ be fitted from a rectangular grid with center at $\left(x_{1}^{0}, x_{2}^{0}\right)$. Then if $e=0$,

$$
\left.\frac{\partial \eta}{\partial x_{i}} \right\rvert\,\left(x_{1}^{0}, x_{2}^{0}\right)=\stackrel{\rightharpoonup}{\beta}_{i} . \quad i=1,2
$$

Proof: Let the design matrix be

$$
x=\left|\begin{array}{ccc}
1 & x_{1}^{0}-h & x_{2}^{0}-k  \tag{2.21}\\
1 & x_{1}^{0}-h & x_{2}^{0}+k \\
1 & x_{1}^{0}+h & x_{2}^{0}-k \\
1 & x_{1}^{0}+h & x_{2}^{0}+k
\end{array}\right|
$$

Translate the center of the design to the origin by means of the translation

$$
\begin{align*}
& x_{1}^{\prime}=x_{1}-x_{1}^{0} \\
& x_{2}^{\prime}=x_{2}-x_{2}^{0} . \tag{2.22}
\end{align*}
$$

Solving for $x_{1}$ and $x_{2}$ in (2.22) and substituting in (2.14),
there results
$\eta=\beta_{0}^{\prime}+\beta_{1}{ }^{\prime} x_{1}^{\prime}+\beta_{2}{ }^{\prime} x_{2}^{\prime}+\beta_{11}{ }^{i} x_{1}^{\prime}{ }^{2}+\beta_{12}{ }^{\prime} x_{1}{ }^{\prime} x_{2}^{\prime}+\beta_{22^{\prime}} x_{2}^{\prime^{2}}$
where

$$
\beta_{1}^{\prime}=\beta_{1}+2 \beta_{11} x_{1}^{0}+\beta_{12} x_{2}^{0}
$$

and

$$
\begin{equation*}
\beta_{2}^{1}=\beta_{2}+\beta_{12} \mathrm{x}_{1}^{0}+2 \beta_{22} \mathrm{x}_{2}^{0} \tag{2.24}
\end{equation*}
$$

The design (2.21) becomes the design (2.16) and hence, by theorem 2.1,

$$
\begin{equation*}
\hat{\beta}_{1}^{\prime}=\beta_{1}^{\prime} \text { and } \hat{\beta}_{2}^{\prime}=\beta_{2}^{\prime} . \tag{2.25}
\end{equation*}
$$

If the partial derivatives $(2.19)$, be evaluated at the point $\left(\mathrm{x}_{1}{ }^{0}, \mathrm{x}_{2}{ }^{0}\right)$, the result is the right side of (2.24).

In the presence of error, the coefficients of the fitted plane are not equal to the first order partial derivatives evaluated at the center of the design. However, that the first order partial derivatives are estimated without bias by the linear fitted coefficients will be shown in theorems 2.3 and 2.4 .

Theorem 2.3. Let $\hat{y}$ be fitted by a rectangular grid whose center is at $(0,0)$. Then

$$
\left.E\left(\hat{\beta}_{L}\right)=\frac{\partial \eta}{\partial x_{i}} \right\rvert\,(0,0) . \quad i=1,2
$$

Proof: From theorem 2.1, it is known that

$$
\left.\frac{\partial \eta}{\partial x_{i}} \right\rvert\,(0,0) \quad=\beta_{i} . \quad i=1,2
$$

It remains to show that $E\left(\hat{\beta}_{i}\right)=\beta_{i}$ but this is the content of equation (2.13).

That the first order partial derivatives are estimated without bias at an arbitrary point is the content of theorem 2.4 .

Theorem 2.4. Let $y$ be fitted by means of a rectangular grid whose center is at the point $\left(x_{1}^{0}, x_{2}{ }^{0}\right)$. Then

$$
\left.E\left(\hat{\beta}_{i}\right)=\frac{\partial \eta}{\partial x_{i}} \right\rvert\,\left(x_{1}^{0}, x_{2}^{0}\right) . \quad i=1,2
$$

Proof: The proof is accomplished by a translation to the origin as in theorem 2.2 followed by the application of theonem 2.3.


Frequently in what is to follow the response will be considered without error and the linear coefficients of the fitted plane will be used on that basis. Although in practice there will always be error, theorems 2.2 and 2.4 are comforting in that the linear coefficients are estimated without bias when error is taken into account.

Definition 2.8. The line through the center of the design in a direction determined by the fitted plane in the presence of error will be called the steepest ascent Line.

For example, for two variables $x_{1}$ and $x_{2}$, the fitted plane has the form

$$
\hat{y}=\hat{\beta}_{0}+\hat{\beta}_{1} x_{1}+\hat{\beta}_{2} x_{2}
$$

and the steepest ascent line through the origin is

$$
\begin{equation*}
x_{2}=\frac{\hat{\beta}_{2}}{\hat{\beta}_{1}} \quad x_{1} \tag{2.26}
\end{equation*}
$$

Definition 2.9. The line through a point in a direction deter mined by the first order partial derivatives at that point will be called the gradient line.

The gradient line through the origin for a quadratic surface is

$$
\begin{equation*}
x_{2}=\frac{\beta_{2}}{\beta_{1}} \quad x_{1} \tag{2.27}
\end{equation*}
$$

Responses at the Points of a Square Grid

Consider the response
$y=\beta_{0}+\beta_{1} x_{1}+\beta_{2} x_{2}+\beta_{11} x_{1}^{2}+\beta_{12} x_{1} x_{2}+\beta_{22^{x}}{ }^{2}+e$
$=\beta_{0}+\left(x_{1}, x_{2}\right)\binom{\beta_{1}}{\beta_{2}}+\left(x_{1}, x_{2}\right)\left(\begin{array}{ll}\beta_{11} & \beta_{12 / 2} \\ \beta_{12 / 2} & \beta_{22}\end{array}\right)\binom{x_{1}}{x_{2}}+\mathrm{e}$
$=B_{0}+Z B_{1}+Z B_{2} Z^{\prime}+e$
where $\quad Z=\left(x_{1}, x_{2}\right), B_{0},=\beta_{0}, B_{1}=\binom{\beta_{1}}{\beta_{2}}$
and $\quad B_{2}=\left(\begin{array}{ll}\beta_{11} & \beta_{12 / 2} \\ \beta_{12 / 2} & \beta_{22}\end{array}\right)$.

In this notation the predicted response of equation (2.8),
$\hat{y}=\hat{\beta}_{0}+\hat{\beta}_{1} x_{1}+\hat{\beta}_{2} x_{2}$, becomes $\hat{y}=\hat{B}_{0}+\mathrm{ZB}_{1}$.

Now let $B_{2}=I_{2}$ and suppose that the responses $y_{1}, y_{2}, y_{3}$, and $y_{4}$ are measured without error at the grid points $(-h,-h)$, $(-h, h),(h,-h)$ and $(h, h)$. Then at any of these grid points the response as given in (2.28) may be written as

$$
\begin{equation*}
y=B_{0}+Z B_{1}+2 h^{2} \tag{2,30}
\end{equation*}
$$

According to least squares
$\hat{\beta}_{0}=1 / 4\left(y_{1}+y_{2}+y_{3}+y_{4}\right)=1 / 4\left(4 B_{0}+4 Z I_{2} Z^{\prime}\right)=B_{0}+2 h^{2}$.
Also, since the responses are measured without error, $\hat{B}_{1}=B_{1}$ so that equation (2.29) becomes

$$
\begin{equation*}
\hat{y}=B_{0}+2 h^{2}+Z B_{1} \tag{2.32}
\end{equation*}
$$

which is seen to be the same as $y$. If the center of the grid is not at the origin, a translation of the axes to the center of the grid results in new variables for which again $y=\hat{y}$ at the grid points. These results are summarized in theorem 2.5.

Theorem 2.5. Let $y$ be measured without error and be given by (2.28) wherein $B_{2}=I_{2}$. Then the least squares plane passes through the responses at the points of a square grid.

Suppose now that the matrix $B_{2}$ is allowed to be of the form $B_{2}=\left|\begin{array}{ll}b_{11} & 0 \\ 0 & b_{22}\end{array}\right| \quad$ so that the circular contours considered
in theorem 2.5 become elliptical.

Theorem 2.6. Let $y$ be measured without error and be given by ( 2.28 ) wherein $B_{2}=$ diag. $\left(b_{11}, b_{22}\right)$. Then the least squares plane passes through the responses at the points of a square grid.

Proof: The responses at the grid points are given by

$$
\begin{equation*}
\mathrm{y}=\mathrm{B}_{0}+\mathrm{ZB}_{1}+\mathrm{h}^{2} \mathrm{~b}_{11}+\mathrm{h}^{2} \mathrm{~b}_{22} \tag{2.33}
\end{equation*}
$$

Now

$$
\begin{aligned}
\hat{B}_{0} & =1 / 4\left(y_{1}+y_{2}+y_{3}+y_{4}\right)=1 / 4\left(4 B_{0}+4 z B_{2} z^{1}\right) \\
& =1 / 4\left[4 B_{0}+4\left(h^{2} b_{11}+h^{2} b_{22}\right)\right]=B_{0}+h^{2} b_{11}+h^{2} b_{22} \cdot(2.34)
\end{aligned}
$$

Since $\hat{B}_{1}=B_{1}$,

$$
\begin{equation*}
\hat{\mathrm{y}}=\mathrm{B}_{0}+\mathrm{h}^{2} \mathrm{~b}_{11}+\mathrm{h}^{2} \mathrm{~b}_{22}+\mathrm{ZB}_{1} \tag{2.35}
\end{equation*}
$$

so that on the grid points, $y=\hat{y}$.

It can be readily verified that $y \neq \hat{y}$ when the matrix $B_{2}$ has off-diagonal elements. Theorems 2.5 and 2.6 can both be generalized but since 2.5 is a special case of 2.6 it will suffice to generalize the latter.

Theorem 2.7. Let $y$ be measured without error and be of the form $y=\eta+e$ where $\eta=B_{0}+Z B_{1}+Z B_{2} Z^{\prime}(2,36)$ and where the dimensions of $B_{0}, Z, B_{1^{\prime}}$ and $B_{2}$ are respectively $1 \times 1,1 \times p, p x+$ and $p x p$. Let $B_{2}=\operatorname{diag} \cdot\left(b_{11}, b_{22^{\prime}} \cdots b_{p p}\right)$.

Then the least squares hyperplane, $\hat{y}=\hat{B}_{0}+\hat{Z B}_{1}$, passes through the responses at the points of a p-dimensional square grid.

Proof: The responses at the grid points are given by

$$
\begin{equation*}
y=B_{0}+Z B_{1}+h^{2}\left(b_{11}+b_{22}+\ldots+b_{p p}\right) . \tag{2,37}
\end{equation*}
$$

Now

$$
\begin{align*}
\hat{\beta}_{0} & =1 / N\left(y_{1}+y_{2}+\ldots+y_{N}\right)=1 / N\left(N B_{0}+N Z B_{2} Z^{\prime}\right) \\
& =B_{0}+h^{2}\left(b_{11}+b_{22}+\ldots+b_{p p}\right) \tag{2.38}
\end{align*}
$$

where $N=2^{\mathrm{p}}$.
Since $\hat{B}_{1}=B_{1}$.

$$
\hat{y}=B_{0}+h^{2}\left(b_{11}+b_{22}+\ldots+b_{p p}\right)+z B_{1}=y
$$

This relationship between the least squares plane fitted from the grid points and the responses at these points is surprising. It is remarkable that the responses at the grid points should lie in a plane. Apparently the grid points in the factor space are members of a set of points for which the response is the same as the predicted response. What is the locus of these points in the factor space? Evidently the answer may be found by setting. $y=\hat{y}$. In order to visualize the geometry, consider the two dimensional case of theorem 2.6 with $b_{11}$ and $b_{22}$ replaced by $\beta_{11}$ and $\beta_{22}$. Then $\left\{\mathrm{x}_{1}, \mathrm{x}_{2} \mid \mathrm{y}=\hat{\mathrm{y}}\right\}$ satisfies the equation

$$
\begin{equation*}
\beta_{11} \mathrm{x}_{1}^{2}+\beta_{22} \mathrm{x}_{2}^{2}=\mathrm{h}^{2}\left(\beta_{11}+\beta_{22}\right) \tag{2.39}
\end{equation*}
$$

which may be written as


Suppose that the contours are elliptical and that $y$ attains a maximum value at some point in the factor space so that $\beta_{11}$ and $\beta_{22}$ are both negative. Then $\left\{x_{1}, x_{2} \mid y>y\right\}$ lies within the ellipse given in (2.40) while $\left\{\mathrm{x}_{1}, \mathrm{x}_{2} \mid y<\hat{y}\right\}$, lies outside the ellipse given in (2.40).

If in $(2.40), \beta_{11}=\beta_{22},(2,40)$ may be written as

$$
\begin{equation*}
x_{1}^{2}+x_{2}^{2}=2 h^{2} \tag{2,41}
\end{equation*}
$$

which is a circle of radius $\sqrt{2 h}$. Suppose that an experiment is performed at a non-grid point on this circle. If the observed response, $y$, and the predicted response, $\hat{y}$, are nearly the same at this point, it indicates that $\beta_{11}$ and $\beta_{22}$ are approximately of the same magnitude. If in addition, the experimenter has reason to believe that $\beta_{12}=0$, then circular contours are indicated.

## Choice of Design Size

Suppose that the $\beta^{\prime} s$ of the fitted plane are estimated by means of a square grid of size $h$ whose center is at ( 0,0 ).. Then

$$
X=\left(\begin{array}{ccc}
1 & -h & -h  \tag{2.42}\\
1 & -h & h \\
1 & h & -h \\
1 & h & h
\end{array}\right)
$$

and

$$
\left(X^{\prime} X\right)^{-1}=\text { diag. }\left(\frac{1}{4}, \frac{1}{4 h^{2}}, \frac{1}{4 h^{2}}\right)
$$

so that var. $\hat{\beta}_{1}=$ var. $\hat{\beta}_{2}=\frac{\sigma^{2}}{4 h^{2}}$.

It would seem advisable to take $h$ as large as possible so as to reduce the variance. However, enlargement of the design, while reducing the variance, generally increases the bias. The failure of the estimator $\hat{y}$ is due to the unavoidable variability of the data and to the inadequacy of the assumed model in representing the true model. A measure which takes into account both of these factors is the mean-square-error, hereafter referred to as $m$, $s$. e,

By definition,

$$
\begin{equation*}
\text { m.s.e. }\left(\hat{\beta}_{1}\right)=E\left(\hat{\beta}_{1}-\beta_{1}\right)^{2} \text {. } \tag{2,43}
\end{equation*}
$$

This may be written so as to display the relationship between m.s.e., variance and bias as

$$
\begin{equation*}
E\left(\hat{\beta}_{1}-\beta_{1}\right)^{2}=E\left[\hat{\beta}_{1}-E\left(\hat{\beta}_{1}\right)\right]^{2}+\left[E\left(\hat{\beta}_{1}\right)-\beta_{1}\right]^{2} \tag{2.44}
\end{equation*}
$$

where the first term on the right is the var $\hat{\beta}_{1}$ and the second term is the $\left(\text { bias in } \hat{\beta}_{1}\right)^{2}$.

The following discussion attempts to give an indication of the grid
size required in order to minimize the $m$.s.e.
Suppose that the experimenter has chosen the model

$$
\begin{equation*}
y=\eta+e \tag{2.45}
\end{equation*}
$$

where

$$
\begin{equation*}
\eta=\beta_{0}+b_{1} x_{1}+\beta_{2} x_{2}+\beta_{11} x_{1}^{2}+\beta_{12} x_{1} x_{2}+\beta_{22} x_{2}^{2} \tag{2.46}
\end{equation*}
$$

but in reality,

$$
\begin{equation*}
\mathrm{y}=\eta_{\mathrm{l}}+\mathrm{e} \tag{2.47}
\end{equation*}
$$

where

$$
\begin{equation*}
\eta_{1}=\eta+\beta_{111} x_{1}^{3} \tag{2.48}
\end{equation*}
$$

Now

$$
\left(\begin{array}{l}
\hat{\beta}_{0}  \tag{2.49}\\
\hat{\beta}_{1} \\
\hat{\beta}_{2}
\end{array}\right)=\operatorname{diag} \cdot\left(\frac{1}{4}, \frac{1}{4 h^{2}}, \frac{1}{4 h^{2}}\right)\left(\begin{array}{cccc}
1 & 1 & 1 & 1 \\
-h & -h & h & h \\
-h & h & -h & h
\end{array}\right)\left(\begin{array}{l}
y_{1} \\
y_{2} \\
y_{3} \\
y_{4}
\end{array}\right)
$$

so that

$$
\begin{equation*}
E\left(\hat{\beta}_{1}\right)=\frac{1}{4 h}\left(4 h \beta_{1}+4 h^{3} \beta_{111}-e_{1}-e_{2}+e_{3}+e_{4}\right)=\beta_{1}+h^{2} \beta_{111} . \tag{2.50}
\end{equation*}
$$

Hence,
$(\text { bias in } \hat{\beta})^{2}=\left[E\left(\hat{\beta}_{1}\right)-\beta_{1}\right]^{2}=\left(\beta_{1}+h^{2} \beta_{111}-\beta_{1}\right)^{2}=h^{4} \beta_{111}^{2}$.
Substituting the var $\hat{\beta}_{1}$ and the $\left(\text { bias in } \hat{\beta}_{1}\right)^{2}$ into $(2.44)$ there results,

$$
\begin{equation*}
\text { m.s.e. }\left(\hat{\beta}_{1}\right)=\frac{\sigma^{2}}{4 h^{2}}+h^{4} \beta_{111}^{2} \tag{2.52}
\end{equation*}
$$

Now $\beta_{\text {lll }}$ is generally unknown but may be expressed as

$$
\begin{equation*}
\beta_{111}=c \sigma^{2} \tag{2.53}
\end{equation*}
$$

and on substitution in (2.52) there results

$$
\begin{equation*}
\text { m.s.e. }\left(\hat{\beta}_{1}\right)=\frac{\sigma^{2}}{4 h^{2}}+h^{4} c^{2} \sigma^{4} \tag{2.54}
\end{equation*}
$$

Differentiating (254) with respect to $h$ and solving for the $h$ which minimizes m.s.e. ( $\hat{\beta}_{1}$ ),

$$
\begin{equation*}
\frac{\mathrm{d} \text { m.s.e. }}{\mathrm{dh}}\left(\hat{\beta}_{1}\right)=-\frac{2 \sigma^{2}}{4 h^{3}}+4 h^{3} c^{2} \sigma^{4}=0 \tag{2.55}
\end{equation*}
$$

and

$$
\begin{equation*}
h=\frac{\sqrt{2}}{2(c \sigma)^{1 / 3}} \tag{2.56}
\end{equation*}
$$

If (2.56) is differentiated with respect to $h$ the result is

$$
\begin{equation*}
\frac{d^{2} m \text { s.e. }\left(\hat{\beta}_{1}\right)}{d h^{2}}=\frac{3 \sigma^{2}}{2} h^{-4}+12 h^{2} c^{2} \sigma^{4} \tag{2.57}
\end{equation*}
$$

Since all of the exponents are even, (2.57) is positive which insures that the $h$ of (2.56) minimizes the m.s.e. ( $\hat{\beta}_{\mathrm{L}}$ ).

Now suppose that instead of (2.47), the actual response has the form

$$
\begin{equation*}
y=\eta_{2}+e \tag{2.58}
\end{equation*}
$$

where

$$
\begin{equation*}
\eta_{2}=\eta+\beta_{222} x_{2}^{3} \tag{2.59}
\end{equation*}
$$

In this case, it is readily verified that $E\left(\hat{\beta}_{1}\right)=\beta_{1}$ so that $\left(\right.$ bias in $\left.\hat{\beta}_{1}\right)=0$.

The use of similar calculations produces Table $I$. Table $I$ is a summary of the values of $h$ required to minimize the associated mean square errors in the presence of all possible combinations of cubic terms. These combinations are listed in the far left column. In this column the number 1 stands for, $\beta_{111} x_{1}{ }^{3}, 2$ for $\beta_{222^{\prime}}{ }^{3}$, 3 for $\beta_{112} x_{l}{ }^{2} x_{2}$ and 4 for $\beta_{122} x_{1} x_{2}^{2}$. With this coding the mixture $\beta_{111} \mathrm{x}_{1}{ }^{3}+\beta_{112} \mathrm{x}_{1}^{2} \mathrm{x}_{2}+\beta_{122} \mathrm{x}_{1} \mathrm{x}_{2}{ }^{2}$ is written simply as $1,3,4$.

DESIGN SIZE NECESSARY TO MINIMIZE M. S. E. $\hat{\beta}_{1}$

$$
\begin{aligned}
& \text { Cubic Terms } \\
& 1 \\
& \text { M. S. E. } \hat{\beta}_{1} \\
& \sigma^{2} / 4 h^{2}+h^{4} \beta_{1 l l}^{2} \\
& 2 \quad \sigma^{2} / 4 h^{2} \\
& 3 \quad \sigma^{2} / 4 h^{2} \\
& 4 \\
& \sigma^{2} / 4 h^{2}+\beta_{122}{ }^{2} h^{4} \\
& 1,2 \quad \sigma^{2} / 4 h^{2}+h^{4} \beta_{111}^{2} \\
& 1,3 \quad \sigma^{2} / 4 h^{2}+h^{4} \beta_{111}^{2} \\
& {\left.\underline{(B i a s ~ i n ~} \hat{\beta}_{1}\right)^{2}}^{h} \text { to min M.S.E. } \\
& h^{4} \beta_{111}{ }^{2} \\
& \sqrt{2 / 2(c \sigma)^{1 / 3}} \\
& 0 \text { h large } \\
& 0 \quad h \text { large } \\
& h^{4} \beta_{122^{2}} \frac{\sqrt{2}}{2(c \sigma)^{1 / 3}} \text {. } \\
& h^{4} \beta_{111}^{2} \quad \frac{\sqrt{2}}{2(c \sigma)^{1 / 3}} \\
& 1,4 \quad \sigma^{2} / 4 h^{2}+h^{4}\left(\beta_{11 L^{+}} \beta_{122}\right)^{2} \\
& h^{4}\left(\beta_{111}+\beta_{12}\right)^{2} \quad \frac{2^{1 / 6}}{2(c \sigma)^{1 / 3}} \\
& 2,3 \quad \sigma^{2} / 4 h^{2} \\
& 2,4 \quad \sigma^{2} / 4 h^{2}+\beta_{122^{2} h^{4}} \\
& h^{4} \beta_{112}{ }^{2} \\
& \text { h large } \\
& 3,4 \quad \sigma^{2} / 4 h^{2}+\beta_{122}{ }^{2} h^{4} \\
& h^{4} \beta_{112} \quad \frac{\sqrt{2}}{2(c \sigma)^{1 / 3}} \\
& 1,2,3 \quad \sigma^{2} / 4 h^{2}+h^{4} \beta_{111}^{2} \\
& h^{4} \beta_{111}^{2} \quad \frac{\sqrt{2}}{2(c \sigma)^{1 / 3}} \\
& \text { 1, 3,4 Like } 1,4 \\
& \text { 2, 3, } 4 \text { Like } 4 \\
& \text { 1, 2, 3, } 4 \text { Like } 1,4
\end{aligned}
$$

## CHAPTER III

## PROPERTIES OF THE METHOD OF <br> STEEPEST ASCENT

In accordance with the method of steepest ascent, an initial set of experiments are performed in order to determine the steepest ascent path. Additional experiments are performed at points along the path until it is felt that further appreciable gain cannot be realized. No use is made of information available at these points other than to decide whether or not additional experimentation should be carried on along the path.

Now it would seem that an acceptable strategy would be to run the first additional point somewhere on the steepest ascent line and use this point together with the original design points to calculate a new path of steepest ascent. This process could be repeated, altering the path on the basis of information obtained at each additional point. The rather surprising fact for certain designs is that, if additional points are taken on the path and used together with the original points to calculate a new path, the new path turns out to be the same as the original one. This result will be proved after the introduction of formulae which simplify the calculation of the path based on additional points.

Suppose that $N_{l}$ initial observations are taken at points in the factor space where the response is given by $y=\eta+e$ where $\eta$ is quadratic in the $p$ variables, $x_{1}, x_{2}, \ldots, x_{p}$. In matrix notation the $N_{l}$ responses at the $N_{l}$ design points can be written

$$
\begin{equation*}
Y_{1}=X_{1} \beta+E_{1} \tag{3.1}
\end{equation*}
$$

where $X_{1}$ is the $N_{1} x q$ design matrix, $q$ is the number of parameters in the model, $\beta$ is a $q$ x 1 vector of unknown parameters, $E_{l}$ is an $N_{1} \times l$ vector of random errors, and $Y_{1}$ is the $N_{1} \times 1$ vector of responses. Let

$$
\begin{equation*}
\hat{y}=x \hat{\beta} \tag{3.2}
\end{equation*}
$$

be the least squares prediction of $y$ based on the initial design $X_{1}$,
Now suppose that $N_{2}$ additional observations are taken.
Denote the associated $N_{2} \times q$ design matrix by $X_{2}$ and the response vector by $Y_{2}$. Let $X$ and $Y$ be the design matrix and response vector associated with the $N_{1}+N_{2}=N$. observations and let $\hat{a}$ be the least-squaresestimate of $\hat{\beta}$ where $\hat{a}$ is based upon all of the $N$ observations. Further let

$$
\begin{align*}
& \left(X_{1}^{\prime} X_{1}\right)=S_{1} \\
& \left(X_{2}^{\prime} X_{2}\right)=S_{2} \\
& \left(X^{\prime} X\right)=S . \tag{3.3}
\end{align*}
$$

Then

$$
\begin{align*}
S^{-1} & =\left(X_{1}^{\prime}, X_{2}^{\prime}\right)\binom{X_{1}}{X_{2}}^{-1}=\left(X_{1}^{\prime}: X_{1}+X_{2} X_{2}\right)^{-1}=\left(S_{1}+S_{2}\right)^{-1} \\
& =\left(I+S_{2} S_{1}^{-1}\right) S_{1}=S_{1}^{-1}\left(I+S_{2} S_{1}^{-1}\right)^{-1}=S_{1}^{-1}\left(I+X_{2} X_{2} S_{1}^{-1}\right)^{-1} .
\end{align*}
$$

Now by using the relationship

$$
\begin{equation*}
(I+A B)^{-1}=I-A(I+B A)^{-1} B \tag{3.5}
\end{equation*}
$$

in (3.4) and associating $X_{2}^{\prime}$ with $A$ and $X_{2} S_{1}^{-1}$ with $B_{1}$ there results

$$
\begin{equation*}
S^{-1}=S_{1}^{-1}-S_{1}^{-1} X_{2}^{\prime}\left(I+X_{2} S_{1}^{-1} X_{2}^{\prime}\right)^{-1} X_{2} S_{1}^{-1} \tag{3.6}
\end{equation*}
$$

Now let

$$
\begin{align*}
& J=X_{2} S_{1}^{-1} \\
& R=X_{2} S_{1}^{-1} X_{2}^{\prime} \\
& G=(I+R)^{-1} \\
& Q=Y_{2}-X_{2} \hat{\beta} \tag{3,7}
\end{align*}
$$

Then (3.6) may be written

$$
\begin{equation*}
S^{-1}=S_{1}^{-1}-J^{\prime} G J . \tag{3.8}
\end{equation*}
$$

Now

$$
\hat{a}=S^{-1} X^{\prime} Y=\left(S_{1}^{-1}-J^{\prime} G J\right)\left(X_{1}^{\prime}, X_{2}^{\prime}\right)\binom{Y_{1}}{Y_{2}}
$$

$$
\begin{align*}
& =\left(S_{1}^{-1}-J^{\prime} G J\right)\left(X_{1}^{\prime} Y_{1}+X_{2}^{\prime} Y_{2}\right) \\
& =S_{1}^{-1} X_{1}^{\prime} Y_{1}+S_{1}^{-1} X_{2}^{\prime} Y_{2}-J^{\prime} G J X_{1}^{\prime} Y_{1}-J^{\prime} G J X_{2}^{\prime} Y_{2} \tag{3.9}
\end{align*}
$$

The second term of (3.9) may be written

$$
\begin{equation*}
S_{1}^{-1} X_{2}^{\prime} Y_{2}=J^{\prime} Y_{2}=J^{\prime} G G^{-1} Y_{2}=J^{\prime} G\left(I+J X_{2}^{\prime}\right) Y_{2}=J^{\prime} G Y_{2}+J^{\prime} G J X_{2}^{\prime} Y_{2} \tag{3,10}
\end{equation*}
$$

Substituting (3.10) in (3.9) there results

$$
\begin{equation*}
\hat{a}=\hat{\beta}+J^{\prime} G Y_{2}-J^{\prime} G X_{2} \hat{\beta}=\hat{\beta}+J^{\prime} G\left(Y_{2}-X_{2} \hat{\beta}\right)=\hat{\beta}+J^{\prime} G Q \tag{3.11}
\end{equation*}
$$

The advantage in using (3.8) and (3.11) is that in finding $S^{-1}$ and $\hat{a}$, the dimension of the matrix $(I+R)$ which must be inverted is $\mathrm{N}_{2} \times \mathrm{N}_{2}$. These expressions can be used profitably whenever $N_{2}$ is less than $q$.

The method for finding $\hat{a}$ and $S^{-1}$ based on new observations with a minimum of calculation was set forth by R. L. Plackett (5) and the expressions for $\hat{a}$ and $S$ but not their derivation, are to be found in (1).

Theorem 3.1. Let a unit grid with center at the origin be used to determine the steepest ascent line $x_{2}=\frac{\hat{\beta}_{2}}{\hat{\beta}_{1}} x_{1}$. If a second steepest ascent line, $x_{2}=\frac{\hat{a}_{2}}{\hat{a}_{1}} x_{1}$ is determined by means of the initial grid and an additional point on the first line, the second
line will be the same as the first line.

Proof: The grid design is given as

$$
X_{1}=\left(\begin{array}{rrr}
1 & -1 & -1  \tag{3.12}\\
1 & -1 & 1 \\
1 & 1 & -1 \\
1 & 1 & 1
\end{array}\right)
$$

and the fifth point may be taken as

$$
\begin{equation*}
X_{2}=\left(1, h \hat{\beta}_{1}, h \hat{\beta}_{2}\right) \tag{3.13}
\end{equation*}
$$

Now by (3.11)

$$
\hat{a}=\hat{\beta}+J^{\prime} G Q
$$

and since in the present application $N_{2}=1$, bath $G$ and $Q$ are scalars. Now

$$
\begin{equation*}
J^{\prime}=S_{1}^{-1} X_{2}^{\prime}=\frac{1}{4} I_{3} X_{2}^{\prime}=\frac{1}{4} X_{2}^{\prime} \tag{3.14}
\end{equation*}
$$

so that

$$
\left|\begin{array}{l}
\hat{a}_{0}  \tag{3,15}\\
\hat{a}_{1} \\
\hat{a}_{2}
\end{array}\right|=\left|\begin{array}{l}
\hat{\beta}_{0} \\
\hat{\beta}_{1} \\
\hat{\beta}_{2}
\end{array}\right|+\frac{1}{4}\left|\begin{array}{c}
1 \\
h \hat{\beta}_{1} \\
h \hat{\beta}_{2}
\end{array}\right| \quad G Q:
$$

It follows that

$$
\begin{aligned}
& \hat{a}_{0}=\hat{\beta}_{0}+\frac{1}{4} G Q=\hat{\beta}_{0}\left(1+\frac{G Q}{4 \beta_{0}}\right) \\
& \hat{\dot{a}}_{1}=\hat{\beta}_{1}+\frac{1}{4} h \hat{\beta}_{1} G Q=\hat{\beta}_{1}\left(1+\frac{h G Q}{4}\right)
\end{aligned}
$$

$$
\begin{equation*}
\hat{a}_{2}=\hat{\beta}_{2}+\frac{1}{4} \mathrm{~h} \hat{\beta}_{2} G Q=\hat{\beta}_{2}\left(1+\frac{h G Q}{4}\right) . \tag{3.16}
\end{equation*}
$$

Hence

$$
\begin{gathered}
\frac{\hat{a}_{2}}{\hat{a}_{1}}=\frac{\hat{\beta}_{2}}{\hat{\beta}_{1}} \\
\text { so that } x_{2}=\frac{\hat{\beta}_{2}}{\hat{\beta}_{1}} \quad x_{1} \text { is the same line as } x_{2}=\frac{\hat{a}_{2}}{\hat{a}_{1}} \quad x_{1} .
\end{gathered}
$$

It is to be noted that the proportionality of the coefficients in the last proof did not extend to $\hat{\beta}_{0}$ and $\hat{a}_{0}$. The factor $h$ was not specified and it may be seen from (3.15) or (3.16) that if $h$ is chosen so that

$$
\begin{equation*}
h \hat{\beta}_{0}=1, \tag{3.18}
\end{equation*}
$$

then (3.15) becomes

$$
\begin{equation*}
\hat{a}=\hat{\beta}+\frac{1}{4} h G Q \hat{\beta}=\left(I+\frac{1}{4} h G Q\right) \hat{\beta} \tag{3,19}
\end{equation*}
$$

from which it is apparent that the proportionality property does extend to $\quad \hat{a}_{0}$ and $\hat{\beta}_{0}$.

The magnitude of $h$ determines the step size taken along the steepest ascent line. Apparently the choice of $h$ as given by (3.18) has no particular merit other than simplifying the calculations.

Definition 3.1. If the second steepest ascent line is the same as the first, the steepest ascent line will be said to be invariant. Thus the mathematical expression of invariance is

$$
\begin{equation*}
\frac{\hat{a}_{j}}{\hat{a}_{i}}=\frac{\hat{\beta}_{j}}{\hat{\beta}_{i}}, i, j=1,2, \ldots, p \tag{3.20}
\end{equation*}
$$

Theorem 3.1 proves the property of invariance for $p=2$. It is natural to inquire if invariance holds in $p$-dimensional space in general. That it does for certain designs is affirmed in theorem 3.2.

Theorem 3.2. Let $\hat{\beta}$ be determined from the $2^{p}=N_{1}$ points of a p-dimensional unit grid design with center at the origin. Let $\hat{a}$ be determined from these points and the additional point $X_{2}=\left(1, h \hat{\beta}_{1}, \ldots, h \hat{\beta}_{p}\right)$. Then

$$
\frac{\hat{\beta}_{j}}{\hat{\beta}_{i}}=\frac{\hat{a}_{j}}{\hat{a}_{i}}, \quad i, j=1,2, \ldots, p .
$$

Proof:

$$
\begin{array}{r}
\left(X_{1}^{\prime} X_{1}\right)^{-1}=\frac{1}{N_{1}} I_{p+1} \quad \text { so that } \\
J^{\prime}=\frac{1}{N_{1}} I_{p+1} X_{2}^{\prime}=\frac{1}{N_{1}}\left|\begin{array}{c}
1 \\
h \hat{\beta}_{1} \\
\cdot \\
\cdot \\
\dot{h} \hat{\beta}_{p}
\end{array}\right|
\end{array}
$$

Now by (3.11), $\hat{a}=\hat{\beta}+J^{\prime} G Q$ and since $N_{2}=1, G$ and $Q$ are scalars. Substituting in (3.11) there results

$$
\left(\left.\begin{array}{c}
\hat{a}_{0}  \tag{3.21}\\
\hat{a}_{1} \\
\cdot \\
\cdot \\
\hat{a}_{p}
\end{array}\left|=\left|\begin{array}{c}
\hat{\beta}_{0} \\
\hat{\beta}_{1} \\
\cdot \\
\cdot \\
\hat{\beta}_{p}
\end{array}\right|+\frac{G Q}{N_{1}}\right| \begin{array}{c}
1 \\
h \hat{\beta}_{1} \\
0 \\
\\
\hat{h}_{p}
\end{array} \right\rvert\,\right.
$$

so that

$$
\hat{a}_{i}=\left(1+\frac{G Q h}{N_{1}}\right) \hat{\beta}_{i} \quad i=1,2, \ldots, p
$$

from which (3.20) follows.

The relationship between the design matrix, the additional points, and the invariance property is given in theorem 3.3. Consider the following conditions:
(a) $\frac{\hat{a}_{j}}{\hat{a}_{i}}=\frac{\hat{\beta}_{j}}{\hat{\beta}_{i}} \quad i, j=1 ; 2, \ldots, p$
(b) $\quad X_{2}=\left(1, h_{\hat{\beta}}, h \hat{\beta}_{2}, \cdots, h \hat{\beta}_{p}\right)$
(c) $\quad\left(X_{1}^{\prime} X_{1}\right)^{-1}=\left(\begin{array}{cc}u & 0 \\ 0 & v I_{p}\end{array}\right)$.

Condition (a) is the mathematical statement of invariance, Condition (b) insures that $X_{2}$ is in the direction of steepest ascent while condition (c) deals with the form of the design matrix.

Theorem 3.3. Any two of the above three conditions implies the third.

Proof: Let $\hat{a}, \hat{\beta}, S_{1}^{-1}$, and $X_{2}$ be partitioned as indicated.

$$
\hat{a}=\left|\begin{array}{c}
\hat{a}_{0} \\
\hat{a}_{1} \\
\cdot \\
\cdot \\
\cdot \\
\hat{a}_{p}
\end{array}\right|=\binom{A_{0}}{A_{1}}, \quad\left(\begin{array}{c}
\hat{\beta}_{0} \\
\hat{\beta}_{1} \\
\cdot \\
\cdot \\
\cdot \\
\hat{\beta}_{p}
\end{array}\right)=\binom{B_{0}}{B_{1}}
$$

$S_{1}^{-1}=\left(\begin{array}{c|cccc}u & 0 & 0 & \cdots & 0 \\ \hline 0 & v & 0 & \cdots & 0 \\ 0 & 0 & \cdot & & \cdot \\ \cdot & \cdot & & \cdot & \cdot \\ \cdot & \cdot & & \cdot & \cdot \\ \cdot & \cdot & & \cdot & \cdot \\ 0 & 0 & \cdot & \cdots & v\end{array}\left|=\left|\begin{array}{ll}\mu & 0 \\ 0 & v I_{p}\end{array}\right| \quad(3,22)\right.\right.$
and

$$
X_{2}=\left(1 h \hat{\beta}_{1} \cdot \ldots h \hat{\beta}_{p}\right)=\left(X_{20}, X_{21}\right)=\left(1, h B_{1}^{\prime}\right)
$$

To prove that (b) and (c) imply (a), the equation $\hat{\alpha}=\hat{\beta}+J^{\prime} G Q$ is written using the above partitions as

$$
\binom{A_{0}}{A_{1}}=\binom{B_{0}}{B_{1}}+\left(\begin{array}{cc}
u & 0  \tag{3.23}\\
0 & v I_{p}
\end{array}\right) \quad\binom{1}{h B_{1}} \quad G Q
$$

On equating corresponding members and using the values indicated in ( 3.22 ), there results

$$
A_{0}=B_{0}+u G Q \text { and }
$$

$$
\begin{equation*}
A_{1}=B_{1}+v I_{p} h_{1} B_{1} G Q . \tag{3.24}
\end{equation*}
$$

Now $G Q$ is a scalar so that

$$
\begin{equation*}
A_{1}=B_{1}(1+v h G Q) \tag{3.25}
\end{equation*}
$$

from which condition (a) follows.
To prove that (a) and (c) imply (b), by condition (a) , $A_{1}$ may be written as $k \cdot B_{1}$. Hence

$$
\binom{A_{0}}{k B_{1}}=\binom{B_{0}}{B_{1}}+\left(\begin{array}{cc}
u & 0 \\
0 & v I_{p}
\end{array}\right) \quad\binom{1}{X_{21}} \quad G Q
$$

so that

$$
\begin{align*}
A_{0} & =B_{0}+u G Q  \tag{3.26}\\
k B_{1} & =B_{1}+v I_{p} X_{21}^{\prime} G Q \tag{3.27}
\end{align*}
$$

It follows from (3.27) that

$$
(k-1) B_{1}=G Q v X_{21}^{\prime} \text {. Hence } X_{21}=\frac{k-1}{v G Q} B_{1}^{\prime}
$$

which indicates that $X_{2 l}$ is a scalar multiple of $B_{1}^{\prime}$ as in condition (b).

In order to prove that (a) and (b) imply (c), the symmetric matrix $S_{l}^{-1}$ is partitioned as

$$
S_{1}^{-1}=\left(\begin{array}{l|llll}
c_{00} & c_{01} & \cdot & \cdot & \cdot \\
\hline c_{10} & c_{11} & \cdot & \cdot & c_{0 p} \\
\cdot & \cdot & & c_{1 p} \\
\cdot & \cdot & & & \\
\cdot & \cdot & & & \\
c_{p 0} & c_{p l} & \cdot & \cdot & \cdot \\
c_{p p}
\end{array}\right)=\left(\begin{array}{ll}
c_{00} & c_{01} \\
c_{10} & c_{11}
\end{array}\right) .
$$

Now

$$
J=\left(1, h B_{1}^{\prime}\right)\left(\begin{array}{ll}
C_{00} & C_{01} \\
C_{10} & C_{11}
\end{array}\right)
$$

so that

$$
J=\left(\begin{array}{ll}
c_{00} & c_{01} \\
c_{10} & c_{11}
\end{array}\right)\binom{1}{h B_{1}}=\left(\begin{array}{lll}
c_{00} & + & c_{01}^{h B_{1}} \\
c_{10} & + & c_{11}^{h B_{1}}
\end{array}\right)
$$

Hence

$$
\begin{equation*}
\binom{A_{0}}{A_{1}}=\binom{B_{0}}{B_{1}}+\binom{C_{00}+C_{01}{ }^{h} B_{1}}{C_{10}+C_{11} h_{1} B_{1}} \quad G Q \tag{3,28}
\end{equation*}
$$

On equating corresponding elements of (3.28), there results

$$
\begin{equation*}
A_{0}=B_{0}+C_{00} G Q+C_{01}^{h B_{1}} G Q \tag{3.29}
\end{equation*}
$$

and

$$
\begin{equation*}
A_{1}=B_{1}+C_{10} G Q+C_{11} \mathrm{hB}_{1} G Q \tag{3,30}
\end{equation*}
$$

Now by (a), $A_{1}$ may be written as $k B_{1}$ sothat (3.30) becomes

$$
\begin{equation*}
k B_{1}=B_{1}+C_{10} G Q+C_{11} h B_{1} G Q \tag{3.31}
\end{equation*}
$$

Since (3.31) holds for all $B_{1}, C_{10}=0$ so that (3.31) becomes

$$
\begin{equation*}
\left(k I_{p}-h_{Q Q C_{L 1}}\right) B_{1}=B_{1} . \tag{3,32}
\end{equation*}
$$

It follows that

$$
\begin{equation*}
\left(k I_{p}-h_{\mu Q C}\right)=I_{p} \tag{3.33}
\end{equation*}
$$

so that

$$
c_{11}=\frac{k I_{p}-I_{p}}{h G Q}=\frac{(k-1) I_{p}}{h G Q}
$$

Hence $\quad C_{11}$ is a diagonal matrix as in condition (c). Since $C_{10}=0$, the expression for $A_{0}$ in (3.29) becomes

$$
\begin{equation*}
A_{0}=B_{0}+C_{00} G Q \tag{3.34}
\end{equation*}
$$

so that

$$
\begin{equation*}
C_{00}=\frac{A_{0}-B_{0}}{G Q} \tag{3.35}
\end{equation*}
$$

Note that if the invariance property is extended to $\hat{a}_{0}$ and $\hat{\beta}_{0}$, then $A_{0}=k B_{0}$ and, as in (3.18), $B_{0}=1 / h$. Hence (3.35) becomes

$$
C_{00}=\frac{k B_{0}-B_{0}}{G Q}=\frac{(k-1) B_{0}}{G Q}=\frac{k-1}{G Q h}
$$

which may be recognized as the scalar multiplier of $I_{p}$ occurring in the expression for $C_{11}$.

Suppose that instead of just one additional point, $\mathrm{N}_{2}$ additional points were taken in the steepest ascent direction. Would the new steepest ascent direction based on the original grid and the $N_{2}$ additional points be the same as the old? Again the answer is in the affirmative and this is the content of theorem 3.4.

> Theorem 3.4. Let $\hat{\beta}$ be determined from the $2^{p}=N_{1}$ points of a p-dimensional unit grid. Let $\hat{a}$ be determined from the $N_{1}$ initial points and $N_{2}$ additional points which lie in the steepest ascent direction. Then

$$
\frac{\hat{\beta}_{j}}{\hat{\beta}_{i}}=\frac{\hat{a}_{j}}{\hat{a}_{i}} \quad i, j=1,2, \ldots, p
$$

Proof: The $N_{2} \times p+1$ matrix, $X_{2}$, may be written and partitioned as

$$
X_{2}=\left|\begin{array}{lllll}
1  \tag{3.36}\\
1 \\
\cdot & h_{1} \hat{\beta}_{1} & h_{1} \hat{\beta}_{2} & \cdots & h_{1} \hat{\beta}_{p} \\
h_{2} \hat{\beta}_{1} & h_{2} \hat{\beta}_{2} & \cdots & h_{2} \hat{\beta}_{p} \\
\cdot & \cdot & & \cdot \\
\cdot & \cdot & & \cdot \\
\cdot & \cdot & & \cdot \\
{ }_{h_{N}} \hat{\beta}_{1}{ }^{h_{1}}{ }_{N} \hat{\beta}_{2} & & h_{N_{2}} \hat{\beta} p
\end{array}\right| \cdot
$$

Now $W_{L}$ may be written as

$$
W_{1}=-\left|\begin{array}{c}
h_{1} B_{1}^{\prime}  \tag{3.37}\\
h_{2} B_{1}^{\prime} \\
\cdot \\
\cdot \\
\cdot \\
h_{N_{2}} B_{1}^{\prime}
\end{array}\right|=\left|\begin{array}{c}
h_{1} \\
h_{2} \\
\cdot \\
\cdot \\
\cdot \\
h_{N_{2}}
\end{array}\right| B_{1}^{\prime}=H B_{1}^{\prime}
$$

where

$$
H=\left|\begin{array}{c}
h_{1} \\
h_{2} \\
\cdot \\
\cdot \\
\cdot \\
h_{N_{2}}
\end{array}\right| \quad \text { and } \quad B_{1}=\left|\begin{array}{c}
\hat{\beta}_{1} \\
\hat{\beta}_{2} \\
\cdot \\
\cdot \\
\hat{\beta}_{p}
\end{array}\right|
$$

Then (3.36) may be written as

$$
X_{2}=\left(W_{0}, W_{1}\right)=\left(W_{0}, H B_{1}^{\prime}\right) .(3,38)
$$

Now

$$
\begin{equation*}
J=X_{2} S_{1}^{-1}=\left(W_{0}, H B_{1}^{\prime}\right) S_{1}^{-1} \tag{3.39}
\end{equation*}
$$

and
$\mathrm{S}_{1}^{-1}$ may be partitioned as

Then using $\mathrm{S}_{\mathrm{l}}{ }^{-1}$ as given by (3.40) in (3.39), J may be written as

$$
J=\left(W_{0}, H B_{1}^{\prime}\right)\left(\begin{array}{ll}
C_{00} & 0 \\
0 & C_{11}
\end{array}\right)=\left(W_{0} C_{00}, H B_{1}^{\prime} C_{11}\right)
$$

so that

$$
J^{\prime}=\left(\begin{array}{ll}
C_{00} & W^{\prime}  \tag{3,41}\\
C_{11} B_{1} H^{\prime}
\end{array}\right) \text {. }
$$

Although explicit expressions may be written for $G$ and $Q$, it is only necessary to note their dimensions. Now $G=(I+R)^{-1}=$ $\left(I+X_{2} S_{1}^{-1} X_{2}^{\prime}\right)^{-1}$ so that the dimension of $G$ is that of $X_{2} S_{1}^{-1} X_{2}^{\prime}$ which is $N_{2} \times N_{2}$. Also, $Q=Y_{2}-X_{2} \hat{\beta}$ has dimension $N_{2} \times 1$ so that the product $G Q$ has dimension $N_{2} \times 1$. Now substituting the expression (3.41) for $J^{\prime}$ in the equation for $\hat{a}$ and
partitioning $\hat{\alpha}$ as indicated there obtains

$$
\hat{a}=\left(\begin{array}{c}
\hat{a}_{0} \\
\frac{\hat{a}_{1}}{a_{1}} \\
\cdot \\
\cdot \\
\hat{a}_{p}
\end{array}\right)=\binom{A_{0}}{A_{1}}=\binom{B_{0}}{B_{1}}+\binom{C_{00} W_{0}^{\prime}}{C_{11} B_{1^{\prime}} H^{\prime}} G Q
$$

Hence,

$$
A_{0}=B_{0}+C_{00} W_{0}^{\prime} G Q
$$

and

$$
\begin{equation*}
A_{1}=B_{1}+C_{11} B_{1} H^{\prime} G Q \tag{3.43}
\end{equation*}
$$

Now since $H^{\prime}$ is $1 \times N_{2}$ and $G Q$ is $N_{2} \times 1, H!G Q$ is a scalar. Also $C_{11}=\frac{1}{N_{1}} I_{p}$ so that

$$
A_{1}=B_{1}+\frac{H^{\prime} G Q}{N_{1}} B_{L}=B_{1}\left(1+\frac{H^{\prime} G Q}{N_{L}}\right) \cdot(3,44)
$$

Let $1+\frac{H^{\prime} G Q}{N_{l}}=k$, so that (3.41) becomes $A_{L}=k B_{L}$, that is

$$
\left|\begin{array}{c}
\hat{a}_{1} \\
\hat{a}_{2} \\
\cdot \\
\cdot \\
\hat{a}_{p}
\end{array}\right|=k \quad\left|\begin{array}{c}
\hat{\beta}_{1} \\
\hat{\beta}_{2} \\
\cdot \\
\cdot \\
\hat{\beta}_{p}
\end{array}\right|
$$

from which the invariance property follows.

In view of the generalization of theorem 3.2 as given in theorem 3.4, It is to be expected that theorem 3.3 may also be extended. The generalization is given in theorem 3.5. Consider the following conditions:
(a) $\frac{\hat{a}_{j}}{\hat{a}_{i}}=\frac{\hat{\beta}_{j}}{\hat{\beta}_{i}} \quad$ i, $j=1,2, \ldots p$
(b) $\left.\quad X_{2}=\left|\begin{array}{l}1 \\ 1 \\ \cdot \\ \cdot \\ 1\end{array}\right| \begin{array}{llll}h_{1} \hat{\beta}_{1} & h_{1} \hat{\beta}_{2} & \cdots & h_{1} \hat{\beta} p \\ h_{2} \hat{\beta}_{1} & h_{2} \hat{\beta}_{2} & \cdots & h_{2} \hat{\beta}_{p} \\ \cdot & \cdot & & \\ \cdot & \cdot & & \\ h_{N_{2}} \hat{\beta}_{1}{ }^{h_{N_{2}}} \hat{N}_{2} & \cdots & h_{N_{2}} \hat{\beta}\end{array} \right\rvert\,$
(c) $\quad\left(X_{1}^{t} X_{1}\right)^{-1}=\left(\begin{array}{ll}u & 0 \\ 0 & v I_{p}\end{array}\right)$.

These conditions are the same as those preceding theorem 3.3 except that condition (b) implies that there are $N_{2}$ additional points in the direction of steepest ascent instead of just one additional point.

Theorem 3.5. Any two of the above three conditions imply
the third.

Proof: $X_{2}, \hat{a}$, and $\hat{\beta}$ may be partitioned as in the orem 3.4. $\left(X_{l}^{\prime} X_{l}\right)^{-1}$ may be partitioned as

$$
\left(X_{L}^{\prime} X_{1}\right)^{-1}=\left(\left.\begin{array}{c|cccc}
u & 0 & , \ldots & \cdot & 0 \\
\hline 0 & v & \cdot & \cdot & \cdot \\
\cdot & \cdot & & & 0 \\
\cdot & \cdot & & & \\
\cdot & \cdot & & & \\
0 & 0 & . & \cdot & \cdot
\end{array} \right\rvert\,=\left(\begin{array}{ll}
C_{00} & 0 \\
0 & C_{11}
\end{array}\right)\right.
$$

The proof that (b) and (c) imply (a) is the content of theorem 3.4 with $\frac{1}{\mathrm{~N}_{1}}$ replaced by $u$ and $v$ as indicated in the above partitioning of $\left(X_{1}^{\prime} X_{1}\right)^{-1}$.

To prove that (a) and (c) imply (b) the equation $\hat{a}=\hat{\beta}+J^{\prime} G Q$ is written as
$\binom{A_{0}}{k B_{1}}=\binom{B_{0}}{B_{1}}+\left(\begin{array}{cc}u & 0 \\ 0 & v I_{p}\end{array}\right)\left(\begin{array}{cc}W & 0^{\prime} \\ \vdots_{21}^{\prime}\end{array}\right) \quad G Q$
(3.45)
where $A_{l}$ has been replaced by $k B_{l}$ according to (a). Equating corresponding elements of (3.45) there results

$$
\begin{equation*}
A_{0}=B_{0}+u W_{0}^{\prime} G Q \tag{3.46}
\end{equation*}
$$

and

$$
\begin{equation*}
k B_{1}=B_{1}+v I_{p} X_{21}^{\prime} G Q \tag{3.47}
\end{equation*}
$$

Equation (3.47) may be written as

$$
\begin{equation*}
\left(\frac{k-1}{v}\right) B_{1}=X_{21}^{\prime} G Q \tag{3.48}
\end{equation*}
$$

or

$$
\begin{equation*}
M X_{21}=\ell B_{1}^{\prime} \tag{3.49}
\end{equation*}
$$

where $M=(G Q)^{\prime}$ has dimension $1 \times N_{2}$ and $\ell$ is the scalar
$\frac{k-1}{v}$. Evidently, the elements $x_{i j}$ of $X_{21}$ are of the form $x_{i, j}=a_{i j} \hat{\beta}_{i}$. Let $M=\left(m_{1}, m_{2}, \ldots m_{N_{2}}\right)$. Then (3.49)
may be written as

Now (3.50) must be true for all M and in particular must hold as $M$ ranges over the unit vectors

$$
\begin{aligned}
& u_{1}=(1,0, \ldots 0) \\
& u_{2}=(0,1, \ldots 0)
\end{aligned}
$$

and

$$
{ }^{\mathrm{u}} \mathrm{~N}_{2}=(0,0, \ldots 1)
$$

When $M=u_{1}$, equating the corresponding elements of (3.50) yields

$$
\begin{aligned}
& a_{11} \hat{\beta}_{1}=\ell \hat{\beta}_{1} \\
& a_{12} \hat{\beta}_{2}=\ell \hat{\beta}_{\cdot 2} \\
& \cdot \\
& a_{1 p} \hat{\beta}_{p}=\ell \hat{\beta}_{p}
\end{aligned}
$$

so that $a_{11}=a_{12}=. \quad=a_{1 p}$

More generally, when $M=u_{i}$,

$$
\begin{aligned}
a_{i 1} \hat{\beta}_{1} & =\ell \hat{\beta}_{1} \\
a_{i 2} \hat{\beta}_{2} & =\ell \hat{\beta}_{2} \\
& \cdot \\
& \cdot \\
a_{i p} \hat{\beta}_{p} & =\ell \hat{\beta}_{p}
\end{aligned}
$$

so that $a_{i 1}=a_{i 2}=. \quad . a_{i p}$.
Since this is true for $i=1,2, \ldots . N_{2}$, the $a_{i j}$ in any row are alike. Hence $X_{21}$ has the form

$$
\left(\begin{array} { c c c c c } 
{ a _ { 1 } \hat { \beta } _ { 1 } } & { a _ { 1 } \hat { \beta } _ { 2 } } & { \cdot } & { \cdot } & { \cdot } \\
{ a _ { 1 } \hat { \beta } _ { p } } \\
{ a _ { 2 } \hat { \beta } _ { 1 } } & { a _ { 2 } \hat { \beta } _ { 2 } } & { \cdot } & { \cdot } & { \cdot } \\
{ \cdot } & { a _ { 2 } \hat { \beta } _ { p } } \\
{ \cdot } & { \cdot } & { } & { } & { } \\
{ \cdot } & { \cdot } & { } & { } \\
{ a _ { N _ { 2 } } \hat { \beta } _ { 1 } } & { a _ { N _ { 2 } } \hat { \beta } _ { 2 } } & { \cdot } & { \cdot } & { \cdot } \\
{ a _ { N _ { 2 } } \hat { \beta } _ { p } }
\end{array} \left|=\left|\begin{array}{c}
a_{1} B_{1}^{\prime} \\
a_{2} B_{1}^{\prime} \\
\\
a_{N_{2}} B_{1}^{\prime}
\end{array}\right|=A B_{1}^{\prime}\right.\right.
$$

which is of the form of (b).
In order to prove that (a) and (b) imply (c), the symmetric matrix $S_{1}^{-1}$ is partitioned as

$$
\left|\begin{array}{c|ccc}
C_{00} & C_{01} & \cdots & C_{0 p} \\
\hline C_{10} & C_{11} & \cdots & C_{1 p} \\
\cdot & \cdot & & \\
\cdot & \cdot & & \\
\cdot & \cdot & & \\
C_{p 0} & C_{p l} & \cdots & C_{p p}
\end{array}\right|=\left(\begin{array}{cc}
C_{00} & C_{01} \\
C_{10} & C_{11}
\end{array}\right)
$$

$$
\begin{aligned}
& \text { Now } J=\left(W_{0}, H B_{1}^{\prime}\right)\left(\begin{array}{ll}
C_{00} & C_{01} \\
C_{10} & C_{11}
\end{array}\right) \text { so that } \\
& J^{\prime}=\left(\begin{array}{ll}
C_{00} & C_{01} \\
C_{10} & C_{11}
\end{array}\right) \cdot\binom{W_{0}^{\prime}}{B_{1} H^{\prime}} \cdot \text { Hence }
\end{aligned}
$$

$$
\begin{equation*}
\binom{A_{0}}{A_{1}}=\binom{B_{0}}{B_{1}}+\binom{C_{00} W_{0}^{\prime}+C_{01} B_{1} H^{\prime}}{C_{10} W_{0}^{\prime}+C_{11} B_{1} H^{\prime}}, \quad G Q \tag{3.51}
\end{equation*}
$$

It follows that

$$
\begin{equation*}
A_{0}=B_{0}+C_{00} W_{0}^{\prime} G Q+C_{01} B_{1} H^{\prime} G Q \tag{3.52}
\end{equation*}
$$

and that

$$
\begin{equation*}
A_{1}=B_{1}+C_{10} W_{0}^{\prime} G Q+C_{11} B_{1} H^{\prime} G Q \tag{3.53}
\end{equation*}
$$

Now by (a), $A_{1}=k B_{1}$ so that $(3.53)$ becomes

$$
\begin{equation*}
k B_{1}=B_{1}+C_{10} W_{0}^{\prime} G Q+C_{11} B_{1} H^{\prime} G Q \tag{3.54}
\end{equation*}
$$

Since ( 3.54 ) must be true for all vectors $B_{1}, C_{10}=0$ so that ( 3.54 ) becomes

$$
\begin{equation*}
\mathrm{kB}_{1}=\mathrm{B}_{1}+\mathrm{C}_{11} \mathrm{~B}_{1} \mathrm{H}^{\prime} \mathrm{GQ} . \tag{3.55}
\end{equation*}
$$

Now H'GQ is a scalar so that ( 3.55 ) may be written

$$
\left(k I-H^{\prime} G Q C_{1 I}\right) B_{1}=B_{1} \text { which }
$$

implies that

$$
\mathrm{kI}-\mathrm{H}^{\prime} \mathrm{GQC}_{11}=\mathrm{I}
$$

It follows that

$$
C_{11}=\frac{(k-1)}{H^{\prime} G Q} \quad I \text { so that } C_{11}
$$

is a scalar matrix as in (c).
If the invariance property is extended to include $\hat{a}_{0}$ and $\hat{\beta}_{0}$ then $H$ must be chosen so that $B_{0} H=W_{0}$. Then since $C_{01}=0$, equation (3.52) becomes

$$
\begin{aligned}
\mathrm{kB}_{0} & =B_{0}+C_{00} W_{0}^{\prime} G Q \text {. Hence } \\
C_{00} & =\frac{(k-1) B_{0}}{W_{0}^{\prime} G Q}=\frac{(k-1) B_{0}}{Q^{\prime} G^{\prime} W_{0}}=\frac{(k-1) B_{0}}{Q^{\prime} G^{\prime} B_{0} H} \\
& =\frac{(k-1)}{H^{\prime} G Q} \quad \text { which is the scalar multiplier of } I \text { in the }
\end{aligned}
$$ expression for $C_{11}$ above.

## The Predicted Response

Thus far the invariance property has been investigated but no expression for the response at an additional point has been given. Suppose that a p-dimensional unit grid has been employed to estimate $\hat{\beta}$ and that in addition to the $N_{1}=2^{\mathrm{P}}$ grid points, an additional point $x=\left(x_{1}, x_{2}, \ldots, x_{p}\right)$ is used. It is convenient then to distinguish between:
(a). $y$, the observed response at $x$,
(b). $\hat{y}_{1}$, the least squares prediction of $\eta$ at $x$ on the basis of the $N_{1}$ points, and
(c). $\hat{y}_{2}$, the least squares prediction of $\eta$ at $x$ on the basis of the $N_{L}+1$ points.

Clearly, $\hat{y}_{1}=\hat{X \beta}$ and $\hat{y}_{2}=\hat{X a}$. Now suppose that the additional point is in the steepest ascent direction so that in matrix notation, $X_{2}=\left(1, h \hat{\beta}_{1}, \ldots, h \hat{\beta}_{p}\right)=\left(1, h B_{1}^{\prime}\right)$. Now

$$
\begin{align*}
\hat{y}_{2} & =X_{2} \hat{a}=X_{2}\left(\hat{\beta}+J^{\prime} G Q\right) \\
& =X_{2} \hat{\beta}+x_{2} S_{1}^{-1} X_{2}^{\prime} G\left(y-X_{2} \hat{\beta}\right) \\
& =\hat{y}_{1}+R(I+R)^{-1}\left(y-\hat{y}_{1}\right) \tag{3.56}
\end{align*}
$$

Now $R$ is a positive scalar written

$$
R=X_{2} S_{1}^{-1} X_{2}^{\prime}=\left(1, h B_{1}^{\prime}\right)\left(\begin{array}{cc}
\frac{1}{N_{1}} & 0 \\
0 & \frac{I_{p}}{N_{1}}
\end{array}\right) \quad\binom{1}{h B_{1}}=\frac{1}{N_{1}}\left(1+h_{1}^{2} B_{1}^{\prime} B_{1}\right)
$$

Denote the quantity $R(I+R)^{-1}=\frac{R}{1+R}$ by $k^{2}$ so that $0<k^{2}<1$. Then by $(3,56)$,

$$
\begin{equation*}
\hat{y}_{2}=\hat{y}_{1}+k^{2}\left(y-\hat{y}_{1}\right) \tag{3.57}
\end{equation*}
$$

so that

$$
\begin{equation*}
\hat{y}_{2}-\hat{y}_{1}=k^{2}\left(y-\hat{y}_{1}\right) \text { and } 0<\frac{\hat{\mathrm{y}}_{2}-\hat{\mathrm{y}}_{1}}{\mathrm{y}-\hat{\mathrm{y}}_{1}}<1 \tag{3.58}
\end{equation*}
$$

It is to be noted that $R=X_{2} S^{-1} X_{2}^{\prime}$ is simply the sum of squares of the elements of the additional point divided by the number of points in the initial design. The result (3.58) is more general than has been indicated since it is valid whenever the design matrix has the form $X_{1}^{\prime} X_{1}=p I$. Further, the additional point is not restricted to the steepest ascent line but may be arbitrary. In these circumstances
$R=X_{2}\left(X_{1}^{\prime} X_{1}\right)^{-1} X_{2}^{\prime}=\frac{X_{2} X_{2}^{\prime}}{p} \quad$ which is a positive scalar.
Then (3.56) becomes
$\hat{y}_{2}=\hat{y}_{1}+\frac{\frac{x_{2} x_{2}}{p}}{\frac{\dot{p}+x_{2} x_{2}^{\prime}}{}}\left(y-\hat{y}_{1}\right)=\hat{y}_{1}+\frac{x_{2} x_{2}^{\prime}}{p+x_{2}^{\prime} x_{2}^{\prime}}\left(y-\hat{y}_{1}\right)$.

Hence (3.58) is valid in the general case.
The expression ( 3.58 ) symbolizes a condition which intuition could have foretold. It states that the observation $y$ and the prediction on the basis of $N_{1}+1$ points, $\hat{y}_{2}$, either are both greater than the prediction on the basis of $N_{l}$ points, $\hat{y}_{l}$, or they are both less than $\hat{y}_{1}$. It also implies that $\hat{y}_{2}$ is always between $y$ and $\hat{y}_{1}$.

These facts can be more elegantly stated if (3.57) is written as

$$
\begin{equation*}
\hat{\mathrm{y}}_{2}=\mathrm{k}^{2} \mathrm{y}+\left(1-\mathrm{k}^{2}\right) \hat{\mathrm{y}}_{1} \tag{3.60}
\end{equation*}
$$

from which $\hat{\mathrm{y}}_{2}$ may be recognized to be a convex combination: of $y$ and $\hat{y}_{l}$. In addition, since $k^{2}=\frac{\mathrm{X}_{2}^{\prime} \mathrm{X}_{2}}{\mathrm{~N}_{\mathrm{l}}}$, the further along the steepest ascent path the additional observation is taken, the greater is the weight placed upon it. Figure 1 indicates the relative positions of $y, \hat{y}_{1}$, and $\hat{y}_{2}$ in a two dimensional factor space.


Figure 1. Relative Magnitude of Observed and Predicted Responses

A judicious choice of the location of the first experiment on the steepest ascent line is a problem to which some space will be given subsequently. Apparently the step size along the steepest ascent line is something about which the experimenter is supposed to have some feeling. To be more concrete, suppose the experimenter wishes to take a step of such a size that the predicted yield will increase by an amount $\Pi$.

Now at the origin, $x=(0, \ldots, 0)$ so that $x=(1,0, \ldots, 0)$ and

$$
\begin{equation*}
\hat{y}=X \hat{\beta}=\hat{\beta}_{0} \tag{3.61}
\end{equation*}
$$

It is desired to determine $h$ so that at $x=\left(h \hat{\beta}_{1}, \ldots, h \hat{\beta}_{p}\right)$ the response has increased by an amount $I I$, that is

$$
\begin{equation*}
\hat{y}=\left(1, h B_{1}^{\prime}\right)\binom{B_{0}}{B_{1}}=\pi+\hat{\beta}_{0} . \tag{3.62}
\end{equation*}
$$

Then $\hat{\beta}_{0}+h B_{1}^{\prime} B_{1}=\Pi+\hat{\beta}_{0}$ so that

$$
\begin{equation*}
h_{2}=\frac{\Pi}{B_{1}^{\prime} B_{1}} \tag{3.63}
\end{equation*}
$$

Hence an increase in response by the amount $I$ is predicted at the point

$$
\begin{align*}
x= & \frac{\pi}{B_{1}{ }^{i} B_{1}}\left(\hat{\beta}_{1}, \ldots, \hat{\beta}_{p}\right)=\frac{\Pi B_{1}^{\prime}}{B_{1}^{\prime} B_{1}}  \tag{3.64}\\
& \text { Upper and Lower Bounds for } \frac{\beta_{11}}{\beta_{22}}
\end{align*}
$$

Usually the experimenter has some prior information regarding the response under consideration. He is quite likely to know whether the response has a maximum or a minimum in the region of interest and he may even know the shape of the factor space contours. The following discussion indicates a position along the steepest ascent line where the first experiment may be run in order to utilize prior information.

Consider the response
$y=\beta_{0}+\beta_{1} x_{1}+\beta_{2} x_{2}+\beta_{11} x_{1}^{2}+\beta_{12} x_{1} x_{2}+\beta_{22^{x}}{ }^{2}+e$
and the least squares prediction

$$
\begin{equation*}
\hat{y}=\hat{\beta}_{0}+\hat{\beta}_{1} x_{1}+\hat{\beta}_{2} x_{2}+\hat{\beta}_{12} x_{1} x_{2} \tag{3.65}
\end{equation*}
$$

where the $\beta^{\prime}$ s are estimated by means of a unit grid. It may
readily be verified that

$$
\begin{aligned}
& E\left(\hat{\beta}_{0}\right)=\beta_{0}+\beta_{11}+\beta_{22} \\
& E\left(\hat{\beta}_{1}\right)=\beta_{1} \\
& E\left(\hat{\beta}_{2}\right)=\beta_{2} \\
& E\left(\hat{\beta}_{12}\right)=\beta_{12}
\end{aligned}
$$

so that

$$
\begin{align*}
E(y-\hat{y})= & E\left[\left(\beta_{0}+\beta_{1} x_{1}+\beta_{2} x_{2}+\beta_{11} x_{1}^{2}+\beta_{12} x_{1} x_{2}+\beta_{22} x_{2}^{2}+e\right)\right. \\
& \left.-\left(\beta_{0}+\beta_{11}+\beta_{22}+\beta_{1} x_{1}+\beta_{2} x_{2}+\beta_{12} x_{1} x_{2}\right)\right] \\
= & \beta_{11} x_{1}^{2}+\beta_{22} x_{2}^{2}-\beta_{11}-\beta_{22} . \tag{3.66}
\end{align*}
$$

Hence the difference between the observed response and the predicted response at a point $\left(x_{1}, x_{2}\right)$ is an unbiased estimate of
$\beta_{11} x_{1}^{2}+\beta_{22} x_{2}^{2}-\beta_{11}-\beta_{22}$.
The gradient line at the origin has the equation $x_{2}=\beta_{2} / \beta_{1} x_{1}$
so that a point on the gradient line may be written as $\left(h \beta_{1}, h \beta_{2}\right)$. Then the difference $y-\hat{y}$ at the point $\left(h \beta_{1}, h \beta_{2}\right)$ is an estimate


That is,

$$
\begin{equation*}
\left.E(y-\hat{y})\right|_{\left(h \beta_{1}, h \beta_{2}\right)}=\beta_{11} h^{2} \beta_{1}^{2}+\beta_{22^{h} \beta_{2}^{2}-\beta_{11}-\beta_{22} .} \tag{3.67}
\end{equation*}
$$

The quantities $\beta_{1}$ and $\beta_{2}$ have been estimated by $\hat{\beta}_{q}$ and $\hat{\beta}_{2}$.
$\beta_{11}$ and $\beta_{22}$ are unknown while the choice of $h$ remains with the investigator: Suppose $h$ is chosen so that $h^{2} \hat{\beta}_{1}^{2}-1>0$ and $h^{2} \hat{\beta}_{2}^{2}-1<0$. Also let the observed response at ( $\mathrm{h} \hat{\beta}_{1}, \mathrm{~h} \hat{\beta}_{2}$ ) be larger than that predicted there so that $\mathrm{y}-\hat{\mathrm{y}}>0$. Then

$$
\begin{equation*}
\beta_{11}\left(h^{2} \tilde{\beta}_{1}^{2}-1\right)+\beta_{22}\left(h^{2} \hat{\beta}_{2}^{2}-1\right)>0 . \tag{3.68}
\end{equation*}
$$

Further, if $\beta_{22}>0,(3.68)$ may be written

$$
\begin{equation*}
\frac{\beta_{11}}{\beta_{22}}>-\frac{h^{2} \hat{\beta}_{2}^{2}-1}{h^{2} \hat{\beta}_{1}^{2}-1} \tag{3.69}
\end{equation*}
$$

whereas if $\beta_{22}<0$,

$$
\begin{equation*}
\frac{\beta_{11}}{\beta_{22}}<-\frac{h^{2} \hat{\beta}_{2}^{2}-1}{h^{2} \hat{\beta}_{1}^{2}-1} \tag{3,70}
\end{equation*}
$$

Now the right side of (3.69) and (3.70) is positive so that in the first case it serves as a lower bound to the ratio $\beta_{11} / \beta_{22}$ while in the second case it serves as an upper bound to the ratio $\beta_{11} / \beta_{22}$.

Now consider the combination $h^{2} \hat{\beta}_{1}^{2}-1>0, h^{2} \hat{\beta}_{2}^{2}-1>0$, $\beta_{11}<0$, and $\beta_{22}<0$. Then $(3.68)$ leads to

$$
\begin{equation*}
\frac{\beta_{11}}{\beta_{22}}<-\frac{h^{2} \hat{\beta}_{2}^{2}-1}{h^{2} \hat{\beta}_{1}^{2}-1} \tag{3.71}
\end{equation*}
$$

Now the right side of (3.71) is negative while the left side is positive
so that evidently this combination of signs cannot occur.

$$
\begin{gather*}
\text { Now take }\left(h^{2} \hat{\beta}_{1}^{2}-1\right)>0,\left(h^{2} \hat{\beta}_{2}^{2}-1\right)>0, \\
\beta_{11}>0 \text {, and } \beta_{22}>0 \text {. Then }(3.68) \text { leads to } \\
\frac{\beta_{11}}{\beta_{22}}>-\frac{h^{2} \hat{\beta}_{2}^{2}-1}{h^{2} \hat{\beta}_{1}^{2}-1} . \tag{3,72}
\end{gather*}
$$

Since it is known at the outset that $\frac{\beta 11}{\beta_{22}}$ is positive and
since the right side is negative, (3.72) yields no new information.

by $t$, Let $U(t)$ denote an
upper bound of $t$ and $U(-t)$ denote an upper bound of $-t$. Let the symbol $L(\underset{\sim}{f})$ be used in a similar fashion. Also let $N$ indicate that no new information is available and let $I$ indicate that a particular combination of signs cannot occur. With this notation, table II summarizes all possible sign combinations for $y>\hat{y}$. When $y<\hat{y}$ it is merely necessary to interchange $N$ with $I$ and $L$ with $U$.

A knowledge of the ratio $\beta_{11} / \beta_{22}$ can be useful in some instances. If interaction is present the axes of the contour system will be rotated through an angle $\theta$ where $\tan \boldsymbol{Z} \theta=\frac{\beta_{12}}{\beta_{11}-\beta_{22}}$. If $\beta_{11} / \beta_{22}$ is close to $L$, then $\theta$ is close to $45^{\circ}$ or $135^{\circ}$
depending on the sign of $\beta_{12}$. If the experimenter feels that no interaction is present then $\beta_{12}=0$. In this case it can be shown that the eccentricity, $\rho$, of the elliptical contours is given by $\rho=\sqrt{\left|d^{2}-1\right|}$ where $d=\beta_{\text {Il }} / \beta_{22}$. Hence a measure of the degree of elongation of the contour system is available.

TABLE II UPPER AND LOWER BOUNDS FOR $\beta_{11} / \beta_{22}$

| $y-\hat{y}$ | $h^{2} \hat{\beta}_{1}^{2}-1$ | $h^{2} \hat{\beta}_{2}{ }^{2}-1$ | $\beta_{11}$ | ${ }^{\beta} 22$ | $\beta_{11 /} \beta_{22}$ |
| :---: | :---: | :---: | :---: | :---: | :---: |
| $+$ | $+$ | $+$ | + | + | N |
| $+$ | $+$ | $t$ | $\pm$ | - | $U(-t)$ |
| $+$ | + | + | - | + | L ( $-t$ ) |
| $+$ | $\dagger$ | + | - | - | I |
| $t$ | $t$ | - | $+$ | + | L (t) |
| $t$ | +- | - | + | - | N |
| $+$ | $+$ | - | - | $\pm$ | I |
| + | + | - | - | - | U ( t ) |
| $t$ | - | $+$ | $+$ | $t$ | U (it ) |
| $t$ | - | + | + | - | I |
| $t$ | - | $t$ | - | $t$ | N |
| $t$ | - | $t$ | - | - | $L$ (t) |
| $t$ | - | - | + | $t$ | I |
| $+$ | - | - | $+$ | - | $L(-t)$ |
| $t$ | - | - | - | $t$ | $U(-t)$ |
| $t$ | - | - | - | - | N |

In the case where $\beta_{12}=0$, equation $(3.74)$ gives the equation of the line from the origin to the true maximum as

$$
x_{2}=\frac{\beta_{2}}{\beta_{1}} \frac{\beta_{11}}{\beta_{22}} x_{1} .
$$

An estimate of the ratio $\beta_{11} / \beta_{22}$ can be used to adjust the slope $\hat{\beta}_{2} / \hat{\beta}_{1}$ of the path of steepest ascent.

## Gradient Properties

$$
\begin{align*}
& \text { Consider the response } y=\eta+e \text { where } \\
& \eta=\beta_{0}+\beta_{1} x_{1}+\beta_{2} x_{2}+\beta_{11} x_{1}^{2}+\beta_{12} x_{1} x_{2}+\beta_{22^{x}}{ }^{2} \tag{3.73}
\end{align*}
$$

Now

$$
\frac{\partial \eta}{\partial x_{1}}=\beta_{1}+2 \beta_{11} x_{1}+\beta_{12} x_{2}
$$

and

$$
\frac{\partial \eta}{\partial x_{2}}=\beta_{2}+\beta_{12} x_{1}+2 \beta_{22} x_{2} .
$$

If these are set equal to 0 and solved for $x_{1}$ and $x_{2}$, there results

$$
x_{1}=\frac{2 \beta_{1} \beta_{22}-\beta_{2} \beta_{12}}{\beta_{12}^{2}-4 \beta_{11} \beta_{22}}
$$

and

$$
\begin{equation*}
x_{2}=\frac{2 \beta_{2} \beta_{11}-\beta_{1} \beta_{12}}{\beta_{12}^{2}-4 \beta_{11} \beta_{22}} \tag{3.74}
\end{equation*}
$$

which are the coordinates of the maximum of $\eta$. Since the coordinates of the maximum will often be referred to, let it be agreed to denote them by $\left(x_{1}^{*}, x_{2}^{*}\right)$.

The path in the factor space which leads to the coordinates of the true maximum is along the line

$$
x_{2}=\frac{x_{2} *}{x_{1} *} \quad x_{1}, \text { provided } \beta_{12}^{2}-4 \beta_{11} \beta 2_{22} \neq 0
$$

The gradient path is normal to the factor space contours at the origin and has the equation $x_{2}=\beta_{2} / \beta_{1} x_{1}$. Unfortunately, only in very special instances will the gradient pass through the coordinates of the maximum. Evidently these paths will coincide whenever

$$
\begin{equation*}
\frac{\beta_{2}}{\beta_{1}}=\frac{2 \beta_{2} \beta_{11}-\beta_{1} \beta_{12}}{2 \beta_{1} \beta_{22}-\beta_{2} \beta_{12}}, \tag{3.75}
\end{equation*}
$$

which reduces to

$$
\begin{equation*}
2 \beta_{1} \beta_{2}\left(\beta_{22}-\beta_{11}\right)+\beta_{12}\left(\beta_{1}^{2}-\beta_{2}^{2}\right)=0 . \tag{3.76}
\end{equation*}
$$

This is satisfied whenever

$$
\begin{aligned}
& \left(\text { a) } \beta_{22}=\beta_{11} \text { and } \beta_{1}=\beta_{2},\right. \\
& (\mathrm{b}) \beta_{22}=\beta_{11} \text { and } \beta_{12}=0, \text { or } \\
& (\mathrm{c}) \beta_{2}=0 \text { and } \beta_{12}=0 .
\end{aligned}
$$

In each of these cases, it is geometrically evident that the gradient line at the origin passes through the center of the contour system. In (a), the line $x_{2}= \pm x_{1}$ passesthrough the center of the system and the line $x_{2}=x_{1}$ is tangent at the origin. In case ( $b$ ), the contour system is circular while in case (c), the center is on the $x_{1}$ axis.

In general, the factor space contour system will not be one of the special cases discussed in the previous paragraph. The experimenter would indeed be extremely fortunate if the contour system were such that the maximum on the line $x_{2}=\frac{\beta_{2}}{\beta_{1}} x_{1}$ was close to the true maximum. In accordance with the steepest ascent technique, experimentation may be continued at the points of a second grid, the center of which is located at the maximum of the first steepest ascent path. Further experimentation is then performed along the line determined from the second grid. It is relevant therefore to find an expression for the maximum response attainable on the gradient and for the coordinates of the point where this maximum occurs. It is also relevant to find the equation of the gradient at this maximum point.

Theorem 3.6. The maximum value of $\eta$ subject to $x_{2}=\mathrm{mx}_{1}$ occurs at the point

$$
\begin{align*}
& x_{1}=-\frac{a}{2 b} \\
& x_{2}=-\frac{m a}{2 b} \tag{3.77}
\end{align*}
$$

and is given by

$$
\begin{equation*}
\eta=\beta_{0}-\frac{a^{2}}{2 b} \tag{3.78}
\end{equation*}
$$

where

$$
a=\beta_{1}+m \beta_{2} \text { and } b=\beta_{11}+m \beta_{12}+m^{2} \beta_{22} \text { (3.79) }
$$

Proof: On substituting $\mathrm{mx}_{1}$ for $\mathrm{x}_{2}$ in (3.73), there results

$$
\eta=\beta_{0}+\beta_{1} x_{1}+\beta_{2} m_{1}+\beta_{11} x_{1}^{2}+\beta_{12} x_{1}^{2} m+\beta_{22^{m}}{ }^{2} x_{1}^{2}
$$

$$
=\beta_{0}+\left(\beta_{-1}+m \beta_{2}\right) x_{1}+\left(\beta_{11}+m \beta_{12}+m^{2} \beta_{22}\right) x_{1}^{2} .(3.80)
$$

Then finding $\frac{d \eta}{d x_{1}}$ and setting it equal to 0 , there results

$$
\frac{d \eta}{d x_{1}}=\beta_{1}+m \beta_{2}+2\left(\beta_{11}+m \beta_{12}+\beta_{22^{m}}{ }^{2}\right) x_{1}=0
$$

so that

$$
\begin{aligned}
& x_{1}=-\frac{\beta_{1}+m \beta_{2}}{2\left(\beta_{11}+m \beta_{12}+m^{2} \beta_{22}\right.} \\
& x_{2}=-\frac{m\left(\beta_{1}+m \beta_{2}\right)}{2 \beta_{11}+m \beta_{12}+m^{2} \beta_{22}}
\end{aligned}
$$

thereby verifying (3.77): Using (3.80),

$$
\begin{equation*}
\eta=\beta_{0}-\frac{a^{2}}{2 b}+\frac{a^{2}}{4 b^{2}} \cdot b=\beta_{0}-\frac{a^{2}}{2 b} \tag{3.81}
\end{equation*}
$$

thereby establishing (3.78).

In particular, the line $x_{2}=m x_{1}$ is the gradient line when $m=\frac{\beta_{2}}{\beta_{1}}$. In this case the maximum response occurs at the point

$$
\begin{array}{r}
x_{1}=-\frac{\beta_{1}\left(\beta_{1}^{2}+\beta_{2}^{2}\right)}{2\left(\beta_{1}^{2} \beta_{11}+\beta_{1} \beta_{2} \beta_{12}+\beta_{2}^{2} \beta_{22}\right.} \\
x_{2}=-\frac{\beta_{2}\left(\beta_{1}^{2}+\beta_{2}^{2}\right)}{2\left(\beta_{1}^{2} \beta_{11}+\beta_{1} \beta_{2}^{\beta}{ }_{12}+\beta_{2}^{2} \beta_{22}\right.} \tag{3.82}
\end{array}
$$

and is given by

$$
\begin{equation*}
\eta=\beta_{0}-\frac{\left(\beta_{1}^{2}+\beta_{2}^{2}\right)^{2}}{2\left(\beta_{1}^{2} \beta_{11}+\beta_{1} \beta_{2} \beta_{12}+\beta_{2}^{2} \beta_{22}\right)} \tag{3.83}
\end{equation*}
$$

Theorem 3.7. The gradient line at the point on the line $x_{2}=m x_{1}$ where $\eta$ is maximized, is perpendicular to the line $\mathrm{x}_{2}=\mathrm{mx}_{1}$.

Proof: The maximum point may be shifted to the origin by the translation

$$
\begin{align*}
& z_{1}=x_{1}-h \\
& z_{2}=x_{2}-k \tag{3.84}
\end{align*}
$$

where

$$
\begin{equation*}
h=-\frac{a}{2 b} \quad \text { and } \quad k=-\frac{m a}{2 b} \tag{3.85}
\end{equation*}
$$

are the coordinates of the maximum point as given in theorem 3.6. Then equation (3.73) transforms to

$$
\begin{equation*}
\eta=\beta_{0}^{\prime}+\beta_{1}^{\prime z_{1}}+\beta_{2}^{\prime z_{2}}+\beta_{11} z_{1}^{2}+\beta_{12} z_{1} z_{2}+\beta_{22_{2}}^{2} \tag{3.86}
\end{equation*}
$$

where $\beta_{1}^{\prime}$ and $\beta_{2}^{\prime}$ are given by

$$
\begin{align*}
& \beta_{1}^{\prime}=\beta_{1}+2 \beta_{11} h+\beta_{12} k \\
& \beta_{2}^{\prime}=\beta_{2}+\beta_{12} 2^{h}+2 \beta_{22} k \tag{3.87}
\end{align*}
$$

Hence the gradient line at the origin in the $z$ space is

$$
\begin{equation*}
z_{2}=\frac{\beta_{2}+\beta_{12^{h}+2 \beta_{22}}}{\beta_{1}+2 \beta_{11}^{h}+\beta_{12^{k}}} z_{1} \tag{3.88}
\end{equation*}
$$

Using (3.85) and (3.79), the numerator of this equation may be written

$$
\begin{gathered}
\beta_{2}+\beta_{12} h+2 \beta_{22} k=\beta_{2}-\frac{\beta_{12}}{2 b}-\frac{2 \beta_{22} m a}{2 b} \\
=\frac{2 \beta_{2} b-a \beta_{12}-2 m a \beta_{22}}{2 b} \\
=\frac{2 \beta_{2}\left(\beta_{11}+m \beta_{12}+\beta_{22^{m}}\right)-\left(\beta_{1}+m \beta_{2}\right) \beta_{12}-2 m\left(\beta_{1}+m \beta_{2}\right) \beta_{22}}{2 b} \\
=\frac{\left(2 \beta_{2} \beta_{11}+m \beta_{2} \beta_{12}-\beta_{1} \beta_{12}-2 m \beta_{22} \beta_{1}\right)}{2 b}
\end{gathered}
$$

The denominator of ( 3.88 ) may be written

$$
\begin{gathered}
\beta_{1}-\frac{2 \beta_{11} a}{2 b}-\frac{\beta_{12} m a}{2 b}=\frac{2 \beta_{1} b-2 \beta_{11} a-\beta_{12} \mathrm{ma}}{2 b} \\
=\frac{2 \beta_{1}\left(\beta_{11}+m \beta_{12}+\beta_{22^{m}}{ }^{2}\right)-2 \beta_{11}\left(\beta_{1}+m \beta_{2}\right)-m \beta_{12}\left(\beta_{1}+m \beta_{2}\right)}{2 b} \\
=\frac{m\left(\beta_{1} \beta_{12}+2 \beta_{1} \beta_{22} m-2 \beta_{2} \beta_{11}-m \beta_{2} \beta_{12}\right)}{2 b} .
\end{gathered}
$$

Therefore equation (3.88) becomes
$z_{2}=\frac{2 \beta_{2} \beta_{11}+m \beta_{2} \beta_{12}-\beta_{1} \beta_{12}-2 m \beta_{22} \beta_{1}}{-m\left(2 \beta_{2} \beta_{11}+m \beta_{2} \beta_{12}-\beta_{1} \beta_{12}-2 m \beta_{22} \beta_{1}\right)} z_{1}$
and provided the denominator is not zero,

$$
\begin{equation*}
z_{2}=-\frac{1}{m} z_{1} . \tag{3.90}
\end{equation*}
$$

The condition that the denominator is not zero implies that $m \neq 0$ and $\mathrm{m} \neq \frac{\beta_{1}{ }_{1} 1_{2}-2 \beta_{2} \beta_{11}}{\beta_{2}{ }_{12}-2 \beta_{1} \beta_{22}}$. As $\mathrm{m} \rightarrow 0, \frac{1}{\mathrm{~m}} \rightarrow \infty$, so that the gradient lines approach perpendicularity. As may be seen by (3.75), if $m=\frac{\beta_{1} \beta_{12}-2 \beta_{2} \beta_{11}}{\beta_{2} \beta_{12}-2 \beta_{1} \beta_{22}}$, then the line $x_{2}=m x_{1}$ goes through
the true maximum. Substituting $x_{1}$ and $x_{2}$ as given by (3.84) in (3.90), there results

$$
\begin{equation*}
x_{2}+\frac{m a}{2 b}=-\frac{1}{m}\left(x_{1}+\frac{a}{2 b}\right) \tag{3.91}
\end{equation*}
$$

Since the slope of this line and the line $x_{2}=m x_{1}$ are negative reciprocals, the lines are perpendicular.

In particular, theorem 3.7 applies when $m=\frac{\beta_{2}}{\beta_{1}}$, so that the second gradient line is perpendicular to the first.

An examination of equation (3.75) has shown that only under very special circumstances will the gradient at the origin pass through the true maximum. Now the same can be said of any other point, and in particular, the point where the response is maximized on the gradient line through the origin, the codrdinates of which, are given by ( 3.77 ) . Hence, a second gradient line at the maximum of the first could be expected to pass through the true maximum only under special conditions.

It may be asked if there are points on the gradient line through the
origin where a second gradient line passes thr ough the true maximum. That the answer is in the affirmative can be seen from figure 2. Evidently, there are two such points, those two points where the gradient line intersects the axes of the elliptical contour system.


Figure 2. Intersection of the Gradient With the Axes of the Contour System .

If a point on the gradient be denoted by $\left(h \beta_{1}, h_{2}\right)$, then the slope of the line joining this point to the true maximum may be equated to the slope of the gradient at the point $\left(h_{1} \beta_{1}, h_{2} \beta_{2}\right)$. The resulting quadratic equation when solved for $h$ gives

$$
\begin{equation*}
h=\frac{\left(\beta_{11}+\beta_{22}\right) \pm \sqrt{\left(\beta_{11}-\beta_{22}\right)^{2}+\beta_{12}^{2}}}{\beta_{12}^{2}-4 \beta_{11} \beta_{22}} \tag{3.92}
\end{equation*}
$$

In general, the gradient line at the maximum point of the gradient line through the origin, does not pass through the true maximum. Is there any line such that the gradient at the maximum passes through the true maximum? That there are three such lines is the content of theorem 3.8.

Theorem 3.8. Let $\left(x_{1}{ }^{*}, x_{2}^{*}\right)$ be the coordinates of the maximum response attainable on the line $x_{2}=m x_{1}$. If the gradient at $\left(x_{1}{ }^{*}, x_{2}^{*}\right)$ passes through the coordinates of the true maximum, then $m$ has one of the values

$$
\begin{align*}
& m_{1}=\frac{2 \beta_{11} \beta_{2}-\beta_{12} \beta_{1}}{2 \beta_{22^{\beta} 1}-\beta_{12} \beta_{2}},  \tag{3.93}\\
& m_{2}=\frac{\left(\beta_{22}-\beta_{11}\right)+\sqrt{\left(\beta_{11}-\beta_{22}\right)^{2}+\beta_{12}^{2}}}{\beta_{12}}  \tag{3.94}\\
& m_{3}=\frac{\left(\beta_{22}-\beta_{11}\right)-\sqrt{\left(\beta_{11}-\beta_{22}\right)^{2}+\beta_{12}^{2}}}{\beta_{12}}
\end{align*}
$$

Proof: In accordance with (3.77), the maximum response attainable on the line $x_{2}=m x_{1}$ occurs at the point

$$
-\frac{a}{2 b},-\frac{m a}{2 b} \text { where } a \text { and } b \text { are given by (3.79). }
$$

According to theorem (3.7), the gradient at the point ( $\mathrm{x}_{1}{ }^{*}, \mathrm{x}_{2}{ }^{*}$ )
is perpendicular to the line $x_{2}=m x_{1}$. The equation of the line
perpendicular to the line $x_{2}=m x_{1}$ at the point $-\frac{a}{2 b} \cdot-\frac{m a}{2 b}$
has the equation

$$
\begin{equation*}
x_{2}=-\frac{1}{m} x_{l}-\frac{a\left(m^{2}+1\right)}{2 m b} \tag{3.96}
\end{equation*}
$$

If this line is to pass through the true maximum, then the point

$$
\begin{aligned}
& x_{1}^{*}=\frac{2 \beta_{1} \beta_{22}-\beta_{2}^{\beta} 12}{\beta_{12}^{2}-4 \beta_{11} \beta_{22}} \\
& x_{2}^{*}=\frac{2 \beta_{2} \beta_{11}-\beta_{1} \beta_{12}}{\beta_{12}^{2}-4 \beta_{11} \beta_{22}}
\end{aligned}
$$

must satisfy (3.96). Imposing this condition and substituting in the values of $a$ and $b$, there results

$$
\begin{align*}
& m^{3}\left[\beta_{12}\left(\beta_{2} \beta_{12}-2 \beta_{22} \beta_{1}\right)\right]+ \\
& m^{2}\left[4 \beta_{1} \beta_{22}\left(\beta_{22}-\beta_{11}\right)+2 \beta_{2} \beta_{12}\left(2 \beta_{11}-\beta_{22}\right)-\beta_{1} \beta_{12}^{2}\right]+ \\
& m\left[4 \beta_{11} \beta_{2}\left(\beta_{11}-\beta_{22}\right)+2 \beta_{1} \beta_{12}\left(2 \beta_{22}-\beta_{11}\right)-\beta_{2} \beta_{12}^{2}\right]+ \\
& \beta_{12}\left(\beta_{1} \beta_{12}-2 \beta_{11} \beta_{2}\right)=0 \tag{3.97}
\end{align*}
$$

Consider the line determined by $(0,0)$ and $\left(x_{1}, * x_{2}^{*}\right)$ 。
Obviously the coordinates of the maximum on this line are the coordinates of the true maximum so that trivially the gradient at this point goes through the true maximum. Hence one of the roots of (3.97) is $m_{1}$ as given in (3.93). Using $m_{l}$ in the synthetic division process yields the depressed equation

$$
\beta_{12} m^{2}+2\left(\beta_{11}-\beta_{22}\right) m-\beta_{12}=0,(3.98)
$$

the roots of which are $\mathrm{m}_{2}$ and $\mathrm{m}_{3}$ as given by (3.94) and (3.95).

Since $m_{l}$ is the slope of the line from the origin to the true maximum in the factor space, it is natural to expect $m_{2}$ and $m_{3}$ to have geometrical significance. It is easily verified that $m_{2} m_{3}=-1$ so that the lines $x_{2}=m_{2} x_{1}$ and $x_{2}=m_{3} x_{1}$ must be mutually perpendicular. Since the axes of the contour system in the factor space are also mutually perpendicular, it is reasonable to look for some relationship between the axes of the contour system and the lines which have slopes $m_{2}$ and $m_{3}$. The relationship is given in theorem 3.9.

Theorem 3.9. The lines $x_{2}=m_{2} x_{1}$ and $x_{2}=m_{3} x_{1}$
are parallel to the axes of the contour system.

Proof: By analytic gepmetry the angle of rotation of a conic is given by :

$$
\begin{equation*}
\tan \theta=\frac{1-\cos 2 \theta}{\sin 2 \theta} \tag{3.99}
\end{equation*}
$$

where

$$
\tan 2 \theta=\frac{\beta_{12}}{\beta_{11}^{-}-\beta_{22}}
$$

Now

$$
\begin{equation*}
\sin 2 \theta=\frac{\beta_{12}}{\sqrt{\left(\beta_{11}-\beta_{22}\right)^{2}+\beta_{12}^{2}}} \tag{3,100}
\end{equation*}
$$

and


Substituting (3.100) and (3.101) in (3.99) there results

$$
\tan \theta=\frac{\beta_{22}-\beta_{11}+\sqrt{\left(\beta_{11}-\beta_{22}\right)^{2}+\beta_{12}^{2}}}{\beta_{12}}=m_{2} .
$$

Since $m_{2}$ has the same slope as one of the axes, $m_{3}$ must have the same slope as the other.

The geometrical significance of the roots of the cubic equation (3.97) now becomes clear. The roots $\mathrm{m}_{2}$ and $\mathrm{m}_{3}$ are the slopes of lines passing through the origin which are parallel to the axes of the contour system. The maximum response is attained at the point of intersection of the line with an axis of the contour system. The gradient is coincident with the axis and hence passes through the center of the contour system. The situation is illustrated in figure 3.

An interesting situation arises when the axis of rotation passes through the origin. Suppose that the contour system in figure 3 is rotated until $x_{1}{ }^{\prime}$ coincides with the line $x_{2}=m_{1} x_{1}$. Then $m_{1}$ and $m_{2}$ are equal so that the lines $x_{2}=m_{1} x_{1}$ and $x_{2}=m_{2} x_{1}$ coincide. Further, since the $x_{1}^{\prime}$ axis passes through the origin and is normal to the contours, it must coincide with the gradient at the origin. Hence, $\beta_{2} / \beta_{1}=m_{2}=m_{1}$. These statements, which are geometrically evident, may readily be verified by the methods of analytic geometry.


Figure 3. Relationship Between $\mathrm{m}_{1}, \mathrm{~m}_{2}, \mathrm{~m}_{3}$, and the Contour System Axes .

## CHAPTER IV

## THE CHOICE OF SCALE

The most important decision that must be made by an experimenter using the steepest ascent technique occurs at the very outset of the investigation. He must decide where to run the experiments. The seemingly trivial nature of this statement may account in part for the lack of attention that it has received.

For purposes of illustration, suppose it is felt that the maximum yield obtained in a certain chemical process depends upon the controllable operating conditions of temperature and pressure. Obviously, the combinations of temperature and pressure should be chosen so that differences in the responses will be large enough to be clearly recognized. It is to be expected that, since investigators are human and possess varying degrees of experience, they will make different choices of the factor levels for the initial experiments. Unfortunately, as will be shown later, each different choice of factor levels leads to a different path for further experimentation.

Once the levels are chosen, the pressure and temperature variables are coded so that the rectangular array of experimental points map into the corner points of a unit grid. Hence the initial choice of levels determines the coding. The path of steepest ascent is calculated in terms of the coded variables. It is then decoded and written in terms of temperature and pressure. The investigator then proceeds
to vary the temperature and pressure in additional experiments along the decodedipath. This, in effect, is the procedure used in the field and is putlined by Box and Wilson (1), Davies (2), and Cochran and Cox (3). The crucial nature of the choice of levels and the effect of further experimentation along the decoded path is treated in the following discussion.

## The Lack of Scale Invariance of the

Steepest Ascent Path

Definition 4.1. The space of the uncoded variables in which the experiments are performed will be called the process space and will be denoted by $P$.

In this section the following notation will be used:
$z=\left(z_{1}, z_{2},, \ldots, z_{p}\right)$ will denote a point in $P$,
$Z$ will denote the design matrix of the model,
$Z_{1}$ will denote the design matrix of the initial experimental points, and
$Z_{2}$ will denote the design matrix of an additional experimental point or additional experimental points.

As an example, suppose that the response in $P$ is
$y=\delta_{0}+\delta_{1} z_{1}+\delta_{2}^{z_{2}}+\delta_{11} z_{1}^{2}+\delta_{12^{z} 1_{2}}+\delta_{22_{2}}{ }^{2}+$ e.

This may be written as

$$
\left.\begin{aligned}
& y=\left(1, z_{1}, z_{2}, z_{1}^{2}, z_{1}, z_{2}, z_{2}^{2}\right) \\
& \delta_{1} \\
& \delta_{2} \\
& \delta_{11} \\
& \delta_{12} \\
& \delta_{22}
\end{aligned} \right\rvert\,+\begin{gathered}
\delta_{0} \\
\\
=Z \Delta+e .
\end{gathered}
$$

If the initial experimental points in $P$ are $(-2,-3),(-2,3)$, $(2,-3)$ and $(2,3)$, then

$$
Z_{1}=\left(\begin{array}{rrr}
1 & -2 & -3 \\
1 & -2 & 3 \\
1 & 2 & -3 \\
1 & 2 & 3
\end{array}\right)
$$

If an additional experiment is run at $z=(5,6)$, then $Z_{2}=(1,5,6)$.

Definition 4.2 . The space of the coded process space variables will be called the design space and will be denoted by $D$.

The symbols $x, X, X_{1}$ and $X_{2}$ will be usedin $D$ in the same fashion as the symbols $Z, Z, Z_{1}$, and $Z_{2}$ are used in $P$.

In order to map the points of a p-dimensional rectangle in $P$ into a unit grid whose center is at the origin in $D$, a translation followed by a change of scale must be used. No loss of generality
will occur if it is assumed that the center of the rectangle in $P$ has already been translated to the origin in $P$. With this assumption the mapping of the rectangular grid in $P$ to the unit grid in $D$ can be accomplished by a change in scale alone.

Definition 4.3. The transformation which maps the unit grid points of $D$ into the grid points of $P$ will be denoted by $T_{1}$.

As an example suppose that $p=2$ and that the process points are $(-40,-20),(-40,20),(40,-20)$ and $(40,20)$, then the mapping required to transform the unit grid points into the process points is represented by the equation
$\left(1, z_{1}, z_{2}\right)=\left(1, x_{1}, x_{2}\right)\left(\begin{array}{ccc}1 & 0 & 0 \\ 0 & 40 & 0 \\ 0 & 0 & 20\end{array}\right)$. In this example $T_{1}$ is the matrix $\left(\begin{array}{ccc}1 & 0 & 0 \\ 0 & 40 & 0 \\ 0 & 0 & 20\end{array}\right)$. It is to be noticed that $T_{1}$ is associated with the decoding and is specified at the moment that the experimenter chooses the process space grid points.

Whereas $T_{1}$ will be used with the $1, x p+1$ matrices $X_{1}$ and $Z_{1}, T$ will be used with the $1 \times q$ matrices $X$ and $Z$. For example, suppose that

$$
\begin{aligned}
y & =\beta_{0}+\beta_{1} x_{1}+\beta_{2} x_{2}+\beta_{11} x_{1}^{2}+\beta_{12} x_{1} x_{2}+\beta_{22^{x}}{ }^{2}+e \\
& =x \beta+e .
\end{aligned}
$$

In this case the transformation $Z=X T$ may be written out as $\left(1, z_{1}, z_{2}, z_{11}, z_{1} z_{2}, z_{22}\right)=\left(1, x_{1}, x_{2}, x_{11}, x_{1} x_{2}, x_{22}\right) \operatorname{diag} \cdot\left(1, t_{1}, t_{2}, t_{1}^{2}, t_{1} t_{2}, t_{2}^{2}\right)$.

Evidently the matrix $\mathrm{T}_{1}$ determines the matrix $\quad \mathrm{T}$.
It is necessary to distinguish between the steepest ascent line fitted to the unit grid in $D$, its map in $P$, and the steepest ascent line in $P$ calculated from the process grid points.

Definition 4.4. The steepest ascent line calculated in $D$
from a unit grid will be denoted by $L_{d}$. The transform of $L_{d}$
will be denoted by $L_{t}$. The steepest ascent line in $P$ calcu-
lated from the experimental grid points will be denoted by $L_{p}$.
It is indeed unfortunate that it is generally necessary to distinguish between $L_{t}$ and $L_{p}$. That this is true is the content of theorem 4.1.

Theorem 4.1. A necessary and sufficient condition that $L_{t}=L_{p}$ is that

$$
\begin{equation*}
T_{1}=\operatorname{diag} .\left(1, t I_{p}\right) \tag{4.1}
\end{equation*}
$$

Proof: In order to prove necessity let the response in $D$ be given by

$$
\begin{equation*}
y=x \beta+e \tag{4.2}
\end{equation*}
$$

where $X$ is $1 \times q$ and $\beta$ is $q \times t .: \operatorname{Let} T_{1}=\operatorname{diag} .\left(1, t_{1}, \ldots, t_{p}\right)$.
Now $L_{d}$ is determined by the equations

$$
\begin{equation*}
x_{j}=\frac{\hat{\beta}_{j}}{\hat{\beta}_{l}} x_{1}, j=2,3, \ldots, p \tag{4.3}
\end{equation*}
$$

According to (4.1), $x_{j}=z_{j} / t_{j}, j=0,1, \ldots, p$
so that (4.3) becomes

$$
\frac{z_{j}}{t_{j}}=\frac{\hat{\beta}_{j}}{\hat{\beta}_{1}} \frac{z_{1}}{t_{1}} \text { or } z_{j}=\frac{t_{j}}{t_{1}} \frac{\hat{\beta}_{j}}{\hat{\beta}_{1}} z_{1} \quad j=2,3, \ldots p(4.4)
$$

Equation (4.4) determines the line $L_{t}$. It is necessary to find the expression for $L_{p}$. Since $Z=X T, X=Z T^{-1}$ which may be substituted in (4.2) to obtain

$$
\begin{equation*}
y=Z T^{-1} \beta+e=Z \Delta+e \tag{4.5}
\end{equation*}
$$

where $\Delta=T^{-1} \beta$, It is desired to find $\hat{\Delta}$ in the equation

$$
\begin{equation*}
\hat{y}=\mathrm{z} \hat{\Delta} \tag{4.6}
\end{equation*}
$$

where $Z$ is $1 \times p+1$ and $\hat{\Delta}$ is $p+1 \times 1$. According to least \$quares.

$$
\begin{aligned}
\hat{\Delta} & =\left(Z_{1}^{\prime} Z_{1}\right)^{-1} Z_{1}^{\prime} Y_{1} \\
& =\left(T_{1} X_{1}^{\prime} X_{1} T_{1}\right)^{-1} T_{1} X_{1}^{\prime} Y_{1} \\
& =\left(N T_{1}^{2}\right)^{-1} T_{1} X_{1}^{\prime} Y_{1}=N^{-1} T_{1}^{-1} X_{1}^{\prime} Y_{1} \\
& =T_{1}^{-1} \hat{\beta}
\end{aligned}
$$

where $N=2^{P}$, Then the line $L_{p}$ is determined by the equations

$$
\begin{equation*}
z_{j}=\frac{\frac{\hat{\beta}_{j}}{t_{j}}}{\frac{\hat{\beta}_{1}}{t_{1}}} \quad z_{1}=\frac{t_{1} \hat{\beta}_{j}}{t_{j} \hat{\beta}_{1}} z_{1}, \quad j=2,3, \ldots ., p \tag{4.8}
\end{equation*}
$$

Hence if the line $L_{p}$ is to be the same as the line $L_{t}$, then by (4.4) and (4.8) ,

$$
\frac{t_{j}}{t_{1}} \frac{\hat{\beta}_{j}}{\hat{\beta}_{l}}=\frac{t_{1} \hat{\beta}_{j}}{t_{j} \hat{\beta}_{l}}, \quad j=2,3, \ldots, p
$$

which implies that

$$
t_{1}=t_{2}=\cdots,=t_{p}
$$

That $T_{1}=$ diag. $\left(1, t I_{p}\right)$ is a sufficient condition follows from equations (4.4) and (4.8).

It is instructive to examine the steepest ascent procedure in the light of theorem 4. 1 . The experimenter determines $T_{1}$ at the outset. In theory, experimentation along the path of steepest ascent is an attempt to proceed in a direction which is normal to the factor space contours at the center of the experiment. Now the lines $L_{p}$ and $L_{t}$ cannot both be normal to the contours at the same point in P. A little reflection convinces one that it is the line $L_{p}$ which is trying to estimate the gradient at the center of the rectangle in $P$. However, the method of steepest ascent calls for further experimentation on $L_{t}$, The situation is shown in figure 4 ,

> Experimentation Along the Transformed Path of Steepest Ascent

The lack of scale invariance exhibited by the steepest ascent line and the invariance property of the steepest ascent line, previously examined in chapter III, combine in the present instance to make matters still more bizarre.


Figure 4. The Effect of a Change of Scale on the Steepest Ascent Path .

It is known from chapter III that if an additional point is run on the line $L_{d}$ and a new path calculated, the new path will still be the line $L_{d}$. The additional point on $L_{d}$ maps by $T_{i}$ to an additional point on $L_{t}$. Suppose that in $P$ a second steepest ascent path is calculated by use of the process points and the additional point on $L_{t}$. What is the new steepest ascent line in the space $P$ ? In order to answer this question it is necessary to examine the property of invariance using a $p$-dimensional rectangular grid instead of a p - dimensional unit grid as in chapter III.

Theorem 4,2. The steepest ascent line $L_{p}$ is invariant if the additional point is on the line $L_{t}$.

Proof: The additional point on $L_{t}$ is the map of an additional point on $L_{d}$. Let the point on $L_{d}$ be written in matrix form as $X_{2}=\left(1, h \hat{\beta}_{1}, \ldots h \hat{\beta}_{p}\right)=\left(1, h B_{1}^{\prime}\right)$.

The additional point is mapped from $D$ to $P$ by

$$
\begin{equation*}
Z_{2}=X_{2} T_{1} \tag{4,10}
\end{equation*}
$$

$\mathrm{T}_{1}$ may be partitioned as

$$
T_{1}=\left(\begin{array}{cc}
1 & 0  \tag{4.11}\\
0 & \text { diag. }\left(t, \ldots t_{p}\right)
\end{array}\right)=\left(\begin{array}{cc}
1 & 0 \\
0 & T_{11}
\end{array}\right)
$$

Substituting $\mathrm{X}_{2}$ as given in (4.9) in (4.10) and partitioning $\mathrm{T}_{1}$ as in (4.11), there results

$$
Z_{2}=\left(1, \mathrm{hB}_{1}^{\prime}\right)\left(\begin{array}{cc}
1 & 0  \tag{4.12}\\
0 & \mathrm{~T}_{11}
\end{array}\right)=\left(1, \mathrm{hB}_{1}^{\mathrm{t}} \mathrm{~T}_{11}\right) .
$$

As in theorem 4.1, the equation $y=X \beta+e$ in $D$ becomes the equation $y=Z \Delta+e$ in $P$ where

$$
\Delta=\mathrm{T}^{-1} \cdot \beta
$$

Let

$$
\vec{\Delta}=\left(\left.\begin{array}{c}
\hat{\delta}_{0} \\
\frac{\hat{\delta}_{1}}{\cdots} \\
\cdots \\
\cdots \\
\hat{\delta_{p}}
\end{array} \right\rvert\,=\binom{\hat{\Delta}_{0}}{\hat{\Delta}_{1}}\right. \text { be calculated from the initial experimental }
$$

points. Let $\hat{a}=\left|\begin{array}{c}\hat{a}_{0} \\ \frac{\hat{a}_{0}}{a_{1}} \\ \vdots \\ \vdots \\ \hat{a}_{p}\end{array}\right|=\binom{\hat{A}_{0}}{\hat{A}_{1}}$ be the revision of $\hat{\Delta}$
based on the initial experimental points and the additional point $Z_{2}$ on $L_{t}$. Then according to (3.11),

$$
\begin{equation*}
\hat{a}=\hat{\Delta}+J^{\prime} G Q \tag{4.13}
\end{equation*}
$$

where the symbols $J, G$, and $Q$ are used as in theorem 3,7. By definition
$J=Z_{2}\left(Z_{, 1}^{\prime} Z_{1}\right)^{-1}$ so that $J^{\prime}=\left(Z_{1}^{\prime} \cdot Z_{1}\right)^{-1} Z_{2}^{\prime}$.
Now $\left(Z_{1}^{\prime} Z_{1}\right)^{-1}=\left(T_{1} X_{1}^{\prime} X_{1} T_{1}\right)^{-1}=\frac{\left(T_{1}{ }^{2}\right)^{-1}}{N}$

$$
=\frac{1}{N}\left(\begin{array}{cc}
1 & 0  \tag{4.15}\\
0 & \left(T_{11}^{2}\right)-1
\end{array}\right) \text { where } N=2^{p}
$$

Substituting the expression for $Z_{2}$ as given in (4.12) and the expression for $\left(Z_{1}^{\prime} Z_{1}\right)^{-1}$ as given in (4.15) into the expression for $J^{\prime}$ as given in (4.14), there results

$$
\begin{align*}
& J^{\prime}=\frac{1}{N}\left(\begin{array}{cc}
1 & 0 \\
0 & \left(T_{11}{ }^{2}\right)^{-1}
\end{array}\right)\binom{1}{h E_{11} B_{1}}=\frac{1}{N}\binom{1}{h T_{11}-B_{B}} \\
& =\frac{1}{\mathrm{~N}}\left(\begin{array}{ccc}
1 \\
\mathrm{~h} & \hat{\Delta}_{1}
\end{array}\right) . \tag{4.16}
\end{align*}
$$

Substituting the value of $\mathrm{J}^{\prime}$ as given in (4.16) into the expression for $\bar{a}$ given in (4.13) and noting that $G Q$ is a scalar, there results

$$
\begin{equation*}
\binom{\hat{A}_{0}}{\hat{A}_{1}}=\binom{\hat{\Delta}_{0}}{\hat{\Delta}_{1}}+\frac{G Q}{N}\binom{1}{h \hat{\Delta}_{1}} \tag{4.17}
\end{equation*}
$$

It follows that
$\hat{A}_{1}=\hat{\Delta}_{1}+\frac{G Q}{N h} \hat{\Delta}_{1}=\hat{\Delta}_{1}\left(1+\frac{G Q}{N h}\right)=k_{1} \hat{\Delta}_{1}$
where the scalar $k=1+\frac{G Q}{N h}$. Hence $\hat{A}_{1}$ isa scalar multiple of $\hat{\Delta}_{l}$ so that the path $L_{p}$ remains unchanged

The theorem 4.2 generalizes so that the line $L_{p}$ remains invariant regardless of how many additional points are taken on $L_{t}$.

Theorem 4.3. The steepest ascent line $L_{p}$ is invariant if $N_{2}$ additional points are taken on $L_{t}$.

Proof: The $\mathrm{N}_{2}$ additional points on $\mathrm{L}_{\mathrm{d}}$, written in matrix notation as

$$
\left.X_{2}=\left|\begin{array}{l}
1 \\
1 \\
\cdot \\
\cdot \\
1
\end{array}\right| \begin{array}{cccc}
h_{1} \hat{\beta}_{1} & h_{1} \hat{\beta}_{2} & \cdots & h_{1} \hat{\beta}_{p} \\
h_{2} \hat{\beta}_{1} & h_{2} \hat{\beta}_{2} & \cdots & h_{2} \hat{\beta}_{p} \\
\cdot & & & \\
{ }^{h_{N_{2}}} \hat{\beta}_{1} & h_{N_{2}} \hat{\beta}_{2} & \cdots & h_{N_{2}} \hat{\beta}_{p}
\end{array} \right\rvert\,=\left(W_{0}, H_{1}^{\prime}\right)
$$

map to $\mathrm{N}_{2}$ additional points on $\mathrm{L}_{\mathrm{t}}$. This may be written as $Z_{2}=X_{2} T_{1}=\left(W_{0}, H B_{1}^{\prime}\right)\left(\begin{array}{cc}1 & 0 \\ 0 & T_{11}\end{array}\right)=\left(W_{0}, H B_{1}^{\prime} T_{11}\right)$.

Let $\hat{a}$ and $\hat{\Delta}$ be as in theorem 4.2. Then

$$
\begin{align*}
J^{\prime} & =\left(Z_{1}^{\prime} Z_{1}\right)^{-1} Z_{2}^{\prime}=\frac{1}{N}\left(\begin{array}{cc}
1 & 0 \\
0 & \left(T_{11}^{2}\right)^{-1}
\end{array}\right)\binom{W_{0}^{\prime}}{T_{11} B_{1} H^{1}} \\
& =\frac{1}{N}\binom{w_{0}^{\prime}}{T_{11}^{-1} B_{1} H^{\prime}} . \tag{4.20}
\end{align*}
$$

Substituting this expression for $J^{\prime}$ into the equation $\hat{a}=\hat{\Delta}+J^{\prime} G Q$, there results
$\binom{\hat{A}_{0}}{\hat{A}_{1}}=\binom{\hat{\Delta}_{0}}{\hat{\Delta}_{1}}+\frac{1}{\mathrm{~N}}\binom{W_{0}^{t}}{T_{11}^{-1} \mathrm{~B}_{1} \mathrm{H}^{s}} \quad \mathrm{GQ}$.
Hence

$$
\begin{equation*}
\hat{A}_{1}=\hat{\Delta}_{1}+\frac{1}{N} T_{11}^{-1} B_{1} H^{\prime} G Q . \tag{4.22}
\end{equation*}
$$

Now $H^{\prime} G Q$ is a scalar so that (4.22) becomes

$$
\begin{align*}
\hat{A}_{1} & =\hat{\Delta}_{1}+\frac{H^{\prime} G Q}{N} T_{11}^{-1} B_{1}=\hat{\Delta}_{1}+\frac{H^{\prime} G Q}{N} \hat{\Delta}_{1} \\
& =\hat{\Delta}_{1}\left(1+\frac{H^{\prime} G Q}{N}\right)=k \hat{\Delta}_{1} \tag{4.23}
\end{align*}
$$

where $k=1+\frac{H^{t} G Q}{N}$. Hence $\hat{A}_{1}$ is a scalar multiple of $\hat{\Delta}_{1}$ and the path $L_{p}$ remains unchanged.

The Choice Between the Lines $L_{d}$ and $L_{p}$
It is clear from the preceding theorems that in $P$, the path of steepest ascent and the path upon which the experiments are being run, are generally distinct. Since experimentation along the path $L_{t}$ is customary rather than mandatory, it is fair to ask which path comes closest to the true maximum. The answer depends upon the structure of $\eta$ in the equation $y=\eta+e$. The complexity of the problem increases with the dimension $p$ and the degree of $\eta$. In order to simplify the problem and visualize the geometry involved, let $\eta$ be of second degree in the variables $x_{1}$ and $x_{2}$. Under the transformation

$$
\begin{aligned}
& z_{1}=t_{1} x_{1} \\
& z_{2}=t_{2} x_{2}
\end{aligned}
$$

the response in - D, which is written
$y=\beta_{0}+\beta_{1} x_{1}+\beta_{2} x_{2}+\beta_{11} x_{1}^{2}+\beta_{12} x_{1} x_{2}+\beta_{22} x_{2}^{2}+e$,
is written in $P$ as
$y=\delta_{0}+\delta_{1} z_{1}+\delta_{2} z_{2}+\delta_{11} z_{1}^{2}+\delta_{12} z_{1} z_{2}+\delta_{22^{z}}{ }^{2}+e$ where $\delta_{0}=\beta_{0}, \delta_{1}=\beta_{1} / t_{1}, \cdots \delta_{22}=\beta_{22} / t_{2}^{2}$.

The coordinates of the maximum of $\eta$ in $P$ are by (3.74),

$$
z_{1}^{*}=\frac{2 \delta_{1} \delta_{22}-\delta_{2} \delta_{12}}{\delta_{12}^{2}-4 \delta_{11} \delta_{22}}
$$

$$
\begin{equation*}
z_{2}^{*}=\frac{2 \delta_{2} \delta_{11}-\delta_{1} \delta_{12}}{\delta_{12}^{2}-4 \delta_{11} \delta_{22}} . \tag{4.24}
\end{equation*}
$$

The equation of the line joining $(0,0)$ to $\left(z_{1}{ }^{*}, z_{2}{ }^{*}\right)$ is then

$$
\begin{equation*}
z_{2}=\frac{2 \delta_{2} \delta_{11}-\delta_{1} \delta_{12}}{2 \delta_{1} \delta_{22}-\delta_{2} \delta_{12}} z_{1} . \tag{4.25}
\end{equation*}
$$

Definition 4.5. Let the line joining the origin to the point in the $P$ space where the response is maximized be called $L_{m}$. There have been three lines defined in $P$, the lines $L_{t}, L_{p}$, and $L_{m}$. Their equations in terms of the $\beta^{\prime} s$ are

$$
\begin{align*}
& L_{t}: \quad z_{2}=\frac{t_{2} \hat{\beta}_{2}}{t_{1} \hat{\beta}_{1}} z_{1},  \tag{4.26}\\
& L_{p}: \quad z_{2}=\frac{t_{1} \hat{\beta_{2}}}{t_{2} \hat{\beta}_{1}} \quad z_{1}, \tag{4.27}
\end{align*}
$$

and

$$
\begin{equation*}
L_{m}: \quad z_{2}=\frac{t_{2}\left(2 \beta_{2} \beta_{11}-\beta_{1} \beta_{12}\right)}{\mathrm{t}_{1}\left(2 \beta_{1} \beta_{22}-\beta_{2} \beta_{12}\right)} z_{1} . \tag{4,28}
\end{equation*}
$$

The various angles that these lines make with one another may be found by use of the formula

$$
\tan \theta=\frac{m_{2}-m_{1}}{1+m_{1} m_{2}} \text { where } \theta m_{1} \text {, and } m_{2} \text { are the }
$$

conventional symbols of analytic geometry. One might be inclined to choose the path for which the angle between the path and the line to the
maximum is minimized, Alternatively, with the aid of equation (3.78), a comparis on could be made between the maximum value of $\eta$ attainable on $L_{t}$ and the maximum value of $\eta$ attainable on $L_{p}$. Each of the above approaches are algebraically involved and shed little light upon the strategy that the experimenter should adopt. For this reason the discussion will be of a geometrical nature taking full advantage of the simplicity of a two dimensional factor space. Frequently the investigator will feel that he knows, in a general way, the orientation of the contour system. Naturally he will choose the initial experimental points to be in harmony with his conjecture. Thus, if the experimenter's conception of the contour system is as in figure 5 , he would be quite likely to feel that the rectangular grid shown there is appropriate.


Figure 5. An Appropriate Grid for Horizontal Elongated Contours.

The mapping of the unit grid points in $D$ to the rectangle in $P$ may be written

$$
\begin{aligned}
& z_{1}=t_{1} x_{1} \\
& z_{2}=t_{2} x_{2}
\end{aligned}
$$

where, in accordance with figure $5,{ }^{\mathrm{t}}{ }_{1}>\mathrm{t}_{2}$.
Let the true maximum of a horizontal elongated contour system be in the first quadrant and assume that the slope of the major axis is positive. Under these conditions there are four distinct cases which may arise. The first case is shown in figure 6.


Figure
for a Decision.

The line $L_{t}$ can have either of the positions indicated in figure 6 by $L_{t}^{\prime}$, or $L_{t}{ }^{\prime \prime}$. The choice as to which path to take is left to the judgement of the experimenter.

In figure 7 the line $L_{p}$ lies between the line to the true maximum and the line $L_{t}$. Hence, the experimenter should choose the line $L_{p}$ for further experimentation.


Figure 7 . The Slope of $L_{p}$ is Positive; the Line $L_{p}$ Should
be Chosen.

In the third case, the slopes of the lines $L_{t}$ and $L_{p}$ are both negative as shown in figure 8.


Figure 8. The Slope of $L_{p}$ is Negative; the Line $L_{t}$ Should be Chosen.

The line $L_{t}$ should be chosen as a path for further experimentation since it lies between $L_{p}$ and $L_{m}$.

In figure 9 the slopes of $L_{t}$ and $L_{p}$ are negative. $L_{p}$ should be chosen as a path for further experimentation since it lies between $L_{t}$ and $L_{m}$.

Similar analyses are applicable to the various other orientations of the contour system axes which the experimenter may suspect.



The foregoing remarks are applicable to a p-dimensional factor space. Should information regarding the orientation of any one of the $p$ axes become available, the preceding arguments remain valid.

## CHAPTER V

## SUMMARY

The method of steepest ascent is part of a general technique which has been proposed to locate a maximum, minimum, or stationary point in the factor space. The use of the steepest ascent method is often accompanied by scant knowledge of its properties. In this study various properties of the steepest ascent path are investigated. Particular attention is given to the estimation of the coefficients of the fitted plane which determine the path of steepest ascent, the direction of additional paths calculated by the use of points on the first path, and the behavior of the steepest ascent path when subjected to a change of scale.

## Conclusions

In chapter II several theorems are presented which are concerned with the coefficients of the plane fitted to a quadratic response function. This plane is fitted by the method of least squares from a rectangular grid in the factor space. It is shown that the linear coefficeints of the fitted plane are unbiased estimates of the first order partial derivatives of the response function evaluated at the center of the grid.

It is shown that, in the absence of experimental error, a p-dimensional hyperplane passes through the responses at the $2^{p}$ points in the factor space provided that there are no cross product
terms in the quadratic response function. This surprising fact is used to divide the factor space into three-point sets, the set where the predicted response $\hat{y}$ equals the observed response $y$, the set where $y>\hat{y}$, and the set where $y<\hat{y}$.

Also it is shown that decreasing the variance of $\hat{\beta}_{1}$ or $\hat{\beta}_{2}$ by increasing the design size results in an increase in the bias of these estimators. Since mean square error is dependent on both bias and variance, a study is made of the relationship between design size and mean square error. A design size is given which will minimize the mean square error in the case where a quadratic response model is chosen, when in truth, the response function is cubic.

In chapter III it is shown that the revised path of steepest ascent, calculated by means of the initial grid and additional points on the initial steepest ascent line, is the same as the initial path of steepest ascent. This property of the steepest ascent line, termed invariance, is shown to hold when an arbitrary number of additional points are added to the initial steepest ascent line. It is shown that the three conditions, invariance, location on the path of steepest ascent, and the use of a square grid, are so related that any two of them imply the third. These results make it impossible to obtain a different path of steepest ascent by this approach.

An expression is obtained for $\hat{y}_{2}$ in terms of $\hat{y}_{1}$ and $y$ where $\hat{y}_{2}$ is the response predicted on the basis of $N+1$ points, $\hat{y}_{1}$ is the response predicted on the basis of $N$ points, and $y$ is the observed response. It is shown that $\hat{y}_{2}$ is a convex combination
of $y$ and $\hat{y}_{1}$.
If the contours are elliptic and a point exists in the factor space where the response is maximized, a method is given whereby upper and lower bounds to the ratio $\beta_{11} / \beta_{22}$ may be estimated. Some properties of the gradient line are examined. The coordinates of the maximum response attainable on the gradient line and the value of the response at that point are given.

Chapter IV deals with the behavior of the steepest ascent line when it is subjected to a change of scale. Necessary and sufficient conditions are given for the steepest ascent line to be invariant under a change of scale.

It is shown that if points along the transform of the steepest ascent path in the space of the coded variables are used, together with the initial experimental points, to calculate a revised path of steepest ascent, the revised path will be the same as the initial path.

Examples are given which demonstrate the manner in which prior knowledge of the general orientation of the contour system may be used in making a choice between the transform of the coded steepest ascent line and the steepest ascent line in the uncoded variables.

## Areas for Future Research

As additional experiments are performed along the path of steepest ascent, the transformed path of steepest ascent, or indeed any line, additional information becomes available at each point.

This information is, of course, contained in the value of the observation
at the additional point. A reasonable strategy would be to use this information to revise the path before proceeding to the next point. Although the unexpected invariance property seems to be an obstacle to such a strategy, further work should be done in this area.

Instead of experimenting along some path believed to be desirable, it might be just as well to estimate all of the unknown parameters initially and proceed along the path to the maximum calculated from the estimated parameters. For example, in the case of a quadratic response in two variables, the six unknown parameters could be estimated by means of an initial six point design in the factor space. As a competing strategy, the unknown parameters could be estimated by means of the initial design points and points along the line of steepest ascent. In the example of the six unknown parameters, the full quadratic could be estimated upon the completion of the second experiment on the steepest ascent line. The difficulty lies in finding an adequate basis for comparing these strategies.

The discussion in chapter IV concerning the choice between the lines $L_{p}$ and $L_{t}$ was restricted severely. The more general case, in which the location of the true maximum and the orientation of the contour system are arbitrary, needs attention. This problem should be analyzed both from an algebraic and a geometric viewpoint.
(L) Box, G. E. P., and K. B. Wilson. 'On the Experimental Attainment of Optimum Conditions. " Journal of the Royal Statistical Society. Vol. 13, 1-45, 1951- —
(2) Davies, O. L. The Design and Analysis of Industrial Experiments. New York: Hofner Publishing Company, 1954 .
(3) Cochran, W.C. and G.M. Cox. Experimental Designs. New York: John Wiley and Sons, 1957.
(4) Kempthorne, O., B. V. Shah, and R. J. Buehler. Some Properties of Steepest Ascent and Related Procedures for Finding Optimum C onditions. Technic al Report $\overline{\text { Number 1, }}$ Office of $\overline{\text { Naval Research, }}$ NR-042-207, April, 1961.
(5) Plackett, R. L., "Some Theorems In Least Squares." Biometrika. Vol 37, 149-157, 1950.

VITA

Charles Henry Johns on

Candidate for the Degree of

Doctor of Philosophy

Thesis: SOME PROPERTIES OF THE METHOD OF STEEPEST ASCENT

Major Field: Statistics
Biographical:
Personal Data: Born at Chicago, Illinois, June 12, 1925, the son Emil and Abbey Johnson.

Education: Attended grade school in Chicago, Illinois; graduated from Lane Technical High School in 1943; received the Associate of Arts degree from North Park Junior College in June, 1947; received the Bachelor of Science degree from Bradley University, with a major in Mathematics, in June, 1949; received the Master of Science degree from Bradley University, with a major in Mathematics, in June, 1950; completed the requirements for the Doctor of Philosophy degree in August, 1963.

Professional Experience: Entered the United States Army in January, 1944; held Graduate Assistantship at University of Pittsburgh, 1950-52; worked as an Assistant Actuary, Continental Casualty Company, 1952-1955; taught as an instructor of mathematics at De Pauw University, 19551958; Assistant Professor of Mathematics and Astronomy, 1958-1963.

