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Techniques of Water-Resources Investigations of the United States Geological Survey

Chapter B4

# REGRESSION MODELING OF GROUND-WATER FLOW 

By Richard L. Cooley and Richard L. Naff

Book 3
APPLICATIONS OF HYDRAULICS

# UNITED STATES DEPARTMENT OF THE INTERIOR <br> MANUEL LUJAN, JR., Secretary 

GEOLOGICAL SURVEY
Dallas L. Peck, Director

## Preface

Scientists and engineers have been using ground-water flow models to study ground-water flow systems for more than 20 years. The basic modeling process seems to be relatively straightforward. Initially, a sound conceptual model is formed and is translated into a tractable, mathematical model. Contributing to (and following) this conceptualization process is the collection of field information, such as (1) location and extent of hydrostratigraphic units, recharge areas, discharge areas, and system boundaries; (2) hydraulic head measurements; and (3) pumping discharges. These data form the basis for input to the flow model. Finally, the model is run, and the desired information such as head distribution or flux rates is extracted. However, people engaged in modeling usually observe that two pervasive problems considerably complicate the situation. One problem is that good, general methods of measuring (or computing) some of the variables that characterize the flow system and its geologic framework do not exist. One example is measurement of groundwater recharge. No direct ways of measuring recharge exist, and the accuracy of indirect methods is often unknown. Furthermore, many indirect methods are applicable only to unique situations. The second problem relates to errors in the measurements and their propagation into model results. No error-free measurement (or computation) methods for obtaining data on the flow system exist. Thus, even the variables that can be estimated will contribute to error, so that model results will always be unreliable to some extent. As a consequence of these two problems, measurement (or computation) of the necessary input variables, application of them to an adequate model, and calculation of the desired results to an acceptable accuracy generally are not possible. Other methods that recognize and deal with the problems of incompleteness and (or) inaccuracy of data must also be applied. The present text has been designed to teach these methods to scientists and engineers engaged in ground-water modeling.

The basic methodology is multiple, nonlinear regression, in which the regression model is some type of ground-water flow model. As seen subsequently, this methodology is consistent with known aspects of the physical systems to be analyzed and requires relatively few assumptions. Even though the present text is directed specifically toward ground-water modeling, the procedures to be discussed are applicable to a number of different types of modeling problems. Thus, the methods are usually discussed in a general context; in other words, without reference to any specific model.

Material in the present text evolved from notes developed for training courses in parameter estimation for ground-water flow models taught by the authors and others at the U.S. Geological Survey National Training Center, Denver Federal Center, Lakewood, Colo. The philosophy of these courses, and of this text, is to teach general methods that are applicable to a wide range of problems and to teach these methods in sufficient depth so that students can apply them to many problem situations not considered in the courses or text.

The main body of the text is organized into six major sections. The first section is an introduction that discusses the general topic of modeling ground-water flow. This section shows that groundwater modeling problems are an incomplete combination of direct-type problems (solution for hydraulic head given values of flow system and framework variables) and inverse-type problems (solution for flow system and framework variables given values of hydraulic head) that commonly require solution by optimization procedures which give the best fit between observed and calculated results. Because the specific optimization approach employed here is regression and regression procedures are based on statistical concepts, the second section is included to provide the student with the necessary statistical background material. It is not designed to be an exhaustive review of basic statistics; rather, it presents material essential to understanding the following sections. The third section presents detailed material on linear and nonlinear regression. Although most
of the material on linear regression is fairly standard, some of the material on nonlinear regression is not. In particular, specific modifications presented to induce convergence of the iterative solution procedure for nonlinear regression have not, to the writers' knowledge, been presented elsewhere in the form given here. The fourth section applies the nonlinear regression method to the specific problem of developing a general finite-difference model of steady-state ground-water flow. In the fifth section, statistical procedures are given to analyze and use general linear and nonlinear regression models. The tests and analytical procedures presented are not exhaustive; they are the ones that the writers have found to be most useful for analyzing the real systems examined to date. The sixth section is designed to be supplemental to the preceding sections. Specialized procedures presented include nonlinear regression for models that cannot be solved directly for the dependent variable, a measure of model nonlinearity called Beale's measure, and a statistical test for compatibility of prior information on parameters and parameter estimates derived from sample (observed head) information.

A number of exercises have been included, and a complete discussion of the answers can be found in the seventh major section at the end of the text. These problems exercise the student on nearly all methods presented. In addition, three computer programs are documented and listed: the program for nonlinear-regression solution of ground-water flow problems of section four, a program to calculate Beale's measure, and a program to calculate simulated errors in computed dependent variables such as hydraulic head.

The mathematical background necessary to use this text includes basic mathematics through differential and integral calculus, including partial derivatives, and matrix algebra. A background in elementary statistics would be useful but is not essential. In addition, a sound knowledge of ground-water hydrology and ground-water flow modeling are needed to effectively apply the methods presented.

References for cited material are given at the end of each major section. Good supplemental sources for the unreferenced material not peculiar to this text are presented as "Additional Reading" at the end of each reference list. It is expected that students who have difficulty with the material in this text will consult the more expanded developments in these supplemental sources.

Several people, in addition to the writers, contributed extensively to this text. Charles R. Faust wrote earlier sections on statistical review and basic regression and contributed several exercises, Steven P. Larson wrote an earlier version and documentation of the nonlinear-regression flow program of section four and contributed earlier versions of several exercises, James V. Tracy contributed to the documentation of the nonlinear-regression flow program, and Thomas Maddock III wrote the first version of the statistics review section. In addition, all of these people helped teach the training courses from which the present text evolved. Finally, the writers would like to thank the technical reviewers, Brent M. Troutman and Allan L. Gutjahr, for their many hours of review work and the secretaries, Anita Egelhoff, Evelyn R. Warren, and Patricia A. Griffith, for their patience and care in typing the manuscript.

## TECHNIQUES OF WATER-RESOURCES INVESTIGATIONS OF THE UNITED STATES GEOLOGICAL SURVEY

The U.S. Geological Survey publishes a series of manuals describing procedures for planning and conducting specialized work in water-resources investigations. The manuals published to date are listed below and may be ordered by mail from the U.S. Geological Survey, Books and OpenFile Reports, Federal Center, Box 25425, Denver, Colorado 80225 (an authorized agent of the Superintendent of Documents, Government Printing Office).

Prepayment is required. Remittance should be sent by check or money order payable to U.S. Geological Survey. Prices are not included in the listing below as they are subject to change. Current prices can be obtained by writing to the USGS, Books and Open File Reports. Prices include cost of domestic surface transportation. For transmittal outside the U.S.A. (except to Canada and Mexico) a surcharge of 25 percent of the net bill should be included to cover surface transportation. When ordering any of these publications, please give the title, book number, chapter number, and "U.S. Geological Survey Techniques of Water-Resources Investigations."

TWI 1-D1. Water temperature-influential factors, field measurement, and data presentation, by H.H. Stevens, Jr., J.F. Ficke, and G.F. Smoot. 1975.65 pages.
TWI 1-D2. Guidelines for collection and field analysis of ground-water samples for selected unstable constituents, by W.W. Wood. 1976. 24 pages.

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TWI 3-B3. Type curves for selected problems of flow to wells in confined aquifers, by J.E. Reed. 1980. 106 p .
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TWI 5-A2. Determination of minor elements in water by emission spectroscopy, by P.R. Barnett and E.C. Mallory, Jr. 1971. 31 pages.

TWI 5-A3. Methods for the determination of organic substances in water and fluvial sediments, edited by R.L. Wershaw, M.J. Fishman, R.R. Grabbe, and L.E. Lowe. 1987. 80 pages. This manual is a revision of "Methods for Analysis of Organic Substances in Water" by Donald F. Goerlitz and Eugene Brown, Book 5, Chapter A3, published in 1972.
TWI 5-A4. Methods for collection and analysis of aquatic biological and microbiological samples, edited by P.E. Greeson, T.A. Ehlke, G.A. Irwin, B.W. Lium, and K.C. Slack. 1977. 332 pages.

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## COMPUTER DISKETTE INFORMATION

This report contains a computer diskette (in the pocket at the back of the report) that has the source codes and data sets to be used with the report. The source codes were developed using the Microsoft Fortran Compiler, Version 3.3, with the DOS 2.0 operating system on an IBM PC/XT computer with the IBM 8088 Math Coprocessor and 256KB memory. Except for the OPEN statements near the beginning of the codes, Fortran 66 was used throughout to make the codes as machine independent as possible. For more information concerning the contents of the diskette, insert the diskette in drive A of the computer, type README, and push the enter key.

## METRIC CONVERSION FACTORS

For those readers who prefer to use metric units, conversion factors for terms used in this report are listed below:

| Multiply | By | To obtain |
| :---: | :---: | :---: |
| foot (ft) | 0.3048 | meter (m) |
| square foot ( $\mathbf{f t}^{2}$ ) | 0.09290 | square meter ( $\mathrm{m}^{2}$ ) |
| foot per day (ft/d) | 0.3048 | meter per day (m/d) |
| square foot per second ( $\mathrm{ft}^{2} / \mathrm{s}$ ) | 0.09290 | square meter per second ( $\mathrm{m}^{2} / \mathrm{s}$ ) |
| square foot per day (ft ${ }^{2} / \mathrm{d}$ ) | 0.09290 | square meter per day ( $\mathrm{m}^{2} / \mathrm{d}$ ) |
| gallon per day per foot (gal/d/ft) | 0.01242 | square meter per day ( $\mathrm{m}^{2} / \mathrm{d}$ ) |
| cubic foot per second ( $\mathrm{ft}^{3} / \mathrm{s}$ ) | 0.02832 | cubic meter per second (m ${ }^{3} / \mathrm{s}$ ) |
| cubic foot per day ( $\mathrm{ft}^{3} / \mathrm{d}$ ) | 0.02832 | cubic meter per day ( $\mathrm{m}^{2} / \mathrm{d}$ ) |
| gallon per minute (gal/min) | 0.00006309 | cubic meter per second ( $\mathrm{m}^{3} / \mathrm{s}$ ) |
| ounces per ton | 31.25 | grams per metric ton ${ }^{1}$ |

[^0]
# REGRESSION MODELING OF GROUND-WATER FLOW 

By Richard L. Cooley and Richard L. Naff

## 1 Introduction

### 1.1 Flow Equation and Boundary Conditions

The most general form of the ground-water flow equation that we consider here is given as

$$
\frac{\partial}{\partial x}\left(T_{x x} \frac{\partial h}{\partial x}\right)+\frac{\partial}{\partial y}\left(T_{y y} \frac{\partial h}{\partial y}\right)+R(H-h)+W
$$

$$
\begin{equation*}
+\sum_{\ell=1}^{N} \delta\left(x-a_{\ell}\right) \delta\left(y-b_{\ell}\right) Q_{\ell}=S \frac{\partial h}{\partial t} \tag{1.1-1}
\end{equation*}
$$

where

$$
\begin{aligned}
& T_{\xi \xi}(x, y)=\text { transmissivity }\left(K_{\xi \xi} b\right) \text { in the } \xi=x \\
& \text { or } y \text { direction; } \\
& K_{\xi \xi}(x, y)=\text { hydraulic conductivity of the } \\
& \text { aquifer in the } \xi \text { direction; } \\
& b(x, y)=\text { thickness of the aquifer; } \\
& R(x, y)=\text { hydraulic conductance (hydraulic } \\
& \text { conductivity divided by thick- } \\
& \text { ness) of sediments underlying a } \\
& \text { stream or of an aquitard underly- } \\
& \text { ing or overlying the aquifer; } \\
& W(x, y, t)=\text { source-sink term (positive for a } \\
& \text { source), distributed areally; } \\
& \sum_{\rho=1}^{N} \delta\left(x-a_{p} \delta \delta\left(y-b_{\rho}\right) Q_{\rho}=\right.\text { Dirac delta designation for } \\
& N \text { wells, each one pumping at rate } \\
& Q_{\ell}(t) \text { (positive for injection) and } \\
& \text { located at ( } a_{p}, b_{p} \text { ); } \\
& S(x, y)=\text { storage coefficient; } \\
& h(x, y, t)=\text { hydraulic head in the aquifer; }
\end{aligned}
$$

$H(x, y, t)=$ head at the stream bottom or at the distal side of the aquitard;
$x, y=$ Cartesian coordinates;
$t=$ time;
and $T_{x x}$ and $T_{y y}$ are continuous functions of $x$ and $y$.
With suitable internal boundary conditions, the region can be zoned with respect to $T_{\xi \xi}$ : Such boundary conditions involve head and specific discharge multiplied by thickness normal to the boundary $\left(q_{n}\right)$ and can be stated for a boundary between $T_{\xi \xi}$ zones $k$ and $\ell$ as

$$
\begin{gather*}
(h)_{k}=(h)_{\ell}  \tag{1.1-2}\\
\left(q_{n}\right)_{k}=\left(q_{n}\right)_{\ell} \tag{1.1-3}
\end{gather*}
$$

where $(\cdot)_{k}$ indicates that the quantity in parentheses is evaluated just within the $k$ side of the boundary and similarly for $\ell$. Zonation with respect to $R, W$, or $S$ requires no internal boundary conditions.

External boundary conditions applying at the periphery of the domain being modeled are given as

$$
\begin{equation*}
T_{x x} \frac{\partial h}{\partial x} n_{x}+T_{y y} \frac{\partial h}{\partial y} n_{y}+\alpha h=\beta \tag{1.1-4}
\end{equation*}
$$

where $\alpha(x, y, t)$ and $\beta(x, y, t)$ are given functions, and $\left\{n_{x}(x, y), n_{y}(x, y)\right\}$ is the outward-pointing unit normal at the boundary. The sum of the first two terms is the flux $q_{B}$ normal to the boundary (positive for outflow). Equation 1.1-4
incorporates the standard boundary conditions of specified flux ( $q_{B}$ ) and specified head ( $h_{B}$ ) but also allows for linear combinations to be given.

### 1.2 Types of Solutions

### 1.2.1 Direct Solution for Head

The classical problem of mathematical physics (and, by assumption, of ground-water hydrology) is to directly solve equations 1.1-1 through 1.1-4 for $h=h(x, y, t)$. Given that any specific problem is properly posed, such a solution will always exist. The conditions for properly posing a problem are the following.

1. The positions of all internal boundaries are known exactly. Examples of internal boundaries are abrupt changes in $T_{\xi \xi}, R$, $S$, or $W$; internal known flux $\left(q_{n}\right)$ boundaries; and internal known head boundaries. Note that a river is often treated as either an internal known head boundary where the river is assumed to have no width, or a zone of differing $R$ where each bank is a zone boundary.
2. The positions and types of all external boundaries are known exactly. External boundaries frequently are known flux $\left(q_{B}\right)$ types or known head ( $h_{B}$ ) types. Sometimes some linear combination is known.
3. Hydrogeologic variables $T_{\xi \xi}, R$, and $S$ and hydrologic variables $W^{W}$ and $Q_{\ell}$ are known at all points in the region.
4. All boundary-condition variables $H, \alpha$, and $\beta$ are known. The initial head (at $t=0$ ) is a boundary condition and must also be known.
Obviously, ground-water flow problems are not actually of the classical type because none of the conditions cited above ever are met exactly. Conditions 1 and 2 are often most closely fulfilled, but estimates (often crude) usually must suffice for the variables in conditions 3 and 4. Any errors in these input variables are propagated directly into the solution. However, reasonable (but incorrect) estimates of the variables can be shown to yield errors in predicted $h(x, y, t)$ that have the characteristic of being bounded (that is, they do not tend to plus or minus infinity). Also, as the errors in the input
variables tend to zero, the errors in computed head do also.

### 1.2.2 Inverse Solution for Parameters

An inverse solution involves solving equations 1.1-1 through 1.1-4 for one or more of the variables $T_{\xi \xi}, R, S, W, Q_{q}, \alpha$, or $\beta$, over the region; these variables are termed parameters here. Because $R, S, W, Q_{f}, \alpha$, and $\beta$ are not involved in derivatives, theoretically they may be solved for algebraically. Unless $T_{\xi \xi}$ is constant, it is involved in derivatives and, thus, must be obtained by solving a differential equation. To understand this, note that equation 1.1-1 may be rearranged to give

$$
\begin{equation*}
a \frac{\partial T_{x x}}{\partial x}+b \frac{\partial T_{y y}}{\partial y}+c T_{x x}+d T_{y y}-F=0 \tag{1.2-1}
\end{equation*}
$$

where

$$
\begin{gathered}
a(x, y, t)=\frac{\partial h}{\partial x}, \quad b(x, y, t)=\frac{\partial h}{\partial y}, \\
c(x, y, t)=\frac{\partial^{2} h}{\partial x^{2}}, \quad d(x, y, t)=\frac{\partial^{2} h}{\partial y^{2}}, \text { and } \\
F(x, y, t)=S \frac{\partial h}{\partial t}-R(H-h)-W-\sum_{f=1}^{N} \delta\left(x-a_{\ell}\right) \delta\left(y-b_{l}\right) Q_{q} .
\end{gathered}
$$

In general, if $T_{x x}$ and $T_{y y}$ are known, conditions for finding $R, S, W, Q_{\hat{\beta}} \alpha$, or $\beta$ are:

1. Conditions 1 and 2 for the classical direct solution are met.
2. Head distribution $h(x, y, t)$ is known exactly.
3. The solution for the desired combination of parameters to be obtained is unique. The latter condition is completely problem dependent. Because solution involves only algebraic manipulations, the condition reduces to the requirement that the system of algebraic equations involving the desired parameters has a unique solution. Generally, solution involves picking the required number of points spatially and through time to yield the necessary number of equations.
To find $T_{x x}$ and $T_{y y}$, more conditions are
required than for finding $R, S, W, Q_{\beta} \alpha$, or $\beta$. These conditions are:
4. Conditions 1 and 2 for finding $R, S, W, Q_{p}$ $\alpha$, or $\beta$ must be met.
5. The direction of the velocity vector must be known everywhere, or $T_{x x} / T_{y y}$ must be known everywhere, or quantities $a, b, c$, $d$, and $F$ in equation 1.2-1 must be known at two (or more) points in time to give a unique solution to equation 1.2-1 written in the form of a pair of simultaneous differential equations. These requirements result because $1.2-1$ is one equation in two unknowns. Hence, an additional relationship is required. If the velocity direction is known everywhere, then by employing Darcy's law the additional relationship is derived as

$$
\begin{equation*}
\frac{T_{x x}}{T_{y y}}=\frac{q_{x} b}{q_{y} a} \tag{1.2-2}
\end{equation*}
$$

where $a, b$, and $q_{x} / q_{y}$ (the ratio of the $x$ and $y$ direction fluxes) are known.
3. If either the direction of the velocity vector or $T_{x x} / T_{y y}$ is known, then either $T_{x x}$ or $T_{y y}$ must be known on a possibly discontinuous curve crossing all flowlines. If solution is to be obtained by solving a simultaneous pair of differential equations, then $T_{x x}$ must be known on a possibly discontinuous curve that spans the range of $y$, and $T_{y y}$ must be known on a possibly discontinuous curve that spans the range of $x$. These are extensions of the Cauchy boundary condition for a firstorder differential equation involving a single dependent variable and are required for solution of the problem.
4. The function $F$ in equation $1.2-1$ must be known everywhere. This means that all quantities in $F$ must be known or that a mathematical form for $F$ can be assumed.
Ground-water flow modeling does not fit into the category of inverse solutions, although a significant part of most model studies is to find values of the parameters that allow values of calculated head to match those observed in the field. The difficulty is that the required conditions are almost never met. Head distribution is never known exactly because measurements
do not exist at all points and, where these measurements do exist, they are not exact. Furthermore, some measure of $T_{\xi \xi}$ is virtually never available on the required curves, and information on directions of flow vectors for even scattered locations usually is nonexistent. Assumptions concerning zonations in which $T_{x x} / T_{y y}$ and (or) $T_{x x}$ and $T_{y y}$ may be considered constant simplify the problem, but the fact that $h$ must be known still remains.

Because the head distribution is not known exactly, coefficients $a, b, c, d$, and $F$ in equation 1.2-1 are in error. Furthermore, head appears as a derivative in all of these quantities. Hence, any error in $h$ is propagated into the inverse solution as a derivative of error. The effects of this propagation are often disastrous because, if $\epsilon_{h}$ is defined as error in head, $\epsilon_{h} \rightarrow 0$ does not imply that $\partial \epsilon_{h} / \partial \xi \rightarrow 0$. Also, $\left|\partial \epsilon_{h} / \partial \xi\right| \gg\left|\epsilon_{h}\right|$ is common, and it can happen that $\left|\partial \epsilon_{h} / \partial \xi\right| \rightarrow \infty$ even if $\epsilon_{h}$ is bounded. Therefore, the error in computed $T_{\xi \xi}$ (or other parameter) may not approach zero as $\epsilon_{h} \rightarrow 0$, and may, in fact, be quite large (Neuman, 1980, p. 342-344).

### 1.2.3 Solution Using Real Data

In the previous section, we argued that problems involving ground-water flow modeling of real field systems are neither of the classical nor inverse type, because the data necessary for the problems to be classified as either type are usually lacking. An estimate of the hydraulic head distribution based on measurements (that are in error with respect to the model) taken at selected points usually exist. Estimates of the parameters are usually either completely unknown or have been obtained by spot measurements, few of which are directly useful for construction of appropriate effective values for use in equation 1.1-1. That modeling problems in ground-water hydrology involve an incomplete combination of several types of data in which error and error propagation are important considerations is evident.

### 1.3 Sources of Error in GroundWater Data

Uncertainty (or errors) in ground-water data may have many sources, and enumeration of all
possible sources would be a nearly impossible task. However, a consideration of some of the more important sources of error serves to illustrate the importance of the error component.

### 1.3.1 Sources of Error in Head Data

Some major potential sources of randomappearing error in head data with respect to the model (equations 1.1-1 through 1.1-4) are:

1. Areal ground-water models assume that the head used is the average over the vertical. However, wells may not be open over the entire interval modeled, and if they are, they may not measure the average. Flow into and (or) out of a well distorts the hydraulic head field in the vicinity of the well so that the recorded water level does not represent the average head.
2. Permeability varies from point to point, which causes water levels to vary from values they would have if permeability were uniform. However, models usually do not take this detailed variation into account. This phenomenon has been extensively studied during the last 10 years, and literature reviews are contained in the works of Dagan (1986) and Gelhar (1984, 1986).
3. Water levels measured in wells in use may contain unknown amounts of residual drawdown. In addition, unused wells may be near wells that are in use, with resulting unknown drawdown in the unused well.
4. Measurement of well-head elevation may be in error.
Actual total error from the above sources is highly problem dependent, but it is easy to imagine errors of several feet. It should be noted that measurement error in water levels was not mentioned as a major source of error because it commonly amounts to one- or two-tenths of one foot or less. Finally, major model error in equations 1.1-1 through 1.1-4 (for example, head dependence in one or more parameters or three-dimensional flow) was also not mentioned because error resulting from this source is bias and should be detected and eliminated by analysis of model results.

### 1.3.2 Sources of Error in Parameter Data

Because there are several different parameters to be considered, and each can be estimated or measured in several different ways, a large number of sources of error exist in parameter data. Model error is not considered here, but other types of bias are potentially important and are often difficult to detect. Some examples of errors in parameter data illustrating the nature of the problem are:

1. Too few estimates of parameters are available to compute stable estimates of statistics, such as mean and variance.
2. Results of point sampling are often biased because a large amount of data does not necessarily allow computation of nearly true or effective values of a parameter and its variance. For example, permeability values from core analyses often are not representative of regional values, because flow through large fractures is not reproduced by core analyses. Also, effective values of a parameter and its variability are usually not directly given by standard mean and variance formulas.
3. Transmissivities estimated from specificcapacity data collected by drillers are subject to numerous sources of error. Common sources include (1) mismeasuring water levels or pumping rates, (2) allowing the water level to recover after bailing, (3) clogging of the slots or screen, and (4) inaccurate reporting. There are so many sources of error that the errors may often appear to be random. A persistent source of bias results because drillers drill wells in favorable locations and only screen (or slot) the most productive zones.
4. Transmissivities and storage coefficients estimated from pumping-test analysis are subject to many of the same errors as above, but the more carefully controlled tests should reduce their frequency and magnitude. In addition, a single test may not be representative of an entire hydrostratigraphic unit.
5. Transmissivities estimated from lithological data are usually biased to an unknown degree.

### 1.4 Model Construction

Ground-water models are constructed by using the types of data alluded to in the previous section. Hence, measured or estimated parameter data, either reliable or complete enough to employ directly in a model to reproduce measured head data with an acceptable model fit, are rare. As a result, adjustment of parameter values, and sometimes basic model structure, is used to improve model fit. Two basic groups of methods currently in use to accomplish this are: (1) trial and error procedures and (2) optimization methods that minimize a formal objective function.

### 1.4.1 Trial and Error Methods

Trial and error is the method of repeated simulation until the calculated head distribution obtained with a reasonable set of parameters fits closely enough to satisfy the analyst. Sometimes an objective measure of goodness of fit, such as $\Sigma\left(h^{\text {calc }}-h^{o b s}\right)^{2}$, is used to aid the analyst in deciding whether or not a change in parameters (or model structure) has improved the overall model fit. However, no matter how the method is applied, it has several inherent critical deficiencies:

1. No methodology exists to guarantee that the simulations will proceed in a direction that could lead to the best set of parameters.
2. Determining when that best set has been reached is difficult.
3. No practical way of determining how many other sets of parameters could yield similar correspondence between $h^{c a l c}$ and $h^{o b s}$ exists.
4. Deciding whether or not additional parameters or a more refined model would significantly improve model fit is difficult.
5. No way of quantitatively assessing the predictive reliability of the model exists.
A method of model construction that addresses these deficiencies would allow construction and use of a model with a much greater degree of confidence than that provided by trial and error methods. Hence, attention is turned to formal optimization procedures.

### 1.4.2 Formal Optimization Procedures

Optimization procedures utilize a formal criterion of goodness of fit, often called an objective function. This function is minimized (or sometimes maximized, depending on the form of the function) with respect to the parameters to yield an optimum or best-fit solution. Minimization (or maximization) sometimes is subject to certain other criteria regarding values that the parameters, or pertinent functions of the parameters related to the model, may take on. These criteria are called constraints.
Examples of objective functions are:

$$
\begin{aligned}
& \sum_{\ell=1}^{n_{s}}\left(h_{\tau} \hat{h}_{\ell}\right)^{2}, \\
& \sum_{\ell=1}^{n_{s}}\left|h_{\ell} \hat{h}_{\ell}\right|, \\
& \max _{\ell}\left|h_{\ell} \hat{h}_{\ell}\right|, \text { and } \\
& \sum_{\ell=1}^{n_{S}} w_{\ell}\left(h_{\ell} \hat{h}_{\ell}\right)^{2}+\sum_{m=1}^{n_{p}} k_{m}\left(p_{m}-\left.\hat{p}_{m}\right|^{2}\right.
\end{aligned}
$$

where
$h_{\rho}=$ observed head,
$\hat{h}_{\ell}=$ calculated head,
$p_{m}=$ observed or prior estimate of a parameter,
$\hat{p}_{m}=$ calculated parameter value that, when used in the model, produces $\hat{h}_{\ell}$,
$w_{\ell}=$ weight related to the reliability of the observation $h_{p}$,
$k_{m}=$ similar weight applied to $p_{m}$,
$n_{s}=$ number of observations of head, and
$n_{p}=$ number of observations of parameters.
The last example is called a compound objective function because it contains both head and parameters explicitly. Note that minimization of each of the functions with respect to the parameters of the model produces a solution that is overall a best fit to the data, according to the objective function. If the signs of the
functions were changed, maximization would produce the same result.

Examples of constraints are:

$$
\begin{aligned}
& p_{m}^{L}<\hat{p}_{m}<p_{m}^{U} \\
& a \hat{p}_{k}+b \hat{p}_{m}+c \hat{p}_{n}=f \\
& a \hat{p}_{k}+b \hat{p}_{m}+c \hat{p}_{n}<f
\end{aligned}
$$

where $a, b, c$, and $f$ are constants or known functions; superscript $L$ refers to a lower limit; and superscript $U$ refers to an upper limit. The bestfit solution obtained by minimizing (or maximizing) the appropriate objective function must simultaneously satisfy the appropriate constraints.

Because the solution obtained by an optimization procedure has known properties, it may be analyzed. The exact procedures used and the extent to which the model may be analyzed depend on the type of optimization method selected for use. Statistical regression procedures handle, on a probabilistic basis, the propagation of data errors (with respect to the model) into the estimates of parameters and predictive capability of the model. Methods have been developed for estimating parameters, testing assumptions made during development
of techniques, testing model fit, determining the reliability and significance of the model and the parameters contained in it, effecting corrective measures for violation of some assumptions, and estimating the reliability of predictions to be made with the model. These procedures and the statistical background necessary to apply them are detailed in the remainder of the text.

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## Additional Reading

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## 2 Review of Probability and Statistics

Casual observation of our environment indicates that many phenomena are not strictly predictable. We cannot, for instance, exactly say what the maximum air temperature at any particular location will be tomorrow, although we might be able to give a probable range. This probable range might be based on our past experience, which would enable us to say that tomorrow's high, considering the location and season, will probably fall within a specified interval. A more sophisticated forecasting model may enable us to reduce the range within which we think tomorrow's high will fall, but random elements in the forecasting procedure would preclude giving an exact answer. As another example of randomness, consider the toss of a coin. Prior to the toss, we can only give the possible outcomes, either a head or a tail, and, if the coin is fair, say that either have equal likelihood of occurring. However, this ability to state precisely that any future outcome of this experiment can, with equal probability, result in either a head or a tail is an important advantage over that offered for predicting tomorrow's maximum temperature. In this latter case, because of the complex nature of the processes resulting in tomorrow's maximum, the likelihood that we could give a precise statement concerning the probability that it will fall in our predicted interval is remote. Instead of attempting to untangle these complexities, we might opt to study the history of maximum temperatures at the location and annual date in question. By assuming that this history will extrapolate into the future (that is, that weather dynamics in future years will remain essentially unchanged from those in previous years), we could give an estimate of the likelihood that tomorrow's maximum will fall in a particular interval. However, tools need to be developed to carry out this investigation.

### 2.1 Basic Concepts

Randomness itself can be considered to be centered around an experiment; the outcome of the experiment will have a random quality
attached to it. For example, in a coin-toss experiment, the outcome is dominated by the random element (either a head or a tail). On the other hand, many experiments have a large deterministic factor. For example, in a chemical titration experiment we measure the unknown and the amount of titrant used, then calculate the amount of a specific substance in the unknown. However, measurement error creeps into our technique, and results vary from realization to realization of the experiment. Some experiments, such as annual, peak river flows, are not ours to perform but only to observe. This experiment is an example of an event in nature that has a large random component which nature provides. As we attempt to measure these flows, we introduce additional randomness, which we generally ignore. Hydraulic conductivities measured from core samples are similar to peak flows; nature has already provided for randomness, which is constrained by certain deterministic factors, such as type of source material, distance of transport, climate, and diagenesis. Again, for every realization of this experiment, measurement error is introduced, which may not be small.

All possible outcomes of an experiment are known as its sample space. The sample space of a coin-toss experiment consists of either a head (H) or tail (T):

$$
S=\{H, T\} .
$$

If the experiment consists of the toss of two coins, then the sample space consists of

$$
S_{1}=\{(H, H),(H, T),(T, H),(T, T)\} .
$$

On the other hand, if we are only interested in the total number of heads which might result from a single toss of two coins, we could define the experiment as this sum, which would result in the sample space

$$
S_{2}=\{0,1,2\} .
$$

In the case of $S_{1}$, every member of the sample space is equally likely to occur, whereas for $S_{2}$, a one is twice as likely to occur as either zero or two, provided that the coin is fair.

The sample space for a hydraulic-conductivity experiment could be defined as all positive real
numbers; that is, measurements from cores might result in values (outcomes) which could be as small as zero or, if we stretch our imaginations, infinitely large. This space could be considered to be a continuous equivalent of the $S_{2}$ space for the two-coin experiment. That is, a porous medium is an extremely complex random process itself. By conducting hydraulicconductivity measurements on cores, we quantify this randomness in much the same way that counting heads quantifies an outcome of the two-coin experiment. However, by quantifying the randomness of the porous medium in this manner, we have never investigated the possible existence of more basic, perhaps nonnumeric sample spaces similar to $S_{1}$ of the two-coin experiment for a porous medium. Even if we were to discover the existence of such a space, we would then need to find a rule, or algorithm, which would allow us to connect the two spaces. We shall not worry about the possibilities of an $S_{1}$-like space for many processes; however, when they are available, they provide an excellent mechanism for investigating the characteristics of $S_{2}$-like spaces.
An event is defined as any subset of the sample space. The investigator is usually interested in the relative frequency of occurrence of an event. In the case of the $S_{2}$ space and the two-coin experiment, it is apparent that half the time a realization experiment should result in a one. This event is equivalent to the event in the $S_{1}$ space corresponding to the union of ( $H, T$ ) and ( $T, H$ ), which occurs with a relative frequency of one-half. Thus, the relative frequency of a head occurrence for the two-coin experiment is not dependent upon the definition of the sample space, but on the basic randomness controlling the experiment.
The investigator is frequently confronted with the problem of needing a numerical result for the outcome of a random, but not necessarily numerical, experiment. In the case of coin-toss experiments, the basic outcome is seen to be a particular arrangement of heads and (or) tails. By assigning a head a value of one and a tail a value of zero and then summing, it is possible to translate these basic results into something measurable. This process of assigning a numerical value to a nonnumerical outcome leads to the definition of a random variable.

Definition: A random variable is a function whose value is a real number determined by each element in a sample space.

When the outcome of the experiment is numerical, then this result can be considered to be the random variable (this statement is merely a special case of the above definition). From the above definition, we see that a mathematical transformation of a random variable is also a random variable. (Throughout this review, a random variable is indicated by an upper case English or Greek letter, whereas a value that it may take on is indicated by another letter, usually lower case of the same type as used for the random variable.)
The concepts of a random experiment, sample space, and random variable are flexible. For instance, if in the case of the toss of two coins, the experiment is defined as the total number of heads appearing, then the $S_{2}$ sample space is an automatic result, and the random variable can also be considered to be this result. However, if the experiment is defined to be the arrangement of heads and (or) tails resulting from a toss (that is, the $S_{1}$ space), then the same effect can be obtained by letting the random variable over the $S_{1}$ space be a function that assigns a one to a head and a zero to a tail and then sums the result. The investigator usually defines the sample space, or experiment, to suit a particular objective. As a matter of convenience, the space is usually selected such that the relative frequencies of occurrence of events within the space are definable. Access to such basic sample spaces as $S_{1}$ for the two-coin experiment allow for the calculation of relative frequencies for events in both $S_{1}$ and $S_{2}$. Without the existence of a space like $S_{1}$, determining the true relative frequency of occurrence for an event in $S_{2}$ is difficult, if not impossible. This situation is also evident from the hydraulicconductivity experiment, where only an $S_{2}$-like sample space is available to the investigator.

A random variable can also be described as either being discrete, as in the coin-toss experiment, or continuous, as represented by the hydraulic-conductivity experiment. A discrete random variable is defined over a sample space whose elements are discrete, although there may be as many as there are whole numbers
(mathematicians refer to this phenomenon as being countably infinite). A continuous random variable is defined over a continuous sample space whose elements are infinite in number (therefore these elements are uncountably infinite).

### 2.2 Frequencies and Distributions

### 2.2.1 Discrete Random Variables

Although frequencies of occurrence are usually associated with events in a sample space, they are also associated with values of random variables, since random variables are functions of the elements in a sample space. That is, particular values of a random variable correspond to particular events in the sample space and, therefore, have frequencies of occurrence. Even though we will speak of the relative frequency of occurrence for particular values of a random variable, we are, in reality, speaking of a corresponding event in the sample space. In fact, we frequently use a range of values of a random variable to define an event in a sample space, thus avoiding the task of describing which elements of the sample space compose the event.

Frequencies of occurrence for events in many discrete sample spaces can be deduced from the following axiomatic premise: If an experiment can result in any one of $N$ different equally likely outcomes, and if exactly $n$ of these outcomes correspond to event $A$, then the relative frequency of occurrence of $A$ is $n / N$. As a simple example of employment of this premise, consider an experiment consisting of a toss of a die. The sample space consists of the integers 1 through 6 and, for any realization of the experiment, each element of the sample space has equal likelihood of occurrence. By considering each element of the sample space to be an event, one can calculate the frequency of occurrence, $f\left(x_{i}\right)$, with which a random variable takes on the value $x_{i}$. For this experiment, only the integer values 1 through 6 of $x_{i}$ have frequencies of occurrence other than $0 ; f x_{i}$ ) can be graphically represented as shown in figure $2.2-1$. In this case, $f\left(x_{i}\right)$ is referred to as the discrete density function of the discrete random variable


Figure 2.2-1
consisting of the outcome of a toss of a single die.
When two dice are cast, the experiment can be defined either as the sum that results from the toss or simply as all possible arrangements that could appear on the dice. If the sum is chosen, then the sample space consists of the integers 2-12, which would also be the range of values that the random variable could take on. The elements of this space, however, are not equally likely to occur. The sample space consisting of all arrangements of the numbers appearing on the two dice, presented graphically in table 2.2-1, has elements which are equally likely to occur.

Table 2.2-1

| Second <br> die | 1 | 2 | 3 | 4 | 5 | 6 |
| :---: | :---: | :---: | :---: | :---: | :---: | :---: |
|  | $(1,1)$ | $(2,1)$ | $(3,1)$ | $(4,1)$ | $(5,1)$ | $(6,1)$ |
|  | $(1,2)$ | $(2,2)$ | $(3,2)$ | $(4,2)$ | $(5,2)$ | $(6,2)$ |
|  | $(1,3)$ | $(2,3)$ | $(3,3)$ | $(4,3)$ | $(5,3)$ | $(6,3)$ |
|  | $(1,4)$ | $(2,4)$ | $(3,4)$ | $(4,4)$ | $(5,4)$ | $(6,4)$ |
|  | $(1,5)$ | $(2,5)$ | $(3,5)$ | $(4,5)$ | $(5,5)$ | $(6,5)$ |
|  | $(1,6)$ | $(2,6)$ | $(3,6)$ | $(4,6)$ | $(5,6)$ | $(6,6)$ |

The relative frequency of occurrence of an event corresponding to any subset of elements in this space can be calculated by using the premise concerning equally likely outcomes.

A random variable, consisting of the sum that results from any outcome of the two dice experiment, takes on the integer values 2-12 over the sample space represented by table $2.2-1$. The discrete density function for this random variable can now be derived from the basic premise concerning outcomes that are equally likely, since each value for this discrete random variable corresponds to a particular event consisting of a particular subset of elements in the
sample space indicated by table $2.2-1$. Thus, the value of $x_{i}=3$ corresponds to the event containing the elements $(2,1)$ and ( 1,2 ) and has a relative frequence of occurrence of $2 / 36$. Letting $x_{i}$ represent the integer values that this random variable can obtain, its density function, $f\left(x_{i}\right)$, can be represented as shown in figure 2.2-2.

Note that had the first definition of the experiment been used, then every element of the sample space consisting of the integers 2-12 would have frequencies of occurrence, when considering each element as an event, equivalent to those shown in figure 2.2-2.

Frequencies of occurrence, or deduced frequencies of occurrence as indicated in figures 2.2-1 and 2.2-2, are indications of the future. We can make probability statements concerning the possibility of a random variable taking on future values from such knowledge. In a craps (two-dice) game, we know that the probability of rolling a natural, an outcome of 7 or 11 on the first cast, is $2 / 9$ simply because these values of the random variable for the two-dice experiment correspond to elements in the sample space which occur with a relative frequency of 2/9. Formally, the statement that this discrete random variable $X$ take on the values of 7 or 11 with a probability of $2 / 9$ is written

$$
P(X=7 \text { or } X=11)=2 / 9
$$

The probability that this random variable takes on any integer value between 2 and 12 is
obtainable directly from its frequency density, figure 2.2-2.

A probability statement that is frequently encountered concerns the probability that a random variable is less than or equal to a specific value. For the random variable corresponding to the sum of outcomes of the cast of two dice, we may ask, what is the probability that the random variable $X$ is less than or equal to 5 ? The probability of this event is equal to the probability that $X$ take on any integer value 2 through 5:

$$
P(X \leq 5)=P(X=2 \text { or } X=3 \text { or } X=4 \text { or } X=5) .
$$

This probability is the sum of the probabilities of the individual events that $X$ take on the integer values 2 through 5:

$$
P(X \leq 5)=1 / 36+2 / 36+3 / 36+4 / 36=5 / 18
$$

(If the student is not convinced of this relationship, he or she should examine the elements of the sample space represented by table 2.2-1 to ascertain that it holds.) Note that $P(X \leq 12)$ is unity; that is, an event which occurs with a probability of one will, undoubtedly, take place. A probability of zero indicates, on the other hand, that the event of concern cannot possibly occur.

The probability statement $P(X \leq a)$, where $a$ is any real number, is given a special definition for both discrete and continuous random variables. That is, $F(a)=P(X \leq a)$ is known as the


Figure 2.2-2
cumulative distribution function of the random variable $X$. For the case of the sum of outcomes for two dice, $F(a)$ appears as illustrated in figure 2.2-3. Because a random variable represents a functional mapping from the sample space to the real number space, we can be assured that the probability of the event $X \leq a$ exists and is equal to the sum of the probabilities of all events corresponding to values of the random variable which are less than or equal to $a$. In general, for discrete random variables, the cumulative distribution function can be evaluated by summing the appropriate relative frequencies of occurrence:

$$
\begin{equation*}
F(a)=\sum_{x_{i} \leq a} f\left(x_{i}\right) . \tag{2.2-1}
\end{equation*}
$$

The cumulative distribution function for all random variables, discrete or continuous, has the following properties:

1. $F(a)$ is a nondecreasing function of $a$,
2. $\lim F(a)=1$, $a \rightarrow \infty$
3. $\lim F(a)=0$.
$a \rightarrow-\infty$
These properties will be demonstrated in detail for continuous random variables in a later section. For a discrete random variable, these properties reflect the fact that, by definition, the discrete density function can never have a negative frequency of occurrence and that the sum of frequencies must equal one.
In the next section, an estimator for the density function of continuous random variables is developed, which will eventually allow us to explore the nature of density and cumulative distribution functions of continuous random variables.

## Problem 2.2-1

An urn contains one red, one white, and two blue balls, all of equal dimensions. A ball is


Figure 2.2-3
drawn from the urn, replaced, and then another draw is made.
a. What possible arrangements (red, white, and (or) blue) of the two balls, considering order of selection, could occur (see, for example, table 2.2-1 for two dice)?
b. What is the frequency of occurrence of any of the above events? (Hint: let the balls be represented by the symbols $R, W$, $B 1$, and B2.)
c. A value of one is assigned to a blue ball, two to a red ball, and three to a white. A random variable consists of the sum of any outcome consisting of two draws with replacement. Develop a discrete density function for this random variable.
d. What is the probability that this random variable takes on a value of 4? What arrangements of balls correspond to this value of the random variable?

### 2.2.2 Histograms

In many cases, we do not have access to all values of random variables in a sample space (in particular, for many continuous random variables). We sample the population consisting of all possible values of the random variable and hope to draw inferences from this sample. The inferences we draw are usually in the form of statistics, which we refer to as sample statistics. We like to think that sample statistics estimate values of population parameters, which are constants reflecting the true frequency distribution of the random variable.- This is frequently the case if the observations composing samples are made randomly and without bias. Samples composed of such observations are referred to as random samples and are expected to be representative of the population.

Estimates of density functions for random variables are frequently made from random. samples. Although certain experiments, such as a coin toss, allow for the deduction of frequencies of occurrence of events, other experiments defy a theoretical calculation, forcing us to estimate from a random sample. These estimates, known as histograms, are generally constructed by repeating the experiment a large number of times (thus, sampling the population
of all possible outcomes), dividing the range of these outcomes into class intervals, and calculating the relative number of points that fall in each interval. We might imagine, for example, that we could watch a craps game and note the outcome of each roll of dice. After a thousand rolls, we would calculate the relative percentage of each integer, 2-12, which occurred. If these sample frequencies of occurrence were not close to that shown previously for the theoretical result, we would suspect that the dice had been tampered with.

As an example of a histogram constructed from observed values of a continuous random variable, consider the transmissivity data shown in table 2.2-2. Figure 2.2-4 represents a histogram constructed directly from these data, which constitute a random sample from the population of transmissivities as determined from specific capacities of wells in carbonate rocks of central Pennsylvania. A second histogram, figure 2.2-5, was constructed from a logarithmic transformation of these data as shown in table 2.2-3. The first histogram was constructed by using a class interval of 50,000 $\mathrm{gal} / \mathrm{d} / \mathrm{ft}$, and the second is based upon an interval of one-half a $\log _{10}$ cycle. The first histogram is not very illustrative because most of the wells have transmissivities less than $50,000 \mathrm{gal} / \mathrm{d} / \mathrm{ft}$ (the underlying population frequency is probably heavily skewed to the right). By logarithmically transforming of the random variable, we scale the abscissa so as to remove the skewness in the histogram, causing it to be more bell shaped. This type of transformation is frequently used on random variables that have a zero lower bound, causing the transformed variable to have tails that tend to infinity in both directions. The transformation also tends to remove any right skewness in the frequency distribution of these random variables. With regard to the transformed variate, the histogram in figure 2.2-5 suggests a bell-shaped population frequency distribution. More data and smaller class intervals, as suggested in the following paragraphs, should cause the histogram shown in figure 2.2-5 to approach its population shape, which we may suspect to be a normal distribution; the untransformed random variable would then result from a log-normal distribution.

Toble 2.2-2
[From Siddiqui (1969, p. 433-436)]

| Transmissivity $\mathrm{gal} / \mathrm{d} / \mathrm{ft}$ | $\log _{10} T$ | Trangmisgivity $\mathrm{ga} / \mathrm{d} / \mathrm{ft}$ | $\log _{10} T$ |
| :---: | :---: | :---: | :---: |
| 15.0 | 1.176 | 2,370.0 | 3.375 |
| 18.0 | 1.255 | 2,440.0 | 3.387 |
| 21.0 | 1.322 | 2,540.0 | 3.405 |
| 29.0 | 1.462 | 2,800.0 | 3.447 |
| 32.0 | 1.505 | 2,820.0 | 3.450 |
| 35.0 | 1.544 | 3,380.0 | 3.529 |
| 50.0 | 1.699 | 4,410.0 | 3.644 |
| 52.0 | 1.716 | 4,520.0 | 3.655 |
| 56.0 | 1.748 | 5,500.0 | 3.740 |
| 62.0 | 1.792 | 5,650.0 | 3.752 |
| 84.0 | 1.924 | 6,030.0 | 3.780 |
| 92.0 | 1.964 | 6,240.0 | 3.795 |
| 106.0 | 2.025 | 6,340.0 | 3.802 |
| 118.0 | 2.072 | 7,290.0 | 3.863 |
| 142.0 | 2.152 | 8,130.0 | 3.910 |
| 160.0 | 2.204 | 11,000.0 | 4.041 |
| 175.0 | 2.243 | 13,100.0 | 4.117 |
| 184.0 | 2.265 | 13,700.0 | 4.137 |
| 202.0 | 2.305 | 14,500.0 | 4.161 |
| 264.0 | 2.422 | 17,200.0 | 4.236 |
| 354.0 | 2.549 | 17,700.0 | 4.248 |
| 370.0 | 2.568 | 19,700.0 | 4.294 |
| 374.0 | 2.573 | 23,100.0 | 4.364 |
| 455.0 | 2.658 | 24,200.0 | 4.384 |
| 463.0 | 2.666 | 26,400.0 | 4.422 |
| 515.0 | 2.712 | 33,400.0 | 4.524 |
| 528.0 | 2.723 | 34,700.0 | 4.540 |
| 615.0 | 2.789 | 42,400.0 | 4.627 |
| 705.0 | 2.848 | 46,300.0 | 4.666 |
| 753.0 | 2.877 | 52,000.0 | 4.716 |
| 800.0 | 2.903 | 66,500.0 | 4.823 |
| 984.0 | 2.993 | 68,400.0 | 4.835 |
| 1,150.0 | 3.059 | 132,000.0 | 5.121 |
| 1,290.0 | 3.111 | 152,000.0 | 5.182 |
| 1,500.0 | 3.176 | 423,000.0 | 5.626 |
| 1,580.0 | 3.199 | 423,000.0 | 5.626 |
| 1,670.0 | 3.223 | 528,000.0 | 5.723 |
| 1,850.0 | 3.267 | 528,000.0 | 5.723 |
| 2,310.0 | 3.364 | 528,000.0 | 5.723 |

Table 2.2-3

| Class <br> interval | Number of <br> occurrences | Relative <br> frequency | Cumulative <br> frequency |
| :--- | :---: | :---: | ---: |
| $1.0-1.5$ | 4 | 0.051 | 0.051 |
| $1.5-2.0$ | 8 | .103 | .154 |
| $2.0-2.5$ | 8 | .103 | .257 |
| $2.5-3.0$ | 12 | .154 | .411 |
| $3.0-3.5$ | 12 | .154 | .565 |
| $3.5-4.0$ | 10 | .128 | .693 |
| $4.0-4.5$ | 10 | .128 | .821 |
| $4.5-5.0$ | 7 | .090 | .911 |
| $5.0-5.5$ | 2 | .026 | .937 |
| $5.5-6.0$ |  | 5 | .064 |
|  |  |  | 1.001 |

[^1]TECHNIQUES OF WATER-RESOURCES INVESTIGATIONS


Figure 2.2-4


Figure 2.2-5

We are now in position to estimate the probability of occurrence of an event associated with the log-transformed random variable. Let $X=\log _{10} T$ represent the transformed variate plotted in figure $2.2-5$, and assume that the histogram for $X$ is representative of the population frequency. The estimated probability that $X$ is less than or equal to 5.0 but greater than $4.5, P(4.5<X \leq 5.0)$, then, is the sample frequency of occurrence of this event (equal to 0.09 ). The probability that $X$ is less than or equal to 5.0 can be estimated by summing the frequencies of occurrence of all events smaller than 5.0 ; thus $P(X \leq 5.0) \approx 0.911$. Thus, the chances are about 91 in 100 that the transmissivity of the carbonate rocks in central Pennsylvania, as determined by any random well, will be less than or equal to $1 \times 10^{5} \mathrm{gal} / \mathrm{d} / \mathrm{ft}$. The reader should realize that these results are only approximate, as the histogram is an approximation of the true population frequency distribution.
An estimate of the cumulative distribution function can also be constructed from a random sample. Let $F_{n}(a)$ represent this estimate, known as the sample distribution function; an appropriate estimator for $F_{n}(a)$ is the sum of all estimated relative frequencies for values of the random variable $X$ less than $a$ :

$$
\begin{equation*}
P(X \leq a) \approx F_{n}(a)=\sum_{i \leq a / \Delta x} f_{i}^{*} \tag{2.2-2}
\end{equation*}
$$

where
$f_{i}^{*}=n_{i} / n=$ sample frequency of occurrence of an event represented by the $i$ th class interval,
$\Delta x=$ size of class interval,
$n_{i}=$ number of outcomes having values in interval $i$, and $n=$ size of random sample.

An application of this procedure for the logarithmic transformation of transmissivity is shown in figure 2.2-6.

### 2.2.3 Continuous Random Variables

The definition of frequency $f_{i}^{*}$ used in equation 2.2-2 suffers from the deficit that it is dependent upon the size of the class interval; that is, if $\Delta x$ decreases in size while $n$ remains constant, then $f_{i}^{*}$ must also decrease, as we are also effectively decreasing the value of $n_{i}$ within this interval. Indeed, even if $n$ were allowed to become large as $\Delta x$ decreases, thus causing $n_{i}$ for any arbitrary interval to be large, $f_{i}^{*}$ could still be made arbitrarily small by decreasing the interval size sufficiently. However, this phenomenon would prevent us from defining a frequency for a single point in a continuous


Figure 2.2-6
random variable, unless we are content to associate it with some arbitrary class interval. To overcome this problem, probabilists have defined a different measure of frequency for continuous random variables that consists of the frequency of occurrence $f_{i}^{*}$ scaled by its class interval:

$$
\begin{equation*}
f_{i}=f_{i}^{*} / \Delta x . \tag{2.2-3}
\end{equation*}
$$

This normalized frequency, referred to as the sample density, should be relatively stable for reasonable choices of $\Delta x$ and $n$, and in the limiting case of $n$ approaching infinity and $\Delta x$ approaching zero, $f_{i}$ should be constant. An additional ordinate has been added to figure 2.2-5 to show the sample-density distribution of the $\log _{10} T$ data.
The sample distribution function of equation 2.2-2 can now be redefined in terms of equation $2.2-3$ as follows:

$$
\begin{equation*}
F_{n}(a)=\sum_{i \leq a / \Delta x} f_{i} \Delta x \tag{2.2-4}
\end{equation*}
$$

This definition lends itself to an exploration of the population equivalents of $F_{n}(a)$ and $f_{i}$. If the random sample is of sufficient size to sample every member of the sample space and $\Delta x$ is taken infinitely small, then the population equivalents of $F_{n}(a)$ and $f_{i}$ should be approached. By letting $n$ become large and $\Delta x$ small, we see that

$$
\begin{equation*}
\sum_{i \leq a / \Delta x} f_{i} \Delta x \approx \int_{-\infty}^{a} f(x) d x \tag{2.2-5}
\end{equation*}
$$

where $f(x)$, the population equivalent of $f_{i}$, is known as the probability density function. Because $f(x)$ is the population equivalent of $f_{i}$, then the integral representation in equation 2.2-5 of summing these scaled frequencies must be the population equivalent of $F_{n}(a)$, which of course is the same cumulative distribution function defined earlier in section 2.2.1:

$$
\begin{equation*}
F(a)=\int_{-\infty}^{a} f(x) d x=P(X \leq a) . \tag{2.2-6}
\end{equation*}
$$

However, because a random sample, whether it be finite or infinite, is countable, equation 2.2-5 must be given a special interpretation. Note that, because $f(x)$ is the continuous analog of $f_{i}$, it is always a non-negative quantity.

A stronger statement than equation 2.2-5 can be made concerning the equivalence of $F(a)$ and $F_{n}(a)$ for large sample sizes by noting that $F_{n}(a)$, prior to sampling, is a random variable. That is, if we were to collect different samples of the same size $n$ from the same population, we would not expect that $F_{n}(a)$, computed from each random sampling, would have the same value. We would only hope that, as $n$ becomes large, these different values would approach some constant. Indeed, probabilists have shown that, with a probability of one, $F_{n}(a)$ becomes the constant $F(a)$ as $n$ goes to infinity. This result is particularly remarkable if we first consider that $F_{n}(a)$ can only take on a countable number of values $k / n, 0 \leq k \leq n$, where $k$ is an integer (see equation $2.2-2$ ). Thus, although the values of $F(a)$ are uncountably infinite (continuous), $F_{n}(a)$ can only be, in the case that the random sample is infinitely large, at most, countably infinite. We will use this result loosely by allowing equation 2.2-5 to take on the indicated limits,

$$
\lim _{\substack{\Delta x \rightarrow 0 \\ n \rightarrow \infty}} F_{n}(a)=F(a)
$$

and noting that this result only can occur with a probability of one.
Both $f(x)$ and $F(x)$ are continuous functions of values of the random variable $X$. For the previously illustrated case of $X=\log _{10} T$, the density function might appear as in figure 2.2-7. Figure 2.2-7 represents the population equivalent of figure $2.2-5$, as if all possible outcomes of the random variable were available to us. Similarly, the cumulative frequency distribution, the population equivalent of figure $2.2-6$, for this random variable might appear as in figure 2.2-8.

Because of equation 2.2-6, the density function $f(x)$ can be defined in terms of the cumulative distribution function $F(x)$ by differentiation:


Figure 2.2-7


Figure 2.2-8

$$
\begin{equation*}
f(a)=\frac{d F(a)}{d a}=\frac{d}{d a} \int_{-\infty}^{a} f(x) d x . \tag{2.2-8}
\end{equation*}
$$

This result follows directly from the fundamental theorem of integral calculus, and is applicable only to density functions of continuous random variables. Equation $2.2-8$ is one of three concepts which defines density functions of continuous random variables. The other two state that $f(x)$ must be greater than or equal to zero for any possible value of the random variable and, as will be demonstrated in the next section, that the total mass under the frequency curve must be unity. All density functions of continuous random variables have these concepts in common.

## Problem 2.2-2

a. Construct histograms for the following specific-conductance data using class intervals of 100 and $200 \mu \mathrm{mho} / \mathrm{cm}$, such that the abscissa and ordinate of both histograms are scaled equally. What is the effect of changing the class interval?
b. Construct a cumulative frequency distribution from your $100 \mu \mathrm{mho} / \mathrm{cm}$ class-interval results. Let $X$ represent the specific-conductance random variable; what is

$$
\begin{aligned}
& P(X \leq 600) ? \\
& P(X>400) ? \\
& P(400<X \leq 600) ? \\
& P(X \leq 1300) ?
\end{aligned}
$$

Ordered specific-conductance data

| 63 | 423 | 501 | 582 | 685 | 836 |
| :---: | :---: | :---: | :---: | :---: | :---: |
| 76 | 433 | 504 | 596 | 697 | 839 |
| 168 | 439 | 509 | 598 | 700 | 876 |
| 278 | 440 | 512 | 600 | 704 | 882 |
| 301 | 440 | 518 | 604 | 710 | 895 |
| 304 | 440 | 518 | 617 | 721 | 897 |
| 310 | 444 | 527 | 620 | 723 | 904 |
| 315 | 452 | 529 | 627 | 724 | 906 |
| 319 | 452 | 533 | 629 | 726 | 915 |
| 323 | 452 | 537 | 632 | 728 | 948 |
| 332 | 456 | 538 | 636 | 740 | 968 |
| 347 | 462 | 542 | 641 | 750 | 969 |
| 357 | 469 | 552 | 647 | 750 | 982 |
| 359 | 471 | 562 | 659 | 764 | 997 |
| 363 | 473 | 564 | 659 | 765 | 1,030 |
| 389 | 477 | 564 | 661 | 779 | 1,080 |
| 407 | 487 | 565 | 664 | 783 | 1,106 |
| 408 | 490 | 566 | 665 | 789 | 1,120 |
| 411 | 492 | 570 | 670 | 808 | 1,170 |
| 413 | 493 | 575 | 673 | 808 | 1,230 |
| 417 | 493 | 578 | 675 | 815 |  |
| 418 | 499 | 582 | 677 | 820 |  |

### 2.2.4 Properties of Cumulative Distribution Functions

In the previous section, the cumulative distribution function $F(a)$, defined by the probability statement $P(X \leq a)$, was noted to have the integral form of equation 2.2-6 for continuous random variables. We state all manner of probability statements in terms of the cumulative distribution function, as this is a standard form. For this purpose, properties of cumulative distribution functions, with applications to other probability statements, are developed in this section.

The probability that a random variable $X$ takes on a value in the interval $(a, b]$ can be expressed in terms of cumulative distribution functions as

$$
\begin{equation*}
P(a<X \leq b)=\int_{a}^{b} f(x) d x=F(b)-F(a) \tag{2.2-9}
\end{equation*}
$$

This statement is a direct result of integral calculus, whereby integration is used to sum all the frequencies of occurrences of values of the random variable between $a$ and $b$. From
equation 2.2-9 one sees that the cumulative distribution function is a nondecreasing function of $x$, because

$$
0 \leq P(a<X \leq b) \leq 1
$$

That the total mass under the sample density curve $f_{i}$ is unity is evident from equation 2.2-4; that is,

$$
\begin{equation*}
\lim _{a \rightarrow \infty} F_{n}(a)=\lim _{a \rightarrow \infty} \sum_{i \leq a / \Delta x} f_{i} \Delta x=1 \tag{2.2-10}
\end{equation*}
$$

Because the probability density function $f(x)$ of a continuous random variable $X$ is a limiting form of the sample density $f_{i}$, the mass under its curve is also unity:

$$
\begin{equation*}
\lim _{a \rightarrow \infty} F(a)=\lim _{a \rightarrow \infty} \int_{-\infty}^{a} f(x) d x=1 \tag{2.2-11}
\end{equation*}
$$

Equation 2.2-11 is a property of all cumulative distribution functions. Similarly,

$$
\begin{equation*}
\lim _{a \rightarrow-\infty} F(a)=\lim _{a \rightarrow-\infty} \int_{-\infty}^{a} f(x) d x=0, \tag{2.2-12}
\end{equation*}
$$

which follows from integral calculus, is also a property of cumulative distribution functions.

Equation 2.2-11 allows one to express $P(X>a)$ as

$$
P(X>a)=\int_{a}^{\infty} f(x) d x=1-\int_{-\infty}^{a} f(x) d x=1-F(a),(2.2-13)
$$

which is also a result of Riemannian integration. An alternate statement of equation $2.2-13$ is that $P(X>a)=1-P(X \leq a)$.
By considering equation 2.2-9 in a limit form, we can also find the probability that $X=a$ :

$$
\begin{align*}
P(X=a) & =\lim _{\Delta x \rightarrow 0} P(a<X \leq a+\Delta x), \\
& =\lim _{\Delta x \rightarrow 0}[F(a+\Delta x)-F(a)]=0 . \tag{2.2-14}
\end{align*}
$$

This result is unique to continuous random variables, in contradistinction to discrete random variables. From equation 2.2-14 one sees that $P(X \leq a)$ is equivalent to $P(X<a)$ for continuous random variables, as the endpoint, $a$, of the semi-infinite interval does not contribute mass to the probability statement.
Equations 2.2-9, 2.2-11, and 2.2-13 can be demonstrated for discrete random variables by using the summation form of the cumulative distribution function (equation 2.2-1). In contradistinction to continuous random variables, the endpoint in $P(X \leq a)$ for a discrete random variable can contribute significant mass to the statement.

A number of frequency densities that result from randomness in nature, or probabilistic models of random events, have been investigated and published. Cumulative distributions of these densities are frequently tabulated and are found in many reference books on probability and statistics. Equations 2.2-9 and 2.2-13 are especially useful in evaluating probability statements of tabulated random variables.

### 2.2.5 An Example: The Normal Distribution

Let the random variable $Y$ represent the amount of titrant used in a titration experiment to neutralize measured amounts of the unknown $x$. A scatter diagram of titrant versus unknown might appear as in figure 2.2-9. The solid line represents the true stoichiometric balance between titrant and unknown. The dots, representing repetitions of the experiment, deviate from this line by an amount $\epsilon$, which represents a value of the measurement error $\varepsilon$. These errors represent a continuous random variable that could theoretically vary from $-\infty$ to $+\infty$ (the graphed points only represent a random sample from the population). If the experimental apparatus is functioning properly, however, we would expect these dots to be concentrated in the general vicinity of the solid line.
A distribution that is frequently used to model errors that are symmetrically distributed about some common point is the normal


Figure 2.2-9
distribution. The density of the normal distribution is a bell-shaped curve, symmetric about its mean $\mu_{\mathcal{E}}$, and with most of the mass concentrated within one standard deviation $\sigma_{\varepsilon}$ of the mean (see figure 2.2-10). In the case of the titration experiment, we would hope that the most frequently found value of the error would be near-zero and expect that $\mu_{\mathcal{E}}$ would equal zero. The standard deviation $\sigma_{\mathcal{E}}$ is a measure of the dispersion, or spread, of the errors about the mean and is equal to the distance from the mean to an inflection point on the curve $f(\epsilon)$. The mean and standard deviation will be formally defined in a later section.

A normal random variable is frequently standardized with its mean and standard deviation by the following transformation:

$$
\begin{equation*}
Z=\left(\varepsilon-\mu_{\varepsilon}\right) / \sigma_{\varepsilon} \tag{2.2-15}
\end{equation*}
$$

The cumulative distribution for this standard normal random variable is tabulated (table 2.10-1) for use by the investigator, since its probability density function, $f_{Z}(z)$, is parameter free:

$$
\begin{equation*}
f_{Z}(z)=\frac{e^{-z^{2} / 2}}{\sqrt{2 \pi}} \tag{2.2-16}
\end{equation*}
$$

Given the density function for the standard normal random variable, it is natural to inquire about the form of density, $f_{\mathcal{E}}(\epsilon)$, of the unnormalized random variable $\mathcal{E}$. Consider the cumulative frequency distribution for $Z$. By making the change of variables $z=\left(s-\mu_{\mathcal{E}}\right) / \sigma_{\mathcal{E}}$,

$$
\begin{align*}
F_{Z}(a) & =\int_{-\infty}^{a} \frac{e^{-z^{2} / 2}}{\sqrt{2 \pi}} d z \\
& =\frac{1}{\sqrt{2 \pi} \sigma_{\varepsilon}} \int_{-\infty}^{\epsilon} \exp \left[-\left(\frac{s-\mu_{\varepsilon}}{\sigma_{\varepsilon}}\right)^{2} / 2\right] d s \tag{2.2-17}
\end{align*}
$$

results where $\epsilon=a \sigma_{\mathcal{E}}+\mu_{\mathcal{E}}$ is a value of the unnormalized random variable. Since differentiation is the inverse operator of integration, equation $2.2-17$ is differentiated with respect to $\epsilon$ to find $f_{\mathcal{E}}(\epsilon)$ (see also equation $2.2-8$ ):

$$
\begin{align*}
f_{\varepsilon}(\epsilon) & =\frac{d}{d \epsilon} F_{Z}(a(\epsilon)) \\
& =\frac{1}{\sqrt{2 \pi \sigma_{\varepsilon}}} \exp \left[-\left(\frac{\epsilon-\mu_{\mathcal{E}}}{\sigma_{\mathcal{E}}}\right)^{2} / 2\right] . \tag{2.2-18}
\end{align*}
$$



Figure 2.2-10

Note that equation 2.2-18 is not parameter free, as this density is a function of the parameters $\mu_{\mathcal{E}}$ and $\sigma_{\mathcal{E}}$.

### 2.3 Expectation and the Continuous Random Variable

The discussion in this section is largely presented with continuous random variables in mind. All the results, however, are applicable to discrete random variables; whenever a quantity is defined by an integration over a probability density function for the continuous case, this same quantity can almost invariably be defined by a summation over the discrete density function for the discrete case. The reader should demonstrate the veracity of this statement.

### 2.3.1 The Mean

The mean is a measure of central tendency of a population. As an estimator of this central tendency, consider a finite random sample consisting of $n$ values $x_{i}$ of the random variable $X$. If the sample frequency of occurrence $f_{i}^{*}$ is estimated from this random sample, then a logical estimator of the central tendency is to sum the product of the central value $\bar{x}_{i}$ of each class interval and the frequency of occurrence for that interval:

$$
\begin{equation*}
\bar{x}=\sum_{i \leq x_{m} / \Delta x} f_{i}^{*} \bar{x}_{i}=\sum_{i \leq x_{m} / \Delta x} f_{i} \bar{x}_{i} \Delta x \tag{2.3-1}
\end{equation*}
$$

where $x_{m}$ is the upper limit of the largest class interval necessary to construct $f_{i}^{*}$. The frequencies of occurrence $f_{i}{ }^{*}$ in equation 2.3-1 can be looked upon as weights that sum to one, and the quantities $\bar{x}_{i}$ as equally spaced values of the random variable. The values of the random variable that occur more frequently, as indicated by the random sample, receive larger weights through equation $2.3-1$ and will have a greater influence on $\bar{x}$.
Equation 2.3-1 should be recognized by the reader as also being the definition for the center of mass of physical weights distributed along a line. That is, if $m f_{i}^{*}$ represents the mass of a weight located at $\bar{x}_{i}$, where $m$ is the total mass of all the weights, then equation $2.3-1$ would give us the center of mass of the line with respect to the origin. In the case of a histogram, the role of the weights is played by the sample frequency of occurrence for an interval, which gives us the approximate relative likelihood that any future value of the random variable will occur in that interval. For calculation purposes, this distributed weight over any interval $i$ is replaced by a point weight having the same mass as the distributed weight, but located at the center $\bar{x}_{i}$ of the interval. The sum of the products of the relative masses of these point weights, $f_{i}^{*}$, with their relative distances from the origin, $\bar{x}_{i}$, gives us the center of mass, which is also a measure of the central tendency. Of course, if the sample size $n$ were to become very large, then $\Delta x$ could be made very small, refining equation $2.3-1$ as an estimator of the central tendency of a random variable.
Reasoning similar to that leading to $\bar{x}$ as an estimator of the population mean can be applied directly to defining this parameter. First, given that the density function $f(x)$ is known, then the approximate frequency of occurrence of an event corresponding to an interval of size $\Delta x$ that has as its central value $\bar{x}_{i}$ is $f\left(\bar{x}_{i}\right) \Delta x$. Thus, assuming that these relative frequencies are centered at each $\bar{x}_{i}$, an approximate measure of the population central tendency, $\mu_{X}$, is

$$
\begin{equation*}
\mu_{\mathrm{X}} \approx \sum_{(\mathrm{all} i)} \bar{x}_{i} f\left(\bar{x}_{i}\right) \Delta x \tag{2.3-2}
\end{equation*}
$$

where the values $\bar{x}_{i}$ are equally spaced by $\Delta x$ from each other. Of course, by letting $\Delta x$
become smaller, a more accurate measure of $\mu_{X}$ is developed, until $\mu_{X}$, also known as the expected value, $E[X]$, of the random variable $X$, is defined by the following integral expression:

$$
\begin{equation*}
\mu_{X}=E[X]=\int_{-\infty}^{\infty} x f(x) d x \tag{2.3-3}
\end{equation*}
$$

This equation is the standard form for the expected value of a univariate random variable.

Equation 2.3-3 can also be developed directly from equation 2.3-1 by letting $n \rightarrow \infty$ and $\Delta x \rightarrow 0$ :

$$
\mu_{X}=\lim _{\substack{n \rightarrow \infty \\ \Delta x \rightarrow 0}} \bar{x}=\lim _{\substack{n \rightarrow \infty \\ \Delta x \rightarrow 0}} \sum_{i \leq x_{m}} \sum_{i x} \bar{x}_{i} f_{i} \Delta x=\int_{-\infty}^{\infty} x f(x) d x .(2.3-4)
$$

That is, as $\Delta x$ becomes smaller and as the number of observations becomes very large, $\bar{x}_{i}$ becomes a unique continuous value of $X, f_{i}$ becomes the continuous function $f(x)$, and the summation can be replaced by an integration. As in the case of equation 2.2-7, we can only say that the limit indicated in equation 2.3-4 is reached with a very high probability as $n$ becomes large; however, this probability should be unity as $n$ becomes infinite.
Note that $\mu_{X}$ is a population parameter that is characteristic of the random variable $X$, while $\bar{x}$, being derived from values of a finite random sample from the population of $X$, is only an estimate for $\mu_{X}$. Estimators such as $\bar{x}$ will be developed in greater detail in a later section.

## Problem 2.3-1

a. Find $\bar{x}$ from $100 \mu \mathrm{mho} / \mathrm{cm}$ histogram of problem 2.2-2.
b. Find $\mu_{X}$ for the random variable of problem 2.2-1.

### 2.3.2 Generalization and Application of the Expectation Operator

The operation of finding an expected value can be generalized by considering a function $g(X)$ of continuous random variable $X$. If we
wish to find the average effect of the function $g(x)$ over the outcomes of a random sampling, we again resort to the approximation

$$
\begin{equation*}
\overline{g(x)}=\sum_{i \leq x_{m} / \Delta x} f_{i}^{*} g\left(\bar{x}_{i}\right) . \tag{2.3-5}
\end{equation*}
$$

That is, we weight $g(x)$, where $g(x)$ is evaluated at the center of every class interval, by the frequency of occurrence of that interval and sum all the weighted values of $g\left(\bar{x}_{i}\right)$. To obtain the population equivalent of $\overline{g(x)}, n$ is taken to be large, while $\Delta x$ is taken to be small; this equivalent is denoted by the expectation symbol $E[g(X)]$ :

$$
\begin{equation*}
E[g(X)]=\lim _{\substack{n \rightarrow \infty \\ \Delta x \rightarrow 0}} \overline{g(x)}=\int_{-\infty}^{\infty} g(x) f(x) d x \tag{2.3-6}
\end{equation*}
$$

Equation 2.3-6 represents the general form of the expectation operator for a univariate distribution when the random variable is continuous. A similar form exists for discrete random variables, in which the integration has been replaced by summation.

A trivial but useful property of the expectation operator is that the expected value of any constant $c$ is that constant; for the continuous case, this is easily demonstrated as

$$
\begin{equation*}
E[c]=\int_{-\infty}^{\infty} c f(x) d x=c \int_{-\infty}^{\infty} f(x) d x=c \tag{2.3-7}
\end{equation*}
$$

where equation $2.2-11$ has been invoked. A more important property of $E$ is that it is a linear operator; that is,

$$
E\left[a g_{1}(X)+b g_{2}(X)\right]=a E\left[g_{1}(X)\right]+b E\left[g_{2}(X)\right] .(2.3-8)
$$

This property results because in the continuous case, integration itself is a linear operator:

$$
\begin{align*}
\int_{-\infty}^{\infty}\left(a g_{1}(x)+b g_{2}(x)\right) f(x) d x & =a \int_{-\infty}^{\infty} g_{1}(x) f(x) d x \\
& +b \int_{-\infty}^{\infty} g_{2}(x) f(x) d x \tag{2.3-9}
\end{align*}
$$

As a practical example of finding the expected value of a random variable, consider the problem of finding the mean of $X$ where $X$ is a normal random variable with mean $\mu_{X}$ and standard deviation $\sigma_{X}$ :
$E[X]=\frac{1}{\sqrt{2 \pi \sigma_{X}}} \int_{-\infty}^{\infty} x \exp \left[-\left(\frac{x-\mu_{X}}{\sigma_{X}}\right)^{2} / 2\right] d x .(2.3-10)$
By a change of variable $z=\left(x-\mu_{X}\right) / \sigma_{X}$, we see that equation 2.3-10 becomes

$$
\begin{equation*}
E[X]=\frac{1}{\sqrt{2 \pi}} \int_{-\infty}^{\infty}\left(\mu_{X}+\sigma_{X^{z}} z\right) e^{-z^{2} / 2} d z=\mu_{X} \tag{2.3-11}
\end{equation*}
$$

because $z \cdot \exp \left(-z^{2} / 2\right)$ is an odd function of $z$ in the interval $(-\infty, \infty)$, and equation 2.3-7 holds ( $\mu_{X}$, being a population parameter, is constant). Equation 2.3-11 is the reason why $\mu_{X}$ is defined to be $E[X]$.

### 2.3.3 The Variance, Standard Deviation, and Coefficient of Variation

Although the mean $\mu_{X}$ is a measure of the central tendency of a random variable, it gives no information as to how frequently a random variable will be encountered in its vicinity. The variance $\sigma_{X}^{2}$, defined as the expected value of the function $g(X)=\left(X-\mu_{X}\right)^{2}$, is a population parameter that quantifies this concept. The variance can also be looked upon as an operator that is defined in terms of another operator (the expectation operator) as follows:

$$
\begin{align*}
\sigma_{X}^{2} & =\operatorname{Var}[X]=E\left[\left(X-\mu_{X}\right)^{2}\right], \\
& =\int_{-\infty}^{\infty}\left(x-\mu_{X}\right)^{2} f(x) d x \tag{2.3-12}
\end{align*}
$$

where $\operatorname{Var}[X]$ represents an operator that operates on $X$. The intuitive sense of $\sigma_{X}^{2}$ is that it is the sum of the frequency weighted deviations, which have been squared, from the mean. As such, $\sigma_{X}^{2}$ represents the amount of dispersion of the random variable about the mean: when $\sigma_{X}^{2}$ is relatively large, then a random variable is less likely to have values in the
immediate vicinity of the mean. The standard deviation $\sigma_{X}$ is simply the square root of the variance: $\sigma_{X}=(\operatorname{Var}[X])^{1 / 2}$.
By exercising the linear property of the expectation operator, equation 2.3-12 can be expressed in an alternate form:

$$
\begin{align*}
\sigma_{X}^{2} & =\operatorname{Var}[X]=E\left[\left(X-\mu_{X}\right)^{2}\right]=E\left[X^{2}-2 X \mu_{X}+\mu_{X}^{2}\right], \\
& =E\left[X^{2}\right]-\mu_{X}^{2} . \tag{2.3-13}
\end{align*}
$$

The variance operator, like the expectation operator, can be generalized to operate on any function $g(X)$ :

$$
\begin{equation*}
\operatorname{Var}[g(X)]=E\left[g^{2}(X)\right]-(E[g(X)])^{2} . \tag{2.3-14}
\end{equation*}
$$

The variance operator, however, is not a linear operator, as demonstrated with the function $g(X)=a+b X:$

$$
\begin{align*}
\operatorname{Var}[g(X)] & =E\left[(a+b X)^{2}\right]-(E[a+b X])^{2}, \\
& =b^{2} E\left[X^{2}\right]-b^{2} \mu_{X}^{2}=b^{2} \sigma_{X}^{2}, \tag{2.3-15}
\end{align*}
$$

because $E[a+b X]=a+b \mu_{X}$. By letting $b=0$ in the above example, one can demonstrate that the variance of a constant, as expected, is zero.
When the standard deviation is normalized by the mean of the random variable ( $\mu_{X} \neq 0$ ), it is referred to as the coefficient of variation $V_{X}$ :

$$
\begin{equation*}
V_{X}=\sigma_{X} / \mu_{X} . \tag{2.3-16}
\end{equation*}
$$

Estimators for this population parameter, as well as the variance and standard deviation, are discussed in a later section of this report.
As an example of an application of the variance operator, consider an application on the standard normal random variable $Z=\left(X-\mu_{X}\right) / \sigma_{X}$ :

$$
\begin{equation*}
\operatorname{Var}\left[\left(X-\mu_{X}\right) / \sigma_{X}\right]=\operatorname{Var}[X] / \sigma_{X}^{2}=1 \tag{2.3-17}
\end{equation*}
$$

which results by analogy with equation 2.3-15. Thus, if $X$ is a normal random variable with mean $\mu_{X}$ and variance $\sigma_{X}^{2}$, then $Z$ is a zero-mean random variable with a variance of unity, which is commonly denoted $N(0,1)$.

## Problem 2.3-2

a. For $f(x)= \begin{cases}2 x & 0 \leq x \leq 1 \\ 0 & \text { otherwise }\end{cases}$
i. Plot $f(x)$.
ii. Derive and plot $F(x)$.
iii. Calculate $E[X]$ and $\operatorname{Var}[X]$.
b. An estimator of the variance $\sigma_{X}^{2}$ of the random variable $X$ can be developed directly from equation $2.3-12$. First, $f(x) d x$ is estimated by $f_{i}^{*}$ of equation 2.2-2. Then $x$ is replaced by $\bar{x}_{i}$, the center of each class interval corresponding to $f_{i}^{*}$. Finally, $\mu_{X}$ is estimated by $\bar{x}$ from equation 2.3-1. Then

$$
s_{X}^{* 2}=\sum_{i \leq x_{m}^{\prime}(\Delta x}\left(\bar{x}_{i}-\bar{x}\right)^{2} f_{i}^{*}
$$

gives an estimate of $\sigma_{X}^{2}$. Apply this estimator to the log-transmissivity data of table 2.2-3.

### 2.4 Jointly Distributed Random Variables

The investigator frequently encounters the problem that he or she has to deal with two (or more) random variables in the same probability statement. As an example, in the case of random variables $X$ and $Y$, where $X$ and $Y$ are possibly correlated, one might desire the probability that $X$ is less than or equal to $a$, and $Y$ is less than or equal to $b$. If the investigator should know the form of the joint probability density function $f(x, y)$ for these two random variables, then this probability statement is definable:
$P(X \leq a$ and $Y \leq b)$

$$
\begin{equation*}
=\int_{-\infty}^{a} \int_{-\infty}^{b} f(x, y) d y d x=F(a, b) \tag{2.4-1}
\end{equation*}
$$

where $F(a, b)$ is the equivalent cumulative distribution function. (The statement $P(X \leq a$ and $Y \leq b)$ is also denoted frequently as $P(X \leq a$, $Y \leq b)$; the more explicit form will be used in this discussion.) As in the univariate case, it is required that the mass under the joint probability density function equal unity:

$$
\begin{equation*}
\int_{-\infty}^{\infty} \int_{-\infty}^{\infty} f(x, y) d x d y=1 \tag{2.4-2}
\end{equation*}
$$

The concept of joint probability density functions applies to any number of multiple random variables; the following discussion, however, is largely restricted to the bivariate case.
As an example of an experiment yielding jointly distributed random variables, consider the results from a simple nonsteady-state pumping test of a confined aquifer: When the Theis equation is used to evaluate data from these tests, information concerning the storativity and transmissivity of the aquifer results. Indeed, we can easily imagine that these quantities are random variables, varying from location to location in response to the local distribution of materials composing the aquifer. More important, however, would be the manner in which they vary with regard to each other: Should the clay content of the aquifer increase at some point, it might be expected that the transmissivity will decrease while the storativity, reflecting the compressibility of the aquifer, would increase. Thus, quite possibly these quantities, with regard to the aquifer in question, could be treated as jointly distributed random variables which are, in some manner, interdependent.
Assume for the moment that we have determined the form of the joint density function of storativity and transmissivity. For argument's sake, let $X$ represent the transmissivity random variable (or its logarithmic transformation) and $Y$ represent the storativity (or a functional transformation thereof) and then denote the joint density as $f(x, y)$. Now assume that we are interested in the probability that $X$ is less than or equal to $a$, regardless of the value of $Y$; that is, we wish to evaluate the probability that our measure of the transmissivity will take on a specific range of values, whereas the exact value of storativity is unimportant to us. For our probability statement regarding $X$ to be meaningful, all values of $Y$ which influence the joint density function must be taken into consideration, for different values of $Y$ would surely influence a statement on $X$ alone. To obtain the total contribution of $Y$ to the joint density function, we allow that $Y$ may take on any value in the inter-$\operatorname{val}(-\infty, \infty)$ and write our probability statement as

$$
\begin{align*}
P(X<a \text { and }-\infty<Y<\infty) & =\int_{-\infty}^{a} \int_{-\infty}^{\infty} f(x, y) d y d x \\
& =\int_{-\infty}^{a} f_{X}(x) d x \tag{2.4-3}
\end{align*}
$$

in which the evaluation of the inner integral with respect to $y$ results in a function $f_{X}(x)$ that meets all requirements to be a probability density function. Thus, in general, univariate density functions can be recovered from joint density functions by integration, and this integration has the effect of summing the total contribution of one random variable in the bivariate joint density onto the axis of the other variate, the second variate giving the relative frequency of occurrence of the event in question. These densities are referred to as marginal probability density functions and, with respect to the bivariate joint density $f(x, y)$, they are defined as

$$
\begin{align*}
& f_{X}(x)=\int_{-\infty}^{\infty} f(x, y) d y,  \tag{2.4-4}\\
& f_{Y}(y)=\int_{-\infty}^{\infty} f(x, y) d x, \tag{2.4-5}
\end{align*}
$$

where $f_{X}(x)$ is the marginal density for the $X$ random variable and similarly $f_{Y}(y)$ for the $Y$ random variable. The marginal-density concept is easily extended to multiple random variables when they are jointly distributed.

### 2.4.1 Expectation of Jointly Distributed Random Variables

The expectation operator for jointly distributed random variables is defined in the same manner as in the univariate case. Thus, if $X$ and $Y$ are jointly distributed, and $g(X, Y)$ is a function of these two random variables, then a general definition of the expectation operator is

$$
\begin{equation*}
E[g(X, Y)]=\int_{-\infty}^{\infty} \int_{-\infty}^{\infty} g(x, y) f(x, y) d y d x \tag{2.4-6}
\end{equation*}
$$

If on the other hand, we desire the expected value of $h(X)$, which is a function of $X$ only, we set $g(x, y)$ equal to $h(x)$ and proceed as in equation 2.4-6). The result,

$$
\begin{align*}
E[h(X)] & =\int_{-\infty}^{\infty} \int_{-\infty}^{\infty} h(x) f(x, y) d y d x \\
& =\int_{-\infty}^{\infty} h(x) f_{X}(x) d x \tag{2.4-7}
\end{align*}
$$

shows that, in such cases, finding the expected value reduces to finding the marginal density and integrating. By letting $h(X)$ equal $X$; one realizes that the mean $\mu_{X}$ is equal to the integral of the product of $x$ and the marginal density $f_{X}(x)$, as might be expected.

Consider the case where $g(X, Y)$ equals the product $\left(X-\mu_{X}\right)\left(Y-\mu_{Y}\right)$. The expected value of this product gives an indication of how $X$ and $Y$ vary together. If the absolute value of the expected value of this product is exceptionally large, then one would expect that $X$ and $Y$ are highly correlated. This expected value of $X$ and $Y$ is referred to as the covariance of $X$ and $Y$, and is denoted $\operatorname{Cov}[X, Y]$ or $\sigma_{X Y}$ :

$$
\begin{align*}
\operatorname{Cov}[X, Y] & =E\left[\left(X-\mu_{X}\right)\left(Y-\mu_{Y}\right)\right] \\
& =E[X Y]-\mu_{Y} \mu_{X} \tag{2.4-8}
\end{align*}
$$

Note that the covariance of $X$ with itself is $\operatorname{Cov}[X, X]=\operatorname{Var}[X]$.

Returning to the example of transmissivities and storativities of the previous section, we see that the covariance provides a measure of the degree of interdependence between random variables. That is, because $X$ and $Y$ are both random, we would not expect observations of $X$ and $Y$ to show a perfect relationship; rather, the relationship will be clouded with noise. Because the expected value implies a frequencyweighted average of the function in question, and because the frequency distribution will reflect the amount of relationship between $X$ and $Y$, summing the product of these frequency weights with $\left(X-\mu_{X}\right)\left(Y-\mu_{Y}\right)$ over the total variate space will give the average relationship between $X$ and $Y$. It will be demonstrated in the next section that if $X$, the measure of transmissivity, and $Y$, the measure of storativity, were independent, then the covariance would theoretically be zero. However, if our intuition is correct, we would not expect this; rather we might expect, should the aquifer have a rather high clay content, that the two variables will be negatively correlated.

If the covariance is normalized with the standard deviations of the two random variables, then it is referred to as the correlation coefficient $\rho_{X Y}:$

$$
\begin{equation*}
\rho_{X Y}=\operatorname{Cov}[X, Y] /\left(\sigma_{X} \sigma_{Y}\right) \tag{2.4-9}
\end{equation*}
$$

The correlation coefficient, as a measure of the linear relationship between $X$ and $Y$, has the property that its absolute value is less than or equal to unity:

$$
\begin{equation*}
\left|\rho_{X Y}\right| \leq 1 \tag{2.4-10}
\end{equation*}
$$

That is, when $X$ and $Y$ are precisely linearly related, then $\left|\rho_{X Y}\right|$ will equal unity. If there is no relationship between $X$ and $Y$, as shown in the next section, $\operatorname{Cov}[X, Y]$ and therefore $\rho_{X Y}$ will be zero. This property is demonstrated in appendix 2.11.1, but this appendix requires some knowledge of the next section.

### 2.4.2 Independent Random Variables

Two random variables $X$ and $Y$ are said to be independent if, for all $a$ and $b$,

$$
\begin{align*}
& P(X \leq a \text { and } Y \leq b) \\
&=P(X \leq a) P(Y \leq b) \\
&=\int_{-\infty}^{a} \int_{-\infty}^{b} f_{X}(x) f_{Y}(y) d y d x, \tag{2.4-11}
\end{align*}
$$

where $f_{X}(x)$ and $f_{Y}(y)$ are the densities of $X$ and $Y$, respectively. Equation 2.4-11 implies that the joint density function of two independent random variables is the product of their individual densities, that is,

$$
\begin{equation*}
f(x, y)=f_{X}(x) f_{Y}(y) \tag{2.4-12}
\end{equation*}
$$

Of course, an event corresponding to $X \leq a$ and $Y \leq b$ would be expected to occur with equal or less frequency than an event corresponding to either $X \leq a$ or $Y \leq b$ separately. Only in the case of a complete lack of dependence between these events can we say that $P(X \leq a$ and $Y \leq b)$ $=P(X \leq a) P(Y \leq b)$. This is a somewhat intuitive result that has already been used in connection with the two-dice experiment; if $X$ is an outcome of the first die and $Y$ the second, then $P(X=1$ and $Y=2)=P(X=1) P(Y=2)=1 / 36$.
A random sample is, ideally, a collection of independent random variables. That is, prior to their observation, each element of a random sample is a random variable; its value is not known until after the observation process is
completed. These outcomes should not have any interdependence which might affect the sample density. This generally requires careful design of the experiment from which the observations result so that all $X_{i}, i=1, \ldots, n$, are independent.

The question of independence of two random variables $X$ and $Y$ has important implications on their covariance, for if $X$ and $Y$ are independent, then

$$
\begin{align*}
\operatorname{Cov}[X, Y] & =\int_{-\infty}^{\infty} \int_{-\infty}^{\infty}\left(x-\mu_{X}\right)\left(y-\mu_{Y}\right) f_{X}(x) f_{Y}(y) d x d y \\
& =\int_{-\infty}^{\infty}\left(x-\mu_{X}\right) f_{X}(x) d x \int_{-\infty}^{\infty}\left(y-\mu_{Y}\right) f_{Y}(y) d y \\
& =E\left[X-\mu_{X}\right] E\left[Y-\mu_{Y}\right]=0 \tag{2.4-13}
\end{align*}
$$

However, if the covariance of two random variables is zero, it does not necessarily follow that they are independent. One may only suspect that independence is the cause of a zero covariance.

### 2.4.3 Conditional Probabilities

The marginal probability density function, as developed in equation 2.4-3, can be considered to be a special case of a more general concept referred to as conditioning. Generally speaking, a multivariate probability statement is subject to conditioning when a subset of the random variables pertaining to an experiment falls under some restriction, causing the remaining variables to be conditioned by this restriction. In the case of the marginal density function, we examined the probability that $X$ is less than $a$, given that $Y$ can take on any value in the interval $(-\infty, \infty)$. Thus, the restriction that $Y$ take on a specific set of values conditions the probability that $X$ is less than $a$. Formally, we state this as

$$
P(X<a \mid-\infty<Y<\infty)
$$

In general, the restriction can be applied to any interval ( $b, c$ ), where $b<c$, and need not be limited to the interval $(-\infty, \infty)$. However, as in the case of the marginal density, the variable or variables subject to restriction are effectively removed from the probability statement; the variable or variables being conditioned are the ones over which the frequency of occurrence of a specific event may be questioned.

When an experiment which results in a bivariate random variable is conditioned over a range other than $(-\infty, \infty)$, a reduction of the potential sample space available to the experiment results. In the previous example of jointly varying transmissivities $X$ and storativities $Y$, if we were interested in the conditioned results that $X$ is less than $a$, given that we are only interested in a specific range of values ( $b, c$ ) for storativities, then the specific value which $Y$ takes on does not interest us, as long as it falls between $b$ and $c$. One could proceed as in equation 2.4-3 to evaluate this probability, except for an obvious pitfall: The resulting probability statement over $X$, where $X$ can take on any value less than $a$, would not necessarily have the property of cumulative distribution functions noted in equation 2.2-11. That is, as $X$ is the remaining active random variable in the probability statement, its probability of occurrence over the interval $(-\infty, \infty)$ should be unity:

$$
\lim _{a \rightarrow \infty} P(X<a \mid b<y<c)=1
$$

However, by restricting $Y$ to a specific interval ( $b, c$ ), then as $a$ goes to infinity, an integral of the form of equation 2.4-3 will most probably have a lesser value than unity when the inner integral over $y$ is restricted to a range of something less than $(-\infty, \infty)$. Thus, an integration with the form of equation 2.4-3 alone will not produce a form suitable to serve as a cumulative distribution function for the conditioned variable $X$.

So that a probability statement resulting from conditioning has the limiting value of unity, these statements must be appropriately normalized. If, as in the bivariate case, we desire $P(X<a \mid b<Y<c)$, then we must normalize by $P(-\infty<X<\infty$ and $b<Y<c)$; that is,

$$
\begin{align*}
P(X<a \mid b<Y<c) & =\frac{P(X<a \text { and } b<Y<c)}{P(-\infty<X<\infty \text { and } b<Y<c)} \\
& =\frac{\int_{-\infty}^{a} \int_{b}^{c} f(x, y) d y d x}{\int_{-\infty}^{\infty} \int_{b}^{c} f(x, y) d y d x} \\
& =\frac{\int_{-\infty}^{a} \int_{b}^{c} f(x, y) d y d x}{\int_{b}^{c} f(y) d y} . \tag{2.4-14}
\end{align*}
$$

Thus, the conditional probability density function, $f(x \mid b<Y<c)$, for the conditioned random variable $X$ may be defined as

$$
f(x \mid b<Y<c)=\int_{b}^{c} f(x, y) d y / \int_{b}^{c} f_{Y}(y) d y,(2.4-15)
$$

which of course gives the limiting value of unity when integrated with respect to $x$ over the interval $(-\infty, \infty)$. Note that when $b=-\infty$ and $c=\infty$, then $f(x \mid b<Y<c)=f_{X}(x)$, as indicated by the previous discussion of marginal densities.

Remarkably, the conditional density exists even when the restriction is that, in the example of the bivariate case, $Y$ take on a specific value. To see this easily, consider $P(X<a \mid Y=c)$; then equation 2.4-15 may be written as

$$
\begin{align*}
f(x \mid Y=c) & =\lim _{\delta \rightarrow 0} \frac{\int_{c}^{c+\delta} f(x, y) d y}{\int_{c}^{c+\delta} f_{Y}(y) d y} \\
& =\lim _{\delta \rightarrow 0} \frac{[F(x, c+\delta)-F(x, c)] / \delta}{\left[F_{Y}(c+\delta)-F_{Y}(c)\right] / \delta} \\
& =\left.\frac{d F(x, y) / d y}{d F_{Y}(y) / d y}\right|_{y=c} \\
& =\frac{f(x, c)}{f_{Y}(c)} \tag{2.4-16}
\end{align*}
$$

Thus, we may recover the density function for $X$ for any particular slice, $Y=c$, through the joint density function $f(x, y)$. If $f(x, y)$ were defined for the example of transmissivity and storativity random variables, equation 2.4-16 would enable us to predict the probability of events concerning transmissivity $X$ for any given value of storativity $Y$.
The student should also note that some remarkable simplifications result if $X$ and $Y$ are independent random variables. That is, if $X$ and $Y$ are independent, then from equations 2.4-12 and $2.4-14$ we see that

$$
\begin{equation*}
P(X<a \mid b<Y<c)=P(X<a) . \tag{2.4-17}
\end{equation*}
$$

Indeed, this is yet another way in which we can define independence of random variables.

The following problem is intended to familiarize the student with the concept of conditioning; it is not intended to be rigorous. The key to understanding conditioning, especially for discrete random variables, is to understand how it restricts the sample space and realize that the probability of occurrence of an event which contains the entire remaining sample space must be unity.

## Problem 2.4-1

a. Given two dice that are thrown sequentially, what is the probability that the first is a three and the second is a two? That is,

$$
P(X=3 \text { and } Y=2) ?
$$

b. What is the probability that the sum of the dice is five? That is,

$$
P(X+Y=5) ?
$$

c. Given that the first die is three, what is the probability that the second is two? That is,

$$
P(Y=2 \mid X=3) ?
$$

d. Given that the first die is three, what is the probability that the sum of the two dice is five? That is,

$$
P(X+Y=5 \mid X=3) ?
$$

e. Given that the first die is three, what is the probability that the sum of the two dice is less than or equal to five? That is,

$$
P(X+Y \leq 5 \mid X=3) ?
$$

Parts c, d, and e are conditional probability statements; that is, the probability statement is conditioned by prior information.

### 2.4.4 Variance of a Column Vector

Our purpose in this section is to develop a representation for the variance of a column vector. As a vehicle to this end, consider the linear equation

$$
\begin{equation*}
Y=a_{1} X_{1}+a_{2} X_{2}+a_{3} X_{3} \tag{2.4-18}
\end{equation*}
$$

where $Y, X_{1}, X_{2}$, and $X_{3}$ are random variables and $a_{1}, a_{2}$ and $a_{3}$ are constants. The variance of $Y$ is

$$
\begin{align*}
\operatorname{Var}[Y] & =E\left\{\left[\left(a_{1} X_{1}+a_{2} X_{2}+a_{3} X_{3}\right)\right.\right. \\
& \left.\left.-E\left(a_{1} X_{1}+a_{2} X_{2}+a_{3} X_{3}\right)\right]^{2}\right\} \\
& =a_{1}^{2} \sigma_{X_{1}}^{2}+a_{2}^{2} \sigma_{X_{2}}^{2}+a_{3}^{2} \sigma_{X_{3}}^{2} \\
& +2 a_{1} a_{2} \sigma_{X_{1} X_{2}}+2 a_{1} a_{3} \sigma_{X_{1} X_{3}} \\
& +2 a_{2} a_{3} \sigma_{X_{2} X_{3}} \tag{2.4-19}
\end{align*}
$$

where correlations between $X_{1}, X_{2}$, and $X_{3}$ have been allowed for. A vector representation for equation 2.4-19 is ${ }^{1}$

$$
\begin{gather*}
\operatorname{Var}[Y]=E\left[Y^{2}\right]-(E[Y])^{2}=E\left[\underline{a X X^{T}} \underline{a}^{T}\right] \\
-E[\underline{a X}] E[\underline{a X}] \tag{2.4-20}
\end{gather*}
$$

where

$$
\underline{a}=\left[\mathrm{a}_{1}, \mathrm{a}_{2}, \mathrm{a}_{3}\right] \text { and } \underline{X}=\left[\begin{array}{l}
X_{1} \\
X_{2} \\
X_{3}
\end{array}\right]
$$

Because expectation is a linear operator, the right side of equation 2.4-20 can be expressed as

$$
\begin{gathered}
E\left[a X X^{T} \underline{a}^{T}\right]-E[a X] E\left[\underline{X}^{T} \underline{a}^{T}\right] \\
=\underline{a} E\left[\left(\underline{X}-E[\underline{X})(\underline{X}-E[\underline{X}])^{T}\right] \underline{a}^{T}(2.4-21)\right.
\end{gathered}
$$

where $a X=\underline{X}^{T} \underline{a}^{T}$. The expected value of a matrix is the matrix of expected values of each element. Thus,

$$
E\left[(\underline{X}-E(\underline{X}))(\underline{X}-E(\underline{X}))^{T}\right]=
$$

$$
\left[\begin{array}{lll}
\sigma_{X_{1}}^{2} & \sigma_{X_{1} X_{2}} & \sigma_{X_{1} X_{3}}  \tag{2.4-22}\\
\sigma_{X_{1} X_{2}} & \sigma_{X_{2}}^{2} & \sigma_{X_{2} X_{3}} \\
\sigma_{X_{1} X_{3}} & \sigma_{X_{2} X_{3}} & \sigma_{X_{3}}^{2}
\end{array}\right]
$$

This matrix is defined to be the variance of a $3 \times 1$ column vector $\underline{X}$, and allows one to express equation 2.4-19 in matrix notation as

$$
\begin{equation*}
\operatorname{Var}[Y]=\underline{a} \operatorname{Var}[X] \underline{a}^{T} \tag{2.4-23}
\end{equation*}
$$

[^2]If the variances $\sigma_{X_{1}}^{2}, \sigma_{X_{2}}^{2}$, and $\sigma_{X_{3}}^{2}$ are all equal, the matrix $2.4-22$ becomes

$$
\operatorname{Var}[X]=\left[\begin{array}{lll}
1 & \rho_{X_{1} X_{2}} & \rho_{X_{1} X_{3}}  \tag{2.4-24}\\
\rho_{X_{1} X_{2}} & 1 & \rho_{X_{2} X_{3}} \\
\rho_{X_{1} X_{3}} \rho_{X_{2} X_{3}} & 1
\end{array}\right] \sigma^{2}
$$

where $\rho_{X_{i} X_{j}}$ is the correlation coefficient for $X_{i}$ and $X_{j}$, and $\sigma^{2}$ is the common variance. A further reduction in equation 2.4-22 occurs if $X_{1}$, $X_{2}$, and $X_{3}$ are uncorrelated, causing the correlation coefficients in equation 2.4-24 to be zero. In this case,

$$
\begin{equation*}
\operatorname{Var}[X]=\underline{I} \sigma^{2} \tag{2.4-25}
\end{equation*}
$$

where $I$ is a $3 \times 3$ identity matrix. These forms have practical importance in regression.

## Problem 2.4-2

a. Carry out the expectation indicated and show that equation 2.4-19 holds.
b. Demonstrate that equation 2.4-21 holds and that

$$
\operatorname{Var}[\underline{a}]=\underline{a} \operatorname{Var}[X] \underline{a}^{T}
$$

c. Let $Y_{i}=\underline{a}_{i} \underline{X}$, where $\underline{a}_{i}=\left[a_{i 1}, a_{i 2}, a_{i 3}\right]$, and $\underline{X}$, defined as in equation 2.4-20, is a column vector of random variables. Further, let $\underline{Y}=A X$, where
$\underline{Y}=\left|\begin{array}{l}Y_{1} \\ Y_{2} \\ \cdot \\ \cdot \\ Y_{p}\end{array}\right| \quad$ and $\quad \underline{A}=\left|\begin{array}{l}\underline{a}_{1} \\ \underline{a}_{2} \\ \cdot \\ \cdot \\ \cdot \\ \underline{a}_{p}\end{array}\right|$
That is, $A$ is a $p \times 3$ matrix composed of the row vectors $a_{i}, i=1, \ldots, p$. Show that $\operatorname{Var}[\underline{Y}]=\operatorname{Var}[A X]=A \operatorname{Var}[\underline{X}] A^{T}$. (Hint: Equation 2. $\dot{4}-22$ still defines the variance of a column vector;

$$
\begin{aligned}
\sigma_{Y_{i} Y_{j}} & =E\left[Y_{i} Y_{j}\right]-E\left[Y_{i}\right] E\left[Y_{j}\right] \\
& \left.=E\left[\underline{a}_{i} \underline{X X}^{T} \underline{a}_{j}^{T}\right]-E\left[\underline{a}_{i} X\right] E\left[\underline{X}^{T} \underline{a}_{j}^{T}\right] .\right)
\end{aligned}
$$

### 2.5 Estimators of Population Parameters

A statistic is defined as any computation from a random sample resulting in a specific value. As such, a statistic is considered to be a random variable, since it is highly probable that the computed value would change from random sampling to random sampling. Note that this definition precludes that a statistic contain any unknown parameters. Estimators of population parameters are considered to be statistics and, therefore, random variables. Consider equation 2.3-1 as an estimator of the mean:

$$
\begin{equation*}
\bar{x}=\sum_{i \leq x_{m} / \Delta x} f_{i}^{*} \bar{x}_{i} \tag{2.5-1}
\end{equation*}
$$

The estimated frequency $f_{i}^{*}$ is computed from values of observations $x_{i}$ originating from a random sampling of the sample space. However, prior to sampling, a random sample is merely an abstract collection of random variables $X_{i}, i=1, \ldots, n$. Any function of random variables, as equation $2.5-1$ would be prior to sampling, is also a random variable, perhaps having a completely different distribution than those individuals composing the collection.

Our discussion of statistics will largely be from the a priori viewpoint; that is, in the case of equation 2.5-1, $\bar{x}$ is the value of the random variable $\bar{X}$, which is an estimator for the population mean, as developed from some arbitrary random sample.

### 2.5.1 Mean Estimator

As an estimator for the population mean, equation $2.5-1$, in addition to being cumbersome to compute, has the debility that it is dependent upon an arbitrary selection of a class interval. That is, since $f_{i}^{*}$ is dependent upon $\Delta x$, the value of $\bar{X}$ will depend upon the choice of $\Delta x$ used in the computation. As a means of pursuing this problem, assume that we have at our disposal a random sample consisting of $n$ observations, and at some point their distribution appears as in figure 2.5-1. Because $\Delta x$ is arbitrary, it can be reduced to $\delta$, the minimum of all differences in neighboring values of the random
variable. Then $f_{i}^{*}$ would take on only two values, $1 / n$ or 0 , and would have the ragged sawtoothed shape shown in figure 2.5-2. The teeth in figure 2.5-2 will be concentrated in regions where $f_{i}^{*}$ in figure 2.5-1 is larger. Of course, that a repeat value of a random variable could occur is highly unlikely, as the probability of such an occurrence is essentially zero for a continuous random variable (see equation 2.2-14). By using the class interval $\delta$ of figure 2.5-2 in our computation, we see that the problem of estimating frequency weights for $\bar{x}_{i}$, the central location in each class interval, from a discrete data set has essentially been removed, as these weights now take on only two specific values for this and any other smaller class interval.

The use of the smaller class interval, $\delta$, is expected to produce a better estimator of the population mean because a value of $\bar{X}$ would contain less measurement error associated with the arbitrary selection of the class interval. It is still cumbersome, however, to calculate the central value $\bar{x}_{i}$ of these possibly very small class intervals, especially when one considers that many do not contribute to the estimator. We ask ourselves if it is not possible to use the observations $x_{i}$ in their place. By reducing $\Delta x$ even further until every observation is isolated in the center of its own infinitesimally small class interval, in which case $f_{i}^{*}$ would remain at the $1 / n$ level, observations $x_{i}$ can be used in place of $\bar{x}_{i}$ in equation 2.5-1 without significantly altering the basis of our estimator. Using future observations $X_{i}$ of the process in place of central-interval values, $\bar{x}_{i}$, an estimator based on an infinitesimally small interval would appear as

$$
\begin{equation*}
\bar{X}=\frac{1}{n} \sum_{i=1}^{n} X_{i} \tag{2.5-2}
\end{equation*}
$$

where $n$ is the size of the random sample, and $X_{i}, i=1, \ldots, n$, is the collection of random variables from the random sample. Equation 2.5-2 is the preferred estimator for the population mean.

A sample statistic is said to be unbiased if its expected value is equal to the population parameter that it estimates. Consider the expected value of the sample mean, derived from random variables $X_{i}, i=1, \ldots, n$. Since $E\left[X_{i}\right]=\mu_{X}$,


Figure 2.5-1


Figure 2.5-2
$E[\bar{X}]=\frac{1}{n} \sum_{i=1}^{n} E\left[X_{i}\right]=\frac{1}{n} \sum_{i=1}^{n} \mu_{X}=\mu_{X}$.

Hence, $\bar{X}$ is an unbiased estimator of $\mu_{X}$. Although examining an estimator for unbiased qualities is important, it does not necessarily insure that the estimator is the most efficient (or best) in the sense that the variance of the estimator is the smallest. It is, however, an important quality, and the variance estimator is examined for this quality in the next section.

## Problem 2.5-1

a. Recompute the sample mean for the data set in problem 2.2-2 using equation 2.5-2 as the mean estimator. How does this result vary from that of problem 2.3-1?

How do you, in light of equation 2.2-14, explain the repeat values in the data set (note that this data set represents a random sampling of a continuous random variable)?
b. With regard to a large regional aquifer, well data such as that in table $2.2-2$ represent point estimates of transmissivities. The best estimate of the effective transmissivity (the one to use in modeling the flow field) is generally considered to be the geometric mean of these point estimates. The statistic for the geometric mean is defined as

$$
\bar{T}_{\mathrm{g}}=\left(T_{1} \cdot T_{2} \cdot T_{3} \cdots \cdots T_{n}\right)^{1 / n}
$$

where $T_{i}, i=l, \ldots, n$, is a random sample from the sample space of the $T$ random variable. Letting $X_{i}=\log _{10} T_{i}$, we see that

$$
\begin{aligned}
\log _{10}\left[\bar{T}_{g}\right] & =\frac{1}{n} \log _{10}\left(T_{1} \cdot T_{2} \cdot T_{3} \cdots \cdot T_{n}\right) \\
& =\frac{1}{n} \sum_{i=1}^{n} X_{i}=\bar{X}
\end{aligned}
$$

Therefore,

$$
\bar{T}_{g}=10^{\bar{X}} .
$$

What is the geometric mean of the transmissivity data in table 2.2-2?
As a measure of the dispersion about the geometric mean, one could use the estimator

$$
D_{g}=10^{\left(\tilde{X} \pm S_{X}\right)}
$$

What is the dispersion $D_{g}$ about the geometric mean? Use results of problem 2.3-2, part b, as values for $\bar{X}$ and $S_{X}^{2}$. Considering the dispersion, how do you feel about $\bar{t}_{g}$ being the effective transmissivity of the carbonate rocks of central Pennsylvania?

### 2.5.2 Variance Estimator

As an estimator of the variance $\sigma_{X}^{2}$, consider using an estimator $S_{X}^{* 2}$ whose value $s_{X}^{* 2}$ is calculated from the equation

$$
\begin{equation*}
s_{X}^{* 2}=\sum_{i \leq x_{n} / \Delta x} f_{i}^{*}\left(\bar{x}_{i}-\mu_{X}\right)^{2}, \tag{2.5-4}
\end{equation*}
$$

which is analogous to equation 2.3-12 for the population parameter. If the class interval is taken to be small enough so as to isolate every future observation $X_{i}$ in a class interval, then equation 2.5-4 can be rewritten in terms of random observations $X_{i}, i=1, \ldots, n$, as

$$
\begin{equation*}
S_{X}^{* 2}=\frac{1}{n} \sum_{i=1}^{n}\left(X_{i}-\mu_{X}\right)^{2}, \tag{2.5-5}
\end{equation*}
$$

because $f_{i}^{*}=1 / n$ and $X_{i}=x_{i}$ prior to sampling. When the underlying population $X_{i}, i=1, \ldots, n$, is normally distributed, it can be shown that equation $2.5-5$ is the most efficient, unbiased
estimator of the variance $\sigma_{X}^{2}$ in the sense that its variance is the least of all possible unbiased estimators for $\sigma_{X}^{2}$.

On occasion, the population mean $\mu_{X}$ can be determined from other considerations, as was done in the titration experiment in section 2.2.5. However, usually $\mu_{X}$ is also unknown, requiring that $\mu_{X}$ be replaced by $\bar{X}$ in equation 2.5-5:

$$
\begin{equation*}
S_{X}^{* 2}=\frac{1}{n} \sum_{i=1}^{n}\left(X_{i}-\bar{X}\right)^{2} . \tag{2.5-6}
\end{equation*}
$$

To test whether equation 2.5-6 is an unbiased estimator of $\sigma_{X}^{2}$, the expected value of $S_{X}^{* 2}$ is determined. The actual mechanics of this operation are presented in appendix 2.11.2; only the result is presented here:

$$
\begin{equation*}
E\left(S_{X}^{* 2}\right)=\frac{n-1}{n} \sigma_{X}^{2} \tag{2.5-7}
\end{equation*}
$$

Thus, $S_{X}^{* 2}$ is a biased estimator of $\sigma_{X}^{2}$. To produce an unbiased estimator of $\sigma_{X}^{2}, S_{X}^{* 2}$ is multiplied by the ratio $n /(n-1)$ :

$$
\begin{equation*}
S_{X}^{2}=\frac{n}{n-1} S_{X}^{* 2}=\frac{1}{n-1} \sum_{i=1}^{n}\left(X_{i}-\bar{X}\right)^{2} \tag{2.5-8}
\end{equation*}
$$

This estimator is unbiased but less efficient than equation 2.5-6. However, it is the preferred estimator for small sample sizes.
Heuristically, one can argue that this adjustment to the estimator is necessary, because the population mean $\mu_{X}$ is being estimated by the sample statistic $\tilde{X}$. The sample mean will be located at the centroid of the random sample, regardless of whether its value is near that of the population mean. Thus, an equation that estimates the variance about this centroid will produce a smaller value than if the estimate were made about the population mean. The adjustment, then, merely compensates for the smaller deviates produced by using $\bar{X}$ in place of $\mu_{X}$.

Equation 2.5-8 can be rewritten, with the aid of some algebraic manipulation, to produce a slightly more useful form for hand calculations:

$$
\begin{align*}
S_{X}^{2} & =\frac{1}{n-1} \sum_{i=1}^{n}\left(X_{i}-\bar{X}\right)^{2} \\
& =\frac{1}{n-1}\left[\sum_{i=1}^{n} X_{i}^{2}-2 \bar{X} \sum_{i=1}^{n} X_{i}+n \bar{X}^{2}\right] \\
& =\frac{1}{n-1}\left[\sum_{i=1}^{n} X_{i}^{2}-n \bar{X}^{2}\right] . \tag{2.5-9}
\end{align*}
$$

The estimator for the standard deviation is taken to be the square root of the variance estimator. Values $x_{i}, i=1, \ldots, n$, obtained by sampling the population of $X$ randomly, are used in place of $X_{i}$ in equation 2.5-9 to obtain a value $s_{X}^{2}$ for the sample statistic $S_{X}^{2}$.

### 2.5.3 Estimator of Correlation Coefficient

In a manner analogous to the variance, an estimator for the covariance, and therefore the correlation coefficient, can be derived. Let $R_{X Y}$ represent the estimator for the correlation coefficient $\rho_{X Y}$; then, for paired data,

$$
\begin{equation*}
R_{X Y}=\frac{\sum_{i=1}^{n}\left(X_{i}-\bar{X}\right)\left(Y_{i}-\bar{Y}\right)}{\left[\sum_{i=1}^{n}\left(X_{i}-\bar{X}\right)^{2} \sum_{i=1}^{n}\left(Y_{i}-\bar{Y}\right)^{2}\right]^{1 / 2}} \tag{2.5-10}
\end{equation*}
$$

or, provided that $S_{X}$ and $S_{Y}$ originate from the paired data,
$R_{X Y}=\frac{1}{(n-1) S_{X} S_{Y}} \sum_{i=1}^{n}\left(X_{i}-\bar{X}\right)\left(Y_{i}-\bar{Y}\right)$,
which can be written, for purposes of hand calculation, as
$R_{X Y}=\frac{1}{(n-1) S_{X} S_{Y}}\left(\sum_{i=1}^{n} X_{i} Y_{i}-n \bar{Y} \bar{X}\right)$
where $S_{X}$ and $S_{Y}$ are calculated by taking the square root of either equations 2.5-8 or 2.5-9. The actual value $r_{X Y}$ of $R_{X Y}$ is obtained by
using values $x_{i}, i=1, \ldots, n$, from a random sample in place of $X_{i}$.

### 2.5.4 Summary

In summary, population parameters and equivalent sample statistics can be tabulated as follows:

| Population <br> parameter | Sample <br> statistic |
| :---: | :---: |
| $\mu$ | $\bar{X}=\frac{1}{n} \sum_{i=1}^{n} X_{i}$ |
| $\sigma_{X}^{2}$ | $S_{X}^{2}=\frac{1}{n-1}\left[\sum_{i=1}^{n} X_{i}^{2}-n \bar{X}^{2}\right]$ |
| $V_{X}$ | $C_{X}=S_{X} / \bar{X}$ |
| $\rho_{X Y}$ | $R_{X Y}=\frac{1}{(n-1) S_{X} S_{Y}}\left(\sum_{i=1}^{n} X_{i} Y_{i}-n \bar{Y} \bar{X}\right)$ |

These estimators can also be stated in matrix form. For instance, let $d_{i}=x_{i}-\bar{x}$; then a value for $S_{X}^{2}$ is

$$
\begin{equation*}
s_{X}^{2}=\frac{1}{n-1} \underline{d}^{T} \underline{d} \tag{2.5-13}
\end{equation*}
$$

where $\underline{d}$ is a column vector of deviates and $\underline{d}^{T}$ is its transpose. If $e_{i}=y_{i}-\bar{y}$, then a value for $R_{X Y}$ is

$$
\begin{equation*}
r_{X Y}=\frac{\underline{e}^{T} \underline{d}}{(n-1) s_{X} s_{Y}} \tag{2.5-14}
\end{equation*}
$$

Forms similar to equations 2.5-13 and 2.5-14 are commonly encountered in linear regression.

## Problem 2.5-2

Using the following data set, calculate the sample mean, variance, and standard deviation of both dissolved solids and specific conductance; then calculate their correlation coefficient.

Specific conductance and dissolved solids data for wells in carbonate rocks of Maryland [From Nutter, 1973, p. 63-68]

| Specific conductance <br> $(\mu \mathrm{mho} / \mathrm{cm})$ | Dissolved solids <br> $($ ppma) |
| :---: | :---: |
| 278 | 257 |
| 1,120 | 610 |
| 533 | 338 |
| 723 | 458 |
|  |  |
| 462 | 264 |
| 1,030 | 562 |
| 357 | 231 |
| 304 | 175 |
|  |  |
| 469 | 268 |
| 641 | 388 |
| 969 | 638 |
| 876 | 532 |
|  |  |
| 721 | 405 |
| 895 | 610 |
| 501 | 304 |
| 323 | 171 |
|  |  |
| 310 | 201 |
| 1,230 | 736 |
| 504 | 290 |
| 319 | 208 |
|  |  |
| 704 | 464 |
| 1,130 | 688 |
| 600 | 342 |

### 2.6 Transformation of Random Variables

As we have noted previously, statistics are combinations of random variables and, as such, must be random variables themselves. If the population from which the random sample is selected can be identified, then one can frequently identify the probability density functions of statistics, which are estimators of the population parameters. If a density function is identified, then one should be able to develop criteria for testing the accuracy of these estimators. With these objectives in mind, we proceed to identify density functions that result from the several types of transformations that produce statistics.
Before proceeding with this identification process, we make note of two general results from expectation which are applicable to all
random variables, regardless of their distribution. In general, if $X_{1}, X_{2}, \ldots, X_{n}$ are independent variables with identical mean $\mu_{X}$ and identical standard deviation $\sigma_{X}$, then

$$
\begin{equation*}
\bar{X}=\frac{1}{n} \sum_{i=1}^{n} X_{i} \tag{2.6-1}
\end{equation*}
$$

is also a random variable with mean

$$
\begin{equation*}
\mu_{\bar{X}}=E[\bar{X}]=\frac{1}{n} \sum_{i=1}^{n} E\left[X_{i}\right]=\mu_{X} \tag{2.6-2}
\end{equation*}
$$

and variance

$$
\begin{equation*}
\sigma_{\bar{X}}{ }^{2}=\operatorname{Var}[\bar{X}]=\frac{1}{n^{2}} \sum_{i=1}^{n} \operatorname{Var}\left[X_{i}\right]=\sigma_{X}^{2} / n . \tag{2.6-3}
\end{equation*}
$$

Equation 2.6-2 was used previously to show that $\bar{X}$ is an unbiased estimator of $\mu_{X}$, and equation $2.6-3$ is demonstrated more fully in appendix 2.11.2. Note that equation $2.6-3$ only succeeds because $\operatorname{Cov}\left[X_{i}, X_{j}\right]=0, i \neq j$; that is, the $X_{i}$ 's are independent.

The square root of equation 2.6-3, $\sigma_{\bar{X}}$, is also known as the standard error of $\bar{X}$. The standard deviation of any statistical measure is referred to as the standard error of that statistic.

### 2.6.1 Sum of Independent Normal Random Variables

Let $X_{1}$ and $X_{2}$ be independent normal random variables, $X_{1}$ with mean zero and variance one ( $N(0,1)$ ) and $X_{2}$ with mean zero and variance $k$ ( $N(0, k)$ ). How, then, is their sum distributed? To answer this question, consider

$$
P(Y \leq y)=P\left(X_{1}+X_{2} \leq y\right),
$$

where $Y=X_{1}+X_{2}$. By noting that $P\left(X_{1}+X_{2} \leq\right.$ $y)=P\left(X_{1} \leq y-X_{2}\right.$ and $\left.-\infty \leq X_{2} \leq \infty\right)$, comparison with equation 2.4-11 shows that

$$
\begin{align*}
F_{Y}(y) & =\int_{-\infty}^{\infty} \frac{y-x_{2}}{\frac{e^{-x_{1}^{2} / 2}}{\sqrt{2 \pi}} \frac{e^{-x_{2}^{2} /(2 k)}}{\sqrt{2 \pi k}} d x_{1} d x_{2},} \\
& =\int_{-\infty}^{\infty} F_{X_{1}}\left(y-x_{2}\right) \frac{e^{-x_{2}^{2} /(2 k)}}{\sqrt{2 \pi k}} d x_{2} . \tag{2.6-4}
\end{align*}
$$

To find the probability density function of $Y$, we differentiate $F_{Y}(y)$ with respect to $y$; that is

$$
\begin{aligned}
f_{Y}(y) & =\frac{d}{d y} F_{Y}(y)=\int_{-\infty}^{\infty} \frac{d}{d y} F_{X_{1}}\left(y-x_{2} \frac{e^{-x_{2}^{2} /(2 k)}}{\sqrt{2 \pi k}} d x_{2}\right. \\
& =\int_{-\infty}^{\infty} \frac{e^{-\left(y-x_{2}\right)^{2} / 2}}{\sqrt{2 \pi}} \frac{e^{-x_{2}^{2} /(2 k)}}{\sqrt{2 \pi k}} d x_{2}
\end{aligned}
$$

$$
=\frac{1}{2 \pi \sqrt{k}}
$$

$$
\begin{equation*}
\cdot \int_{-\infty}^{\infty} \exp \left[-\left(y^{2}-2 y x_{2}+(k+1) x_{2}^{2} / k\right) / 2\right] d x_{2} \tag{2.6-5}
\end{equation*}
$$

which, after some algebraic manipulation, yields

$$
\begin{align*}
& f_{Y}(y)=\frac{\exp \left[-y^{2} /(2 k+2)\right]}{2 \pi \sqrt{k}} \\
& \int_{-\infty}^{\infty} \exp \left[-\left(x_{2} \sqrt{\frac{k+1}{2 k}}-y \sqrt{\frac{k}{2 k+2}}\right)^{2}\right] d x_{2} . \quad(2.6-6)  \tag{2.6-6}\\
& \text { By letting } u=\sqrt{2}\left(x_{2} \sqrt{\frac{k+1}{2 k}}-y \sqrt{\frac{k}{2 k+2}}\right) \text {, then } \\
& f_{Y}(y)=\frac{\exp \left[-y^{2} /(2 k+2)\right]}{\sqrt{2 \pi(k+1)}} \frac{1}{\sqrt{2 \pi}} \int_{-\infty}^{\infty} e^{-u^{2} / 2} d u \\
& \quad=\frac{\exp \left[-y^{2} /(2 k+2)\right]}{\sqrt{2 \pi(k+1)}} \tag{2.6-7}
\end{align*}
$$

which follows from equation 2.2-11. Thus, the sum of two independent zero-mean normal random variables, one with variance unity and the other with variance $k$, is a normal random variable with variance $k+1$; that is, $N(0, k+1)$.
If, in the previous problem, $k$ were to equal one, then we see that the sum of two $N(0,1)$ independent random variables is a $N(0,2)$ random variable. By adding yet another independent $N(0,1)$ random variable to the previous two $N(0,1)$ random variables, induction tells us that a $N(0,3)$ random variable results. Thus, in general, the sum of $n$ independent $N(0,1)$ random variables results in a $N(0, n)$ random variable.

We are now in a position to determine the distribution of the statistic $\bar{X}$, as shown in equation
2.6-1, if $\bar{X}$ is determined from a random sample in which all the observations $X_{i}, i=1, \ldots, n$, are independent normal random variables with common mean $\mu_{X}$ and variance $\sigma_{X}^{2}$; that is, $N\left(\mu_{X}, \sigma_{X}^{2}\right)$. We note from equation $2.6-3$ that $\bar{X}$ has the standard deviation $\sigma_{X} / \sqrt{n}$. If we standarize $\bar{X}$ by its mean and standard deviation, and multiply this result by $\sqrt{n}$, then
$\sqrt{n}\left|\frac{\bar{X}-\mu_{X}}{\sigma_{X} / \sqrt{n}}\right|=\sqrt{n}\left|\frac{\sum_{i=1}^{n} X_{i}-n \mu_{X}}{\sqrt{n} \sigma_{X}}\right|=\sum_{i=1}^{n}\left|\frac{X_{i}-\mu_{X}}{\sigma_{X}}\right|$
results. We see that this new statistic is the sum of $n$ normal random variables, each with mean zero and variance one. From the previous paragraph, equation 2.6-8 must be a normal random variable with mean zero and variance $n$. To obtain a random variable with mean zero and variance unity, one would divide equation 2.6-8 by the square root of $n$. By inspection, then, the quantity $\left(\bar{X}-\mu_{X}\right) /\left(\sigma_{X} / \sqrt{n}\right)$ must be a standard normal random variable, and $\bar{X}$ must be normal with mean $\mu_{X}$ and variance $\sigma_{X}^{2} / n$. Thus, if one knew that a random sample were composed of normal random variables with a particular mean and variance, then one could investigate the probability that a future determination of the sample mean could take on a particular range of values.

### 2.6.2 The Chi-Square Distribution

We are frequently concerned with the square of a random variable and may wish to know its density function. Assuming that the random variable $X$ is normally distributed with mean zero and variance one ( $N(0,1)$ ), we may inquire as to the nature of the distribution of its square, $Y=X^{2}$. Proceeding as in the previous section, we find the cumulative distribution of $Y$ :

$$
\begin{align*}
F_{Y}(y) & =P(Y \leq y)=P\left(X^{2} \leq y\right) \\
& =P(-\sqrt{y} \leq X \leq \sqrt{y}) \\
& =\int_{-\sqrt{y}}^{\sqrt{y}} \frac{e^{-x^{2} / 2}}{\sqrt{2 \pi}} d x . \tag{2.6-9}
\end{align*}
$$

By taking the derivative of $F_{Y}(y)$, one finds the density function of $Y$ :

$$
\begin{equation*}
f_{Y}(y)=\frac{d}{d y} F_{Y}(y)=\frac{1}{\sqrt{2 \pi}} \frac{e^{-y / 2}}{\sqrt{y}}, y \geq 0 \tag{2.6-10}
\end{equation*}
$$

which is the chi-square density function with 1 degree of freedom.
Chi-square random variables have a useful additive property similar to that exhibited by independent normals. Namely, if $Y_{1}$ and $Y_{2}$ are independent chi-square random variables with degrees of freedom $\nu_{1}$ and $\nu_{2}$, then $Y_{1}+Y_{2}$ is a chi-square random variable with degrees of freedom $\nu_{1}+\nu_{2}$. Consequently, if $Y_{1}, Y_{2}, \ldots, Y_{n}$ are independent chi-square random variables each with 1 degree of freedom, then $Y_{1}+Y_{2}$ is a chi square with 2 degrees of freedom, $\left(Y_{1}+Y_{2}\right)+Y_{3}$ is a chi square with 3 degrees of freedom and, in general, $\sum_{i=1}^{n} Y_{i}$ is a chi square with $n$ degrees of freedom. Values for the cumulative distribution function of the chi-square distribution with $\nu$ degrees of freedom are to be found in table 2.10-2.

If $X_{i}, i=1, \ldots, n$, are independent normal random variables, each with mean $\mu_{X}$ and variance $\sigma_{X}^{2}$, then $\sum_{i=1}^{n}\left(\left(X_{i}-\mu_{X}\right) / \sigma_{x}\right)^{2}$ must be a chi-square random variable with $n$ degrees of freedom. This follows from the previous argument by letting $Y_{i}=\left(X_{i}-\mu_{X}\right)^{2} / \sigma_{X}^{2}$ and noting that $Y_{i}$ is the square of $N(0,1)$ random variable. Furthermore, because

$$
\sum_{i=1}^{n} \frac{\left(X_{i}-\mu_{X}\right)^{2}}{\sigma_{X}^{2}}=n \frac{\left(\bar{X}-\mu_{X}\right)^{2}}{\sigma_{X}^{2}}+\sum_{i=1}^{n} \frac{\left(X_{i}-\bar{X}\right)^{2}}{\sigma_{X}^{2}},(2.6-11)
$$

the statistic $S_{X}^{2}$ can be written in terms of this sum as

$$
\begin{equation*}
\frac{(n-1) S_{X}^{2}}{\sigma_{X}^{2}}=\sum_{i=1}^{n} \frac{\left(X_{i}-\mu_{X}\right)^{2}}{\sigma_{X}^{2}}-\frac{\left(\bar{X}-\mu_{X}\right)^{2}}{\left(\sigma_{X} / \sqrt{n}\right)^{2}} \tag{2.6-12}
\end{equation*}
$$

Under the condition that the underlying population is independent and normal, it was demonstrated in section 2.6.1 that $\bar{X}$ is normal with mean $\mu_{X}$ and standard deviation $\sigma_{X} I \sqrt{n}$. Thus, $\left(\bar{X}-\mu_{X}\right)^{2} /\left(\sigma_{X}^{2} / n\right)$, under this condition, is chi square with 1 degree of freedom. One might reasonably expect, then, that

$$
\begin{equation*}
\frac{(n-1) S_{X}^{2}}{\sigma_{X}^{2}} \sim \chi^{2}(\nu) \tag{2.6-13}
\end{equation*}
$$

is a chi-square random variable with $\nu=n-1$ degrees of freedom, which is indeed the case when the underlying population of $X_{i}$ 's are independent normal random variables.

### 2.6.3 The F Distribution

The density function for the ratio of two independent chi-square random variables can be calculated rather easily by the method used in the previous sections. However, because we have little need of the actual form of this density function, known as the $F$ distribution, we relieve the student of working through the actual calculation if he or she will accept the following statement: If $X_{1}$ is a chi-square random variable with $\nu_{1}$ degrees of freedom, and $X_{2}$ is a chi-square random variable with $\nu_{2}$ degrees of freedom, and $X_{1}$ and $X_{2}$ are independent, then

$$
\begin{equation*}
\frac{X_{1} / \nu_{1}}{X_{2} / \nu_{2}} \sim F\left(\nu_{1}, \nu_{2}\right) \tag{2.6-14}
\end{equation*}
$$

defines the $F$ distribution with $\nu_{1}$ and $\nu_{2}$ degrees of freedom.
Table $2.10-3$ is a tabulation of values of $F\left(\nu_{1}, \nu_{2}\right)$ which satisfies the probability statement

$$
\begin{equation*}
P\left(F\left(\nu_{1}, \nu_{2}\right) \leq F_{\alpha}\left(\nu_{1}, \nu_{2}\right)\right)=1-\alpha, \tag{2.6-15}
\end{equation*}
$$

where $\alpha$ equals 0.05 ; the meaning of equation $2.6-15$ is illustrated in figure 2.6-1. Note that the reciprocal of an entry $F_{\alpha}\left(\nu_{1}, \nu_{2}\right)$ in table $2.10-3$ is equal to $F_{1-\alpha}\left(\nu_{2}, \nu_{1}\right)$. That is, if equation 2.6-15 holds, then

$$
\begin{align*}
& P\left(F\left(\nu_{2}, \nu_{1}\right)>F_{\beta}\left(\nu_{2}, \nu_{1}\right)\right) \\
&= P\left(1 / F\left(\nu_{2}, \nu_{1}\right) \leq 1 / F_{\beta}\left(\nu_{2}, \nu_{1}\right)\right) \\
&= P\left(F\left(\nu_{1}, \nu_{2}\right) \leq 1 / F_{\beta}\left(\nu_{2}, \nu_{1}\right)\right) \\
&=\beta, \tag{2.6-16}
\end{align*}
$$

as, by equation 2.6-14, $1 / F\left(\nu_{2}, \nu_{1}\right)$ is an $F\left(\nu_{1}, \nu_{2}\right)$ random variable. By comparing equation 2.6-15 with the third line in equation $2.6-16$, we see that, when $\beta$ equals $1-\alpha$,

$$
\begin{equation*}
F_{\alpha}\left(\nu_{1}, \nu_{2}\right)=\frac{1}{F_{1-\alpha}\left(\nu_{2}, \nu_{1}\right)} \tag{2.6-17}
\end{equation*}
$$



Figure 2.6-1

Thus, if we wish to evaluate $F_{1-\alpha}\left(n_{1}, n_{2}\right)$ for the statement

$$
\begin{equation*}
P\left(F\left(n_{1}, n_{2}\right) \leq F_{1-\alpha}\left(n_{1}, n_{2}\right)\right)=\alpha \tag{2.6-18}
\end{equation*}
$$

where $\alpha$ is the relative mass indicated in figure 2.6-2, then we need only find $F_{\alpha}\left(\nu_{1}, \nu_{2}\right)$, where $\nu_{1}=n_{2}$ and $\nu_{2}=n_{1}$, in a table of values for the $F$ distribution and calculate its reciprocal.

As an example of a practical statistic associated with the $F\left(\nu_{1}, \nu_{2}\right)$ random variable, consider two random samples of size $n_{1}$ and $n_{2}$, which have been selected from two normal populations with variances $\sigma_{1}^{2}$ and $\sigma_{2}^{2}$, respectively. Let

$$
\begin{equation*}
X_{1}=\left(n_{1}-1\right) S_{1}^{2} / \sigma_{1}^{2} \tag{2.6-19}
\end{equation*}
$$

and

$$
\begin{equation*}
X_{2}=\left(n_{2}-1\right) S_{2}^{2} / \sigma_{2}^{2} \tag{2.6-20}
\end{equation*}
$$

where $S_{1}^{2}$ and $S_{2}^{2}$ are sample variances that are independent, because they originate from separate random samples. From equations 2.6-13 and 2.6-14, we see that

$$
\begin{equation*}
\frac{S_{1}^{2} / \sigma_{1}^{2}}{S_{2}^{2 / \sigma_{2}^{2}}} \sim F\left(\nu_{1}, \nu_{2}\right) \tag{2.6-21}
\end{equation*}
$$

is an $F\left(\nu_{1}, \nu_{2}\right)$ random variable with $\nu_{1}=n_{1}-1$ and $\nu_{2}=n_{2}-1$ degrees of freedom.

If $\sigma_{1}^{2}$ were to equal $\sigma_{2}^{2}$, then equation 2.6-21 would undergo an obvious simplification. As a case in point, consider the ratio of $\left(\bar{X}-\mu_{X}\right)^{2} /\left(\sigma_{X}^{2} / n\right)$, which is the square of a $N(0,1)$ random variable, and $(n-1) S_{X}^{2} / \sigma_{X}^{2}$, which is a
$\chi^{2}(n-1)$ random variable, where $\bar{X}$ and $S_{X}^{2}$ are statistics developed from the same population. One can show (rather arduously) that $\bar{X}$ and $S_{X}^{2}$ are independent, even though they originate from the same random sample. Thus,
$\frac{\left(\bar{X}-\mu_{X}\right)^{2} /\left(\sigma_{X}^{2} / n\right)}{S_{X}^{2} / \sigma_{X}^{2}}=\frac{\left(\bar{X}-\mu_{X}\right)^{2}}{S_{X}^{2} / n} \sim F(1, n-1)$.
The square root of equation $2.6-22$ is also known as a $T$ random variable with $n-1$ degrees of freedom. However, as the $T$ random variable is, in general, equal to the square root of an $F\left(\nu_{1}, \nu_{2}\right)$ random variable with $\nu_{1}=1$, no additional time will be devoted to it.

## Problem 2.6-1

Residuals $\varepsilon$ from a titration experiment (section 2.2.5) have the following values in moles of acid:

$$
\begin{aligned}
& -0.011,+0.003,+0.004,-0.01,+0.005 \\
& +0.014,+0.004,+0.001,-0.01,+0.003
\end{aligned}
$$

Calculate $\bar{\epsilon}$ and $s_{\varepsilon}^{2}$ from this random sample. Assume $\mu_{\varepsilon}=0$; from equation 2.6-22 derive the probability statement
$\left.P\left(-\sqrt{F_{\alpha}(1, n-1)} \leq \bar{\varepsilon} /\left(S_{\varepsilon} / \sqrt{n}\right) \leq \sqrt{F_{\alpha}(1, n-1}\right)\right)=1-\alpha$.
(Hint: $a^{2} \leq b$ is equivalent to $-\sqrt{b} \leq a \leq \sqrt{b}$ ). Find the interval corresponding to this statement when $\alpha=0.05$ (that is, go to table 2.10-3 and find $F_{\alpha}(1, n-1)$ and then calculate $-\sqrt{F_{\alpha}(1, n-1)}$ and $\sqrt{\left.F_{\alpha}(1, n-1)\right)}$. How does the value of the statistic


Figure 2.6-2
$\bar{\varepsilon} /\left(S_{\varepsilon} / \sqrt{n}\right)$, calculated from the above random sample, compare with this interval? Would you expect this result? Would it bother you if the value fell outside the interval?

### 2.7 Central Limit Theorem

An interesting and difficult-to-prove theorem of statistics and probability, known as the Central Limit Theorem, concerns the sum of random variables:

Let $X_{1}, X_{2}, \ldots, X_{n}$ be a sequence of identically distributed, independent random variables each with mean $\mu_{X}$ and variance $\sigma_{X}^{2}$. Then, the distribution of

$$
\frac{\bar{X}-\mu_{X}}{\sigma_{X} / \sqrt{n}}
$$

tends to a standard normal random variable as $n$ goes to infinity. That is,
$\lim _{n \rightarrow \infty} P\left(\frac{\bar{X}-\mu_{X}}{\sigma_{X} / \sqrt{n}} \leq a\right)=\frac{1}{\sqrt{2 \pi}} \int_{-\infty}^{a} e^{-t^{2} / 2} d t$
regardless of the distribution of $X_{i}, i=1, \ldots, n$.
The key to understanding the impact of this theorem is to realize that the underlying distribution of the $X_{i}$ random variables can be that of any random variable. Thus, for instance, a chi-square random variable with $n$ degrees of freedom is the sum of other chi-square random variables, and, if $n$ becomes large enough, the chi-square random variable approaches the
normal random variable (in fact, the normal tables are used to approximate the chi-square random variable for large values of $n$ ).
Another result of this theorem concerns the robustness of some of the distributions developed in the previous section. In particular, if the underlying distribution of $X_{i}$ 's was not normal in equation 2.6-22, this statistic would still be approximately an $F(1, n-1)$ random variable, provided the sample size $n$ were large. The argument for this statement proceeds as follows: Because the numerator of equation 2.6-22 is, when $n$ is large, the square of an approximately normal random variable (by the central limit theorem), it will tend to be a chi-square distributed random variable with 1 degree of freedom. The denominator, on the other hand, will approach unity for large $n$, because another law of probability dictates that, as $n$ becomes large, $S_{X}^{2}$ approaches $\sigma_{X}^{2}$ (this phenomenon occurs regardless of the underlying distribution). The net result is that, regardless of the distribution of the $X_{i}$ random variables, equation $2.6-22$ tends, for large $n$, to be a chi-square distributed random variable with 1 degree of freedom. However, one can show that, for large $\nu_{2}$, the $F\left(\nu_{1}, \nu_{2}\right)$ random variable (equation $2.6-14$ ) tends to a chi-square random variable whose value has been diminished by a factor of $1 / \nu_{1}$. Thus, regardless of the underlying distribution of the $X_{i}$ random variables, we say that equation 2.6-22 behaves asymptotially as an $F(1, n-1)$ random variable when $n$ is large, as both equations 2.6-22 and 2.6-14 have the same distribution for the limiting case where the degrees of freedom in the denominator become large.

### 2.8 Confidence Limits

We have already noted that statistics are random variables themselves. Now we wish to use the information developed in the previous sections concerning the form of these random variables to make a statement about the reliability of these statistics as estimators. We attempt to define an interval, based upon the statistic, such that a certain percentage of all such intervals, as constructed from different random samples, contain the population parameter that the statistic is thought to estimate. For example, if $5 / 6$ of all possible intervals constructed from repeated sampling contain the population parameter $\theta$, then there is a probability of $5 / 6$ that the interval we construct from any given random sample actually contains $\theta$ (see figure 2.8-1).

As an example of the interval-construction process, consider the statistic $\bar{X}$ and the population parameter $\mu_{X}$. We know, from the central limit theorem, that this statistic is approximately normally distributed with mean $\mu_{X}$ and standard deviation $\sigma_{X} / \sqrt{n}$ and that

$$
\begin{equation*}
\frac{\bar{X}-\mu_{X}}{\sigma_{X} / \sqrt{n}} \sim N(0,1) \tag{2.8-1}
\end{equation*}
$$

is approximately true. Of course, when $\bar{X}$ is based on a random sampling of a normal population, then equation 2.8-1 is exactly true. This standard normal random variable will be used to devise a $(1-\alpha) 100 \%$ confidence interval for
$\mu_{X}$. This objective is achieved by first looking at the probability statement

$$
\begin{equation*}
P\left(-N_{\alpha / 2}(0,1) \leq N(0,1) \leq N_{\alpha / 2}(0,1)\right)=1-\alpha \tag{2.8-2}
\end{equation*}
$$

and finding the values $\pm N_{\alpha / 2}(0,1)$ which correspond to $1-\alpha$. This probability statement says that, $(1-\alpha) 100 \%$ of the time, a value of $N(0,1)$, obtained from a repetition of the experiment, will fall between $-N_{\alpha / 2}(0,1)$ and $N_{\alpha / 2}(0,1)$. Assuming that $\sigma_{X}$ is known, and with a little help from equation $2.8-1$, equation $2.8-2$ can be rewritten as

$$
\begin{align*}
& P\left(\bar{X}-N_{\alpha / 2}(0,1) \frac{\sigma_{X}}{\sqrt{n}} \leq \mu_{X} \leq \bar{X}+\right.\left.N_{\alpha / 2}(0,1) \frac{\sigma_{X}}{\sqrt{n}}\right) \\
&=1-\alpha \tag{2.8-3}
\end{align*}
$$

This probability statement says that, $(1-\alpha) 100 \%$ of the time, the interval $\left(\bar{X}-N_{\alpha / 2}(0,1) \sigma_{X} / \sqrt{n}, \bar{X}\right.$ $\left.+N_{\alpha / 2}(0,1) \sigma_{X} / \sqrt{n}\right)$ constructed with a value of $\bar{X}=\bar{x}$ from a particular random sample will contain $\mu_{X}$. Thus,

$$
\bar{x}-N_{\alpha / 2}(0,1) \frac{\sigma_{X}}{\sqrt{n}} \leq \mu_{X} \leq \bar{x}+N_{\alpha / 2}(0,1) \frac{\sigma_{X}}{\sqrt{n}}(2.8-4)
$$

is a $(1-\alpha) 100 \%$ confidence interval for a random sample of size $n$, whose variance is known and whose sample mean $\bar{x}$ can be calculated. The investigator would be able to say that the probability is $1-\alpha$ that this interval contains $\mu_{X}$; however, in interpreting this statement, one


Figure 2.8-1
should realize that the interval is random, not $\mu_{X}$.

As an actual numeric example of an application of equation $2.8-4$, consider the following data: $\sigma_{X}=0.3, \bar{x}=2.6, n=36$. Find a $95 \%$ confidence interval. From table 2.10-1, for $\alpha / 2=0.025$, we see that $N_{\alpha / 2}(0,1)=1.96$. Hence a $95 \%$ confidence interval is

$$
2.6-(1.96)(0.3 / \sqrt{36}) \leq \mu_{X} \leq 2.6+(1.96)(0.3 / \sqrt{36})
$$

or

$$
2.50 \leq \mu_{X} \leq 2.70
$$

Thus, as 95 of 100 intervals so constructed contain the mean, there is a $95 \%$ probability that this one contains $\mu_{X}$.

If $\sigma_{X}$ is not known, equation 2.8-4 cannot be used. However, if the underlying population is nearly normal, then $\sigma_{X}^{2}$ can be estimated by $S_{X}^{2}$ as discussed in section 2.6 .2 and one can use either the $T$ distribution (with $n-1$ degrees of freedom) or the $F$ distribution as given in equation 2.6-22 to make an appropriate probability statement that can be converted to an interval on $\mu_{X}$.

## Problem 2.8-1

a. Seven gold assays from stockpiled ore are: $9.8,10.2,10.4,9.8,10.0,10.2$, and 9.6 grams per metric ton. Find a $95 \%$ confidence interval for the mean grade of the ore assuming an approximate normal distribution (hint: use equation 2.6-22; why?).
b. Write an interpretation of this interval.

### 2.9 Hypothesis Testing

Assume that you have determined by a method which you consider to be very good that a population parameter $\theta$ should take on a particular range of values. On the other hand, another independent source suggests that the parameter $\theta$ should take on a value $b$, which lies outside this range. This discrepancy is disconcerting, and you need some method of testing this independent estimate of $\theta$. You construct a hypothesis, referred to as the null hypothesis,
$H_{0}$, that $b$ is the true value of the parameter; symbolically this may be stated:

$$
H_{0}: \theta=b .
$$

Ideally, of course, you wish to reject this hypothesis, but a procedure is needed whereby you can approach the problem objectively. When a random sample is available to the investigator, hypothesis tests can provide this procedure.

If, from a random sample $X_{1}, X_{2}, \ldots, X_{n}$, a test statistic $\psi$ can be constructed which, in some manner, is a measure of $\theta$, then often a statistical method can be devised to test the probable veracity of the null hypothesis. It is assumed that the distribution of the test statistic is known under the assumption of the null hypothesis or at least can be approximated. Of course, by definition a statistic cannot contain any unknown parameters. Any population parameters which it may contain must be known either by hypothesis or some other means; otherwise $\psi$ would cease to be a statistic. The statistical test will consist of finding a critical interval with a low probability of occurrence under the null hypothesis such that, should a value of $\psi$ as determined from a random sample fall into this interval, the null hypothesis would be rejected and the alternate hypothesis $H_{1}$, which usually consists of one of the following, would be accepted:

$$
\begin{aligned}
& H_{1}: \theta>b \\
& H_{1}: \theta<b \\
& H_{1}: \theta \neq b .
\end{aligned}
$$

The alternate hypothesis chosen depends on the nature of the test. A method of intelligently selecting critical intervals must be devised before the test can be completed, for not any arbitrary interval with a small probability of occurrence will do.

### 2.9.1 Type I Error

In hypothesis-testing procedures, one is ultimately concerned with the possibility of rejecting the null hypothesis when it is true. The objective of hypothesis testing is to make as
small as possible the probability of committing this error, referred to as a type I error. That is, the probability statement

$$
\begin{equation*}
P\left(\text { reject } H_{0} \mid H_{0} \text { true }\right)=\alpha \tag{2.9-1}
\end{equation*}
$$

is constructed, and $\alpha$, the level of significance of the test, is chosen as small as the investigator deems reasonable. For continuous random variables, the probability $\alpha$ must be associated with some interval about the test statistic $\psi$, the statistic fulfilling the requirements of the null hypothesis. Generally speaking, the test statistic $\psi$ will contain an estimator $\Theta$ of the population parameter $\Theta$. If $\hat{\theta}$, as evaluated from some arbitrary random sampling of the experiment, were to have a value close to $b$, the assumed value of $\theta$ under the null hypothesis, we would not expect to reject the null hypothesis. Rather, only when this value of $\theta$ was distant from $b$ would the null hypothesis be rejected. Thus, the logical choice of an interval in $\psi$ would be one in which all possible values of $\theta$ used in the calculation of $\psi$ would be as distant as possible from $b$. When the distribution of $\psi$ has infinite tails, then this procedure will cause the interval to include one or both tails, depending on the nature of the test. This interval, whose exact starting and (or) ending point(s) will be determined by the significance level $\alpha$ of the test, will correspond to the critical region where $H_{0}$ will be rejected should a calculated value of $\psi$ fall into this region. In most cases, this procedure will cause the critical interval to obtain its maximum length at the chosen significance level $\alpha$.

When transforming equation 2.9-1 into a probability statement over the test statistic $\psi$, it is often preferable to first consider the impact of the alternate hypothesis on $\Theta$. Consider again the null hypothesis where $\theta=b$; then, should $\theta>b$ properly represent the alternate hypothesis, it is useful to consider that, heuristically if not exactly, the probability of committing a type I error can be stated $P\left(\hat{\theta}>a \mid H_{0}\right)$, where $a$ is some value of $\hat{\theta}$ such that $b<a<\infty$, thus giving one an understanding that $\hat{\theta}$ must take on, relatively, a large positive value in order for $H_{0}$ to be rejected. For this alternate hypothesis, a more accurate statement of equation 2.9-1 usually takes the form

$$
\begin{equation*}
P(\psi>c \mid \theta=b)=\alpha, \tag{2.9-2}
\end{equation*}
$$

as the distribution of $\psi$ is always assumed to be known.

To complete the above test, a value for $c$ corresponding to $\alpha$ is obtained from a table of cumulative probabilities. Values of $\psi$ less than $c$ correspond to a region where the probability of committing a type I error may not be small. Therefore, if a value of $\psi$ is less than $c$, we are forced to accept the null hypothesis to avoid committing a type I error. If this value is larger than $c$, then the probability of committing this error is considered small, and we can confidently reject $H_{0}$ at the $\alpha$ significance level.
Examples of hypothesis testing, which should clarify the actual mechanics of the procedure, are presented subsequently; however, before proceeding to these examples, note that we are frequently required to play the role of the devil's advocate in hypothesis testing. Often, we really desire to test the acceptability of a hypothesized value of a parameter. To accomplish this task, we first attempt to reject this value by making it the subject of the null hypothesis. If we cannot reject the null hypothesis, then we must admit that the hypothesized value is indeed a candidate for the true value of the parameter in question.

### 2.9.2 One-Tailed Test

As an example of developing the probability statement associated with equation 2.9-2, assume that we wish to test the hypothesis that the mean of a population is $\mu_{0}$, versus the alternate hypothesis that the population mean is greater than $\mu_{0}$ (assume that the standard deviation $\sigma_{X}$ is known):

$$
H_{0}: \mu_{X}=\mu_{0}
$$

versus

$$
H_{1}: \mu_{X}>\mu_{0}
$$

This test is referred to as a one-tailed test because the alternate hypothesis only allows for a mean greater than that indicated by $H_{0}$.

A random variable is needed whereby we may build a probability statement around the type

I error. Assume that data in the form of a random sample $X_{1}, X_{2}, \ldots, X_{n}$ from the population exist; a natural random variable for this purpose would be the estimator of the mean $\bar{X}$. As $\bar{X}$ is an estimator of $\mu_{X}$, and as the alternate hypothesis presupposes that $\underline{\mu}_{X}$ is large, one would reject $H_{0}$ if a value of $\bar{X}$, as determined from a random sample, were significantly larger than $\mu_{0}$. From the central limit theorem, assume that $\bar{X}$ is approximately normally distributed with mean $\mu_{X}$ and standard deviation $\sigma_{X} / \sqrt{n}$. Equation 2.9-2 can be represented in terms of the statistic $\bar{X}$ as

$$
\begin{equation*}
P\left(\bar{X}>a \mid H_{0}\right)=\alpha \tag{2.9-3}
\end{equation*}
$$

Although the distribution of $\bar{X}$ is known, a statistic which will allow us to incorporate the null hypothesis that $\mu_{X}=\mu_{0}$ is needed. A statistic meeting this requirement and for which values of all the parameters can be supplied is $\left(\bar{X}-\mu_{X}\right) /\left(\sigma_{X} / \sqrt{n}\right)$. With this test statistic, equation 2.9-3 can be restated as

D $P\left(\left.\frac{\bar{X}-\mu_{x}}{\sigma_{X} / \sqrt{n}}>c \right\rvert\, \mu_{X}=\mu_{0}\right)=P\left(\left.\frac{\bar{X}-\mu_{0}}{\sigma_{X} / \sqrt{n}}>N_{\alpha}(0,1) \right\rvert\,=\alpha\right.$
where $\mu_{0}$ is used in place of $\mu_{X}$ to satisfy the null hypothesis. Note that under the null hypothesis $\left(\bar{X}-\mu_{0}\right) /\left(\sigma_{X} / \sqrt{n}\right)$ is a normal random variable with mean zero and variance unity; thus, $N_{\alpha}(0,1)$ becomes the lower limit $c$ of the critical region for this test.

All possible values of the test statistic greater than $N_{\alpha}(0,1)$, where $\alpha$ is the level of significance of the test, constitute the critical region where $H_{0}$ would be rejected. In other words, $N_{\alpha}(0,1)$ is the critical value, corresponding to the limit $c$ in equation 2.9-2, which determines whether we accept or reject the null hypothesis. If a value of the test statistic is greater than $N_{\alpha}(0,1)$, we would reject $H_{0}$ at the $\alpha$ significance level. If the value were less, then we would be forced to accept the null hypothesis for fear of making a type I error.

As a sample application of this procedure, consider the data used to construct the confidence interval at the end of section 2.8: $\sigma_{X}=0.3, \bar{x}=2.6$ and $n=36$. We are told that the
population mean is really zero, a statement that seems rather dubious to us as we believe it to be some positive real number. We set the null hypothesis that $\mu_{X}$ is indeed zero, $H_{0}: \mu_{X}=0$, and hope that we can confidently disallow it. Our alternate hypothesis consists of our own belief, $H_{1}: \mu_{X}>0$. As we wish to be very sure that we do not commit a type I error, we set the level of significance of our test at $\alpha=0.025$. We determine the critical value of our test statistic from table 2.10-1: $N_{\alpha}(0,1)=1.96$. We evaluate the test statistic under the assumption of the null hypothesis: $\bar{x} /\left(\sigma_{\mathrm{X}} / \sqrt{n}\right)=52$. Because this value of the test statistic is considerably larger than the critical value, we reject $H_{0}$ at the 0.025 significance level, realizing that, although we may have committed a type I error, it is highly unlikely.

### 2.9.3 Two-Tailed Test

Suppose that $\sigma_{X}^{2}$ is unknown, but we are given a random sample $X_{1}, X_{2}, \ldots, X_{n}$ from a normal population. We wish to test the hypothesis

$$
H_{0}: \mu_{X}=\mu_{0}
$$

versus

$$
H_{1}: \mu_{X} \neq \mu_{0}
$$

at a significance level $\alpha$. This is referred to as a two-tailed test: We reject $H_{0}$ if a measure of $\mu_{X}$ is either significantly greater or less than $\mu_{0}$.

To construct this test, recall the statistic from equation 2.6-22:

$$
\begin{equation*}
\frac{\left(\bar{X}-\mu_{X}\right)^{2}}{S_{X}^{2} / n} \sim F(1, n-1) \tag{2.9-5}
\end{equation*}
$$

which is the $F(1, n-1)$ random variable. This statistic fulfills our requirement for a test statistic; it can be used to satisfy the null hypothesis, and the remaining statistics or parameters are either known or can be evaluated from a random sample. Now consider the probability of a type I error:

$$
\begin{align*}
& P\left(\text { reject } H_{0} \mid H_{0} \text { true }\right) \\
= & P\left(\bar{X}<a \mid H_{0}\right)+P\left(\bar{X}>b \mid H_{0}\right) \\
= & \alpha \tag{2.9-6}
\end{align*}
$$

where two critical values are necessary as it is possible to reject the null hypothesis if a value of $\bar{X}$ is either larger or smaller than $\mu_{0}$. In terms of the test statistic, under the condition that the null hypothesis holds, we see that

$$
\begin{align*}
& \left.P\left(\left\lvert\, \frac{\bar{X}-\mu_{0}}{S_{X} / \sqrt{n}}\right.\right)^{2}>F_{\alpha}(1, n-1)\right) \\
= & P\left(\frac{\bar{X}-\mu_{0}}{S_{X} / \sqrt{n}}<-\sqrt{F_{\alpha}(1, n-1)}\right. \\
& \text { or } \left.\frac{\bar{X}-\mu_{0}}{S_{X} / \sqrt{n}}>\sqrt{F_{\alpha}(1, n-1)}\right) \\
= & P\left(\frac{\bar{X}-\mu_{0}}{S_{X} / \sqrt{n}}<-\sqrt{F_{\alpha}(1, n-1)}\right) \\
+ & P\left(\frac{\bar{X}-\mu_{0}}{S_{X} / \sqrt{n}}>\sqrt{F_{\alpha}(1, n-1)}\right) \\
= & \alpha \tag{2.9-7}
\end{align*}
$$

which is equivalent to equation 2.9-6.
To complete the test, we need only to evaluate the test statistic with a random sample. If $\left(\bar{x}-\mu_{0}\right)^{2} /\left(s_{X}^{2} / n\right)$ be greater than $F_{\alpha}(1, n-1)$, we would reject the null hypothesis at the $\alpha$ significance level.

### 2.9.4 Type II Error

A test statistic is commonly selected for its ability to determine the probability of committing a type II error, as well as a type I error. A type II error is committed by accepting the null hypothesis when the alternate is true. By calculating the probability that the test statistic does not fall in the critical region, given that $\theta$ takes on any value other than that assumed under the null hypothesis, the probability $\beta$ of committing this error can be evaluated. Thus, for tests indicated previously, $\beta$ is a continuous function of possible values of the population parameter $\theta$, other than the value $b$ assumed under the null hypothesis. For a critical region corresponding to a given $\alpha$, a good test statistic
should produce small values of $\beta$ for hypothetical values of $\Theta$ rather distant from $b$. However, $\beta$ should increase sharply in value as possible values of $\theta$ approach $b$ and obtain a value as close to one as feasible in the immediate vicinity of $b$. Means are available for determining test statistics which, for certain tests, excel at the above characteristics, but a presentation of these methods is beyond the scope of this course. In most cases, a statistic which contains an estimator of the population parameter being tested and for which all other parameters are either known, or estimators of said parameters are contained in the statistic, will suffice as a test statistic; however, it may not be the best test statistic.

Note that if $\alpha$, the probability of committing a type I error, were made extremely small, then the null hypothesis would almost always be accepted. At first glance, one would assume that something was amiss in the hypothesis testing procedure, as it is apparently possible to bias the test by selecting an extreme value for $\alpha$. However, when the value of $\alpha$ is decreased, the probability of committing a type II error, $\beta(\Theta)$, is increased for all values $\theta$. Thus, an investigator who seeks to avoid committing a type I error by intentionally selecting a small value for $\alpha$ runs an increased risk of committing a type II error, which is equally as damaging. If need be, a plot of $\beta(\theta)$ can be made for various hypothetical values of $\alpha$ and $\theta$; this can often be a rather complicated task. A rule-of-thumb value for $\alpha$ is 0.05 , which appears to serve hypothesis test users well in most cases.

### 2.9.5 Summary of Method

To summarize, the steps for testing a hypothesis concerning a population parameter $\theta$ are:

1. Define the null hypothesis $H_{0}: \theta=\theta_{0}$.
2. Decide upon the nature of the test; that is, $H_{1}: \Theta<\Theta_{0}, H_{1}: \Theta>\Theta_{0}$ or $H_{1}: \Theta \neq \Theta_{0}$.
3. Choose a level of significance $\alpha$.
4. Select an appropriate test statistic and establish the critical region.
5. Compute the value of the statistic from a random sample of size $n$.
6. Draw conclusion of test: reject $H_{0}$ if the statistic has a value in the critical region; otherwise accept $H_{0}$.

## Problem 2.9-1

a. Set up problem 2.6-1 as a hypothesis test (do not complete the test).
b. Given two random samples from independent normal populations with the following sample statistics:

| Statistic | Random sample 1 | Random sample 2 |
| :---: | :---: | :---: |
| $n$ | 25 | 16 |
| $\bar{x}$ | 82 | 78 |
| $s_{x}$ | 8 | 7 |

test the following hypotheses at a significance level of $\alpha=0.05$ :

$$
\begin{aligned}
& H_{0}: \frac{\sigma_{1}^{2}}{\sigma_{2}^{2}}=1 \\
& H_{1}: \frac{\sigma_{1}^{2}}{\sigma_{2}^{2}}>1 .
\end{aligned}
$$

c. An outside source informs you that the stockpiled ore of problem 2.8-1 actually only assays an average of 9.8 grams per metric ton. Can you refute this claim at a significance level of 0.05 ? (Construct a hypothesis test for this purpose.)

### 2.10 Tables of Probability Distributions

[All tables modified from Walpole and Myers (1972), with permission from the publisher]

Table 2.10-1


Areas, 1 - a , Under the Normal Curve

| $N_{\alpha}(0,1)$ | 0.00 | 0.01 | 0.02 | 0.03 | 0.04 | 0.05 | 0.06 | 0.07 | 0.08 | 009 |
| :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: |
| - 3.4 | 0.0003 | 0.0003 | 0.0003 | 0.0003 | 0.0003 | 0.0003 | 0.0003 | 00003 | 0.0003 | 0.0002 |
| $-3.3$ | 0.0005 | 0.0005 | 0.0005 | 0.0004 | 0.0004 | 0.0004 | 0.0004 | 0.0004 | 0.0004 | 0.0003 |
| -3.2 | 0.0007 | 0.0007 | 0.0006 | 0.0006 | 0.0006 | 0.0006 | 00006 | 0.0005 | 0.0005 | 0.0005 |
| -3.1 | 0.0010 | 0.0009 | 0.0009 | 0.0009 | 0.0008 | 00008 | 0.0008 | 0.0008 | 0.0007 | 0.0007 |
| - 3.0 | 0.0013 | 0.0013 | 0.0013 | 0.0012 | 0.0012 | 0.0011 | 0.0011 | 0.0011 | 0.0010 | 0.0010 |
| -2.9 | 0.0019 | 0.0018 | 0.0017 | 0.0017 | 0.0016 | 00016 | 0.0015 | 00015 | 0.0014 | 0.0014 |
| $-2.8$ | 0.0026 | 0.0025 | 0.0024 | 0.0023 | 00023 | 0.0022 | 0.0021 | 0.0021 | 0.0020 | 0.0019 |
| $-2.7$ | 0.0035 | 0.0034 | 0.0033 | 0.0032 | 0.0031 | 0.0030 | 0.0029 | 0.0028 | 0.0027 | 0.0026 |
| $-26$ | 0.0047 | 0.0045 | 0.0044 | 00043 | 0.0041 | 0.0040 | 0.0039 | 0.0038 | 0.0037 | 0.0036 |
| $-2.5$ | 0.0062 | 0.0060 | 0.0059 | 0.0057 | 0.0055 | 0.0054 | 0.0052 | 0.0051 | 0.0049 | 0.0048 |
| -- 2.4 | 0.0082 | 0.0080 | 00078 | 0.0075 | 0.0073 | 00071 | 0.0069 | 0.0068 | 0.0066 | 0.0064 |
| $-2.3$ | 00107 | 00104 | 0.0102 | 00099 | 0.0096 | 0.0094 | 0.0091 | 0.0089 | 0.0087 | 0.0084 |
| $-2.2$ | 0.0139 | 00136 | 00132 | 0.0129 | 0.0125 | 0.0122 | 0.0119 | 0.0116 | 0.0113 | 0.0110 |
| - 2.1 | 0.0179 | 00174 | 0.0170 | 0.0166 | 0.0162 | 0.0158 | 0.0154 | 0.0150 | 0.0146 | 0.0143 |
| $-2.0$ | 0.0228 | 0.0222 | 0.0217 | 0.0212 | 0.0207 | 0.0202 | 0.0197 | 0.0192 | 0.0188 | 0.0183 |
| -1.9 | 0.0287 | 0.0281 | 00274 | 0.0268 | 0.0262 | 0.0256 | 0.0250 | 0.0244 | 0.0239 | 0.0233 |
| $-1.8$ | 0.0359 | 0.0352 | 0.0344 | 00336 | 0.0329 | 0.0322 | 0.0314 | 0.0307 | 0.0301 | 0.0294 |
| $-1.7$ | 0.0446 | 0.0436 | 00427 | 00418 | 0.0409 | 0.0401 | 0.0392 | 0.0384 | 0.0375 | 0.0367 |
| -1.6 | 0.0548 | 0.0537 | 00526 | 0.0516 | 0.0505 | 0.0495 | 0.0485 | 0.0475 | 00465 | 0.0455 |
| $-1.5$ | 0.0668 | 00655 | 00643 | 0.0630 | 0.0618 | 00606 | 0.0594 | 0.0582 | 0.0571 | 00559 |
| $-1.4$ | 0.0808 | 0.0793 | 0.0778 | 0.0764 | 00749 | 0.0735 | 0.0722 | 0.0708 | 0.0694 | 0.0681 |
| $-1.3$ | 00968 | 0.0951 | 0.0934 | 0.0918 | 0.0901 | 0.0885 | 0.0869 | 0.0853 | 0.0838 | 0.0823 |
| -12 | 0.1151 | 0.1131 | 0.1112 | 01093 | 0.1075 | 0.1056 | 0.1038 | 01020 | 0.1003 | 0.0985 |
| -1.1 | 0.1357 | 0.1335 | 0.1314 | 0.1292 | 0.1271 | 0.1251 | 0.1230 | 0.1210 | 0.1190 | 01170 |
| $-1.0$ | 0.1587 | 0.1562 | 01539 | 0.1515 | 0.1492 | 0.1469 | 0.1446 | 0.1423 | 0.1401 | 0.1379 |
| -09 | 0.1841 | 0.1814 | 0.1788 | 0.1762 | 0.1736 | 01711 | 0.1685 | 0.1660 | 0.1635 | 0.1611 |
| -0.8 | 0.2119 | 02090 | 0.2061 | 0.2033 | 0.2005 | 01977 | 0.1949 | 01922 | 0.1894 | 0.1867 |
| -0.7 | 0.2420 | 0.2389 | 0.2358 | 0.2327 | 0.2296 | 0.2266 | 0.2236 | 0.2206 | 0.2177 | 0.2148 |
| -06 | 0.2743 | 0.2709 | 0.2676 | 0.2643 | 0.2611 | 0.2578 | 0.2546 | 0.2514 | 0.2483 | 0.2451 |
| -05 | 0.3085 | 0.3050 | 0.3015 | 0.2981 | 0.2946 | 0.2912 | 0.2877 | 0.2843 | 0.2810 | 0.2776 |
| -0.4 | 0.3446 | 0.3409 | 0.3372 | 0.3336 | 0.3300 | 03264 | 0.3228 | 0.3192 | 0.3156 | 0.3121 |
| -0.3 | 0.3821 | 0.3783 | 0.3745 | 0.3707 | 0.3669 | 0.3632 | 0.3594 | 0.3557 | 0.3520 | 0.3483 |
| -0.2 | 0.4207 | 04168 | 0.4129 | 0.4090 | 0.4052 | 0.4013 | 0.3974 | 0.3936 | 0.3897 | 0.3859 |
| -01 | 0.4602 | 0.4562 | 0.4522 | 0.4483 | 0.4443 | 0.4404 | 0.4364 | 0.4325 | 04286 | 0.4247 |
| -0.0 | 05000 | 0.4960 | 04920 | 0.4880 | 0.4840 | 0.4801 | 0.4761 | 04721 | 0.4681 | 0.4641 |
| 0.0 | 0.5000 | 0.5040 | 0.5080 | 05120 | 0.5160 | 0.5199 | 0.5239 | 0.5279 | 0.5319 | 0.5359 |
| 0.1 | 0.5398 | 0.5438 | 05478 | 0.5517 | 0.5557 | 0.5596 | 0.5636 | 0.5675 | 0.5714 | 0.5753 |
| 0.2 | 0.5793 | 0.5832 | 0.5871 | 0.5910 | 0.5948 | 0.5987 | 0.6026 | 0.6064 | 0.6103 | 0.6141 |
| 0.3 | 0.6179 | 0.6217 | 06255 | 0.6293 | 0.6331 | 0.6368 | 0.6406 | 0.6443 | 0.6480 | 0.6517 |
| 04 | 0.6554 | 06591 | 0.6628 | 0.6664 | 0.6700 | 0.6736 | 0.6772 | 0.6808 | 0.6844 | 0.6879 |
| 0.5 | 0.6915 | 0.6950 | 0.6985 | 0.7019 | 0.7054 | 0.7088 | 07123 | 0.7157 | 07190 | 0.7224 |
| 06 | 0.7257 | 0.7291 | 0.7324 | 0.7357 | 0.7389 | 0.7422 | 0.7454 | 0.7486 | 0.7517 | 0.7549 |
| 07 | 0.7580 | 0.7611 | 0.7642 | 0.7673 | 0.7704 | 0.7734 | 0.7764 | 0.7794 | 0.7823 | 0.7852 |
| 0.8 | 07881 | 0.7910 | 07939 | 0.7967 | 0.7995 | 0.8023 | 0.8051 | 0.8078 | 0.8106 | 0.8133 |
| 09 | 0.8159 | 0.8186 | 08212 | 0.8238 | 0.8264 | 0.8289 | 0.8315 | 0.8340 | 0.8365 | 0.8389 |
| 1.0 | 08413 | 0.8438 | 0.8461 | 0.8485 | 0.8308 | 0.8531 | 0.8554 | 0.8577 | 0.8599 | 0.8621 |
| $1.1{ }^{-}$ | 0.8643 | 0.8665 | 0.8686 | 0.8708 0.8907 | 08729 0.8925 | 0.8749 0.8944 | 0.8770 0.8962 | 0.8790 0.8980 | 08810 08997 | 08830 0.9015 |
| 12 | 0.8849 | 0.8869 | 0.8888 0.9066 | 0.8907 | 0.8925 | 0.8944 | 0.8962 | 08980 09147 | 0.8997 | 0.9015 |
| 1.3 1.4 | 0.9032 0.9192 | 0.9049 0.9207 | 0.9066 0.9222 | 0.9082 0.9236 | 0.9099 0.9251 | 0.9115 0.9265 | 0.9131 0.9278 | 09147 0.9292 | 0.9162 0.9306 | 09177 0.9319 |
| 1.5 | 0.9332 | 0.9345 | 0.9357 | 09370 | 0.9382 | 0.9394 | 09406 | 0.9418 | 0.9429 | 0.9441 |
| 1.6 | 0.9452 | 0.9463 | 09474 | 0.9484 | 0.9495 | 09505 | 09515 | 0.9525 | 0.9535 | 0.9545 |
| 1.7 | 0.9554 | 0.9564 | 0.9573 | 09582 | 09591 | 0.9599 | 0.9608 | 09616 | 0.9625 | 0.9633 |
| 1.8 | 0.9641 | 0.9649 | 0.9656 | 0.9664 | 09671 | 0.9678 | 0.9686 | 09693 | 0.9699 | 0.9706 |
| 1.9 | 0.9713 | 0.9719 | 0.9726 | 0.9732 | 0.9738 | 0.9744 | 0.9750 | 0.9756 | 0.9761 | 0.9767 |
| 2.0 | 0.9772 | 0.9778 | 0.9783 | 0.9788 | 0.9793 | 0.9798 | 0.9803 | 0.9808 | 0.9812 | 0.9817 |
| 21 | 0.9821 | 0.9826 | 0.9830 | 0.9834 | 0.9838 | 0.9842 | 0.9846 | 0.9850 | 0.9854 | 0.9857 |
| 22 | 09861 | 0.9864 | 09868 | 0.9871 | 0.9875 | 0.9878 | 09881 | 0.9884 | 0.9887 | 0.9890 |
| 2.3 | 0.9893 | 0.9896 | 0.9898 | 0.9901 | 0.9904 | 09906 | 0.9909 | 0.9911 | 09913 | 0.9916 |
| 2.4 | 0.9918 | 0.9920 | 0.9922 | 0.9925 | 0.9927 | 0.9929 | 0.9931 | 0.9932 | 0.9934 | 0.9936 |
| 2.5 | 09938 | 0.9940 | 0.9941 | 0.9943 | 0.9945 | 0.9946 | 0.9948 | 09949 | 0.9951 | 0.9952 |
| 2.6 | 0.9953 | 0.9955 | 0.9956 | 0.9957 | 0.9959 | 0.9960 | 09961 | 0.9962 | 0.9963 | 0.9964 |
| 2.7 | 0.9965 | 0.9966 | 0.9967 | 0.9968 | 09969 | 0.9970 | 0.9971 | 0.9972 | 0.9973 | 0.9974 |
| 28 | 0.9974 | 0.9975 | 0.9976 | 0.9977 | 0.9977 | 0.9978 | 0.9979 | 0.9979 | 0.9980 | 0.9981 |
| 2.9 | 0.9981 | 0.9982 | 0.9982 | 0.9983 | 0.9984 | 0.9984 | 0.9985 | 0.9985 | 0.9986 | 0.9986 |
| 30 | 0.9987 | 0.9987 | 0.9987 | 0.9988 | 0.9988 | 09989 | 0.9989 | 0.9989 | 0.9990 | 0.9990 |
| 31 | 09990 | 0.9991 | 0.9991 | 0.9991 | 0.9992 | 0.9992 | 0.9992 | 0.9992 | 0.9993 | 0.9993 |
| 3.2 | 0.9993 | 0.9993 | 0.9994 | 0.9994 | 0.9994 | 0.9994 | 0.9994 | 0.9995 | 0.9995 | 0.9995 |
| 3.3 | 0.9995 | 0.9995 | 0.9995 | 0.9996 | 0.9996 | 0.9996 | 0.9996 | 0.9996 | 0.9996 | 0.9997 |
| 3.4 | 0.9997 | 0.9997 | 0.9997 | 0.9997 | 0.9997 | 0.9997 | 0.9997 | 09997 | 0.9997 | 0.9998 |

Table 2.10-2


Critical Values of the Chi-Square Distribution

| $\nu$ | ${ }^{*}$ |  |  |  |  |  |  |  |
| :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: |
|  | 0.995 | 0.99 | 0.975 | 0.95 | 0.05 | 0.025 | 0.01 | 0.005 |
| 1 | 0.04393 | $0.0^{3} 157$ | $0.0^{3} 982$ | $0.0^{2} 393$ | 3.841 | 5.024 | 6.635 | 7.879 |
| 2 | 0.0100 | 0.0201 | 0.0506 | 0.103 | 5.991 | 7.378 | 9.210 | 10.597 |
| 3 | 0.0717 | 0.115 | 0.216 | 0.352 | 7.815 | 9.348 | 11.345 | 12.838 |
| 4 | 0.207 | 0.297 | 0.484 | 0.711 | 9.488 | 11.143 | 13.277 | 14.860 |
| 5 | 0.412 | 0.554 | 0.831 | 1.145 | 11.070 | 12.832 | 15.086 | 16.750 |
| 6 | 0.676 | 0.872 | 1.237 | 1.635 | 12.592 | 14.449 | 16.812 | 18.548 |
| 7 | 0.989 | 1.239 | 1.690 | 2.167 | 14.067 | 16.013 | 18.475 | 20.278 |
| 8 | 1.344 | 1.646 | 2.180 | 2.733 | 15.507 | 17.535 | 20.090 | 21.955 |
| 9 | 1.735 | 2.088 | 2.700 | 3.325 | 16.919 | 19.023 | 21.666 | 23.589 |
| 10 | 2.156 | 2.558 | 3.247 | 3.940 | 18.307 | 20.483 | 23.209 | 25.188 |
| 11 | 2.603 | 3.053 | 3.816 | 4.575 | 19.675 | 21.920 | 24.725 | 26.757 |
| 12 | 3.074 | 3.571 | 4.404 | 5.226 | 21.026 | 23.337 | 26.217 | 28.300 |
| 13 | 3.565 | 4.107 | 5.009 | 5.892 | 22.362 | 24.736 | 27.688 | 29.819 |
| 14 | 4.075 | 4.660 | 5.629 | 6.571 | 23.685 | 26.119 | 29.141 | 31.319 |
| 15 | 4.601 | 5.229 | 6.262 | 7.261 | 24.996 | 27.488 | 30.578 | 32.801 |
| 16 | 5.142 | 5.812 | 6.908 | 7.962 | 26.296 | 28.845 | 32.000 | 34.267 |
| 17 | 5.697 | 6.408 | 7.564 | 8.672 | 27.587 | 30.191 | 33.409 | 35.718 |
| 18 | 6.265 | 7.015 | 8.231 | 9.390 | 28.869 | 31.526 | 34.805 | 37.156 |
| 19 | 6.844 | 7.633 | 8.907 | 10.117 | 30.144 | 32.852 | 36.191 | 38.582 |
| 20 | 7.434 | 8.260 | 9.591 | 10.851 | 31.410 | 34.170 | 37.566 | 39.997 |
| 21 | 8.034 | 8.897 | 10.283 | 11.591 | 32.671 | 35.479 | 38.932 | 41.401 |
| 22 | 8.643 | 9.542 | 10.982 | 12.338 | 33.924 | 36.781 | 40.289 | 42.796 |
| 23 | 9.260 | 10.196 | 11.689 | 13.091 | 35.172 | 38.076 | 41.638 | 44.181 |
| 24 | 9.886 | 10.856 | 12.401 | 13.848 | 36.415 | 39.364 | 42.980 | 45.558 |
| 25 | 10.520 | 11.524 | 13.120 | 14.611 | 37.652 | 40.646 | 44.314 | 46.928 |
| 26 | 11.160 | 12.198 | 13.844 | 15.379 | 38.885 | 41.923 | 45.642 | 48.290 |
| 27 | 11.808 | 12.879 | 14.573 | 16.151 | 40.113 | 43.194 | 46.963 | 49.645 |
| 28 | 12.461 | 13.565 | 15.308 | 16.928 | 41.337 | 44.461 | 48.278 | 50.993 |
| 29 | 13.121 | 14.256 | 16.047 | 17.708 | 42.557 | 45.722 | 49.588 | 52.336 |
| 30 | 13.787 | 14.953 | 16.791 | 18.493 | 43.773 | 46.979 | 50.892 | 53.672 |

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### 2.10 Tables of Probability Distributions-Continued

Table 2.10-3


Critical Values of the F Distribution
$F_{0.05}\left(\nu_{1}, \nu_{2}\right)$

| $\nu_{2}$ | $\nu_{1}$ |  |  |  |  |  |  |  |  |
| :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: |
|  | 1 | 2 | 3 | 4 | 5 | 6 | 7 | 8 | 9 |
| 1 | 161.4 | 199.5 | 215.7 | 224.6 | 230.2 | 234.0 | 236.8 | 238.9 | 240.5 |
| 2 | 18.51 | 19.00 | 19.16 | 19.25 | 19.30 | 19.33 | 19.35 | 19.37 | 19.38 |
| 3 | 10.13 | 9.55 | 9.28 | 9.12 | 9.01 | 8.94 | 8.89 | 8.85 | 8.81 |
| 4 | 7.71 | 6.94 | 6.59 | 6.39 | 6.26 | 6.16 | 6.09 | 6.04 | 6.00 |
| 5 | 6.61 | 5.79 | 5.41 | 5.19 | 5.05 | 4.95 | 4.88 | 4.82 | 4.77 |
| 6 | 5.99 | 5.14 | 4.76 | 4.53 | 4.39 | 4.28 | 4.21 | 4.15 | 4.10 |
| 7 | 5.59 | 4.74 | 4.35 | 4.12 | 3.97 | 3.87 | 3.79 | 3.73 | 3.68 |
| 8 | 5.32 | 4.46 | 4.07 | 3.84 | 3.69 | 3.58 | 3.50 | 3.44 | 3.39 |
| 9 | 5.12 | 4.26 | 3.86 | 3.63 | 3.48 | 3.37 | 3.29 | 3.23 | 3.18 |
| 10 | 4.96 | 4.10 | 3.71 | 3.48 | 3.33 | 3.22 | 3.14 | 3.07 | 3.02 |
| 11 | 4.84 | 3.98 | 3.59 | 3.36 | 3.20 | 3.09 | 3.01 | 2.95 | 2.90 |
| 12 | 4.75 | 3.89 | 3.49 | 3.26 | 3.11 | 3.00 | 2.91 | 2.85 | 2.80 |
| 13 | 4.67 | 3.81 | 3.41 | 3.18 | 3.03 | 2.92 | 2.83 | 2.77 | 2.71 |
| 14 | 4.60 | 3.74 | 3.34 | 3.11 | 2.96 | 2.85 | 2.76 | 2.70 | 2.65 |
| 15 | 4.54 | 3.68 | 3.29 | 3.06 | 2.90 | 2.79 | 2.71 | 2.64 | 2.59 |
| 16 | 4.49 | 3.63 | 3.24 | 3.01 | 2.85 | 2.74 | 2.66 | 2.59 | 2.54 |
| 17 | 4.45 | 3.59 | 3.20 | 2.96 | 2.81 | 2.70 | 2.61 | 2.55 | 2.49 |
| 18 | 4.41 | 3.55 | 3.16 | 2.93 | 2.77 | 2.66 | 2.58 | 2.51 | 2.46 |
| 19 | 4.38 | 3.52 | 3.13 | 2.90 | 2.74 | 2.63 | 2.54 | 2.48 | 2.42 |
| 20 | 4.35 | 3.49 | 3.10 | 2.87 | 2.71 | 2.60 | 2.51 | 2.45 | 2.39 |
| 21 | 4.32 | 3.47 | 3.07 | 2.84 | 2.68 | 2.57 | 2.49 | 2.42 | 2.37 |
| 22 | 4.30 | 3.44 | 3.05 | 2.82 | 2.66 | 2.55 | 2.46 | 2.40 | 2.34 |
| 23 | 4.28 | 3.42 | 3.03 | 2.80 | 2.64 | 2.53 | 2.44 | 2.37 | 2.32 |
| 24 | 4.26 | 3.40 | 3.01 | 2.78 | 2.62 | 2.51 | 2.42 | 2.36 | 2.30 |
| 25 | 4.24 | 3.39 | 2.99 | 2.76 | 2.60 | 2.49 | 2.40 | 2.34 | 2.28 |
| 26 | 4.23 | 3.37 | 2.98 | 2.74 | 2.59 | 2.47 | 2.39 | 2.32 | 2.27 |
| 27 | 4.21 | 3.35 | 2.96 | 2.73 | 2.57 | 2.46 | 2.37 | 2.31 | 2.25 |
| 28 | 4.20 | 3.34 | 2.95 | 2.71 | 2.56 | 2.45 | 2.36 | 2.29 | 2.24 |
| 29 | 4.18 | 3.33 | 2.93 | 2.70 | 2.55 | 2.43 | 2.35 | 2.28 | 2.22 |
| 30 | 4.17 | 3.32 | 2.92 | 2.69 | 2.53 | 2.42 | 2.33 | 2.27 | 2.21 |
| 40 | 4.08 | 3.23 | 2.84 | 2.61 | 2.45 | 2.34 | 2.25 | 2.18 | 2.12 |
| 60 | 4.00 | 3.15 | 2.76 | 2.53 | 2.37 | 2.25 | 2.17 | 2.10 | 2.04 |
| 120 | 3.92 | 3.07 | 2.68 | 2.45 | 2.29 | 2.17 | 2.09 | 2.02 | 1.96 |
| $\infty$ | 3.84 | 3.00 | 2.60 | 2.37 | 2.21 | 2.10 | 2.01 | 1.94 | 1.88 |

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Table 2.10-3 continued

| $\mathrm{F}_{0.05}\left(\nu_{1}, \nu_{2}\right)$ |  |  |  |  |  |  |  |  |  |  |
| :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: |
| $\nu_{2}$ | $\nu_{1}$ |  |  |  |  |  |  |  |  |  |
|  | 10 | 12 | 15 | 20 | 24 | 30 | 40 | 60 | 120 | $\infty$ |
| 1 | 241.9 | 243.9 | 245.9 | 248.0 | 249.1 | 250.1 | 251.1 | 252.2 | 253.3 | 254.3 |
| 2 | 19.40 | 19.41 | 19.43 | 19.45 | 19.45 | 19.46 | 19.47 | 19.48 | 19.49 | 19.50 |
| 3 | 8.79 | 8.74 | 8.70 | 8.66 | 8.64 | 8.62 | 8.59 | 8.57 | 8.55 | 8.53 |
| 4 | 5.96 | 5.91 | 5.86 | 5.80 | 5.77 | 5.75 | 5.72 | 5.69 | 5.66 | 5.63 |
| 5 | 4.74 | 4.68 | 4.62 | 4.56 | 4.53 | 4.50 | 4.46 | 4.43 | 4.40 | 4.36 |
| 6 | 4.06 | 4.00 | 3.94 | 3.87 | 3.84 | 3.81 | 3.77 | 3.74 | 3.70 | 3.67 |
| 7 | 3.64 | 3.57 | 3.51 | 3.44 | 3.41 | 3.38 | 3.34 | 3.30 | 3.27 | 3.23 |
| 8 | 3.35 | 3.28 | 3.22 | 3.15 | 3.12 | 3.08 | 3.04 | 3.01 | 2.97 | 2.93 |
| 9 | 3.14 | 3.07 | 3.01 | 2.94 | 2.90 | 2.86 | 2.83 | 2.79 | 2.75 | 2.71 |
| 10 | 2.98 | 2.91 | 2.85 | 2.77 | 2.74 | 2.70 | 2.66 | 2.62 | 2.58 | 2.54 |
| 11 | 2.85 | 2.79 | 2.72 | 2.65 | 2.61 | 2.57 | 2.53 | 2.49 | 2.45 | 2.40 |
| 12 | 2.75 | 2.69 | 2.62 | 2.54 | 2.51 | 2.47 | 2.43 | 2.38 | 2.34 | 2.30 |
| 13 | 2.67 | 2.60 | 2.53 | 2.46 | 2.42 | 2.38 | 2.34 | 2.30 | 2.25 | 2.21 |
| 14 | 2.60 | 2.53 | 2.46 | 2.39 | 2.35 | 2.31 | 2.27 | 2.22 | 2.18 | 2.13 |
| 15 | 2.54 | 2.48 | 2.40 | 2.33 | 2.29 | 2.25 | 2.20 | 2.16 | 2.11 | 2.07 |
| 16 | 2.49 | 2.42 | 2.35 | 2.28 | 2.24 | 2.19 | 2.15 | 2.11 | 2.06 | 2.01 |
| 17 | 2.45 | 2.38 | 2.31 | 2.23 | 2.19 | 2.15 | 2.10 | 2.06 | 2.01 | 1.96 |
| 18 | 2.41 | 2.34 | 2.27 | 2.19 | 2.15 | 2.11 | 2.06 | 2.02 | 1.97 | 1.92 |
| 19 | 2.38 | 2.31 | 2.23 | 2.16 | 2.11 | 2.07 | 2.03 | 1.98 | 1.93 | 1.88 |
| 20 | 2.35 | 2.28 | 2.20 | 2.12 | 2.08 | 2.04 | 1.99 | 1.95 | 1.90 | 1.84 |
| 21 | 2.32 | 2.25 | 2.18 | 2.10 | 2.05 | 2.01 | 1.96 | 1.92 | 1.87 | 1.81 |
| 22 | 2.30 | 2.23 | 2.15 | 2.07 | 2.03 | 1.98 | 1.94 | 1.89 | 1.84 | 1.78 |
| 23 | 2.27 | 2.20 | 2.13 | 2.05 | 2.01 | 1.96 | 1.91 | 1.86 | 1.81 | 1.76 |
| 24 | 2.25 | 2.18 | 2.11 | 2.03 | 1.98 | 1.94 | 1.89 | 1.84 | 1.79 | 1.73 |
| 25 | 2.24 | 2.16 | 2.09 | 2.01 | 1.96 | 1.92 | 1.87 | 1.82 | 1.77 | 1.71 |
| 26 | 2.22 | 2.15 | 2.07 | 1.99 | 1.95 | 1.90 | 1.85 | 1.80 | 1.75 | 1.69 |
| 27 | 2.20 | 2.13 | 2.06 | 1.97 | 1.93 | 1.88 | 1.84 | 1.79 | 1.73 | 1.67 |
| 28 | 2.19 2 | 2.12 | 2.04 | 1.96 | 1.91 | 1.87 | 1.82 | 1.77 | 1.71 | 1.65 |
| 29 | 2.18 | 2.10 | 2.03 | 1.94 | 1.90 | 1.85 | 1.81 | 1.75 | 1.70 | 1.64 |
| 30 | 2.16 | 2.09 | 2.01 | 1.93 | 1.89 | 1.84 | 1.79 | 1.74 | 1.68 | 1.62 |
| 40 | 2.08 | 2.00 | 1.92 | 1.84 | 1.79 | 1.74 | 1.69 | 1.64 | 1.58 | 1.51 |
| 60 | 1.99 | 1.92 | 1.84 | 1.75 | 1.70 | 1.65 | 1.59 | 1.53 | 1.47 | 1.39 |
| 120 | 1.91 | 1.83 | 1.75 | 1.66 | 1.61 | 1.55 | 1.50 | 1.43 | 1.35 | 1.25 |
| $\infty$ | 1.83 | 1.75 | 1.67 | 1.57 | 1.52 | 1.46 | 1.39 | 1.32 | 1.22 | 1.00 |

### 2.11 Appendices

### 2.11.1 Correlation of Two Linearly Related Random Variables

Consider a linear relationship between two random variables $X$ and $Y$ such that

$$
\begin{equation*}
Y=a+b X+\varepsilon \tag{2.11-1}
\end{equation*}
$$

where $\mathcal{E}$ represents a zero-mean random error (independent of $X$ ). Then

$$
\begin{equation*}
\sigma_{Y}^{2}=b^{2} \sigma_{X}^{2}+\sigma_{\varepsilon}^{2} \tag{2.11-2}
\end{equation*}
$$

because $\sigma_{X \varepsilon}=E\left[\left(X-\mu_{X}\right) \varepsilon\right]=0$. By direct calculation of $\sigma_{X Y}$ from equation 2.11-1, one obtains

$$
\begin{align*}
\sigma_{X Y} & =E\left[\left(X-\mu_{X}\right)\left(Y-\mu_{Y}\right)\right] \\
& =E\left[\left(X-\mu_{X}\right)\left(b\left(X-\mu_{X}\right)+\varepsilon\right)\right] \\
& =b \sigma_{X}^{2} . \tag{2.11-3}
\end{align*}
$$

Upon squaring both sides of equation 2.11-3 and dividing by $\sigma_{X}^{2} \sigma_{Y}^{2}$, one obtains

$$
\begin{equation*}
\rho_{X Y}^{2}=b^{2} \frac{\sigma_{X}^{2}}{\sigma_{Y}^{2}} \tag{2.11-4}
\end{equation*}
$$

which, from equation 2.11-2, can be put in the form

$$
\begin{equation*}
\rho_{X Y}^{2}=1-\frac{\sigma_{\varepsilon}^{2}}{\sigma_{Y}^{2}} \tag{2.11-5}
\end{equation*}
$$

Again from equation 2.11-2, it is seen that

$$
\begin{equation*}
b^{2} \sigma_{X}^{2}=\sigma_{Y}^{2}-\sigma_{\varepsilon}^{2} \tag{2.11-6}
\end{equation*}
$$

and as $b^{2} \sigma_{X}^{2}$ is a nonnegative quantity, $\sigma_{Y}^{2}$ must be greater than or equal to $\sigma_{\varepsilon}^{2}$. This shows that, for a linear relationship, $\rho_{X Y}^{2}$ is either less than unity or equal to one if $\sigma_{\varepsilon}^{2}$ is equal to zero.

### 2.11.2 Expected Value of Variance Estimator

The sample statistic $S_{X}^{* 2}$ is defined as

$$
\begin{equation*}
S_{X}^{* 2}=\frac{1}{n} \sum_{i=1}^{n}\left(X_{i}-\bar{X}\right)^{2} \tag{2.11-7}
\end{equation*}
$$

By taking expected values of both sides, one sees that

$$
\begin{align*}
& E\left[S_{X}^{* 2}\right]=\frac{1}{n} E\left[\sum_{i=1}^{n}\left(X_{i}-\bar{X}\right)^{2}\right] \\
= & \frac{1}{n} E\left[\sum_{i=1}^{n}\left(X_{i}-\mu_{X}-\left(\bar{X}-\mu_{X}\right)\right)^{2}\right] \\
= & \frac{1}{n}\left[\sum_{i=1}^{n} E\left[\left(X_{i}-\mu_{X}\right)^{2}\right]-E\left[n\left(\bar{X}-\mu_{X}\right)^{2}\right]\right] \tag{2.11-8}
\end{align*}
$$

where use is made of the fact that $\sum_{i=1}^{n}\left(X_{i}-\mu_{X}\right)$ $=n\left(\bar{X}-\mu_{X}\right)$. Now the second expected value in equation 2.11-8 becomes

$$
\begin{align*}
E\left[n\left(\bar{X}-\mu_{X}\right)^{2}\right] & =\frac{1}{n} E\left[\sum_{i=1}^{n}\left(X_{i}-\mu_{X}\right) \sum_{j=1}^{n}\left(X_{j}-\mu_{X}\right)\right] \\
& =\frac{1}{n} E\left[\sum_{i=1}^{n} \sum_{j=1}^{n}\left(X_{i}-\mu_{X}\right)\left(X_{j}-\mu_{X}\right)\right] \\
& =\frac{1}{n} \sum_{i=1}^{n} \sum_{j=1}^{n} \operatorname{Cov}\left[X_{i}, X_{j}\right] . \tag{2.11-9}
\end{align*}
$$

As $X_{i}$ and $X_{j}$ are randomly selected and therefore independent,

$$
\operatorname{Cov}\left[X_{i}, X_{j}\right]=\left\{\begin{array}{ll}
\sigma_{X}^{2} & i=j  \tag{2.11-10}\\
0 & i \neq j
\end{array} .\right.
$$

Thus, equation 2.11-9 becomes

$$
\begin{equation*}
E\left[n\left(\bar{X}-\mu_{X}\right)^{2}\right]=\sigma_{X}^{2} \tag{2.11-11}
\end{equation*}
$$

which allows us to write equation 2.11-8 as

$$
\begin{align*}
E\left[S_{X}^{* 2}\right] & =\frac{1}{n}\left[\left(\sum_{i=1}^{n} \sigma_{X}^{2}\right)-\sigma_{X}^{2}\right] \\
& =\frac{n-1}{n} \sigma_{X}^{2} \tag{2.11-12}
\end{align*}
$$

and demonstrates the desired result. Note that equation 2.11-11 also demonstrates that
$\sigma_{\bar{X}}{ }^{2}=\sigma_{X}^{2} / n$
as $\sigma_{\bar{X}}{ }^{2}=E\left[\left(\bar{X}-\mu_{X}\right)^{2}\right]$.

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## Additional Reading

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# 3 Regression Solution of Modeling Problems 

### 3.1 Introduction and Background

Ground-water flow models are members of a class of models known as mathematical models, in which the physical model of the flow system is replaced by mathematical expressions containing mathematical variables, parameters, and constants (Krumbein and Graybill, 1965, p. 15). Mathematical models always involve simplification of the actual (true) physical system. Krumbein and Graybill (1965, p. 15) argued that mathematical models can be classified into several types, including deterministic models, statistical models, and stochasticprocess models.

A deterministic model is one in which the dependent variable(s) can be exactly computed from an expression involving independent variables, parameters, and constants. Note that deterministic models do not have to be physically based but may instead be completely empirical. The classical and inverse flow models discussed in section 1 are of the physically based deterministic type. In contrast, a statistical model is a deterministic model that has one or more random components added. These random components frequently involve measurement or other errors, but may involve separate sources of random variability as well. Incorporation of the errors in both observed heads and estimated parameters discussed in section 1 converts the deterministic flow model into a statistical model.

The term "stochastic model" can be considered to be synonymous with the term "statistical model" (Krumbein and Graybill, 1965, p. 19). A stochastic-process model may consider random effects such as those contained in the statistical model but in addition has a stochastic process built into it. Generation of a spatially varying permeability field in an aquifer has been considered to be a stochastic process by Bakr and others (1978), Gutjahr and others (1978), and Smith and Freeze (1979a, 1979b). Recently, this type of process has been incorporated into a parameter estimation scheme for a steady-state ground-water flow model
(Kitanidis and Vomvoris, 1983). Stochasticprocess models are not considered further here.

### 3.1.1 Assumed Model Structure

Consider an experiment where two variables, $\xi$ and $Y$, are measured repeatedly. The independent variable, $\xi$, is considered to be a precisely defined quantity, whereas the dependent variable, $Y$, whose values depend upon values of the independent variable, contains some error resulting from the experimental process. A scatter diagram of the data might appear as in figure 3.1-1.
From the scatter diagram or from physical considerations, the experimentalist may decide that an appropriate model equation for the data is

$$
\begin{equation*}
Y_{i}=\beta_{1}+\beta_{2} \xi_{i}+\epsilon_{i} \tag{3.1-1}
\end{equation*}
$$

where $\beta_{1}$ and $\beta_{2}$ are the intercept and slope of the equation for a straight line, subscript $i$ represents the $i$ th observation of $(\xi, Y)$, and $\epsilon_{i}$ is the true error in $Y$ for observation $i$. The quantity $\beta_{1}+\beta_{2} \xi$ is the deterministic part of the equation (the computed value of the dependent variable), and, because $\epsilon_{i}$ is the true error, parameters $\beta_{1}$ and $\beta_{2}$ are the true parameters representing the deterministic part of the model response. True error $\epsilon_{i}$, often called a disturbbance, is a random variable and, thus, represents the stochastic part of the model response. Note that if the model is correct and no other source of bias in $\epsilon_{i}$ exists, $E\left(\epsilon_{i}\right)=0$.
Equation 3.1-1 is linear in parameters $\beta_{1}$ and $\beta_{2}$. Another example of a model equation that is linear in the parameters is

$$
\begin{equation*}
Y_{i}=\beta_{1}+\beta_{2} \xi_{i}+\beta_{3} \xi_{i}^{2}+\epsilon_{i} \tag{3.1-2}
\end{equation*}
$$

which is the equation for a second degree polynomial. In equation 3.1-2 there is still only one independent variable, $\xi$, although the equation has two terms containing $\xi$. An alternate form for equation 3.1-2 is

$$
\begin{equation*}
Y_{i}=X_{i 1} \beta_{1}+X_{i 2} \beta_{2}+X_{i 3} \beta_{3}+\epsilon_{i}, \tag{3.1-3}
\end{equation*}
$$

where

$$
\left.\begin{array}{l}
X_{1}=1  \tag{3.1-4}\\
X_{2}=\xi \\
X_{3}=\xi^{2}
\end{array}\right\}
$$



Figure 3.3-1

In general, any equation that is linear in parameters $\beta_{1}, \beta_{2}, \ldots, \beta_{p}$, where there are $p$ parameters for the system, can be written in the form

$$
\begin{equation*}
Y_{i}=X_{i 1} \beta_{1}+X_{i 2} \beta_{2}+\ldots+X_{i p} \beta_{p}+\epsilon_{i} \tag{3.1-5}
\end{equation*}
$$

where

$$
\begin{equation*}
X_{i j}=X_{i j}\left(\xi_{i 1}, \xi_{i 2}, \ldots, \xi_{i k}\right) \tag{3.1-6}
\end{equation*}
$$

is a function of $k$ independent variables that multiplies the $j$ th parameter and does not contain the parameters. Because

$$
\begin{equation*}
\frac{\partial}{\partial \beta_{j}}\left(X_{i 1} \beta_{1}+X_{i 2} \beta_{2}+\ldots+X_{i p} \beta_{p}\right)=X_{i j} \tag{3.1-7}
\end{equation*}
$$

the $X$ terms are often called sensitivity coefficients or, simply, sensitivities. They indicate the change in the model response the computed value of the dependent variable) at observation point $i$ for a unit change in parameter $\beta_{j}$. Equation 3.1-5 can be written compactly in matrix form as

$$
\begin{equation*}
\underline{Y}=\underline{X} \underline{\beta}+\underline{\epsilon} \tag{3.1-8}
\end{equation*}
$$

where

$$
\underline{Y}=\left|\begin{array}{c}
Y_{1}  \tag{3.1-9}\\
Y_{2} \\
\vdots \\
Y_{n}
\end{array}\right|
$$

$$
\begin{align*}
& \underline{X}=\left|\begin{array}{cccc}
X_{11} & X_{12} & \ldots & X_{1 p} \\
X_{21} & X_{22} & \ldots & X_{2 p} \\
\vdots & & & \vdots \\
X_{n 1} & \ldots & & X_{n p}
\end{array}\right|  \tag{3.1-10}\\
& \underline{\beta}=\left[\begin{array}{c}
\beta_{1} \\
\beta_{2} \\
\vdots \\
\beta_{p}
\end{array}\right]  \tag{3.1-11}\\
& \underline{\epsilon}=\left[\begin{array}{l}
\epsilon_{1} \\
\epsilon_{2} \\
\vdots \\
\epsilon_{n}
\end{array}\right] \tag{3.1-12}
\end{align*}
$$

and $n$ observations are assumed.
Often we are faced with models where the equation is no longer linear in the parameters, $\underline{\beta}$. For example, suppose that the model equation is

$$
\begin{equation*}
Y_{i}=\sqrt{\beta_{1}+\beta_{2}} \tan \left(\left.\frac{\sqrt{\beta_{1}+\beta_{2}}}{\beta_{1}} \xi_{i} \right\rvert\,+\epsilon_{i} .\right. \tag{3.1-13}
\end{equation*}
$$

Equation 3.1-13 cannot be reduced to the form of equation 3.1-5, and, thus, is not linear in the parameters. Equations of this type are written in the general vector form

$$
\begin{equation*}
\underline{Y}=\left\{\mid \xi_{1}, \xi_{2}, \ldots, \xi_{k} ; \beta_{1}, \beta_{2}, \ldots, \beta_{p}\right)+\underline{\epsilon} \tag{3.1-14}
\end{equation*}
$$

or, in more compact form,

$$
\begin{equation*}
\underline{Y}=f(\xi, \beta)+\underline{\epsilon} \tag{3.1-15}
\end{equation*}
$$

where $f$ is an $n$-vector, each element $f_{i}$ of which is a general function of the $k$ independent variables, $\xi_{i \ell}(\ell=1,2, \ldots, k)$, and $p$ parameters, $\beta_{j}(j=1,2, \ldots, p)$. Equation 3.1-15 incorporates equation $3.1-8$ because $3.1-8$ is simply the special case where $(\underline{\xi} \underline{\underline{\xi}}, \beta)=\underline{\underline{X}}(\underline{\underline{\xi}}) \beta$.

Some equations may be nonlinear in $\underline{\beta}$ but linear in some transformation of $\beta$. For example, the model

$$
\begin{equation*}
Y_{i}=\beta_{1}\left(\beta_{2}\right)^{X_{i}} \epsilon_{i} \tag{3.1-16}
\end{equation*}
$$

in which the error $\epsilon_{i}$ is multiplicative, is nonlinear in $\beta_{1}$ and $\beta_{2}$. However, equation 3.1-16 may be written as

$$
\begin{equation*}
\log Y_{i}=\log \beta_{1}+X_{i} \log \beta_{2}+\log \epsilon_{i} \tag{3.1-17}
\end{equation*}
$$

which is linear in $\log \beta_{1}$ and $\log \beta_{2}$ and has an additive error term. Thus, it is of the standard linear form. Equations such as 3.1-16 are frequently best utilized in their transformed, thus linearized, form. However, all model analyses (to be discussed further on) would probably be made in terms of the transformed variables, and this would have to be remembered when results were interpreted.

Types of models other than the linear and nonlinear ones discussed above also exist. Some types involve a complex model equation that cannot be solved explicitly for the dependent variable. In other cases the function $f(\underline{\xi}, \underline{\beta})$, which is assumed to be a known function of $\xi$ and $\underline{\beta}$, cannot be obtained and, therefore, must be replaced by a numerical formulation. However, the basic model structure of equation 3.1-15, where the error $\underline{\epsilon}$ in $\underline{Y}$ is assumed to be additive to a deterministic dependent variable vector, is always assumed. Additional complexities of the other types of models are handled by auxiliary equations appended to equation 3.1-15. The other models are introduced at appropriate places further on.

### 3.1.2 Least-Squares Estimation

Because the true parameter set $\underline{\beta}$ and true error set $\epsilon$ are generally unknown, the true model equation 3.1-15 must be regarded as unknown, even though the form of the model is known (or, at least assumed). We.do, however, have measurements to make up the independent variable set $\xi$ and observation set $\underline{Y}$. We would like to use these measurements and the form of the model to obtain estimates of $\underline{\beta}$ and $\underline{\epsilon}$. The method explained in the following paragraphs is based on the idea that, if estimates of $\underline{\beta}$ and $\epsilon$ can be found such that the error structure of the true model is duplicated as closely as possible, then the resulting model should, in some sense, be the best possible approximation of the true model.

Assume that all $\epsilon_{i}(i=1,2, \ldots, n)$ as random
variables have finite common variance $\sigma^{2}$ and that $\epsilon_{i}$ and $\epsilon_{j}, i \neq j$, are uncorrelated. Then

$$
\begin{equation*}
\operatorname{Var}(\underline{\underline{E}})=\underline{I} \sigma^{2} . \tag{3.1-18}
\end{equation*}
$$

The scalar variance $\sigma^{2}$ can be solved for by taking the trace of both sides of equation 3.1-18:

$$
\operatorname{tr}[\operatorname{Var}(\underline{\epsilon}]]=\operatorname{tr}(\underline{\underline{I}}) \sigma^{2}
$$

or

$$
\operatorname{tr}\left[E\left((\underline{\epsilon}-E(\underline{\epsilon}))(\underline{\epsilon}-E(\underline{\epsilon}))^{T}\right)\right]=n \sigma^{2}
$$

or

$$
E\left[\left(\underline{\epsilon}-E(\underline{\underline{( })})^{T}(\underline{\epsilon}-E(\underline{\epsilon})]\right]=n \sigma^{2}\right.
$$

from which

$$
\begin{equation*}
\sigma^{2}=\frac{E\left[(\underline{\epsilon}-E(\underline{\epsilon}))^{T}(\underline{\epsilon}-E(\underline{\epsilon}))\right]}{n} . \tag{3.1-19}
\end{equation*}
$$

Ordinarily the assumptions would be made that the model being used is the correct one and that no other source of bias in $\underline{\varepsilon}$ exists, so that $E(\underline{\epsilon})=\underline{0}$ and

$$
\begin{equation*}
\sigma^{2}=\frac{E\left(\underline{\epsilon}^{T} \underline{\epsilon}\right)}{n} \tag{3.1-20}
\end{equation*}
$$

Equation 3.1-20 indicates that the sum of squared disturbances over all observations, averaged over many sets of observations, divided by $n$ yields $\sigma^{2}$.
As indicated previously, the investigator only has available the data and the form of the model, so that $\underline{\epsilon}, \sigma^{2}$, and $\underline{\beta}$ must all be considered as unknowns. However, a good approximation of the true model would produce estimates of $\epsilon$ that, for many observations, would yield a variance approaching $\sigma^{2}$. Let $\underline{b}$ be an estimator of $\beta$. Then a linear model incorporating $\underline{b}$ is

$$
\begin{equation*}
\underline{Y}=\underline{X} \underline{b}+\underline{e} \tag{3.1-21}
\end{equation*}
$$

where the vector $\underline{e}$ is an estimate of $\underline{\epsilon}$ called the residual vector. From equations $3.1-20$ and $3.1-21$, an estimate of $\sigma^{2}$ is

$$
\begin{equation*}
\tilde{\sigma}^{2}=\frac{e^{T_{e}} \underline{e}}{n} \tag{3.1-22}
\end{equation*}
$$

Most arbitrary parameter sets are expected to yield values of $\tilde{\sigma}^{2}$ that are larger than $\sigma^{2}$ because they would yield models that do not fit the data well. Of all possible parameter sets $\underline{b}$, the one that fits the data the best and at the same time minimizes $\tilde{\sigma}^{2}$ is the set that minimizes the sum of squares function $S(b)$,

$$
\begin{equation*}
S(\underline{b})=\underline{e}^{T} \underline{e} \tag{3.1-23}
\end{equation*}
$$

with respect to $\underline{b}$. The process of finding estimates of $\sigma^{2}$ and $\underline{\beta}$ by minimizing $S(\underline{b})$ is termed least squares estimation. It is developed fully in sections 3.2 and 3.3.
Recall that equation 3.1-18 and, hence, equations 3.1-22 and 3.1-23 assume that the $\epsilon_{i}$ all come from the same distribution having variance $\sigma^{2}$ and that $\epsilon_{i}$ and $\epsilon_{j}, i \neq j$, are uncorrelated. In some instances $\epsilon_{i}$ and $\epsilon_{j}$ have different variances, say $\sigma_{i}^{2}$ and $\sigma_{j}^{2}$, and they may even be correlated so that they have nonzero covariance $\sigma_{i j}$. In this case, equation 3.1-18 must be written in the more general form

$$
\begin{equation*}
\operatorname{Var}(\underline{\epsilon})=\underline{V} \sigma^{2} \tag{3.1-24}
\end{equation*}
$$

where $\underline{\underline{V}} \sigma^{2}$ is a symmetric, positive definite variance-covariance matrix defined as

$$
\underline{\underline{V}} \sigma^{2}=\left|\begin{array}{ccccc}
\sigma_{1}^{2} & \sigma_{12} & \sigma_{13} & \ldots & \sigma_{1 n}  \tag{3.1-25}\\
\sigma_{21} & \sigma_{2}^{2} & \sigma_{23} & \ldots & \sigma_{2 n} \\
& & \ldots & & \\
\sigma_{n 1} & \sigma_{n 2} & \sigma_{n 3} & \ldots & \sigma_{n}^{2}
\end{array}\right| .
$$

In equation 3.1-24, $\sigma^{2}$ is no longer the common variance of all $\epsilon_{i}$ but is instead another type of common variance. Its exact meaning can be discerned as follows. Define $\underline{V}^{1 / 2}$ as the nonsingular symmetric matrix such that $\underline{V}^{1 / 2} \underline{\underline{V}}^{1 / 2}=\underline{\underline{V}}$. Then, from equation 3.1-24,

$$
\underline{V}^{-1 / 2} \operatorname{Var}(\underline{\epsilon}) \underline{V}^{-1 / 2}=\underline{V}^{-1 / 2} \underline{V} V^{-1 / 2} \sigma^{2}
$$

or

$$
\begin{equation*}
\operatorname{Var}\left(\underline{V}^{-1 / 2} \underline{\epsilon}\right)=\underline{I}^{2} \sigma^{2} \tag{3.1-26}
\end{equation*}
$$

from which $\sigma^{2}$ is seen to be the constant or common variance of the transformed disturbances
$V^{-1 / 2} \underline{\varepsilon}$. Equation $3.1-26$ shows that these disturbances are uncorrelated.

The more general conditions represented by equation 3.1-24 may easily be incorporated into the least squares procedure. All developments through equation 3.1-20 are repeated using $\underline{V}^{-1 / 2} \underline{\underline{E}}$ instead of $\underline{\epsilon}$. The result is

$$
\begin{align*}
\sigma^{2} & =\frac{E\left[\left(\underline{V}^{-1 / 2} \underline{\epsilon}\right)^{T}\left(\underline{V}^{-1 / 2} \epsilon\right]\right.}{n} \\
& =\frac{E\left(\underline{\epsilon}^{T} \underline{V}^{-1} \underline{\epsilon}\right)}{n} \tag{3.1-27}
\end{align*}
$$

which suggests that

$$
\begin{equation*}
S(\underline{b})=\underline{e}^{T} \underline{V}^{-1} \underline{e} \tag{3.1-28}
\end{equation*}
$$

be minimized instead of equation 3.1-23.
If equation 3.1-23 is minimized to find the parameter estimates when the more general error structure given by equation $3.1-24$ is correct, then the incorrect error structure will be reflected in parameter estimates that are less precise than if equation 3.1-28 were used. The proper sum of squares function to minimize when equation $3.1-24$ represents the correct error structure is equation 3.1-28.
Even more general cases can be postulated to yield $S(\underline{b})$ in the form

$$
\begin{equation*}
S(\underline{b})=\underline{e}^{T} \underline{\omega}_{\underline{\omega}} \tag{3.1-29}
\end{equation*}
$$

where $\underline{\omega}$ is a general symmetric positive definite weight matrix that subsumes $\underline{V}^{-1}$. To apply equation 3.1-29, $\underline{\omega}$ does not necessarily reflect the error structure of $\epsilon$. Instead it may reflect the investigator's desire to emphasize (or deemphasize) certain components of $S(\underline{b})$. Equation 3.1-29 is used in all developments to follow in which the general form is applicable.
As a final note, least-squares estimation should be viewed as more than simply a parameter estimation procedure. The development given in this section is intended to show that the procedure is an attempt to reproduce the true model structure: the variance, $\sigma^{2}$, the distribution of $\epsilon$, and $\beta$. Although it is possible to use least-squares estimation as just an algebraic process, making no assumptions
about structure, considerably more information can be gained by taking the more general approach to make model analysis, including analysis of assumptions initially made concerning model structure, an integral part of the regression process.

### 3.1.3 Inclusion of Prior Information

The model structure given by equation 3.1-15 is general. Nothing is implied about the nature of $\underline{Y}$ except that it is a dependent variable vector, in error by the amount $\epsilon$. A very general interpretation of equation $3.1-15$ is to assume that $\underline{Y}, f$, and $\underline{\epsilon}$ are each composed of two partitions, one giving sample information and one giving prior information on parameters. This viewpoint amounts to an expansion of the original formulation given in the previous section where only sample information was considered (Theil, 1963).

For example, suppose that an investigator collects data $(\xi, Y)$ on a process for which the model equation is given by equation 3.1-2. However, suppose that he or she also has developed methods to collect some data directly on parameters $\beta_{1}$ and $\beta_{3}$ and suppose that these data can be represented by the equations

$$
\left.\begin{array}{l}
P_{1}=a_{11} \beta_{1}+u_{1}  \tag{3.1-30}\\
P_{2}=a_{21} \beta_{1}+a_{23} \beta_{3}+u_{2}
\end{array}\right\}
$$

where $a_{i j}$ is a constant and $u_{i}$ is a random error. If $a_{11}=1$, then $P_{1}$ is a direct observation, subject to error $u_{1}$, of $\beta_{1}$. The entire set of equations representing the system, then, can be written
$Y_{1}=\beta_{1}+\xi_{1} \beta_{2}+\xi_{1}^{2} \beta_{3}+\epsilon_{1}$
$Y_{2}=\beta_{1}+\xi_{2} \beta_{2}+\xi_{2}^{2} \beta_{3}+\epsilon_{2}$
$Y_{n}=\beta_{1}+\xi_{n} \beta_{2}+\xi_{n}^{2} \beta_{3}+\epsilon_{n}$
$P_{1}=a_{11} \beta_{1}+u_{1}$
$P_{2}=a_{21} \beta_{1}+a_{23} \beta_{3}+u_{2}$
If $\underline{Y}, \underline{X}$, and $\underline{\epsilon}$ are augmented to include the prior information, then equation 3.1-31 is of the form 3.1-8 where
$\underline{Y}=\left|\begin{array}{l}Y_{1} \\ Y_{2} \\ \cdots \\ Y_{n} \\ P_{1} \\ P_{2}\end{array}\right| \quad \underline{\epsilon}=\left|\begin{array}{c}\epsilon_{1} \\ \epsilon_{2} \\ \cdots \\ \epsilon_{n} \\ u_{1} \\ u_{2}\end{array}\right|$

$$
\underline{X}=\left|\begin{array}{lll}
1 & \xi_{1} & \xi_{1}^{2} \\
1 & \xi_{2} & \xi_{2}^{2} \\
& \cdots & \\
1 & \xi_{n} & \xi_{n}^{2} \\
a_{11} & 0 & 0 \\
a_{21} & 0 & a_{23}
\end{array}\right|
$$

Note that the number of observations is now the number of equations giving sample information $(n)$ plus the number of equations giving prior information (2).

Although the equations giving prior information are often linear, they do not need to be. Hence, a general form of equation 3.1-15 to include the prior information may be assumed:

$$
\begin{equation*}
\underline{Y}=\{(\xi, \underline{\beta})+\underline{\epsilon} \tag{3.1-32}
\end{equation*}
$$

where

$$
\underline{Y}=\left[\begin{array}{l}
\underline{Y}_{s}  \tag{3.1-33}\\
Y_{p}
\end{array}\right] \quad\left\{=\left[\begin{array}{l}
\underline{f}_{s} \\
f_{p}
\end{array}\right] \quad \underline{\epsilon}=\left[\begin{array}{l}
\epsilon_{s} \\
\epsilon_{p}
\end{array}\right]\right.
$$

and subscripts $s$ and $p$ indicate partitions of the respective vectors pertaining to sample and prior information, respectively. Corresponding to these partitions, it is convenient to redefine $n$ as the total number of observations, $n=n_{s}+n_{p}$, where $n_{s}$ is the number of items (or equations) of sample information and $n_{p}$ is the number of items (or equations) of prior information.

To apply the least squares procedure to equation 3.1-15 as augmented by the prior information, it is assumed for now that

$$
\begin{align*}
& \operatorname{Var}\left(\epsilon_{s}\right)=\underline{V}_{s} \sigma^{2}  \tag{3.1-34}\\
& \operatorname{Var}\left(\epsilon_{p}\right)=\underline{V}_{p} \sigma^{2} \tag{3.1-35}
\end{align*}
$$

$$
\begin{equation*}
\operatorname{Cov}\left(\underline{\epsilon}_{s}, \underline{\epsilon}_{p}\right)=0 \tag{3.1-36}
\end{equation*}
$$

where $\underline{V}_{s}$ is symmetric, positive definite and of order $n_{s}$, and $V_{p}$ is symmetric, positive definite and of order $n_{p}$. Equation 3.1-36 indicates that sample disturbances $\epsilon_{s}$ and prior information disturbances $\epsilon_{p}$ are not correlated with each other. With use of equations 3.1-34 through 3.1-36, equation 3.1-24 becomes

$$
\begin{align*}
\operatorname{Var}\left[\begin{array}{l}
\underline{\epsilon}_{s} \\
\underline{\epsilon}_{p}
\end{array}\right] & =\left[\begin{array}{cc}
\operatorname{Var}\left(\underline{\epsilon}_{s}\right) & \underline{0} \\
\underline{\underline{0}} & \operatorname{Var}\left(\underline{\epsilon}_{p}\right)
\end{array}\right] \\
& =\sigma^{2}\left[\begin{array}{cc}
\underline{\underline{V}} & \underline{\underline{0}} \\
\underline{\underline{0}} & \underline{\underline{V}}_{p}
\end{array}\right] \\
& =\underline{\underline{V} \sigma^{2}} \tag{3.1-37}
\end{align*}
$$

With use of equation 3.1-37, 3.1-28 becomes

$$
\begin{array}{r}
S(\underline{b})=\underline{e}^{T} \underline{\underline{V}}^{-1} \underline{e} \\
=\left[\underline{e}_{s}^{T}, \underline{e}_{p}^{T}\right]\left[\begin{array}{ll}
\underline{V}_{s}^{-1} & \underline{\underline{0}} \\
\underline{\underline{V}} & \underline{p}
\end{array}\right]\left[\begin{array}{l}
\frac{e_{s}}{e} \\
\underline{e}_{p}
\end{array}\right] \\
=\underline{e}_{s}^{T} \underline{V}_{s}^{-1} \underline{e}_{s}+\underline{e}_{p}^{T} \underline{V}_{p}^{-1} \underline{e}_{p} \tag{3.1-38}
\end{array}
$$

where the residual vector is defined as

$$
\underline{e}=\left[\begin{array}{l}
e_{s}  \tag{3.1-39}\\
e_{p}
\end{array}\right]
$$

and

$$
\underline{V}^{-1}=\left[\begin{array}{ll}
\underline{V}_{s}^{-1} & \underline{0}  \tag{3.1-40}\\
\underline{0} & \underline{V}_{p}^{-1}
\end{array}\right]
$$

The least-squares procedure may be generalized even further by using equation $3.1-29 \mathrm{in}$ stead of equation 3.1-28 to define $S\left(\frac{b}{)}\right.$. In this case the weight matrix $\omega$ is defined $\bar{b} y$

$$
\underline{\underline{\omega}}=\left[\begin{array}{ll}
\underline{\underline{\omega}}_{s} & \underline{\underline{0}}  \tag{3.1-41}\\
\underline{\underline{0}} & \underline{\underline{\omega}}_{p}
\end{array}\right]
$$

where $\omega_{\text {s }}$ is a symmetric positive definite submatrix of order $n_{s}$ that pertains to the sample information and $\underline{\omega}_{p}$ is a symmetric positive definite submatrix of order $n_{p}$ that pertains to the prior information. Because equation 3.1-41 is of block diagonal form like equation 3.1-37, zero correlation of sample and prior information is again assumed. Thus, $S(\underline{b})$ may be written in the same form as equation 3.1-38, or as

$$
\begin{equation*}
S(\underline{b})=\underline{e}_{s}^{T} \underline{\underline{\omega}}_{s} \underline{e}_{s}+\underline{e}_{p}^{T} \underline{\underline{\omega}}_{p} \underline{e}_{p} \tag{3.1-42}
\end{equation*}
$$

## Problem 3.1-1

You are charged with a ground-water study in the vicinity of Lake Ohpupu (figure 1). Estimates of transmissivity and recharge for the confined aquifer surrounding the lake are necessary for the completion of your report. Taking advantage of the unusually colinear equipotential contours on the west side of the lake (constructed from an unbiased source, of course), you decide that estimates based on a uniform stream tube will suffice. Recharge to the aquifer is largely from precipitation and is uniform over the region. Assume that the boundary heads at the range front and the lake are imprecisely known; estimates of these parameters will also be necessary. Your project has limited funds to bore $n_{s}$ holes along the stream tube and obtain measurements of head at $n_{s}$ locations of distance, $s$, from the range front.
The steady-state flow equation for a stream tube is

$$
\begin{equation*}
\frac{d}{d s}\left(T D \frac{d h}{d s}\right)+W D=0 \tag{1}
\end{equation*}
$$

where

$$
T=\text { transmissivity }\left(\mathrm{ft}^{2} / \mathrm{d}\right) ;
$$

$W=$ recharge ( $\mathrm{ft} / \mathrm{d}$ );
$D=$ width of stream tube ( ft );
$h=$ hydraulic head (ft); and
$s=$ distance along tube from the range front (ft).

The boundary conditions are taken to be

$$
\left.\begin{array}{l}
h=h_{\mathrm{o}} \text { at } s=0  \tag{2}\\
h=h_{b} \text { at } s=s_{b}
\end{array}\right\}
$$ LAKE OHPUPU



Figure 1

By integrating the flow equation twice with respect to $s$ (assuming $T, D$, and $W$ to be constant) and using the above boundary conditions, the solution for hydraulic head $h$ along the stream tube is found to be:

$$
\begin{equation*}
h=\frac{W}{2 T}\left(s_{b}-s\right) s+h_{b} \frac{s}{s_{b}}+h_{o}\left(\frac{s_{b}-s}{s_{b}}\right) . \tag{3}
\end{equation*}
$$

Let

$$
\left.\begin{array}{lll}
\beta_{1}=h_{o} & \beta_{3}=W / T & X_{2}=\frac{s}{s_{b}}  \tag{4}\\
\beta_{2}=h_{b} & X_{1}=\frac{s_{b}-s}{s_{b}} & X_{3}=\frac{\left(s_{b}-s\right) s}{2}
\end{array}\right\}
$$

and write the above solution (equation 3) to the flow equation using the definitions of $\beta_{i}$ and $X_{i}$ ( $i=1,2,3$ ). Then write the system of $n_{s}$ linear regression equations in the three unknown parameters using matrix notation and indicate the contents of each matrix. Identify dependent variable(s), independent variable(s), sensitivities, and parameters.
Let $b_{i}$ be an estimate of $\beta_{i}$, and:
a. Assume that $\operatorname{Var}(\underline{\underline{G}})=I \sigma^{2}$. Write $S(\underline{b})$ using the matrix form of the model equation 3.1-21 with the estimated parameters $b_{i}$. Write a few terms of $\boldsymbol{S}(\underline{b})$ using algebraic notation.
b. Assume that a unique estimate of the variance of the error associated with every head observation is available and that these errors are uncorrelated. Indicate the contents of the resulting weight matrix $\underset{\omega}{\omega}=\underline{V}^{-1}$. Write $S(\underline{b})$ using the matrix form of the model equation with parameters $b_{i}$. Write a few terms of $S(\underline{b})$ using weights $\omega_{j j}$ and algebraic notation.
c. Assume case a above, except that there is a prior estimate of $h_{b}$ having a standard deviation of $\sigma_{h_{b}}$. Indicate the contents of the resulting weight matrix $\omega=V^{-1}$. Write $S(b)$ using the matrix form of the model equation with parameters $b_{i}$. Write a few terms of $S(\underline{b})$ using algebraic notation. Include the term involving the prior information.

### 3.2 Regression When the Model is Linear

### 3.2.1 Derivation of Solution

The linear model assumed is

$$
\begin{equation*}
Y=X_{1} \beta_{1}+X_{2} \beta_{2}+\ldots+X_{p} \beta_{p}+\epsilon \tag{3.2-1}
\end{equation*}
$$

where the $X_{i}$ are not functions of the parameters. If $n$ observations are used, then an equation of the form of equation $3.2-1$ is written for each observation, so that the system can be written in matrix form as

$$
\begin{equation*}
\underline{Y}=\underline{X} \underline{\beta} \underline{\underline{\epsilon}} \underline{\underline{\epsilon}} . \tag{3.2-2}
\end{equation*}
$$

To find estimates of $\underline{\beta}$ and $\underline{\epsilon}$, the weighted error sum of squares $S(\underline{b})$,

$$
\begin{align*}
S(\underline{b}) & =\underline{e}^{T} \underline{\omega} \\
& =(\underline{Y}-\underline{X} b)^{T} \underline{\underline{\varphi}}(\underline{Y}-\underline{X} b) \tag{3.23}
\end{align*}
$$

is minimized with respect to $\underline{b}$.
To minimize $S(\underline{b})$ with respect to $\underline{b}$ means to take the derivative of $S(\underline{b})$ with respect to each element of $b, b_{j}(j=1,2, \ldots, p)$, and set the results to zero, or

$$
\begin{align*}
\left.\frac{\partial}{\partial b_{j}} S(\underline{b})\right|_{\underline{b}}=\underline{\hat{b}} & \left.=\frac{\partial}{\partial b_{j}}\left[(\underline{Y}-\underline{X} \underline{b})^{T} \underline{\omega}(\underline{Y}-\underline{X})\right]\right]_{\underline{b}}=\underline{\hat{b}} \\
& =0, j=1,2, \ldots, p, \tag{3.2-4}
\end{align*}
$$

where $\left.\right|_{\underline{b}=\hat{b}}$ signifies that $\underline{\hat{b}}$ is the set of parameters that causes the derivatives of $S(\underline{b})$ to be zero. By employing the rule of differentiating a product and noting that $\omega$ is independent of $\underline{b}$, it can be seen that

$$
\left.\begin{array}{rl}
\frac{\partial}{\partial b_{j}}\left[(\underline{Y}-\underline{X b})^{T} \underline{\omega}(\underline{Y}-\underline{X} b\right.
\end{array}\right]=\left[\frac{\partial}{\partial b_{j}}(\underline{Y}-\underline{X b})^{T}\right] \underline{\omega}(\underline{Y}-\underline{X b}) .
$$

To evaluate $\frac{\partial}{\partial b_{j}}(\underline{Y}-\underline{X b})^{T}$, note that taking the derivative of a vector, $\underline{Y}-\underline{X b}$, with respect to a scalar, $b_{j}$, means taking the derivative of each
entry of the vector and, thus, yields a new vector. Therefore, because observations $\underline{Y}$ are independent of $\underline{b}$,

$$
\begin{align*}
\frac{\partial}{\partial b_{j}}\left(\underline{Y}^{-X}\right)^{T} & \left.=\frac{\partial}{\partial b_{j}}\left(\left\lvert\, \begin{array}{l}
Y_{1} \\
Y_{2} \\
\vdots \\
Y_{n}
\end{array}\right.\right]-\left[\underline{X}_{1}, \underline{X}_{2}, \ldots, \underline{X}_{p}\right]\left(\begin{array}{l}
b_{1} \\
b_{2} \\
\vdots \\
b_{p}
\end{array}\right)\right)^{T} \\
& =-\underline{X}_{j}^{T} \tag{3.2-6}
\end{align*}
$$

where $\underline{X}_{j}$ is the $j$ th column vector of matrix $\underline{X}$. Also, the first and second terms on the right side of equation 3.2-5 are equal because the transpose of a scalar is the scalar, and $\omega$ is symmetric so that $\omega^{\omega} \underline{\omega}^{T}$. Hence, transposing the second term on the right side of equation 3.2-5 gives

$$
\begin{align*}
& (\underline{Y}-\underline{X} b)^{T} \underline{\omega} \frac{\partial}{\partial b_{j}}(\underline{Y}-\underline{X})=\left[(\underline{Y}-\underline{X} \underline{b})^{T} \underline{\omega} \frac{\partial}{\partial b_{j}}(\underline{Y}-\underline{X b})\right]^{T} \\
& =\left[\frac{\partial}{\partial b_{j}}(\underline{Y}-\underline{X b})^{T}\right] \underline{\omega}(\underline{Y}-\underline{X b}) . \tag{3.2-7}
\end{align*}
$$

The combination of equstions 3.2-4 through 3.2-7 yields

$$
\begin{equation*}
-2 \underline{X}_{j}^{T}(\underline{Y}-X \hat{b})=0 \tag{3.2-8}
\end{equation*}
$$

or

$$
\begin{equation*}
\underline{X}_{j}^{T} \omega X \hat{b}=X_{j \omega}^{T} Y, j=1,2, \ldots, p . \tag{3.2-9}
\end{equation*}
$$

The system of equations implied by equation $3.2-9$ can be written as

$$
\left|\begin{array}{l}
\underline{X}_{1}^{T}  \tag{3.2-10}\\
\frac{X_{2}^{T}}{T} \\
\vdots \\
\underline{X_{p}^{T}}
\end{array}\right| \underline{\omega X \hat{b}}=\left[\left.\begin{array}{l}
\underline{X}_{1}^{T} \\
\frac{X_{2}^{T}}{\vdots} \\
\underline{X_{p}^{T}}
\end{array} \right\rvert\, \underline{\omega}\right.
$$

Each vector $\underline{X}_{j}^{T}$ in equation 3.2-10 is a row vector, so that, by definition,

$$
\left|\begin{array}{c}
\underline{X}_{1}^{T}  \tag{3.2-11}\\
\underline{X}_{2}^{T} \\
\vdots \\
\underline{X}_{p}^{T}
\end{array}\right|=\underline{X}^{T}
$$

and equation $3.2-10$ becomes

$$
\begin{equation*}
\underline{X}^{T} \underline{\omega X} \hat{\underline{b}}=\underline{X}^{T} \underline{\omega} \underline{Y} \tag{3.2-12}
\end{equation*}
$$

The set of equations symbolized in matrix form by equation 3.2-12 are called the normal equations, and parameters $\hat{b}$ are called the estimates of $\underline{\beta}$. The estimates are found from

$$
\begin{equation*}
\underline{\hat{b}}=\left(\underline{X}^{T} \underline{\omega}\right)^{-1} \underline{X}^{T} \underline{\omega} Y \tag{3.2-13}
\end{equation*}
$$

Students not comfortable with the preceding development should read Draper and Smith (1981, p. 5-17, 70-80, 85-87). This material covers fitting a straight line without and with matrix nomenclature, and then extends the results to the general linear situation. Weighted least squares (where $\omega \neq \underline{I}$ ) is covered in Draper and Smith (1981, p. 108-116).
Elements of $X$ are often of vastly differing magnitudes. Thus, when working with a calculator or computer, round-off error can cause serious errors to develop when solving equation 3.2-13. It is often useful to scale equation 3.2-12 with respect to a matrix $C$, which is a diagonal matrix defined as follows: Let $\underline{X}^{T} \underline{\omega}=\underline{A}$. Then $\underline{C}=\operatorname{diag}\left\{1 / A_{11}^{1 / 2}, 1 / A_{22}^{1 / 2}, \ldots\right.$, $1 / A_{p p}^{1 / 2}$, where $A_{i i}$ is a diagonal entry of $A$. Thus, equation 3.2-12 can be transformed to become

$$
\begin{equation*}
\underline{C}^{T} \underline{X}^{T} \underline{\omega C C^{-1}} \underline{\underline{b}}=\underline{C}^{T} \underline{X}^{T} \underline{\omega} \tag{3.2-14}
\end{equation*}
$$

or

$$
\begin{equation*}
\underline{S}^{T} \underline{\omega S} \hat{A}=\underline{S}^{T} \underline{\omega} Y \tag{3.2-15}
\end{equation*}
$$

where

$$
\begin{array}{r}
\underline{S}=X \underline{C} \\
\underline{\hat{a}}=\underline{C}^{-1} \underline{\hat{b}} \tag{3.2-17}
\end{array}
$$

The effect of the scaling is to preserve the symmetry of $X^{T} \omega X$ while at the same time to produce a matrix having all diagonal entries equal to unity. Thus, variability from entry to entry of the $S^{T} \omega S$ matrix is usually reduced considerably over that of $\underline{X}^{T} \omega X$.

It is also sometimes useful to transform equation 3.2-2 and equation 3.2-12 to incremental form. By definition

$$
\begin{equation*}
f(\underline{\xi}, \underline{b})=\underline{X b} . \tag{3.2-18}
\end{equation*}
$$

Then, subtracting equation $3.2-18$ from equation $3.2-2$ results in

$$
\begin{equation*}
\underline{Y}-f\left(\xi^{\prime}, \underline{b}\right)=X(\underline{\beta}-\underline{b})+\underline{\epsilon} \tag{3.2-19}
\end{equation*}
$$

which is an incremental linear model. To obtain the analog to equation $3.2-12$, premultiply equation $3.2-18$ by $X^{T} \omega$ and subtract the result from equation $3 . \overline{2}-\overline{12}$ to obtain

$$
\begin{equation*}
\underline{X}^{T} \underline{\omega X}(\underline{\hat{b}}-\underline{b})=\underline{X}^{T} \underline{\omega}(\underline{Y}-f(\underline{\xi}, \underline{b})) \tag{3.2-20}
\end{equation*}
$$

Equation 3.2-20 can be transformed to obtain a result analogous to equation 3.2-15:

$$
\begin{equation*}
\underline{S}^{T} \underline{\underline{\omega} S \hat{\delta}}=\underline{S}^{T} \underline{\omega}(\underline{Y}-[\mid \underline{\xi}, \underline{b})) \tag{3.2-21}
\end{equation*}
$$

where

$$
\begin{equation*}
\underline{\delta}=\underline{C}^{-1}(\underline{\hat{b}}-\underline{b}) . \tag{3.2-22}
\end{equation*}
$$

If $\underline{\hat{b}}$ as calculated initially using equations $3.2-15$ and $3.2-17$ is in error because of roundoff, then $\underline{\hat{b}}$ can be used to calculate $(\hat{\{ } \hat{\xi}, \hat{b})$ which then can be substituted into equation $3.2-21$ to calculate $\hat{\delta}$. By using equation $3.2-22$, a new improved estimate of $\hat{\hat{b}}$ can be obtained. Writing equation 3.2-2 in incremental form also provides a basis for procedures, involving statistical analysis of the model, that apply for both linear and (with restrictions) nonlinear models. These are discussed later on.

### 3.2.2 Solution Algorithm

Sequential steps to follow are:

1. Form $X^{T} \omega X$ and $X^{T} \omega Y$.
2. Transform equation $\overline{3} .2-12$ to $3.2-15$.
3. Solve equation 3.2-15 (or 3.2-21) for $\underline{\hat{a}}$.
4. Solve equation $3.2-17$ for $\hat{\hat{b}}$.

## Problem 3.2-1

a. By using $\underset{\sim}{\omega}$ from case c, problem 3.1-1, write out the normal equations used to estimate parameters $\underline{\beta}$ explicitly in sum-of-product (algebraic) form. (Use equation 3.2-12.)
b. By using either data set 1 (table 1) or 2 (table 2), generate the least squares coefficient matrix ( $\underline{X}^{T} \underline{\omega}$ ) and then compute its inverse. Do not round off any intermediate calculations or the final inverse. To aid in the calculations, table 3 gives the sums of products for the sample information from the two data sets. You must add the prior information to complete the sums of products.
c. Find the vector $\hat{\underline{b}}$. Do not round off the results.

Table 1.——Data sel 1

| $s_{j}(\mathrm{ft})$ | $X_{j 1}$ | $X_{j 2}$ | $X_{j 3}$ | Observed head, <br> $Y_{j}(\mathrm{ft})$ |
| ---: | ---: | ---: | ---: | ---: |
| 50 | 0.95 | 0.05 | 23,750 | 48.33 |
| 150 | .85 | .15 | 63,750 | 45.76 |
| 250 | .75 | .25 | 93,750 | 42.08 |
| 350 | .65 | .35 | 113,750 | 38.34 |
| 450 | .55 | .45 | 123,750 | 35.30 |
| 550 | .45 | .55 | 123,750 | 31.00 |
| 650 | .35 | .65 | 113,750 | 25.85 |
| 750 | .25 | .75 | 93,750 | 21.76 |
| 850 | .15 | .85 | 63,750 | 16.11 |
| 950 | .05 | .95 | 23,750 | 12.48 |

Assume $\sigma^{2}=0.25 \mathrm{ft}^{2}, s_{b}=1,000 \mathrm{ft}$, and prior information as follows: $h_{b}=11 \mathrm{ft}$ and $\sigma_{h_{b}}=1.1 \mathrm{ft}$.

Table 2.-Data set 2

| $s_{j}(\mathrm{ft})$ | $X_{j 1}$ | $X_{j 2}$ | $X_{j 3}$ | Observed head, <br> $Y_{j}(f \mathrm{t})$ |
| :--- | ---: | ---: | ---: | ---: |
| 100 | 0.9 | 0.1 | 45,000 | 47.13 |
| 200 | .8 | .2 | 80,000 | 44.14 |
| 300 | .7 | .3 | 105,000 | 39.89 |
| 400 | .6 | .4 | 120,000 | 36.36 |
| 500 | .5 | .5 | 125,000 | 32.48 |
| 600 | .4 | .6 | 120,000 | 29.70 |
| 700 | .3 | .7 | 105,000 | 24.33 |
| 800 | .2 | .8 | 80,000 | 19.10 |
| 900 | .1 | .9 | 45,000 | 14.96 |

Assume $\sigma^{2}=0.25 \mathrm{ft}^{2}, s_{b}=1,000 \mathrm{ft}$, and prior information as follows: $h_{b}=9.5 \mathrm{ft}$ and $\sigma_{h_{b}}=0.95 \mathrm{ft}$.

Table 3.-Sums of products of sample information

|  | Data Set 1 | Data Set 2 |
| :--- | :--- | :--- |
| $\sum_{j} X_{j 1} X_{j 1}$ | 3.3250 | 2.8500 |
| $\sum_{j} X_{j 1} X_{j 2}$ | 1.6750 | 1.6500 |
| $\sum X_{j 1} X_{j 3}$ | 418,750 | 412,500 |
| $\sum_{j} X_{j 2} X_{j 2}$ | 3.3250 | 2.8500 |
| $\sum_{j} X_{j 2} X_{j 3}$ | 418,750 | 412,500 |
| $\sum_{j} X_{j 3} X_{j 3}$ | $83,340,625,000$ | $83,325,000,000$ |
| $\sum_{j} X_{j 1} Y_{j}$ | 192.18350 | 168.2030 |
| $\sum \sum X_{j 2} Y_{j}$ | 124.82650 | 119.8870 |
| $\sum_{j} X_{j 3} Y_{j}$ | $26,879,687.5$ | $26,583,550$ |

### 3.2.3 Singularity and Conditioning

Singularity of the least-squares coefficient matrix occurs whenever columns of the sensitivity matrix, $\underline{X}$, are linearly dependent because this causes rows (or columns) of the coefficient matrix $\underline{X}^{T} \omega X$ to be linearly dependent. Linear dependence in $\underline{X}$ may be stated as

$$
\begin{equation*}
\underline{X} \underline{c}=\underline{0} \tag{3.2-23}
\end{equation*}
$$

where not all components of the vector $\underline{c}$ of order $p$ are zero. By premultiplying equation $3.2-23$ by $\underline{X}^{T} \underline{\omega}$,

$$
\begin{equation*}
\underline{X}^{T} \underline{\omega X c}=\underline{0} \tag{3.2-24}
\end{equation*}
$$

which shows that columns of $\underline{X}^{T} \omega X$ (or rows since $\underline{X}^{T} \omega X$ is symmetric) are linearly dependent. Note that transformation of $\underline{X}^{T} \omega \underline{X}$ to $\underline{S}^{T} \underline{\omega}$ alters only the form of $\underline{c}$.
Near-singularity, also referred to as illconditioning, occurs whenever the columns of $\underline{X}$ (or $\underline{S}$ ) are almost linearly dependent. Often, this condition is indicated by a high degree of correlation among two or more parameter estimates. This correlation reflects the redundancy in the problem. As a result of ill-conditioning, computed parameters can be affected greatly by accumulation of round-off error generated by solving the normal equations. Also, computed variances of the parameters, which are proportional to the diagonal elements of $\left(\underline{X}^{T} \underline{\omega}\right)^{-1}$, will be large.
A common form of ill-conditioning results if
a column of $\underline{X}$ approaches zero so that $\underline{c}=[0,0, \ldots, 1,0, \ldots, 0]^{T}$, where the one appears in the row corresponding to the zero column in $\underline{X}$. This condition indicates that the model is insensitive to the parameter corresponding to the zero column in $\underline{X}$ and that the parameter should be eliminated from the model. The problem is readily detected by examining the $\underline{X}$ matrix. Another readily detected form of ill-conditioning results if two columns of $\underline{X}$ are nearly proportional, or

$$
\begin{equation*}
\alpha \underline{X}_{i} \cong \underline{X}_{j} \tag{3.2-25}
\end{equation*}
$$

so that $\underset{\sim}{\underline{c}}=[0,0, \ldots, \alpha, 0, \ldots,-1,0, \ldots, 0]^{T}$, where $\alpha$ appears in row $i$ and -1 appears in row $j$ of c. In this case

$$
\begin{aligned}
X_{i}^{T} \underline{\omega} X_{j} & \cong \alpha X_{i}^{T} \omega X_{i} \\
& \cong \frac{1}{\alpha} \underline{X}_{j}^{T} \underline{\omega} X_{j}
\end{aligned}
$$

so that

$$
\begin{align*}
\underline{S}_{i}^{T} \underline{\omega}_{j} & =\frac{\underline{X}_{i}^{T} \underline{\omega}_{j}}{\left.\sqrt{\left(\underline{X}_{i}^{T} \underline{\omega X}_{i}\right)\left(\underline{X}_{j}^{T} \underline{\omega} \underline{X}_{j}\right.}\right)} \\
& \cong \pm 1 \tag{3.2-26}
\end{align*}
$$

Thus, to detect this problem, one need only examine the $\underline{S}^{T} \underline{\omega} \boldsymbol{S}$ matrix for an off-diagonal entry nearly equal to $\pm 1$. This type of linear dependence indicates that parameters $b_{i}$, and $b_{j}$ should be combined because the model can be written as

$$
\begin{aligned}
Y & =X_{1} b_{1}+\ldots+X_{i} b_{i}+\ldots+X_{j} b_{j}+\ldots+X_{p} b_{p} \\
& \cong X_{1} b_{1}+\ldots+X_{i}\left(b_{i}+\alpha b_{j}\right)+\ldots+X_{p} b_{p} \\
& =X_{1} b_{1}+\ldots+X_{i} b_{i}^{*}+\ldots+X_{p-1} b_{p-1}(3.2-27)
\end{aligned}
$$

where $b_{i}^{*}$ replaces $b_{i}+\alpha b_{j}$, and all subsequent variables are shifted by one so that the last variable number is $p-1$.

An excellent way to detect general illconditioned (or completely linearly dependent) problems is to orthogonalize the columns of the scaled sensitivity matrix, $\underline{S}$ (Draper and Smith, 1981, p. 275-278). If the columns are all linearly independent, then they can all be transformed
so as to be orthogonal to one another; that is, such that

$$
\begin{equation*}
\underline{Q}^{T} \underline{Q}=\underline{D} \tag{3.2-28}
\end{equation*}
$$

where $Q$ is the nonzero transformation of $S$, and $D$ is a full-rank diagonal matrix. If linear dependence exists in $S$, then equation $3.2-28$ is replaced by a similar diagonal form except that one or more diagonal entries will be zero. The technique is to successively transform columns such that each new column is orthogonal to all of the previously transformed columns. If column dependence exists, then eventually a column will be calculated that exists entirely of very small numbers (theoretically all zeros for a linearly dependent problem). This column, then, is almost (or completely) linearly dependent on one or more of the previous columns.

The transformation procedure is called GramSchmidt orthogonalization and takes the following form:

$$
\left.\begin{array}{l}
\underline{Q}_{1}=\underline{S}_{1}  \tag{3.2-29}\\
\underline{Q}_{j}=\underline{S}_{j}-\sum_{i=1}^{j-1} c_{i j} Q_{i}, j=2,3, \ldots, p
\end{array}\right\}
$$

where

$$
c_{i j}=\frac{Q_{i}^{T} \underline{S}_{i}}{\underline{Q}_{i}^{T} \underline{Q}_{i}}
$$

$Q_{i}=$ the transformed vector orthogonal to vectors already in $\underline{Q}$,
$\underline{S}_{j}=$ the next column vector of $\underline{S}$ to be transformed.

### 3.3 Regression When the Model is Nonlinear

### 3.3.1 Modified Gauss-Newton Mothod

If the model is nonlinear in the parameters but is linear in the dependent variable, then the model may be written in the standard form for nonlinear regression:

$$
\begin{equation*}
Y=f\left(\xi_{1}, \xi_{2}, \ldots, \xi_{k} ; \beta_{1}, \beta_{2}, \ldots \beta_{p}\right)+\epsilon \tag{3.3-1}
\end{equation*}
$$

Because of the nonlinearity, $f$ cannot be written in the form $f=X_{1} \beta_{1}+X_{2} \beta_{2}+\ldots+X_{p} \beta_{p}$. The case more complicated than equation 3.3-1, where the model is nonlinear in both the parameters and the dependent variable, is treated in section 6.1. When there are $n$ observations, equation 3.3-1 may be written in matrix form as

$$
\begin{equation*}
\underline{Y}=f(\xi, \underline{\beta})+\underline{\epsilon} \tag{3.3-2}
\end{equation*}
$$

or, in terms of general estimate $\underline{b}$ of $\underline{\beta}$ and estimate $\underline{e}$ of $\underline{\epsilon}$,

$$
\begin{equation*}
\underline{Y}=f(\underline{\xi}, \underline{b})+\underline{e} . \tag{3.3-3}
\end{equation*}
$$

As for the linear case, the regression solution of equation 3.3-3 is obtained by minimizing the weighted error sum of squares:

$$
\begin{align*}
S(\underline{b}) & =e^{T} \underline{\omega e} \\
& =(\underline{Y}-f(\underline{\xi}, \underline{b}))^{T} \underline{\omega}(\underline{Y}-f(\underline{\xi}, \underline{b})) \tag{3.3-4}
\end{align*}
$$

However, because equation 3.3-3 is nonlinear, solution of the problem is not as direct as it was for the linear case.
A convenient and robust solution method is obtained by linearizing equation 3.3-3 around an initial estimate of parameters, then proceeding as if the problem were linear. This yields a new set of parameters that minimizes equation 3.3-4 where $f$ is replaced by the linear approximation. The new parameters are then substituted for the initial set, and the process is repeated to yield a better set of parameters. The iterative process stops whenever the change in calculated parameters from one iteration to the next is small. At that point the minimum of equation 3.3-4 has been found.
To derive the method, first $f(\underline{\xi}, \underline{b})$ is expanded about the initial set of parameters $\underline{b}_{0}$ by using a truncated Taylor series to obtain a linear approximation for $f(\xi, b)$ :

$$
\begin{equation*}
f(\underline{\xi}, \underline{b}) \cong f\left(\xi, \underline{b}_{0}\right)+\underline{X}_{0}\left(\underline{b}-\underline{b}_{0}\right) \tag{3.3-5}
\end{equation*}
$$

where

$$
\begin{equation*}
\underline{X}_{0}=\left\{X_{i j}^{0}\right\}=\left\{\left.\frac{\partial f_{i}}{\partial b_{j}}\right|_{\underline{b}=b_{0}}\right\}(n \times p) \tag{3.3-6}
\end{equation*}
$$

and $f_{i}$ is $f$ calculated at the $i$ th observation point. The components of $\underline{X}_{0}$ are recognized as sensitivity coefficients, or simply sensitivities.

By using equation 3.3-5, equation 3.3-3 may be written as

$$
\begin{equation*}
\underline{Y}-f\left(\xi, \underline{b}_{0}\right) \cong \underline{X}_{0}\left(\underline{b}-\underline{b}_{0}\right)+\underline{e} \tag{3.3-7}
\end{equation*}
$$

which is of the form of the incremental linear model of section 3.2. Note that if the model is linear so that $f(\underline{\xi}, \underline{b})=\underline{X}(\underline{\xi}) \underline{b}$, then the truncated Taylor series and hence equation 3.3-7 are exact. In this case expansion in a Taylor series is another way of deriving the incremental linear model. If the model is nonlinear, equation 3.3-7 is the approximate (linearized) model for parameters in the vicinity of $\underline{b}_{0}$, as illustrated by a simple one-parameter example in figure 3.3-1.
An approximate best estimate of $\underline{\beta}$ (which is exact for a linear model) can be obtained by minimizing $S(\underline{b})=e^{T} \underline{\omega} \cong\left(\underline{Y}-f\left(\underline{\xi}, \underline{b}_{0}\right)-\underline{X}_{0}\left(\underline{b}-\underline{b}_{0}\right)\right)^{T}$ $\cdot \omega\left(\underline{Y}-f\left(\xi_{,}, \underline{b}_{0}\right)-X_{0}\left(\underline{b}-\bar{b}_{0}\right)\right)$ with respect to $\underline{b}$. This process is carried out exactly like it was for the linear model and yields the set of normal equations

$$
\begin{equation*}
\underline{X}_{0}^{T} \omega \underline{X}_{0} \underline{d}_{1}=\underline{X}_{0}^{T} \underline{\underline{\omega}}\left(\underline{Y}-f\left(\xi \underline{\xi}, \underline{b}_{0}\right)\right) \tag{3.3-8}
\end{equation*}
$$

where subscript 1 indicates the first approximate solution and

$$
\begin{equation*}
\underline{d}_{1}=\underline{b}_{1}-\underline{b}_{0} . \tag{3.3-9}
\end{equation*}
$$

To reduce round-off error in nonlinear regression, it is generally useful to scale equation


Figure 3.3-1
3.3-8 in the same manner as for linear regression (equations 3.2-14 through 3.2-17), although another type of scaling that is useful for some types of problems is introduced in problem 3.3-1. Scaling equation 3.3-8 produces

$$
\begin{equation*}
\underline{S}_{0}^{T} \underline{\omega}^{\omega} \underline{S}_{0} \underline{\delta}_{1}=\underline{S}_{0}^{T} \omega\left(\underline{Y}-f\left(\xi, \underline{b}_{0}\right)\right) \tag{3.3-10}
\end{equation*}
$$

where

$$
\begin{align*}
& \underline{S}_{0}=\underline{X}_{0} \underline{C}_{0}  \tag{3.3-11}\\
& \underline{\delta}=\underline{C}_{0}^{-1} \underline{d}_{1}  \tag{3.3-12}\\
& \underline{C}_{0}=\operatorname{diag}\left\{\left(A_{11}^{0}\right)^{-1 / 2},\left(A_{22}^{0}\right)^{-1 / 2}, \ldots,\left(A_{p p}^{0}\right)^{-1 / 2}\right\}  \tag{3.3-13}\\
& \underline{A}_{0}=\underline{X}_{0}^{T} \underline{\omega} X_{0}
\end{align*}
$$

and $A_{i i}^{0}$ is a diagonal component of $A_{0}$.
Because equation 3.3-8 is not exact, equation 3.3-4 is not truly minimized, and $\underline{b}_{1}$ is not actually the optimal set. Hence, $\underline{b}_{1}$ is substituted for $\underline{b}_{0}$, and the entire process is repeated to yield another, hopefully improved, estimate. As a general iteration equation, 3.3-10 may be written in the form

$$
\begin{equation*}
S_{T}^{T} \operatorname{cS}_{j} \delta_{r+1}=S_{r}^{T} \underline{\omega}\left(\underline{Y}-f\left(\xi, \underline{b}_{r}\right)\right) \tag{3.3-15}
\end{equation*}
$$

where

$$
\begin{gather*}
\underline{S}_{r}=\underline{X}_{r} C_{r}  \tag{3.3-16}\\
\underline{\delta}_{r+1}=\underline{C}_{r}^{-1}\left(\underline{b}_{r+1}-\underline{b}_{r}\right)=\underline{C}_{r}^{-1} \underline{d}_{r+1}  \tag{3.3-17}\\
\underline{C}_{r}=\operatorname{diag}\left\{\left(A_{11}^{r}\right)^{-1 / 2},\left(A_{22}^{r}\right)^{-1 / 2}, \ldots,\left(A_{p p}^{r}\right)^{-1 / 2}\right\} \tag{3.3-18}
\end{gather*}
$$

and $A_{i i}^{r}$ is defined analogously to $A_{i i}^{0}$. As the process converges, $\delta_{r+1} \rightarrow 0$ and equation 3.3-7 becomes $\underline{Y}-f(\underline{\xi}, \underline{b})=\underline{e}$. At the same point $\underline{\hat{b}} \mathrm{~min}$ imizes $S(\underline{b})$ in equation $3.3-4$, or $S(\underline{b})=(\underline{Y}-f$ $(\underline{\xi}, \underline{b}))^{T} \underline{\omega}(\underline{Y}-f(\xi, \underline{b}, \underline{b})$, which is a minimum for the nonlinear equation. This process for finding the minimum of $S(\underline{b})$ is known as the Gauss-Newton method.

A sketch of progression of the iterations to the minimum for a hypothetical two-parameter problem is given in figure 3.3-2. As diagrammed


Figure 3.3-2
in figure 3.3-2, the solution often does not progress directly toward the minimum.
Modifications to the Basic Procedure. It is well known that the Gauss-Newton method as defined by equation $3.3-15$, does not always converge. To help induce convergence, a damping parameter, $\rho$, is introduced by modifying equation 3.3-17 to give

$$
\begin{equation*}
\underline{b}_{r+1}=\rho \underline{d}_{r+1}+\underline{b}_{r} . \tag{3.3-19}
\end{equation*}
$$

where $\underline{d}_{r+1}=\underline{C}_{r} \underline{\delta}_{r+1}$. If $0<\rho<1$, the changes in computed parameters are less than would result for $\rho=1$; thus the method is an interpolation method. Similarly, if $\rho>1$, the method is an extrapolation method.
Inspection of equation 3.3-19 reveals that $\rho$ changes the magnitude of the displacement from $\underline{b}_{r}$ to $\underline{b}_{r+1}$. However, because all components of the displacement vector $\underline{d}_{r+1}$ are scaled by the same multiplier $\rho$, the direction of the displacement vector is not altered. If the displacement vector $\underline{d}_{r+1}$ is oriented in a direction nearly parallel to a contour in the sum of squares surface ( $S(\underline{b})$ ), then little, if any, improvement (in terms of reducing $S(b)$ ) can result from solution of equation 3.3-15. In this case it would be desirable to alter the direction of $\underline{a}_{r+1}$ to point closer to a down-gradient direction. For example, in figure 3.3-3 vector $d_{r+1}$
yields no improvement in estimates $\underline{b}_{r}$, but vector $\underline{d}_{r+1}^{\prime}$ yields a significant improvement in the estimates.

A modification that accomplishes the desired alteration of direction of $\underline{d}_{r+1}$, and reduces its magnitude as well, consists of adding a positive parameter $\mu$, known as the Marquardt parameter (Marquardt, 1963), to the main diagonal of the coefficient matrix $S_{-}^{T} S_{S}$, of equation 3.3-15. Scaling is needed so that $\mu$ can have the same effect on each entry of the main diagonal of the coefficient matrix. The scaling accomplishes this effect because each entry of the main diagonal of the scaled matrix $S_{-}^{T} \omega_{S}$, is unity. Mathematically, the Marquardt modification can be stated as follows.

$$
\begin{equation*}
\left(\underline{S}_{r}^{T} \underline{\underline{\omega}} \underline{S}_{r}+\mu \bar{I}\right) \delta_{r+1}=S_{r}^{T} \underline{\omega}\left(\underline{Y}-f\left(\underline{\xi}, \underline{b_{r}}\right)\right) \tag{3.3-20}
\end{equation*}
$$

Solution Algorithm. The sequential steps implementing the modified Gauss-Newton procedure are:

1. Calculate $f\left(\xi_{2}, b_{r}\right)$ and $\underset{\sim}{S}$, using initial parameters $\underline{b}_{r}$ and the combination of equation 3.3-6 (with index $r$ replacing 0 ) and equation 3.3-16.
2. Solve equation $3.3-20$ for $\hat{\delta}_{+1}$.
3. Solve equation $3.3-17$ for $\underline{d}_{r+1}$.
4. Solve equation 3.3-19 for $\underline{b}_{r+1}$.
5. Test to determine if $\left|d_{i}^{r+1} / c\right|>\epsilon$, where $\epsilon$ is


Figure 3.3-3
a small number such as $0.01, c=b_{i}^{r}$ for $b_{i}^{r} \neq 0$, and $c=1$ for $b_{i}^{r}=0$.
6. If $\left|d_{i}^{r+1} / c\right|>\epsilon$, increment $r$ by one and return to 1 . If not, then the process has converged.

## Problem 3.3-1

The Theis equation for flow to a well in a confined, nonleaky aquifer is

$$
\begin{equation*}
s=\frac{Q}{4 \pi T} \int_{r_{0}^{2} S / 4 T t}^{\infty} \frac{e^{-z}}{z} d z \tag{1}
\end{equation*}
$$

where $t=$ time (d);
$s=$ drawdown ( ft );
$r_{0}=$ radial distance to observation well (ft);
$Q=$ discharge ( $\mathrm{ft}^{3} / \mathrm{s}$ );
$T=$ transmissivity ( $\mathrm{ft}^{2} / \mathrm{s}$ ); and
$S=$ storage coefficient.
The integral can be evaluated by summing the infinite series

$$
\begin{align*}
\int_{u}^{\infty} \frac{e^{-z}}{z} d z= & -0.577216-\ln u+u \\
& -\frac{u^{2}}{2 \cdot 2!}+\frac{u^{3}}{3 \cdot 3!}-\frac{u^{4}}{4 \cdot 4!}+\ldots \tag{2}
\end{align*}
$$

where $n!=n \cdot(n-1) \cdot(n-2) . . .3 \cdot 2 \cdot 1$.

The Theis equation is nonlinear in the parameters $T$ and $S$. Using the information in section 3.3.1, develop an algorithm for solving this equation for $T$ and $S$, given time and drawdown data.
a. Let
$f\left(t, r_{0} ; T, S\right)=\frac{Q}{4 \pi T} W(u)$
where $u=r_{0}^{2} S / 4 T t$ and $\mathrm{W}(u)$ is the integral in equation 2 . Find the sensitivities for $T$ and $S$. (See equation 3.3-6. Hint:

$$
\left.\frac{d}{d \alpha} \int_{u(\alpha)}^{\infty} f(x) d x=-f(u(\alpha)) \frac{d u(\alpha)}{d \alpha} .\right)
$$

b. Assume that initial estimates of transmissivity and storage coefficient, $T_{0}$ and $S_{0}$, exist. In equation 3.3-6 let $j=T$ indicate the sensitivity for $T$, and $j=S$ indicate the sensitivity for $S$. Then note that $X_{i T}^{0}$ can be scaled to become $Z_{i T}^{0}=X_{i T}^{0} T_{0}^{0}$, and $X_{i S}^{0}$ can be scaled to become $Z_{i S}^{0}=X_{i S}^{0} S_{0}$. Modify the functions for $X_{i T}^{0}$ and $X_{i S}^{0}$ computed in step a to become the scaled functions $Z_{i T}^{0}$ and $Z_{i S}^{0}$. Do you see any resulting simplifications in arithmetic? Do you think that the scaled sensitivities $Z_{i T}^{0}$ and $Z_{i S}^{0}$ might be more nearly uniform in value for any fixed $i$ than $X_{i T}^{0}$ and $X_{i S}^{0}$ ? What do you think this uniformity in value accomplishes?
c. Construct, for $n$ time observations at a single spatial location, the incremental linear model (equation 3.3-7). Then transform this model so that scaled sensitivities $Z_{i T}$ and $Z_{i S}$ are used instead of $X_{i T}$ and $X_{i S}$. How does the parameter displacement vector $\underline{d}$ transform? How can you recover $\underline{d}$ from the transformed displacement vector? (Hint: equations 3.3-10 through 3.3-14 are analogous to the present scaling problem.)
Figure 1 is a flow diagram for programming the steps indicated in section 3.3.1. Obtain a coding sheet and proceed to write code according to this flow diagram. Helpful hints, numbered on the flow diagram, follow:


Figure 1

1. $n=$ number of observations;
$r_{0}=$ radial distance to observation well;
$t_{i}=$ time of each observation;
$s_{i}=$ observed drawdown for each time $t_{i}$;
$T_{0}=$ initial guess, transmissivity;
$S_{0}=$ initial guess, storage coefficient;
$Q=$ discharge;
$\rho=$ damping parameter;
$\epsilon=$ convergence criterion;
$r_{m x}=$ maximum number of iterations.
2. $f_{i}\left(r_{0}, t_{i} ; T_{r}, S_{r}\right)=\frac{Q}{4 \pi T_{r}} W\left(r_{0}, t_{i} ; T_{r}, S_{r}\right)$.

You will have to program a finite number of terms of the infinite series to evaluate $W\left(r_{0}, t_{i} ; T_{r}, S_{r}\right)$. Note that one computation of $f_{i}$ will occur for every time observation $t_{i}$. These computations form the column vector $f$.
3. The sensitivities will form an $n \times 2$ array (one sensitivity for each parameter $T$ and $S$; sensitivities are evaluated for each observation). Use the information in parts $a$ and $b$ above.
4. Matrix multiplication is done with the basic algorithm

$$
c_{i j}=\sum_{k=1}^{n} a_{i k} b_{k j}
$$

where $a_{i k}, b_{k j}$, and $c_{i j}$ are elements of matrices $A, \underset{B}{B}$, and $\underset{C}{ }$, respectively, and $C=A B$. If $\underline{\underline{A}}=\underline{B}^{T}$, then

$$
c_{i j}=\sum_{k=1}^{n} b_{k i} b_{k j}
$$

because $b_{i k}^{T}=b_{k i}$. Similarly, if $A=\underline{D}^{T}$, then $\underline{C}=A B$ becomes

$$
c_{i j}=\sum_{k=1}^{n} d_{k i} b_{k j}
$$

Note that the model error variance $\operatorname{Var}(\underline{\epsilon})$ is assumed to be $I \sigma^{2}$.
5. Use the definition of the inverse to construct the inverse.
6. See part $c$ above. Also, note that $\underline{d}_{r+1}$ $=\left[\Delta T_{r+1}, \Delta S_{r+1}\right]$.
7. Include the damping parameter $\rho$ in your calculation of the new regression parameters. If convergence does not occur, then you may be best advised to set $\rho$ to a value less than one.
8. This is the end of the iteration loop.
9. You may wish to compute the final estimate of residuals $\underline{e}=f-\underline{s}$, then print them also.

Because convergence problems may arise, limit the number of passes that can be made through the algorithm. Test your code using the following data.

## 36-Hour Pumping Test

Test starts: 12 February 1976, 0805 h . Test ends: 13 February 1976, 2005 h. Production well was pumped varying between $517 \mathrm{gal} / \mathrm{min}$ and $530 \mathrm{gal} / \mathrm{min}$. Time and drawdown data at observation well 175 ft from production well appear in table 1.

Table 1
[Data from S.P. Larson, presently of S.S. Papadopulos and Associates (formerly U.S. Geological Survey), 1978]

| Time (s) | Drawdown $(\mathbf{f t})$ |
| ---: | :---: |
| 480 | 1.71 |
| 1,020 | 2.23 |
| 1,500 | 2.54 |
| 2,040 | 2.77 |
| 2,700 | 3.04 |
| 3,720 | 3.25 |
| 4,920 | 3.56 |

To aid in debugging your computer program, some example calculations follow. Assume that $T_{0}=0.1 \mathrm{ft}^{2} / \mathrm{s}, S_{0}=0.0005, Q=1.16 \mathrm{ft}^{3} / \mathrm{s}$, and $r_{o}=175 \mathrm{ft}$. Then
$u_{i}^{0}=\frac{r_{0}^{2} S_{0}}{4 T_{0} t_{i}}=\frac{(175)^{2}(0.0005)}{4(0.1) t_{i}}=38.28125 / t_{i}$
$f_{i}^{0}=\frac{Q}{4 \pi T_{0}} W\left(u_{i}^{0}\right)=\frac{1.16}{4 \pi(0.1)} W\left(u_{i}^{0}=0.923099 W\left(u_{i}^{0}\right)\right.$
Use of these values for $u_{i}^{0}$ and $f_{i}^{0}$ allows for calculation of the following data.

| $t_{i}$ | $u_{i}^{0}$ | $W\left(u_{i}^{0}\right)$ | $f_{i}^{0}$ |
| ---: | ---: | ---: | ---: |
| 480 | 0.0797526 | 2.02980 | 1.87371 |
| 1,020 | .0375306 | 2.74256 | 2.53165 |
| 1,500 | .0255208 | 3.11640 | 2.87675 |
| 2,040 | .0187653 | 3.41721 | 3.15442 |
| 2,700 | .0141782 | 3.69296 | 3.40897 |
| 3,720 | .0102907 | 4.00957 | 3.70123 |
| 4,920 | .00778074 | 4.28665 | 3.95700 |

Data for $u_{i}^{0}$ and $f_{i}^{0}$ are used to compute $Z_{i T}^{0}$ and $Z_{i S}^{0}$.

$$
\begin{aligned}
& Z_{i T}^{0}=-f_{i}^{0}+\frac{Q}{4 \pi T_{0}} e^{-u_{i}^{0}=-f_{i}^{0}+0.923099 e^{-u_{i}^{0}}} \\
& Z_{i S}^{0}=-\frac{Q}{4 \pi T_{0}} e^{-u_{i}^{0}=-0.923099 e^{-u_{i}^{0}}}
\end{aligned}
$$

Note that $Z_{i T}^{0}=X_{i T}^{0} T_{0}$ and $Z_{i S}^{0}=X_{i S}^{0} S_{0}$ are calculated directly, without first calculating sensitivities $X_{i T}^{0}$ and $X_{i S}^{0}$, then formally making the multiplication. These scaled sensitivities result from the following scaled linearized model.

$$
\begin{aligned}
f_{i}^{1} & =f_{i}^{0}+X_{i T}^{0} \Delta T_{1}+X_{i S}^{0} \Delta S_{1} \\
& \left.\left.=f_{i}^{0}+X_{i T}^{0} T o \left\lvert\, \frac{\Delta T_{1}}{T_{0}}\right.\right)+X_{i S}^{0} S_{0} \left\lvert\, \frac{\Delta S_{1}}{S_{0}}\right.\right) \\
& =f_{i}^{0}+Z_{i T}^{0} \frac{\Delta T_{1}}{T_{0}}+Z_{i S}^{0} \frac{\Delta S_{1}}{S_{0}}
\end{aligned}
$$

The computed values for $Z_{i T}^{0}$ and $Z_{i S}^{0}$ are:

| $z_{t T}^{0}$ | $z_{i S}^{0}$ |
| :---: | ---: |
| -1.02137 | -0.852339 |
| -1.64255 | -.889097 |
| -1.97691 | -.899839 |
| -2.24848 | -.905938 |
| -2.49887 | -.910103 |
| -2.78758 | -.913648 |
| -3.04106 | -.915944 |

Use of the computed values for $Z_{i T}^{0}$ and $Z_{i S}^{0}$ yields the entries in the $\underline{\underline{Z}}^{T} \underline{\underline{Z}}$ matrix:

$$
\sum_{i}\left(Z_{i T}^{0}\right)^{2}=35.96800209 \sum_{i}\left(Z_{i S}^{0}\right)^{2}=5.649402682
$$

$\sum_{i} Z_{i T}^{0} Z_{i S}^{0}=13.75336059$

Elements of the $\underline{Z}^{T}\left(\underline{s}-f_{0}\right)$ vector are computed in the following manner.

| $z_{i T}^{0}$ | $s_{i}$ | $f_{i}^{0}$ | $z_{i T}^{0}\left(s_{i}-f_{i}^{0}\right)$ |
| :---: | :---: | :---: | :---: |
| -1.02137 | 1.71 | 1.87371 | 0.1672084827 |
| -1.64255 | 2.23 | 2.53165 | .4954752075 |
| -1.97691 | 2.54 | 2.87675 | .6657244425 |
| -2.24848 | 2.77 | 3.15442 | .8643606816 |
| -2.49887 | 3.04 | 3.40897 | .9220080639 |
| -2.78758 | 3.25 | 3.70123 | 1.257839723 |
| -3.04106 | 3.56 | 3.95700 | 1.20730082 |

$\sum_{i} Z_{i T}^{0}\left(s_{i}-f_{i}\right)=5.579917421$

| $\mathrm{z}_{s}^{0}$ | $s_{i}$ | $f_{i}^{0}$ | $\mathrm{z}_{s S_{i}}^{0}\left(s_{i}-f_{i}{ }^{0}\right)$ |
| :---: | :---: | :---: | ---: |
| -0.852339 | 1.71 | 1.87371 | 0.1395364177 |
| -.889097 | 2.23 | 2.53165 | .2681961101 |
| -.899839 | 2.54 | 2.87675 | .3030207833 |
| -.05938 | 2.77 | 3.15442 | .3482606860 |
| -.910103 | 3.04 | 3.40897 | .3358007039 |
| -.913648 | 3.25 | 3.70123 | .412653870 |
| -.915944 | 3.56 | 3.95700 | .3636297680 |

$$
\sum_{i} Z_{i S}^{0}\left(s_{i}-f_{i}^{0}\right)=2.170709856
$$

Finally, the two elements of the scaled displacement vector ( $\Delta T_{1} / T_{0}, \Delta S_{1} / S_{0}$ ) are computed and used to compute $T_{1}$ and $S_{1}$, the new estimates for $T$ and $S$.

$$
\begin{aligned}
\frac{\Delta S_{1}}{S_{0}}= & (2.170709856-(13.75336059) \\
& (5.579917421) /(55.96800209) \\
& /\left(55.649402682-(13.75336059)^{2}\right. \\
= & \frac{0.03707407158}{0.3904247981}=0.09495829097 \\
\frac{\Delta T_{1}}{T_{0}}= & (5.579917421 \\
- & (0.09495829097)(13.75336059) \\
& / 35.96800209 \\
= & 0.1188256660 \\
S_{1}= & (1+0.0949583)(0.0005)=0.000547479 \\
T_{1}= & (1+0.118826)(0.1)=0.111883
\end{aligned}
$$

Some similar calculations for the second iteration are:

$$
\begin{aligned}
& u_{i}^{1}=\frac{(175)^{2}(0.000547479)}{4(0.111883) t_{i}}=37.4645 / t_{i} \\
& f_{i}^{1}=\frac{1.16}{4 \pi(0.111883)} W\left(u_{i}^{1}\right)=0.825057 W\left(u_{i}^{1}\right)
\end{aligned}
$$

| $t_{i}$ | $u_{i}^{1}$ | $W\left(u_{i}{ }^{1}\right)$ | $f_{i}{ }^{1}$ |
| ---: | ---: | :---: | :---: |
| 480 | 0.0780510 | 2.04973 | 1.69114 |
| 1020 | .0367299 | 2.76334 | 2.27991 |
| 1500 | .0249763 | 3.13743 | 2.58856 |
| 2040 | .0183650 | 3.43838 | 2.83686 |
| 2700 | .0138757 | 3.71423 | 3.06445 |
| 3720 | .0100711 | 4.03091 | 3.32573 |
| 4920 | .00761474 | 4.30805 | 3.55439 |


| $Z_{i T}^{1}=-f_{i}^{1}+0.825057 e^{-u_{i}^{1}}$ |  |
| :--- | :---: |
| $Z_{i S}^{1}=-0.825057 e^{-u_{i}^{1}}$ |  |
| $z_{i T}^{1}$ | $z_{i S}^{1}$ |
| -0.928031 | -0.763109 |
| -1.48461 | -.795303 |
| -1.78385 | -.804705 |
| -2.02682 | -.810043 |
| -2.25076 | -.813688 |
| -2.50894 | -.816789 |
| -2.73559 | -.818798 |

### 3.3.2 Nonlinear Regression When the Model Is Numerical

The basic model equations assumed in all previous developments have been of the closed form or analytical type where the dependent variable $f$ is a known function of $\xi$ and $\underline{\beta}$. In many cases such models may either not exist or be too complicated for practical use. In these cases the basic equation relating the dependent variable to the independent variables and parameters may be a numerical solution that can be stated in the general form

$$
\begin{equation*}
\underline{D}(\underline{h}, \underline{\xi}, \underline{\beta}) \underline{h}=q(\underline{h}, \underline{\xi}, \underline{\beta}) . \tag{3.3-21}
\end{equation*}
$$

Equation 3.3-21 is a nonlinear matrix equation in which $\underline{h}$ is the solution (dependent variable) vector of order $m ; D$ is a nonsingular coefficient matrix of order $m$ that is a function of $h, \xi$, and $\beta$; and $q$ is a vector of order $m$ that is a function
of $\underline{h}, \underline{\xi}$, and $\underline{\beta}$. Order $m$ is not related to the number of observations $n$, but instead is simply the order required to give a good numerical approximation to the solution of the problem.

If equation $3.3-21$ is linear in $\underline{h}$ so that $\underline{D}$ and $q$ are not functions of $h$, then equation $\overline{3.3}-21$ may be solved directly for $\underline{h}$. In this case the Gauss-Newton method may be used to obtain the regression solution. (The nonlinear case is considered in section 6.1.2.) The procedure is as follows. First write equation 3.3-21 in the form

$$
\begin{equation*}
\underline{h}=\underline{D}^{-1} q \tag{3.3-22}
\end{equation*}
$$

which is explicit in the dependent variable $h$. Next, note that $\underline{h}$ in equation 3.3-22 and $f \overline{\text { in }}$ equation 3.3-2 (or, as an estimate, equation 3.3-3) are expressions of the same quantity, the only difference between them being that elements of $\underline{h}$ are values of the dependent variable computed at points defined by the numerical solution, and elements of $f$ are values of the dependent variable computed at observation points. If all $\boldsymbol{n}_{\boldsymbol{s}}$ observation points are contained in the set of points required for the numerical solution, which implies that $m \geq n_{s}$, then $f$ is obtained from $\underline{h}$ simply by eliminating those entries in $\underline{h}$ not corresponding to observation points. In other instances the points in $m$ may not correspond to those in $n_{s}$. For these instances an interpolation scheme would be used to obtain $f$ from $h$. In either case, the vector $f\left(\underline{\xi}, \underline{b}_{r}\right)$ is obtained by using $\underline{h}$ computed from equation 3.3-22 in which $\underline{b}_{r}$ was used to evaluate $D$ and $q$.

The final step in forming the Gauss-Newton solution is to derive the sensitivity matrix $X$. To accomplish this step, write equation $3.3-\overline{21}$ in terms of a general parameter set $\underline{b}$, then differentiate it with respect to $\underline{b}$ to yield

$$
\begin{equation*}
\underline{D} \frac{\partial \underline{h}}{\partial b_{j}}+\frac{\partial \underline{\underline{D}}}{\partial b_{j}} \underline{h}=\frac{\partial \underline{q}}{\partial b_{j}}, j=1,2, \ldots, p \tag{3.3-23}
\end{equation*}
$$

or

$$
\frac{\partial \underline{h}}{\partial b_{j}}=\underline{D}^{-1}\left(\frac{\partial q}{\partial b_{j}}-\frac{\partial \underline{\underline{D}}}{\partial b_{j}} \underline{h}\right), j=1,2, \ldots, p .(3.3-24)
$$

The quantity $\partial \underline{h} / \partial b_{j}$ forms a column of the sensitivity matrix for points in $m$. Sensitivity
matrix $X_{r}$ would be found by first computing $\partial h / \partial b_{j}$ using equation 3.3-24 written in terms of $\underline{h}_{r}$ and $b_{r}$, then following one of the two procedures described above for obtaining $f$ from $\underline{h}$. By incorporating the procedures to compute numerical estimates of $f\left(\underline{\xi}, \underline{b}_{r}\right)$ and $\underline{X}_{r}$, the algorithm given for the Gauss-Newton method may be followed exactly to find the regression solution of equation 3.3-22.

By studying the sequence of calculations in the solution algorithm carefully, it will be noted that to calculate $\underline{h}_{r}$ and $\underline{X}_{r}, \underline{D}_{r}(\underline{D}$ computed using $\underline{b}_{r}$ ) must already have been computed. This requires one of two possible calculation schemes. Either $D_{r}$ and $p$ matrices of the form $\left(\partial \underline{D} / \partial \underline{b}_{j}\right)_{r}$ must be formed at the same time and stored before $\underline{h}_{r}$ and $\underline{X}_{r}$ are calculated, or $\underline{D}_{r}$ must be formed before $\underline{h}_{r}$ is calculated and then each matrix $\left(\partial D / \partial b_{j}\right)_{r}$ formed as needed to calculate each column (that is, $\underline{X}_{j}^{r}$ ) of $\underset{X}{r}$. The first alternative could require a considerable amount of computer memory or the use of scratch files, whereas the second alternative could require repetitive calculation because many arithmetic operations could be the same for forming both $D_{r}$ and $\left(\partial \underline{D} / \partial b_{j}\right)_{r}$. Often, however, matrices $\left(\partial \underline{D} / \partial b_{j}\right)_{r}$ can be written in a condensed or decomposed form to conserve computer memory. In this way $D_{r}$ and the decomposed form of $\left(\partial \underline{D} / \partial b_{j}\right)_{r}$ can be computed together without using a significant amount of extra memory. Then each matrix ( $\left.\partial D / \partial b_{j}\right)_{r}$ may be assembled as needed without performing numerous repetitive calculations.

A significant amount of computer memory can also be wasted unless care is taken when forming $\underline{X}_{\text {. }}$. The general procedure is to form the column vector $\left(\partial q / \partial b_{j}\right)_{r}-\left(\partial \underline{D} / \partial b_{j}\right) r_{r}$, then use equation 3.3-24 to form the vector $\left(\partial h / \partial b_{j}\right)_{r}$, which replaces the first vector in central computer memory. From this, vector $\underline{X}_{j}^{r}$ is immediately formed and stored. The matrix composed of vectors ( $\left.\partial \underline{h} / \partial b_{j}\right)_{r}$ should not be stored in central computer memory because it is often large. If desired, it may be stored column by column on a scratch file for later retrieval and printing.

## Problem 3.3-2

Assume the finite difference representation of a flow problem shown in figure 1.


Figure 1
Finite difference equations for this problem can be written as follows.

1. $1 / 2 \Delta y T_{1}\left(h_{2}-h_{1}\right) / \Delta x+1 / 2 \Delta x T_{1}\left(h_{4}-h_{1}\right) / \Delta y$ $=-1 / 2 \Delta x^{1 / 2} \Delta y W_{1}$.
2. $1 / 2 \Delta y T_{2}\left(h_{3}-h_{2}\right) / \Delta x-1 / 2 \Delta y T_{1}\left(h_{2}-h_{1}\right) / \Delta x$ $+1 / 2 \Delta x T_{1}\left(h_{5}-h_{2}\right) / \Delta y+1 / 2 \Delta x T_{2}\left(h_{5}-h_{2}\right) / \Delta y$ $=-1 / 2 \Delta x^{1 / 2} \Delta y W_{1}-1 / 2 \Delta x^{1 / 2} \Delta y W_{2}$.
3. $-1 / 2 \Delta y T_{2}\left(h_{3}-h_{2}\right) / \Delta x+1 / 2 \Delta x T_{2}\left(h_{6}-h_{3}\right) / \Delta y$ $=-1 / 2 \Delta x 1 / 2 \Delta y W_{2}$.
4. $\Delta y T_{1}\left(h_{5}-h_{4}\right) / \Delta x+1 / 2 \Delta x T_{1}\left(h_{B 1}-h_{4}\right) / \Delta y$ $-1 / 2 \Delta x T_{1}\left(h_{4}-h_{1}\right) / \Delta y=-1 / 2 \Delta x \Delta y W_{1}$.
5. $\Delta y T_{2}\left(h_{6}-h_{5}\right) / \Delta x-\Delta y T_{1}\left(h_{5}-h_{4}\right) / \Delta x$ $+1 / 2 \Delta x T_{1}\left(h_{8}-h_{5}\right) / \Delta y+1 / 2 \Delta x T_{2}\left(h_{8}-h_{5}\right) / \Delta y$ $-1 / 2 \Delta x T_{1}\left(h_{5}-h_{2}\right) / \Delta y-1 / 2 \Delta x T_{2}\left(h_{5}-h_{2}\right) / \Delta y$ $=-1 / 2 \Delta x \Delta y W_{1}-1 / 2 \Delta x \Delta y W_{2}$.
6. $-\Delta y T_{2}\left(h_{6}-h_{5}\right) / \Delta x+1 / 2 \Delta x T_{2}\left(h_{9}-h_{6}\right) / \Delta y$ $-1 / 2 \Delta x T_{2}\left(h_{6}-h_{3}\right) / \Delta y=-1 / 2 \Delta x \Delta y W_{2^{-1 / 2}} \Delta y q_{B 1}$.
7. $h_{7}=h_{B 1}$.
8. $1 / 2 \Delta y T_{2}\left(h_{9}-h_{8}\right) / \Delta x-1 / 2 \Delta y T_{1}\left(h_{8}-h_{B 1}\right) / \Delta x$ $-1 / 2 \Delta x T_{1}\left(h_{8}-h_{5}\right) / \Delta y-1 / 2 \Delta x T_{2}\left(h_{8}-h_{5}\right) / \Delta y$ $=-1 / 2 \Delta x^{1 / 2} \Delta y W_{1}-1 / 2 \Delta x x^{1 / 2} \Delta y W_{2}$.
9. $-1 / 2 \Delta y T_{2}\left(h_{9}-h_{8}\right) / \Delta x-1 / 2 \Delta x T_{2}\left(h_{9}-h_{6}\right) / \Delta y$ $=-1 / 2 \Delta x x^{1 / 2} \Delta y W_{2}{ }^{-1 / 2} \Delta y q_{B 1}$.

Or, by assuming that $\Delta x=\Delta y=a$,

1. $T_{1}\left(h_{2}-h_{1}\right)+T_{1}\left(h_{4}-h_{1}\right)=-1 / 2 a^{2} W_{1}$.
2. $T_{2}\left(h_{3}-h_{2}\right)-T_{1}\left(h_{2}-h_{1}\right)+\left(T_{1}+T_{2}\right)\left(h_{5}-h_{2}\right)$ $=-1 / 2 a^{2}\left(W_{1}+W_{2}\right)$.
3. $-T_{2}\left(h_{3}-h_{2}\right)+T_{2}\left(h_{6}-h_{3}\right)=-1 / 2 a^{2} W_{2}$.
4. $2 T_{1}\left(h_{5}-h_{4}\right)+T_{1}\left(h_{B 1}-h_{4}\right)-T_{1}\left(h_{4}-h_{1}\right)$ $=-a^{2} W_{1}$.
5. $2 T_{2}\left(h_{6}-h_{5}\right)-2 T_{1}\left(h_{5}-h_{4}\right)+\left(T_{1}+T_{2}\right)\left(h_{8}-h_{5}\right)$ $-\left(T_{1}+T_{2}\right)\left(h_{5}-h_{2}\right)=-a^{2}\left(W_{1}+W_{2}\right)$.
6. $-2 T_{2}\left(h_{6}-h_{5}\right)+T_{2}\left(h_{9}-h_{6}\right)-T_{2}\left(h_{6}-h_{3}\right)$ $=-a^{2} W_{2}-a q_{B 1}$.
7. $h_{7}=h_{B 1}$.
8. $T_{2}\left(h_{9}-h_{8}\right)-T_{1}\left(h_{8}-h_{B 1}\right)-\left(T_{1}+T_{2}\right)\left(h_{8}-h_{5}\right)$ $=-1 / 2 a^{2}\left(W_{1}+W_{2}\right)$.
9. $-T_{2}\left(h_{9}-h_{8}\right)-T_{2}\left(h_{9}-h_{6}\right)=-1 / 2 a^{2} W_{2}-a q_{B 1}$.
a. Verify the finite difference equations so that you understand their physical basis. (Hint: read appendix, section 4.3.1.)
b. Write the equations in matrix form:

$$
\underline{D h}=q
$$

by explicitly writing out $D, \underline{h}$, and $q$.
c. Let vectors $\left(\partial q / \partial b_{j}\right)_{r}-\left(\partial \underline{D} / \partial b_{j}\right)_{r} \underline{h}_{r}=J_{j}^{r}$, and develop $J_{j}^{r}(j=1,2,3,4)$ for the parameters $\beta_{1}=T_{1}, \beta_{2}=T_{2}, \beta_{3}=W_{1}$, and $\beta_{4}=q_{B 1}$.
d. Using the modified Gauss-Newton method, develop the algorithm (not computer program) to solve for the parameters. Assume that all nodes except number 7 are observation points.

### 3.3.3 Convergence and Conditioning

A value of $\rho, 0<\rho \leq 1$, can be shown to exist for which the Gauss-Newton procedure, as modified using equation $3.3-19$, will converge to the global minimum value of $S(\underline{b})$ provided that:

1. An initial estimate of the parameters can be found such that they lie within a parameter region $R$ bounded by sets of parameters $\underline{b}^{*}$ defined by

$$
\begin{equation*}
S\left(\underline{b}_{0}\right) \leq \min _{\underline{b}^{*}} S\left(\underline{b}^{*}\right) \tag{3.3-25}
\end{equation*}
$$

and the global minimum point lies within this region.
2. For all $\underline{b}$ belonging to $R, X$ is a continuous and unique matrix function.
3. The matrix $S^{T} \omega S$ is nonsingular and is a continuous function of $b$.
Condition 2 is almost always met. Condition 1 requires that the system be well enough understood that intelligent initial estimates of parameters can be made. Difficulties frequently arise in connection with condition 3. Unless the problem is correctly specified, the least squares coefficient matrix (for example, $S^{T} \omega S$ ) can be singular. Moreover, problems often arise because of ill-conditioning (that is, nearsingularity) of the matrix. Although the addition of the Marquardt parameter, $\mu$, is intended to help these cases, convergence can be difficult to obtain. In the following paragraphs, the general question of convergence is considered first. This is followed by discussions of singularity and ill-conditioning.

In general, the rate of convergence has been found to be related to the number of parameters being estimated, as predicted by theory. That is, the greater the number of parameters, the slower the rate, all other things being equal. It is also related to the conditioning of the problem and to the nearness of the initial set of parameters to the optimum set, in that the rate of convergence is usually much faster near a minimum of $S(\underline{b})$. As a rule of thumb, one may often expect convergence within a number of iterations equal to either 5 or twice the number of parameters, whichever is greater. Fewer iterations are required for well-conditioned problems.

A problem that frequently retards the convergence rate, or even causes divergence, is overshoot. This happens when the parameter correction vector $\rho \delta_{r+1}$ has a favorable orientation but is much longer than an ideal value. The result is that the new set of parameters $b_{r+1}$ is almost as far as (or even further than) the old set $b_{r}$ from the optimum value. A twoparameter example is illustrated in figure 3.3-4.

Overshoot is detected as large oscillations with accompanying changes in sign of components of $\rho \delta_{r+1}$ from one iteration to the next. The remedy is to decrease the value of $\rho$ such that $0<\rho<1$. In figure 3.3-4, a good value would be $\rho=0.5$, which would give $\underline{b}_{r+1}$ at point 1 .

Care must be taken not to make $\rho$ too small so that undershoot becomes a problem. Undershoot occurs when $\rho \hat{\delta}_{r+1}$ is too small, and it manifests itself as small steps $\rho \delta_{+1}$, the components of which usually do not change sign. The remedy is to increase the value of $\rho$, in rare occasions such that $\rho>1$. As a practical guide, the best value of $\rho$ is one that causes some


Figure 3.3-4
oscillations in sign of a few components of $\rho \underline{\delta}_{r+1}$ from one iteration to the next.
Sometimes the solution may converge to a local minimum instead of the global minimum. Most commonly this can occur (1) if $\rho$ is too large and the overshoot causes the search vector $\rho \delta_{r+1}$ to escape from the region defined by equation 3.3-25, or (2) if the initial estimate $b_{0}$ is not in the region defined by equation 3.3-25. Detection of this problem is accomplished through adequate knowledge of the system so that the wrong solution can be recognized. If the computed parameter vector $\underline{\hat{b}}$ is physically illogical or the model analysis discussed in section 5 suggests that the model is not correct, then one might suspect that a local minimum has been reached. The remedy to the problem is to decrease $\rho$ and (or) choose another initial estimate $\underline{b}_{0}$. If $\underline{b}_{0}$ is changed, usually the distance between $\underline{b}_{0}$ and the vector $\hat{\hat{b}}$ computed for the local minimum should be increased. Thus, if some components of $\underline{\hat{b}}$ are far too large to be realistic, then the same components of $\underline{b}_{0}$ should be reduced in value. If several attempts at changing $\rho$ and (or) $\underline{b}_{0}$ do not produce a change in $\underline{b}$, then the cause of the poor results is probably not a local minimum.
As for the linear case, singularity of the leastsquares coefficient matrix occurs whenever columns of the sensitivity matrix are linearly dependent. Near-singularity, caused by nearlinear dependence, is a more frequent occurrence. As a result of ill-conditioning, step sizes $\delta_{r+1}$ can be highly erratic, appearing to head toward no well-defined point and can be dominated by overshoot. In addition, some problems may start fairly well conditioned for the initial parameters but may become progressively more poorly conditioned during the iterative solution process.
The same techniques for analysis of poorly conditioned problems as are used for linear problems may be used for nonlinear problems as well. Whenever a problem is poorly conditioned, the sensitivity matrix $X$ may be examined for a near-zero column, and $S^{T} \omega S$ may be examined for off-diagonal components near unity. Also, the orthogonal transformation may be used to indicate that ill-conditioning exists and to point out possible columns where linear dependency occurs.

Use of the Marquardt parameter, $\mu$, is intended to improve conditioning by adding a small quantity to the main diagonal of the leastsquares coefficient matrix. Although conditioning is always artificially improved by employing $\mu$, the parameters resulting from applying the least squares process to a very poorly conditioned problem may be considerably in error unless the actual causes of the poor conditioning are discovered and the conditioning improved without using $\mu$.

### 3.3.4 Computation of $\mu$ and $\rho$

For best efficiency, both $\mu$ and $\rho$ should be recomputed at each iteration, $r$. A number of schemes exist in the literature for making these computations, but virtually all schemes involve assuming several trial values of $\mu$ and $\rho$, then performing all of the calculations for iteration $r$ for each of the trial values. The best values to use are then computed so as to minimize or substantially decrease $S(b)$. The problem with these schemes is that they require so much time that one is often much better off settling for approximate values of $\mu$ and $\rho$ computed by using a much simpler scheme.
The scheme adopted here is derived from the considerations discussed in section 3.3.3. Parameter $\mu$ is used only when the problem is so poorly conditioned that the search direction $\delta$ must be altered. Overshoot and undershoot are controlled primarily through use of $\rho$.
Because the Marquardt parameter is used with a scaled problem formulation, computations must be made using scaled quantities (see equation 3.3-16). By direct computation, it can be verified that the scaled gradient ( $\partial S(\underline{b}) / \partial b_{j}$ ) - $C_{j j}(j=1,2, \ldots, \mathrm{p})$ of the sum of squares $\bar{S}(\underline{b})$ is given by

$$
\begin{equation*}
\underline{C}^{T} \frac{\partial S(\underline{b})}{\partial \underline{b}}=-\underline{g}=-S^{T} \underline{\omega}(\underline{Y}-f(\xi, \underline{b},)) \tag{3.3-26}
\end{equation*}
$$

Thus, by definition (Spiegel, 1959, p. 16), the angle between $g$, which points directly down the sum of squares surface, and displacement vector $\underline{\delta}$ is given for the $r$ th iteration by

$$
\begin{equation*}
\cos \theta=\frac{\delta_{r+1}^{T} \underline{g}_{r}}{\sqrt{\left(\delta_{r+1}^{T} \underline{\delta}_{r+1}\right)\left(g_{r r}^{\left.T_{g r}\right)}\right.}} \tag{3.3-27}
\end{equation*}
$$

For two parameters, $b_{1}$ and $b_{2}$, the relationships given by equation 3.3-27 are illustrated in figure 3.3-5, (Note the use of scaled parameters; see equation 3.3-17.)

If $\theta=90^{\circ}$, then, as discussed in section 3.3.1, no improvement in parameters is likely to result from application of the Gauss-Newton procedure. However, application of the Marquardt parameter, $\mu>0$, will result in $\theta<90^{\circ}$ (Marquardt, 1963) because vector $\underline{\delta}$ is shifted progressively toward $g$ as $\mu$ increases. Thus, a viable scheme for choosing $\mu$ is to define a maximum value of $\theta, \theta_{m x}<90^{\circ}$, and compute $\mu$ so that $\theta$ never exceeds $\theta_{m x}$. This can be accomplished rather simply. At the beginning of the regression set $\mu_{0}=0$. Then at each iteration $r$, check and recompute $\mu$ as necessary:
$\left.\begin{array}{l}\mu_{r}=\mu_{\ell} \\ \text { if } \underline{\delta}_{r+1}^{T} g_{r} \geq \cos \theta_{m x} \sqrt{\left(\underline{\delta}_{r+1}^{T} \delta_{r+1}\right)\left(g_{g_{r}}^{T}\right)} \text { or } \\ \mu_{\ell+1}=\frac{3}{2} \mu_{\ell}+0.001 \\ \text { if } \delta_{r+1}^{T} g_{r}<\cos \theta_{m x} \sqrt{\left(\delta_{r+1}^{T} \delta_{r+1}\right)\left(g_{r g r}^{T}\right)}\end{array}\right)$
At the beginning of iteration $r, \ell=1$ and $\mu_{f}=\mu_{r-1}$. Then equation $3.3-20$ is solved and equation 3.3-28 is applied. If the second part of equation 3.3-28 is employed, equation 3.3-20 is resolved using $\mu_{\ell+1}, \ell$ is incremented by one, and equation 3.3-28 is used again. This process is continued until the first part of equation $3.3-28$ is used, at which point the appropriate value of $\mu$ for iteration $r$ has been found. The formula for computing $\mu_{\ell+1}$ from $\mu_{\ell}$ is empirical but gives what experience has shown to be a good range in values of $\mu$. For each resolution of equation 3.3-20, $S_{r}^{T} \omega_{\underline{\omega}}^{S_{r}}$ and $g_{r}$ are not recomputed. Thus, the calculations are not extensive.

Computation of $\rho$ is designed to prevent disastrous overshoot and to keep $\rho \underline{\delta}$ within the region $R$ defined by equation $3.3-25$. A simple but usually effective scheme is to estimate the maximum fraction that any of the parameters could change and still remain within $R$ and then to prevent any parameter from changing any more than this amount over any iteration. Let $t_{m x}$ be this maximum fractional change. Then at iteration $r, \rho$ is calculated as follows:

$$
\begin{equation*}
t=\max _{i}\left|d_{i}^{r+1} / c\right| \tag{3.3-29}
\end{equation*}
$$



Figure 3.3-5

$$
\left.\begin{array}{c}
\rho=1 \text { if } t \leq t_{m x} \text { or }  \tag{3.3-30}\\
\rho=t_{m x} / t \text { if } t>t_{m x}
\end{array}\right\}
$$

where $c=b_{i}^{r}$ if $b_{i}^{r} \neq 0$ and $c=1$ if $b_{i}^{r}=0$.

### 3.4 Regression Including Prior Information

### 3.4.1 Model Structure

Recall that the standard nonlinear regression model including prior information on the parameters may be written in the form (equation 3.1-32)

$$
\begin{equation*}
\underline{Y}=f(\xi, \beta)+\underline{\epsilon} \tag{3.4-1}
\end{equation*}
$$

where

$$
\begin{gather*}
\underline{Y}=\left[\begin{array}{l}
\underline{Y}_{s} \\
\underline{Y}_{p}
\end{array}\right]  \tag{3.4-2}\\
\left(\underline{\underline{\xi}, \underline{\beta})}=\left[\begin{array}{l}
\underline{f}_{s}(\underline{\xi}, \underline{\beta}) \\
\epsilon_{p}(\underline{\xi}, \underline{\beta})
\end{array}\right]\right.  \tag{3.4-3}\\
\underline{\epsilon}=\left[\begin{array}{l}
\underline{\epsilon}_{s} \\
\epsilon_{p}
\end{array}\right] \tag{3.4-4}
\end{gather*}
$$

and subscripts $s$ and $p$ indicate sample and prior information, respectively. To correspond with the above partitions into sample and prior information, the sensitivity matrix should be written in the form

$$
\underline{X}=\left[\begin{array}{l}
\underline{X}_{s}  \tag{3.4-5}\\
\underline{X}_{p}
\end{array}\right]
$$

where $X$ is a function of $\underline{\beta}$ for a nonlinear model. The transformation of $\bar{X}$ to $S$ used to control round-off error is defined analogously. Finally, recall that the weight matrix is partitioned as

$$
\underline{\omega}=\left[\begin{array}{ll}
\underline{\omega}_{s} & \underline{0}  \tag{3.4-6}\\
\underline{0} & \underline{\omega}_{p}
\end{array}\right]
$$

where $\underline{\omega}_{s}$ and $\underline{\omega}_{p}$ correspond to sample and prior information, respectively.

Often partition $\underline{X}_{p}$ will be obtained in a different manner than $\underline{X}_{s}$. For example, the model for the sample information may be numerical of the type in equation 3.3-22, whereas the model for the prior information may be of the analytical linear or nonlinear form. Thus, $\underline{X}_{s}$ would be obtained as described in section 3.3.2, and $X_{p}$ would be obtained as described in section 3.3.1. Other types of differences are handled in a similar fashion. Obviously, if sample and prior models are of the same type, then $\underline{X}_{s}$ and $\underline{X}_{p}$ are obtained in the same manner.
Despite the possibilities for combinations of linear, nonlinear, analytical, and numerical models, remember that all models have the general form of the incremental linear model when expanded in the Taylor series. Because all models resolve to the incremental linear form, for simplicity subsequent discussions in this section are based on this model only.

### 3.4.2 Solution Procedures

Whenever $\underline{\omega}_{s}$ and $\underline{\omega}_{p}$ are both known, solution for both linear and nonlinear models is unaltered from that given in the previous sections. However, recall that, because of the block diagonal form of equation 3.4-6, $S(\underline{b})$ and, hence, the normal equations can be written in a special form.

By applying the standard minimization technique to $S(\underline{b})$ as given by equation 3.1-42, which can be written in the form

$$
S(\underline{b})=\left[\underline{e}_{s}^{T}, \underline{e}_{p}^{T}\right]\left[\begin{array}{ll}
\underline{\omega}_{s} & \underline{0}  \tag{3.4-7}\\
\underline{0} & \underline{\omega}
\end{array}\right]\left\{\begin{array}{l}
\underline{e}_{s} \\
e_{p}
\end{array}\right]
$$

the normal equations for the incremental linear model become

$$
\left.\begin{array}{rl} 
& {\left[\underline{X}_{s}^{T}, \underline{X}_{p}^{T}\right]}
\end{array}\right]\left[\begin{array}{ll}
\underline{\omega}_{s} & \underline{0} \\
\underline{0} & \underline{\omega}_{p}
\end{array}\right]\left[\begin{array}{l}
\underline{X}_{s}  \tag{3.4-8}\\
\underline{X}_{p}
\end{array}\right] \underline{\hat{\hat{b}}-\underline{b})}\left[\begin{array}{ll}
\left.\underline{X}_{s}^{T}, \underline{X}_{p}^{T}\right]
\end{array}\right]\left[\begin{array}{ll}
\underline{\omega_{s}} & \underline{0} \\
\underline{0} & \underline{\omega}_{p}
\end{array}\right]\left[\begin{array}{l}
\underline{Y}_{s}-t_{s}(\underline{\xi}, \underline{b}) \\
\underline{Y}_{p}-\epsilon_{p}(\underline{\xi}, \underline{b})
\end{array}\right] .
$$

or

$$
\begin{align*}
& \left(\underline{X}_{s}^{T} \underline{\omega}_{s} \underline{X}_{s}+\underline{X}_{p}^{T} \underline{\omega}_{p} \underline{X}_{p}\right)(\underline{\hat{b}}-\underline{b}) \\
& =\underline{X}_{s}^{T} \underline{\omega}_{s}\left(\underline{Y}_{s}-f_{s}(\underline{\xi}, \underline{b})\right)+\underline{X}_{p}^{T} \underline{\omega}_{p}\left(\underline{Y}_{p}-f_{p}(\underline{\xi}, \underline{b})\right) \tag{3.4-9}
\end{align*}
$$

where $f_{s}(\xi, \underline{b})$ and $f_{p}(\xi, b)$ are for sample and prior information, respectively. Equation 3.4-9 is of the same form, and thus applies, as the equation for each iteration of solution of a nonlinear regression problem.

Frequently, the weight matrix is constructed from variance-covariance matrices for $\epsilon_{s}$ and $\epsilon_{p}$ that are given in the form

$$
\begin{align*}
& \operatorname{Var}\left(\epsilon_{s}\right)=\underline{V}_{s} \sigma^{2}  \tag{3.4-10}\\
& \operatorname{Var}\left(\epsilon_{p}\right)=\underline{U} \tag{3.4-11}
\end{align*}
$$

where the usual form $V_{p} \sigma^{2}$ for equation 3.4-11 cannot be used because $\operatorname{Var}\left(\epsilon_{p}\right)$ is not known as a function of $\sigma^{2}$. Thus, with $\underline{\omega}$ defined as

$$
\begin{align*}
\underline{\omega} & =\left[\begin{array}{ll}
{\left[\operatorname{Var}\left(\underline{\epsilon}_{s}\right)\right]^{-1}} & \underline{0} \\
\underline{0} & {\left[\operatorname{Var}\left(\underline{\epsilon}_{p}\right)\right]^{-1}}
\end{array}\right] \sigma^{2} \\
& =\left[\begin{array}{l}
\underline{V}_{s}^{-1} \underline{0} \\
\underline{0} \\
\underline{U^{-1}} \sigma^{2}
\end{array}\right] \tag{3.4-12}
\end{align*}
$$

then equation 3.4-9 becomes
$\left(\underline{X}_{s}^{T} \underline{V}_{s}^{-1} \underline{X}_{s}+\underline{X}_{p}^{T} \underline{U}^{-1} \underline{X}_{p} \sigma^{2}\right)(\underline{\hat{b}}-\underline{b})$
$=\underline{X}_{s}^{T} \underline{V}_{s}^{-1}\left(\underline{Y}_{s}-f_{s}(\underline{\xi}, \underline{b})\right)$
$+\underline{X}_{p}^{T} \underline{U}^{-1} \sigma^{2}\left(\underline{Y}_{p}-\mathcal{L}_{p}(\underline{\xi}, \underline{b})\right)$.
Hence, $\sigma^{2}$ apparently would have to be known to form the regression solution, whereas $\sigma^{2}$ is considered to be an unknown.

Theil (1963) showed that, for a linear model, $\sigma^{2}$ may be estimated for use in equation 3.4-13 by its ordinary least squares estimate (that is, the estimate obtained when prior information is not used). Bias produced by this estimate was shown by Theil (1963) to be of the order of $n_{s}^{-1 / 2}$. The procedure to be followed is to first solve the ordinary least squares problem by omitting all prior information; then find the estimate of $\sigma^{2}$ (to be given further on); finally use this estimate of $\sigma^{2}$ in the normal equations to solve the complete problem, including prior information.

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## 4 Numerical Nonlinear Regression Solution of General Steady-State Ground-Water Flow Problems

### 4.1 Assumed Model and Solution Procedure

A model to solve fairly general steady-state ground-water flow problems by using the regression procedures presented in section 3 is given in this section. A complete description of the method is given first. Documentation and listing of the computer program are given in the appendix, section 4.3.4.

### 4.1.1 Problem Specification

The equation assumed to govern groundwater flow for the class of problems to be analyzed is derived from equation 1.1-1 by letting $\partial h / \partial t \rightarrow 0$, which results in

$$
\begin{align*}
& \frac{\partial}{\partial x}\left(T_{x x} \frac{\partial h}{\partial x}\right)+\frac{\partial}{\partial y}\left(T_{y y} \frac{\partial h}{\partial y}\right)+R(H-h)+W \\
& +\sum_{\ell=1}^{N} \delta\left(x-a_{\ell}\right) \delta\left(y-b_{\ell}\right) Q_{\ell}=0, \tag{4.1-1}
\end{align*}
$$

where the symbols are defined the same as for equation 1.1-1.

Functions $T_{\xi \xi}$ (that is, $T_{x x}$ and $\left.T_{y y}\right), R$, and $W$ are each formulated within the region being modeled as the product of a parameter and a given (or known) function. To provide for spatial variability of parameters, the region is subdivided into a number of discrete zones within each of which the parameters are assumed
constant. Hence, known spatial variability (often, but not necessarily, smooth or continuous) is superimposed upon the discontinuous spatial variability dictated by the parameter zonation. As an example, hydraulic conductivity $K_{\xi \xi}$ may often be considered to be constant within particular rock types, each of which may be considered to be a discrete zone. Thus $K_{\xi \xi}$ may be considered to be a parameter. Thickness $b$ may be known from measurements and may vary continuously. The function $T_{\xi \xi}$ is, of course, given as $K_{\xi \xi} b$. Finally, because the controls that dictate a particular zonation may vary from parameter to parameter, zones for one type of parameter (for example, the parameter contained in $T_{\xi \xi}$ ) do not necessarily correspond to zones for another type (for example, the parameter contained in $W$ ). An example of zonation is given in figure 4.1-1 where the given function is unity so that the parameters are $T_{\xi \xi}$ and $W$.
Internal boundary conditions applying at discontinuities in $T_{\xi \xi}$ are that the specific discharge normal to the boundary and the hydraulic head both remain unchanged as the boundary is crossed. External boundary conditions applying on the periphery of the region being modeled include specified specific discharge normal to the boundary, specified hydraulic head at the boundary, or a mixture of the two types along the boundary.

Specific discharge $q_{B}$ normal to the boundary is assumed to vary along the boundary in a manner similar to that of $T_{\xi \xi}$. It may have discontinuities and may vary smoothly between discontinuities. Discontinuities in $q_{B}$ might often be expected to correspond to discontinuities in $T_{\xi!}$.
Hydraulic head variation along a specified head boundary is a continuous function, $h_{B}$, although the boundary may be subdivided into segments within each of which head $h_{B}$ may vary linearly or curvilinearly with distance.


Figure 4.1-1

Subdivision into segments is often based on the causes of the known head conditions along various segments of the boundary.
Unknown quantities to be determined are $T_{\xi \xi}, R, W, q_{B}$, and $h_{B}$. Single values of $T_{x x}$ or $T_{y y}^{\prime}, R$, and $W$ (or multipliers such as $K_{x x}^{x x}$ or $K_{y y}$ ) are assumed to be parameters in each zone, although any of these parameters may be held constant (assumed known as exact prior information) in each zone. Separate zones (or segments) are specified for the values of $q_{B}$, and single multipliers for the fluxes in each flux zone are assumed to be parameters. Within each of the separate segments (zones) of a specified head boundary the heads are adjusted as a linear function of distance by the regression procedure so that the parameters are the values of head at each end of each segment. Even though the adjustment is linear, the actual shape of the head profile along the boundary may be curvilinear.

### 4.1.2 Matrix Form of Regression Model

For most field problems, equation 4.1-1 with its attendant boundary conditions cannot be solved analytically. Thus, the regression solution must be based on a numerical solution of equation 4.1-1, which is expressed as a matrix equation. The particular numerical solution method is given in the appendix, section 4.3.1.
The matrix equation comprising the numerical solution is given as

$$
\begin{equation*}
\underline{D h}=\underline{q} \tag{4.1-2}
\end{equation*}
$$

where
$D=$ the square coefficient matrix of order $m$, the number of nodes used to discretize the modeled region;
$\underline{h}=$ the hydraulic head vector of order $m$; and $\bar{q}=$ the known vector of order $m$.
Matrix $\underline{D}$ contains parameters for $T_{\xi \xi}$ and $R$, whereas vector $q$ can contain all parameter types. To express any specified head value, say $h_{j}=h_{B j}$,

$$
\begin{align*}
& D_{j j}=1 \quad D_{j i}=D_{i j}=0, i \neq j  \tag{4.1-3}\\
& q_{j}=h_{B j} \tag{4.1-4}
\end{align*}
$$

where $D_{j j}, D_{i j}$, and $D_{j i}$ are components of $D$, and $q_{j}$ is a component of $q$. To accomplish the condition that $D_{i j}=0$ in equation $i, i \neq j$, the term $D_{i j} h_{B j}$ is transferred to the right-hand side of equation $i$ so that $q_{i}$ contains the term $-D_{i j} h_{B j}$. Then $D_{i j}$ in $D$ is set to zero.

The known head value is computed from

$$
\begin{equation*}
h_{B j}=\frac{h_{B j}^{0}\left[L_{j} H_{t}+\left(1-L_{j}\right) H_{s}\right]}{L_{j} H_{t}^{0}+\left(1-L_{j}\right) H_{s}^{0}} \tag{4.1-5}
\end{equation*}
$$

where
$s=$ node at one end of the boundary segment within which node $j$ lies;
$t=$ node at the other end;
$H_{s}=$ head (parameter) at node $s$;
$\boldsymbol{H}_{t}=$ head (parameter) at node $\boldsymbol{t}$;
$L_{j}=$ ratio of distance along the boundary from node $s$ to node $j$, to distance along the boundary from node $s$ to node $t$; and superscript $0=$ an initial or reference value.

Indices $s$ and $t$ can be equal so that $j=s=t$ for the case where only one specified head is present. Also, $H_{s}$ and $H_{t}$ can represent the same parameter, so that the entire specified head boundary behaves as a unit.
Because equation 4.1-2 is a linear matrix equation, the modified Gauss-Newton procedure is employed to solve the regression problem. Equation 4.1-2 is the same as equation 3.3-21, except that in equation 4.1-2 the coefficient matrix $\underline{D}$ and right-side vector $q$ are not functions of dependent variable vector $\underline{h}$. Hence, sensitivities may be calculated using equation 3.3-24.

Prior information is assumed to be given (if available) on each parameter individually so that

$$
\begin{equation*}
\underline{X}_{p}=\left[\underline{I}_{n_{p}}, \underline{0}\right]_{\left(n_{p} \times p\right)} \tag{4.1-6}
\end{equation*}
$$

where $\underline{I}_{n_{p}}$ is the identity matrix of order $n_{p}$. Thus, in equation 4.1-6, direct prior information is assumed to be given on the first $n_{p}$ parameters. Placement of these parameters first in the vector $\underline{\beta}$ simplifies theoretical statement of equation 4.1-6 but is not necessary in practice.

The linearized regression model assumed, then, is of the form of equation 3.3-7 partitioned as suggested by equations $3.4-2$ through 3.4-5:

$$
\left.\begin{array}{l}
\underline{Y}_{s}=f_{s}\left(\xi, \underline{\xi}, \underline{b_{r}}\right)+\underline{X}_{s}^{r}\left(\underline{b}_{r+1}-\underline{b}_{r}\right)+\underline{e}_{s}  \tag{4.1-7}\\
\underline{Y}_{p}=f_{p}\left(\xi, \underline{b}_{r}\right)+\underline{X}_{p}^{r}\left(b_{r+1}-\underline{b}_{r}\right)+\underline{e}_{p}
\end{array}\right\}
$$

where $\underline{Y}_{s}$ is the vector of heads at observation points, $\underline{L}_{s}\left(\xi, b_{r}\right)$ is the vector of computed heads at observation points for iteration $r, \underline{Y}_{p}$ is the vector of prior estimates of the first $n_{p}$ parameters, $f_{p}(\xi, \underline{b})$ is the vector of the first $n_{p}$ elements of the computed parameter vector (on which there is prior information) for iteration $r$, and, from equation 4.1-6, $\underline{X}_{p}^{r}=\left[\underline{I}_{n}, 0\right]$.
No correlation or other coupling is assumed to exist among components of either $\underline{\epsilon}_{s}$ or $\epsilon_{p}$. Matrix $\underline{\omega}$ is assumed to be of the form

$$
\underline{\omega}=\left[\begin{array}{l}
\underline{V}_{s}^{-1}  \tag{4.1-8}\\
\underline{0} \\
\underline{0} \\
\underline{U^{-1}} \sigma^{2}
\end{array}\right]
$$

where

$$
\begin{align*}
& \underline{V}_{s}^{-1}=\left[\operatorname{Var}\left(\underline{\epsilon}_{s}\right)\right]^{-1} \sigma^{2}  \tag{4.1-9}\\
& \underline{U}^{-1} \sigma^{2}=\left[\operatorname{Var}\left(\epsilon_{p}\right)\right]^{-1} \sigma^{2} \tag{4.1-10}
\end{align*}
$$

and both $\underline{V}_{s}^{-1}$ and $\underline{U}^{-1}$ are diagonal.

### 4.1.3 Nonlinear Regression Solution

Nonlinear-regression solution for the model given in section 4.1.2 is accomplished by using the algorithm at the end of section 3.3.1. Nodal sensitivities are calculated as illustrated in the appendix, section 4.3.2. We assume that observation points may be located anywhere within the flow region, so that computed heads, $f_{s}\left(\xi, b_{r}\right)$ and sensitivities, $\underline{X}_{s}^{r}$, at observation points in general must be obtained by interpolation from surrounding nodal values. Standard bilinear interpolation using the four adjacent nodes surrounding an observation is used as the interpolation method (Wang and Anderson, 1982, p. 153-155).

The normal equations used are equation $3.3-20$, as modified to include prior information also (see equation 3.4-13):

$$
\begin{align*}
& {\left[\left(\underline{S}_{j}^{r}{ }^{T} \underline{V}_{s}^{-1} \underline{S}_{s}^{r}+\left(\underline{S}_{p}\right)^{T} \underline{U}^{-1} \underline{S}_{p}^{r} s^{2}\right.\right.} \\
& +\mu I] \underline{\delta}_{r+1}=\left(\underline{S}_{s}^{r}\right)^{T} \underline{V}_{s}^{-1}\left(\underline{Y}_{s}-f_{s}\left(\xi, \xi_{-}\right)\right) \\
& +\left(\underline{S}_{j}^{r}\right)^{T} \underline{U}_{s}^{-1} s^{2}\left(\underline{Y}_{p}-f_{p}\left(\xi, b_{r}\right)\right) \tag{4.1-11}
\end{align*}
$$

where subscripts $s$ and $p$ refer to sample and prior information, respectively;
$S_{p}^{r}=$ the matrix $\left[\underline{I}_{n}, \underline{0}\right]$ from equation 4.1-6, transformed using equation 3.3-16;
$f_{p}\left(\xi, b_{r}\right)=a$ vector composed of the $r$ th estimate of those parameters on which there is prior information; and
$s^{2}=$ the ordinary least-squares estimate of $\sigma^{2}$ (to be developed later).

### 4.2 Singularity and Conditioning

Singularity of the least-squares coefficient matrix can occur whenever (1) no measured flow rates (such as well or spring discharges) are in the model and (2) an attempt is made to compute all parameters. To understand how this occurs, consider first the case where there are no specified head parameters and no prior information, but all other parameters are to be computed. Also, assume for simplicity that all observation points correspond to node points. In this case it can be shown (appendix, section 4.3.3) that

$$
\begin{equation*}
\underline{J} \underline{b}=\underline{0} \tag{4.2-1}
\end{equation*}
$$

where subscripts $r$ were omitted to simplify nomenclature, and $\underline{J}=\left\{\underline{J}_{j}\right\}=\left\{\partial q / \partial b_{j}-\left(\partial \underline{D} / \partial b_{j}\right) \underline{h}\right\}$. By using equation 3.3-23,

$$
\begin{equation*}
\underline{J}_{j}=\frac{\partial \underline{q}}{\partial b_{j}}-\frac{\partial \underline{\underline{D}}}{\partial b_{j}} \underline{h}=\underline{D} \frac{\partial \underline{h}}{\partial b_{j}} \tag{4.2-2}
\end{equation*}
$$

so that

$$
\begin{equation*}
\underline{J b}=\sum_{j=1}^{p} J_{j} b_{j}=\underline{D} \sum_{j=1}^{p} \frac{\partial \underline{h}}{\partial b_{j}} b_{j} . \tag{4.2-3}
\end{equation*}
$$

Because $\underline{D}$ is nonsingular,

$$
\underline{D}^{-1} \underline{D}_{j=1}^{p} \frac{\partial h}{\partial b_{j}} b_{j}=\underline{0}
$$

or, by eliminating those nodes not corresponding to observation points from $\partial \underline{h} / \partial b_{j}$,

$$
\begin{equation*}
\sum_{j=1}^{p} \underline{X}_{s j} b_{j}=\underline{0} \tag{4.2-4}
\end{equation*}
$$

Equation 4.2-4 can also be written in the form

$$
\begin{equation*}
\underline{X}_{s} \underline{c}=\underline{0} \tag{4.2-5}
\end{equation*}
$$

where

$$
\begin{equation*}
\underline{c}=\underline{b}=\left[b_{1}, b_{2}, \ldots, b_{p}\right]^{T} . \tag{4.2-6}
\end{equation*}
$$

Recall that equation $4.2-5$ implies that the leastsquares coefficient matrix is singular.
If a known flow rate $Q_{i}$ is at node $i$, then equation $4.2-1$ becomes

$$
\begin{equation*}
J \underline{b}=\underline{Q} \tag{4.2-7}
\end{equation*}
$$

where

$$
\begin{equation*}
\underline{Q}=\left[0,0, \ldots, Q_{i}, \ldots, 0\right]_{(m)}^{T} \tag{4.2-8}
\end{equation*}
$$

so that, all other things being equal,

$$
\begin{equation*}
\underline{\underline{X}} \mathbf{s} \underline{\boldsymbol{c}} \neq \underline{0} . \tag{4.2-9}
\end{equation*}
$$

In the case where at least one parameter $j$ is fixed, then $J b$ has column $j$ of $J$ and element $j$ of $b$ deleted. Thus, equation $4.2-1$ no longer holds, so that equation $4.2-9$ will hold, if no other source of singularity exists.

Whenever there are specified head parameters, $\underline{J}$ and $\underline{X}_{s}$ both contain columns resulting from these parameters, and $J b \neq 0$. However, for those columns not involving the specified head parameters, $\sum J_{j} b_{j}=\underline{0}$, where $j$ denotes all parameters except specified head parameters. If
$c_{j}=0$ for those columns resulting from specified head parameters and $c_{j}=b_{j}$ for the remaining columns, then equation 4.2-5 holds for the entire sensitivity matrix, which indicates that the problem is again singular. Because addition of specified head parameters has no influence on this type of singularity, it is assumed for simplicity for the remainder of this section that there are no specified head parameters.
Singularity caused by attempting to find all parameters in the absence of known flow rates can be rectified by using prior information. For the case of prior information, equation 4.2-5 can be written

$$
\left[\begin{array}{l}
\underline{X}_{s}  \tag{4.2-10}\\
\underline{X}_{p}
\end{array}\right] \underline{\underline{c}=\underline{0} .}
$$

The only way for equation $4.2-10$ to hold is if $\underline{X}_{s} \underline{c}=\underline{0}$ and $\underline{X}_{p} \underline{c}=\underline{0}$. If the only cause of the singularity is given by equation $4.2-1$, then $\underline{c}=\underline{b}$ is the only linearly independent solution of equation 4.2-10, and, if $\underline{X}_{p}$ is derived from equation 4.1-6, $X_{p i i}=0(i=1, \overline{2}, \ldots, p)$ is the only way that equation 4.2-10 can hold. Hence,prior information on any parameter can theoretically condition the problem so that all parameters can be found.
The maximum number of parameters that can be found for any problem can also be obtained through nondimensionalization (or partial nondimensionalization) to find the smallest number of independent groups. In addition, nondimensionalization also illustrates the idea that solution is actually often best expressed in terms of ratios of the parameters. As an example, consider the case where a region is composed of two zones where $T_{1}, T_{2}, W_{1}$, and $W_{2}$ are parameters. Then the flow equations for each zone are

$$
\left.\begin{array}{l}
T_{1}\left(\frac{\partial^{2} h}{\partial x^{2}}+\frac{\partial^{2} h}{\partial y^{2}}\right)+W_{1}=0  \tag{4.2-11}\\
T_{2}\left(\frac{\partial^{2} h}{\partial x^{2}}+\frac{\partial^{2} h}{\partial y^{2}}\right)+W_{2}=0
\end{array}\right)
$$

and the boundary conditions between zones are

$$
\left.\begin{array}{l}
T_{1}\left(\frac{\partial h}{\partial n}\right)_{1}=T_{2}\left(\frac{\partial h}{\partial n}\right)_{2}  \tag{4.2-12}\\
(h)_{1}=(h)_{2}
\end{array}\right\}
$$

where the notation $(\cdot)_{1}$ indicates that the quantity in parentheses is evaluated just within the 1 side of the boundary and similarly for $(\cdot)_{2}$. If equations 4.2-11 and 4.2-12 are written in the alternative forms

$$
\begin{gather*}
\frac{T_{1}}{T_{2}}\left(\frac{\partial^{2} h}{\partial x^{2}}+\frac{\partial^{2} h}{\partial y^{2}}\right)+\frac{W_{1}}{T_{2}}=0  \tag{4.2-13}\\
\frac{\partial^{2} h}{\partial x^{2}}+\frac{\partial^{2} h}{\partial y^{2}}+\frac{W_{2}}{T_{2}}=0  \tag{4.2-14}\\
\frac{T_{1}}{T_{2}}\left(\frac{\partial h}{\partial n}\right)_{1}=\left(\frac{\partial h}{\partial n}\right)_{2} \\
(h)_{1}=(h)_{2},
\end{gather*}
$$

instead of four independent parameters, there are only three written as ratios: $W_{1} / T_{2}, W_{2} / T_{2}$, $T_{1} / T_{2}$. A known flow rate in zone $1, Q_{1}$, would add the term $Q_{1} / T_{2}$ to the first part of equation 4.1-13. In this case $W_{1} / T_{2}, W_{2} / T_{2}, T_{1} / T_{2}$, and $Q_{1} / T_{2}$ could all be considered parameters. Knowing $Q_{1}$ (either exactly or with uncertainty) would provide unique estimates of the four original parameters. Another approach would be to find $T_{1}, T_{2}, W_{1}, W_{2}$ as parameters, knowing that the problem is not singular because there are four independent ratios for the problem.

Another common way for singularity to occur is if a column of $\underline{X}_{s}$ is zero: $\underline{X}_{s j}=\underline{0}$. This results if measurements are taken at points where there is no sensitivity to the parameter, $b_{j}$, corresponding to the column. If this is the only source of singularity, then $c=\left[0,0, \ldots, b_{j}\right.$, $0, \ldots, 0]^{T}$ is the only linearly independent solution of equation 4.2-5. In this case if $\underline{X}_{p}$ is derived from equation 4.1-6 and there is prior information on $b_{j}$, then $\underline{X}_{\underline{p}} \neq \underline{\text {, }}$, so that the prior information solves the singularity problem. Addition
of prior information on any other parameter alone obviously will not help.

Two sources of singularity result if $\underline{X}_{s j}=\underline{0}$ and $J b=\underline{0}$. In this case one solution of equation $4.2-\overline{5}$ is, as before, given by equation 4.2-6. However, because $\underline{X}_{s j}=0, c_{j}$ can be any arbitrary value less than infinity and so can be set to zero. Hence, $\underline{c}=\left[b_{1}, b_{2}, \ldots, 0, b_{j+1}, \ldots, b_{p}\right]^{T}$, where the zero appears in row $j$ of $\underline{c}$, is another solution to equation 4.2-5. Addition of prior information on parameter $j$ alone does not solve the singularity problem because, even though $X_{p j j}$ does not equal zero, $c_{j}$ and $X_{p i i}, i \neq j$, do equal zero so that $X_{p} \underline{\underline{c}}=\underline{0}$. A third solution of equation 4.2-5 is $\underline{c}=\left[0,0, \ldots, b_{j}, 0, \ldots, 0\right]^{T}$. In this instance addition of prior information on any or all parameters except parameter $j$ yields $\underline{X}_{p} \underline{c}=\underline{0}$. Thus, if $\underline{X}_{s j}=\underline{0}$ and $J \underline{b}=\underline{0}$, then the problem is singular unless prior information is added on parameter $j$ and at least one other parameter.

If the columns of $\underline{X}_{s}$ are almost linearly dependent, then the problem is ill-conditioned. Thus, if either $Q_{i}$ in equation 4.2-8 is almost zero or $\underline{X}_{j} \cong \underline{0}$, then an ill-conditioned problem can result. However, ill-conditioning can occur in a number of ways. The techniques given in section 3.2.3 can be used to detect conditioning problems.

## Problem 4.2-1

Solve problem 3.2-1 with the regression computer program (appendix 4.3.4). Assume that the stream tube is one foot wide and that transmissivities are unity. Place a row of nodes along each side of the stream tube, but specify observed heads only along one row or down the center of the tube (number of observed heads should be the same as in problem 3.2-1). Allow two iterations. What would happen if you were to attempt to estimate both $W$ and $T$ ?

## Problem 4.2-2

Figure 1 gives the zone map for a steady-state ground-water flow system in a hypothetical region. The finite difference mesh and types of boundary conditions also are shown on the map. Use the regression program (appendix 4.3.4) to construct a regression flow model for the region.
dISTANCE, THOUSANDS OF FEET


Figure 1

Prior information exists on the following parameters:

| $T_{2}=420 \mathrm{ft}^{2} / \mathrm{d} \quad$ devi | deviation $=84$ |
| :---: | :---: |
| $W_{1}=0.0004 \mathrm{ft} / \mathrm{d}$ | 0.00012 |
| $W_{3}=0.00017 \mathrm{ft} / \mathrm{d}$ | 0.000051 |
| $R_{2}=0.08 \mathrm{~d}^{-1}$ | 0.008 |
| $Q_{1}=-97,000 \mathrm{ft}^{3} / \mathrm{d}$ | 1940 |
| $Q_{2}=-51,000 \mathrm{ft}^{3} / \mathrm{d}$ | 1020 |
| $h_{B}$ at $(15,16)=10.4 \mathrm{ft}$ | 1.04 |
| $h_{B}$ at $(15,7)=4.8 \mathrm{ft}$ ( ${ }^{\text {a }}$ | eter $\quad 0.48$ |
| $h_{B}$ at (15,6)=4.8 ft $\}$ one parametex | ter 0.48 |
| $h_{B}$ at (15,5) $=5.4 \mathrm{ft}$ | 0.54 |

Assume that $h_{B}$ varies linearly between the estimated values. For one reason or another, the estimates of $h_{B}$ at the four nodes are not observations. (They may have been interpolated from a contour map, for example.)
There is no prior information on the remaining parameters, but probable limits of variation for these parameters are

$$
\begin{gathered}
30 \leq T_{1} \leq 80 \\
10 \leq T_{3} \leq 40 \\
-0.0003 \leq W_{2} \leq-0.00005 \\
0.2 \leq q_{B 1} \leq 0.8 \\
0.15 \leq q_{B 2} \leq 0.4
\end{gathered}
$$

From these ranges, initial estimates of the parameters may be determined.
The observed head data in table 1 were collected. They are of uniform reliability.

Table 1.

| Node | Value $(\mathrm{ft})$ | Node | Value (ft) |
| :--- | ---: | :--- | ---: |
| $(8,2)$ | 60.70 | $(7,11)$ | 6.68 |
| $(14,2)$ | 75.64 | $(13,11)$ | -15.32 |
| $(12,3)$ | 60.27 | $(3,12)$ | 16.88 |
| $(10,4)$ | 29.67 | $(5,12)$ | 15.87 |
|  |  |  |  |
| $(7,5)$ | 4.22 | $(9,12)$ | 4.48 |
| $(11,5)$ | 4.37 | $(11,12)$ | -18.34 |
| $(13,5)$ | 6.07 | $(13,13)$ | -2.47 |
| $(15,5)$ | 5.81 | $(15,13)$ | 8.10 |
|  |  |  |  |
| $(10,7)$ | 4.57 | $(3,14)$ | 54.12 |
| $(8,8)$ | 5.21 | $(5,14)$ | 38.27 |
| $(12,8)$ | 7.01 | $(10,14)$ | 0.053 |
| $(15,8)$ |  | $(12,14)$ | -2.92 |
|  |  |  |  |
| $(7,9)$ | 4.95 | $(7,15)$ | 8.30 |
| $(4,10)$ | 4.04 | $(14,15)$ | 4.54 |
| $(9,10)$ | -89.36 | $(11,16)$ | 85.82 |
| $(11,10)$ |  |  | 2.26 |

The river stage is about 4.5 ft everywhere.
Assuming that $\sigma^{2}=1$, find all possible parameters for the model. First, however, determine how many parameters you can find!

Examine the sensitivity maps. Are there data in relatively high sensitivity areas for all parameters? Do you think that there are places where new data points would improve the results?

### 4.3 Appendices

### 4.3.1 Integrated Finite Difference Model

The numerical solution of equation $4.1-1$ is obtained by using integrated finite difference methods. A rectangular grid of nodes is assumed as indicated in figure 4.3-1. Each node point is enclosed by a subdomain, which is a rectangular region bounded by sides located half way between adjacent node points.
The coordinates of a typical node ( $i, j$ ) are given as $\left(x_{i} y_{j}\right)$. With the nomenclature shown in figure 4.3-1, equation 4.1-1 can be integrated over a subdomain enclosing node ( $i j$ ) to produce

$$
\begin{align*}
& \int_{\Delta y_{j}}\left(T_{x x} \frac{\partial h}{\partial x}\right)_{i+1 / 2} d y-\int_{\Delta y_{j}}\left(T_{x x} \frac{\partial h}{\partial x}\right)_{i-1 / 2} d y \\
& +\int_{\Delta x_{i}}\left(T_{y y} \frac{\partial h}{\partial y}\right)_{j+1 / 2} d x-\int_{\Delta x_{i}}\left(T_{y y} \frac{\partial h}{\partial y}\right)_{j-1 / 2} d x \\
& +\int_{\Delta x_{i}} \int_{\Delta y_{j}} R(H-h) d x d y \\
& +\int_{\Delta x_{i}} \int_{\Delta y_{j}} W d x d y+\sum_{p=1}^{N_{i j}} Q_{p}=0  \tag{4.3-1}\\
& \text { where } \\
& \left.\begin{array}{l}
\Delta x_{i}=1 / 2\left(\Delta x_{i+1 / 2}+\Delta x_{i-1 / 2}\right) \\
\Delta y_{j}=1 / 2\left(\Delta y_{j+1 / 2}+\Delta y_{j-1 / 2}\right)
\end{array}\right\} \tag{4.3-2}
\end{align*}
$$

and $N_{i j}$ is the number of pumping wells in subdomain $\Delta x_{i} \Delta y_{j}$.
If $T_{\xi \xi}, R$, and $W$ are assumed to be constant in each cell ( $A, B, C, D$ ) adjacent to node ( $i, j$ ), then a valid numerical approximation of equation $4.3-1$ is


Figure 4.3-1

$$
\begin{align*}
& T_{x x i+1 / 2, j} \Delta y_{j} \frac{h_{i+1, j}-h_{i, j}}{\Delta x_{i+1 / 2}} \\
& -T_{x x i-1 / 2 j} \Delta y_{j} \frac{h_{i, j}-h_{i-1, j}}{\Delta x_{i-1 / 2}} \\
& +T_{y y i j+1 / 2} \Delta x_{i} \frac{h_{i, j+1}-h_{i, j}}{\Delta y_{j+1 / 2}} \\
& -T_{y y i j-1 / 2} \Delta x_{i} \frac{h_{i, j}-h_{i, j-1}}{\Delta y_{j-1 / 2}} \\
& +R_{i, j} \Delta x_{i} \Delta y_{j}\left(H_{i, j}-h_{i, j}\right) \\
& +W_{i, j} \Delta x_{i} \Delta y_{j}+Q_{i, j}=0 \tag{4.3-3}
\end{align*}
$$

where

$$
\begin{gathered}
T_{x x i+1 / 2 j}=\frac{\Delta y_{j-1 / 2} T_{x x B}+\Delta y_{j+1 / 2} T_{x x C}}{2 \Delta y_{j}} \\
T_{x x i-1 / 2 j}=\frac{\Delta y_{j-1 / 2} T_{x x A}+\Delta y_{j+1 / 2} T_{x x D}}{2 \Delta y_{j}} \\
T_{y y i, j+1 / 2}=\frac{\Delta x_{i-1 / 2} T_{y y D}+\Delta x_{i+1 / 2} T_{y y C}}{2 \Delta x_{i}} \\
T_{y y i j-1 / 2}=\frac{\Delta x_{i-1 / 2} T_{y y A}+\Delta x_{i+1 / 2} T_{y y B}}{2 \Delta x_{i}}
\end{gathered}
$$

$$
\begin{equation*}
R_{i, j}=\frac{\Delta x_{i-1 / 2} \Delta y_{j-1 / 2} R_{A}+\Delta x_{i+1 / 2} \Delta y_{j-1 / 2} R_{B}+\Delta x_{i+1 / 2} \Delta y_{j+1 / 2} R_{C}+\Delta x_{i-1 / 2} \Delta y_{j+1 / 2} R_{D}}{4 \Delta x_{i} \Delta y_{j}} \tag{4.3-8}
\end{equation*}
$$

$$
\begin{equation*}
W_{i, j}=\frac{\Delta x_{i-1 / 2} \Delta y_{j-1 / 2} W_{A}+\Delta x_{i+1 / 2} \Delta y_{j-1 / 2} W_{B}+\Delta x_{i+1 / 2} \Delta y_{j+1 / 2} W_{C}+\Delta x_{i-1 / 2} \Delta y_{j+1 / 2} W_{D}}{4 \Delta x_{i} \Delta y_{j}} \tag{4.3-9}
\end{equation*}
$$



Figure 4.3-2


Figure 4.3-3
where, from equation 4.3-12 for node $k$ not on a specified head boundary,

$$
\begin{align*}
D_{k, k-N C} & =-T_{y y k, 2} \frac{\Delta x_{i}}{\Delta y_{j-1 / 2}},  \tag{4.3-14}\\
D_{k, k-1} & =-T_{x x k, 1} \frac{\Delta y_{j}}{\Delta x_{i-1 / 2}},  \tag{4.3-15}\\
D_{k, k} & =T_{x x k, 3} \frac{\Delta y_{j}}{\Delta x_{i+1 / 2}}+T_{y y k, 4} \frac{\Delta x_{i}}{\Delta y_{j+1 / 2}} \\
& +T_{x x k, 1} \frac{\Delta y_{j}}{\Delta x_{i-1 / 2}}+T_{y y k, 2} \frac{\Delta x_{i}}{\Delta y_{j-1 / 2}} \\
& +R_{k} \Delta x_{i} \Delta y_{j}, \tag{4.3-16}
\end{align*}
$$

$$
\begin{gather*}
D_{k, k+1}=-T_{x x k, 3} \frac{\Delta y_{j}}{\Delta x_{i+1 / 2}},  \tag{4.3-17}\\
D_{k, k+N C}=-T_{y y k, 4} \frac{\Delta x_{i}}{\Delta y_{j+1 / 2}},  \tag{4.3-18}\\
q_{k}=R_{k} \Delta x_{i} \Delta y_{j} H_{k}+W_{k} \Delta x_{i} \Delta y_{j}+Q_{k} . \tag{4.3-19}
\end{gather*}
$$

For node $k$ on a specified head boundary, $D_{k, k-N C}=D_{k, k-1}=D_{k, k+1}=D_{k, k+N C}=0, D_{k, k}=1$, and $q_{k}=h_{B k}$, the specified head. All remaining $D_{k, l}=0$ for equation $k$ in both cases. To preserve symmetry of $D$, equations $\ell, \ell \neq k$, are modified as indicated just after equation 4.1-4.

The flow across specified flow boundaries is incorporated by using the $Q_{k}$ term, so that the total flow crossing the specified flow boundary of the subdomain around node $k$ is added into $Q_{k}$. If $Q_{k}=0$ on a boundary node and the head at the node is not specified, then the boundary for the node is automatically a no-flow type. When computing the total flow to add into a specified flow node, remember that nodes are on boundaries so that subdomains for boundary nodes are only fractions of the full subdomains. For example, see figure 4.3-4.

### 4.3.2 Computation of Nodal Sensitivities for the Integrated Finite Difference Model

Partial derivatives $J_{\ell}$, defined by

$$
\begin{equation*}
J_{l}=\frac{\partial q}{\partial b_{\ell}}-\frac{\partial \underline{\underline{D}}}{\partial b_{\ell}} h_{\ell}=1,2, \ldots, p \tag{4.3-20}
\end{equation*}
$$



Figure 4.3-4
are employed in equation 3.3-24 to compute sensitivities for the nonlinear regression solution of the numerical model given in section 4.3.1. Using the definitions of the elements of $D$ and $q$ given in section 4.3.1, partial derivatives for the various parameters are computed as follows:

1. $b_{p}=T_{x x A}$ at node $k$

$$
\begin{gather*}
\frac{\partial D_{k, k-N C}}{\partial b_{\ell}}=0  \tag{4.3-21}\\
\frac{\partial D_{k, k-1}}{\partial b_{\ell}}=-\frac{\Delta y_{j}}{\Delta x_{i-1 / 2}} \frac{\partial}{\partial T_{x x A}} \\
\cdot\left(\frac{\Delta y_{j-1 / 2} T_{x x A}+\Delta y_{j+1 / 2} T_{x x D}}{2 \Delta y_{j}}\right)=-\frac{\Delta y_{j-1 / 2}}{2 \Delta x_{i-1 / 2}}(4.3-22) \\
\frac{\partial D_{k, k}}{\partial b_{\ell}}=\frac{\Delta y_{j-1 / 2}}{2 \Delta x_{i-1 / 2}}  \tag{4.3-23}\\
\frac{\partial D_{k, k+1}}{\partial b_{\ell}}=0  \tag{4.3-24}\\
\frac{\partial D_{k, k+N C}}{\partial b_{\ell}}=0 . \tag{4.3-25}
\end{gather*}
$$

Similar expressions result for $b_{l}=T_{x x B}, T_{x x C}$, and $T_{x x D}$. If more than one of the cells $A, \ldots, D$ lie in the same zone, then derivatives for the individual expressions are summed to form the final value. For example, if $b_{q}=T_{y y A}=T_{y y B}$ $=T_{y y C}=T_{y y D}$, then

$$
\begin{align*}
& \begin{aligned}
& \frac{\partial D_{k, k-N C}}{\partial b_{p}}=\frac{\partial D_{k, k-N C}}{\partial T_{y y A}}+\frac{\partial D_{k, k-N C}}{\partial T_{y y B}} \\
&+\frac{\partial D_{k, k-N C}}{\partial T_{y y C}}+\frac{\partial D_{k, k-N C}}{\partial T_{y y D}} \\
&=-\frac{\Delta x_{i}}{\Delta y_{j-1 / 2}}\left|\frac{\Delta x_{i-1 / 2}}{2 \Delta x_{i}}+\frac{\Delta x_{i+1 / 2}}{2 \Delta x_{i}}+0+0\right| \\
&=-\frac{\Delta x_{i}}{\Delta y_{j-1 / 2}} .
\end{aligned} .
\end{align*}
$$

Derivatives for any configuration of zone boundary are handled by combinations of the form of equations 4.3-21 through 4.3-26.
An example of application of equations 4.3-21 through 4.3-26 for an irregular zone boundary is given in figure 4.3-5. Let $T_{x x}=T_{y y}=T$ for simplicity. Then, for a two-zone problem, the two transmissivities are $T_{1}$ and $T_{2}$.

For node 7:

$$
\begin{aligned}
& \frac{\partial D_{7,2}}{\partial T_{1}}=-\frac{\Delta x_{2}}{\Delta y_{11 / 2}}, \frac{\partial D_{7,6}}{\partial T_{1}}=-\frac{\Delta y_{2}}{\Delta x_{11 / 2}}, \frac{\partial D_{7,7}}{\partial T_{1}} \\
& =\frac{\Delta y_{2}}{\Delta x_{2^{1 / 2}}}+\frac{\Delta x_{2}}{\Delta y_{2^{1 / 2}}}+\frac{\Delta y_{2}}{\Delta x_{11 / 2}}+\frac{\Delta x_{2}}{\Delta y_{11 / 2}}, \\
& \frac{\partial D_{7,8}}{\partial T_{1}}=-\frac{\Delta y_{2}}{\Delta x_{2^{1 / 2}}}, \frac{\partial D_{7,12}}{\partial T_{1}}=-\frac{\Delta x_{2}}{\Delta y_{21 / 2}} .
\end{aligned}
$$

For node 7 the derivatives of the $D_{i j}$ 's with respect to $T_{2}$ are all zero because $T_{2}$ does not appear in any of the $D_{i j}$ 's. For node 8:

$$
\begin{aligned}
& \frac{\partial D_{8,3}}{\partial T_{1}}=-\frac{\Delta x_{3}}{\Delta y_{11 / 2}}, \frac{\partial D_{8,7}}{\partial T_{1}}=-\frac{\Delta y_{2}}{\Delta x_{21 / 2}}, \frac{\partial D_{8,8}}{\partial T_{1}} \\
& =\frac{\Delta y_{11 / 2}}{2 \Delta x_{3^{1 / 2}}}+\frac{\Delta x_{21 / 2}}{2 \Delta y_{2^{1 / 2}}}+\frac{\Delta y_{2}}{\Delta x_{2^{1 / 2}}}+\frac{\Delta x_{3}}{\Delta y_{1^{1 / 2}}}, \\
& \frac{\partial D_{8,9}}{\partial T_{1}}=-\frac{\Delta y_{1^{1 / 2}}}{2 \Delta x_{31 / 2}}, \frac{\partial D_{8,13}}{\partial T_{1}}=-\frac{\Delta x_{21 / 2}}{2 \Delta y_{2^{1 / 2}}} \\
& \frac{\partial D_{8,3}}{\partial T_{2}}=\frac{\partial D_{8,7}}{\partial T_{2}}=0, \frac{\partial D_{8,8}}{\partial T_{2}}=\frac{\Delta y_{2^{1 / 2}}}{2 \Delta x_{31 / 2}}+\frac{\Delta x_{31 / 2}}{2 \Delta y_{2_{21 / 2}}} \\
& \frac{\partial D_{8,9}}{\partial T_{2}}=-\frac{\Delta y_{2^{1 / 2}}}{2 \Delta x_{3^{1 / 2}}}, \frac{\partial D_{8,13}}{\partial T_{2}}=-\frac{\Delta x_{31 / 2}}{2 \Delta y_{21 / 2}}
\end{aligned}
$$

Figure 4.3-5
2. $b_{\ell}=R_{A}$ at node $k$

$$
\begin{gather*}
\frac{\partial D_{k, k-N C}}{\partial b_{\ell}}=\frac{\partial D_{k, k-1}}{\partial b_{\ell}}=\frac{\partial D_{k, k+1}}{\partial b_{\ell}} \\
=\frac{\partial D_{k, k+N C}}{\partial b_{\lambda}}=0 \\
\frac{\partial D_{k, k}}{\partial b_{\ell}}=\Delta x_{i} \Delta y_{j}\left(\frac{\Delta x_{i-1 / 2} \Delta y_{j-1 / 2}}{4 \Delta x_{i} \Delta y_{j}}\right)=1 / 4 \Delta x_{i-1 / 2} \Delta y_{j-1 / 2} \tag{4.3-28}
\end{gather*}
$$

$$
\begin{equation*}
\frac{\partial q_{k}}{\partial b_{\ell}}=1 / 4 \Delta x_{i-1 / 2} \Delta y_{j-1 / 2} H_{k} \tag{4.3-29}
\end{equation*}
$$

Derivatives for $R_{B}, R_{C}$, and $R_{D}$ are similar.
3. $b_{P}=W_{A}$ at node $k$

$$
\begin{equation*}
\frac{\partial q_{k}}{\partial b_{\ell}}=1 / 4 \Delta x_{i-1 / 2} \Delta y_{j-1 / 2} \tag{4.3-30}
\end{equation*}
$$

D All derivatives of $D_{i j}$ are zero.
4. $b_{\ell}=q_{B 1}$, boundary flux 1 in $Q_{k}$, where, for example, $Q_{k}=q_{B 2} \cdot 1 / 2 \Delta y_{j-1 / 2}+q_{B 1} \cdot 1 / 2 \Delta y_{j+1 / 2}$ (see figure 4.3-6).

$$
\begin{equation*}
\frac{\partial q_{k}}{\partial b_{\ell}}=1 / 2 \Delta y_{j+1 / 2} \tag{4.3-31}
\end{equation*}
$$

A similar expression results for $b_{\ell}=q_{B 2}$, and if $q_{B 1}=q_{B 2}$, the derivative is the sum of two parts, of equation 4.3-31 and its analog for $q_{B 2}$.
5. $b_{\ell}=H_{s}$ or $H_{t}$

The specified head at any node $k$ along a specified head boundary is given as

$$
\begin{equation*}
h_{B k}=A_{k}\left[L_{k} H_{t}+\left(1-L_{k}\right) H_{s}\right] \tag{4.3-32}
\end{equation*}
$$

where

$$
\begin{equation*}
A_{k}=\frac{h_{B k}^{0}}{L_{k} H_{t}^{0}+\left(1-L_{k}\right) H_{s}^{0}} \tag{4.3-33}
\end{equation*}
$$



Figure 4.3-6
and the meanings of the symbols are defined after equation 4.1-5.

If node $m$ (here only $m$ indicates an arbitrary node number) is adjacent to a boundary segment bounded by nodes $s$ and $t$, and node $k$ lies in the segment so that it appears in equation $m$, then for $b_{\ell}=H_{s}$,
$\frac{\partial q_{m}}{\partial b_{\ell}}=-D_{m, k} \frac{\partial h_{B k}}{\partial H_{s}}=-D_{m, k} A_{k}\left(1-L_{k}\right)$.
Similarly, for $b_{\ell}=H_{t}$,
$\frac{\partial q_{m}}{\partial b_{\ell}}=-D_{m, k} \frac{\partial h_{B k}}{\partial H_{t}}=-D_{m, k} A_{k} L_{k}$.
If node $k$ lies on the boundary, then the equation for node $k$ in equation 4.3-13 becomes

$$
\begin{align*}
& h_{k}=h_{B k} \\
& =A_{k}\left[L_{k} H_{t}+\left(1-L_{k}\right) H_{s}\right] \tag{4.3-36}
\end{align*}
$$

and, for $b_{l}=H_{s}$,

$$
\begin{equation*}
\frac{\partial q_{k}}{\partial b_{\ell}}=A_{k}\left(1-L_{k}\right) \tag{4.3-37}
\end{equation*}
$$

and similarly for $b_{f}=H_{t}$.

### 4.3.3 Derivation of Equation 4.2-1

By careful examination of equations 4.3-12 through 4.3-19 it can be seen that, if there are
no specified head parameters, $\underline{q}-\underline{\underline{D} h}=\underline{0}$ can be written in the form

$$
\begin{align*}
& a_{i 1} b_{1}+a_{i 2} b_{2}+\ldots+a_{i p} b_{p}-Q_{i}=0 \\
& i, 1,2, \ldots, m \tag{4.3-38}
\end{align*}
$$

where
$a_{i j}=$ coefficient containing $\Delta x, \Delta y$, and head differences;
$b_{j}=$ any parameter except a specified head parameter; and
$Q_{i}=$ term not containing parameters in $\underline{b}$.
Define

$$
\begin{equation*}
\underline{J}_{j}=\frac{\partial q}{\partial b_{j}}-\frac{\partial \underline{\underline{D}}}{\partial b_{j}} \underline{h} . \tag{4.3-39}
\end{equation*}
$$

Then, by carrying out the differentiations indicated in equation 4.3-39 and comparing the result with equation 4.3-38 it can be seen that

$$
\begin{equation*}
J_{i j} b_{j}=a_{i j} b_{j} \tag{4.3-40}
\end{equation*}
$$

so that

$$
\begin{equation*}
\sum_{j=1}^{p} J_{i j} b_{j}=\sum_{j=1}^{p} a_{i j} b_{j}=Q_{i} \tag{4.3-41}
\end{equation*}
$$

If $\underline{b}$ contains all possible parameters (except specified head parameters) and there are no known fluxes, then $Q_{i}=0$ and

$$
\begin{equation*}
\sum_{j=1}^{p} J_{i j} b_{j}=0 \tag{4.3-42}
\end{equation*}
$$

or

$$
\begin{equation*}
\underline{J} \underline{b}=\underline{0} . \tag{4.3-43}
\end{equation*}
$$

### 4.3.4 Documentation of Program for Nonlinear Regression Solution of Steady-State Ground-Water Flow Problems

Introduction.-This program is designed to obtain a nonlinear regression solution to the
finite-difference model of steady-state groundwater flow given in section 4.3.1. Basic calculation methods are given in sections 4.1 and 4.3.2.
The program was developed using the Microsoft ${ }^{1}$ Fortran Compiler, Version 3.3, with the DOS ${ }^{1} 2.0$ operating system on an IBM $^{1}$ $\mathrm{PC} / \mathrm{XT}$ computer with the IBM $^{1} 8088$ Math Coprocessor and 256 KB memory. Except for the OPEN statements near the beginning of the code, Fortran 66 was used throughout to make the code as machine independent as possible. The source code is contained in files INVFD.FOR and INVSUB.FOR in the $5 \frac{1}{4} \mathrm{in}$. diskette accompanying this report. These two files must be linked or compiled together.
The computer program is composed of a main program and eight subroutines. The main program controls input-output and performs all computations that cannot be accomplished more effectively with subroutines. The eight subroutines (D4SOLV, COEF, LSTSQ, PRTOT, ORDER, ARRAY, ARRAYI, HOBS) perform the following specialized tasks:
D4SOLV Obtains an LDU factorization solution of the set of linear algebraic equations resulting from application of the finite difference methods, assuming the equations are ordered in an alternating diagonal fashion (Price and Coats, 1974).
COEF Computes coefficients necessary for the determination of sensitivities and heads.
LSTSQ Computes the coefficients of the normal equations and solves the system of equations to determine the vectors of parameter changes and parameters.
PRTOT Prints matrices or vectors in a column configuration.
ORDER Computes equation numbers at grid points corresponding to the alternating diagonal ordering scheme.
ARRAY Reads and (or) prints 1- and 2 -dimensional real array variables.
ARRAYI Reads and (or) prints 1- and 2dimensional integer array variables.

[^3]HOBS Reads and prints observed heads and weighting values; computes coefficients for bilinear interpolation of computed heads and sensitivities at observation points.
The basic flow of the program can be described as follows:
A. Data are input and variables are initialized.
B. Using coefficients generated in COEF, an initial solution corresponding to the initial parameter estimates is computed by D4SOLV.
C. In an iterative fashion, the following four steps are taken until the regression technique converges or until the number of iterations exceeds the maximum allowed.
(1) Sensitivities are calculated using coefficients computed in COEF and in the main program.
(2) LSTSQ is employed to form and solve the normal equations.
(3) Parameters are updated in LSTSQ using the parameter change vector generated.
(4) Various coefficients involving the updated parameters are computed in COEF, and current estimates of head are computed using D4SOLV.
D. Various statistics associated with the regression analysis are computed.
Aquifer Property Zonation and Variable Definition.-Basic model geometry is defined by the finite difference grid that is constructed over the region to be modeled. Nodes, consisting of grid intersections, are numbered from the lower left-hand corner of the grid (columns from left to right and rows from bottom to top). Cells, consisting of intragrid areas bounded by four adjacent nodes, are numbered similarly (figure 4.3-7).

The finite difference grid is divided into aquifer property zones, which define zonal values for the aquifer properties $T_{x x}, T_{y y}, R$, and $W$. Each zonal value is constant within the zone. Variation of a property within a zone is accomplished by assigning cell values. The aquifer property at any particular cell is computed as the product of the zonal value and the cell value of the property. Thus, if all cell values for a property within a zone are given a value of unity, the zonal value becomes the value of the property for each cell within that zone.


Figure 4.3-7

Zones are created by subdividing the grid into groups of cells having distinct combinations of zonal values. Cells belonging to these zones are accordingly assigned distinct zone numbers (IZN), reserving IZN $=0$ to indicate groups of cells outside of the model area.

All zonal values are either regarded as regression parameters to be determined by the procedure or are held constant and, thus, are not regarded as regression parameters to be determined, as specified in the input. From a conceptual viewpoint there is no difference between these two designations because zonal values that are held constant can be regarded as regression parameters having exact prior information. However, from a computational viewpoint it is most efficient to eliminate parameters being held constant from the calculations. Hence, these types of parameters do not appear in the normal equations or any of the vectors and matrices derived from them. To simplify nomenclature, in subsequent discussions the term regression parameter refers only to those regression parameters to be determined.

Cell values for an aquifer property are input using rectangular blocks of cells, which are defined for each property for convenience of input only. These blocks need not bear any relationship to the zones. Cell values may be constant or variable within each block. Cell values of a property are unaffected by the regression procedure.

Boundary Conditions and Boundary Parameters.-Two types of boundary conditions may be used: specified flow and (or) specified head. Nonzero specified-flow boundaries where the flow rate is known can be imposed
by assigning the appropriate value of the specified flow rate to the nodal value of variable WELL. Because the no-flow boundary is the default condition, zeros do not have to be assigned to WELL to simulate this condition. Specified-head boundaries where the head is known can be imposed by assigning the value of the known head at the boundary nodes to variable HC, which also describes nodal values of computed head. A negative one must be assigned to variable IN corresponding to each node that is to be considered a specified head node. Segments of the boundary that will be considered as either flow or head regression parameters must not be entered into WELL or HC arrays.
Different variables are used if specified-flow or specified-head boundary conditions are to be considered as regression parameters. Variable QB is used to indicate the zonal value for discharge across groups of nodes that form a specified-flow segment (or zone). The specified volumetric flow per unit width for each cell boundary within a zone is the product of QB for that zone and a multiplier for the cell boundary. By restricting a flow-boundary zone to a single node, point recharge and discharge can be simulated. QB can be a regression parameter and, therefore, can be modified by the regression procedure; the multiplier is unaffected. If the variable IP that specifies the regression parameter number for the boundary zone is set to zero, then QB is held constant, and the segment is treated as a known-flow boundary, thus giving two possible ways (via WELL and via QB) to designate known-flow boundaries.
Specified-head boundary nodes that are to be considered regression parameters are defined by segments composed of a sequence of nodes (variables ILOC and JLOC) along portions of the boundary. The specified heads at the first and last nodes in the sequence can be either different regression parameters or a single regression parameter or held constant, depending on the nature of the problem. Adjustments to heads at these nodes computed by the regression procedure are apportioned to other nodes in the sequence. The proportion is the ratio of the distance (along the sequence of nodes) between the end node and the node of interest, to the distance between the two end nodes. These
factors (PLA and PLB) are computed by the program for a given segment. As in the case with flow-boundary parameters, if the regression parameter number is set to zero, then that parameter is held constant. If the parameter numbers at both ends of the segment are set to zero, then the segment is treated as a knownhead boundary, thus giving two possible ways to designate known-head boundaries. A negative one must be entered into IN for all nodes on the specified-head boundary, whether or not the boundary involves regression parameters.
The definition of some of the more important variables related to aquifer properties and boundary conditions in the computer program are given below.

| Variable name | Definition |
| :---: | :---: |
| PAR(1), PAR(2). PAR(3), PAR(4) | Zonal value for transmissivity ( $T_{x x}$ and $T_{y y}$ ), hydraulic conductance $(R)$, and distributed recharge ( $W$ ), respectively. |
| IZN | An integer array that indicates the zone number of each cell. |
| CX, CY, VL, QR | Cell values for $x$-direction transmissivity, $y$-direction transmissivity, hydraulic conductance, and distributed recharge, respectively. |
| HR | Nodal values of head on the boundary of the confining bed opposite the aquifer. |
| WELL | Nodal (or point) values of known volumetric discharge (or recharge) from a well or other known-rate source-sink phenomena. |
| QBF | Cell-boundary multipliers for specified flow. |
| PLA, PLB | Arrays, which give the proportional distances from any point to either end of a segment where the specified heads at one (or both) end(s) of the segment is a (or are) regression parameter(s). |

Numbering of Regression Parameters.-The three types of regression parameters (aquifer parameters, specified-flow parameters, and specified-head parameters) are numbered consecutively and in any order. Also, any of the zonal values for aquifer properties or specifiedflows and segment-end values for specified head zones can be given a single regression parameter number. In this case the zonal or segmentend values sharing the same number must have
identical regression parameter values. Any differences in property values must be specified by differences in multipliers. A common example is to use only one regression parameter for both the $x$ - and $y$-direction zonal transmissivity values and to fix the degree of anisotropy for the zone by using the CX and CY arrays. As another example, if good estimates of the hydraulic gradient across cell boundaries forming a specified-flow boundary for an aquifer property zone were available, and if transmissivity were isotropic, then these gradient values could be used as multipliers so that the regression parameter QB for the specified-flow zone comprising the specified-flow boundary would be transmissivity. Thus, the $x$ - and $y$-direction transmissivities for the aquifer property zone and QB for the specified-flow zone could all be the same regression parameter. A high degree of flexibility for distributing aquifer properties and boundary conditions while, at the same time, keeping the number of regression parameters to a minimum is achieved with these types of schemes.
Prior Information on Regression Param-eters.-If estimates of the regression parameters and their (less than infinite) reliability are available from other sources (for example, aquifer tests), it is desirable to introduce this information into the regression analysis. For this case, initial values for the parameters are taken to be the prior information. The reliability of each estimate is represented by a standard deviation. Array WP is used to store these values for the regression parameters. However, only the standard deviations for the aquifer regression parameters are read directly into WP. The standard deviations for boundary regression parameters are read in through temporary variables (SDQB for specified-flow regression parameters, and SDHA and SDHB for specified-head regression parameters) and are only subsequently placed into the WP array.
The use of prior information of known reliability requires an estimate of the error variance of the heads (variable, EV) computed using ordinary least squares. If the estimate differs substantially from the value computed by the analysis using prior information, the problem should be resolved using the latter computed value as the estimate of error variance.

In some instances prior information of unknown reliability may be available. Use of this type of information is an advanced topic and is not covered in this report. The papers by Cooley (1982, 1983) cover the method in detail. Variables RP and BP are used to input the additional information needed for this method.
Solution-Only Mode.-To facilitate the calculation of certain statistical measures, the program is capable of bypassing the regression analysis and computing only head distributions for various combinations of parameter values. This is accomplished by specifying the solutiononly option (variable ISO) and providing the various combinations of parameter values for which solutions are desired.
Using the Program.-The computer code has been designed to be as machine independent as possible. Also, to minimize confusion, all arrays have been dimensioned explicitly. The following list summarizes the minimum dimensions required for the program to operate properly for a specific problem. If
$N_{g}$ is the number of node points ( $N_{x} \times N_{y}$ ), $N_{x}^{g}$ is the number of node columns, $N_{y}$ is the number of node rows, $N_{e}^{y}$ is the number of active (nonspecified head) nodes in the grid,
$N_{o}$ is the number of observed heads, $N_{\text {ot }}$ is the total number of observed heads plus the number of parameters on which there is prior information,
$N_{z}$ is the number of aquifer property zones in the model grid,
$N_{p}$ is the number of regression parameters associated with aquifer property zones, $N_{n}$ is the number of zonal aquifer properties that are not regression parameters,
$N_{h s}$ is the number of specified-head zones (or segments),
$N_{q p}$ is the total number of nodes on boundaries where flow is a regression parameter, $N_{h p}$ is the total number of nodes along boundaries where head is composed of one or more specified-head regression parameters, $N_{m h}$ is the maximum number of nodes in any specified-head zone, and
$N_{R}$ is the total number of regression parameters;
then the array variables should be dimensioned as in table 4.3-1.

Table 4.3-1

| Variable name | Dimension |
| :---: | :---: |
| WELL,HR,HC,XV,ILOC,JLOC,IN | $N_{\text {g }}$ |
| CX,CY,VL,QR,CXS,CYS,VLS,QRS,IZN | $\left(N_{x}-1\right)\left(N_{y}-1\right)$ |
| DX |  |
| DY,JPOS | $N_{y}$ |
| HCI,BK,BL,BM,BN,HO,W,KOBS | $N_{0}$ |
| PAR,XS, KN | 4 |
| IPRM | 4, $N_{z}$ |
| QBF,IBNA,IBNB | $N_{q p}$ |
| PLA,PLB,IBHN | $N_{h p}$ or $N_{m h}$, whichever is larger. |
| CXHR,CXHL,CYHT,CYHB | $N_{\text {hp }}$ |
| IHSN | $2 N_{h p}$ |
| IBPA, IBPB | $N_{h s}$ |
| IBZN | $N_{\text {qp }}+N_{h s}$ |
| P,WP,NCBA,NCEA,NCBF,NCEF,NCBH,NCEH | $N_{R}$ |
| X | $N_{R} N_{o}$ |
| S | $N_{R} N_{\text {ot }}$ |
| A | $N_{R} N_{R}$ |
| B | $N_{R}+N_{n}$ |
| V | $N_{e}$ or $3 N_{R}$, whichever is larger. |
| AU,IC | 5, $\mathrm{N}_{e} / 2$ |
| AL | ${ }^{1} N_{m}, N_{e} / 2$ |

[^4]Note that array variables that have a single dimension (CX, CY, QBF, HO, etc.) and are passed to subroutines are dimensioned as unity within the subroutines (only the initial address of an array is actually passed to a subroutine). This unit dimension should not be changed in subroutines when the dimension of the variable is changed in the main program. A similar system is used for multidimensional arrays, and their dimensions within subroutines should not be changed either. To accompany any change in program dimensions, variable NVD must be set equal to the dimensions of $A$ and the first
dimension of X , and variable NAD must be set equal to the first dimension of AL. These variables are defined near the beginning of the main program.

Input Data.-The input is arranged into data sets, each data set being composed of one or more lines of logically related input data, such as cell-by-cell multipliers for $x$-direction transmissivity or zone-by-zone initial parameter values. Each input line in a data set is a maximum of 80 columns, or characters, long. The formats for the data applying for each line are given with the discussions of the data sets.

Data Set A.
Three title lines of user's choice (format, 20A4).
Data Set B.
Problem size information; one line (format, 16I5).

| Line columns | Variable | Definition |
| :---: | :---: | :---: |
| 1-5 | ID | Number of node columns. |
| 6-10 | JD | Number of node rows. |
| 11-15 | NZNS | Number of aquifer property zones. |
| 16-20 | NOBS | Number of observations of head. |
| 21-25 | NPAR | Number of regression parameters associated with aquifer property zones. |
| 26-30 | NVAR | Total number of regression parameters. |
| 31-35 | NWELS | Number of known point flows. |
| 36-40 | NQBZ | Number of specified-flow boundary zones. |
| 41-45 | NHBZ | Number of specified-head boundary zones. |
| 46-50 | NUM | Maximum number of iterations allowed for the regression analysis. |
| 51-55 | IPRX | Additional print sensitivities and orthogonalize-sensitivities option. Code 1 to select the option. |
| 56-60 | IPO | Additional printout option. Code 1 to select the option. |
| 61-65 | ISO | Head-solutions only option. Code 1 to select the option. |

Data Set C.
Special input parameters; one line (format, 8F10.0).

| Line columns | Variable | Definition |
| :---: | :---: | :---: |
| 1-10 | DMX | Maximum fractional change, $t_{m x}$, allowed any regression parameter over any iteration. |
| 11-20 | CSA | Cosine, $\cos \Theta_{m x}$, of the maximum angle allowed between the gradient direction and the search direction (normally set to 0.08). |
| 21-30 | RP | Ridge parameter for regression analysis using prior information of unknown reliability. Code 0.0 if not used. |
| 31-40 | BP | Bias parameter for regression using prior information of unknown reliability. Code 0.0 if not used. |
| 41-50 | EV | Estimated error variance for problems using prior information of known reliability. Code 0.0 if not used. |

Data Sets D through K.
A number of variables are input into the code by first subdividing the grid into rectangular regions (blocks) and then reading the variables for each block. Blocking can be applied to either cells or nodes, depending upon the variable being input. Blocking allows considerable flexibility in the input of certain variables and, once understood, can speed the construction of a model. Block and zone boundaries do not necessarily have to coincide; blocking is basically a convenient way of assigning variable values to every node or cell in the grid.

Data sets $\mathbf{D}$ through K represent the real (floating-point) variables subject to blocking. These variables, in the order they must appear, are listed next.

| Data Set | Variable | Type of Variable |  |
| :---: | :--- | :--- | :--- |
| D | DX | Cell array $\ldots .$. | Distance between node points in $x$ or I direction. |
| E | DY | Cell array $\ldots$. | Distance between node points in $y$ or J direction. |
| F | CX | Cell array $\ldots$. | Multiplier (cell value) for $x$-direction transmissivity. |
| G | CY | Cell array $\ldots$. | Multiplier (cell value) for $y$-direction transmissivity. |
| H | VL | Cell array $\ldots$. | Multiplier (cell value) for hydraulic conductance of confining |
|  |  |  | bed. |
| I | HR | Nodal array $\ldots$ | Head on boundary of confining bed opposite the aquifer. |
| J | QR | Cell array .... | Multiplier (cell value) for recharge rate per unit area. |
| K | HC | Nodal array $\ldots$ | Initial head at active node or fixed head at specified-head |
|  |  |  |  |

Each data set $D$ through $K$ consists of an initial line defining the number of blocks (NOBL) into which the grid has been subdivided, and then a subsequent line or set of lines that define the blocks and the value or values of the variable to be input. The initial line, read with a A4, IX, 215 format, has the following form:

| Line columns | Variable |  |  |
| :--- | :--- | :--- | :--- |
| $1-4$ | NME | $\ldots \ldots$. | Definition |
| $6-10$ | NOBL | $\ldots \ldots$. | Array name for the variable. <br> Number of rectangular input blocks into which the variable <br> has been subdivided. |
| 11-15 | IPRN $\ldots \ldots .$. | Print option for full array. Set to 0 for print. Set to 1 for no <br> output. |  |

The initial line directs the program to seek NOBL blocks of information for a particular variable. If the variable is uniform over the block, then a single line suffices to define the block location and the uniform value to be assigned to every node or cell. If the variable is nonuniform over the block, then by specifying a value of IVAR equal unity on this line the program can be directed to seek additional lines specifying values of the variable for each node or cell in the block. This information is input through the format 4I5, F10.0,I5 as follows:

| Line columns | Variable | Definition |
| :---: | :---: | :---: |
| 1-5 | IB | Beginning column of the rectangular input block. |
| 6-10 | IE | Final column of the rectangular input block. |
| 11-15 | JB | Beginning row of the rectangular input block. |
| 16-20 | JE | Final row of the rectangular input block. |
| 21-30 | FACT | If the array set is uniform for the entire block, FACT is the cell or nodal value that is assigned to each element. If the array set is not uniform, each cell or nodal value on the subsequent data lines will be multiplied by FACT. |
| 31-35 | IVAR | Code 0 if the array set is uniform. Code 1 if it is not uniform. |

If a value of IVAR equal to unity is specified, then the program will seek subsequent node or cell data sufficient to define the variable at every node or cell in the block. This information is input through the temporary variable $\mathrm{A}(\mathrm{I}, \mathrm{J})$ with an 8 F 10.0 format in the following manner:

| Line columns | Variable | Definition |
| :---: | :---: | :---: |
| 1-10 | A(IB,JB) | Temporary variable specifying nodal or cell values for the arrays in data sets $D$ through $K$. Note that the information is read in row by row for the grid, each new row beginning a new line. |
| 11-20 | A(IB+1,JB) |  |
| ! | ! |  |
|  | A(IE,JB) |  |
| 1-10 | A(IB, JB+1) |  |
| , | $\mathrm{A}(\mathrm{IB}+1, \mathrm{JB}+1)$ |  |
|  | ! |  |
|  | A(IE,JB+1) |  |
|  | : |  |
|  | A(IE,JE) |  |

Any variable that is not defined by blocking over a particular part of the grid will be automatically set to zero on that part.
Data sets $D$ and $E$, representing the internal spatial dimensions of the grid, are read in by blocking for convenience to the programmer. Because both the horizontal spacing DX and vertical spacing DY (as measured from the lower left corner of the grid) are, in reality, singly dimensioned arrays, it is necessary to set JB and JE equal to unity for both variables. The variable IE then equals ID-1 in the case of DX, and JD-1 in the case of DY (IB equals one, of course, in both cases). Variable grid spacing can be input by specifying a value of IVAR equal to unity and following this with the necessary array information in an 8 F 10.0 format.
Note that only cell values of $x$ - and $y$-direction transmissivities, hydraulic conductance, and recharge (data sets $F, G, H$, and $J$ ) are read in through the blocking scheme. Data set P contains the zonal values by which these cell values are multiplied. A typical example of usage would be to form transmissivity as the product of hydraulic conductivity and thickness. Data sets F and G would contain the variable thickness of the aquifer, and the variable in data set $P$ would represent the hydraulic conductivity zone by zone. Their product would be the transmissivity.

Data Set L.
Two integer variables also are input by blocking. Both are defined below, as read in by a 1615 format, although only that variable associated with data set L is input at this location.

| Data Set | Variable | Type of Variable | Definition |
| :--- | :--- | :--- | :--- |
| L | IZN | Cell array $\ldots .$. | Zone number of each cell. Each cell having a nonzero zone <br> number must have CX or CY $>0$. |
| S | IN | Nodal array $\ldots$ | Denotes specified head. Set to -1 at nodes where head is <br> specified, including nodes in segments involving specified- <br> head regression parameters, and leave as zero at all remain- <br> ing nodes. |

The initial line for these data is identical to that of the real variable case. The line defining blocks into which integer variables are divided is similar to that of the real variable case with the exception of the variable IFACT, as noted subsequently (format, 6I5):

| Line columns | Variable | Definition |
| :---: | :---: | :---: |
| 1-5 | IB | Beginning column of the rectangular input block. |
| 6-10 | IE | Final columun of the rectangular input block. |
| 11-15 | JB | Beginning row of the rectangular input block. |
| 16-20 | JE | Final row of the rectangular input block. |
| 21-25 | IFACT | If the array set is uniform for the entire block, IFACT is the cell or nodal value that is assigned to each grid point. If the array set is not uniform, each cell or nodal value on the subsequent data lines will not be multiplied by IFACT. |
| 26-30 | IVAR | Code 0 if the array set is uniform. Code 1 if it is not uniform. |

The nonuniform integer input is identical to the real variable case, except that the temporary variable INT(I,J) input with format 1615 is used in place of $A(I, J)$.

## Data set M.

Observed head data; set contains NOBS lines (format 3I5,4F10.0).

| Line columns | Variable |  | Definitions |
| :--- | :--- | :--- | :--- |
| $\mathbf{1 - 5}$ | N | $\ldots \ldots \ldots \ldots$ | Observation number. |
| $\mathbf{6 - 1 0}$ | IL $\ldots \ldots \ldots \ldots$ | Cell column in which observation lies. |  |
| $\mathbf{1 1 - 1 5}$ | JL $\ldots \ldots \ldots \ldots$ | Cell row in which observation lies. |  |
| $\mathbf{1 6 - 2 5}$ | XL $\ldots \ldots \ldots$ | x location of observation. |  |
| $\mathbf{2 6 - 3 5}$ | YL $\ldots \ldots \ldots$ | y location of observation. |  |
| $\mathbf{3 6 - 4 5}$ | HO(N) $\ldots \ldots \ldots$ | Observed value of head. |  |
| $\mathbf{4 6 - 5 5}$ | W(N) $\ldots \ldots \ldots$ | Reliability weight, $\omega$. |  |

Observations are numbered from 1 through NOBS but may be read in any order. Observation N is assumed to lie in cell (IL,JL) at $\boldsymbol{x}$ and $\boldsymbol{y}$ location (XL,YL). If the observation lies on a node point, it may be assigned to any adjacent cell bounded by the node. The origin for $x$ and $y$ is assumed to be node ( 1,1 ). Omit this data set if NOBS $=0$.

Data Set N .
Aquifer regression parameter numbers; set contains NZNS lines (format, 16I5).

| Line columns | Variable | Definition |
| :---: | :---: | :---: |
| 1-5 | I | Zone number. |
| 6-10 | IPRM(1,I) | Parameter number of $x$-transmissivity in zone $I$. Code 0 if it is not a regression parameter. |
| 11-15 | IPRM(2,I) | Parameter number of $y$-transmissivity in zone I. Code 0 if it is not a regression parameter. |
| 16-20 | IPRM(3,I) | Parameter number for hydraulic conductance in zone I. Code 0 if it is not a regression parameter. |
| 21-25 | IPRM(4,I) | Parameter number for distributed recharge in zone I. Code 0 if it is not a regression parameter. |

Lines may appear in any order with respect to zone number, but there must be NZNS lines. Parameters may have any number from 1 through NVAR. Note that parameters of the same (or even different) property in different zones may have identical parameter numbers. Omit this data set if NPAR equals zero.

Data Set 0 .
Standard deviations for aquifer regression parameters; one value per line for a total of NPAR lines (format, I5,F10.0).

| Line columns | Variable |  | Definition |
| :--- | :--- | :--- | :--- |
| $1-5$ | K | $\ldots \ldots \ldots \ldots$ | Aquifer parameter number. |
| $6-15$ | $\mathrm{WP}(\mathrm{K})$ | $\ldots \ldots$ | Standard deviation of each aquifer regression parameter. Code |
|  |  |  | 0.0 if no prior information exists for the parameter. |

Omit data set if NPAR equals zero.
Data Set P.
Zonal aquifer property values; set contains NZNS lines (format, I5,4F10.0).

| Line columns | Variable | Definitions |
| :---: | :---: | :---: |
| 1-5 | I | Zone number. |
| 6-15 | PAR(1) | Zonal $x$-transmissivity value for zone I. |
| 16-25 | PAR(2) | Zonal $\boldsymbol{y}$-transmissivity value for zone I. |
| 26-35 | PAR(3) | Zonal hydraulic conductance value for zone I. |
| 36-45 | PAR(4) | Zonal distributed recharge value for zone I. <br> (PAR(4)*QR has units of volumetric rate per unit area.) |

Lines may appear in any order with respect to zone number, but there must be NZNS lines.

## Data Set Q.

Known point flow rates; set contains NWELS lines (format, 215,F10.0).

| Line columns | Variable |  | Definition |  |
| :--- | :--- | :--- | :--- | :--- |
| $\mathbf{1 - 5}$ | I | $\ldots \ldots \ldots$ | Column location of point flow. |  |
| $6-10$ | J | $\ldots \ldots \ldots$ | $\ldots$ | Row location of point flow. |
| $11-20$ | WELL(I,J) | $\ldots$ | Total volumetric flow to or from node, negative for withdrawal. |  |

If a point flow is located between node points, the total rate can be apportioned among the four adjacent nodes using bilinear or similar interpolation. The apportioned flow rates must be supplied by the user. Omit data set if NWELS equals zero.

## Data Set R.

Specified boundary-flow zones and flow-zone parameters; set contains NQBZ lines, one line for each zone (format, 5I5,3F10.0).

| Line columns | Variable | Definition |
| :---: | :---: | :---: |
| 1-5 | IA | Column location of the $A$ end of the segment (zone). |
| 6-10 | JA | Row location of the $A$ end of the segment (zone). |
| 11-15 | IB | Column location of the $B$ end of segment (zone). |
| 16-20 | JB | Row location of the $B$ end of segment (zone). |
| 21-25 | IP | Regression parameter number. Set equal to zero if QB is not a regression parameter. |
| 26-35 | QB | Zonal flow value. |
| 36-45 | SDQB | Standard deviation for regression parameter. Code as 0.0 if no prior information exists on the parameter. |
| 46-55 | QBM | Multiplier for zonal flow value. |

Note that $I A \leq I B$ and $J A \leq J B$, which define the $A$ and $B$ ends of the segment. If IA equals IB and JA equals JB, the flow is restricted to a single node. In this case, the product QB*QBM equals total volumetric flow into or out of the node. Otherwise, the product is a volumetric rate per unit cell width. When data set is used to model flow boundary conditions, zone must follow either a row or column. Regression parameters can have any numbers from 1 to NVAR. Note that by setting IP equal to zero, a fixed specified-flow condition is simulated. Omit data set when NQBZ equals zero.

## Data Set S.

Specified-head boundary designation; see data set L. This data set, when used in conjunction with data set K or T (following), can be used to construct specified-head boundaries. In particular, nodes designated in this data set by -1 are forced to take on values specified in data set K or T . In addition to peripheral boundary conditions, this data set, in conjunction with data set K or T, can be used to model other constant head conditions such as bodies of open water.

## Data Set T.

This data set defines specified-head boundary parameter zones. Each zone is associated with a subset of the $T$ data and the number of subsets will equal NHBZ. The initial line in each subset contains size and descriptive information about the zone and appears as follows (format, 4I5,2F10.0).

| Line columns | Variable | Definition |
| :---: | :---: | :---: |
| 1-5 | IZ | Segment or zone number. |
| 6-10 | NN | Number of nodes in the segment. |
| 11-15 | IBPA | Regression parameter number for $A$ end of segment. Set equal to zero if head at the $A$ end is not a regression parameter. |
| 16-20 | IBPB | Regression parameter number for $B$ end of segment. Set equal to zero if head at the $B$ end is not a regression parameter. |
| 21-30 | SDHA | Standard deviation of head at A end of segment. Code as 0.0 if no prior information exists on the parameter. |
| 31-40 | SDHB | Standard deviation of head at B end of segment. Code as 0.0 if no prior information exists on the parameter. |

The A and B ends are arbitrary. Segments are numbered from 1 through NHBZ, and regression parameters may have any number from 1 through NVAR. A single head change can be found for all the intermediate specified head nodes by allowing IBPA to equal IBPB.

After reading the initial line, the program then seeks NN subsequent lines in each subset that define the heads along the boundary segment (format, 215,F10.0).

| Line columns | Variable |  |
| :--- | :--- | :--- |
| $1-5$ | ILOC $\ldots \ldots \ldots$ | Definition |
| 6-10 | JLOC $\ldots \ldots \ldots$ | Row location of node. |
| $11-20$ | V $\ldots \ldots \ldots$. | Estimated head at node ILOC, JLOC. |

Note that any shape of head surface can be input along a boundary segment. Because heads at the segment ends are the only regression parameters in each segment, their influence is distributed to the intermediate nodes by linear interpolation. The linear interpolation is based upon distance from the
parameter in question, with a weight of one assigned at the end node occupied by the parameter and zero at the node of the opposite end of the segment. If zeros are assigned to the head variables $V$ at intermediate nodes in a boundary segment, then the program automatically assumes that the head surface along the segment is simply a straight line between the heads specified at the end nodes. Omit the entire data set if NHBZ equals zero.

The following data sets are required only if the solution-only option (ISO, data set $B$ ) is specified. In the following description input variables are loosely termed parameters for convenience, but it must be realized that they are not actually regression parameters because no regression is performed. Also, data set $V$ requires that the parameters being varied have nonzero numbers corresponding to regression parameter numbers defined for the initial solution. Hence, the input for the initial solution must be coded as if it were to be a regression for these parameters although no regression will actually be performed.

## Data Set U.

Additional solution specification, one line (format, I5).

| Line columns | Variable | Definition |  |
| :--- | :--- | :--- | :--- |
| $\mathbf{1 - 5}$ | N | $\ldots \ldots \ldots \ldots$ | Number of solutions required using alternative parameter sets. <br> Code 0 if a solution is desired only for the initial set of <br> parameters. |

Data Set V.
New set of parameters; the set contains NVAR lines (format I5,F10.0).

| Line columns | Variable |  |  |
| :--- | :--- | :--- | :--- |
| $\mathbf{1 - 5}$ | I $\ldots \ldots \ldots \ldots$ | Definition |  |
| 6-15 | B(I) $\ldots \ldots \ldots$ | Parameter number. |  |

Parameters are numbered in the same order as used for the initial solution. This data set is repeated N times, and the program will compute and print the solution corresponding to each set of parameters.

Output.-The following discussion gives the content and order of the output obtained from the program. It should be noted that some of the output is only obtained under certain specified conditions. All output is clearly labeled. However, order numbers in the following discussion do not appear in the output; they are for convenience in listing the order of output only. The statistical measures cited below are described in section 5.

1. Three title lines (data set A).
2. Problem size information (data set B).
3. Special input parameters (data set C).
4. Array sets (data sets $D$ through L). For each variable, input information (variable name, block number, IB, IE, JB, JE, and value of FACT or IFACT) for each block is listed block by block. If specified, this information is then followed by the values of all entries in the array.
5. Observed head data (data set $M$ ).
6. Aquifer regression parameter numbers (data set N ). This input is printed only if NPAR is greater than zero.
7. Coefficients of variation for aquifer regression parameters (data set 0 ). This input is printed only if NPAR is greater than zero.
8. Zonal aquifer properties (data set P).
9. Known point-flow rates (data set $Q$ ). This input is printed only if NWELS is greater than zero.
10. Specified boundary-flow information (data set R). This input is printed only if $N Q B Z$ is greater than zero.
11. Specified boundary-head distribution (data set S). This input is printed in the same manner as indicated under 4.
12. Specified boundary-head information (data set T). For each segment, input information (segment number, number of nodes in the segment, IBPA, IBPB, SDHA, and SDHB) is printed followed by a listing of node locations and input values of specified heads. This input is printed only if NHBZ is greater than zero.
13. Error message. If $I Z N>0, C X \leq 0$, and $C Y \leq 0$ occur at the same cell, then the following error message is printed: "AT

CELL $(i, j), \quad$ IZN $>0, \quad C X=0, A N D$ CY=0."
14. Error message. If the error in 13 happens at one or more cells, then the following message is printed: "PROGRAM ABORTED BECAUSE OF CONFLICT BETWEEN IZN, CX, AND CY." Execution then terminates.
15. Error message. If an active node is isolated from other active nodes, then the following message is printed: ACTIVE NODE ( $i, j$ ) CANNOT BE ISOLATED." Execution then terminates.
16. Information on matrix solution procedure used to compute heads. Message that solution is by LDU factorization is followed by the computed minimum required dimensions of arrays used in the solution.
17. Initial solution for heads.
18. Program branch. If the solution-only option is specified, then output skips to 35.
19. Number of parameters having prior information.
20. Sensitivity matrix $\underset{X}{X}$. This is printed only if optional print-out was selected.
21. Error message. If any diagonal term of the coefficient matrix $\underline{S}_{s}^{T} \underline{V}_{s}^{-1} \underline{S}_{s}+\underline{S}_{p}^{T} \underline{U}^{-1}$ $\cdot S_{p} s^{2}$ of the normal equations is smaller than $10^{-10}$, then the following message is printed: "SENSITIVITIES FOR PARAMETER $i$ EFFECTIVELY ZERO." If this error occurs, the current weighted residuals $\underline{\hat{u}}$ (section 5.5.1) are then printed. If the sensitivity print and orthogonalization option was selected, the sensitivities and orthogonalized sensitivities are also printed in the forms given in 33 and 34. Execution then terminates.
22. Coefficient matrix $\underline{S}_{s}^{T} \underline{V}_{s}^{-1} \underline{S}_{s}+S_{p}^{T} U^{-1} S_{p} s^{2}$ $+\mu I$ and the gradient vector $S_{s}^{T} V_{s}^{-1}\left(Y_{s}\right.$ $\left.-f_{s}(\xi, \underline{b})\right)+S_{p}^{T} \underline{U}^{-1} s^{2}\left(Y_{p}-f_{p}(\xi, \underline{b})\right)$. This output is printed only if the optional printout was selected. (Note: If prior information of unknown reliability was used, then the matrix and vector will be modified to include this information.)
23. Error message. If, during solution of the nurmal equations, it becomes evident
that the problem is singular, then the following message is printed: "LEAST SQUARES COEFFICIENT MATRIX SINGULAR; SOLUTION FOR PARAMETERS NOT UNIQUE." The course of action is then the same as that given under 21.
24. Iteration number, current sum of weighted, squared deviations of computed from observed heads, determinant of the current least squares coefficient matrix defined in 22 above, current value of $\mu$, current value of $\rho$, followed by the current parameter vector $b_{r+1}$ defined by equation 3.3-19.
25. Error message. If a parameter is more than one thousand times smaller in magnitude than initially specified, then, the message "PARAMETER $i$ EFFECTIVELY ZERO" is printed. If this problem occurs, then further iterations are aborted, the current solution is taken as the final one, and the course of action given in 21 is taken.
26. Solution converged message and final number of iterations. If the solution did not converge in the allotted number of iterations, then a message to this effect is written instead, and output skips to 32.
27. Error message. If the coefficient matrix of the normal equations (see 22) is singular when $\mu=0$, and this status has not been detected because $\mu>0$ has been computed and used by the program, then the message given in 23 is printed. In this case the subsequent course of action is the same as in 21.
28. Error variance $\left(s^{2}\right)$ (section 5.4.1), final total sum of squares (sum of weighted, squared deviations of computed from observed heads plus sum of weighted, squared deviations of computed from prior estimates of parameters), and correlation coefficient ( $R_{y}$ ) (section 5.4.2).
29. Final parameter estimates and their estimated standard errors (section 5.4.3).
30. Estimated variance-covariance matrix, $\left(\underline{X}_{\underline{X}}^{T} \underline{V}_{s}^{-1} \underline{X}_{s}+\underline{X}_{p}^{T} \underline{U}^{-1} \underline{X}_{p} s^{2}\right)^{-1} s^{2}$ (section
5.4.3). The ordering of rows and columns matches that for $\underline{b}$ in 24 .
31. Correlation matrix for parameters, $\left\{r_{i j}\right\}$ $=\left\{\operatorname{Cov}\left(b_{i}, b_{j}\right) /\left(\operatorname{Var}\left(b_{i}\right) \cdot \operatorname{Var}\left(b_{j}\right)\right)^{1 / 2}\right\}$ (section 5.4.4). Again, the ordering matches that for $b$.
32. Computed and observed heads, and weighted residuals, $\underline{\hat{u}}$ (section 5.5.1).
33. Nodal sensitivities printed parameter by parameter. These are printed only if the sensitivity print and orthogonalization option was selected or if the solution did not converge.
34. Orthogonalized, scaled sensitivities, $Q$, with the sensitivities for the prior information forming the last $n_{p}$ rows. These are printed only if the sensitivities in 33 above are printed. This completes the output for regression solutions.
35. New parameters (data set $V$ ), all in sequential order; the solution number, and the solution for heads. These are printed for all solutions when the solution-only option is invoked.

Example Problem.-The following example problem illustrates use of most of the program options. As illustrated in figure 4.3-8, the modeled area consists of three aquifer zones bounded by three specified boundary-flow zones containing three boundary-flow regression parameters, a no-flow boundary, and two specified boundary-head segments containing three boundary-head regression parameters. Initial values for the regression parameters are:

```
\(q_{B 1}=8\)
\(q_{B 2}=0.8\)
\(q_{B 3}=1\)
\(h_{B 1}=40\)
\(h_{B 2}=10\)
\(h_{B 3}=16\)
\(T_{1}=4,000\)
\(T_{2}=400\)
```

$$
\begin{aligned}
& T_{3}=1,000 \\
& W_{1}=0.0001 \\
& W_{2}=0.0005 \\
& W_{3}=-0.0001 \\
& R_{1}=0.001 \\
& R_{2}=0.0007 \\
& R_{3}=0.0015
\end{aligned}
$$

Parameter $q_{B 1}$ has prior information with a standard deviation of 0.8 on it. Because all aquifer properties and boundary flows are constant within their respective zones, multipliers for these parameters may be assigned values of unity. Assume that the estimated error variance


Figure 4.3-8
for use with the prior information has a value of unity, $t_{m x}=1.5$, and $\cos \Theta_{m x}=0.08$.

Locations of observed heads are indicated by the small open circles on figure 4.3-8. Values corresponding to these locations are given in table 4.3-2. All observations have a weight $\omega_{i i}$ of unity.

Table 4.3-2

| Node | Value | Node | Value |
| :---: | :---: | :---: | :---: |
| $(1,1)$ | 29.51 | $(10,4)$ | 22.46 |
| $(3,1)$ | 26.38 | (12.4) | 17.16 |
| $(5,1)$ | 25.16 | $(14,4)$ | 15.74 |
| (7,1) | 26.81 | $(16,4)$ | 14.69 |
| $(9,1)$ | 23.82 | $(1,5)$ | 36.05 |
| $(11,1)$ | 23.59 | $(3,5)$ | 31.64 |
| $(13,1)$ | 16.68 | $(5,5)$ | 32.91 |
| $(15,1)$ | 15.31 | $(7,5)$ | 25.05 |
| $(2,2)$ | 28.27 | $(9,5)$ | 25.47 |
| $(4,2)$ | 25.17 | $(11,5)$ | 20.39 |
| $(6,2)$ | 26.94 | $(13,5)$ | 14.40 |
| $(8,2)$ | 23.59 | $(15,5)$ | 14.37 |
| $(10,2)$ | 22.53 | $(2,6)$ | 38.66 |
| $(12,2)$ | 17.89 | $(4,6)$ | 34.51 |
| $(14,2)$ | 15.87 | $(6,6)$ | 27.81 |
| $(16,2)$ | 15.98 | $(8,6)$ | 25.98 |
| $(1,3)$ | 31.25 | $(10,6)$ | 22.16 |
| $(3,3)$ | 27.50 | $(12,6)$ | 15.73 |
| $(5,3)$ | 26.22 | $(14,6)$ | 14.79 |
| $(7,3)$ | 25.65 | $(16,6)$ | 10.48 |
| $(9,3)$ | 24.35 | $(1,7)$ | 41.28 |
| $(11,3)$ | 22.48 | $(3,7)$ | 34.82 |
| $(13,3)$ | 15.85 | $(5,7)$ | 32.00 |
| $(15,3)$ | 15.71 | $(7,7)$ | 26.97 |
| $(2,4)$ | 29.36 | $(9,7)$ | 25.67 |
| $(4,4)$ | 27.76 | $(11,7)$ | 18.12 |
| $(6,4)$ | 25.51 | $(13,7)$ | 18.40 |
| $(8,4)$ | 24.43 | $(15,7)$ | 11.90 |

Heads $H$ on the distal side of the aquitard are computed by first assigning constant values to all cells in an aquifer property zone, then computing nodal values as the average of all adjacent cell values. Cell values are 25 for zone 1, 35 for zone 2 , and 15 for zone 3.

Input data for the example problem are coded on figure 4.3-9 and are contained in file EXPROB.DAT in the diskette accompanying this report. These data should be compared with the data input instructions given above.
Output is given in figure 4.3-10. Statistical measures listed are described in section 5.

## EXAMPLE PROBLEM FOR DATA INPUT



TECHNIQUES OF WATER-RESOURCES INVESTIGATIONS

| 21 | 9 | 3 | 8000 | 2000 | 24.35 | 1 |
| :---: | :---: | :---: | :---: | :---: | :---: | :---: |
| 22 | 11 | 3 | 10000 | 2000 | 22.48 | 1 |
| 23 | 13 | 3 | 12000 | 2000 | 15.85 | 1 |
| 24 | 15 | 3 | 14000 | 2000 | 15.71 | 1 |
| 25 | 2 | 4 | 1000 | 3000 | 29.36 | 1 |
| 26 | 4 | 4 | 3000 | 3000 | 27.76 | 1 |
| 27 | 6 | 4 | 5000 | 3000 | 25.51 | 1 |
| 28 | 8 | 4 | 7000 | 3000 | 24.43 | 1 |
| 29 | 10 | 4 | 9000 | 3000 | 22.46 | 1 |
| 30 | 12 | 4 | 11000 | 3000 | 17.16 | 1 |
| 31 | 14 | 4 | 13000 | 3000 | 15.74 | 1 |
| 32 | 15 | 4 | 15000 | 3000 | 14.69 | 1 |
| 33 | 1 | 5 | 0 | 4000 | 36.05 | 1 |
| 34 | 3 | 5 | 2000 | 4000 | 31.64 | 1 |
| 35 | 5 | 5 | 4000 | 4000 | 32.91 | 1 |
| 36 | 7 | 5 | 6000 | 4000 | 25.05 | 1 |
| 37 | 9 | 5 | 8000 | 4000 | 25.47 | 1 |
| 38 | 11 | 5 | 10000 | 4000 | 20.39 | 1 |
| 39 | 13 | 5 | 12000 | 4000 | 14.40 | 1 |
| 40 | 15 | 5 | 14000 | 4000 | 14.37 | 1 |
| 41 | 2 | 6 | 1000 | 5000 | 38.66 | 1 |
| 42 | 4 | 6 | 3000 | 5000 | 34.51 | 1 |
| 43 | 6 | 6 | 5000 | 5000 | 27.81 | 1 |
| 44 | 8 | 6 | 7000 | 5000 | 25.98 | 1 |
| 45 | 10 | 6 | 9000 | 5000 | 22.16 | 1 |
| 46 | 12 | 6 | 11000 | 5000 | 15.73 | 1 |
| 47 | 14 | 6 | 13000 | 5000 | 14.79 | 1 |
| 48 | 15 | 6 | 15000 | 5000 | 10.48 | 1 |
| 49 | 1 | 6 | 0 | 6000 | 41.28 | 1 |
| 50 | 3 | 6 | 2000 | 6000 | 34.82 | 1 |
| 51 | 5 | 6 | 4000 | 6000 | 32.00 | 1 |
| 52 | 7 | 6 | 6000 | 6000 | 26.97 | 1 |
| 53 | 9 | 6 | 8000 | 6000 | 25.67 | 1 |
| 54 | 11 | 0 | 10000 | 6000 | 18.12 | 1 |
| 55 | 13 | 6 | 12000 | 6000 | 18.40 | 1 |
| 56 | 15 | 6 | 14000 | 6000 | 11.90 | 1 |
| 1 | 1 | 1 | 47 |  |  |  |
| 2 | 2 | 2 | 58 |  |  |  |
| 3 | 3 | 3 | 69 |  |  |  |
| 1 |  | 0 |  |  |  |  |
| 2 |  | 0 |  |  |  |  |
| 3 |  | 0 |  |  |  |  |
| 4 |  | 0 |  |  |  |  |
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| 6 |  | 0 |  |  |  |  |
| 7 |  | 0 |  |  |  |  |
| 8 |  | 0 |  |  |  |  |
| 9 |  | 0 |  |  |  |  |
| 1 |  | 4000 | 4000 | . 001 | . 0001 |  |
| 2 |  | 400 | 400 | . 0007 | . 0005 |  |
| 3 |  | 1000 | 1000 | . 0015 | -. 0001 |  |
| 1 | 1 | 1 | 410 | 8 | . 8 | 1 |
| 1 | 4 | 1 | 711 | . 8 | 0 | 1 |


| 8 | 1 | 16 | 1 | 12 |  | 1 |  | 0 |
| :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: |
| IN | 2 | 0 |  |  |  |  |  |  |
| 16 | 16 | 1 | 7 | -1 |  |  |  |  |
| 1 | 15 | 7 | 7 | -1 |  |  |  |  |
| 1 | 16 | 13 | 14 |  | 0 |  | 0 |  |
| 1 | 7 |  | 40 |  |  |  |  |  |
| 2 | 7 |  |  |  |  |  |  |  |
| 3 | 7 |  |  |  |  |  |  |  |
| 4 | 7 |  |  |  |  |  |  |  |
| 5 | 7 |  |  |  |  |  |  |  |
| 6 | 7 |  |  |  |  |  |  |  |
| 7 | 7 |  |  |  |  |  |  |  |
| 8 | 7 |  |  |  |  |  |  |  |
| 9 | 7 |  |  |  |  |  |  |  |
| 10 | 7 |  |  |  |  |  |  |  |
| 11 | 7 |  |  |  |  |  |  |  |
| 12 | 7 |  |  |  |  |  |  |  |
| 13 | 7 |  |  |  |  |  |  |  |
| 14 | 7 |  |  |  |  |  |  |  |
| 15 | 7 |  |  |  |  |  |  |  |
| 16 | 7 |  | 10 |  |  |  |  |  |
| 2 | 7 | 14 | 15 |  | 0 |  | 0 |  |
| 16 | 7 |  | 10 |  |  |  |  |  |
| 16 | 6 |  |  |  |  |  |  |  |
| 16 | 5 |  |  |  |  |  |  |  |
| 16 | 4 |  |  |  |  |  |  |  |
| 16 | 3 |  |  |  |  |  |  |  |
| 16 | 2 |  |  |  |  |  |  |  |
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Figure 4.3-9—Continued

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EXAMPLE PROBLEM FOR DATA INPUT


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| 3 | . $84556 \mathrm{E}-04$ | . $35640 \mathrm{E}-03$ | -.42379E-04 | -1104.0 | 124.56 | -18.535 | 877.71 |
| :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: |
|  | 91.807 | 11.550 | .54188E-01 | . $12744 \mathrm{E}-01$ | .29967E-01 | . 22834E-01 | .17803E-01 |
|  | .65461E-04 |  |  |  |  |  |  |
| 4 | .20950E-04 | .19118E-03 | -. 11934E-03 | -377.48 | 78.699 | -52.851 | 851.92 |
|  | 49.653 | $32.568$ | .17697E-01 | .44771E-02 | . 10882 | .33990E-01 | .30295E-01 |
|  | .18996E-03 |  |  |  |  |  |  |
| 5 | -.44029E-04 | .83352E-04 | -. 34274E-03 | 195.90 | 36.560 | -155.50 | 774.03 |
|  | 21.594 | 93.554 | .58100E-02 | .15130E-02 | . 37073 | . $34274 \mathrm{E}-01$ | .35670E-01 |
|  | .57555E-03 |  |  |  |  |  |  |
| 6 | .26304E-03 | . 36640E-04 | -. 98740E-03 | 64.792 | 16.440 | -468.15 | 542.22 |
|  | 9.4671 | 265.38 | . 22527E-02 | . $59308 \mathrm{E}-03$ | . 49558 | .23666E-01 | .27636E-01 |
|  | . $17904 \mathrm{E}-02$ |  |  |  |  |  |  |
| 7 | .28057E-04 | . $35984 \mathrm{E}-05$ | . $30068 \mathrm{E}-03$ | 7.7542 | 1.6281 | -698.55 | 52.278 |
|  | . 92848 | 605.89 | .21021E-03 | . $55570 \mathrm{E}-04$ | . 67776 | . $30985 \mathrm{E}-02$ | . $94553 \mathrm{E}-02$ |
|  | .25686E-01 |  |  |  |  |  |  |
| 8 | .25379E-05 | . $31618 \mathrm{E}-06$ | -. 12566E-03 | . 74317 | . 14383 | -249.14 | 4.5450 |
|  | .81492E-01 | 453.52 | .17832E-04 | . $47267 \mathrm{E}-05$ | . 53296 | . $40054 \mathrm{E}-03$ | .24008E-01 |
|  | . 29071 |  |  |  |  |  |  |
| 9 | -.83160E-04 | .57567E-03 | -. 11317E-04 | -2294.8 | 99.794 | -4.9144 | 873.71 |
|  | 147.71 | 3.0815 | . 26739 | .48691E-01 | . $70743 \mathrm{E}-02$ | .11698E-01 | .65538E-02 |
|  | .17175E-04 |  |  |  |  |  |  |
| 0 | .74473E-04 | . $49706 \mathrm{E}-03$ | -. 25302E-04 | -1475.1 | 149.77 | -11.007 | 869.67 |
|  | 126.28 | 6.8912 | .91556E-01 | . $21604 \mathrm{E}-01$ | .16239E-01 | .18505E-01 | . $13061 \mathrm{E}-01$ |
|  | .38574E-04 |  |  |  |  |  |  |
| 1 | .47626E-04 | . $30114 \mathrm{E}-03$ | -.69809E-04 | -697.32 | 120.92 | -30.536 | 855.15 |
|  | 78.330 | 19.029 | .30075E-01 | . $76372 \mathrm{E}-02$ | .47836E-01 | . $33684 \mathrm{E}-01$ | .27683E-01 |
|  | . $10790 \mathrm{E}-03$ |  |  |  |  |  |  |
| 2 | -.78560E-05 | . $13518 \mathrm{E}-03$ | -. 19833E-03 | -50.049 | 58.932 | -87.936 | 812.07 |
|  | 35.078 | 54.179 | . $96800 \mathrm{E}-02$ | . 25227E-02 | . 14360 | . $41772 \mathrm{E}-01$ | . $40022 \mathrm{E}-01$ |
|  | .31695E-03 |  |  |  |  |  |  |
| 3 | .66592E-04 | .54863E-04 | -. 56959E-03 | 331.28 | 24.679 | -258.98 | 679.82 |
|  | 14.175 | 156.44 | . $33205 \mathrm{E}-02$ | . $87589 \mathrm{E}-03$ | . 23450 | . $34954 \mathrm{E}-01$ | .39677E-01 |
|  | . $96774 \mathrm{E}-03$ |  |  |  |  |  |  |
| 4 | .93956E-04 | . $11523 \mathrm{E}-04$ | . $78272 \mathrm{E}-03$ | 26.881 | 5.2279 | -919.11 | 166.18 |
|  | 2.9720 | 533.15 | .66123E-03 | .17509E-03 | . 21476 | . $99077 \mathrm{E}-02$ | .14989E-01 |
|  | .71191E-02 |  |  |  |  |  |  |
| 5 | . $90142 \mathrm{E}-05$ | .11049E-05 | .25398E-03 | 2.7048 | . 50386 | -286.38 | 15.800 |
|  | . 28467 | 587.43 | .61321E-04 | .16276E-04 | . 19256 | .14460E-02 | .23067E-01 |
|  | .78345E-01 |  |  |  |  |  |  |
| 6 | . 00000 | . 00000 | . 00000 | . 00000 | . 00000 | . 00000 | . 00000 |
|  | . 00000 | . 00000 | . 00000 | . 00000 | . 00000 | . 00000 | . 16667 |
|  | . 83333 |  |  |  |  |  |  |
| 7 | -.45126E-03 | . $77585 \mathrm{E}-03$ | -. 93026E-05 | -2315.3 | 79.365 | -4.0344 | 829.78 |
|  | 211.03 | 2.5325 | . 43630 | . 10673 | .56983E-02 | .12821E-01 | . $58605 \mathrm{E}-02$ |
|  | . $14073 \mathrm{E}-04$ |  |  |  |  |  |  |
| 8 | -. 30579E-04 | . 80329E-03 | -. 15146E-04 | -1711.6 | 182.15 | -6.5629 | 828.39 |
|  | 199.47 | 4.1228 | . 14804 | .43134E-01 | . $91359 \mathrm{E}-02$ | .16340E-01 | . $94350 \mathrm{E}-02$ |
|  | .22863E-04 |  |  |  |  |  |  |
| 9 | .28901E-05 | .57360E-03 | -. 39933E-04 | -925.31 | 239.20 | -17.263 | 821.58 |
|  | 150.39 | 10.867 | . $49074 \mathrm{E}-01$ | . $13689 \mathrm{E}-01$ | .22888E-01 | . $32544 \mathrm{E}-01$ | .24242E-01 |
|  | .59929E-04 |  |  |  |  |  |  |

798.48
$.50865 \mathrm{E}-01$
730.34
$.62438 \mathrm{E}-01$
513.48
$.50880 \mathrm{E}-01$
48.803
$.22167 \mathrm{E}-01$
4.1862
.10958
743.34
$.66367 \mathrm{E}-02$
741.46
$.16217 \mathrm{E}-01$
726.82
$.62115 \mathrm{E}-01$
700.61
$.96892 \mathrm{E}-01$
594.94
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143.17
$.47757 \mathrm{E}-01$
13.157
$.70359 \mathrm{E}-01$
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213.74
$.53671 \mathrm{E}-02$
216.49
$.14555 \mathrm{E}-01$
282.90
$.48924 \mathrm{E}-01$
572.98
.14489

| 113.08 | -47.651 |
| :--- | :--- |
| $.56894 \mathrm{E}-01$ | $.59526 \mathrm{E}-01$ |
| 42.275 | -139.00 |
| .11083 | $.58707 \mathrm{E}-01$ |
| 17.063 | -407.58 |
| .11689 | $.40150 \mathrm{E}-01$ |
| 1.6132 | -466.70 |
| $.66737 \mathrm{E}-01$ | $.59267 \mathrm{E}-02$ |
| .13791 | 29.768 |
| $.37527 \mathrm{E}-01$ | $.10161 \mathrm{E}-02$ |
| 139.64 | -4.1385 |
| $.57427 \mathrm{E}-02$ | $.16999 \mathrm{E}-01$ |
| 420.38 | -9.4268 |
| $.12207 \mathrm{E}-01$ | $.26347 \mathrm{E}-01$ |
| 259.77 | -27.335 |
| $.27894 \mathrm{E}-01$ | $.81928 \mathrm{E}-01$ |
| 82.965 | -69.148 |
| $.48881 \mathrm{E}-01$ | .10402 |
| 26.031 | -210.11 |
| $.64502 \mathrm{E}-01$ | $.83153 \mathrm{E}-01$ |
| 4.9790 | -781.62 |
| $.31186 \mathrm{E}-01$ | $.25569 \mathrm{E}-01$ |
| .44848 | -55.665 |
| $.17367 \mathrm{E}-01$ | $.57210 \mathrm{E}-02$ |
| .00000 | .00000 |
| .00000 | .00000 |
| -301.31 | -1.1524 |
| $.15875 \mathrm{E}-02$ | $.78076 \mathrm{E}-01$ |
| 901.99 | -1.9997 |
| $.25531 \mathrm{E}-02$ | $.73485 \mathrm{E}-01$ |
| 2002 | -7.8992 |
| $.75221 \mathrm{E}-02$ | $.94268 \mathrm{E}-01$ |
| 184.86 | -31.070 |
| $.21123 \mathrm{E}-01$ | .18904 |
|  |  |


| $\begin{aligned} & .25177 \mathrm{E}-03 \\ & \mathbf{3 0 . 0 2 7} \end{aligned}$ | $\begin{aligned} & .11028 \mathrm{E}-03 \\ & .14925 \mathrm{E}-01 \end{aligned}$ | $\begin{aligned} & -317.08 \\ & .40309 \mathrm{E}-02 \end{aligned}$ |
| :---: | :---: | :---: |
| .92081E-04 | . 32071 E-03 | 231.85 |
| 87.594 | . $48754 \mathrm{E}-02$ | . 1305 |
| . $37011 \mathrm{E}-04$ | -.93635E-03 | 105.13 |
| 258.23 | . $19044 \mathrm{E}-02$ | . 50911 E |
| . $34962 \mathrm{E}-05$ | .55627E-03 | 8336 |
| 601.69 | .17858E-03 | 772 |
| . 29879 | -.48766E-04 | 83 |
| 451.40 | . 15211 | .40639E-05 |
| .13148E-02 | -.95542E-05 | -1554.3 |
| 2.6005 | . 22761 | . 11082 |
| 880E-02 | -.21865E- | -861 |
| 5.9472 | .76582E-01 | 25953 E |
| . $55544 \mathrm{E}-03$ | -.64379E-04 | -282. |
| 17.470 | .19856E-01 | 530 |
| .17295E | -. 165 | -11.25 |
| 44.753 | . $64338 \mathrm{E}-02$ |  |
| . 55001 E | -. 51144E-03 | 364.06 |
| 138.63 | .22930E-02 | . 62403 E |
| .10587E-04 | . $76168 \mathrm{E}-03$ | 4.853 |
| 517.97 | . $46383 \mathrm{E}-03$ | . 12548 |
| .95788E | . 84 |  |
| 574.43 | . 43482 E | . 11721 |
| . 00000 | . 00000 | . 00000 |
| . 00000 | . 00000 | . 00000 |
| -. 25598E-02 | -.26627E-05 | -477.25 |
| . 72464 | .95183E-01 | 1.2178 |
| -. 22910E-02 | -.46499E-05 | -345.96 |
| 1.2640 | . $40683 \mathrm{E}-$ | 73922 E |
| -.36705E-02 | -. 18796E- | -139.57 |
| 5.0881 | .13808E-01 | .71506E |
| . $35423 \mathrm{E}-03$ | -.75545E-04 | -1 |
| 20.368 | .61410E-02 | .18139E |




| 37-.65253E-05 |  | . $86770 \mathrm{E}-04$ | -. 23158E-03 | 219.88 | 42.569 | -91.671 | 547.07 |
| :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: |
|  |  | 61.963 | . $24664 \mathrm{E}-02$ | .69171E-03 | . $31053 \mathrm{E}-01$ | . 16382 | . 18057 |
| 38 | . $27163 \mathrm{E}-03$ | .28872E-04 | -. $81642 \mathrm{E}-03$ | 112.84 | 13.856 | -307.78 | 389.90 |
|  | $7.3247$ | 216.36 | . 10252E-02 | .28178E-03 | .28030E-01 | . 11166 | . 16798 |
| 39 | . $23516 \mathrm{E}-04$ | .25208E-05 | .43777E-03 | 9.3753 | 1.1989 | -333.52 | 34.509 |
|  | $\begin{aligned} & .64136 \\ & .10225 \mathrm{E}-01 \end{aligned}$ | 552.79 | .98203E-04 | . $26781 \mathrm{E}-04$ | .79262E-02 | . $26022 \mathrm{E}-01$ | . $99611 \mathrm{E}-01$ |
| 40 | .18717E-05 | .20407E-06 | -.38846E-03 | . 72193 | .96428E-01 | 279.22 | 2.8109 |
|  | . $52032 \mathrm{E}-01$ | 421.64 | . $84765 \mathrm{E}-05$ | .22995E-05 | . $30619 \mathrm{E}-02$ | . $70276 \mathrm{E}-02$ | . 25312 |
|  | . 10455 |  |  |  |  |  |  |
| 41 | -.25042E-04 | -. $45069 \mathrm{E}-03$ | -.94596E-06 | -112.33 | -613.46 | -. 40719 | 57.617 |
|  | 931.86 | . 25718 | . $17418 \mathrm{E}-01$ | . 26527 | . $52938 \mathrm{E}-03$ | . 26812 | .21579E-01 |
|  | . $14039 \mathrm{E}-05$ |  |  |  |  |  |  |
| 42 | -. $21672 \mathrm{E}-04$ | -.22841E-02 | -. 33126E-05 | -63.908 | 782.39 | -1.3912 | 70.690 |
|  | 893.31 | . 89634 | .68573E-02 | . $15118 \mathrm{E}-01$ | . $13611 \mathrm{E}-02$ | . 23864 | . $64002 \mathrm{E}-01$ |
|  | .46073E-05 |  |  |  |  |  | .64002E-01 |
| 43 | -.14945E-03 | .67722E-03 | -. 29884E-04 | 23.624 | 298.22 | -12.191 | 327.26 |
|  | 153.08 | 8.0378 | . $31654 \mathrm{E}-02$ | . $11004 \mathrm{E}-02$ | . $80304 \mathrm{E}-02$ | . 34417 | . 20931 |
|  | . $38360 \mathrm{E}-04$ |  |  |  |  |  |  |
| 44 | -.28535E-05 | .11684E-03 | -.71883E-04 | -35.333 | 56.396 | -28.215 | 358.20 |
|  | 28.327 | 19.149 | . $18309 \mathrm{E}-02$ | .53855E-03 | . $11594 \mathrm{E}-01$ | . 31371 | . 27945 |
|  | .82432E-04 |  |  |  |  |  |  |
| 45 | . 19335E-04 | .27719E-04 | -. $273488 \mathrm{E}-03$ | 275.39 | 13.427 | -98.168 | 314.97 |
|  | 6.9334 | 70.983 | . $76354 \mathrm{E}-03$ | .21377E-03 | .13371E-01 | . 23796 | . 33550 |
|  | . $23516 \mathrm{E}-03$ |  |  |  |  | . 237 | . 3355 |
| 46 | . $52326 \mathrm{E}-04$ | . $48415 \mathrm{E}-05$ | .67366E-03 | 21.772 | 2.3279 | -589.79 | 70.055 |
|  | 1.2223 | 374.84 | . $16018 \mathrm{E}-03$ | .44215E-04 | . $49386 \mathrm{E}-02$ | . 10476 | . 26050 |
|  | .15742E-02 |  |  |  |  |  |  |
| 47 | . $40914 \mathrm{E}-05$ | . $41281 \mathrm{E}-06$ | -. $12829 \mathrm{E}-03$ | 1.6529 | . 19695 | 99.608 | 5.7928 |
|  | . 10474 | 418.28 | .15304E-04 | . $41870 \mathrm{E}-05$ | .16250E-02 | .43177E-01 | . 30714 |
|  | .16394E-01 |  |  |  |  |  |  |
| 48 | . 00000 | . 00000 | . 00000 | . 00000 | . 00000 | . 00000 |  |
|  | . 00000 | . 00000 | . 00000 | . 00000 | . 00000 | . 00000 | . 83333 |
|  | . 16667 |  |  |  |  |  |  |
| 49 | . 00000 | . 00000 | . 00000 | . 00000 | . 00000 | . 00000 | . 00000 |
|  | . 00000 | . 00000 | . 00000 | . 00000 | . 00000 | 1.0000 | . 00000 |
|  | . 00000 |  |  |  |  |  |  |
| 50 | . 00000 | . 00000 | . 00000 | . 00000 | . 00000 | . 00000 | . 00000 |
|  | . 00000 | . 00000 | . 00000 | . 00000 | . 00000 | . 86667 | . 13333 |
|  | . 00000 |  |  |  |  |  |  |
| 51 | . 00000 | . 00000 | . 00000 | . 00000 | . 00000 | . 00000 | . 00000 |
|  | . 00000 | .00000 | . 00000 | . 00000 | . 00000 | . 73333 | . 26667 |
|  | . 00000 |  |  |  |  |  |  |
| 52 | . 00000 | . 00000 | . 00000 | . 00000 | . 00000 | . 00000 | . 00000 |
|  | . 00000 | . 00000 | . 00000 | . 00000 | . 00000 | . 60000 | . 40000 |
|  | . 00000 |  |  |  |  |  |  |



SCALED GRADIENT VECTOR
 following is for iteration 6 .
sCaled least squares matrix

$$
\begin{aligned}
& \bullet \\
& \text { in }
\end{aligned}
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$\infty$.25747
.95324
$.91209 \mathrm{E}-02$
.18514
$.45636 \mathrm{E}-01$
$.54336 \mathrm{E}-02$
.22633
1.0000
$.34231 \mathrm{E}-02$
.13920
.64480
$.93450 \mathrm{E}-02$
.26989
$.53228 \mathrm{E}-01$
$.17694 \mathrm{E}-04$
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$$
\begin{aligned}
& .58498 \\
& .43378 \mathrm{E}-01 \\
& -.5165 \mathrm{E}-01 \\
& 1.0000 \\
& .86505 \mathrm{E}-02 \\
& -.3430 \mathrm{E}-01 \\
& -.55596 \\
& -.18514 \\
& .2358 \mathrm{E}-01 \\
& -.57010 \\
& -.77598 \\
& . .29284 \mathrm{E}-01 \\
& -.97069 \mathrm{E}-02 \\
& . .8788 \mathrm{E}-01 \\
& .14509 \mathrm{E}-03
\end{aligned}
$$

\] | 4376 |
| :--- |
| 6090 | .14376

.16090

$.18259 \mathrm{E}-01$
$-.40783 \mathrm{E}-01$
$-.44719 \mathrm{E}-01$
$.18788 \mathrm{E}-01$
$.36389 \mathrm{E}-01$
$-.84249 \mathrm{E}-01$
.1543
$.53228 \mathrm{E}-01$
.14376
$.77942 \mathrm{E}-02$
$.12050 \mathrm{E}-01$
$.33840 \mathrm{E}-01$
.42776
1.0000
.27776

SCALED CRADIENT VECTOR
$\begin{array}{cccccc}\text { SCALED CRADIENT VECTOR } & & \\ -.51657 E-03 & .11084 \mathrm{E}-03 & -.20590 \mathrm{E}-03 & -.28675 \mathrm{E}-03 & -.93586 \mathrm{E}-03 \\ -.47620 \mathrm{E}-03 & .19812 \mathrm{E}-03 & -.25707 \mathrm{E}-03 & -.15761 \mathrm{E}-02 & -.30471 \mathrm{E}-03\end{array}$
$\begin{array}{lllll}\text { ITERATION NO. } & 6 & & & \\ \text { YQQ }=51.970 & \text { DET }(C)= & .45800 E-04 & \text { AMP }= & .00000 \\ \text { REGRESSION PARAMETERS } & & \\ 2865.3 & 117.67 & 497.85 & .99782 E-03 & .965 \\ .14220 E-03 & 8.0000 & .22321 & .89599 & 40 .\end{array}$

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SOLUTION CONVERGED IN 6 ITERATIONS
ERROR VARIANCE $=1.2374$
ESTIMATED SCM OF SQUARED ERRORS $=51.970$
CORRELATION COEFFICIENT $=.99033$
ELATION COEFFICIENT $=.99033$
ESTIMATED PARAMETER DATA
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|  | 1 | 2 |
| :--- | :--- | :--- |
|  |  |  |
| 1 | $.16734 \mathrm{E}+07$ | 92126. |
| 2 | 92126. | 71196. |
| 3 | $.19755 \mathrm{E}+06$ | -12931, |
| 4 | -.12629 | $.18984 \mathrm{E}-01$ |
| 5 | $.81201 \mathrm{E}-01$ | $.52827 \mathrm{E}-01$ |
| 6 | .47075 | $-.17544 \mathrm{E}-01$ |
| 7 | -.10060 | $-.84450 \mathrm{E}-01$ |
| 8 | .25244 | .19080 |
| 9 | $.54361 \mathrm{E}-01$ | $-.55455 \mathrm{E}-02$ |
| 10 | 283.68 | 11.633 |
| 11 | 235.55 | 129.38 |
| 12 | 437.83 | 9.4207 |
| 13 | 22.301 | -5.5050 |
| 14 | 26.987 | 1.8507 |
| 15 | -23.100 | -5.3819 |


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| 26.987 | -23.100 |
| 1.8507 | -5.3819 |
| 16.104 | -10.743 |
| $-.12935 \mathrm{E}-04$ | $.12131 \mathrm{E}-04$ |
| $-.16303 \mathrm{E}-05$ | $-.29205 \mathrm{E}-05$ |
| $.10117 \mathrm{E}-04$ | $-.99534 \mathrm{E}-04$ |
| $-.21829 \mathrm{E}-04$ | $.28963 \mathrm{E}-04$ |
| $.13911 \mathrm{E}-04$ | $-.19299 \mathrm{E}-04$ |
| $-.49939 \mathrm{E}-04$ | $-.89174 \mathrm{E}-04$ |
| $.52352 \mathrm{E}-16$ | $-.19250 \mathrm{E}-15$ |
| $-.20156 \mathrm{E}-02$ | $-.74691 \mathrm{E}-02$ |
| $.80067 \mathrm{E}-01$ | -.16715 |
| -.19602 | $.9971 \mathrm{E}-01$ |
| .44382 | -.23097 |
| -.23097 | 1.3307 |
| Figure $4.3-10$－Continued |  | $\cdots$





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NODAL SENSITIVITY ARRAYS
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## PARAMETER NUMBER <br> 6


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$3451 E-02$
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SCALED AND ORTHOGONALIZED SENSITIVITY MATRIX

| 8 |
| :---: |
| $-.30847 \mathrm{E}-01$ |
| $-.10784 \mathrm{E}-01$ |
| $-.11047 \mathrm{E}-01$ |
| $-.18262 \mathrm{E}-01$ |
| $-.18556 \mathrm{E}-01$ |
| $.38504 \mathrm{E}-02$ |
| $.92882 \mathrm{E}-02$ |
| $.90602 \mathrm{E}-02$ |
| $-.10764 \mathrm{E}-01$ |
| $-.28296 \mathrm{E}-02$ |
| $-.11724 \mathrm{E}-01$ |
| $-.18494 \mathrm{E}-01$ |
| $-. .98763 \mathrm{E}-02$ |
| $-.25158 \mathrm{E}-02$ |
| $.16658 \mathrm{E}-02$ |
| . .00000 |
| $-.41809 \mathrm{E}-02$ |
| $.19606 \mathrm{E}-01$ |
| $. .1134 \mathrm{E}-01$ |
| $-.1128 \mathrm{E}-01$ |
| $-. .14950 \mathrm{E}-01$ |
| $. .12575 \mathrm{E}-02$ |
| $-.12634 \mathrm{E}-02$ |
| $.10264 \mathrm{E}-02$ |
| $.00891 \mathrm{E}-01$ |
| $.66978 \mathrm{E}-01$ |
| $.15586 \mathrm{E}-01$ |
| $-.11471 \mathrm{E}-01$ |
| $.7360 \mathrm{E}-02$ |
| $-.34587 \mathrm{E}-02$ |
| $.82561 \mathrm{E}-03$ |

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# <div class="inline-tabular"><table id="tabular" data-type="subtable">
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</tr>
<tr style="border-top: none !important; border-bottom: none !important;">
<td style="text-align: center; border-left: none !important; border-bottom-style: solid !important; border-bottom-width: 1px !important; border-top: none !important; width: auto; vertical-align: middle; ">5</td>
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<td style="text-align: center; border-left: none !important; border-bottom: none !important; border-top: none !important; width: auto; vertical-align: middle; ">8</td>
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<table-markdown style="display: none">| 1 |
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| 8 |</table-markdown></div> 



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$10-3699 \varepsilon 66^{\circ}$ ．
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$23188 E-01$
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 $.86688 \mathrm{E}-03$
 둥웅
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 $-.48766 \mathrm{E}-01$
$-.306162 \mathrm{E}-02$
$-.1806 \mathrm{E}-03$
.11057
$.73254 \mathrm{E}-01$
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$\begin{array}{llllll}.78650 \mathrm{E}-02 & -.36457 \mathrm{E}-02 & -.61941 \mathrm{E}-01 & -.41596 \mathrm{E}-02 & -.18909 \mathrm{E}-01 & -.26062 \mathrm{E}-01 \\ .6477 \mathrm{E}-04 & .14940 \mathrm{E}-02 & -.16047 \mathrm{E}-02 & -.15188 \mathrm{E}-02 & .32843 \mathrm{E}-01 & -.63188 \mathrm{E}-02\end{array}$ 20-388Lを9$.33027 \mathrm{E}-02$
$.17000 \mathrm{E}-03$

 | $-.14467 \mathrm{E}-01$ |
| :--- |
| $-.20443 \mathrm{E}-01$ | $-.11890 \mathrm{E}-01$

$-.72207 \mathrm{E}-01$ $-.7229 \mathrm{E}-\mathrm{E}$
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$.72429 \mathrm{E}-01$ $.72429 \mathrm{E}-01$
$.40402 \mathrm{E}-01$
$.83755 \mathrm{E}-02$ $.83755 \mathrm{E}-02$
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$-.55678 \mathrm{E}-01$ $-.55678 \mathrm{E}-01$
$-.87705 \mathrm{E}-01$ $-.87705 \mathrm{E}-01$
-.11973 -.11973
.- .15176



 $\begin{array}{ll}-.36457 \mathrm{E}-02 & -.61941 \mathrm{E}-01 \\ .14940 \mathrm{E}-02 & -.1647 \mathrm{E}-02 \\ -.88409 \mathrm{E}-01 & .18257 \mathrm{E}-02 \\ -.77255 \mathrm{E}-01 & -.69040 \mathrm{E}-03 \\ .89476 \mathrm{E}-01 & -.14396 \mathrm{E}-01 \\ -.87212 \mathrm{E}-02 & -.18886 \mathrm{E} \\ -.02 \\ .34025 \mathrm{E}-01 & -.77051 \mathrm{E}-01 \\ -.35402 \mathrm{E}-02 & -.98877 \mathrm{E}-01 \\ -.45556 \mathrm{E}-04 & -.15020 \mathrm{E}-01 \\ .00000 & .00000 \\ .00000 & .00000 \\ .00000 & .00000 \\ .00000 & .00000 \\ .00000 & .0000 \\ .00000 & .00000 \\ .00000 & .00000 \\ .00000 & .00000 \\ .00000 & . .00000 \\ .12415 \mathrm{E}-01 & .52711 \mathrm{E}-01\end{array}$.69553E-02

| $.62057 \mathrm{E}-02$ |
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| $.10056-02$ |

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.12352E-01 . 18326E-02 00000 8 8 88 응 8 8 응 . 79278


## Program Listing

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    C FINITE DIFFERENCE PROGRAM FOR NONLINEAR REGRESSION SOLUTION
    C OF TWO-DIMENSIONAL, STEADY-STATE, GROUND-WATER FLOW PROBLEMS
    BY R. L. COOLEY, USGS, DENVER, COLO.
    IMPLICIT DOUBLE PRECISION (A-H,O-Z)
    DIMENSION TITLE(20),DX(30),DY(30),CX(500),CY(500),VL(500)
    1,QR(500),WELL(500),HR(500),HC(500),HCI (70) , BK(70), BL(70)
    2,BM(70),BN(70),HO(70),W(70),PAR(4),QBF(50),PLA(50),PLB(50)
    3, CXHR(50), CXHL(50),CYHT (50), CYHB (50),AU(5,250),AL(20,250)
    4,V(500),X(20,70),S(20,90),XV(500),XS(4),P(20),WP(20)
    5,A(20,20),B(50)
    DIMENSION JPOS(30),IZN(500),IBZN(50),IPRM(4, 20),IBNA(50)
    1,IBNB (50), IBPA(10),IBPB (10),IBHN (50),IHSN(100), KOBS (70)
    2,LN(4),NCBA (20),NCEA(20),NCBF (20),NCEF (20) ,NCBH(20)
    3,NCEH(20),ILOC(500),JLOC(500),IN(500),IC(5, 250)
    COMMON/INT/NIJ ,NEQ, ICR , ICR1 , IB1 , LH1 , ID , JD , IM , JM , NOBS , NQSD , NBH
    1,NVAR,NVX2,KOUNT, INDT, IPO
        COMMON/LOC/ILOC,JLOC
        COMMON/TNME/IIN, IOUT
        COMMON/FLT/CX,CY,VL,QR,WELL,HR ,HC , BR , BL , BM, BN ,HO,W
        COMMON/SOLV/AU,AL
        COMMON/REG/DMX,ADMX,AP,CSA,AMP ,RP,BP,YSQ
        EQUIVALENCE (TITLE(1),A(1,1),S(1,1),AU(1,1)),(HC(1),XV(1))
        1,(CX(1),HCI(1)),(ILOC(1),IC(1,1))
    OPEN (5,FILE='INVFD.DAT',STATUS='OLD',ACCESS='SEQUENTIAL'
    1,FORM=' FORMATTED')
    OPEN (6,FILE='INVFD.OUT',STATUS='NEW',ACCESS='SEQUENTIAL'
    1,FORM=' FORMATTED')
    OPEN (7,STATUS='NEW',ACCESS='SEQUENTIAL',FORM='UNFORMATTED')
C**DEFINE INPUT FILE, OUTPUT FILE, SCRATCH FILE, AND ARRAY DIMENSIONS
    IIN=5
    IOUT=6
    ITA=7
    NVD=20
    NAD=20
C**READ THREE TITLE LINES
    WRITE(IOUT, 804)
    DO 5 I=1,3
    READ(IIN, 801) (TITLE(J),J=1,20)
    SET A
    5 WRITE(IOUT, 803) (TITLE(J),J=1, 20)
C**READ JOB SPECIFICATION DATA
    READ(IIN, 800) ID,JD,NZNS ,NOBS,NPAR,NVAR,NWELS ,NQBZ,NHBZ,NUM, IPRX
    1,IPO,ISO
    SET B
    WRITE(IOUT, 802) ID,JD,NZNS,NOBS,NPAR,NVAR,NWELS,NQBZ,NHBZ,NUM
    1,IPRX,IPO,ISO
    READ(IIN, 820) DMX,CSA,RP,BP,EV
    SET C
    WRITE(IOUT,806) DMX,CSA,RP,BP,EV
C**READ INITIAL ARRAY DATA
    IM=ID-1
    JM=JD-1
    CALL ARRAY(DX,IM,1,0) SET D
```

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Program Listing-Continued
    CALL ARRAY(DY,JM,1,0)
```

SET E SET F SET G SET H SET I SET J SET K
$C * *$ READ GRID ZONATION
CALL ARRAYI (IZN, IM,JM,0)
SET L
C**READ OBSERVED HEAD DATA
IF (NOBS.LT.1) GO TO 15
CALL HOBS (DX, DY, V (1) , V (ID+1) , BK, BL , BM, BN , HO, W, KOBS)
SET M
C**INITIALIZE PARAMETER NUMBERS AND STANDARD DEVIATIONS TO ZERO
15 DO $25 \mathrm{I}=1$, NZNS
DO $20 \mathrm{~K}=1,4$
$20 \operatorname{IPRM}(\mathrm{~K}, \mathrm{I})=0$
25 CONTINUE
DO $27 \mathrm{I}=1$,NVAR
$27 \mathrm{WP}(\mathrm{I})=0$.
$C \star *$ READ AQUIFER PARAMETER NUMBERS
IF (NPAR.LT.1) GO TO 35
WRITE(IOUT, 816)
DO $30 \mathrm{~J}=1$,NZNS
$\operatorname{READ}(\operatorname{IIN}, 800) \mathrm{I},(\operatorname{IPRM}(\mathrm{K}, \mathrm{I}), \mathrm{K}=1,4)$
SET N
WRITE (IOUT, 818) I , (IPRM ( $\mathrm{K}, \mathrm{I}$ ) , $\mathrm{K}=1,4$ )
30 CONTINUE
C**READ AQUIFER PARAMETER STANDARD DEVIATIONS
WRITE (IOUT, 822)
DO $32 \mathrm{~J}=1$, NPAR
READ (IIN, 812) K,WP(K)
SET 0
32 WRITE(IOUT, 823) K,WP(K)
$C * *$ READ INITIAL AQUIFER PARAMETERS BY ZONE, AND LOAD THEM INTO THE
C PARAMETER VECTOR 35 WRITE(IOUT, 810)

M=NVAR
DO $45 \mathrm{~J}=1$, NZNS
READ (IIN , 812) I, PAR (1), $\operatorname{PAR}(2), \operatorname{PAR}(3), \operatorname{PAR}(4)$
SET P

```
WRITE(IOUT, 814) I,PAR(1),PAR(2),PAR(3),PAR(4)
DO 40 K=1,4
L=IPRM(K,I)
IF(L.GT.0) GO TO 40
M=M+1
L}=
IPRM(K,I)=L
    40 B(L)=PAR(K)
    4 5 \text { CONTINUE}
C**DEFINE JPOS ARRAY SUCH THAT COLUMN+JPOS(ROW)=NODE NUMBER
            JPOS (1)=0
            DO 50 J=2,JD
        50 JPOS(J)=JPOS (J - 1)+ID
C**READ POINT FLOW DATA
        NIJ=ID*JD
```

Program Listing-Continued
DO $55 \mathrm{~N}=1$, NIJ
$55 \operatorname{WELL}(\mathrm{~N})=0$.
IF(NWELS.LT.1) GO TO 61
WRITE(IOUT, 824)
DO $60 \mathrm{~K}=1$, NWELS
READ (IIN, 826) I, J, TMP
SET Q
WRITE (IOUT, 828) I,J,TMP
$\mathrm{L}=\mathrm{I}+\mathrm{J}$ POS ( J )
60 WELL $(\mathrm{L})=$ TMP
C**READ AND FORM ARRAYS FOR SPECIFIED POINT OR LINE FLOWS
61 NQSD=0
IF (NQBZ.LT.1) GO TO 85
WRITE(IOUT, 830)
$\mathrm{N}=0$
DO $80 \mathrm{~J}=1$, NQBZ
READ (IIN, 832) IA, JA, IB, JB, IP, QB , SDQB , QBM
SET R
WRITE (IOUT, 831) IA, JA, IB, JB, IP , QB , SDQB , QBM
$\mathrm{M}=1$
$K=I A-1$
IF(JA.EQ.JB) GO TO 62
$M=$ ID
$\mathrm{K}=\mathrm{JA}-1$
$62 \mathrm{MA}=\mathrm{IA}+\mathrm{JPOS}(\mathrm{JA})$
$M B=I B+J P O S(J B)-M$
IF(MB.GE.MA) GO TO 64
IF (IP.LT.1) GO TO 63
$\mathrm{N}=\mathrm{N}+1$
IBNA ( N ) $=\mathrm{MA}$
$\operatorname{IBNB}(\mathrm{N})=\mathrm{MA}$
$\mathrm{QBF}(\mathrm{N})=.5 * \mathrm{QBM}$
$\operatorname{IBZN}(N)=I P$
GO TO 68
63 WELL (MA) $=$ QB*QBM
GO TO 80
$64 \mathrm{QBM}=.5 * \mathrm{QBM}$
IF(IP.LT.1) GO TO 70
DO $66 \mathrm{~L}=\mathrm{MA}, \mathrm{MB}, \mathrm{M}$
$\mathrm{N}=\mathrm{N}+1$
IBNA (N) $=\mathrm{L}$
$\operatorname{IBNB}(\mathrm{N})=\mathrm{L}+\mathrm{M}$
$\mathrm{K}=\mathrm{K}+1$
TEMP=DX (K)
$I F(M . E Q . I D) T E M P=D Y(K)$
QBF ( N ) $=\mathrm{QBM} * T E M P$
$66 \operatorname{IBZN}(N)=I P$
$68 \mathrm{~B}(\mathrm{IP})=\mathrm{QB}$
$W P(I P)=S D Q B$
GO TO 80
$70 \mathrm{TMP}=\mathrm{QB} * \mathrm{QBM}$
DO $75 \mathrm{~L}=\mathrm{MA}, \mathrm{MB}, \mathrm{M}$
$\mathrm{K}=\mathrm{K}+1$
TEMP $=\mathrm{DX}(\mathrm{K})$

Program Listing-Continued

```
            IF(M.EQ.ID) TEMP=DY(K)
            TEMP=TMP*TEMP
            WELL(L)=WELL(L)+TEMP
    75 WELL(L+M)=WELL(L+M)+TEMP
    8 0 ~ C O N T I N U E ~
        NQSD=N
C**READ SPECIFIED BOUNDARY HEAD POSITIONS AS -1'S
    85 CALL ARRAYI(IN,ID,JD,0)
                                    SET S
C**READ DATA AND FORM ARRAYS FOR SPECIFIED HEADS AND PARAMETERS
    IF(NHBZ.LT.1) GO TO 110
    WRITE(IOUT, 833)
    NBH=0
    DO }108\textrm{KK=1,NHBZ
    READ(IIN, 834) IZ,NN,M,N,SDHA,SDHB
    SET T
    WRITE(IOUT,836) IZ,NN,M,N,SDHA,SDHB
    DO }95\textrm{J}=1,\textrm{NN
    READ(IIN,826) ILOC(J),JLOC(J),V(J) SET T
    95 WRITE(IOUT, 840) ILOC(J),JLOC(J),V(J)
        IBPA(IZ)=M
        IBPB(IZ)=N
        IF(M.LT.1) GO TO }9
        B(M)=V(1)
        WP(M)=SDHA
    97 IF(N.LT.1) GO TO 98
        B(N)=V(NN)
        WP(N)=SDHB
    98 J=JLOC(1)
        K=ILOC(1)+JPOS (J)
        M=M+N
        NBHS=NBH
        IF(IN(K).LT.-1) GO TO 100
        NBH=NBH+1
        IF(M.GT.0) IN(K)=-NBH-1
        IBZN(NBH+NQSD)=IZ
        IBHN (NBH)=K
        PLA(NBH)=1.
        PLB (NBH)=0.
        100 IF(NN.LT.2) GO TO 107
            DIST=0.
            DO 102 KNT=2,NN
            J=JLOC (KNT)
            L=ILOC(KNT)+JPOS(J)
            NBH=NBH+1
            IF(M.GT.0) IN(L)=-NBH-1
            IBZN}(NBH+NQSD)=IZ
            IBHN (NBH)=L
            JM1=JLOC(KNT-1)
            IF(J.EQ.JM1) GO TO 101
            J=MINO(J,JM1)
            DIST=DIST+DY(J)
            GO TO }10
    101 I=MINO(ILOC(KNT),ILOC(KNT-1))
```

Program Listing-Continued
DIST=DIST+DX(I)
$102 \mathrm{PLB}(\mathrm{NBH})=$ DIST
$\mathrm{N}=\mathrm{NBH}-\mathrm{NN}+1$
DO $106 \mathrm{KNT}=2$,NN
$\mathrm{J}=\mathrm{JLOC}(\mathrm{KNT})$
$\mathrm{L}=\mathrm{ILOC}$ (KNT) $+\mathrm{JPOS}(\mathrm{J})$
$\mathrm{N}=\mathrm{N}+1$
TMPA $=$ PLB ( N )/DIST
TMPB=1. - TMPA
TMPC=TMPA*V (NN) + TMPB*V(1)
IF(DABS (V(KNT)).LE.O.) GO TO 104
TMP=V(KNT)/TMPC
TMPA=TMPA*TMP
TMPB=TMPB $*$ TMP
TMPC=V(KNT)
$104 \operatorname{PLA}(N)=T M P B$
$\operatorname{PLB}(\mathrm{N})=\mathrm{TMPA}$
$106 \mathrm{HC}(\mathrm{L})=$ TMPC
$107 \mathrm{IF}(\mathrm{M} . \mathrm{LT} .1) \mathrm{NBH}=\mathrm{NBHS}$
$108 \mathrm{HC}(\mathrm{K})=\mathrm{V}(1)$
C**COMPARE CX AND CY WITH IZN FOR CONFLICT
110 IER=0
$\mathrm{N}=0$
DO $115 \mathrm{~J}=1, \mathrm{JM}$
DO $115 \mathrm{I}=1$, IM
$\mathrm{N}=\mathrm{N}+1$
$\operatorname{IF}(\operatorname{IZN}(\mathrm{N}) . \operatorname{LT} .1)$ Go To 115
IF(CX(N).GT.O..OR.CY(N).GT.O.) GO TO 115
IER=1
WRITE(IOUT, 842) I,J
115 CONTINUE
IF(IER.LT.1) GO TO 120
WRITE(IOUT, 844)
STOP
C**TRANSFER DOMAIN GEOMETRY TO IN(M) AND COMPUTE CELL FLOW-COEFFICIENTS
$120 \mathrm{~N}=0$
DO $122 \mathrm{~J}=1$, JM
$\mathrm{DYN}=.5 * \mathrm{DY}(\mathrm{J})$
DO $122 \mathrm{I}=1, \mathrm{IM}$
$\mathrm{N}=\mathrm{N}+1$
IF(IZN(N).LT.1) GO TO 122
$\mathrm{M}=\mathrm{N}+\mathrm{J}$
$\operatorname{IF}(\operatorname{IN}(M) . \operatorname{GT} .-1) \quad \operatorname{IN}(M)=1$
$\operatorname{IF}(\operatorname{IN}(M-1) . G T .-1) \quad \operatorname{IN}(M-1)=1$
$\operatorname{IF}(\mathrm{IN}(\mathrm{M}+\mathrm{ID}-1) . \mathrm{GT} .-1) \quad \mathrm{IN}(\mathrm{M}+\mathrm{ID}-1)=1$
IF (IN(M+ID).GT.-1) IN(M+ID)=1
$C X(N)=C X(N) * D Y N / D X(I)$
DXN=. $5 * \mathrm{DX}(\mathrm{I})$
CY (N) $-\mathrm{CY}(\mathrm{N}) * \mathrm{DXN} / \mathrm{DY}(\mathrm{J})$
AREA $=D X N * D Y N$
$\mathrm{VL}(\mathrm{N})=\mathrm{VL}(\mathrm{N}) *$ AREA
$\mathrm{QR}(\mathrm{N})=\mathrm{QR}(\mathrm{N}) *$ AREA

Program Listing-Continued
122 CONTINUE
C**SET UP D4 ORDERING
CALL ORDER(JPOS,IN,IC)
C**COMPUTE INITIAL SOLUTION
CALL COEF (WELL, HR, HC, CX, CY, VL, QR, CXHR, CXHL, CYHT, CYHB , QBF , B, AU , AL 1, V, IZN, IBZN, IPRM, IBNA, IBNB, IN, IC, NAD)
CALL D4SOLV(HC,AU, AL, V, IN, IC, NAD)
WRITE(IOUT, 846)
CALL ARRAY(HC, ID, JD, 1)
IF(ISO.EQ.1) GO TO 640
C**COMPUTE AND COUNT PRIOR INFORMATION DATA
NPRIR=0
DO 137 I=1, NVAR
$\mathrm{P}(\mathrm{I})=\mathrm{B}(\mathrm{I})$
IF(WP(I).LE.O.) GO TO 137
$\mathrm{WP}(\mathrm{I})=\mathrm{EV} /(\mathrm{WP}(\mathrm{I}) * \mathrm{WP}(\mathrm{I}))$
NPRIR=NPRIR +1
137 CONTINUE
WRITE(IOUT, 848) NPRIR
C**INITIALIZE BEGINNING AND END POINT ARRAYS
DO $148 \mathrm{I}=1$,NVAR
$\operatorname{NCBA}(\mathrm{I})=0$
$\operatorname{NCEA}(\mathrm{I})=0$
$\operatorname{NCBF}(\mathrm{I})=0$
$\operatorname{NCEF}(\mathrm{I})=0$
$\operatorname{NCBH}(\mathrm{I})=0$
$148 \operatorname{NCEH}(\mathrm{I})=0$
C**DEFINE BEGINNING AND END POINT ARRAYS FOR AQUIFER PARAMETERS
IF (NPAR.LT.1) GO TO 154
$\mathrm{N}=0$
DO $152 \mathrm{~J}=1$, JM
DO $152 \mathrm{I}=1, \mathrm{IM}$
$\mathrm{N}=\mathrm{N}+1$
$\mathrm{L}=\mathrm{I}$ ZN(N)
IF(L.LT.1) GO TO 152
DO $150 \mathrm{M}=1,4$
$K=I P R M(M, L)$
IF(K.GT.NVAR) GO TO 150
$\operatorname{NCEA}(\mathrm{K})=\mathrm{N}$
$\operatorname{IF}(\operatorname{NCBA}(\mathrm{K}) . \operatorname{LT} .1) \operatorname{NCBA}(\mathrm{K})=\mathrm{N}$
150 CONTINUE
152 CONTINUE
C**ORDER IBZN AND CORRESPONDING ARRAYS FOR LINE FLOW PARAMETERS
C FROM SMALLEST TO LARGEST
154 IF(NQSD.LT.1) GO TO 162
DO $158 \mathrm{I}=1$,NQSD
DO $156 \mathrm{~J}=\mathrm{I}$, NQSD
IF(IBZN(J).GE.IBZN(I)) GO TO 156
ITMP $=$ IBZN (I)
$\operatorname{IBZN}(\mathrm{I})=\operatorname{IBZN}(\mathrm{J})$
$\operatorname{IBZN}(\mathrm{J})=I T M P$
ITMP $=$ IBNA ( I )

Program Listing-Continued

$$
\operatorname{IBNA}(\mathrm{I})=\operatorname{IBNA}(\mathrm{J})
$$

IBNA( $J$ ) $=$ ITMP
ITMP=IBNB (I)
$\operatorname{IBNB}(I)=\operatorname{IBNB}(J)$
$\operatorname{IBNB}(\mathrm{J})=I T M P$
TMP=QBF(I)
QBF (I) $=\mathrm{QBF}(\mathrm{J})$
QBF(J) $=$ TMP
156 CONTINUE
158 CONTINUE
C**DEFINE BEGINNING AND END POINT ARRAYS FOR LINE FLOW PARAMETERS DO $160 \mathrm{I}=1$,NQSD
$K=I B Z N$ (I)
NCEF (K) $=1$
IF (NCBF(K).LT.1) $\operatorname{NCBF}(K)=I$
160 CONTINUE
$C * * D E F I N E$ SEQUENCE NUMBERS, AND BEGINNING AND END POINT ARRAYS
C FOR SPECIFIED HEAD PARAMETERS
162 IF (NBH.LT.1) GO TO 174
DO $164 \mathrm{I}=1$, NBH
$J=I B Z N(I+N Q S D)$
$\mathrm{K}=\mathrm{IBPA}(\mathrm{J})$
IF (K.EQ.0) GO TO 163
NCEH (K) $=1$
$\mathrm{IF}(\mathrm{NCBH}(\mathrm{K}) . \operatorname{LT} .1) \mathrm{NCBH}(\mathrm{K})=\mathrm{I}$
3
$163 \mathrm{~K}=\mathrm{IBPB}(\mathrm{J})$
IF(K.EQ.0) GO TO 164
NCEH (K) - I
$\operatorname{IF}(\mathrm{NCBH}(\mathrm{K}) . \operatorname{LT} .1) \mathrm{NCBH}(\mathrm{K})=\mathrm{I}$
164 CONTINUE
$\mathrm{L}=0$
DO $172 \mathrm{~K}=1$, NVAR
IF(NCBH(K).LT.1) GO TO 172
$J B=N C B H(K)$
$\mathrm{JE}=\mathrm{NCEH}$ (K)
NCBH (K) $=\mathrm{L}+1$
DO $170 \mathrm{~J}=\mathrm{JB}, \mathrm{JE}$
$I=I B Z N(J+N Q S D)$
IF (IBPA(I).EQ.K.OR.IBPB(I).EQ.K) GO TO 168
GO TO 170
$168 \mathrm{~L}=\mathrm{L}+1$
IHSN (L) $=\mathrm{J}$
170 CONTINUE
NCEH (K) $=\mathrm{L}$
172 CONTINUE
C**BEGIN ITERATIONS
174 INDT=0
ER=. 01
ERP=. 001
$\mathrm{AMP}=0$.
NVX2 $=$ NVAR+NVAR
KOUNT $=0$

## Program Listing-Continued

176 KOUNT=KOUNT+1
REWIND ITA
C**SOLVE FOR SENSITIVITIES:
DO $260 \mathrm{~N}=1$, NVAR
DO $178 \mathrm{I}=1$,NEQ
$178 \mathrm{~V}(\mathrm{I})=0$.
IF(NCBA(N).LT.1) GO TO 208
C**ASSEMBLE R.H.S. FOR AQUIFER PARAMETERS
$\mathrm{LB}=\mathrm{NCBA}(\mathrm{N})$
LE=NGEA(N)
DO $200 \mathrm{~L}=\mathrm{LB}, \mathrm{LE}$
$\mathrm{J}=\mathrm{IZN}$ (L)
IF(J.LT.1) GO TO 200
$N A=L+(L-1) / I M$
$N B=N A+1$
$N C=N B+I D$
$N D=N A+I D$
$\operatorname{INA}=I N$ (NA)
$\operatorname{INB}=\mathrm{IN}(\mathrm{NB})$
$\operatorname{INC}=\mathrm{IN}(\mathrm{NC})$
IND=IN(ND)
IF(IPRM(1,J).NE.N) GO TO 180
IF (INA.GT.0) V(INA) $=\mathrm{V}($ INA $)+\mathrm{CX}(\mathrm{L}) *(\mathrm{HC}(\mathrm{NB})-\mathrm{HC}(\mathrm{NA}))$
IF (INB.GT.0) V(INB) $=\mathrm{V}($ INB $)+C X(L) *(H C(N A)-H C(N B))$
IF (INC.GT.0) $V($ INC $)=V(I N C)+C X(L) *(H C(N D)-H C(N C))$
IF (IND.GT.0) V(IND) $=\mathrm{V}($ IND $)+\mathrm{CX}(\mathrm{L}) *(\mathrm{HC}(N \mathrm{C})-\mathrm{HC}(\mathrm{ND}))$
$180 \operatorname{IF}(\operatorname{IPRM}(2, J) . N E . N)$ GO TO 185
IF (INA.GT.0) $V($ INA $)=V($ INA $)+C Y(L) *(H C(N D)-H C(N A))$
IF (INB.GT.0) $V($ INB $)=V($ INB $)+C Y(L) *(H C(N C)-H C(N B))$
IF (INC.GT.0) V(INC) $=\mathrm{V}(\mathrm{INC})+\mathrm{CY}(\mathrm{L}) *(\mathrm{HC}(\mathrm{NB})-\mathrm{HC}(\mathrm{NC}))$
IF(IND.GT.0) V(IND) $=\mathrm{V}($ IND $)+\mathrm{CY}(\mathrm{L}) *(\mathrm{HC}(N A)-\mathrm{HC}(\mathrm{ND}))$
$185 \operatorname{IF}(\operatorname{IPRM}(3, J) . N E . N)$ GO TO 190
IF(INA.GT.0) V(INA) $=\mathrm{V}($ INA $)+V L(L) *(H R(N A)-H C(N A))$
IF (INB.GT.0) V(INB) $=\mathrm{V}(\mathrm{INB})+\mathrm{VL}(\mathrm{L}) *(\mathrm{HR}(\mathrm{NB})-\mathrm{HC}(\mathrm{NB}))$
$\mathrm{IF}(\mathrm{INC.GT.0}) \mathrm{V}(\mathrm{INC})=\mathrm{V}(\mathrm{INC})+\mathrm{VL}(\mathrm{L}) *(\mathrm{HR}(\mathrm{NC})-\mathrm{HC}(\mathrm{NC}))$
IF(IND.GT.0) V(IND)=V(IND)+VL(L)*(HR(ND)-HC(ND))
$190 \operatorname{IF}(\operatorname{IPRM}(4, \mathrm{~J}) . N E . N)$ GO TO 200
$\mathrm{IF}($ INA. GT.0) $\quad V(I N A)=V(I N A)+Q R(L)$
$I F(I N B . G T .0) V(I N B)=V(I N B)+Q R(L)$
$\operatorname{IF}(I N C . G T .0) \quad V(I N C)=V(I N C)+Q R(L)$
IF (IND.GT.0) V(IND) $=V($ IND $)+Q R(L)$
200 CONTINUE
$C * *$ ASSEMBLE R.H.S. FOR SPECIFIED LINE FLOW PARAMETERS
208 IF (NCBF(N).LT.1) GO TO 212
$L B=\operatorname{NCBF}(\mathrm{N})$
LE=NCEF(N)
DO $210 \mathrm{~L}=\mathrm{LB}, \mathrm{LE}$
$I=I B N A(L)$
$J=I N(I)$
$\operatorname{IF}(\mathrm{J} . \mathrm{GT} .0) \mathrm{V}(\mathrm{J})=\mathrm{V}(\mathrm{J})+\mathrm{QBF}(\mathrm{L})$
$\mathrm{I}=\mathrm{IBNB}(\mathrm{L})$
$\mathrm{J}=\mathrm{IN}$ (I)
$212 \operatorname{IF}(\operatorname{NCBH}(N) . L T .1)$ GO TO 216
$\mathrm{LB}=\mathrm{NCBH}(\mathrm{N})$
$L E=N C E H(N)$
DO 214 L-LB,LE
$\mathrm{K}=\mathrm{IHSN}(\mathrm{L})$
$I=I B Z N(K+N Q S D)$
TMP=0.
IF (IBPA (I).EQ.N) TMP=PLA(K)
$\operatorname{IF}(\operatorname{IBPB}(I) . E Q \cdot N) T M P=T M P+P L B(K)$
$\mathrm{I}=\mathrm{IBHN}(\mathrm{K})+1$
IF (I.LE.NIJ) J=IN(I)
IF (J.GT.0) $V(J)=V(J)+C X H R(K) * T M P$
$I=I B H N(K)-1$
IF (I.GT.0) J=IN(I)
IF (J.GT.0) $\mathrm{V}(\mathrm{J})=\mathrm{V}(\mathrm{J})+$ CXHL $(\mathrm{K}) * T M P$
$I=I B H N(K)+I D$
IF(I.LE.NIJ) J=IN(I)
IF (J.GT.0) V(J) $=V(\mathrm{~J})+$ CYHT (K) *TMP
I=IBHN(K) - ID
IF(I.GT.0) J=IN(I)
IF (J.GT.0) $\mathrm{V}(\mathrm{J})=\mathrm{V}(\mathrm{J})+\mathrm{CYHB}(\mathrm{K}) * T M P$
214 CONTINUE

## 1

C**MODIFY R.H.S.--UPPER HALF
216 DO $220 \mathrm{~J}=1$, ICRI
$I I=I C(1, J)$
DO $218 \mathrm{I}=2$,II
$\mathrm{LR}=\mathrm{IC}(\mathrm{I}, \mathrm{J})$
$V(L R)=V(L R)-A U(I, J) * V(J)$
218 CONTINUE
$220 \mathrm{~V}(\mathrm{~J})=\mathrm{V}(\mathrm{J}) / \mathrm{AU}(1, \mathrm{~J})$
C**MODIFY R.H.S. - -LOWER HALF
$\mathrm{JJ}=\mathrm{NEQ}-\mathrm{ICR}$
DO $224 \mathrm{~J}=1$, JJ
$J R=J+I C R 1$
LR=JR
DO $222 I=2$,IB1
LR=LR+1
$\mathrm{IF}(\mathrm{AL}(\mathrm{I}, \mathrm{J}) . \mathrm{NE} .0 \mathrm{O}) \mathrm{V}(\mathrm{LR})=\mathrm{V}(\mathrm{LR})-\mathrm{AL}(\mathrm{I}, \mathrm{J}) * \mathrm{~V}(\mathrm{JR})$
222 CONTINUE
$224 \mathrm{~V}(\mathrm{JR})=\mathrm{V}(\mathrm{JR}) / \mathrm{AL}(1, \mathrm{~J})$
$C * * B A C K$ SOLVE--LOWER HALF
$V(N E Q)=V(N E Q) / A L(1, N E Q-I C R 1)$
DO $230 \mathrm{~J}=1$, JJ
KK=NEQ-J
KL=KK-ICR1
LmKK
DO $226 \mathrm{I}=2$, IB1
$\mathrm{L}=\mathrm{L}+1$
$\mathrm{IF}(\mathrm{AL}(\mathrm{I}, \mathrm{KL}) . \mathrm{NE} .0) .\mathrm{V}(\mathrm{KK})=\mathrm{V}(\mathrm{KK})-\mathrm{AL}(\mathrm{I}, \mathrm{KL}) * \mathrm{~V}(\mathrm{~L})$

Program Listing-Continued

## 226 CONTINUE

230 CONTINUE
C**BACK SOLVE--UPPER HALF
DO 250 J-1,ICR1
KK=ICR-J
II-IC ( $1, \mathrm{KK}$ )
DO $240 \mathrm{I}=2$, II
$\mathrm{L}=\mathrm{IC}(\mathrm{I}, \mathrm{KK})$
$\mathrm{V}(\mathrm{KK})=\mathrm{V}(\mathrm{KK})-\mathrm{AU}(\mathrm{I}, \mathrm{KK}) * \mathrm{~V}(\mathrm{~L})$
240 CONTINUE
250 CONTINUE
WRITE(ITA) (V(I), I=1,NEQ)
C**COMPUTE SENSITIVITIES AT OBSERVATION POINTS
DO 255 I-1, NOBS
k-Kobs (I)
$\mathrm{LN}(1)=\mathrm{IN}(\mathrm{K})$
$\operatorname{LN}(2)=\operatorname{IN}(\mathrm{K}+1)$
$\mathrm{LN}(3)=\mathrm{IN}(\mathrm{K}+\mathrm{ID})$
$\mathrm{LN}(4)=\mathrm{IN}(\mathrm{K}+\mathrm{ID}+1)$
DO $253 \mathrm{~J}=1,4$
$\mathrm{L}=\mathrm{LN}(\mathrm{J})$
IF(L.GT.0) GO TO 252
Xs(J)-0.
IF(L.GT.-2) GO TO 253
L=-L-1
$I Z=I B Z N(L+N Q S D)$
TMP $=0$.
IF(IBPA(IZ).EQ.N) TMP=PLA(L)
IF (IBPB(IZ).EQ.N) TMP=TMP+PLB(L)
$\mathrm{XS}(\mathrm{J})=\mathrm{TMP}$
GO то 253
$252 \mathrm{XS}(\mathrm{J})=\mathrm{V}(\mathrm{L})$
253 CONTINUE
$255 \mathrm{X}(\mathrm{N}, \mathrm{I})=\mathrm{BK}(\mathrm{I}) * \mathrm{XS}(1)+\mathrm{BL}(\mathrm{I}) * \mathrm{XS}(2)+\mathrm{BM}(\mathrm{I}) * \mathrm{XS}(3)+\mathrm{BN}(\mathrm{I}) * \mathrm{XS}(4)$
260 CONTINUE
IF(IPO.NE.1) GO TO 270
WRITE(IOUT, 850)
DO $265 \mathrm{~N}=1$, NOBS
WRITE (IOUT, 852) $\mathrm{N},(\mathrm{X}(\mathrm{K}, \mathrm{N}), \mathrm{K}=1$, NVAR $)$
265 CONTINUE
C**CALL LEAST SQUARES
270 CALL LSTSQ(HC , BK , BL , BM, BN , HO , W, P, WP , X , A , B , V, KOBS , IN , NVD)
IF(INDT.GT.0) GO TO 515
C**COMPUTE NEW SPECIFIED HEADS
IF (NBH.LT.1) GO TO 310
DO $300 \mathrm{~N}=1$, NBH
$M=I B Z N(N+N Q S D)$
$\mathrm{K}=\mathrm{IBPA}(\mathrm{M})$
TMPA=0.
$\operatorname{IF}(\mathrm{K} . \mathrm{GT} .0) \quad$ TMPA $=\mathrm{PLA}(\mathrm{N}) * \mathrm{~V}(\mathrm{~K})$
$\operatorname{L=IBPB}(M)$
TMPB=0.

```
    IF(L.GT.0) TMPB=PLB(N)*V(L)
        J=IBHN(N)
        HC (J) =HC(J)+TMPA+TMPB
        300 CONTINUE
C**CHECK FOR CONVERGENCE
    310 IF(ADMX.LT.ER) GO TO 350
C**CHECK FOR PARAMETERS GOING TO ZERO
    IND=0
    DO 335 I=1,NVAR
    IF(DABS(B(I)).GT.DABS(ERP*P(I))) GO TO 335
    WRITE(IOUT,858) I
    IND=1
    335 CONTINUE
        IF(IND.GT.0) GO TO 515
    IF(KOUNT.EQ.NUM) GO TO 340
C**COMPUTE NEW HEADS AT GRID POINTS
    CALL COEF(WELL, HR,HC, CX, CY,VL,QR,CXHR, CXHL, CYHT, CYHB ,QBF, B,AU,AL
    1,V,IZN,IBZN, IPRM, IBNA, IBNB,IN, IC,NAD)
    CALL D4SOLV(HC,AU,AL,V,IN,IC,NAD)
    GO TO }17
    340 WRITE(IOUT,860) NUM
    GO TO 515
    350 WRITE(IOUT, 862) KOUNT
C**COMPUTE FINAL ESTIMATES OF HEAD
    REWIND ITA
    DO 366 K=1,NVAR
    READ(ITA) (CX(I),I=1,NEQ)
    DO 364 J=1,NIJ
    L=IN(J)
    IF(L.GT.0) HC(J)=HC(J)+CX(L)*V(K)
    364 CONTINUE
    366 CONTINUE
C**CORRECT A FOR MARQUARDT PARAMETER
    IF(NVAR.EQ.1) GO TO 420
    IF(AMP.LE.O.) GO TO }38
    DO 380 I=1,NVAR
    A(I,I)=1.+RP
    DO 375 JmI,NVAR
    375 A(J,I)=A(I,J)
    380 CONTINUE
    AMP=-1.
    CALL LSTSQ(HC,BK,BL, BM, BN,HO,W, P,WP,X,A,B,V,KOBS,IN,NVD)
    IF(INDT.GT.0) GO TO 515
C**COMPUTE A-INVERSE
    385 A(NVAR,NVAR)=1./A(NVAR,NVAR)
        NM1=NVAR-1
        DO 410 K-1,NM1
        KP1=K+1
        DO 395 I=KP1,NVAR
        SUM=0 .
        IM1=I-1
        DO 390 J=K,IM1
```

Program Listing-Continued

```
    390 SUM=SUM+A(I,J)*A(J,K)
        A(K,I)=-SUM
    395 A(I,K)=-SUM*A(I, I)
        DO 405 I=1,K
        SUM=A(K,I)
        DO 400 J=KP1,NVAR
    400 SUM=SUM+A(J,K)*A(I,J)
        A(K,I)=SUM
    405 A(I, K)=A(K,I)
    410 CONTINUE
        DO 415 J=1,NVAR
    415 A(J,NVAR)=A(NVAR,J)
        GO TO 425
    420 A(1,1)=1./(1.+RP)
C**COMPUTE TR((A-INVERSE)**2) AND A-INVERSE - RP*(A-INVERSE)**2
    425 TRACE=0.
        IF(RP.LE.O.) GO TO 448
        DO 445 N=1,NVAR
        DO 430 J=1,NVAR
    430 V(J)=A(J,N)
        SUMA=0.
        DO 440 J=N,NVAR
        SUM=0.
        DO 435 I=1,NVAR
    4 3 5 ~ S U M = S U M + V ( I ) * A ( I , J )
        V(J+NVAR)=SUM
    440 A(J,N)=A(J,N)-RP*SUM
    445 TRACE=TRACE+V (N+NVAR)
C**COMPUTE SUM OF SQUARED ERRORS
    448 YSQ=0.
        DO 450 N=1,NOBS
        K=KOBS (N)
        HCI (N)=BK(N)*HC(K)+BL(N)*HC(K+1)+BM(N)*HC (K+ID) +BN(N)*HC(K+ID+1)
    450 YSQ = YSQ+(HO(N)-HCI(N))*W(N)*(HO(N)-HCI(N))
        DO 455 I=1,NVAR
    455 YSQ =YSQ+(P(I)-B(I))*WP(I)*(P(I)-B(I))
C**COMPUTE ERROR VARIANCE
        TEMP=NPRIR-NVAR
        OBS=NOBS
        VAR=YSQ/(OBS+TEMP+RP*RP*TRACE)
C**COMPUTE CORRELATION COEFFICIENT
        SUMA=0.
        SUMB=0.
        SUMC=0.
        SUMD=0.
        SUM=0.
        DO 460 N=1,NOBS
        TMP=W(N)**. }
        W(N)=TMP
        TEMP=TMP*HO(N)
        TMP=TMP*HCI (N)
        SUMA=SUMA+TEMP
```

```
3
3
Program Listing-Continued
    SUMB=SUMB+TMP
    SUMC=SUMC+TEMP*TEMP
    SUMD=SUMD+TMP*TMP
    SUM=SUM+TEMP*TMP
    4 6 0 \text { CONTINUE}
    R=(OBS*SUM-SUMA*SUMB)/( (OBS*SUMC-SUMA*SUMA)*(OBS*SUMD -SUMB*SUMB))
    1**.5
C**PRINT ERROR VARIANCE, ESTIMATED SUM OF SQUARED ERRORS, AND
C CORRELATION COEFFICIENT
    WRITE(IOUT,864) VAR,YSQ,R
C**COMPUTE VARIANCE-COVARIANCE MATRIX
    DO 463 J=1,NVAR
    TEMP=V(J+NVX2)
    DO 462 I=J,NVAR
    A(I,J)=VAR*A(I,J)/(V(I+NVX2)*TEMP)
    462 A(J,I)=A(I,J)
    463 V(J)=A(J,J)**. }
C**PRINT PARAMETERS AND STANDARD ERRORS
    WRITE(IOUT, 870)
    DO 480 J=1,NVAR
    480 WRITE(IOUT, 856) J, B(J),V(J)
C**PRINT VARIANCE-COVARIANCE MATRIX
    490 WRITE(IOUT, 874)
    CALL PRTOT(A,NVAR,NVD,0)
C**COMPUTE AND PRINT CORRELATION MATRIX
    DO 510 J=1,NVAR
    TEMP=V(J)
    DO 500 I=J,NVAR
    A(I,J)=A(I,J)/(V(I)*TEMP)
    500 A(J,I)=A(I,J)
    510 CONTINUE
    WRITE(IOUT, 876)
    CALL PRTOT(A,NVAR,NVD,0)
C**PRINT COMPUTED AND OBSERVED HEADS, AND COMPUTE AND PRINT RESIDUALS
    GO TO 518
    5 1 5 \text { DO 516 N-1,NOBS}
    K-KOBS(N)
    HCI (N)=BK(N)*HC(K)+BL(N)*HC(K+1)+BM(N)*HC(K+ID)+BN(N)*HC(K+ID+1)
    516 W(N)=W(N)**. }
    518 WRITE(IOUT,878)
    DO 520 N=1,NOBS
    RES=W(N)*(HCI (N)-HO(N))
    WRITE(IOUT,880) N,HCI(N),HO(N),RES
    520 CONTINUE
C**PRINT HYDRAULIC HEADS FOR EACH NODE
    WRITE(IOUT, 881)
    CALL ARRAY(HC,ID,JD,1)
C**PRINT SENSITIVITIES FOR EACH NODE
    IF(IPRX.LT.1.AND.KOUNT.LT.NUM) STOP
    WRITE(IOUT,882)
    REWIND ITA
    DO 530 KK=1,NVAR
```


## Program Listing-Continued

```
    READ(ITA) (CX(I),I-1,NEQ)
    WRITE(IOUT,884) KK
    DO 525 N=1,NIJ
    L=IN(N)
    IF(L.GT.0) GO TO 523
    XV(N)=0.
    IF(L.GT.-2) GO TO 525
    L=-L-1
    IZ=IBZN(L+NQSD)
    TMP=0.
    IF(IBPA(IZ).EQ.KK) TMP=PLA(L)
    IF(IBPB(IZ).EQ.KK) TMP=TMP+PLB(L)
    XV(N)=TMP
    GO TO 525
    523 XV(N)=CX(L)
    525 CONTINUE
    530 CALL ARRAY(XV,ID,JD,1)
        IF(NVAR.LT.2) STOP
C**SCALE AND ORTHOGONALIZE COLUMNS OF SENSITIVITY MATRIX, X:
C**SCALE X AND AUGMENT X TO INCLUDE PRIOR
    DO 535 N=1,NOBS
    DO 535 K=1,NVAR
    535 S(K,N)=X(K,N)*W(N)/V(K+NVX2)
        IF(NPRIR.LT.1) GO TO 539
        N=NOBS
        DO 538 I=1,NVAR
        IF(WP(I).LT.1.E-10) GO TO 538
        N=N+1
        DO 537 J=1,NVAR
    537 S(J,N)=0.
    S(I,N)=WP(I)**.5/V(I+NVX2)
    538 CONTINUE
C**ORTHOGONALIZE S
    539 NTMP=NOBS+NPRIR
        DO 540 I=1,NTMP
    540 CY(I)=S(1,I)
        DO 600 N=2,NVAR
        NM1=N-1
        SUM=0.
        DO 550 I=1,NTMP
        SUM=SUM+CY(I)*CY(I)
        S(NM1,I)=CY(I)
    550 CONTINUE
        IF(SUM.LT.1.E-20) GO TO 610
        V(NM1)=1./SUM
        DO 570 J=1,NM1
        SUM=0.
        DO 560 K=1,NTMP
    560 SUM=SUM+V(J)*S(J,K)*S(N,K)
    570 CX(J)=SUM
        DO 590 K=1,NTMP
        SUM=0.
```

```
3
Program Listing-Continued
    DO 580 I=1,NM1
    50 SUM=SUM+S(I,K)*CX(I)
    590 CY(K)=S (N,K)-SUM
    6 0 0 ~ C O N T I N U E ~
C**PRINT ORTHOGONALIZED S
    610 WRITE(IOUT,886)
        K=1
        L=8
        DO 630 M=1,NVAR,8
        IF(L.GT.NVAR) L=NVAR
        WRITE(IOUT, 888) (I, I=K,L)
        DO 620 J=1,NTMP
        S(NVAR,J)=CY(J)
        WRITE(IOUT,890) J,(S (I,J),I=K,L)
    6 2 0 ~ C O N T I N U E ~
    WRITE(IOUT , 890)
    K=K+8
    L-L+8
    6 3 0 \text { CONTINUE}
    STOP
C**READ, PRINT, AND EXECUTE FOR ALTERNATE SOLUTIONS
    640 READ(IIN,800) N
        IF(N.LT.1) STOP
        DO }690\textrm{KNT=1,N
        WRITE(IOUT,891) KNT
        DO }670\mathrm{ L=1,NVAR
        READ(IIN,812) I,PR SET V
        V(I)=PR-B(I)
    670 B(I)=PR
    WRITE(IOUT, 892)
    CALL PRTOT(B,NVAR,0,1)
    IF(NBH.LT.1) GO TO 685
    DO 680 N=1,NBH
    M=IBZN(N+NQSD)
    K=IBPA(M)
    TMPA=0.
        IF(K.GT.0) TMPA=PLA(N)*V(K)
        L=IBPB(M)
        TMPB=0.
        IF(L.GT.0) TMPB=PLB(N)*V(L)
        J=IBHN(N)
    680 HC(J)=HC(J)+TMPA+TMPB
    685 CALL COEF(WELL,HR,HC,CX,CY,VL,QR,CXHR,CXHL, CYHT , CYHB, QBF, B, AU ,AL
        1,V,IZN, IBZN, IPRM, IBNA, IBNB, IN ,IC,NAD)
            CALL D4SOLV(HC,AU,AL,V,IN,IC,NAD)
            WRITE(IOUT,898)
            CALL ARRAY(HC,ID,JD,1)
    6 9 0 \text { CONTINUE}
            STOP
C
    800 FORMAT (16I5)
    801 FORMAT (20A4)
```


## Program Listing-Continued

802 FORMAT (56HONUMBER OF COLUMNS (ID)
$\$=$, I5
\$/56H NUMBER OF ROWS (JD)
, I5
\$/56H NUMBER OF AQUIFER ZONES (NZNS) …..................................... I5
\$/56H NUMBER OF OBSERVATIONS (NOBS) -...................................... I5
\$/56H NUMBER OF AQUIFER PARAMETERS (NPAR) ........................... I5
\$/56H TOTAL NUMBER OF PARAMETERS (NVAR) -................................ I5
\$/56H NUMBER OF KNOWN POINT FLOWS (NWELS) ….......................... I5



$\$ / 56 \mathrm{H}$ SENSITIVITY PRINT AND ORTHOGONALIZATION OPTION (IPRX) $=$, I5


803 FORMAT (1H ,20A4)
804 FORMAT (1H1)
806 FORMAT (50H MAXIMUM ALLOWABLE PARAMETER CORRECTION (DMX) - =
\$,G11. 5
\$/50H SEARCH DIRECTION ADJUSTMENT PARAMETER (CSA) $-=$, G11. 5
\$/50H RIDGE PARAMETER FOR REGRESSION (RP) $-\cdots-.-.=$, G 11.5
\$/50H BIAS PARAMETER FOR REGRESSION (BP) $\ldots-\ldots$....... $=$, G11. 5
\$/50H ESTIMATED ERROR VARIANCE (EV) .-..................... $=$, G11.5)
810 FORMAT (1H0,12X,34HINITIAL AQUIFER PARAMETERS BY ZONE/6H ZONE
$1,5 \mathrm{X}, 5$ HTRANX, $8 \mathrm{X}, 5 \mathrm{HTRANY}, 8 \mathrm{X}, 5$ HVLEAK , $8 \mathrm{X}, 5 \mathrm{HQDIST}$ )
812 FORMAT (I5,4F10.0)
814 FORMAT (1H , I4, 2X,4(2X,G11.5))
816 FORMAT (1H0,11X,25HAQUIFER PARAMETER NUMBERS $/ 1 \mathrm{H}, 5 \mathrm{X}, 4 \mathrm{HZONE}, 4 \mathrm{X}$
1, 5HTRANX, 3X, 5HTRANY, 3X, 5HVLEAK, 3X, 5HQDIST)
818 FORMAT (1H , 8I8)
820 FORMAT ( 8 F 10.0 )
822 FORMAT ( $1 \mathrm{HO}, 4 \mathrm{X}, 19 \mathrm{HSTANDARD}$ DEVIATIONS $/ 1 \mathrm{H}, 3 \mathrm{X}, 22 \mathrm{HFOR}$ AQUIFER PARAME
1TERS/1H ,6X,4HPAR.,6X,4HSTD./1H ,7X,3HNO.,6X,4HDEV.)
823 FORMAT (1H ,5X,I4,4X,G11.5)
824 FORMAT (1HO,12X,11HPOINT FLOWS $/ 1 \mathrm{H}, 7 \mathrm{X}, 1 \mathrm{HI}, 7 \mathrm{X}, 1 \mathrm{HJ}, 4 \mathrm{X}, 9 \mathrm{HVOL} . \operatorname{RATE})$
826 FORMAT (2I5,F10.0)
828 FORMAT (1H , 218, 4X,G11.5)
830 FORMAT (1H0,22X,27HINITIAL SPECIFIED FLOW DATA/1H ,6X,9HNODE NO.S
$1,7 \mathrm{X}, 4 \mathrm{HPAR} ., 6 \mathrm{X}, 4 \mathrm{HFLOW}, 10 \mathrm{X}, 4 \mathrm{HSTD} . / 1 \mathrm{H}, 19 \mathrm{H}$ IA JA IB JB,4X
2, 3HNO. , 4X, 9HPARAMETER, 7X,4HDEV., $5 \mathrm{X}, 10 \mathrm{HMULTIPLIER)}$
831 FORMAT (1H ,4(1X,I3,1X),2X,I3,3X,3(2X,G11.5))
832 FORMAT (5I5, 3F10.0)
833 FORMAT (1HO,18X,27HINITIAL SPECIFIED HEAD DATA)
834 FORMAT (4I5,2F10.0)
836 FORMAT (1HO,22H NO. OF NODES IN ZONE, $\mathrm{I} 4,3 \mathrm{H}=, \mathrm{I} 3 / 1 \mathrm{H}$
$1,20 \mathrm{H}$ PAR. NO. HEAD $A=, I 3,12 \mathrm{X}, 18 \mathrm{HPAR}$. NO. HEAD $\mathrm{B}=, \mathrm{I} 3 / 1 \mathrm{H}$
$2,21 \mathrm{H}$ STD. DEV. HEAD $A=$, G11.5,22H STD. DEV. HEAD B $=$, G11. 5
$3 / 1 \mathrm{H}, 20 \mathrm{X}, 22 \mathrm{HINITIAL}$ VALUES OF $\mathrm{HEAD} / 1 \mathrm{H}, 21 \mathrm{X}, 1 \mathrm{HI}, 5 \mathrm{X}, 1 \mathrm{HJ}, 8 \mathrm{X}, 4 \mathrm{HHEAD}$ )
840 FORMAT (1H ,19X,2(I3,3X),2X,G11.5)
842 FORMAT (9HOAT CELL , 1H(, I3, 1H, $13,1 \mathrm{H}), 23 \mathrm{H}, \mathrm{I} Z N>0, \mathrm{CX}=0$, AND $\mathrm{CY}=0$ )
844 FORMAT (60HOPROGRAM ABORTED BECAUSE OF CONFLICT BETWEEN IZN, CX, A
IND CY)
846 FORMAT (18HO INITIAL SOLUTION)

## Program Listing-Continued

848 FORMAT (46HONO. OF PARAMETERS HAVING PRIOR INFORMATION $=$,I4)
850 FORMAT ( $1 \mathrm{HO} 0,3 \mathrm{X}, 55 \mathrm{HOBSERVATION}$ NUMBER AND SENSITIVITIES FOR EACH PA 1RAMETER)
852 FORMAT (1H ,1X,14,7(1X,G11.5)/(1H ,(5X,7(1X,G11.5))))
856 FORMAT ( $1 \mathrm{H}, 1 \mathrm{X}, \mathrm{I} 4,3 \mathrm{X}, 4(\mathrm{G} 11.5,4 \mathrm{X})$ )
858 FORMAT (11HOPARAMETER , I3,17H EFFECTIVELY ZERO)
860 FORMAT (//32HOSOLUTION FAILED TO CONVERGE IN , I3,11H ITERATIONS)
862 FORMAT (//23HOSOLUTION CONVERGED IN ,I3,11H ITERATIONS)
864 FORMAT (18HOERROR VARIANCE $=$, G11.5/35H ESTIMATED SUM OF SQUARED E 1RRORS $=$, G11.5/27H•CORRELATION COEFFICIENT $=$, G11.5)
870 FORMAT ( $1 \mathrm{H} 0,5 \mathrm{X}, 24 \mathrm{HESTIMATED}$ PARAMETER DATA/1H , 6H PAR.
$1 / 1 \mathrm{H}, 6 \mathrm{H}$ NO., $5 \mathrm{X}, 4 \mathrm{HPAR} ., 9 \mathrm{X}, 9 \mathrm{HSTD}$. DEV.)
874 FORMAT (28HO VARIANCE-COVARIANCE MATRIX)
876 FORMAT (20HO CORRELATION MATRIX)
878 FORMAT ( $1 \mathrm{HO} 0,21 \mathrm{X}, 14 \mathrm{HHEAD}$ RESIDUALS $/ 1 \mathrm{H}, 7 \mathrm{H}$ OBS.,5X,9HPREDICTED 1, $7 \mathrm{X}, 8 \mathrm{HOBSERVED}, 8 \mathrm{X}, 8 \mathrm{HWEIGHTED} / 1 \mathrm{H}, 4 \mathrm{X}, 3 \mathrm{HNO} ., 7 \mathrm{X}, 5 \mathrm{HVALUE}, 10 \mathrm{X}$ 2,5HVALUE, 10X, 8HRESIDUAL)
880 FORMAT (1H , 2X, 14, 1X, 3(5X, G11.5))
881 FORMAT (33HO FINAL COMPUTED NODAL HEAD ARRAY)
882 FORMAT (26HO NODAL SENSITIVITY ARRAYS)
884 FORMAT (19HO PARAMETER NUMBER ,I5)
886 FORMAT ( $1 \mathrm{H} 0,3 \mathrm{X}, 44 \mathrm{HSCALED}$ AND ORTHOGONALIZED SENSITIVITY MATRIX $1 / 1 \mathrm{H}, 49 \mathrm{H}$ OBSERVATION NUMBER AND VALUES FOR EACH PARAMETER)
888 FORMAT (7H OBS.,13X,14HPARAMETER NOS./1H ,3X,3HNO.,4X,8(I3,9X))
890 FORMAT ( $1 \mathrm{H}, 1 \mathrm{X}, \mathrm{I4}, 8(1 \mathrm{X}, \mathrm{G11} .5)$ )
891 FORMAT (24HOADDITIONAL SOLUTION NO., I5)
892 FORMAT (1H0,30X,14HNEW PARAMETERS/1H ,3(4X,3HNO., $8 \mathrm{X}, 5 \mathrm{HVALUE}, 4 \mathrm{X})$ )
898 FORMAT (16HO COMPUTED HEADS) END

## Program Listing-Continued

SUBROUTINE ARRAY(A, IND, JND, IT)
IMPLICIT DOUBLE PRECISION (A-H,0-Z)
DIMENSION A(IND, JND)
COMMON/TNME/IIN,IOUT
C**IF IT=0, LOAD 1 AND 2 DIMENSIONAL ARRAYS
C**IF IT=1, PRINT 2 DIMENSIONAL ARRAYS
IF(IT.EQ.1) GO TO 55
DO $5 \mathrm{~J}=1, \mathrm{JND}$
DO $5 \mathrm{I}=1$, IND
$5 \mathrm{~A}(\mathrm{I}, \mathrm{J})=0$.
READ (IIN,65) NME, NOBL, IPRN
WRITE (IOUT,75) NME
DO $50 \mathrm{~K}=1$, NOBL
READ (IIN, 70) IB, IE, JB, JE, FACT, IVAR
WRITE(IOUT, 80) K, IB, IE, JB, JE, FACT
IF(IVAR.GT.0) GO TO 20
DO $10 \mathrm{~J}=\mathrm{JB}, \mathrm{JE}$
DO $10 \mathrm{I}=\mathrm{IB}, \mathrm{IE}$
$10 \mathrm{~A}(\mathrm{I}, \mathrm{J})=\mathrm{FACT}$
GO TO 50
20 DO $40 \mathrm{~J}=\mathrm{JB}, \mathrm{JE}$
$\operatorname{READ}(\mathrm{IIN}, 90)(\mathrm{A}(\mathrm{I}, \mathrm{J}), \mathrm{I}=\mathrm{IB}, \mathrm{IE})$
DO $40 \mathrm{I}=\mathrm{IB}, \mathrm{IE}$
$40 \mathrm{~A}(\mathrm{I}, \mathrm{J})=\mathrm{A}(\mathrm{I}, \mathrm{J}) * \mathrm{FACT}$
50 CONTINUE
IF(IPRN.GT.0) RETURN
WRITE(IOUT,100) NME
55 DO $60 \mathrm{~K}=1$, IND, 10
$110=\mathrm{K}+9$
IF(IIO.GT.IND) $110=1 N D$
WRITE(IOUT,110) (I, I=K,II0)
WRITE(IOUT,105)
DO $60 \mathrm{~J}=1$, JND
$\mathrm{JR}=\mathrm{JND}-\mathrm{J}+1$
60 WRITE(IOUT,120) JR,(A(I,JR), $\mathrm{I}=\mathrm{K}, \mathrm{I} 10$ )
RETURN
C
65 FORMAT (A4,1X,2I5)
70 FORMAT (4I5,F10.0,3I5)
75 FORMAT ( $1 \mathrm{HO}, \mathrm{A} 4$ )
80 FORMAT ( $1 \mathrm{H}, \mathrm{I} 3,2 \mathrm{X}, 5 \mathrm{HIB}=, \mathrm{I} 5,2 \mathrm{X}, 5 \mathrm{HIE}=, \mathrm{I} 5,2 \mathrm{X}, 5 \mathrm{HJB}=, \mathrm{I} 5,2 \mathrm{X}$
$1,5 \mathrm{HJE}=, \mathrm{I} 5,2 \mathrm{X}, 7 \mathrm{HFACT}=, \mathrm{G11.5}$ )
90 FORMAT (8F10.0)
100 FORMAT (1HO,1X,A4,6H ARRAY)
105 FORMAT (1H )
110 FORMAT ( $1 \mathrm{HO} 0,10(9 \mathrm{X}, 13)$ )
120 FORMAT (1H ,I3,1X,10(1X,G11.5))
END
SUBROUTINE ARRAYI (INT, IND, JND,IT)
DIMENSION INT(IND,JND)
COMMON/TNME/IIN, IOUT
C**IF IT=0, LOAD 1 AND 2 DIMENSIONAL INTEGER ARRAYS
C**IF IT=1, PRINT 2 dIMENSIONAL INTEGER ARRAYS

Program Listing-Continued

```
    IF(IT.EQ.1) GO TO 45
    DO 5 J=1,JND
    DO 5 I=1,IND
    5 INT (I,J)=0
    READ(IIN,55) NME,NOBL,IPRN
    WRITE(IOUT,65) NME
    DO 40 K=1,NOBL
    READ(IIN,60) IB,IE,JB,JE,IFACT,IVAR
    WRITE(IOUT,70) K,IB,IE,JB,JE,IFACT
    IF(IVAR.GT.0) GO TO 20
    DO 10 J=JB,JE
    DO 10 I=IB,IE
    10 INT(I,J)=IFACT
    GO TO 40
    20 DO 30 J=JB,JE
    READ(IIN,60) (INT(I,J),I=IB,IE)
    30 CONTINUE
    40 CONTINUE
    IF(IPRN.GT.0) RETURN
    WRITE(IOUT,80) NME
    45 DO 50 K=1,IND,30
    I30=K+29
    IF(I30.GT.IND) I 30=IND
    WRITE(IOUT, 90) (I, I=K,I30)
    WRITE(IOUT, 100)
    DO 50 J=1,JND
    JR=JND-J+1
    50 WRITE(IOUT, 110) JR,(INT(I,JR), I=K,I30)
    RETURN
    C
    55 FORMAT (A4,1X,2I5)
    60 FORMAT (16I5)
    65 FORMAT (1H0,A4)
    70 FORMAT (1H ,I3,2X,5HIB = ,I5,2X,5HIE = ,I5,2X,5HJB = ,I5,2X
    1,5HJE = ,I5,2X,8HIFACT = ,I5)
    80 FORMAT (1H0,1X,A4,6H ARRAY)
    90 FORMAT (1H0,3X,30(1X,I3))
    100 FORMAT (1H)
    110 FORMAT (1H ,31(I3,1X))
        END
        SUBROUTINE ORDER(JPOS,IN,IC)
        DIMENSION JPOS(1),IN(1),IC(5,1)
        COMMON/INT/NIJ,NEQ,ICR,ICR1,IB1,LH1,ID, JD,IM, JM,NOBS,NQSD,NBH
        1,NVAR,NVX2,KOUNT, INDT,IPO
        COMMON/TNME/IIN, IOUT
C**COMPUTE EQUATION NUMBERS FOR D4 ORDERING:
        NXP=ID+JD-1
        K=0
C**ORDER--LEFT TO RIGHT, BOTTOM TO TOP
    DO 20 I=1,NXP,2
    DO 20 J=1,JD
    IK=I-J+1
    IF(IK.LT.1.OR.IK.GT.ID) GO TO 20
```

Program Listing-Continued

```
    N=IK+JPOS(J)
    IF(IN(N).LT.1) GO TO 20
    K=K+1
    IN(N)=K
    20 CONTINUE
    ICR=K+1
    DO 30 I=2,NXP,2
    DO 30 J=1,JD
    IK=I-J+1
    IF(IK.LT.I.OR.IK.GT.ID) GO TO 30
    N=IK+JPOS(J)
    IF(IN(N).LT.1) GO TO 30
    K=K+1
    IN(N)=K
    30 CONTINUE
C**COMPUTE BAND WIDTH AND DETERMINE CONNECTING EQUATION NUMBERS:
    MNO=9999
    MXO=0
    N=0
    IND=0
    DO 80 J=1,JD
    DO }80\textrm{I}=1,\textrm{ID
    N=N+1
    JR=IN(N)
    IF(JR.LT.l.OR.JR.GE.ICR) GO TO }8
    IU=1
C**BELOW
    IF((J-1).LT.1.OR.IN(N-ID).LT.1) GO TO 40
    IU=IU+1
    IC(IU,JR)=IN(N-ID)
    MM=IN (N-ID)-JR
    MXO=MAXO (MM,MXO)
    MNO=MINO (MM,MNO)
C**LEFT
    40 IF((I-1).LT.1.OR.IN(N-1).LT.1) GO TO 50
    IU=IU+1
    IC(IU,JR)=IN(N-1)
    MM=IN(N-1)-JR
    MNO=MINO (MM,MNO)
    MXO=MAXO (MM,MXO)
C**RIGHT
    50 IF((I+1).GT.ID.OR.IN(N+1).LT.1) GO TO 60
            IU=IU+1
            IC(IU,JR)=IN(N+1)
            MM=IN(N+1)-JR
            MXO=MAXO (MM,MXO)
            MNO=MINO (MM,MNO)
C**ABOVE
    60 IF((J+1).GT.JD.OR.IN(N+ID).LT.1) GO TO 70
            IU=IU+1
            IC(IU,JR)=IN(N+ID)
            MM=IN(N+ID)-JR
            MXO=MAXO (MM,MXO)
```Program Listing-Continued
```

```
    MNO=MINO (MM,MNO)
    70 IC (1,JR)=IU
        IF(IU.GT.1) GO TO }8
        WRITE(IOUT,110) I,J
        IND=1
        8 0 ~ C O N T I N U E ~
        IF(IND.GT.0) STOP
        NEQ=K
        ICRI=ICR-1
        IB1=MXO-MNO+1
        LHl=NEQ-ICRI
        WRITE(IOUT,90)
        WRITE(IOUT,100) ICR1,IB1,LH1,ICR1,NEQ
        RETURN
    C
        90 FORMAT (51HOSOLUTION BY LDU FACTORIZATION ASSUMING D4 ORDERING)
        100 FORMAT (65H MINIMUM DIMENSIONS FOR ARRAYS USED BY THIS METHOD ARE
        1AS FOLLOWS/1H ,12H AU: 5 BY,I5/1H ,4H AL:,I5,3H BY,I5/1H
        2,12H IC: }5\mathrm{ BY,I5/1H ,4H V:,I5)
        110 FORMAT (1H0,13HACTIVE NODE (,I3,1H,,I3,2OH) CANNOT BE ISOLATED)
        END
        SUBROUTINE COEF(WELL,HR,HC,CX,CY,VL,QR,CXHR,CXHL,CYHT,CYHB,QBF
        l, B,AU,AL,V,IZN, IBZN, IPRM, IBNA, IBNB, IN , IC ,NAD)
        IMPLICIT DOUBLE PRECISION (A-H,O-Z)
        DIMENSION WELL(1),HR(1),HC(1),CX(1),CY(1),VL(1),QR(1),CXHR(1)
        1,CXHL(1),CYHT (1) ,CYHB (1),QBF (1) , B(1),AU(5,1),AL (NAD , 1),V(1)
        DIMENSION IZN(1),IBZN(1),IPRM(4,1),IBNA(1),IBNB(1),IN(1),IC(5,1)
        COMMON/INT/NIJ,NEQ, ICR, ICR1, IB1, LH1, ID , JD, IM , JM, NOBS , NQSD , NBH
        1,NVAR,NVX2,KOUNT, INDT,IPO
    C**INITIALIZE ARRAYS
        DO 10 J=1,ICR1
        DO 10 I=1,5
        10 AU(I,J)=0.
        DO 20 J=1,LH1
        DO 20 I=1,IB1
    20 AL(I,J)=0.
        DO 40 I=1,NIJ
        N=IN(I)
        IF(N.GT.0) V(N)=WELL(I)
    4 0 ~ C O N T I N U E ~
        IF(NBH.LT.1) GO TO 45
        DO 42 I=1,NBH
        CXHR(I)=0.
        CXHL (I) =0.
        CYHT(I)=0.
    42 СYHB(I)=0.
C**CALCULATE V FOR SPECIFIED FLOW PARAMETERS
    45 IF(NQSD.LT.1) GO TO 52
        DO 50 I=1,NQSD
        IZ=IBZN(I)
        TMP=B(IZ)*QBF(I)
        INA=IBNA(I)
        L=IN(INA)
```

Program Listing-Continued

```
    IF(L.GT.0) V(L)=V(L)+TMP
    INB=IBNB (I)
    L=IN(INB)
    IF(L.GT.0) V(L)=V(L)+TMP
    5 0 ~ C O N T I N U E ~
C**BEGIN MAIN LOOP
    52 N=0
    DO 150 J=1,JM
    DO 150 I=1,IM
    N=N+1
    M=IZN(N)
    IF(M.LT.1) GO TO 150
    LTX=IPRM(1,M)
    LTY=IPRM(2,M)
    LVL=IPRM(3,M)
    LQD=IPRM(4,M)
    NB=N+J
    NA=NB-1
    NC=NB+ID
    ND=NA+ID
    INA=IN(NA)
    INB=IN(NB)
    INC=IN(NC)
    IND=IN(ND)
    CXT=B(LTX)*CX(N)
    CYT=B(LTY)*CY(N)
    VT.T=B (LVL)*VL(N)
    QRT=B(LQD)*QR(N)
    E=CXT+CYT+VLT
C**CALCULATE AU, AL, V, AND COEFFICIENT ARRAYS FOR SPECIFIED HEAD
C PARAMETERS
    K=-INA-1
    IF(K) 60,75,53
    53 CXHR(K)=CXHR(K)+CXT
    CYHT (K)=CYHT (K)+CYT
    GO TO 75
    60 IF(INA.GE.ICR) GO TO 65
    AU (1, INA) =AU (1, INA)+E
    AU(4,INA) =AU(4,INA)-CXT
    AU (5,INA) =AU (5, INA) -CYT
    GO TO 70
    65 AL (1,INA-ICR1)=AL (1,INA-ICR1)+E
    70 V(INA)=V(INA)+QRT+VLT*(HR (NA)-HC(NA))+CXT* (HC(NB)-HC(NA))
    1+CYT*(HC (ND)-HC (NA))
    75 K=-INB-1
        IF(K) 85,100,77
    77 CXHL (K)=CXHL (K)+CXT
        CYHT (K)=CYHT (K)+CYT
        GO TO 100
    85 IF(INB.GE.ICR) GO TO 90
    AU(1,INB)=AU(1,INB)+E
    AU(3,INB)=AU(3,INB)-CXT
    AU(5,INB)=AU(5,INB)-CYT
```


## Program Listing-Continued

GO TO 95
$90 \mathrm{AL}(1, \mathrm{INB}-\mathrm{ICR} 1)=\mathrm{AL}(1, \mathrm{INB}-\mathrm{ICR} 1)+\mathrm{E}$
$95 \mathrm{~V}(\mathrm{INB})=\mathrm{V}(\mathrm{INB})+\mathrm{QRT}+\mathrm{VLT} *(\mathrm{HR}(\mathrm{NB})-\mathrm{HC}(\mathrm{NB}))+\mathrm{CXT} *(\mathrm{HC}(\mathrm{NA})-\mathrm{HC}(\mathrm{NB}))$
$1+\mathrm{CYT}$ ( $\mathrm{HC}(\mathrm{NC})-\mathrm{HC}(\mathrm{NB}))$
$100 \mathrm{~K}=-$ INC -1
IF (K) $110,125,102$
102 CXHL (K) $=$ CXHL (K) + CXT
CYHB (K) $=$ CYHB ( $K$ ) + CYT
GO TO 125
110 IF (INC.GE.ICR) GO TO 115
$\mathrm{AU}(1, \mathrm{INC})=\mathrm{AU}(1, \operatorname{INC})+E$
$\operatorname{AU}(2, \operatorname{INC})=\operatorname{AU}(2, \operatorname{INC})-\mathrm{CYT}$
$\mathrm{AU}(3, \operatorname{INC})=\mathrm{AU}(3, \mathrm{INC})-\mathrm{CXT}$
GO TO 120
$115 \mathrm{AL}(1, \mathrm{INC}-\mathrm{ICR} 1)=\mathrm{AL}(1, \mathrm{INC}-\mathrm{ICR} 1)+\mathrm{E}$
$120 \mathrm{~V}(\mathrm{INC})=\mathrm{V}(\mathrm{INC})+\mathrm{QRT}+\mathrm{VLT} \mathrm{T}^{*}(\mathrm{HR}(\mathrm{NC})-\mathrm{HC}(\mathrm{NC}))+\mathrm{CXT} *(\mathrm{HC}(\mathrm{ND})-\mathrm{HC}(\mathrm{NC}))$
$1+\mathrm{CYT}$ ( $\mathrm{HC}(\mathrm{NB})-\mathrm{HC}(\mathrm{NC}))$
$125 \mathrm{~K}=-\mathrm{IND}-1$
IF(K) $135,150,127$
$127 \operatorname{CXHR}(\mathrm{~K})=$ CXHR (K) +CXT
CYHB (K) $=$ CYHB (K) + CYT
GO TO 150
135 IF(IND.GE.ICR) GO TO 140
$\mathrm{AU}(1, \operatorname{IND})=\mathrm{AU}(1, \mathrm{IND})+E$
$\mathrm{AU}(2, \mathrm{IND})=\mathrm{AU}(2, \mathrm{IND})-\mathrm{CYT}$
$\operatorname{AU}(4, I N D)=\operatorname{AU}(4, I N D)-C X T$
GO TO 145
$140 \mathrm{AL}(1$, IND - ICR1 $)=\mathrm{AL}(1$, IND-ICR1) +E
$145 \mathrm{~V}(\mathrm{IND})=\mathrm{V}(\mathrm{IND})+\mathrm{QRT}+\mathrm{VLT} *(\mathrm{HR}(\mathrm{ND})-\mathrm{HC}(\mathrm{ND}))+\mathrm{CXT} *(\mathrm{HC}(\mathrm{NC})-\mathrm{HC}(\mathrm{ND}))$
$1+\mathrm{CYT}$ ( $\mathrm{HC}(\mathrm{NA})-\mathrm{HC}(\mathrm{ND}))$
150 CONTINUE
C**COMPRESS AU
$\mathrm{N}=0$
D0 $190 \mathrm{~J}=1, \mathrm{JD}$
D0 $190 \mathrm{I}=1$, ID
$\mathrm{N}=\mathrm{N}+\mathrm{I}$
$\mathrm{K}=\mathrm{IN}$ ( N )
IF(K.LT.1.OR.K.GT.ICR1.OR.IC(1,K).EQ.5) GO TO 190
IU=1
IF ((J-1).LT.1.OR.IN(N-ID).LT.1) GO TO 160
$\mathrm{IU}=\mathrm{IU}+1$
$\mathrm{AU}(\mathrm{IU}, \mathrm{K})=\mathrm{AU}(2, \mathrm{~K})$
$160 \mathrm{IF}((\mathrm{I}-1) . \operatorname{LT} .1 . \mathrm{OR} . \operatorname{IN}(\mathrm{N}-1) . \operatorname{LT} .1) \mathrm{GO}$ TO 170
$I U=I U+1$
$\operatorname{AU}(I U, K)=A U(3, K)$
170 IF ( $(\mathrm{I}+1)$. GT.ID.OR.IN(N+1).LT.1) GO TO 180
$\mathrm{IU}=\mathrm{IU}+1$
$\mathrm{AU}(\mathrm{IU}, \mathrm{K})=\mathrm{AU}(4, \mathrm{~K})$
180 IF ( (J+1).GT.JD.OR.IN(N+ID).LT.1) GO TO 190
$\mathrm{IU}=\mathrm{IU}+1$
$A U(I U, K)=A U(5, K)$
190 CONTINUE
RETURN

```
Program Listing-Continued
    END
    SUBROUTINE D4SOLV(HC,AU,AL,V,IN,IC,NAD)
    IMPLICIT DOUBLE PRECISION (A-H,O-Z)
    DIMENSION HC (1), AU (5,1), AL (NAD,1),V(1)
    DIMENSION IN(1),IC(5,1)
    COMMON/INT/NIJ,NEQ,ICR, ICR1, IB1, LH1,ID,JD, IM, JM, NOBS,NQSD,NBH
    1,NVAR,NVX2,KOUNT, INDT,IPO
C**DECOMPOSE TO FILL AL
    DO 280 J=1,ICR1
    II=IC(1,J)
    DO 270 I=2,II
    LR=IC(I,J)
    L=LR-ICR1
    C=AU(I,J)/AU(1,J)
    DO 260 K=I,II
    KL=IC(K,J)-LR+1
    AL(KL,L)=AL (KL,L)-C*AU (K,J)
    260 CONTINUE
            AU(I,J)=C
            V(LR)=V(LR)-C*V(J)
    270 CONTINUE
    280 V(J)=V(J)/AU(1,J)
C**DECOMPOSE AL
    JJ=NEQ-ICR
            DO 310 J=1,JJ
            JR=J+ICR1
            L=J
            DO 300 I=2,IB1
            L=L+1
            IF(AL(I,J).EQ.O.) GO TO 300
            LR=L+ICR1
            C=AL(I,J)/AL(1,J)
            KL=0
            DO 290 K=I,IB1
            KL=KL+1
            IF(AL(K,J).NE.0.) AL(KL,L)=AL(KL,L)-C*AL(K,J)
        290 CONTINUE
            AL(I,J)=C
            V(LR)=V(LR)-C*V(JR)
    300 CONTINUE
    310 V(JR)=V(JR)/AL(1,J)
C**BACK SOLVE--LOWER HALF
            V(NEQ)=V(NEQ)/AL(1,NEQ-ICR1)
            DO 330 J=l,JJ
            K=NEQ-J
            KL=K-ICR1
            L=K
            DO 320 I=2,IB1
            L=L+1
            IF(AL(I,KL).NE.0.) V(K)=V(K)-AL(I,KL)*V(L)
        320 CONTINUE
    330 CONTINUE
C**BACK SOLVE--UPPER HALF
```

```
DO 350 J=1,ICR1
K=ICR-J
II=IC(1,K)
DO 340 I=2,II
L=IC(I,K)
V(K)=V(K)-AU(I,K)*V(L)
    340 CONTINUE
    350 CONTINUE
C**COMPUTE HC+DELTHC
DO 360 N=1,NIJ
L=IN (N)
IF(L.LT.1) GO TO 360
HC(N)=HC(N)+V(L)
    360 CONTINUE
        RETURN
        END
        SUBROUTINE LSTSQ(HC,BK,BL,BM,BN,HO,W,P,WP,X,C,B,V,KOBS,IN,NVD)
        IMPLICIT DOUBLE PRECISION (A-H,O-Z)
        DIMENSION HC(1),BK(1),BL(1),BM(1),BN(1),HO(1),W(1),P(1),WP(1)
        1,X(NVD,1),C(NVD,1),B(1),V(3)
            DIMENSION KOBS(1),IN(1)
            COMMON/INT/NIJ,NEQ,ICR,ICR1,IB1,LH1,ID,JD,IM, JM,NOBS,NQSD,NBH
            1,NVAR,NVX2,KOUNT,INDT,IPO
            COMMON/TNME/IIN, IOUT
            COMMON/REG/DMX, ADMX,AP, CSA, AMP, RP, BP,YSQ
C**CHECK FOR NONZERO MARQUARDT PARAMETER
    NMI=NVAR-1
    IF(AMP.LT.-.5) GO TO 105
C**INITIALIZE
            DO 20 J=1,NVAR
            DO 10 I=1,NVAR
        10 C(I,J)=0.
        20 V(J)=0.
            YSQ=0.
C**FORM COEFFICIENT MATRIX AND R.H.S. VECTOR
            DO 70 N=1,NOBS
            K=KOBS (N)
            TEMP=HO(N)-BK(N)*HC(K)-BL(N)*HC(K+1)-BM(N)*HC(K+ID)-BN(N)
            1*HC(K+ID+1)
            DO 60 J=1,NVAR
            TMP=W(N)*X(J,N)
            DO }50\mathrm{ I=J,NVAR
        50 C(I,J)=X(I,N)*TMP+C(I,J)
        60 V(J)=TMP*TEMP+V (J)
            YSQ=YSQ+TEMP*W(N)*TEMP
        70 CONTINUE
            IF(NVAR.EQ.1) GO TO 190
            DO }80\textrm{I}=1\mathrm{ ,NVAR
            TEMP=C(I,I)+WP(I)
            IF(TEMP.GT.1.E-10) GO TO 78
            WRITE(IOUT, 260) I
            INDT=1
            GO TO }8
```

Program Listing-Continued
$78 \mathrm{C}(\mathrm{I}, \mathrm{I})=$ TEMP**. 5
80 CONTINUE
IF (INDT.GT.0) RETURN
DO $100 \mathrm{~J}=1$, NMI
TEMP $=\mathrm{C}(\mathrm{J}, \mathrm{J})$
$\mathrm{JP} 1=\mathrm{J}+1$
DO $90 \mathrm{I}=\mathrm{JP} 1$, NVAR
$C(I, J)=C(I, J) /(C(I, I) * T E M P)$
$90 \mathrm{C}(\mathrm{J}, \mathrm{I})=\mathrm{C}(\mathrm{I}, \mathrm{J})$
$V(J)=(V(J)+W P(J) *(P(J)-B(J))) / T E M P+R P * T E M P *(B P * P(J)-B(J))$
$V(J+N V A R)=V(J)$
$V(J+N V X 2)=T E M P$
$100 \mathrm{C}(\mathrm{J}, \mathrm{J})=1 .+\mathrm{RP}+\mathrm{AMP}$
TEMP $=$ C (NVAR, NVAR)
V (NVAR) $=(\mathrm{V}($ NVAR $)+$ WP $($ NVAR $) *(P(N V A R)-B(N V A R))) / T E M P$
1+RP*TEMP* (BP*P (NVAR) -B (NVAR))
$V($ NVX2 $)=V(N V A R)$
$V(N V A R+N V X 2)=T E M P$
$C($ NVAR , NVAR $)=1 .+R P+A M P$
IF (IPO.NE.1) GO TO 105
WRITE (IOUT, 250)
CALL PRTOT (C,NVAR,NVD,0)
WRITE (IOUT, 255)
WRITE (IOUT, 230) (V(I), I=1,NVAR)
C**SOLVE FOR V USING LDU FACTORIZATION:
C**DECOMPOSITION AND FORWARD SUBSTITUTION
$105 \mathrm{DET}=1$.
DO $140 \mathrm{~K}=1$,NM1
PIV $=\mathrm{C}(\mathrm{K}, \mathrm{K})$
DET=DET*PIV
IF (DABS (PIV). GT. 1.E-10) GO TO 110
WRITE(IOUT, 210)
INDT=1
RETURN
110 PIV=1./PIV
$\mathrm{KPl}=\mathrm{K}+1$
DO $130 \mathrm{~J}=\mathrm{KP} 1$, NVAR
$\mathrm{TMP}=\mathrm{C}(\mathrm{J}, \mathrm{K}) * \mathrm{PIV}$
DO $120 \mathrm{I}=\mathrm{J}, \mathrm{NVAR}$
$120 C(I, J)=C(I, J)-T M P * C(I, K)$
$130 \mathrm{~V}(\mathrm{~J})=\mathrm{V}(\mathrm{J})-\mathrm{TMP} * \mathrm{~V}(\mathrm{~K})$
$C(K, K)=P I V$
140 CONTINUE
$D E T=D E T * C$ (NVAR,NVAR)
IF (DABS (C(NVAR,NVAR)).GT.1.E-10) GO TO 150
WRITE (IOUT, 210)
INDT=1
RETURN
150 IF(AMP.LT.-.5) RETURN
C**BACK SUBSTITUTION
$\mathrm{V}(\mathrm{NVAR})=\mathrm{V}$ (NVAR) $/ \mathrm{C}($ NVAR, NVAR $)$
I=NVAR
160 I=I-1

Program Listing-Continued

```
    IF(I.LT.1) GO TO 175
    IP1=I+1
    SUM=0.
    DO 170 J=IP1,NVAR
    170 SUM=SUM+C(J,I)*V(J)
    V(I)=(V(I)-SUM)*C(I,I)
    GO TO 160
    C**CHECK SOLUTION AND ADD MARQUARDT PARAMETER IF NEEDED
    175 TMPA=0.
            TMPB=0.
            TMPC=0.
            DO 176 I=1,NVAR
            TMPA=TMPA+V (I)*V(I)
            TMPB=TMPB+V(I+NVAR)*V(I+NVAR)
    176 TMPC=TMPC+V(I)*V(I+NVAR)
            IF(TMPC.GT.CSA*DSQRT(TMPA*TMPB)) GO TO 200
            AMP=1.5*AMP+.001
            IF(AMP.GT.1.) GO TO 200
            DO 180 I=1,NVAR
            V(I) =V(I+NVAR)
            C(I,I)=1.+RP+AMP
            DO 178 J=I,NVAR
    178 C(J,I)=C(I,J)
    180 CONTINUE
            GO TO 105
C**SOLUTION WHEN NVAR=1
    190 TEMP=C (1,1)+WP(1)
            IF(TEMP.GT.1.E-10) GO TO 195
            I=1
            WRITE(IOUT, 260) I
            INDT=1
            RETURN
    195 V(3)=TEMP**. }
            V(2)=(V(1)+WP(1)* (P(1)-B(1)))/V(3)+RP*V(3)*(BP*P(1)-B(1))
            C(1, 1)=1.+RP
            DET=C(1,1)
            V(1)=V(2)/DET
            IF(IPO.NE.1) GO TO 200
            WRITE(IOUT, 250)
            CALL PRTOT(C,1,NVD,0)
            WKITE(IOUT,255)
            WRITE(IOUT,230) V(1)
C**COMPUTE AND PRINT PARAMETERS
    200 ADMX=0.
        DO 203 J=1,NVAR
        V(J)=V (J)/V(J+NVX2)
        TMPA=1.
        IF(B(J)) 201,202,201
    201 TMPA=B(J)
    202 TMP=DABS (V(J)/TMPA)
        IF(TMP.GT.ADMX) ADMX=TMP
    203 CONTINUE
        AP}=1
```

Program Listing-Continued

```
        IF(ADMX.GT.DMX) AP=DMX/ADMX
        DO 204 J=1,NVAR
        V(J)=AP*V(J)
    204 B(J)=V(J)+B(J)
        WRITE(IOUT,220) KOUNT,YSQ,DET,AMP,AP
        WRITE(IOUT,230) (B(J),J=1,NVAR)
        RETURN
C
    210 FORMAT (42HOLEAST SQUARES COEFFICIENT MATRIX SINGULAR
    1/35H SOLUTION FOR PARAMETERS NOT UNIQUE)
    220 FORMAT (1HO,14HITERATION NO. ,I3/1H ,6HYSQ = ,G11.5,2X
    1,9HDET(C) = ,G11.5,2X,6HAMP = ,G11.5,2X,5HAP = ,G11.5
    2/1H ,21HREGRESSION PARAMETERS)
    230 FORMAT ((1H ,8(G11.5,2X)))
    250 FORMAT (29H0 SCALED LEAST SQUARES MATRIX)
    255 FORMAT (24HO SCALED GRADIENT VECTOR)
    260 FORMAT (29HOSENSITIVITIES FOR PARAMETER ,I4,17H EFFECTIVELY ZERO)
        END
        SUBROUTINE PRTOT(C,NO,NOD,IT)
        IMPLICIT DOUBLE PRECISION (A-H,O-Z)
        DIMENSION C(1)
        COMMON/TNME/IIN,IOUT
C**IF IT=0, PRINT SYMMETRIC MATRIX DIVIDED VERTICALLY INTO TEN-COLUMN
C BLOCKS
C**IF IT=1, PRINT VECTOR IN THREE COLUMNS
            IF(IT.EQ.1) GO TO 25
            DO 20 L=1,NO,10
            J10=L+9
            IF(J10.GT.NO) J10=NO
            WRITE(IOUT,30) (J,J=L,J10)
            WRITE(IOUT,50)
            K=-NOD
            DO 10 I=1,NO
            K=K+NOD
    10 WRITE (IOUT,40) I, (C(J+K),J=L,J10)
            WRITE(IOUT,60)
    20 CONTINUE
            RETURN
    25 NR=NO/3
            IF((3*NR).NE.NO) NR=NR+1
            DO 26 K=1,NR
    26 WRITE(IOUT,80) (L,C(L),L=K,NO,NR)
            RETURN
c
    30 FORMAT (1HO,8X,T3,9(9X,I3))
    40 FORMAT (1H ,I3,10(1X,G11.5)
    50 FORMAT (1H)
    60 FORMAT (1HO)
    80 FORMAT (1H ,3X,3(13,7X,G11.5,3X))
        END
```

Program Listing-Continued
SUBROUTINE HOBS (DX,DY, X,Y,BK, BL, BM, BN , HO, W, KOBS)
IMPLICIT DOUBLE PRECISION (A-H,O-Z)
DIMENSION DX(1),DY(1),X(1),Y(1),BK(1),BL(1),BM(1),BN(1),HO(1)
1,W(1)
DIMENSION KOBS(1)
COMMON/INT/NIJ,NEQ, ICR, ICR1, IB1, LH1, ID, JD, IM, JM, NOBS, NQSD, NBH
1,NVAR,NVX2,KOUNT, INDT, IPO
COMMON/TNME/IIN, IOUT
C**COMPUTE X-LOCATIONS OF NODE POINTS
$\mathrm{X}(1)=0$.
DO $10 \mathrm{I}=2$,ID
$10 \mathrm{X}(\mathrm{I})=\mathrm{X}(\mathrm{I}-1)+\mathrm{DX}(\mathrm{I}-1)$
C**COMPUTE Y-LOCATIONS OF NODE POINTS
$\mathrm{Y}(1)=0$.
DO $20 \mathrm{~J}=2$, JD
$20 \mathrm{Y}(\mathrm{J})=\mathrm{Y}(\mathrm{J}-1)+\mathrm{DY}(\mathrm{J}-1)$
WRITE(IOUT,40)
DO $30 \mathrm{I}=1$, NOBS
C**READ OBSERVED HEAD DATA
READ (IIN, 50) N, TL, JL, XL, YL, HO (N), W(N)
WRITE(IOUT,60) N, IL, JL, XL, YL, HO (N), W(N)
C**COMPUTE LOCATION OF FIRST NODE IN CELL (IL,JL)
K=IL+IM* (JL-1)
KOBS (N) $=\mathrm{K}+(\mathrm{K}-1) / \mathrm{IM}$
3
C**COMPUTE WEIGHTS FOR BILINEAR INTERPOLATION
AREA $=D X(I L) * D Y$ (JL)
$\mathrm{BK}(\mathrm{N})=(\mathrm{X}(\mathrm{IL}+1)-\mathrm{XL}) *(\mathrm{Y}(\mathrm{JL}+1)-\mathrm{YL}) /$ AREA
$\operatorname{BL}(\mathrm{N})=(\mathrm{XL}-\mathrm{X}(\mathrm{IL})) *(\mathrm{Y}(\mathrm{JL}+\mathrm{l})-\mathrm{YL}) /$ AREA
BM $(\mathrm{N})=(\mathrm{X}(\mathrm{IL}+1)-\mathrm{XL}) *(\mathrm{YL}-\mathrm{Y}(\mathrm{JL})) /$ AREA
$30 \operatorname{BN}(\mathrm{~N})=(\mathrm{XL}-\mathrm{X}(\mathrm{IL})) *(\mathrm{YL}-\mathrm{Y}(\mathrm{JL})) /$ AREA
C
40 FORMAT (1H0,25X,18HOBSERVED HEAD DATA/1H ,6H OBS., 2X,9HCELL LOC.
$1,3 \mathrm{X}, 6 \mathrm{HX}$ LOC., $7 \mathrm{X}, 6 \mathrm{HY}$ LOC. $, 8 \mathrm{X}, 4 \mathrm{HOBS} ., 8 \mathrm{X}, 6 \mathrm{HWEIGHT} / 1 \mathrm{H}, 6 \mathrm{H}$ NO., 4 X 2,1HI, 5X,1HJ, 30X,4HHEAD)
50 FORMAT ( $315,4 \mathrm{Fl} 0.0$ )
60 FORMAT (1H ,1X,3(I4,2X),4(G11.5,2X))
RETURN
END

## References Cited

Cooley, R.L., 1982, Incorporation of prior information on parameters into nonlinear regression groundwater flow models, 1-Theory: Water Resources Research, v. 18, no. 4, p. 965-976.

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Price, H.S., and Coats, K.H., 1974, Direct methods in reservoir simulation: Society of Petroleum Engineers Journal, June 1974, p. 295-308.
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## Additional Reading

Cooley, R.L., 1977, A method of estimating parameters and assessing reliability for models of steady-state groundwater flow, 1-Theory and numerical properties: Water Resources Research, v. 13, no. 2, p. 318-324.

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## 5 Elementary Analysis and Use of the Regression Model

### 5.1 Assumed Forms of Model Equations

As a purely algebraic process, regression contains no assumptions other than those already mentioned. However, to statistically analyze results of, and predictions to be made by, the method, additional assumptions must be made. Based on these assumptions, an effective methodology has been developed (see, for example, Draper and Smith, 1981) to analyze and use a linear regression model.

The statistical methods also may be applied to a nonlinear model, provided the model is close enough to being linear. Fortunately, whether or not the model is close enough can usually be determined. All statistics and procedures are, accordingly, derived for a linear, or effectively linear, model. To make the equations applicable for both a linear and nonlinear model, the basic types of models assumed are the incremental linear model and the nonlinear model as linearized using the Taylor series expansion.

The model assumed, then, is of the form

$$
\begin{equation*}
f_{\beta}-f_{0} \cong \underline{\underline{X}\left(\underline{\beta}-\underline{b}_{0}\right)} \tag{5.1-1}
\end{equation*}
$$

where strict equality applies for a linear model and, for a nonlinear model, $\underline{X}$ is assumed to be evaluated at $\underline{b}_{0}$. Also, for simplicity of notation, define

$$
\begin{align*}
& f_{\beta}=f(\xi, \beta)  \tag{5.1-2}\\
& f_{0}=\left(\xi, b_{0}\right) . \tag{5.1-3}
\end{align*}
$$

Based on equation 5.1-1, the true regression model is

$$
\begin{equation*}
\underline{Y}-f_{0} \cong \underline{\cong}\left(\underline{\beta}-\underline{b}_{0}\right)+\underline{\epsilon} \tag{5.1-4}
\end{equation*}
$$

where strict equality only holds for a linear model because $\underline{\epsilon}$ is the true vector of disturbances. The estimated regression model derived from equation 5.1-4 is

$$
\begin{equation*}
\underline{Y}-f_{0} \cong \underline{=} \underline{\underline{X}}\left(\underline{b}-\underline{b}_{0}\right)+\underline{e} \tag{5.1-5}
\end{equation*}
$$

where, as for equation 5.1-4, strict equality only holds for a linear model because $\underline{e}$ is assumed to be the true vector of residuals defined by $\underline{e}=\underline{Y}-f(\xi, \underline{b})$.

By minimizing $S(\underline{b})=e^{T} \underline{\omega e}$ with respect to $\underline{b}$ using the standard procedure, exact best-fit estimates $\underline{\hat{b}}$ of $\underline{\beta}$ and $\underline{\hat{e}}=Y-f(\underline{\xi}, \underline{\hat{b}})$ of $\underline{\epsilon}$ are obtained. For a linear model $\underline{b}$ is obtained exactly by using equation $5.1-\overline{5}$ as the estimated regression model. For a nonlinear model, use of the linearized model leads to an approximate relationship to find $\hat{b}$. Thus, by minimizing $S(\underline{b})$ using equation $5.1-5$ as the estimated regression model, normal equations

$$
\begin{equation*}
\underline{X}^{T} \underline{\underline{\omega}}\left(\underline{\hat{b}}-\underline{b}_{0}\right) \cong \underline{X}^{T} \underline{\omega}\left(\underline{\underline{Y}}-f_{0}\right) \tag{5.1-6}
\end{equation*}
$$

that are approximate for a nonlinear model are derived. The regression model obtained by replacing general estimates $\underline{b}$ and $\underline{e}$ in equation 5.1-5 with best-fit estimates $\underline{\underline{b}}$ and $\underline{\hat{e}}$ is

$$
\begin{equation*}
\underline{Y}-\mathcal{E}_{0} \cong \underline{X}\left(\underline{\hat{b}}-\underline{b}_{0}\right)+\underline{\hat{e}} . \tag{5.1-7}
\end{equation*}
$$

By utilizing the definition of $\underline{\hat{e}}(\underline{\hat{e}}=\underline{Y}-\hat{t}$, where $\hat{f}=f(\xi, \hat{b})$, in equation 5.1-7, a predictive model

$$
\begin{equation*}
\left.\hat{t}-\underline{f}_{0} \cong \underline{=} \underline{\underline{\hat{b}}}-\underline{b}_{0}\right) \tag{5.1-8}
\end{equation*}
$$

is obtained.
A final point is the establishment of the general condition for a minimum in $S(\underline{b})$. If $\underline{b}_{0}$ is set equal to $\underline{\hat{b}}$ in equation 5.1-6, then $f_{0}=\hat{f}$, and

$$
\begin{equation*}
\underline{X}^{T} \underline{\omega}(\underline{Y}-\hat{f})=\underline{0} . \tag{5.1-9}
\end{equation*}
$$

Because any approximation inherent in equation 5.1-6 is removed as $\underline{b}_{0} \rightarrow \underline{\hat{b}}$, equation 5.1-9 is exact for both linear and nonlinear models. The left side of equation 5.1-9 is the negative of the gradient of $S(\underline{b})$. Thus, equation 5.1-9 states that the gradient of $S(\underline{b})$ is zero at a minimum point of $S(\underline{b})$.

From here on, for simplicity the approximate equality sign in regression models and normal equations (for example, 5.1-1, 5.1-4 through $5.1-8$ ) is replaced by an equal sign. However,
remember that all relationships derived by using the linearized model are approximate for a nonlinear model.

### 5.2 Assumptions of Regression Modeling

Some of the assumptions listed below have already been mentioned; they are discussed more completely here.

1. A true model exists:

$$
\begin{equation*}
Y=f\left(\xi_{1}, \xi_{2}, \ldots, \xi_{k} ; \beta\right)+\epsilon \tag{5.2-1}
\end{equation*}
$$

The model response, $Y$, consists of two parts, a deterministic part, $f$, and an additive stochastic part, $\epsilon$.
2. The disturbances, $\epsilon$, have the following properties

$$
\begin{align*}
& E(\underline{\epsilon})=\underline{0}  \tag{5.2-2}\\
& \operatorname{Var}(\underline{\epsilon})=\underline{V} \sigma^{2} \tag{5.2-3}
\end{align*}
$$

where the structure, $\underline{V}$, of the variancecovariance matrix $V \sigma^{2}$ is assumed to be symmetric positive definite and known. Alternative forms for equations $5.2-2$ and $5.2-3$ are obtained by premultiplying equation $5.2-2$, and pre- and postmultiplying equation $5.2-3$, by $\underline{V}^{-1 / 2}$ to obtain

$$
\begin{align*}
& E\left(\underline{V}^{-1 / 2} \underline{\epsilon}\right)=\underline{0}  \tag{5.2-4}\\
& \operatorname{Var}\left(\underline{V}^{-1 / 2} \underline{\epsilon}\right)=\underline{I} \sigma^{2} \tag{5.2-5}
\end{align*}
$$

The assumptions given by equations 5.2-2 through 5.2-5 indicate that $\underline{\varepsilon}$ is considered to be a vector of random variables with zero mean and variance-covariance matrix $V \sigma^{2}$. Furthermore, weighted disturbances $\underline{V}^{-1 / 2} \underline{E}$ have constant variance $I \sigma^{2}$ and are uncorrelated. To require the expected value of $\underline{\epsilon}$ to be zero is to require that equation $5.2-1$ be the true (or unbiased) model and to require in addition that $\epsilon$ be unbiased. Although imperfections in most physical theories prevent the former assumption from holding strictly true, a model should be constructed so that the absolute value of any $E\left(\epsilon_{j}\right)$ is as small as possible. From the practical
point of view, it is required that the bias not be significant. Criteria for this are developed later on.
The full form of $V$ is usually very difficult to obtain from the type of data usually available. However, if $\underline{V}$ is assumed to be diagonal so that there is no correlation among the $\epsilon_{j}$, then $\underline{V}$ can often be found by using graphical methods of analyzing residuals, to be discussed later on.
3. The matrices $\underline{\omega}$ and $\underline{V}^{-1 / 2}$ are equivalent; that is,

$$
\begin{equation*}
\underline{\omega}=\underline{V}^{-1} . \tag{5.2-6}
\end{equation*}
$$

For a linear model at least, the Gauss-Markov theorem (Beck and Arnold, 1977, p. 232-234) establishes that the variance of $b_{j}, \operatorname{Var}\left(\hat{b}_{j}\right)$, is a minimum if equation 5.2-6 is true. Furthermore, to compute $\operatorname{Var}(\underline{\hat{b}})$ correctly, whether or not equation 5.2-6 is true, $V$ must be known. Hence, assumption of another form for $\underline{\omega}$ (such as $\underline{\underline{I}}$, for example) would not avoid the problem of having to know $V$ to analyze the model. However, for a linear model, equation 5.2-6 is not essential to compute an unbiased estimate of $\underline{\beta}$. This fact may be demonstrated as follows. Solve equation 5.1-6 for $\underline{\hat{b}}-\underline{b}_{0}$ and take the expected value of it to obtain:.

$$
\begin{align*}
E\left(\underline{\hat{b}}-\underline{b}_{0}\right) & =\left(\underline{\underline{X}}^{T} \underline{\omega X}\right)^{-1} \underline{X}^{T} \underline{\omega} E\left(\underline{Y}-\epsilon_{0}\right) \\
& =\left(\underline{X}^{T} \underline{\omega X}\right)^{-1} \underline{X}^{T} \underline{\omega X}\left(\underline{\beta}-\underline{b}_{0}\right)=\underline{\beta}-\underline{b}_{0} \tag{5.2-7}
\end{align*}
$$

where the fact that $E(\underline{\epsilon})=0$ was used. From equation $5.2-7$ it is seen that

$$
\begin{equation*}
E(\underline{\hat{b}})=\underline{\beta} . \tag{5.2}
\end{equation*}
$$

If $\underline{V}$ is diagonal, then

$$
\underline{V}=\left[\begin{array}{ll}
1 / \omega_{1} 1 / \omega_{2} &  \tag{5.2-9}\\
& \ddots \\
& \\
& \\
& \\
& \\
& \\
& \\
& \\
\end{array}\right.
$$

where the double subscripts on $\omega$ have been replaced by single subscripts to indicate the diagonal nature of $V$.
4. The disturbances are normally distributed:

$$
\begin{equation*}
\underline{\epsilon} \sim N\left(\underline{0}, \underline{V} \sigma^{2}\right) \tag{5.2-10}
\end{equation*}
$$

or

$$
\begin{equation*}
V^{-1 / 2} \underline{\varepsilon} \sim N\left(\underline{0}, I I^{2} \sigma^{2}\right) . \tag{5.2-11}
\end{equation*}
$$

Assumption of either equation $5.2-10$ or $5.2-11$ is only necessary if investigations using the $F$ distribution are to be performed.
That $\underline{\epsilon}$ (or $\underline{\underline{V}}^{-1 / 2} \underline{\epsilon}$ ) be normally distributed implies that the elements of $\underline{V}^{-1 / 2} \underline{\epsilon}$ are neither systematic nor constant but are equally likely to be positive or negative. In addition, small errors are more frequent than large ones. Many types of models are subject to a number of sources of error, any one of which may or may not be normally distributed. However, in the case where a resultant error is the sum of a number of components, Central Limit Theorem implies that $\epsilon$ (or $\underline{V}^{-1 / 2} \epsilon$ ) could be normally distributed even if its component vectors were not.

Because $\epsilon$ and $\underline{\beta}$ are unknown, the assumptions (1 through 4) discussed cannot be checked directly. However, they may often be checked indirectly, which is a subject of model analysis.

### 5.3 Relationships Between Residuals and Disturbances

Many of the investigations involving the regression model are based either directly or indirectly on relationships between residuals $\underline{\hat{e}}$ and disturbances $\epsilon$. Residuals may be written in terms of disturbances by employing equations 5.1-6 and 5.1-7. First, equation 5.1-7 is written in the form

$$
\begin{equation*}
\underline{\hat{e}}=\underline{Y}-t_{0}-\underline{X}\left(\underline{\hat{b}}-\underline{b}_{0}\right) . \tag{5.3-1}
\end{equation*}
$$

Then equation $5.1-6$ is solved for $\hat{\underline{b}}-\underline{b}_{0}$ and substituted into equation $5.3-1$ to obtain

$$
\begin{align*}
\underline{\hat{e}} & =\underline{Y}-f_{0}-\underline{X}\left(\underline{X}^{T} \omega X\right)^{-1} \underline{X}^{T} \underline{\omega}\left(\underline{Y}-f_{0}\right) \\
& =\left(\underline{I}-X\left(\underline{X}^{T} \omega X\right)^{-1} \underline{X}^{T} \underline{\omega}\right)\left(\underline{Y}-f_{0}\right) . \tag{5.3-2}
\end{align*}
$$

If $\underline{b}_{0}$ is set equal to $\underline{\beta}$, then $\underline{Y}-\underline{f}_{0}=\underline{Y}-f_{\beta}=\underline{\epsilon}$ and equation 5.3-2 gives

$$
\begin{equation*}
\underline{\hat{e}}=\left(\underline{I}-\underline{X}\left(\underline{X}^{T} \omega \underline{\underline{\omega}}\right)^{-1} \underline{X}^{T} \underline{\omega}\right) \underline{\underline{E}} . \tag{5.3-3}
\end{equation*}
$$

It is frequently more convenient to work with weighted residuals, $\underline{\underline{1}}^{1 / 2} \underline{\hat{e}}$, and weighted disturbances, $\underline{\omega}^{1 / 2} \underline{\epsilon}$. In this case equation 5.3-3 becomes

$$
\begin{equation*}
\underline{\omega}^{1 / 2} \underline{\hat{e}}=\left(\underline{I}-\underline{\omega}^{1 / 2} \underline{X}\left(\underline{X}^{T} \underline{\omega} \underline{X}\right)^{-1} \underline{X}^{T} \underline{\omega}^{1 / 2}\right) \underline{\underline{\omega}}^{1 / 2} \underline{\underline{\epsilon}} . \tag{5.3-4}
\end{equation*}
$$

An interesting and useful property of the matrix appearing in either equation 5.3-3 or 5.3-4 is displayed, for equation $5.3-4$ for example, as follows:

$$
\begin{align*}
& {\left[\underline{\omega}^{1 / 2} \underline{X}\left(\underline{X}^{T} \underline{\omega X}\right)^{-1} \underline{X}^{T} \underline{\omega}^{1 / 2}\right]\left[\underline{\omega}^{1 / 2} \underline{X}\left(\underline{X}^{T} \underline{\omega} \underline{X}\right)^{-1} \underline{X}^{T} \underline{\omega}^{1 / 2}\right]} \\
& =\underline{\omega}^{1 / 2} \underline{X}\left(\underline{X}^{T} \underline{\omega}^{\underline{ }}\right)^{-1} \underline{X}^{T} \underline{\omega}^{1 / 2} \underline{\omega}^{1 / 2} \underline{X} \underline{X}^{T} \underline{\omega} \underline{X}^{-1} \underline{X}^{T} \underline{\omega}^{1 / 2} \\
& =\underline{\omega}^{1 / 2} \underline{X}\left(\underline{X}^{T} \underline{\omega X}^{-1} \underline{X}^{T} \underline{\omega}^{1 / 2} .\right. \tag{5.3-5}
\end{align*}
$$

In other words, the matrix times itself yields the original matrix. This result is true for both equations 5.3-3 and 5.3-4. For equation 5.3-4 it is also true that the matrix is symmetric (as can be seen in the derivation of equation 5.3-5), so that the matrix times its transpose yields the original matrix. This type of matrix is known as a symmetric idempotent matrix.

Another useful property can be derived based on the idempotency discussed above. For convenience let

$$
\begin{equation*}
\underline{R}=\underline{\omega}^{1 / 2} \underline{X}\left(\underline{X}^{T} \underline{\omega}^{\underline{X}}\right)^{-1} \underline{X}^{T} \underline{\underline{\omega}}^{1 / 2} \tag{5.3-6}
\end{equation*}
$$

Then because $\underline{R}$ is idempotent,

$$
\begin{align*}
& (\underline{I}-\underline{R})(\underline{I}-\underline{R})=\underline{I}-\underline{R}-\underline{R}+\underline{R} \\
& =\underline{I}-\underline{R} . \tag{5.3-7}
\end{align*}
$$

Hence, $\underline{\underline{I}-\underline{R}}$ is also idempotent.

### 5.4 Some Statistical Measures

The first step in model analysis should always be to examine some statistical measures that indicate (1) goodness of fit of the model to the data and (2) model conditioning as it affects reliability of the computed parameters. With the background given above, useful statistical measures can be derived.

### 5.4.1 The Error Variance, $s^{2}$

For a linear model this measure is an unbiased estimate of $\sigma^{2}$. For a nonlinear model, $s^{2}$ is biased. A sketch of the derivation follows: The sum of squares, $S(\underline{\hat{b}})$, is defined as

$$
\begin{align*}
S(\underline{\hat{b}}) & =\underline{e}^{T} \underline{\omega} \hat{e} \\
& =\left(\underline{\omega}^{1 / 2} \underline{\hat{e}}\right)^{T}\left(\underline{\omega}^{1 / 2} \underline{\hat{e}}\right) . \tag{5.4-1}
\end{align*}
$$

Using equations 5.3-4,5.3-7, and the fact that tr(scaler) $=$ scaler, equation $5.4-1$ becomes

$$
\begin{align*}
S(\underline{\hat{b}}) & =\left(\underline{\omega}^{1 / 2} \underline{\epsilon}\right)^{T}(\underline{\underline{I}}-\underline{R})^{T}(\underline{\underline{I}}-\underline{\underline{R}})\left(\underline{\omega}^{1 / 2} \underline{\epsilon}\right) \\
& =\operatorname{tr}\left[\left(\underline{\omega}^{1 / 2} \underline{\epsilon}\right)^{T}(\underline{\underline{I}-\underline{R}})\left(\underline{\omega}^{1 / 2} \underline{\epsilon}\right)\right] \\
& =\operatorname{tr}\left[(\underline{I}-\underline{\underline{R}})\left(\underline{\omega}^{1 / 2} \underline{\epsilon}\right)\left(\underline{\omega}^{1 / 2} \underline{\epsilon}\right)^{T}\right] . \tag{5.4-2}
\end{align*}
$$

The expected value of equation $5.4-2$ is

$$
\begin{align*}
E(S(\underline{\hat{b}})) & \left.=\operatorname{tr}\left\{(\underline{\underline{I}-R}-\underline{\underline{e}}) E\left[\underline{\omega}^{1 / 2} \underline{\epsilon}\right)\left(\underline{\underline{\omega}}^{1 / 2} \underline{\epsilon}\right)^{T}\right]\right\} \\
& =\operatorname{tr}\left[(\underline{\underline{I}}-\underline{R}) \operatorname{Var}\left(\underline{\underline{\omega}}^{1 / 2} \underline{\epsilon}\right]\right. \\
& =\operatorname{tr}\left[\left(\underline{\left.\underline{\underline{I}}-R) \sigma^{2}\right]}\right.\right. \\
& =(n-p) \sigma^{2} \tag{5.4-3}
\end{align*}
$$

from which

$$
\begin{equation*}
\sigma^{2}=\frac{E(S(\underline{\hat{b}}))}{n-p} \tag{5.4-4}
\end{equation*}
$$

The fact that $\operatorname{tr}(R)=p$ can be demonstrated by rearranging the matrices within $\underline{R}$. The estimate, $s^{2}$, of $\sigma^{2}$ is

$$
\begin{equation*}
s^{2}=\frac{S(\underline{\hat{b}})}{n-p}=\frac{(\underline{Y}-\hat{f})^{T} \underline{\omega}(\underline{Y}-\hat{f})}{n-p} \tag{5.4-5}
\end{equation*}
$$

A useful modification of equation 5.4-5 that is exact for a linear model and almost exact for a nonlinear one is obtained by choosing $\underline{b}_{0}$ in equation $5.1-7$ to be very near $\hat{b}$. Then

$$
\begin{array}{rlr}
s^{2} & =\frac{\left(\underline{Y}-f_{0}-\underline{X}\left(\underline{\hat{b}}-\underline{b}_{0}\right)\right)^{T} \underline{\underline{\omega}}\left(\underline{Y}-f_{0}-\underline{X}\left(\underline{\hat{b}}-\underline{b}_{0}\right)\right)}{n-p} \quad \text { This m } \\
& =\frac{\left(\underline{Y}-f_{0}\right)^{T} \underline{\omega}\left(\underline{Y}-f_{0}\right)-2\left(\underline{\underline{b}}-\underline{b}_{0}\right)^{T} \underline{X}^{T} \underline{\omega}\left(\underline{Y}-f_{0}\right)+\left(\underline{\hat{b}}-\underline{b}_{0}\right)^{T} \underline{X}^{T} \underline{\omega} X\left(\underline{\hat{b}}-\underline{b}_{0}\right)}{n-p}
\end{array}
$$

$$
\begin{equation*}
=\frac{\left(\underline{Y}-f_{0}\right)^{T} \underline{\omega}\left(\underline{Y}-f_{0}\right)-\left(\underline{\hat{b}}-\underline{b}_{0}\right)^{T} \underline{X^{T}} \underline{\omega}\left(\underline{Y}-f_{0}\right)}{n-p} \tag{5.4-6}
\end{equation*}
$$

where use was made of equation 5.1-6. For a linear model $\underline{b}_{0}$ may be chosen to be $\underline{0}$. For a nonlinear model a bias exists in equations 5.4-5 or 5.4-6 that results from the fact that development is based on assumption of a linear model.
Even when biased, $s^{2}$ gives a useful measure of overall goodness of fit of the model. The standard deviation or scatter is given by $s$. In general, $s / \Delta Y_{s}$ should be small, where $\Delta Y_{s}$ is the difference between maximum and minimum values of $Y_{s}$.

### 5.4.2 The Correlation, $\boldsymbol{R}_{\boldsymbol{y}}$, Between $\underline{\omega}^{1 / 2} \underline{Y}$ and $\underline{\omega}^{1 / 2} \underline{\underline{f}}$

This measure is defined as

$$
\begin{equation*}
R_{y}=\frac{\underline{d}_{y}^{T} \underline{d} \hat{y}}{\left[\left(\underline{d}_{y}^{T} \underline{d}_{y}\right)\left(\underline{d}_{f}^{T} \underline{d} f\right)\right]^{1 / 2}} \tag{5.4-7}
\end{equation*}
$$

where

$$
\begin{align*}
& \underline{d}_{y}=\underline{\omega}^{1 / 2} \underline{Y}-m_{y} \underline{l}  \tag{5.4-8}\\
& \underline{d} \hat{f}=\underline{\omega}^{1 / 2} \hat{\underline{L}} m \underline{f} \underline{l}  \tag{5.4-9}\\
& m_{y}=\sum_{i=1}^{n} \underline{\omega}_{i}^{1 / 2} \underline{Y} \mid n  \tag{5.4-10}\\
& m_{f}=\sum_{i=1}^{n} \omega_{i}^{1 / 2} \hat{f} \mid n \tag{5.4-11}
\end{align*}
$$

$\underline{l}=$ vector of ones,
and $\omega_{i}^{1 / 2}$ is row $i$ of $\omega^{1 / 2}$. The correlation $R_{y}$ is another measure of goodness of fit. Usually it should be greater than about 0.9 to indicate a good fit to the reliable data.

### 5.4.3 The Variance-Covariance Matrix for $\underline{\hat{b}}$

This measure may be derived directly from equation 5.1-6 and is

$$
\begin{align*}
\operatorname{Var}(\underline{\hat{b}}) & =\operatorname{Var}\left[\left(\underline{X}^{T} \underline{\omega X}\right)^{-1} \underline{X}^{T} \underline{\omega}\left(\underline{Y}-f_{0}\right)+\underline{b}_{0}\right] \\
& =\left(X^{T} \underline{\omega}\right)^{-1} X^{T} \underline{\omega} \operatorname{Var}(\underline{Y}) \omega X\left(X^{T} \omega X\right)^{-1} \\
& =\left(X^{T} \omega X\right)^{-1} \sigma^{2} \tag{5.4-12}
\end{align*}
$$

where assumptions from equations $5.2-3$ and 5.2-6 and the fact that $\operatorname{Var}(\underline{Y})=\operatorname{Var}\left(\underline{Y}-f_{\beta}\right)=\omega^{-1} \sigma^{2}$ were employed. The estimate of $\operatorname{Var}(\underline{b})$ is

$$
\begin{equation*}
\operatorname{Var}(\underline{\hat{b}})=\left(\underline{X}^{T} \underline{\underline{\omega X}}\right)^{-1} s^{2} \tag{5.4-13}
\end{equation*}
$$

The standard error of the estimate for the $i$ th parameter is given by the square root of the $i$ th diagonal component of $\left(\underline{X}^{T} \underline{\omega X}\right)^{-1} s^{2}$. This estimate is a measure of the range over which the respective parameter may be varied to produce a similar solution for the dependent variable as that obtained using $\underline{b}$.

### 5.4.4 The Correlation, $r_{1 j}$, Between any Two Parameters $\hat{b}_{\text {, }}$ and $\hat{b}_{\text {j }}$

By definition

$$
\begin{equation*}
r_{i j}=\frac{\operatorname{Cov}\left(\hat{b}_{i}, \hat{b}_{j}\right)}{\left[\operatorname{Var}\left(\hat{b}_{i}\right) \operatorname{Var}\left(\hat{b}_{j}\right)\right]^{1 / 2}} \tag{5.4-14}
\end{equation*}
$$

where the variance and covariance terms are components of either $\left(X^{T} \omega X\right)^{-1} s^{2}$ or $\left(X^{T} \omega X\right)^{-1} \sigma^{2}$. This measure gives an estimate of the degree of linear dependence of one parameter on another throughout the course of repeated experiments if such experiments were to be carried out. As discussed earlier, it is an indication of the degree of linear dependency in the sensitivity matrix.

## Problem 5.4-1

This problem is concerned with preliminary analysis of the linear regression solution of problem 3.2-1. The measures (except $s / \Delta Y_{s}$ ), to be computed in $a, b$, and c below, also are calculated by the computer model of problem 4.2-1. Check your computations against the computergenerated results.
a. Using equation 5.4-6, compute $s^{2}$. Compute $s / \Delta Y_{s}$. Would you say that the fit is very good?
b. Using equation 5.4-13, compute $\operatorname{Var}(\underline{\hat{b}})$. Are the parameters determined very precisely?
c. Using equation 5.4-14, determine $r$, the correlation matrix. Are there any evident problems with conditioning?

### 5.5 Analysis of Residuals

Examination of the statistics discussed in the previous section should give a preliminary indication of general model conditioning and model fit to the data. However, a thorough analysis of residuals is necessary in order to examine the validity of the assumptions given in section 5.2. Interest is focused primarily on indications of nonrandomness of the residuals and on indications that the residuals are not distributed normally. The analysis should include both sample and prior information partitions of the residuals so that any incompatibility between the two partitions can be detected as differences between the two partitions. Although the techniques given in the present section are usually adequate to detect any incompatibility, a formal test given in section 6.3 also can be applied if desired.

Analytical methods used here are graphical. Draper and Smith (1981, p. 141-192) give a number of methods for examining residuals, and they emphasize that graphical procedures involving visual analysis are the most valuable tools because violations of assumptions serious enough to require corrective action generally are apparent on the various plots. However, to use the procedures effectively it is necessary to determine the properties that the residuals should be expected to exhibit under ideal conditions.

### 5.5.1 Distribution of Residuals

Investigation of the distribution of weighted residuals $\underline{\omega}_{j}^{1 / 2} \underline{\hat{e}}$, where $\underline{\omega}_{j}^{1 / 2}$ is a row of $\underline{\omega}^{1 / 2}$, made in order to infer the distribution of weighted disturbances $\omega_{j}^{1 / 2} \underline{\epsilon}$, is difficult because, even if the assumption given by equation 5.2-5 holds so that the elements $\underline{\omega}_{j}^{1 / 2} \underline{\epsilon}$ are uncorrelated and have equal variance, elements of $\omega_{j}^{1 / 2} \underline{\hat{e}}$ are correlated and have unequal variance. To show this for the linearized model, equations 5.3-4 and 5.3-6 can be combined to give

$$
\begin{equation*}
\underline{\underline{\omega}}^{1 / 2} \underline{\hat{e}}=(\underline{(\underline{I}}-\underline{R}) \underline{\omega}^{1 / 2} \underline{\underline{E}} \tag{5.5-1}
\end{equation*}
$$

from which

$$
\begin{align*}
E\left(\underline{\omega}^{1 / 2} \underline{\hat{e}}\right) & =(\underline{\underline{I}}-\underline{R}) \underline{\omega}^{1 / 2} E(\underline{\epsilon}) \\
& =\underline{0} \tag{5.5-2}
\end{align*}
$$

and

$$
\begin{align*}
\operatorname{Var}\left(\underline{\omega}^{1 / 2} \hat{e}\right) & =\operatorname{Var}\left[(\underline{I}-\underline{R}) \underline{\underline{1}}^{1 / 2} \underline{\underline{E}}\right] \\
& =(\underline{I}-\underline{R}) \operatorname{Var}\left(\underline{( }^{1 / 2} \underline{\epsilon}\right)(\underline{I}-\underline{R})^{T} \\
& =(\underline{I}-\underline{R})(\underline{I}-\underline{R})^{T} \sigma^{2} \\
& =(\underline{I}-\underline{R}) \sigma^{2} . \tag{5.5-3}
\end{align*}
$$

Hence, if assumptions given by equations 5.2-4 through 5.2-6, and 5.2-11 hold,

$$
\begin{equation*}
\underline{\omega}^{1 / 2} \underline{\hat{e}}=\underline{\hat{u}} \sim N\left(\underline{0},(\underline{I}-\underline{R}) \sigma^{2}\right) \tag{5.5-4}
\end{equation*}
$$

where for convenience, the definition is made that $\underline{\omega}^{1 / 2} \underline{\hat{e}}=\underline{\hat{u}}$.

It can be shown that $\underline{I}-\underline{R} \rightarrow \underline{I}$ as $n-p \rightarrow \infty$. Whenever $n-p$ becomes small, correlation and unequal variance become significant. Most tests for distribution of residuals assume equal (or, a common) variance and no correlation because all residuals are assumed to have come from the same univariate distribution. Therefore, correlation and unequal variance of the $\hat{u}_{j}$ are serious problems with regard to testing for normality when the number of parameters is not small compared to the number of observations.

Another difficulty concerns the determination of whether or not the model fits the data. If the model fits the data and correlation of the values of $\hat{u}_{j}$ is not significant, then these residuals should appear to be nearly random. However, if correlation is significant, then the correlation will be reflected in the residual values. Patterns could develop in some of the plots (to be discussed), and these patterns could be mistakenly attributed to lack of model fit.

### 5.5.2 Graphical Procedures

The first step in using graphical procedures is to develop a control group. Several sets of
simulated residuals distributed as in equation 5.5-4 form the control group. These sets are then compared graphically with the true weighted residuals $\hat{\hat{u}}$ to help decide whether the distribution of $\underline{\hat{u}}$ differs to a visually detectable extent from a normal distribution and whether correlation could cause an apparently nonrandom (or non-normal) pattern of residuals to develop in the residual plots.

A set of simulated residuals may be generated by generating a set of uncorrelated random normal deviates $\underline{d}$ so that $E(\underline{d})=\underline{0}$ and $\operatorname{Var}(\underline{d})$ $=I s^{2}$, then forming linear combinations of these deviates that have the covariance given by equation $5.5-3$. The method of generating the simulated residuals from the uncorrelated random deviates can be derived as follows. Assume, as a working hypothesis, that

$$
\begin{equation*}
g=\underline{\Omega} d \tag{5.5-5}
\end{equation*}
$$

where $g$ is the set of simulated residuals, and $\underline{\Omega}$ is a symmetric and nonstochastic matrix to $\bar{b} e$ determined. Vector $g$ must be generated so that $E(g)=\underline{0}$ and $\operatorname{Var}(\underline{g})=(\underline{I}-\underline{R}) s^{2}$ (equations 5.5-2 and 5.5-3). From equation 5.5-5

$$
\begin{align*}
E(\underline{g}) & =\underline{\Omega} E(\underline{d}) \\
& =\underline{0} \tag{5.5-6}
\end{align*}
$$

as required. By definition

$$
\begin{align*}
\operatorname{Var}(\underline{g}) & =\underline{\underline{\Omega}} \operatorname{Var}(\underline{d}) \underline{\underline{\Omega}}^{T} \\
& =\underline{\Omega}^{2} s^{2} \tag{5.5-7}
\end{align*}
$$

where the definition of $\operatorname{Var}(\underline{d})$ and the symmetry of $\underline{\underline{\Omega}}$ were used. Hence, $\underline{\underline{\Omega}}$ must be defined so that

$$
\begin{equation*}
\underline{\mathbf{\Omega}}^{2}=\underline{I}-\underline{R} . \tag{5.5-8}
\end{equation*}
$$

However, because $I-\underline{R}$ is idempotent, ( $(\underline{I}-\underline{R})$ $=(I-R)^{2}$, and equation $5.5-8$ may be simplified to become

$$
\begin{equation*}
\underline{\Omega}=\underline{I}-\underline{R} \tag{5.5-9}
\end{equation*}
$$

so that equation 5.5-5 assumes as its final form

$$
\begin{equation*}
g=(\underline{I}-\underline{R}) \underline{d} . \tag{5.5-10}
\end{equation*}
$$

To generate a set of simulated residuals, $g$, it is a simple matter to generate a set of uncorrelated random normal deviates, $\underline{d}$, then use equation 5.5-10. This procedure is followed for the number of sets (usually at least three) desired to form the control group.

Normal probability plots.-These are graphs of cumulative frequency, $F$, versus values of the elements of vectors $\underline{\hat{u}}, \underline{d}$, or $\underline{g}$. Cumulative frequency corresponding to the $i$ th element of one of the vectors ( $\underline{\hat{u}}$, for example) is computed from the formula

$$
\begin{equation*}
F_{i}=m_{i} /(n+1), i=1,2, \ldots, n \tag{5.5-11}
\end{equation*}
$$

where $m_{i}$ is the number of values of $\underline{\hat{u}}$ (for example) smaller than or equal to $\hat{u}_{i}$, and $n$ is the number of observations. Use of $n+1$ in the denominator adjusts for the fact that $F_{n}$ cannot be equal to 1 because the tail of the normal distribution extends to infinity. An example of a normal probability plot for $\underline{\hat{u}}$ is illustrated in figure 5.5-1.

To determine the effects of correlation and unequal variance, normal probability plots for the sets $\underline{d}$ and the sets $g$ may be compared. If the plots for $\underline{d}$ and $g$ are very similar, then a normal probability plot of $\underline{\hat{u}}$ would not be expected to be affected to a great extent by

correlation and unequal variance. Other types of plots involving $\underline{\hat{u}}$ probably would not be affected much by correlation and unequal variance either. Significant departures of the plots for $g$ from those for $\underline{d}$ suggest serious correlation and unequal variance effects in $\underline{\hat{u}}$ also.

Whether or not $\underline{\hat{u}}$ departs to a visually detectable extent (which is considered to be synonymous with significarit here) from a normal distribution can be determined by comparing plots for the generated set $g$ with the plot for $\underline{\hat{u}}$. If the plot for $\underline{\hat{u}}$ has a trend similar to the set of curves for $g$, then the distribution for $\underline{\hat{u}}$ probably does not differ enough from equation 5.5-4 to consider abandoning the normality assumption. When examining the plots, it must be remembered that, because of correlation and unequal variance, the plots will not necessarily exhibit the linear trend expected for a univariate normal distribution.

Other residual plots.-In the following discussion it is assumed that the effects of correlation and unequal variance resulting from equation 5.5-3 are negligible so that other effects may be investigated. This assumption might hold true even if a normal probability plot is affected by correlation and unequal variance. However, if one or more patterns (or trends) appear to be present in one or more of the residual plots, then analogous plots using $g$ instead of $\hat{\underline{u}}$ also should be prepared and examined. If the suspicious patterns also are present in the plots using $g$, then the patterns probably result from correlation and unequal variance inherent in $(\underline{I}-R) \mathrm{s}^{2}$ and not from model error.
$=$ Three types of plot are often useful: (1) Plot of $\hat{u}_{j}$ vs. $\hat{f}_{j}$; (2) plots of $\hat{u}_{j}$ vs. independent variables $\left(\xi_{i}\right) ;(3)$ plot of $\hat{u}_{j}$ vs. Cartesian coordinates of point $j$. If $\underline{\omega}=\underline{I}$ was employed in the regression, then $\hat{u}_{j}=\overline{\hat{e}}_{j}$. Additional discussion of the first two types of plots may be found in Draper and Smith (1981, p. 147-148).

1. Plot of $\hat{u}_{j}$ vs. $\hat{f}_{j}$. This type of plot is illustrated in figure 5.5-2. Under the given assumptions, the plot should display a roughly horizontal band of residuals having no apparent trend as sketched in figure 5.5-3. In this and succeeding figures $5.5-4$ through 5.5-6, the dashed line outlines the limits of the data. A standard runs test (Draper and Smith, 1981,


Figure 5.5-2


Figure 5.5-3
p. 157-162) could be used to test for randomness of signs of $\hat{u}_{j}$ along the $\hat{f}_{j}$ axis.

Three principal types of abnormalities in the plot of $\hat{u}_{j}$ vs. $\hat{f}_{j}$ are often apparent:
a. Unequal band width (figure 5.5-4).

This type of trend (or one opposite to it) generally indicates that the variance of $\hat{u}_{j}$ is not constant. In figure 5.5-4 the observations would appear to be less reliable as $\hat{f}_{j}$ increases.

If the abnormal plot resulted from a least squares analysis where $\omega=I$ had been assumed, then a diagonal form of $\underline{\omega} \neq \underline{\underline{\omega}}$ might be indicated. In the illustration, $\omega_{j}$ should decrease with $\hat{f}_{j}$. However, if some form of $\omega \neq I$ had originally been assumed, then an abnormal plot, of the form shown in figure 5.5-4, involving $\hat{u}_{j}$ would suggest the $\hat{u}_{j}$ are not of equal reliability and that $\omega$ is not correct. Hence, $\omega$ should be modified, $\overline{\text { and }}$ the regression performed again. If a full form of $\omega$ is required by the true model,


Figure 5.5-4
this form cannot generally be discovered by this type of analysis. Before a weighted least squares is performed, the underlying cause of the variable reliability problem should be investigated so that the diagonal weights can be added according to a rational criterion.
b. Sloping band (figure 5.5-5). This type of problem often indicates model error. Typically, it is caused by omitting an intercept from the model. For a linear model having an intercept, $\Sigma \hat{u}_{j} \hat{f}_{j}=0$ always. Hence in this case an overall slope such as depicted in figure 5.5-5 cannot result from correlation and unequal variance in $(I-R) s^{2}$. Also, if the model is linear (or effectively so) and $E(\underline{\epsilon})=\underline{0}$, then $\operatorname{Cov}(\hat{\underline{u}}, \hat{\varphi}=\underline{0}$ even if there is no intercept. Thus, it might be expected that a plot such as shown in figure 5.5-5 would usually not result from correlation and unequal variance if the model is correct and if the observed data are adequate.
c. Curved or irregularly shaped band (figure 5.5-6). This is another indication of model error. The model is inadequate because it does not account for all of the sources of variability in the observed data. More, or different, terms should be added to the regression equation. For the plot shown, it is possible that a quadratic term should be added.
2. Plots of $\hat{u}_{j}$ vs. independent variables ( $\xi_{i}$ ). These plots are interpreted in much the same way as the first type of plot; they simply provide a different viewpoint.
3. For trend surface types of regression: plot of $\hat{u}_{j}$ vs. Cartesian coordinates of point $j$ (for one or two-dimensional systems). Threedimensional systems can sometimes be reduced to two dimensions by using cross sections. Systematic highs and lows in the residual pattern usually suggest that the model does not fit the data well. In general, the easier it is to contour the residual map, the more nonrandom the residuals are likely to be.

## Problem 5.5-1

To the program for computing $T$ and $S$ using the Theis equation (problem 3.3-1), add the code necessary to compute $s^{2}$ using equation 5.4-5 and $\operatorname{Var}(\underline{b})$ using equation 5.4-13. You will have to unscale the entries of $\left(Z^{T} Z\right)^{-1}$ to obtain $\widehat{\operatorname{Var}}(\underline{\hat{b}})$. Examination of the calculation


Figure 5.5-5


Figure 5.5-6
procedures for obtaining $\underline{Z}^{T} \underline{Z}$ and $\left(Z^{T} \underline{Z}^{-1}\right.$ will indicate how the unscaling should be accomplished.

Conduct a graphical analysis of residuals $\underline{\hat{e}}$ resulting from the Theis equation problem (problem 3.3-1). Compute five sets of random normal deviates $\underline{d}$ and simulated residuals $g$ using the residuals analysis program (appendix 5.8.1). Plot $\underline{d}$ and $g$ on normal probability paper. Does the plot of $\hat{e}$ differ significantly from the plots of $g$ ? What can you conclude about the distribution of $\underline{e}$ ? Does the Theis model appear to be adequate?

The code given in appendix 5.8 .1 is designed to read $\operatorname{COV}(\mathrm{I}, \mathrm{J})$ (which is $\left.\left(X^{T} \omega X\right)^{-1} s^{2}\right), W(I)$ (which is $\omega_{s}$ ) and $X(I, J)$ (which is $X_{s}$ ) in unformatted form from a file labeled ITB. Normally, these data would be read to file ITB from the numerical nonlinear regression program of appendix 4.3.4. However, the Theis data are not obtained from the numerical nonlinear regression program. The easiest way to read the Theis data is to modify the residuals analysis program by replacing READ(ITB) in statements reading data sets B, C, and D with READ(IIN,2). Data for the program should then be coded as explained in appendix 5.8.1, except that the data for data sets $B, C$, and $D$ will now be coded in format 8F10.0.

## Problem 5.5-2

Use the output from the two-dimensional flow model to analyze the model results of problem 4.2-2.
a. What is the value of $s^{2}$ ? $R_{y}$ ? Compute $s / \Delta Y_{s}$.
b. Is there any evidence of ill-conditioning in the results? If there is, what is the problem? Which parameters are welldetermined and which are not? Why?
c. Conduct a graphical analysis of residuals. Develop four sets of random normal deviates $\underline{d}$ and simulated residuals $g$, computed by using the residuals analysis program (appendix 5.8.1). Plot $\underline{d}$ and $g$ on normal probability paper using equation 5.5-11. Are correlation effects evident? Plot $\underline{\hat{u}}$ (why $\underline{\hat{u}}$ instead of $\hat{e}$ ?) on normal probability paper. Does the plot differ significantly from the plots of $g$ ? Plot $\hat{u}_{j}$ versus $\hat{f}_{j}$. Is there an abnormal pattern? Plot $\hat{e}_{j}$ versus Cartesian coordinates of point $j$ (omitting the prior information). Again, is there a pattern to the residuals? What do you conclude about the adequacy of the model?

### 5.6 Investigation of Alternative Parameter Sets

### 5.6.1 Generalized W Statistic

Suppose we want to test the null hypothesis that some subset $\underline{\beta}_{2}$ of parameter set $\underline{\beta}$ cannot be distinguished from some corresponding given subset $\underline{\underline{\beta}}_{2}$. That is, test

$$
H_{0}: \underline{\beta}_{2}=\tilde{\underline{\beta}}_{2} \text { versus } H_{1}: \underline{\beta}_{2} \neq \tilde{\underline{\beta}}_{2}
$$

The linearized model assumed is

$$
\begin{align*}
\underline{Y} & =f_{0}+\underline{X}\left(\underline{\beta}-\underline{b}_{0}\right)+\underline{\epsilon} \\
& =f_{0}+\underline{X}_{1}\left(\underline{\beta}_{1}-\underline{b}_{01}\right)+\underline{X}_{2}\left(\underline{\beta}_{2}-\underline{b}_{02}\right)+\underline{\epsilon} \tag{5.6-1}
\end{align*}
$$

where $X, \underline{\beta}$, and $\underline{b}_{0}$ are conformably partitioned as follows:

$$
\begin{equation*}
\underline{X}=\left[\underline{X}_{1}, \underline{X}_{2}\right] \tag{5.6-2}
\end{equation*}
$$

$$
\begin{align*}
& \underline{\beta}=\left[\begin{array}{l}
\underline{\beta}_{1} \\
\underline{\beta}_{2}
\end{array}\right] \\
& \underline{b}_{0}=\left[\begin{array}{l}
\underline{b}_{01} \\
\underline{b}_{02}
\end{array}\right] \tag{5.6-3}
\end{align*}
$$

Based on equation 5.6-1, we may state a predictive model of the form

$$
\begin{equation*}
\tilde{t}=f_{0}+\underline{X}_{1}\left(\underline{\underline{b}}_{1}-\underline{b}_{01}\right)+\underline{X}_{2}\left(\tilde{\vec{\beta}}_{2}-\underline{b}_{02}\right) \tag{5.6-4}
\end{equation*}
$$

where $\underline{\underline{b}}_{1}$ is an estimate of $\underline{\beta}_{1}$ and, under the null hypothesis, $\underline{\beta}_{2}$ is assumed to be given by $\tilde{\beta}_{2}$.

The $W$ statistic, which is stated explicitly later on, is based on a comparison of the restricted sum of squares $(\underline{Y}-\tilde{\tilde{f}})^{T} \omega(\underline{Y}-\tilde{f})$ and the unrestricted sum of squares $(Y-\hat{f})^{T} \omega(Y-\hat{f})$. The unrestricted sum of squares is obtained from the standard least squares analysis. The restricted sum of squares is obtained by minimizing

$$
\begin{equation*}
S_{0}\left(\tilde{\underline{b}}_{1}\right)=(\underline{Y}-\tilde{f})^{T} \underline{\omega}(\underline{Y}-\tilde{f}) \tag{5.6-5}
\end{equation*}
$$

with respect to $\underline{\tilde{b}}_{1}$ while holding $\tilde{\underline{\beta}}_{2}$ constant. For a linear model this results in the normal equations:
$\underline{\underline{\tilde{b}}}_{1}-\underline{\underline{b}}_{01}=\left(\underline{X}_{\underline{X}}^{T} \omega X_{1}\right)^{-1} \underline{X}_{\underline{X}}^{T} \omega\left(\underline{Y}-\underline{f}_{0}-\underline{X}_{2}\left(\tilde{\underline{B}}_{2}-\underline{b}_{02}\right)(5.6-6)\right.$
where $\underline{b}_{01}$ and $\underline{b}_{02}$ can be set to zero if desired. If the model is nonlinear, then equation 5.6-5 is minimized with respect to $\tilde{\underline{b}}_{1}$ by using the standard procedures, except that $\tilde{\underline{\beta}}_{2}$ is held constant.

If many repeat experiments using linear normal equation 5.6-6 were performed, and if assumptions given by equations $5.2-5,5.2-6$, and 5.2-10 held true, then it would be found that

$$
\begin{align*}
& W=\frac{\left.\left[(\underline{Y}-\tilde{f})^{T} \underset{\underline{\omega}}{ }(\underline{Y}-\tilde{f})^{-(\underline{Y}-\hat{f}}\right)^{T} \underset{\underline{\omega}}{ }(\underline{Y}-\hat{f})\right] / q}{(\underline{Y}-\hat{f})^{T} \underline{\omega}(\underline{Y}-\hat{f}) /(n-p)} \\
& =\frac{\left[(\underline{Y}-\tilde{f})^{T} \underset{\omega}{\omega}\left(\underline{Y}-\tilde{\hat{t}}-(\underline{Y}-\hat{f})^{T} \underline{\omega}(\underline{Y}-\hat{f})\right] / q\right.}{s^{2}} \\
& \sim F(q, n-p) \tag{5.6-7}
\end{align*}
$$

where $q$ is the order of $\tilde{\underline{\beta}}_{2}$, which is the number of restrictions in $H_{0}$, and equation 5.4-5 was used. The symbol $W$ stands for a random variable. For each experiment, a value $w$ of the random variable $W$ could be computed. According to equation $5.6-7$, by repeating the experiment many times, the probability of the ratio $W$ having a value of $w$ or a smaller value would be found to be given by the cumulative density function $F$.

Because equation 5.6-7 is proportional to the difference between the restricted and unrestricted sums of squares divided by the unrestricted sum of squares, one might suspect the null hypothesis $H_{0}: \underline{\beta}_{2}=\tilde{\underline{\beta}}_{2}$ to be true if $w$ is small. However, if $w$ is large, then one might suspect that $H_{0}$ is incorrect. The rejection region for the hypothesis test is determined by the probability statement $P\left(W>F_{\alpha}(q, n-p)\right)=\alpha$, where $\alpha$ is the significance level of the test and $F_{\alpha}(q, n-p)$ is the upper $100 \alpha \%$ point of the $F$ distribution with $q$ and $n-p$ degrees of freedom. If the ratio $w$ is greater than $F_{\alpha}(q, n-p)$, as found in any table of critical values for the $F$ distribution, then the null hypothesis is rejected because values of $w$ greater than $F_{\alpha}(q, n-p)$ form the rejection region.

An alternative form for the numerator of equation $5.6-7$ may be derived by manipulating the linearized models. The result, after extensive algebra, is

$$
\begin{align*}
& (\underline{Y}-\tilde{\tilde{f}})^{T} \omega(\underline{\underline{Y}}-\tilde{f})-(\underline{Y}-\hat{f})^{T} \omega(\underline{\underline{Y}}-\hat{f}) \\
= & \left.\left(\tilde{\underline{\beta}}_{2}-\hat{\underline{b}}_{2}\right)^{T} \underline{\underline{H}}\left(\underline{\underline{X}}^{T} \underline{\omega X}\right)^{-1} \underline{\underline{H}}^{T}\right]^{-1}\left(\tilde{\underline{\beta}}_{2}-\hat{\underline{b}}_{2}\right) \tag{5.6-8}
\end{align*}
$$

where

$$
\begin{equation*}
\underline{\underline{H}}=[\underline{\underline{0}}, \underline{\underline{I}}]_{(q \times p)}, \tag{5.6-9}
\end{equation*}
$$

$\underline{I}_{q}=$ identity matrix of order $q$, and
$\hat{b}_{2}=$ the partition corresponding to $\underline{\beta}_{2}$ found from the standard (unrestricted) least squares analysis.
(5.6-10)

Thus, the alternative form for equation $5.6-7$ is

$$
\begin{equation*}
W=\frac{\left(\underline{\underline{\beta}}_{2}-\underline{\underline{b}}_{2}\right)^{T}\left[\underline{\underline{H}}\left(\underline{\underline{X}}^{T} \underline{\underline{\omega} X}\right)^{-1} \underline{\underline{H}}^{T}\right]^{-1}\left(\tilde{\underline{\beta}}_{2}-\hat{\underline{b}}_{2}\right) / q}{s^{2}} . \tag{5.6-11}
\end{equation*}
$$

If the model is linear, either equation 5.6-7 or equation 5.6-11 may be used to compute $w$. Both equations require an unrestricted regression to obtain $\hat{b}$ and $\hat{f}$, but equation 5.6-7 requires, in addition, a restricted regression using equation 5.6-6 to obtain $\tilde{\underline{b}}_{1}$ and $\tilde{f}$. Hence, for a linear model, equation $5.6-11$ is often more efficient to use than equation 5.6-7 for practical computations. If equation $5.6-1$ is a linearized equation system, derived from a nonlinear model, then neither equation $5.6-7$ nor equations $5.6-8$ nor $5.6-11$ is exact. However, if equation $5.6-1$ behaves in a way that is close enough to being linear, then equations $5.6-7$, 5.6-8, and 5.6-11 are good approximations. When working with a nonlinear model, $w$ should be computed using both equations $5.6-7$ and $5.6-11$. If the conclusions reached by using the two different expressions for $w$ are different, then the model may be too nonlinear for investigations using the $W$ statistic. Further investigation of model nonlinearity may be performed by employing the modified Beale's measure, which is discussed in section 6.2.

In summary, the procedure for testing

$$
H_{0}: \beta_{2}=\tilde{\beta}_{2} \text { vs. } H_{1}: \tilde{\beta}_{2} \neq \tilde{\beta}_{2}
$$

in the model $\underline{Y}=f_{0}+\underline{X}\left(\underline{\beta}-\underline{b}_{0}\right)+\underline{\epsilon}$ is:

1. Carry out a regression without any restriction to find $\hat{b}$ for the full model.
2. When using equation 5.6-7, fix $\tilde{\underline{®}}_{2}$ and use equation 5.6-6 to find $\underline{\tilde{b}}$, for the restricted model. When using equation 5.6-11, skip this step.
3. Form the ratio $w$ using equation 5.6-7 or equation 5.6-11.
4. Compare $w$ with the appropriate value of $F_{\alpha}(q, n-p)$.

### 5.6.2 Joint Confidence Region for $\underline{\beta}_{\mathbf{2}}$

Equations 5.6-7 or 5.6-11 also may be used to obtain a joint confidence region on $\underline{\beta}_{2}$. The confidence region interpretation is based on fixing a probability of occurrence, $P\left(W<F_{\alpha}(q, n-p)\right)$ $=1-\alpha$, then finding those vectors $\tilde{\beta}_{2}$ that would yield the specified $F$ or a smaller value. The joint confidence region may be written as
$(\underline{Y}-\tilde{f})^{T}{ }_{\omega}(\underline{Y}-\tilde{f})-(\underline{Y}-\hat{f})^{T}{ }_{\omega}^{\omega}(\underline{Y}-\hat{f}) \leq s^{2} q F_{\alpha}(q, n-p)(5.6-12)$
or

$$
\begin{align*}
& \left(\underline{\beta}_{2}-\hat{b}_{2}\right)^{T}\left[\underline{H}\left(\underline{X}^{T} \underline{\omega X}\right)^{-1} \underline{H}^{T}\right]^{-1}\left(\underline{\beta}_{2}-\hat{b}_{2}\right) \\
& \leq s^{2} q F_{\alpha}(q, n-p) \tag{5.6-13}
\end{align*}
$$

Equation 5.6-13 plots as a family of $q$ dimensional ellipsoids (which for $q=1$ is a line segment bounded by two values for $\underline{\beta}_{2}$ ) in parameter space, and these ellipsoids are centered on $\hat{\hat{b}}$. All ellipsoids corresponding to probability levels smaller than 1- $\alpha$ lie within the outermost ellipsoid, which is defined by strict equality in equation 5.6-13. Hence, the specified probability is the probability that $\underline{\beta}_{2}$ lies within the ellipsoid. An equivalent statement is that, if many experiments were conducted, then ( $1-\alpha$ ) $100 \%$ of the ellipsoids would contain the true parameter set $\underline{\beta}_{2}$ Hence, the outermost ellipsoid may be considered to be a joint confidence region on $\beta_{2}$.

Points (in parameter space) on the edge of the confidence region corresponding to the maximuin and minimum (or extreme) values that some parameter $\tilde{\beta}_{2 i}$ may attain and remain in the confidence region are given by

$$
\begin{equation*}
\underline{\tilde{b}}=\underline{\hat{b}} \pm \frac{\sqrt{q F_{\alpha}(q, n-p)}}{s_{b i}} \underline{V}_{b i} \tag{5.6-14}
\end{equation*}
$$

where $\tilde{\tilde{b}}^{T}=\left[\tilde{b}_{1}^{T}, \tilde{B}_{2}^{T}\right]$ and $\underline{V}_{b i}$ is the $i$ th column of $V_{b}=\left(\bar{X}^{T} \underline{\omega}\right)^{-1} s^{2}$. Note that $s_{b i}$ is the square root of $v_{b i i}$.

Equation $i$ in 5.6-14 gives the extreme values of $\tilde{\beta}_{2 i}$. The parameter vector $\underline{\tilde{b}}$ computed using equation 5.6-14 is the same vector that would result if (1) $\tilde{\beta}_{2 i}$ were computed using equation $i$ in $5.6-14$, (2) then the remaining values in $\tilde{B}_{2}$ were computed to satisfy equation $5.6-13$ (with strict equality applying to give points on the edge of the confidence region), and (3) finally, partition $\tilde{b}_{1}$ were computed using equation 5.6-6.

Parameter sets computed using equation 5.6-14 will exactly satisfy the relationship $s^{2} q F_{\alpha}(q, n-p)=(\underline{Y}-\tilde{f})^{T} \omega(\underline{Y}-\tilde{f})-(\underline{Y}-\hat{f})^{T} \underline{\omega}(\underline{Y}-\hat{f})=\left(\tilde{\tilde{\beta}}_{2}\right.$ $\left.-\underline{b}_{2}\right)^{T}\left[\underline{H}\left(X^{T} \omega X\right)^{-1} \underline{H}^{T}\right]^{-1}\left(\underline{\tilde{\beta}}_{2}-\hat{\underline{b}}_{2}\right)$ if the model is linear. However, as discussed in section 5.6.1, this relationship is not exact for a nonlinear model. Therefore, if the regression model is nonlinear, the parameter sets computed using equation $5.6-14$ should be substituted into the
nonlinear model and $(\underline{Y}-\tilde{f})^{T} \omega(\underline{Y}-\tilde{f})-(\underline{Y}-\hat{f})^{T} \omega(\underline{Y}-\hat{f})$ should be computed. If this value is different enough from $s^{2} q F_{\alpha}(q, n-p)$ to change any conclusions, then the model is too nonlinear to use to generate linearized confidence regions. The modified Beale's measure discussed in section 6.2 also can be used to gauge nonlinearity.

Two end-member cases involving the $W$ statistic are often considered separately. In one case $q=p$ so that $\underline{\beta}_{2}=\underline{\beta}$. All parameters are thus included in any test of $H_{0}$, and the confidence region is on all parameters simultaneously. This confidence region is called a joint confidence region on all parameters. An example for two parameters is diagrammed in figure 5.6-1. In the other case, $q=1$ and $\underline{\beta}_{2}=\beta_{p}$. Thus, only one parameter is considered in any test of $H_{0}$, and the confidence region is on only one parameter. The confidence region for this case is termed an individual confidence interval on parameter $\boldsymbol{\beta}_{\boldsymbol{p}}$.

## Problem 5.6-1

a. A method for estimating recharge in Nevada is known as the Maxey-Eakin method. Using this method, recharge rate $W$ in the vicinity of Lake Ohpupu (problem 3.2-1) was estimated to be 0.0003 $\mathrm{ft} / \mathrm{day}$. Also, by using specific capacity estimates from well-log analysis, $T$ was found to be $10 \mathrm{ft}^{2} /$ day. Using these estimates, test the null hypothesis that there is no significant difference at $\alpha=0.05$ between $W / T$ as estimated above


Figure 5.6-1
and the regression estimate of $W / T$. (Hint: Use equation 5.6-11.) Based on the result of the hypothesis test, would you consider using the prior estimate of $W / T$ as prior information in the regression model? What other information would you need if you did use it?
b. In addition to asking whether or not regression estimate $W / T$ is significantly different from another, independent, estimate, one might ask whether or not $W / T$ is even a significant variable in the regression model. Set up and conduct a test to answer this question, then interpret the result.
c. Using equation 5.6-14, find the bounds of the confidence interval on $W / T$ for $\alpha=0.05$. (Hint: You need only consider the equation corresponding to parameter $W / T$ in the system implied by equation 5.6-14.)

## Problem 5.6-2

Use your Theis equation program (problems 3.3-1 and 5.5-1) and equation 5.6-14 to find the sets of parameters corresponding to extreme values of $T$ and to extreme values of $S$, assuming $q=2$.

## Problem 5.6-3

Using equation 5.6-14 and the results of problem 4.2-2, find the sets of parameters corresponding to extreme values of $T_{3}$ and to extreme values of $q_{B 1}$ with $q=2$.

### 5.7 Investigation of Predictive Reliability

### 5.7.1 The Variance-Covariance Matrix for $\hat{\mathbf{f}}$

Equation 5.1-8 is used to obtain

$$
\begin{align*}
\operatorname{Var}(\hat{f}) & =\operatorname{Var}\left[\underline{X}\left(\underline{\hat{b}}-\underline{b}_{0}\right)+f_{0}\right] \\
& =X \operatorname{Var}(\underline{\hat{b}}) \underline{X}^{T} . \tag{5.7-1}
\end{align*}
$$

An analogous measure for weighted values of $\hat{f}$ is

$$
\begin{align*}
\operatorname{Var}\left(\underline{\omega}^{1 / 2} \hat{\underline{t}}\right) & =\operatorname{Var}\left[\underline{\omega}^{1 / 2} \underline{X}\left(\hat{\underline{b}}-\underline{b}_{0}\right)+\underline{\omega}^{1 / 2} \hat{L}_{0}\right] \\
& =\underline{\omega}^{1 / 2} \underline{\underline{X}} \operatorname{Var}(\hat{\underline{b}}) \underline{X}^{T} \underline{\omega}^{1 / 2} \tag{5.7-2}
\end{align*}
$$

By using equations 5.3-6 and 5.4-12, equation 5.7-2 can be written

$$
\begin{equation*}
\operatorname{Var}\left(\underline{\underline{\omega}}^{1 / 2} \hat{f}\right)=\underline{\underline{R}} \sigma^{2} \tag{5.7-3}
\end{equation*}
$$

Estimates corresponding to equations 5.7-1, 5.7-2, and 5.7-3 are

$$
\begin{align*}
\widehat{\operatorname{Var}}(\hat{f}) & =\widehat{\underline{X} \operatorname{Var}}(\hat{b}) \underline{X}^{T}  \tag{5.7-4}\\
\widehat{\operatorname{Var}}\left(\underline{\underline{\omega}}^{1 / 2} \hat{f}\right) & =\underline{\underline{\omega}}^{1 / 2} \underline{\underline{X}} \widehat{\operatorname{Var}}(\underline{\hat{b}}) \underline{X}^{T} \underline{\underline{\omega}}^{1 / 2}  \tag{5.7-5}\\
& =\underline{\underline{\underline{D}}} s^{2} . \tag{5.7-6}
\end{align*}
$$

The standard error of $\hat{f}_{i}$ is given by the
 This estimate gives a measure of potential variability in $\hat{f}_{i}$ resulting from $\widehat{\operatorname{Var}(\hat{b})}$.

It is important to note that equations 5.7-1, $5.7-2,5.7-4$, and $5.7-5$ are valid for prediction vectors $\hat{f}$ having entries that are not necessarily at observation points. This fact may be understood by observing that prediction equation 5.1-8 used to derive $\operatorname{Var}(\hat{f})$ or $\operatorname{Var}\left(\underline{\omega}^{1 / 2} \hat{f}\right)$ (or their estimates) is valid for any set of points, not just observation points. However, entries in $\underline{\underline{X}}$ and, if either equation $5.7-2$ or $5.7-5$ is used, $\omega$ must be available for all points in $\hat{f}$. Matrix $\left(X^{T} \underline{\omega}\right)^{-1} s^{2}$ used for $\widehat{\operatorname{Var}(\underline{b})}$ is, of course, the standard one based on entries in $\underline{\underline{X}}$ only at the observation points.

### 5.7.2 Confidence Interval for $\boldsymbol{f}_{\boldsymbol{\beta}}$

If all parameters are allowed to vary over the confidence region given by equation $5.6-12$ or equation 5.6-13 with $q=p$, then the maximum and minimum values produced for $\tilde{f}_{j}$ form a corresponding confidence interval for $\tilde{f}_{j}$. In this case, because $q=p, \underline{\beta}_{2}=\underline{\beta}$ and $\tilde{f}_{j}=f_{\beta j}$. The resulting confidence interval for $f_{\beta j}$ is

$$
\begin{equation*}
f_{\beta j}=\hat{f}_{j} \pm \sqrt{p F_{\alpha}(p, n-p)} s_{y j} \tag{5.7-7}
\end{equation*}
$$

where

$$
\begin{equation*}
s_{y j}=s \sqrt{\underline{X}_{j}\left(\underline{\underline{X}}^{T} \underline{\underline{\omega X}}^{-1} \underline{X}_{j}^{T}\right.} \tag{5.7-8}
\end{equation*}
$$

and $\underline{X}_{j}$ is a row of a sensitivity matrix corresponding to the point $j$. The point $j$ need not correspond to an observation point.

Equation 5.7-8 gives a simultaneous confidence interval. That is, the probability is $1-\alpha$ that $f_{\beta j}$ lies within the interval indicated by equation 5.7-7, and that $f_{\beta}$ for all other possible points lies within similar intervals simultaneously. If intervals on a number of consecutive points are computed and plotted, the result is called a confidence band.

### 5.7.3 Prediction Interval for Predicted Observation $Y_{j}^{\text {pred }}$

Equation 5.7-7 gives a confidence interval on a computed value $f$ or, in other words, the mean of $Y$, which is a fixed, nonrandom quantity. In some instances, a corresponding interval on a predicted observation, which is a random quantity, is desired, and this interval is termed a prediction interval. Prediction intervals on $k$ predicted values of $Y$ simultaneously can be readily computed if $\omega=V^{-1}$ is diagonal, and they are given by Lieberman (1961):

$$
\begin{align*}
Y_{J}^{\text {pred }}=\hat{f}_{j} \pm \sqrt{k F_{\alpha}(k, n-p)} & \sqrt{s^{2} / \omega_{j}+s_{y j}^{2}}, \\
& j=1,2, \ldots, k \tag{5.7-9}
\end{align*}
$$

As in equation 5.7-7, point $j$ need not (and, in general, probably would not) correspond to an observation point. However, $\omega_{j}$ for the prediction point has to be known.

The term $s^{2} / \omega_{j}+s_{y j}^{2}$ is the total variance in predicted observation $Y_{j}^{p r e d}$. This form results because $Y_{j}^{p r e d}=\left(Y_{j}^{p r e d}-\hat{f}_{j}\right)+\hat{f}_{j}$, where $Y_{j}^{p r e d} \hat{f}_{j}$ is statistically independent of $\hat{f}_{j}$, so that, as an estimate, $\operatorname{Var}\left(Y_{j}^{\text {pred }}\right)=\operatorname{Var}\left(Y_{j}^{\text {pred }}-\hat{f}_{j}\right)+\operatorname{Var}\left(\hat{f}_{j}\right)$ or $\operatorname{Var}\left(Y_{j}^{p r e d}\right)=s^{2} / \omega_{j}+s_{y j}^{2}$.

Equation 5.7-9 does not give prediction intervals on all $\underline{Y}_{\text {pred }}$ simultaneously. Furthermore, as $k$ increases, the prediction interval increases without bound. This result is because the normal distribution, which the errors in $Y$ are assumed to follow, has infinite tails. Hence, even though the probability of an error that is large in magnitude is small, as the number of
values of $Y$ considered simultaneously increases, the probability of an arbitrarily large error in at least one of them increases also. Usually the prediction interval is computed using $k=1$.

## Problem 5.7-1

a. Write out explicitly the form for general entry (ij) of $X\left(X^{T} \omega X\right)^{-1} X^{T}$, which is used in equation $\overline{5} \cdot \overline{7}-1$. For this exercise let ( $X^{T} \omega X^{-1}=A$ and write the result in terms of $\bar{A}$ to simplify the expression. Select and compute a diagonal entry of this matrix at an observation point used for the linear regression solution of problem 3.2-1. Note that you can replace $\underline{X}$ by $\underline{S}$ in $\underline{X}\left(X^{T} \omega X\right)^{-1} X^{T}$, and that this replacement leaves the result unaltered. Can you show this?
b. Using equation 5.7-4 and the diagonal entry computed in part $a$, determine $\operatorname{Var}\left(\hat{f}_{j}\right)$, where $j$ is the selected diagonal entry.
c. Using equation 5.7-7 and the results of b , find the confidence interval on $f_{\beta j}$.

### 5.8 Appendix

### 5.8.1 Documentation of Program to Compute Vectors $\underline{d}$ and $\underline{g}$ of Section 5.5.2.

This program computes vector $d$ of random normal deviates, vector $g$ of correlated normal deviates as defined by equation $5.5-10$, and other useful information related to the distributions of $\underline{\hat{e}}$ and $\underline{\hat{u}}$. Sample and direct prior information are assumed to be given in the form of equation 3.4-12. The sensitivity matrix for the direct prior information is theoretically of the form $\underline{\underline{X}}_{p}=\left[\underline{I}_{p}, \underline{0}\right.$. However, matrix $\left[\underline{I}_{p}, \underline{\underline{0}}\right]$ may be rearranged to conform with any parameter ordering.

The program was developed using the Microsoft Fortran Compiler, Version 3.3, with the DOS 2.0 operating system on an IBM PC/XT computer with the IBM 8088 Math Coprocessor and 256 KB memory. Except for the OPEN statements near the beginning of the
code, Fortran 66 was used throughout to make the code as machine independent as possible. The source code is contained in file RESAN.FOR in the diskette accompanying this report. A random $(0,1)$ number generator is employed as a function subroutine. This routine assumes an integer computer word length equal to at least $1,077,109,141$.
As coded, the contents of data sets $\mathrm{B}, \mathrm{C}$, and $D$ are assumed to be stored in unformatted form on file ITB $=8$. This is so that these sets do not have to be input manually. The listing appended contains the code (file RESINS.FOR on the diskette) and instructions for insertion into the program of appendix 4.3 .4 so that the required data are stored in the proper form and order for use in the present program. The user will have to supply the job control language necessary to store the data and retrieve them for use.
Two variables, NVD and NTD defined near the beginning of the program, must be redefined each time the dimensions of the program are changed. NVD must be set equal to dimensions of COV and the first dimension of S , all three of which are at least NVAR, and NTD must be set equal to the dimensions of $R$, which are at least NTOT $=$ NOBS + NPRIR.

Input Data.-Data Set A.
Problem size information; one line (format 5I5, F10.0).

| Line columns | Variable | Definition |  |
| :--- | :---: | :---: | :---: |
| $1-5$ | NVAR $\ldots \ldots$ | Number of parameters, <br> $p$. |  |
| 6-10 | NOBS | $\ldots$. | Number of sample obser- <br> vations, $n_{s}$. |
| $11-15$ | NPRIR $\ldots$. | Number of regression <br> parameters having direct <br> prior information, $n_{p}$. |  |
| 21-25-20 | NSETS | $\ldots$. | Number of sets of $\underline{d}$ and $g$ <br> vectors to be computed. |
| $26-35$ | NRAN $\ldots . .$Seed for random number <br> generator: any odd <br> number between 1 and |  |  |

Data Set B.
Covariance matrix, $\left(\underline{X}^{T} \underline{\omega}\right)^{-1} s^{2}$ (unformatted; stored in file ITB).

| Variable | Definition |
| :--- | :--- |
| $\operatorname{COV}(1,1)$ | Covariance matrix, entered se- <br> quentially from the diagonal |
| COV(2,1) element through NVAR for <br> each new regression parameter  |  |
| COV(NVAR,1) | number. Each new diagonal <br> COV(2,2) |
| $\quad$element begins a new record. |  |
| $\operatorname{COV}(N V A R, 2)$ |  |
| $\quad \vdots$ |  |
| COV(NVAR,NVAR) |  |

Data Set C.
Weight matrix for sample information, $\underline{V}_{s}^{-1}$ (unformatted; stored in file ITB).

| Variable | Definition |
| :--- | :---: |
| W(1) | Diagonal weight matrix for |
| W(2) | sample information, entered |
| $\vdots$ | sequentially from 1 through |
| W(NOBS) | NOBS. |

Data Set D.
Sensitivity matrix for sample information, $\underline{X}_{s}$ (unformatted; stored in file ITB).

| Variable | Definition |
| :---: | :---: |
| $\mathbf{X ( 1 , 1 )}$ | Sensitivity matrix for sample |
| $\mathrm{X}(2,1)$ | information, entered sequen- |
| $\vdots$ | tially 1 through NVAR for |
| X(NVAR,1) | each observation. Each new |
| X(1,2) | observation begins a new |
| $\vdots$ | record, for a total of NOBS |
| X(NVAR,2) |  |
| $\vdots$ |  |
| X(NVAR,NOBS) |  |

Data Set E.
Estimated error variance used with prior information of known reliability; one line (format F10.0).

| Line columns | Variable | Definition |
| :--- | :---: | :---: |
| $1-10$ | EV | The initial estimate of $\mathrm{s}^{2}$ used <br> in conjunction with prior <br> information. |

Omit data set if $\mathrm{NPRIR}=0$.

Data Set F.
Parameter numbers having prior information (format 16I5).

| Line columns | Variable | Definition |
| :--- | :---: | :--- |
| 1-5 | IPR(1) | Array subscript numbers <br> for regression parameters |
| 6-10 | IPR(2) | in b having prior infor- <br> mation. For use with |
| $\vdots$ | $\vdots$ | the regression ground- <br> water program, the <br> array subscript numbers <br> must be the subscript |
|  |  | numbers in the parameter <br> vector computed by that <br> program. |

Omit data set if NPRIR $=0$.
Data Set G.
Standard deviation matrix for prior information $\underline{U}^{1 / 2}$ (format 8F10.0).

| Line columns | Variable | Definition |
| :--- | :--- | :---: |
| $1-10$ | $F(1)$ | Diagonal standard devia- |
| $11-20$ | $F(2)$ | tion matrix for prior in |
| $\vdots$ | $\vdots$ | formation, entered in the |
|  | F(NPRIR) | same order as IPR(I) <br> from I through NPRIR. |

Omit data set if NPRIR $=0$.

Output.-Output is all clearly labeled. It is ordered as follows:

1. Data sets A through G.
2. Set number of vectors $d$ and $g$. Data for numbers 2 through 5 below are printed sequentially for each set.
3. Vector d. This vector is ordered from smallest to largest entry, and each entry is paired with its theoretical frequency as computed by using equation 5.5-11.
4. Vector g. Each entry is printed in its natural position corresponding to its position in a row or column of $\underline{R}$. Rows and columns of $R$ are ordered $\bar{b}$ y first sample observation numbers followed by prior information numbers, which are the subscripts I of IPR(I).
5. Vector g. This vector is ordered and paired with its theoretical frequency in the same way as $d$ is.
6. Covariance matrix $(\underline{I}-\underline{R}) s^{2}$. This matrix, which is an estimate of the one defined by equation $5.5-3$, is for weighted residuals $\underline{\hat{\hat{u}}}$ composed of both sample and prior information, with the prior information occupying the last $n_{p}$ rows and columns.
7. Correlation matrix for $\hat{\underline{\hat{u}}}$. This matrix is derived from $(\underline{I}-\underline{R}) s^{2}$.
Program Listing
```
    C RESIDUALS ANALYSIS PROGRAM BY R. L. COOLEY, USGS, DENVER, COLO.
    $LARGE: R
            IMPLICIT DOUBLE PRECISION (A-H,O-Z)
            DIMENSION X (20,70), COV (20,20),W(70),WP(20),IPR(20),R(90,90),D(90)
            1,G(90),F(90)
            COMMON/ITP/IIN, LOUT
            COMMON/FLT/X,R
            EQUIVALENCE (X(1,1),D(1)),(W(1),F(1),WP(1)),(COV(1, 1),G(1))
            OPEN (5,FILE='RESAN.DAT',STATUS='OLD',ACCESS='SEQUENTIAL'
            1,FORM=' FORMATTED')
            OPEN (6,FILE='RESAN.OUT',STATUS='NEW',ACCESS='SEQUENTIAL'
            1,FORM=' FORMATTED')
            OPEN (8,FILE='RESAN.IN',STATUS='OLD',ACCESS='SEQUENTIAL'
            1,FORM='UNFORMATTED')
C**FORMAT LIST
            l FORMAT (5I5,F10.0)
            2 FORMAT (8F10.0)
            3 FORMAT (9H1NVAR =,I4/9H NOBS =,I4/9H NPRIR =,I4
            1/9H NSETS =,I4/9H NRAN = ,I4/9H VAR =,G11.5)
            4 FORMAT (1H0,14X,42HRELIABILITY WEIGHTS FOR SAMPLE INFORMATION
            1/1H ,3X,3(3HNO.,11X,1HW,9X))
            5 FORMAT (16I5)
            6 FORMAT (1H0,14X,43HNO.S OF PARAMETERS HAVING PRIOR INFORMATION
            1/1H ,3X,3(3HNO.,8X,3HIPR,10X))
            7 FORMAT (1H0,14X,40HSTANDARD DEVIATIONS OF PRIOR INFORMATION
            1/1H ,3X,3(3HNO.,10X,2HWP,9X))
            8 FORMAT (19HO COVARIANCE MATRIX)
            9 FORMAT (38HO SENSITIVITIES FOR OPTIMUM PARAMETERS)
            10 FORMAT (6HOEV = ,Gll.5)
            11 FORMAT (1H0,20X,31HORDERED, RANDOM NORMAL DEVIATES/1H ,3X,2(3HNO.
                1,8X,1HD,14X,1HF,10X))
    12 FORMAT (1H0,18X,35HORDERED, CORRELATED NORMAL DEVIATES/1H ,3X
        1,2(3HNO., 8X,1HG,14X,1HF,10X))
    13 FORMAT (33HO COVARIANCE MATRIX FOR RESIDUALS)
    14 FORMAT (43HODATA GENERATED FROM RANDOM NUMBER SET NO. ,I3)
    15 FORMAT (34HO CORRELATION MATRIX FOR RESIDUALS)
    16 FORMAT (1HO,22X,26HCORRELATED NORMAL DEVIATES/1H ,3X,3(3HNO.,11X
        1,1HG,9X))
C**DEFINE INPUT FILES, OUTPUT FILE, AND ARRAY DIMENSIONS FOR PRTOT
        IIN=5
        ITB=8
        IOUT=6
        NVD=20
        NTD=90
    C**READ AND PRINT INPUT DATA THEN CONVERT IT INTO FORMS NEEDED
C FOR CALCULATIONS
C**NOTE: NRAN MUST BE ODD AND MUST LIE BETWEEN 1 AND 1048575
        READ(IIN, 1) NVAR,NOBS,NPRIR,NSETS,NRAN,VAR
                            SET A
        WRITE(IOUT,3) NVAR,NOBS,NPRIR,NSETS,NRAN, VAR
        DO 25 J=1,NVAR
        READ (ITB) (COV (I,J),I=J,NVAR) SET B
        DO 20 I=J,NVAR
```


## Program Listing-Continued

```
    20 COV(J,I)=COV(I,J)
    25 CONTINUE
        WRITE(IOUT,8)
        CALL PRTOT(COV,NVAR,NVAR,NVD)
        READ(ITB) (W(I),I=1,NOBS)
        SET C
        WRITE(IOUT,4)
        CALL PRTOT(W,NOBS,1,0)
        DO 35 J=1,NOBS
        READ(ITB) (X(I,J),I=1,NVAR) SET D
    35 CONTINUE
        WRITE(IOUT,9)
        CALL PRTOT(X,NVAR,NOBS,NVD)
        DO 45 J=1,NOBS
        WT=W(J)**. }
        DO 40 I=1,NVAR
    40 X(I,J)=X(I,J)*WT
    4 5 ~ C O N T I N U E ~
        IF(NPRIR.LT.1) GO TO 55
        READ(IIN,2) EV
        SET E
        WRITE(IOUT,10) EV
        READ(IIN,5) (IPR(I),I=1,NPRIR)
        SET F
        WRITE(IOUT,6)
        CALL PRTOTC(IPR,NPRIR)
        READ (IIN,2) (WP(I),I=1,NPRIR) SET G
        WRITE(IOUT,7)
        CALL PRTOT(WP,NPRIR,1,0)
        SIGMA=EV**. }
        DO 50 I=1,NPRIR
    50 WP(I)=SIGMA/WP(I)
C**COMPUTE (I-R)*VAR MATRIX
    55 DO 80 K=1,NOBS
        DO 70 J=1,NVAR
        SUM=0.
        DO 60 I=1,NVAR
    60 SUM=SUM+X (I,K)*COV (I,J)
    70 R(J,K)=SUM
    8 0 ~ C O N T I N U E ~
    IF(NPRIR.LT.1) GO TO 90
    DO }84\textrm{K}=1\mathrm{ ,NOBS
    DO }82\textrm{I}=1\mathrm{ ,NPRIR
    J=IPR(I)
    82 R(I+NOBS,K)=-WP(I)*R(J,K)
    84 CONTINUE
    DO 88 J=1,NPRIR
    L=IPR (J)
    DO }86\mathrm{ I=J,NPRIR
    K=IPR(I)
    86 R(I+NOBS,J+NOBS)=-WP(I)*COV(K,L)*WP(J)
    8 8 \text { CONTINUE}
    90 DO 110 K=1,NOBS
        DO 100 J=K,NOBS
        SUM=0.
        DO 95 I=1,NVAR
```

Program Listing-Continued
$95 \operatorname{SUM}=\operatorname{SUM}+X(\mathrm{I}, \mathrm{K}) * R(\mathrm{I}, \mathrm{J})$
$100 \mathrm{R}(\mathrm{J}, \mathrm{K})=-\mathrm{SUM}$
110 CONTINUE
NTOT=NOBS+NPRIR
DO $130 \mathrm{~J}=1$, NTOT
DO $120 \mathrm{I}=\mathrm{J}, \mathrm{NTOT}$
$120 \mathrm{R}(\mathrm{J}, \mathrm{I})=\mathrm{R}(\mathrm{I}, \mathrm{J})$
$130 \mathrm{R}(\mathrm{J}, \mathrm{J})=\operatorname{VAR}+\mathrm{R}(\mathrm{J}, \mathrm{J})$
C**COMPUTE THEORETICAL FREQUENCIES FOR DATA SETS
TMP $=\mathrm{NTOT}+1$
DO $135 \mathrm{I}=1, \mathrm{NTOT}$
TEMP=I
135 F(I) $=$ TEMP $/$ TMP
SIGMA=VAR** 5
DO $180 \mathrm{~K}=1$, NSETS
WRITE (IOUT,14) K
$C * * C O M P U T E$ RANDOM NORMAL DEVIATES D AND CORRELATED NORMAL DEVIATES G
DO $150 \mathrm{I}=1$, NTOT
SUM $=-6$.
DO $140 \mathrm{~J}=1,12$
140 SUM=SUM+RANUM (NRAN)
$150 \mathrm{D}(\mathrm{I})=$ SIGMA* SUM
DO $170 \mathrm{~J}=1$, NTOT
SUM=0.
DO $160 \mathrm{I}=1, \mathrm{NTOT}$
$160 \operatorname{SUM}=\operatorname{SUM}+\mathrm{R}(\mathrm{I}, \mathrm{J}) * D(\mathrm{I})$
$170 \mathrm{G}(\mathrm{J})=\mathrm{SUM} / \mathrm{VAR}$
$C * * O R D E R$ AND PRINT RANDOM NORMAL DEVIATES AND CORRELATED NORMAL DEVIATES
DO $174 \mathrm{I}=1$, NTOT
DO $172 \mathrm{~J}=\mathrm{I}, \mathrm{NTOT}$
IF(D(J).GE.D(I)) GO TO 172
TMP $=D(I)$
$D(I)=D(J)$
$D(J)=T M P$
172 CONTINUE
174 CONTINUE
WRITE(IOUT, 11)
CALL PRTOTA(D,F,NTOT)
WRITE (IOUT, 16)
CALL PRTOT (G,NTOT, 1,0)
DO $178 \mathrm{I}=1$, NTOT
DO $176 \mathrm{~J}=\mathrm{I}, \mathrm{NTOT}$
IF(G(J).GE.G(I)) GO TO 176
TMP $=G(I)$
$G(I)=G(J)$
$G(J)=T M P$
176 CONTINUE
178 CONTINUE
WRITE (IOUT, 12)
CALL PRTOTA (G,F,NTOT)
180 CONTINUE
D
C**PRINT COVARIANCE MATRIX (I-R)*VAR WRITE (IOUT, 13)

## Program Listing-Continued

CALL PRTOT(R,NTOT,NTOT,NTD)
C**COMPUTE AND PRINT CORRELATION MATRIX DO $210 \mathrm{I}=1$, NTOT
$210 \mathrm{D}(\mathrm{I})=\mathrm{R}(\mathrm{I}, \mathrm{I}) * * .5$
DO $230 \mathrm{~J}=1$, NTOT
TMP $=\mathrm{D}(\mathrm{J})$
DO 220 I=J,NTOT
$R(I, J)=R(I, J) /(T M P * D(I))$
$220 \mathrm{R}(\mathrm{J}, \mathrm{I})=\mathrm{R}(\mathrm{I}, \mathrm{J})$
230 CONTINUE
WRITE (IOUT, 15)
CALL PRTOT(R,NTOT,NTOT,NTD)
STOP
END
SUBROUTINE PRTOT(C,NR,NC,NRD)
C**PRINT MATRICES AND VECTORS
IMPLICIT DOUBLE PRECISION (A-H,O-Z)
DIMENSION C(1)
COMMON/ITP/IIN, IOUT
IF (NC.EQ.1) GO TO 25
DO $20 \mathrm{~L}=1, \mathrm{NC}, 10$
$\mathrm{J} 10=\mathrm{L}+9$
IF(J10.GT.NC) J10=NC
WRITE(IOUT, 35) (J, J=L, J10)
WRITE (IOUT,50)
$\mathrm{KBC}=(\mathrm{L}-1)$ *NRD
KEC=(J10-1)*NRD
DO $10 \mathrm{I}=1$, NR
$K B=K B C+I$
$\mathrm{KE}=\mathrm{KEC}+\mathrm{I}$
10 WRITE (IOUT, 40) I, (C (K) , K=KB, KE, NRD)
20 CONTINUE
RETURN
$25 \mathrm{~N}=\mathrm{NR} / 3$
IF ((3*N).NE.NR) $N=N+1$
DO $30 \mathrm{~K}=1, \mathrm{~N}$
$30 \operatorname{WRITE}($ IOUT, 80$)(\mathrm{L}, \mathrm{C}(\mathrm{L}), \mathrm{L}=\mathrm{K}, \mathrm{NR}, \mathrm{N})$
RETURN
35 FORMAT ( $1 \mathrm{HO} 0,10(9 \mathrm{X}, \mathrm{I} 3)$ )
40 FORMAT ( $1 \mathrm{H}, \mathrm{I} 3,1 \mathrm{X}, 10(1 \mathrm{X}, \mathrm{Gl1} .5)$ )
50 FORMAT ( 1 H )
80 FORMAT (1H , 2X,3(I3,7X,G11.5,3X))
END
SUBROUTINE PRTOTA(VALA, VALB, NO)
C**PRINT VALUES IN TWO GROUPS OF THREE COLUMNS
IMPLICIT DOUBLE PRECISION (A-H,O-Z)
dIMENSION VALA(NO), VALB(NO)
COMMON/ITP/IIN, IOUT
$\mathrm{NR}=\mathrm{NO} / 2$
IF ( $2 *$ NR.NE.NO) NR=NR+1
DO $10 \mathrm{~K}=1$, NR
WRITE (IOUT, 20) (L, VALA(L), VALB(L), L=K, NO,NR)
10 continue

```
    Program Listing-Continued
        RETURN
        20 FORMAT (1H ,2X,2(I3,4X,Gll.5,4X,G11.5,4X))
        END
        SUBROUTINE PRTOTC(IVAL,NO)
C**PRINT INTEGERS IN THREE GROUPS OF TWO COLUMNS
        DIMENSION IVAL(NO)
        COMMON/ITP/IIN,IOUT
        NR=NO/3
        IF(3*NR.NE.NO) NR=NR+1
        DO 10 K=1,NR
        WRITE(IOUT, 20) (L,IVAL(L), L=K,NO,NR)
    10 CONTINUE
        RETURN
        20 FORMAT (1H ,2X,3(I3,8X,I4,9X))
        END
        FUNCTION RANUM(IRAN)
        IMPLICIT DOUBLE PRECISION (A-H,O-Z)
        DATA MODU,MULT,NADD/1048576,1027,221589/
        IRAN=MULT*IRAN+NADD
        IRAN=IRAN- (IRAN/MODU) *MODU
        RANUM=FLOAT (IRAN)/FLOAT (MODU)
        RETURN
        END
    Listing of Inserts to the Regression Ground-Water Flow Program.
C
C**INSERT AFTER STATEMENT 480
        OPEN (8,FILE='RESAN.IN',STATUS='NEW',ACCESS='SEQUENTIAL'
        1, FORM='UNFORMATTED')
            ITB=8
        REWIND ITB
        DO 1000 J=1,NVAR
    1000 WRITE(ITB) (A(I,J),ImJ,NVAR)
        DO 1100 I=1,NOBS
    1100 VL(I)=W(I)*W(I)
        WRITE(ITB) (VL(I),I=1;NOBS)
        DO 1200 J=1, NOBS
    1200 WRITE(ITB) (X (I,J),I=1,NVAR)
C
```


## References Cited

Beck J.V., and Arnold K.J., 1977, Parameter estimation in engineering and science: New York, John Wiley, 501 p.
Draper, N.R., and Smith, H., 1981, Applied regression analysis [2d ed]: New York, John Wiley, 709 p.
Lieberman, G.J., 1961, Prediction regions for several predictions from a single regression line: Technometrics, v. 3, p. 21-27.

## Additional Reading

Beck, J.V., and Arnold, K.J., 1977, Parameter estimation in engineering and science: New York, John Wiley, 501 p.
Graybill, F. A., 1976, Theory and application of the linear model: North Scituate, Massachusetts, Duxbury, 704 p.

## 6 Some Advanced Topics

### 6.1 Advanced Models

### 6.1.1 Regression When the Dependent Variable is Implicit

In all cases considered in previous sections, the assumption has been made that the deterministic part of the model equation can be solved explicitly for the dependent variable. However, this may not be true in some cases. An example of such a model written in terms of the true value for the dependent variable $f=f\left(\xi, \beta_{1}, \beta_{2}\right)$ is

$$
\begin{gather*}
\frac{\beta_{1} k}{3\left(\beta_{1}+\beta_{2}\right)}\left\{\operatorname{1/2} \ln \left(\frac{(k+f)^{2}}{k^{2}-k f+f^{2}}\right)+\sqrt{3}\left[\tan ^{-1}\left(\frac{2 f-k}{k \sqrt{3}}\right)\right.\right. \\
\left.\left.-\tan ^{-1}\left(\frac{-1}{\sqrt{3}}\right)\right]\right\}-\xi=0 \tag{6.1-1}
\end{gather*}
$$

where $k=\sqrt[3]{\beta_{1}+\beta_{2}}$ and $f=f\left(\xi, \beta_{1}, \beta_{2}\right)$ is the exact solution of equation 6.1-1. As can be seen, $f$ is implicit in the model equation and cannot be directly solved for. A general deterministic form for an exact model (that is, a model that does not contain $\epsilon$ ) where the dependent variable is implicit is

$$
\begin{equation*}
g[f(\xi, \beta), \xi, \beta]=0 \tag{6.1-2}
\end{equation*}
$$

where $\xi$ and $\underline{\beta}$ are defined as usual.
Based on equation 6.1-2, a true regression model can be written in terms of observation vector $\underline{Y}$ and disturbances $\underline{\epsilon}$ in the usual form

$$
\begin{equation*}
\underline{Y}=f(\underline{\underline{\xi}}, \underline{\beta})+\underline{\epsilon} \tag{6.1-3}
\end{equation*}
$$

where $f(\xi, \underline{\beta})$ is the vector of order $n$ that is computed from

$$
\begin{equation*}
g[f(\underline{\xi}, \underline{\beta}), \underline{\xi}, \underline{\beta}]=\underline{0} . \tag{6.1-4}
\end{equation*}
$$

In equation 6.1-4 vector $g$ of order $n$ represents $n$ equations, each of which has the form of equation 6.1-2 written for an observation point. As an example, equation 6.1-1 would be written in the form

$$
\begin{align*}
& g_{i}= \\
& \frac{\beta_{1} k}{3\left(\beta_{1}+\beta_{2}\right)}\left\{1 / 2 \ln \left\lvert\, \frac{\left(k+f_{i}\right)^{2}}{k^{2}-k f_{i}+f_{i}^{2}}\right.\right)+\sqrt{3}\left[\tan ^{-1}\left(\frac{2 f_{i}-k}{k \sqrt{3}}\right)\right. \\
& \left.\left.-\tan ^{-1}\left(\frac{-1}{\sqrt{3}}\right)\right]\right\}-\xi_{i}=0, i=1,2, \ldots, n . \quad(6.1-5) \tag{6.1-5}
\end{align*}
$$

Note that in equation 6.1-5 only $f$ evaluated at point $i$ (that is, $f_{i}$ ) appears in the equation to compute $g_{i}$. However, in general this equation could contain values of $f$ evaluated at any number of the possible points $j=1,2, \ldots, n$. An example of this type of model is the numerical model discussed in the next section.

The estimated regression model derived from equations 6.1-3 and 6.1-4 is

$$
\begin{equation*}
\underline{Y}=f(\xi, \underline{b})+\underline{e} \tag{6.1-6}
\end{equation*}
$$

and

$$
\begin{equation*}
g[f(\xi, \underline{z}, \underline{b}), \underline{\xi}, b]=\underline{0} \tag{6.1-7}
\end{equation*}
$$

where $\underline{b}$ and $\underline{e}$ are, as usual, general estimates of $\underline{\beta}$ and $\epsilon$, respectively.

The general approach to solving the implicitvariable problem is very similar to that followed in section 3.3.1. First, the dependent variable values are written using a Taylor series expansion about an initial set of parameters. Then, from this, the linearized regression problem is set up and solved recursively to give the final solution to the nonlinear problem.

Taylor series expansion of $f$ about an arbitrary initial set of parameters $\underline{b}_{0}$ can be written in the form of equation 3.3-5,

$$
\begin{equation*}
f(\xi, \underline{b}) \cong f_{0}+\underline{X}_{0}\left(\underline{b}-\underline{b}_{0}\right) \tag{6.1-8}
\end{equation*}
$$

where

$$
\begin{align*}
& \epsilon_{0}=f\left(\xi, b_{0}\right)  \tag{6.1-9}\\
& \underline{X}_{0}=\left\{\left.\frac{\partial f_{i}}{\partial b_{j}}\right|_{\underline{b}=\underline{b}_{0}}\right\} \tag{6.1-10}
\end{align*}
$$

By using equation 6.1-6, equation 6.1-8 can be modified to give the estimated linearized regression model

$$
\begin{equation*}
\underline{Y}-f_{0}=\underline{X}_{0}\left(\underline{b}-\underline{b}_{0}\right)+\underline{e} \tag{6.1-11}
\end{equation*}
$$

which is exactly the same as the model used for the standard nonlinear case discussed in section 3.3.1.

To solve the linearized regression model based on equation 6.1-11, $f_{0}$ can be computed from equation 6.1-7 with $\underline{b}=\underline{b}_{0}$, and $\underline{X}_{0}$ can be computed by implicitly differentiating equation 6.1-7 with respect to $b_{j}(j=1,2, \ldots, p)$, setting $\underline{b}=\underline{b}_{0}$, then solving for $\underline{X}_{0}$. To accomplish this computation of $X_{0}$, note that for any differential change $d \underline{b}$ in parameter vector $\underline{b}$ to produce a new solution $f$ of $g=0$, the total differential $d g$ must equal zero because $g$ is always zero. Hence, by employing the chain rule of calculus, there results

$$
\begin{equation*}
d g=\left(\frac{\partial g}{\partial b_{j}}+\underline{M} \frac{\partial f}{\partial b_{j}}\right) d b_{j}=\underline{0}, j=1,2, \ldots, p \tag{6.1-12}
\end{equation*}
$$

where $\underline{M}=\left\{M_{i j}\right\}=\left\{\partial g_{i} / \partial f_{j}\right\}$. Note that if $g_{i}$ contains only $f_{\mathrm{i}}$, then $M$ is diagonal. Equation 6.1-12 can be evaluated using $\underline{b}=\underline{b}_{0}$ and $t_{0}=f\left(\xi, \xi_{2}, \underline{b}_{0}\right)$ to give

$$
\left(\frac{\partial f}{\partial b_{j}}\right)_{0}=\underline{X}_{j}^{0}=-M_{0}^{-1}\left(\left.\frac{\partial g}{\partial b_{j}}\right|_{0}, j=1,2, \ldots, p(6.1-13)\right.
$$

where

$$
\begin{equation*}
M_{0}=\left\{\left.\frac{\partial g_{i}}{\partial f_{j}}\right|_{f=f_{0}}\right\}_{(n \times n)} \tag{6.1-14}
\end{equation*}
$$

and subscript (or superscript) 0 means that the quantity is evaluated using $\underline{b}=\underline{b}_{0}$ and $f=f_{0}$. Solution of equation 6.1-7 for $f_{0}$ (using, for example, Newton iteration) followed by solution of equation 6.1-13 for $X_{0}$ provides a convenient method of obtaining initial values $f_{0}$ and ${\underset{X}{0}}^{0}$ from initial parameter estimates $\underline{b}_{0}$. However, for subsequent iterations this method can be time consuming because it involves solving $g=\underline{0}$ each time a new vector $f$ and a new set of sensitivity vectors $\underline{X}_{j}$ are to be computed from an updated parameter set.

A good method for computing good approximate values of $f$ and $X$ corresponding to some arbitrary parameter set $\underline{b}$ that is close to $\underline{b}_{o}$
involves approximating $g$ with another Taylor series expansion. If $\underline{b}$ is close enough to $\underline{b}_{0}$ to allow dropping all terms except linear terms, then Taylor series expansion of equation 6.1-7 about initial set of dependent variable values $f_{0}$ can be written as

$$
\begin{align*}
& g[f(\underline{\xi}, \underline{b}), \underline{\xi}, \underline{b}]=\underline{0} \\
& \cong g_{0}(\underline{b})+\underline{M}_{0}(\underline{b})\left(f-f_{0}\right) \tag{6.1-15}
\end{align*}
$$

where $g_{0}(\underline{b})=g\left(f_{0}, \underline{\xi}, \underline{b}\right), \underline{M}_{0}(\underline{b})=\left\{\partial g_{i} /\left.\partial f_{j}\right|_{f=f_{0}}\right\}$ and $\ell \cong f(\xi, \underline{b})$. By knowing $\underline{b}$, equation 6.1-15 may be solved for $f$. Corresponding values of $\underline{X}$ are obtained as follows. If equation 6.1-15 is implicitly differentiated with respect to $b_{j}(j=1,2, \ldots, p)$, there results
$\frac{\partial g_{0}(\underline{b})}{\partial b_{j}}=-\underline{M}_{0}(\underline{b}) \frac{\partial f}{\partial b_{j}}-\frac{\partial \underline{M}_{0}(\underline{b})}{\partial b_{j}}\left(f-f_{0}\right)$
or,

$$
\begin{array}{r}
\frac{\partial f}{\partial b_{j}}=\underline{X}_{j}=-\underline{M}_{0}^{-1}(\underline{b})\left|\frac{\partial g_{0}(\underline{b})}{\partial b_{j}}+\frac{\partial \underline{M}_{0}(\underline{b})}{\partial b_{j}}\left(f-f_{0}\right)\right|, \\
j=1,2, \ldots, p . \quad(6.1-17 \tag{6.1-17}
\end{array}
$$

By using the above results, solution of the nonlinear regression problem is obtained by a procedure analogous to the procedure followed for the standard nonlinear problem. As indicated previously, to begin the first iteration assume an initial set of parameters $\underline{b}_{0}$ and solve equation 6.1-7 for $f_{0}=f\left(\xi, b_{0}\right)$. Then solve equation 6.1-13 for $\underline{X}_{0}$. Next, form and solve normal equation $3.3-10$ by minimizing $S(\underline{b})$ (given by equation 3.3-4) with respect to $\underline{b}$, then scaling the resulting equations with $\underline{C}_{0}$. That is, form and solve for $\underline{\delta}_{1}$

$$
\begin{equation*}
\underline{S}_{0}^{T} \omega S_{0} \underline{\delta}_{1}=\underline{S}_{0}^{T} \omega\left(\underline{Y}-t_{0}\right) \tag{6.1-18}
\end{equation*}
$$

where

$$
\begin{align*}
& \underline{S}_{0}=\underline{X}_{0} \underline{C}_{0}  \tag{6.1-19}\\
& \underline{\delta}_{1}=\underline{C}_{0}^{-1}\left(\underline{b}_{1}-\underline{b}_{0}\right) . \tag{6.1-20}
\end{align*}
$$

For the second iteration, solve equation 6.1-15, written using $\underline{b}=\underline{b}_{1}$, for $f-f_{0}=f_{1}-f_{0}$. The result can be stated as

$$
\begin{equation*}
\underline{u}_{1}=-\underline{\underline{M}}_{1}^{-1} g_{1} \tag{6.1-21}
\end{equation*}
$$

where subscript 1 on $\underline{\underline{M}}$ and $g$ indicates evaluation using the most recent values available for $f$ and $\underline{b}$ (that is, $f_{0}$ and $\underline{b}_{1}$ ), and

$$
\begin{equation*}
\underline{u}_{1}=f_{1}-f_{0} \tag{6.1-22}
\end{equation*}
$$

Next, solve equation 6.1-17 for $\underline{X}_{1}$ :

$$
\begin{gather*}
X_{j}^{1}=-\underline{M}_{1}^{-1}\left[\left(\frac{\partial g_{0}(\underline{b})}{\partial b_{j}}\right){ }_{1}+\left(\frac{\partial \underline{M}_{0}(\underline{b})}{\partial b_{j}}\right) 1_{1} \underline{u}_{1}\right], \\
j=1,2, \ldots, p \tag{6.1-23}
\end{gather*}
$$

where $(\cdot)_{1}$ indicates evaluation using $\underline{b}=\underline{b}_{1}$. Finally, form and solve the normal equations, written in terms of $\underline{S}_{1}$ and $f_{1}$, for $\delta_{2}$.

For general iteration $r$, the equations to solve are

$$
\begin{array}{cc}
\underline{u}_{r}=-\underline{\underline{M}}_{r}^{-1} \underline{g}_{r} & (6.1-24 \\
\underline{X}_{j}^{r}=-\underline{\underline{M}}_{r}^{-1}\left[\left(\left.\frac{\partial g_{r-1}(\underline{b})}{\partial b_{j}}\right|_{r}+\left(\frac{\partial \underline{M}_{r-1}(\underline{b})}{\partial b_{j}}\right)_{r} \underline{u}_{r}\right],\right. \\
& j=1,2, \ldots, p \\
t_{r}=\underline{u}_{r}+f_{r-1} & (6.1-25 \\
S_{r}^{T} \underline{\omega}_{r} \delta_{r+1}=\underline{S}_{r}^{T} \underline{\underline{\omega}}\left(\underline{Y}-f_{r}\right) \\
\underline{b}_{r+1}=\underline{C}_{r} \underline{\delta}_{r+1}+\underline{b}_{r} & (6.1-26  \tag{6.1-28}\\
&
\end{array}
$$

where

$$
\begin{equation*}
\underline{S}_{r}=\underline{X}_{r} \underline{C}_{r} \text {, } \tag{6.1-29}
\end{equation*}
$$

$f_{0}=f\left(\xi, \underline{b}_{0}\right)$ so that $g_{0}=\underline{0}$, and $f_{-1}=f_{0}$.
At convergence of the solution $g_{r}, u_{r}$, and $\delta_{\gamma+1}$ all tend to zero so that $f_{r}=f(\xi, \hat{b})$ where $\underline{\underline{\hat{b}}}=\underline{b}_{r+1} \underline{\underline{b}}_{\underline{b}}^{r}$. At this point $S(\underline{\hat{b}})=(\underline{Y}-f(\underline{\xi}, \underline{b}))^{T} \underline{\omega}$ $\cdot(\underline{Y}-f(\underline{\xi}, \underline{b}))$ is at a minimum, and the nonlinear regression problem has been solved.

The solution procedure given by equations 6.1-24 through 6.1-29 can actually be considered to be a generalization of the GaussNewton procedure discussed in section 3.3.1, because if the standard nonlinear model is stated in the form

$$
\begin{equation*}
g=f(\xi, \underline{b})-f=\underline{0} \tag{6.1-30}
\end{equation*}
$$

then

$$
\begin{gather*}
g_{r}=f\left(\xi_{\underline{j}}, \underline{b_{r}}\right)-f_{r-1}  \tag{6.1-31}\\
M_{r}=-\underline{I}  \tag{6.1-32}\\
\left(\left.\frac{\partial g_{r-1}(\underline{b})}{\partial b_{j}}\right|_{r}=\left(\frac{\partial f}{\partial b_{j}}\right)_{r}=\underline{X}_{j}^{r}\right. \tag{6.1-33}
\end{gather*}
$$

and equations 6.1-24 through 6.1-29 become the standard Gauss-Newton algorithm.

Iteration parameters $\rho$ and $\mu$ should be applied to the present method in the same manner as they are for the Gauss-Newton method. Use of $\rho$ to modify step size $\underline{\delta}_{r+1}$ leads to equation 3.3-19 $\left(\underline{b}_{r+1}=\rho \underline{C}_{f_{r}} \underline{\delta}_{r+1}+\underline{b}_{r}\right)$ to compute $\underline{b}_{r+1}$. To employ $\mu$, equation 6.1-27 is transformed to

$$
\begin{equation*}
\left(S_{\omega}^{T} \underline{S}_{r}+\mu \underline{I}\right) \underline{\delta}_{r+1}=S_{-}^{T}{ }_{\omega}^{\omega}\left(\underline{Y}-f_{r}\right) . \tag{6.1-34}
\end{equation*}
$$

The method for solution of the implicitvariable model given here requires the same three conditions to guarantee convergence to a global minimum as discussed for the modified Gauss-Newton method in section 3.3.3. However, in addition, the method requires that $M$ and $\partial g / \partial b_{j}(j=1,2, \ldots, p)$ be continuous and unique for all $\underline{b}$ belonging to region $R$ (see equation 3.3-25).
Solution Algorithm.

1. Before the first iteration, solve equation 6.1-7 for $f_{0}$ using an initial estimate $\boldsymbol{b}_{0}$ for $\underline{b}$, and set $t_{-1}=f_{0}$.
2. Solve equation 6.1-24 and equation 6.1-25 for $\underline{u}_{r}$ and $\underline{X}_{j}^{r}(j=1,2, \ldots, p)$.
3. Solve equation 6.1-26 for $\mathcal{t}$.
4. Solve equation 6.1-34 for $\hat{\delta}_{r+1}$.
5. Solve equation 3.3-17 for $\underline{d}_{r+1}$.
6. Solve equation 3.3-19 for $\underline{b}_{r+1}$.
7. Test to determine if $\left|d_{i}^{r+1} / \mathrm{c}\right|>\epsilon(i=1,2, \ldots, p)$. (See step 5, Gauss-Newton algorithm.)
8. If $\left|d_{i}^{r^{+1}} / c\right|>\epsilon$, increment $r$ by one and return to 2 . If not, then the process has converged.

### 6.1.2 Regression When the ImplicitVariable Model is Numerical

If the numerical model assumes the general form of equation 3.3-21, which for convenience is restated here as

$$
\begin{equation*}
\underline{D}(\underline{h}, \underline{\xi}, \underline{\beta}) \underline{h}=q(\underline{h}, \underline{\xi}, \underline{\beta}), \tag{6.1-35}
\end{equation*}
$$

then the method derived in the previous section can be applied. The solution can be conceptually developed in two stages, first making the assumption that numerical solution points coincide exactly with observation points, which implies $m=n$, then, second, relaxing the assumption by following either of the two procedures described to obtain $f$ from $h$ for the Gauss-Newton method in section 3.3.2.

To develop the first stage of the solution, first note that because $m=n$, equation 6.1-35 may be written in the form $g=\underline{0}$ analogous to equation 6.1-7:

$$
\begin{equation*}
g=q-\underline{D} h=\underline{0} \tag{6.1-36}
\end{equation*}
$$

where $h=f$. Next, expand equation 6.1-36 in a Taylor series to give equations exactly analogous to equations 6.1-24 and 6.1-25. Pertinent quantities in these equations are given by

$$
\begin{gathered}
M_{i}^{T}=\left(\left.\frac{\partial q}{\partial h_{i}}\right|_{r}-\left(\left.\frac{\partial \underline{D}}{\partial h_{i}}\right|_{r} \underline{h}_{r-1}-\underline{D}_{i}^{r}, i=1,2, \ldots, m(6.1-37)\right.\right. \\
\left(\frac{\partial g_{r-1}(\underline{b})}{\partial b_{j}}\right)_{r}=\left(\frac{\partial \underline{q}}{\partial b_{j}}\right)_{r}-\left(\frac{\partial \underline{D}}{\partial b_{j}}\right)_{r} h_{r-1} \quad(6.1-38) \\
\left(\frac{\partial \underline{M}_{i}^{r-1}(\underline{b})}{\partial b_{j}}\right)_{r}=\left(\frac{\partial^{2} q}{\partial b_{j} \partial h_{i}}\right)_{r}-\left(\frac{\partial^{2} \underline{D}}{\partial b_{j} \partial h_{i}}\right)_{r} \underline{h}_{r-1} \\
-\left(\frac{\partial \underline{D}_{i}}{\partial b_{j}}\right)_{r}
\end{gathered}
$$

where subscript $i$ on a matrix denotes column $i$ of the matrix. By using equations 6.1-38 and 6.1-39, an equation analogous to 6.1-25 can be written

$$
\begin{gather*}
\left(\frac{\partial \underline{h}}{\partial b_{j}}\right)_{r}=-\underline{M}_{r}^{-1}\left\{\left(\frac{\partial \underline{q}}{\partial b_{j}}\right)_{r}-\left(\frac{\partial \underline{\underline{D}}}{\partial b_{j}}\right)_{r} \underline{h}_{r-1}+\right. \\
\sum_{i=1}^{m}\left\{\left(\frac{\partial^{2} q}{\partial b_{j} \partial h_{i}}\right)_{r}-\left(\frac{\partial^{2} \underline{\underline{D}}}{\partial b_{j} \partial h_{i}}\right)_{r} \underline{h}_{r-1}-\left(\frac{\partial \underline{D}_{i}}{\partial b_{j}}\right)_{r}\right] \\
\left.\cdot u_{i}^{r}\right\}, j=1,2, \ldots, p \tag{6.1-40}
\end{gather*}
$$

The second stage results from using one of the two procedures for obtaining $f$ from $\underline{h}$ described in section 3.3 .2 to obtain $f$ from $h$ and $X_{j}$ from $\partial \underline{h} / \partial b_{j}$ in the present case. With $f$ and $X$ defined, the solution algorithm of section 6.1.1 can be applied directly.

### 6.2 Modified Beale's Measure of Nonlinearity

Most of the methodology discussed here to analyze regression models is based on the assumption that the model is linear in the parameters. In the case that the model is nonlinear, Beale (1960, p. 54-55) developed an empirical measure of degree of nonlinearity with respect to the confidence regions on parameters. However, Guttman and Meeter (1965, p. 635) showed that if the degree of nonlinearity is high, Beale's measure seriously underestimates it. To correct for this underestimation problem, Linssen (1975) modified Beale's measure. More recently, Bates and Watts (1980) developed measures of nonlinearity based on the concepts of differential geometry. Although these measures are based on a much more extensive theory than Beale's (1960) measure or Linssen's (1975) modification, they also require extensive computation. Thus, here Beale's empirical measure as modified by Linssen is extended to give an approximate indication of the degree of average model nonlinearity.
To develop the measure, consider a linearized model of the form of equation 5.1-1, where, for convenience, general estimate $\underline{b}_{\rho}$ replaces $\beta$, and $\underline{b}_{0}$ is set equal to $\underline{b}$ to result in

$$
\begin{equation*}
\hat{f}_{\ell}^{0}=\hat{f}+\underline{X}\left(\underline{b}_{-}-\underline{b}\right) \tag{6.2-1}
\end{equation*}
$$

where $\hat{f}=f(\xi, \underline{\hat{b}})$. To emphasize that $f_{\ell}^{\circ}$ is a linearized estimate of $f\left(\xi_{2}, b_{l}\right)=f_{p}$, a superscript o is appended to $f_{p}$. If $f_{\ell}^{0}$ is calculated from equation 6.2-1 for $m$ sets of parameter vectors $b_{f}(\ell=1,2, \ldots, m)$, then a measure $c^{2}$ of model nonlinearity in the region covered by the varied parameter sets is

$$
\begin{equation*}
c^{2}=\sum_{\ell=1}^{m}\left(f_{l}-f_{\ell}^{0}\right)^{T} \underline{\underline{\omega}}\left(f_{l}-f_{\ell}^{0}\right) \tag{6.2-2}
\end{equation*}
$$

Equation 6.2-2 is the sum of squared distances (that is, squared lengths of vectors) between points $\underline{\omega}^{1 / 2} f_{\ell}^{0}$ and $\underline{\omega}^{1 / 2} f_{\ell}$ in observation space. (Recall that the distance between two points is the length of the vector joining the points, and that the squared length of a vector is given by the sum of squared lengths of its components.)

As explained further on, the most useful measure of nonlinearity is obtained by multiplying equation $6.2-2$ by the quantity

$$
q s^{2} / \sum_{\ell=1}^{m}\left[\left(f_{\ell}^{0}-\hat{f}\right)^{T} \omega\left(f_{\ell}^{0}-\hat{f}\right)\right]^{2}
$$

to obtain
which is an extension for $q \leq p$ of Linssen's (1975) modification of Beale's measure of nonlinearity (Beale, 1960, p. 54-55).

Equation 6.2-3 can be justified as follows (see also Guttman and Meeter, 1965, and Linssen, 1975). The weighted distance between $f_{\ell}^{0}$ and $f_{\ell}$ is designated $\epsilon d$ so that the geometric relationships among weighted vectors $\hat{f}, f_{\beta}$ and $f_{p}^{\circ}$ can be diagrammed in observation space as shown in figure 6.2-1. Now,

$$
\begin{align*}
& \epsilon^{2}=\frac{(\epsilon d)^{2}}{d^{2}}=\frac{d^{2}(\epsilon d)^{2}}{\left(d^{2}\right)^{2}} \\
& =d^{2} \frac{\left(f_{\ell}-f_{\ell}^{0}\right)^{T} \underline{\omega}\left(f_{\ell}-f_{\ell}^{0}\right)}{\left[\left(\hat{f}_{\ell}^{0}-\hat{f}\right)^{T} \underline{\omega}\left(f_{\ell}^{0}-\hat{f}\right)\right]^{2}} \tag{6.2-4}
\end{align*}
$$



Figure 6.2-1
where, by definition,

$$
\begin{align*}
d^{2} & =\left(f_{\ell}^{0}-\hat{f}\right)^{T} \underline{\omega}\left(f_{\ell}^{0}-\hat{f}\right) \\
& =\left[\underline{X}\left(\underline{b}_{\ell}-\hat{b}\right)\right]^{T} \underline{\omega}\left[\underline{X}\left(\underline{b}_{f}-\underline{\hat{b}}\right)\right] \\
& =\left(\underline{b}_{\beta}-\hat{b}^{T} \underline{X}^{T} \underline{\omega} \underline{b^{\prime}}-\underline{\hat{b}}\right) \tag{6.2-5}
\end{align*}
$$

To obtain a more convenient form for $d^{2}$, note that

$$
\begin{align*}
& \left(\underline{Y}-\mathcal{E}_{\ell}^{o}\right)^{T} \underline{\omega}\left(\underline{Y}-\mathcal{L}^{0}\right)-(\underline{Y}-\hat{f})^{T} \underline{\omega}(\underline{Y}-\hat{\hat{E}}) \\
& =\left(\underline{Y}-\hat{t}-\underline{X}\left(\underline{b}_{f}-\underline{\hat{b}}\right)\right)^{T} \underline{\omega}\left(\underline{Y}-\hat{t}-\underline{X}\left(\underline{b}_{f}-\underline{\hat{b}}\right)\right) \\
& -(\underline{Y}-\hat{t})^{T} \underline{\omega}(\underline{Y}-\hat{t}) \\
& =-2\left(\underline{b}_{l}-\underline{\hat{b}}\right)^{T} \underline{X}^{T} \underline{\omega}(\underline{Y}-\hat{f})+\left(\underline{b}_{\ell}-\underline{\hat{b}}\right)^{T} \underline{X}^{T} \underline{\omega} \underline{X}\left(\underline{b}_{\ell}-\underline{\hat{b}}\right) \\
& =\left(\underline{b}_{\boldsymbol{b}}-\underline{\hat{b}}\right)^{T} \underline{X}^{T} \underline{\omega} \underline{X}\left(\underline{b}_{\ell}-\underline{\hat{b}}\right) \tag{6.2-6}
\end{align*}
$$

where equation 5.1-9 was used. The combination of equations $6.2-5$ and $6.2-6$ shows that

$$
d^{2}=\left(\underline{Y}-f_{\ell}^{0}\right)^{T} \underline{\omega}\left(\underline{Y}-f^{0}\right)-(\underline{Y}-\hat{f})^{T} \underline{\omega}(\underline{Y}-\underline{b}) \cdot(6.2-7)
$$

If $f_{\ell}^{0}$ is assumed to lie on the edge of the confidence region given by equation 5.6-12 so that $\underline{b}_{P}^{T}=\left[\tilde{b}_{18}^{T} \tilde{\underline{\beta}}_{2}^{T}\right]$, then from equations 5.6-12 and 6.2-7 it can be seen that

$$
\begin{equation*}
d^{2}=q s^{2} F_{\alpha}(q, n-p) \tag{6.2-8}
\end{equation*}
$$

Hence, if both numerator and denominator of equation 6.2-4 are averaged over $m$ sets of parameters, there results

$$
\begin{equation*}
E\left(\epsilon^{2}\right) \cong \hat{N}_{b} F_{\alpha}(q, n-p) \tag{6.2-9}
\end{equation*}
$$

Based on equation 6.2-9, Beale (1960, p. 60) ranked the degrees of nonlinearity as follows: The model is highly nonlinear if

$$
\begin{equation*}
\hat{N}_{b}>1 / F_{\alpha}(q, n-p) \tag{6.2-10}
\end{equation*}
$$

because in this case $E\left(\epsilon^{2}\right)>1$, and the discrepancy is actually greater than $d$. If

$$
\begin{equation*}
\hat{N}_{b}<0.01 / F_{\alpha}(q, n-p), \tag{6.2-11}
\end{equation*}
$$

then the model is classed as being effectively linear because $E\left(\epsilon^{2}\right)<0.01$. For points in between, Beale (1960, p. 60) stated that the linear theory is adequate to give a rough idea of significance but may not bring out full implications of the analysis. However, Guttman and Meeter (1965, p. 636) noted that equation 6.2-11 may be overly conservative to define a maximum value of $\hat{N}_{b}$ for an approximately linear model. Experiments conducted by the authors indicate that, if

$$
\begin{equation*}
\hat{N}_{b}<0.09 / F_{\alpha}(q, n-p), \tag{6.2-12}
\end{equation*}
$$

then confidence intervals given by linear theory are fairly good approximations of the exact ones as given by Vecchia and Cooley (1987). Thus, equation 6.2-12 is used to define the maximum value of $\hat{N}_{b}$ to consider the model to be roughly linear.

Because equation 6.2-9 is justified by assuming that the points $f_{\ell}^{0}$ lie on the edge of the confidence region, a reasonable way to obtain the points is to choose them from equation 5.6-14, although, as noted by Beale (1960, p. 55), the points do not have to lie on the edge of the confidence region. Thus, one could use $m \leq 2 q$ sets of parameters $\underline{b}_{\rho}^{T}=\left[\hat{b}_{1 \ell} \hat{e}^{T}, \tilde{\tilde{B}}_{2}^{T}\right]$. Note that whether or not the model is linear, $\underline{\underline{b}}_{1 \ell}$ and $\tilde{\underline{ }}_{2}$ corresponding to the partition of $\underline{b}$ given in equation 5.6-3 are properly chosen without the necessity of performing additional least squares solutions to obtain each set $\underline{\tilde{b}}_{19}$. This fact is true because subset $\tilde{\underline{b}}_{1 \rho}$ is required to lie on the edge of the linearized confidence region.

Rigorous use of equations 6.2-10 through 6.2-12 theoretically requires that disturbances be distributed normally. However, it would be convenient to be able to gauge the degree of nonlinearity of the model irrespective of the
properties of $\underline{\epsilon}$. If the confidence region in equation $5.6-12$ were large enough to encompass virtually all physically plausible sets of parameters, then model nonlinearity as assessed using equations $6.2-10$ through $6.2-12$ would be meaningful. Based upon past experience, $F$ values generated using $\alpha=0.05$ have been found to yield such a confidence region and thus to be adequate to gauge nonlinearity.

## Problem 6.2-1

Four sets of parameters that correspond to four points on the edge of the linearized confidence region in equation $5.6-13$ result from problem 5.6-2. These four sets of parameters can be subdivided into two groups of two. Pick two different parameter sets from the two groups and compute two corresponding sets of drawdowns at the observation points using the nonlinear (Theis) model. Then, using the modified Beale's measure program (appendix 6.4.1), find the modified Beale's measure. Is the model nearly linear?

## Problem 6.2-2

Use the four parameter sets resulting from problem 5.6-3 in the nonlinear regression flow program of appendix 4.3.4, as augmented by the inserts of appendix 6.4.1, to compute the modified Beale's measure. Are the various statistical measures obtained from the linearized model approximately valid (at least as determined from the four parameter sets employed)?

### 6.3 Compatibility of Prior and Regression Estimates of Parameters

If the regression model contains prior information on the parameters, an important part of the analysis to determine whether or not the model is correct is to test the null hypothesis that the prior and sample information are in agreement; in other words,

$$
\begin{aligned}
H_{0}: E\left(\underline{Y}_{p}-f_{p}\left(\xi, \underline{b}_{0}\right)\right)= & \underline{X}_{p}\left(\underline{\beta}-\underline{b}_{0}\right) \text { versus } \\
& H_{1}: E\left(\underline{Y}_{p}-f_{p}\left(\underline{\xi}, \underline{b}_{0}\right)\right) \neq \underline{X}_{p}\left(\underline{\beta}-\underline{b}_{0}\right) .
\end{aligned}
$$

As indicated in section 5.5, graphical analysis of residuals can usually detect an incompatibility between sample and prior information. However, in some cases an additional test might be desired. Theil (1963) showed that the test statistic

$$
\begin{align*}
& \Gamma=\left(\underline{Y}_{p}-f_{p}\left(\underline{\xi}, \underline{b}_{0}\right)-\underline{X}_{p}\left(\underline{\hat{b}}^{*}-\underline{b}_{0}\right)\right)^{T} \\
& {\left[\sigma^{2} \underline{X}_{p}\left(\underline{X}_{s}^{T} \underline{\omega}_{s} \underline{X}_{s}\right)^{-1} \underline{X}_{p}^{T}+\sigma^{2} \underline{\omega}_{p}^{-1}\right]^{-1}} \\
& \cdot\left(\underline{Y}_{p}-\underline{f}_{p}\left(\underline{\xi}, \underline{b}_{0}\right)-\underline{X}_{p}\left(\hat{b}^{*}-\underline{b}_{0}\right)\right) \tag{6.3-1}
\end{align*}
$$

where vector $\underline{b}^{*}$ is the ordinary least squares estimate of vector $\beta$, is Chi square distributed with $n_{p}$ degrees of freedom $\left(\chi^{2}\left(n_{p}\right)\right.$ ) provided that all of the assumptions given by equations $5.2-1$ through $5.2-3,5.2-6$, and $5.2-10$ hold true, $\sigma^{2}$ is known, and $\underline{\omega}$ is of the form

$$
\underline{\underline{\omega}}=\left[\begin{array}{ll}
\underline{\omega}_{s} & \underline{0}  \tag{6.3-2}\\
\underline{0} & \underline{\omega}_{p}
\end{array}\right]
$$

where $\omega_{s}$ and $\underline{\omega}_{p}$ are known and symmetric positive definite of order $n_{s}$ and $n_{p}$, respectively. If $\underline{\omega}$ is of the form

$$
\underline{\omega}=\left[\begin{array}{ll}
\underline{V}_{s}^{-1} & \underline{0}  \tag{6.3-3}\\
\underline{0} & \underline{U}^{-1} \sigma^{2}
\end{array}\right]
$$

and $\sigma^{2}$ is unknown, then the test statistic

$$
\begin{align*}
& \hat{\Gamma}=\left(\underline{Y}_{p}-f_{p}\left(\underline{\xi}, \underline{b}_{0}\right)-\underline{X}_{p}\left(\hat{b}^{*}-\underline{b}_{0}\right)\right)^{T} \\
& \cdot\left[s^{2} \underline{X}_{p}\left(\underline{X}_{s}^{T} \underline{V}_{s}^{-1} \underline{X}_{s}\right)^{-1} \underline{X}_{p}^{T}+\underline{U} \underline{U}^{-1}\right. \\
& \cdot\left(\underline{Y}_{p}-\mathcal{f}_{p}\left(\underline{\xi}, \underline{b}_{0}\right)-\underline{X}_{p}\left(\underline{b}^{*}-\underline{b}_{0}\right)\right) \tag{6.3-4}
\end{align*}
$$

is asymptotically $\chi^{2}\left(n_{\lambda p}\right)$ distributed. If the computed value of $\Gamma, \gamma$, is greater than $\chi_{\alpha}^{2}\left(n_{p}\right)$, where $\alpha$ indicates significance level, then the null hypothesis is rejected.

## Problem 6.3-1

Using equation 6.3-4, test the compatibility
of the prior estimate of the boundary head, $f_{p 2}$ and the pure regression estimate, $\hat{b}_{2}{ }^{*}$, of problem 3.2-1. To conduct this test you will have to do an ordinary least-squares regression. The model of appendix 4.3 .4 may be employed for this in the same manner as for problem 4.2-1. Use the model output to obtain the necessary quantities in equation 6.3-4.

### 6.4 Appendix

### 6.4.1 Documentation of Program to Compute the Modified Beale's Measure

This program performs a straightforward computation of the modified Beale's measure, equation 6.2-3. Vectors $f_{l}, f_{l}^{0}$, and $\hat{L}$ are assumed to be composed of sample information and direct prior information on some or all parameters. The weight matrix for sample and prior information is assumed to be given in the form of equation 3.4-12, and the sensitivity matrix for the prior information is assumed to be given by equation 4.1-6.

There are two versions of the program. One is for general use, and all variables needed for the calculation must be read in. The other version is designed to be an integral part of the regression ground-water model documented in appendix 4.3.4 and requires only $q$ and the extra sets of parameters needed for the Beale's measure calculation as input in addition to input already required for the regression solution.

The programs were developed using the Microsoft Fortran Compiler, Version 3.3, with the DOS 2.0 operating system on an IBM PC/XT computer with the IBM 8088 Math Coprocessor and 256 KB memory. Except for the OPEN statements near the beginning of the general code, Fortran 66 was used throughout to make the codes as machine independent as possible. The general source code is contained in file BEALE.FOR, and the version designed to be inserted into the regression code is contained in file BLEINS.FOR, both of which are in the diskette accompanying this report.

Input data for General Version.-Data Set A. Problem size information; one card (format 4I5, F10.0).

| Line Columns | Varriable | Definition |
| :---: | :---: | :---: |
| 1-5 | NVAR | Number of parameters, $p$. |
| 6-10 | NRES | Number of restrictions, $q$. |
| 11-15 | NOBS | Number of sample observations, $n$. |
| 16-20 | NPRIR | Number of regression parameters having direct prior information, $n_{p}$. |
| 21-25 | NPTS | Number of data sets to compute the modified Beale's measure, $m$. |
| 26-35 | VAR | Error variance, $s^{2}$. |

Data Set B.
Estimated regression parameters, $\underline{\hat{b}}$ (format 8F10.0).

| Line columns | Variable |  |
| :---: | :--- | :--- |
| $1-10$ | BOPT(1) | Estimated regression parameters, entered sequentially |
| $11-20$ | BOPT(2) | from 1 through NVAR. |
| $\vdots$ | $\vdots$ |  |
|  | BOPT(NVAR) |  |

Data Set C.
Dependent variable vector for sample information, $\hat{t}_{s}$, computed using $\underline{\hat{b}}$ (format 8F10.0).

| Line columns | Variable | Definition |
| :--- | :---: | :---: |
| $1-10$ | FOPT(1) | Computed dependent variable values, entered sequen- |
| $11-20$ | FOPT(2) | tially from 1 through NOBS. |
| $\vdots$ | $\vdots$ |  |
|  | FOPT(NOBS) |  |

Data Set D.
Weight matrix for sample information, $\underline{V}_{s}^{-1}$ (format 8F10.0).

| Line columns | Variable | Definition |
| :--- | :---: | ---: |
| $1-10$ | W(1) | Diagonal weight matrix for sample information, |
| $11-20$ | W(2) | entered sequentially from 1 through NOBS. |
| $\vdots$ | $\vdots$ |  |
|  | W(NOBS) |  |

Data Set E.
Sensitivity matrix for sample information, $\underline{X}_{s}$ (format 8F10.0)

| Line columns | Variable | Definition |
| :---: | :---: | :---: |
| 1-10 | X $(1,1)$ | Sensitivity matrix for sample information, entered sequentially 1 through NVAR for each observation. Each new observation begins a new line, for a total of NOBS observations. |
| 11-20 | X $(2,1)$ |  |
| $\vdots$ | $\begin{gathered} \vdots \\ \text { X(NVAR,1) } \end{gathered}$ |  |
| 1-10 | X $(1,2)$ |  |
| : | $\begin{gathered} \vdots \\ \text { X(NVAR,2) } \end{gathered}$ |  |
|  | X(NVAR,NOBS) |  |

Data Set F.
Parameter numbers having prior information (format 16I5).

| Line columns | Variable | Definition |
| :--- | :---: | :---: |
| $1-5$ | IPR(1) | Array subscript numbers for regression parameters in |
| 6-10 | IPR(2) | BOPT(1) having prior information, entered in any |
| $\vdots$ | $\vdots$ | order from 1 through NPRIR. |
|  | IPR(NPRIR) |  |

Omit data set of NPRIR=0.
Data Set G.
Standard deviation matrix for prior information, $\underline{U}^{1 / 2}$ (format 8F10.0).

| Line columns | Variable | Definition |
| :--- | :---: | :--- |
| $1-10$ | WP(1) | Diagonal standard deviation matrix for prior informa- |
| $\vdots$ | $\vdots$ | tion, entered in the same order as IPR(I) from 1 |
|  | WP(NPRIR) | through NPRIR. |
| Omit data set if NPRIR $=0$ |  |  |

## Data Set H.

Alternate parameters sets, $\underline{b}_{\rho}$ (format 8F10.0). This data set and the next one are read in sequence ( $H, I, H, I, \ldots$ ) a total of NPTS times.

| Line columns | Variable | Definition |
| :--- | :--- | :--- |
| $1-10$ | $\mathrm{~B}(1)$ | Alternate parameter sets, entered sequentially 1 |
| $11-20$ | $\mathrm{~B}(2)$ | through NVAR. Order must be the same as for |
| $\vdots$ | $\vdots$ | BOPT(I). |
|  | $\mathrm{B}($ NVAR $)$ |  |

Data Set I.
Alternate dependent variable vectors for sample information, $f_{s}$, computed using $\underline{b}_{\ell}$ (format 8F10.0).

| Line columns | Variable |  |
| :--- | :---: | :---: |
| $\mathbf{1 - 1 0}$ | FC(1) | Alternate sample dependent variable values computed |
| $11-20$ | FC(2) | using the nonlinear model, entered sequentially 1 |
| $\vdots$ | $\vdots$ | through NOBS. Order must be the same as for |
|  | FC(NOBS $)$ | FOPT(I). |

Output for General Version.-Output is all clearly labeled; it is ordered as follows:

1. Data sets A through G.
2. Data sets $\mathbf{H}$ and I. Data for numbers 2 through 4 below are printed sequentially for each data set $\ell(\ell=1,2, \ldots, m)$.
3. Dependent variable vector, $L_{s l}^{\circ}$, for sample information, computed using the linearized model.
4. Total sums of squared differences $\left(\mathcal{L}_{-}-\hat{L}\right)^{T}$ $-\underline{\omega}\left(f_{l}-\hat{f}\right)$ and $\left(\hat{f}_{l}^{0}-\hat{L}\right)^{T} \underline{\omega}\left(\mathscr{f}_{i}^{0}-\hat{f}\right)$, where


$$
\left.\underline{\omega}=\left\lvert\, \begin{array}{ll}
\underline{V}_{s}^{-1} & \underline{0} \\
\underline{0} & \underline{U}^{-1} s^{2}
\end{array}\right.\right] .
$$

5. Beale's measure, $\hat{N}_{b}=B N$.

Use of Version Integral with the Regression Ground-Water Program.-This version consists of sets of statements to be inserted into the program of appendix 4.3.4, as indicated on the appended listing. Input is the same as if a regression solution were to be obtained, except that the initial set of parameters must be the optimum set $\underline{\hat{b}}$, and extra data relating to the modified Beale's measure is required. After entering data set $T$, use data sets $U$ and $V$ to enter the data for $\ell=1,2, \ldots, m$ alternate solutions. Follow these data with a final line to input $q$ and $s^{2}$ with format I5,F10.0. A complete regression solution is not obtained; only computations through the calculation of sensitivities on the first iteration are completed before proceeding to calculate the modified Beale's measure. Thus, output consists of regression output through number 19 (see "Output" in appendix 4.3.4) plus output analogous to numbers 2 through 4 of the general version of the modified Beale's measure program.

## 3

Program Listing for General Version.
C MODIFIED BEALE'S MEASURE PROGRAM BY R. L. COOLEY, USGS, DENVER, C COLO.

DIMENSION BOPT (20) , $\operatorname{FOPT}(70), \mathrm{B}(20), \mathrm{FC}(70), \mathrm{FL}(70), \mathrm{X}(20,70)$
$1, \mathrm{~W}(70), \operatorname{IPR}(20), \mathrm{WP}(20)$
COMMON/ITP/IIN, IOUT
COMMON/FLT/X
OPEN (5,FILE='BEALE.DAT',STATUS='OLD',ACCESS-'SEQUENTIAL'
1, FORM='FORMATTED')
OPEN ( 6, FILE=' BEALE.OUT', STATUS='NEW', ACCESS='SEQUENTIAL'
1, FORM=' FORMATTED')
C**FORMAT LIST
1 FORMAT (5I5, F10.0)
2 FORMAT (8F10.0)
3 FORMAT (9H1NVAR $=$,I4/9H NRES $=, 14 / 9 H$ NOBS $=, I 4$
$1 / 9 \mathrm{H}$ NPRIR $=, \mathrm{I} 4 / 9 \mathrm{H}$ NPTS $=$.I4 $/ 9 \mathrm{H}$ VAR $=$, G11.5)
4 FORMAT (1H0, 26X, 18HOPTIMUM PARAMETERS
1/1H , 3X, 3(3HNO., 9X,4HBOPT, 8X))
5 FORMAT ( $1 \mathrm{Ho} 0,9 \mathrm{X}, 52 \mathrm{hDEPENDENT}$ Variables COMPUTED WITH OPTIMUM PaRAME
1TERS/1H ,3X,3(3HNO.,9X,4HFOPT, 8X))
6 FORMAT (1H0,21X, 26HPARAMETERS FOR SAMPLE NO. ,I3
1/1H , 3X, 3(3HNO., 11X, 1HB, 9X))
7 FORMAT ( $1 \mathrm{HO}, 12 \mathrm{X}, 44 \mathrm{HDEPENDENT}$ VARIABLES COMPUTED FOR SAMPLE NO. ,I3
1/1H , 3X, 3(3HNO.,10X,2HFC, 9X))
8 FORMAT (38HO SENSITIVITIES FOR OPTIMUM PARAMETERS)
9 FORMAT ( 1 H0, $6 \mathrm{X}, 55$ HLINEARIZED DEPENDENT VARIABLES COMPUTED FOR SAMP
1LE NO. ,13/1H ,3X,3(3HNO.,10X,2HFL, 9X))
10 FORMAT ( $1 \mathrm{HO} 0,5 \mathrm{HBN}=$, G11.5)
11 FORMAT (23HOSS ((FC-FOPT)*W**.5) $=$, G11.5
$1 / 23 \mathrm{H}$ SS ((FL-FOPT) $*$ W**.5) $=$, G11.5)
12 FORMAT (1H0,14X,42HRELIABILITY WEIGHTS FOR SAMPLE INFORMATION
1/1H ,3X,3(3HNO.,10X,1HW,10X))
13 FORMAT (16I5)
14 FORMAT ( $1 \mathrm{HO}, 12 \mathrm{X}, 43 \mathrm{HNO} \mathrm{S}$ OF PARAMETERS HAVING PRIOR INFORMATION
1/1H , 3X, 3(3HNO. , 8X, 3HIPR, 10X))
15 FORMAT ( $1 \mathrm{HO}, 14 \mathrm{X}, 40 \mathrm{HSTANDARD}$ DEVIATIONS OF PRIOR INFORMATION
1/1H , 3X, 3(3HNO.,10X, 2HWP, 9X))
16 FORMAT (6HOEV $=$, G11.5)
C**DEFINE INPUT FILE, OUTPUT FILE, AND ARRAY DIMENSION IIN-5
IOUT $=6$
NVD=20
C**READ BASE DATA
READ (IIN,1) NVAR, NRES, NOBS, NPRIR,NPTS, VAR SET A
WRITE(IOUT, 3) NVAR, NRES, NOBS, NPRIR, NPTS, VAR
$\operatorname{READ}(\operatorname{IIN}, 2)$ ( $\operatorname{BOPT}(\mathrm{J}), \mathrm{J}=1$, NVAR)
SET B
WRITE(IOUT,4)
CALL PRTOTB (BOPT, NVAR)
$\operatorname{READ}(\operatorname{IIN}, 2)$ ( $\operatorname{FOPT}(\mathrm{I}), \mathrm{I}=1, \mathrm{NOBS})$ SET C
WRITE (IOUT, 5)
CALL PRṪOTB(FOPT,NOBS)
READ (IIN,2) (W(I), I=1,NOBS) SET D
WRITE(IOUT,12)
CALL PRTOTB ( $\mathrm{W}, \mathrm{NOBS}$ )Program Listing for General Version-Continued

```
DO \(20 \mathrm{~J}=1\), NOBS
```

            READ(IIN,2) (X(I,J),I=1,NVAR) SET E
    20 CONTINUE
    WRITE(IOUT, 8)
    CALL PRTOT(X,NVAR,NOBS,NVD)
    IF(NPRIR.LT.1) GO TO 45
    READ(IIN,2) EV
    WRITE(IOUT,16) EV
    READ(IIN,13) (IPR(I),I=1,NPRIR) SET G
    WRITE(IOUT,14)
    CALL PRTOTC(IPR,NPRIR)
    READ(IIN,2) (WP(I),I=1,NPRIR) SET H
    WRITE(IOUT,15)
    CALL PRTOTB(WP,NPRIR)
    DO40 I=1,NPRIR
    40 WP(I)=EV/(WP(I)*WP(I))
    C**READ DATA FOR EACH SAMPLE AND COMPUTE MODIFIED BEALE'S MEASURE, BN
45 SUMA=0.
SUMB=0.
DO }80\textrm{M}=1,NPT
READ(IIN,2) (B(J),J=1,NVAR) SET I
WRITE(IOUT,6) M
CALL PRTOTB(B,NVAR)
READ(IIN,2) (FC(I),I=1,NOBS) SET J
WRITE(IOUT,7) M
CALL PRTOTB(FC,NOBS)
SUMC=0.
SUMD=0.
DO 60 J=1,NOBS
SUM=FOPT(J)
DO 50 I=1,NVAR
50 SUM=SUM+X(I,J)*(B(I)-BOPT(I))
FL(J)=SUM
TMP=FC(J) - SUM
SUMA =SUMA+TMP*W (J)*TMP
TMP=FC(J)-FOPT (J)
SUMC=SUMC+TMP*W (J)*TMP
TMP=SUM-FOPT (J)
SUMD=SUMD+TMP*W(J)*TMP
6 0 ~ C O N T I N U E ~
IF(NPRIR.LT.1) GO TO }7
DO 70 J=1,NPRIR
I=IPR(J)
TMP=B(I)-BOPT(I)
TMP=TMP*WP(J)*TMP
SUMC=SUMC+TMP
70 SUMD=SUMD+TMP
75 WRITE(IOUT,9) M
CALL PRTOTB(FL,NOBS)
WRITE(IOUT,11) SUMC,SUMD
80 SUMB=SUMB+SUMD*SUMD
TMP=NRES
BN=TMP*VAR*SUMA/SUMB

```

Program Listing for General Version-Continued
WRITE(IOUT,10) BN
STOP
END
SUBROUTINE PRTOTB(VAL,NO)
c**PRINT VALUES IN THREE GROUPS OF TWO COLUMNS
DIMENSION VAL(NO)
COMMON/ITP/IIN, IOUT
\(\mathrm{NR}=\mathrm{NO} / 3\)
IF ( \(3 *\) NR.NE.NO) \(N R=N R+1\)
DO \(10 \mathrm{~K}=1\), NR
WRITE(IOUT, 20) (L,VAL(L), L=K,NO,NR)
10 continue
RETURN
20 FORMAT (1H ,2X,3(I3,7X,G11.5,3X))
END
SUBROUTINE PRTOTC(IVAL,NO)
C**PRINT INTEGERS IN THREE GROUPS OF TWO COLUMNS
DIMENSION IVAL(NO)
COMMON/ITP/IIN, IOUT
\(\mathrm{NR}=\mathrm{NO} / 3\)
IF ( \(3 *\) NR.NE.NO) \(N R=N R+1\)
DO \(10 \mathrm{~K}=1\), NR
WRITE(IOUT, 20) (L, IVAL(L), L=K,NO,NR)
10 CONTINUE
RETURN
20 FORMAT ( \(1 \mathrm{H}, 2 \mathrm{X}, 3(\mathrm{I} 3,8 \mathrm{X}, \mathrm{I4}, 9 \mathrm{X})\) )
END
SUBROUTINE PRTOT(C,NR,NC,NRD)
C**PRINT MATRICES DIVIDED VERTICALLY INTO TEN-COLUMN BLOCKS
DIMENSION C(NRD,NC)
COMMON/ITP/IIN, IOUT
DO \(60 \mathrm{~K}=1, \mathrm{NC}, 10\)
J10 \(=\mathrm{K}+9\)
IF(J10.GT.NC) J10=NC
WRITE(IOUT, 70) (J, J=K, J10)
WRITE(IOUT,90)
DO \(30 \mathrm{I}=1\), NR
\(30 \operatorname{WRITE}(\) IOUT, 80\() \mathrm{I},(\mathrm{C}(\mathrm{I}, \mathrm{J}), \mathrm{J}=\mathrm{K}, \mathrm{J} 10)\)
60 CONTINUE
70 FORMAT(1HO,10(9X,13))
80 FORMAT (1H , I3,1X,10(1X,G11.5))
90 FORMAT (1H)
RETURN
END

\section*{Listing of Inserts to the Regression Ground-Water Flow Program.}
```

C
C**INSERT JUST BEFORE EQUIVALENCE STATEMENT FOR MODIFIED BEALE'S MEASURE
DIMENSION BOPT(20),HOPT(70)
EQUIVALENGE (P(1),BOPT(1)),(HO(1),HOPT(1))
C
C
C**INSERT AFTER STATEMENT LABEL 260 FOR MODIFIED BEALE'S MEASURE
DO 1000 J=1,NVAR
1000 BOPT (J)=B(J)
DO 1100 I=1,NOBS
K=KOBS (I)
1100 HOPT(I)=BK(I)*HC(K)+BL(I)*HC(K+1)+BM(I)*HC(K+ID)+BN(I)*HC(K+ID+1)
SUMA=0.
SUMB=0.
GO TO }64
C
C
C**INSERT JUST BEFORE STATEMENT LABEL 690 FOR MODIFIED BEALE'S MEASURE
SUMC=0.
SUMD=0.
WRITE(IOUT, 2000)
2000 FORMAT (1HO,3X,28HCOMPUTED AND LINEARIZED HEADS/1H , 3X,3HNO.,7X
1,2HHC,13X,2HHL)
DO 2200 J=1,NOBS
K=KOBS (J)
HCJ=BK(J)*HC(K)+BL(J)*HC(K+1)+BM(J)*HC(K+ID)+BN(J)*HC(K+ID+1)
HL=HOPT(J)
DO 2100 I=1,NVAR
2100 HL=HL+X(I,J)*(B(I)-BOPT (I))
TMP=HCJ-HL
SUMA=SUMA+TMP*W(J)*TMP
TMP=HCJ-HOPT(J)
SUMC=SUMC+TMP*W(J)*TMP
TMP=HL-HOPT(J)
SUMD=SUMD+TMP*W(J)*TMP
WRITE(IOUT,856) J,HCJ,HL
2200 CONTINUE
IF(NPRIR.LT.1) GO TO 2400
DO 2300 J=1,NVAR
IF(WP(J).LT.1.E-10) GO TO 2300
TMP=B(J)-BOPT (J)
TMP=TMP*WP(J)*TMP
SUMC=SUMC+TMP
SUMD=SUMD+TMP
2300 CONTINUE
2400 WRITE(IOUT,2500) SUMC,SUMD
2500 FORMAT (23HOSS((HC-HOPT)*W**.5) = ,G11.5
1/23H SS((HL-HOPT)*W**.5) = ,G11.5)
SUMB=SUMB+SUMD*SUMD
C
C
C**INSERT AFTER STATEMENT LABEL 690 FOR MODIFIED BEALE'S MEASURE

```

\title{
Listing of Inserts to the Regression Ground-Water Flow Program—Continued
}
```

    READ(IIN,812) NRES,VAR
    WRITE(IOUT,2600) NRES,VAR
    2600 FORMAT ( 8HONRES = ,I4/7H VAR = ,G11.5)
    TMP=NRES
    BLN=TMP*VAR*SUMA/SUMB
    WRITE(IOUT,2700) BLN
    2700 FORMAT (1H0,5HBN = ,G11.5)
    C

```

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Theil, H., 1963, On the use of incomplete prior information in regression analysis: American Statistical Association Journal, v. 58, no. 302, p. 401-414.

Vecchia, A.V., and Cooley, R.L., 1987, Simultaneous confidence and prediction intervals for nonlinear regression models with application to a groundwater flow model: Water Resources Research, v. 23, no. 7, p. 1237-1250.

\section*{Additional Reading}

Cooley, R.L., 1977, A method of estimating parameters and assessing reliability for models of steady-state groundwater flow, 1-Theory and numerical properties: Water Resources Research, v. 13, no. 2, p. 318-324.

1982, Incorporation of prior information on parameters into nonlinear regression groundwater flow models, 1-Theory: Water Resources Research, v. 18, no. 4, p. 965-976.

\section*{7 Answers to Exercises}

In this section answers to the exercises from the first six sections are given. For brevity, the exercises are not restated.

\section*{Problem 2.2-1}
a. (RR),(RW),(RB),(WW),(WR),(WB),(BB), (BW),(BR).
b.
\begin{tabular}{l|llll}
\multicolumn{1}{c}{} & \multicolumn{1}{c}{ R } & W & B1 & B2 \\
\cline { 2 - 6 } & RR & RW & RB1 & RB2 \\
W & WR & WW & WB1 & WB2 \\
W1 & B1R & B1W & B1B1 & B1B2 \\
B2 & B2R & B2W & B2B1 & B2B2
\end{tabular}
\[
\begin{array}{lll}
P(\mathrm{RR})=1 / 16 & P(\mathrm{WW})=1 / 16 & P(\mathrm{BB})=1 / 4 \\
P(\mathrm{RW})=1 / 16 & P(\mathrm{WR})=1 / 16 & P(\mathrm{BW})=1 / 8 \\
P(\mathrm{RB})=1 / 8 & P(\mathrm{WB})=1 / 8 & P(\mathrm{BR})=1 / 8
\end{array}
\]
c.

d. \(P(X=4)=5 / 16\)
(RR),(WB),(BW).

\section*{Problem 2.2-2}
a. The effect of increasing the class interval size is to increase the frequency value for each class. This increase tends to make the histogram appear more peaked for larger intervals. Increasing the class interval also tends to smooth out irregularities.


b. \(P(X \leq 600)=0.53\)
\(P(X>400)=0.88\)
\(P(400<X \leq 600)=0.41\)
\(P(X \leq 1300)=1.0\).
Problem 2.3-1
a. \(\bar{x}=\sum_{\text {all }} \bar{x}_{i} f_{i}^{*}=(2 \times 50+1 \times 150+1 \times 250+12\) (all \(i\) )
\(\times 350+28 \times 450+25 \times 550+21 \times 650+16\)
\(\times 750+10 \times 850+8 \times 950+2 \times 1050+3\) \(\times 1150+1 \times 1250) / 130\) \(=612\).
b. By replacing the integration in equation 2.3-3 with a summation, it is seen that the population mean \(\mu_{X}\) for this random variable is represented by
\[
\begin{aligned}
\mu_{X}= & \sum_{i=2}^{6} i f_{i} \\
= & (2 \times 1 / 4+3 \times 1 / 4+4 \times 5 / 16+5 \times 1 / 8 \\
& +6 \times 1 / 16) \\
= & 3 \frac{1}{2} .
\end{aligned}
\]

\section*{Problem 2.3-2}
a. Triangular density problem.
(i)

(ii)
\(F(x)=\int_{-\infty}^{x} f(s) d s= \begin{cases}1 & , x>1 \\ 2 \int_{0}^{x} s d s=x^{2}, & , 0 \leq x \leq 1 \\ 0 & , x<0\end{cases}\)
(iii)

\[
\begin{gathered}
E[X]=L_{\infty}^{\infty} s f(s) d s=2 \int_{0}^{1} s^{2} d s=\frac{2}{3} \\
E\left[X^{2}\right]=\int_{-\infty}^{\infty} s^{2} f(s) d s=2 \int_{0}^{1} s^{3} d s=\frac{1}{2} \\
\sigma_{X}^{2}=E\left[X^{2}\right]-(E[X])^{2}=\frac{1}{2}-\frac{4}{9}=\frac{1}{18} .
\end{gathered}
\]
b. \(\bar{x}=1.25 \times 0.051+1.75 \times 0.103+2.25 \times 0.103\)
\(+2.75 \times 0.154+3.25 \times 0.154+3.75 \times 0.128\)
\(+4.25 \times 0.128+4.75 \times 0.090+5.25 \times 0.026\) \(+5.75 \times 0.064=3.36\).
\[
s_{x}^{* 2}=\sum_{i=1}^{10}\left(\bar{x}_{i}-\bar{x}\right)^{2} f_{i}^{*}
\]
\[
\begin{aligned}
& =(-2.11)^{2} \times 0.051+(-1.61)^{2} \times 0.103 \\
& +(-1.11)^{2} \times 0.103+(-0.61)^{2} \times 0.154 \\
& +(-0.11)^{2} \times 0.154+(0.39)^{2} \times 0.128 \\
& +(0.89)^{2} \times 0.128+(1.39)^{2} \times 0.090 \\
& +(1.89)^{2} \times 0.026+(2.39)^{2} \times 0.064=1.43 .
\end{aligned}
\]

\section*{Problem 2.4-1}
a. \(P(X=3\) and \(Y=2)=P(X=3) P(Y=2)=1 / 36\).
b. \(P(X+Y=5)=1 / 9\).
c. \(P(Y=2 \mid X=3)=1 / 6\).
d. \(P(X+Y=5 \mid X=3)=1 / 6\).
e. \(P(X+Y \leq 5 \mid X=3)=1 / 3\).

\section*{Problem 2.4-2}
a. \(E[Y]=a_{1} \mu_{X_{1}}+a_{2} \mu_{X_{2}}+a_{3} \mu_{X_{3}}\)
\[
\begin{aligned}
\operatorname{Var}[Y] & =E\left\{[Y-E(Y)]^{2}\right\}=E\left\{\left[a_{1}\left(X_{1}-\mu_{X_{1}}\right)\right.\right. \\
& \left.\left.+a_{2}\left(X_{2}-\mu_{X_{2}}\right)+a_{3}\left(X_{3}-\mu_{X_{3}}\right)\right]^{2}\right\} \\
& =E\left[a_{1}^{2}\left(X_{1}-\mu_{X_{1}}\right)^{2}+a_{2}^{2}\left(X_{2}-\mu_{X_{2}}\right)^{2}\right. \\
& +a_{3}^{2}\left(X_{3}-\mu_{X_{3}}\right)^{2} \\
& +2 a_{1} a_{2}\left(X_{1}-\mu_{X_{1}}\right)\left(X_{2}-\mu_{X_{2}}\right) \\
& +2 a_{1} a_{3}\left(X_{1}-\mu_{X_{1}}\right)\left(X_{3}-\mu_{X_{3}}\right) \\
& \left.+2 a_{2} a_{3}\left(X_{2}-\mu_{X_{2}}\right)\left(X_{3}-\mu_{X_{3}}\right)\right] \\
& =a_{1}^{2} \sigma_{X_{1}}^{2}+a_{2}^{2} \sigma_{X_{2}}^{2}+a_{3}^{2} \sigma_{X_{3}}^{2} \\
& +2 a_{1} a_{2} \sigma_{X_{1} X_{2}}+2 a_{1} a_{3} \sigma_{X_{1} X_{3}} \\
& +2 a_{2} a_{3} \sigma_{X_{2} X_{3}}
\end{aligned}
\]
because \(\sigma_{X_{i}}^{2}=E\left[\left(X_{i}-\mu_{X_{i}}\right)^{2}\right]\) and \(\sigma_{X_{i} X_{j}}=E\left[\left(X_{i}-\mu_{X_{i}}\right)\right.\) \(\left.\cdot\left(X_{j}-\mu_{X}\right)\right]\).
b. \(\underline{a} \operatorname{Var}\left[\underline{X} \underline{a}^{T}=\underline{a} E\left[(\underline{X}-E[\underline{X}])(\underline{X}-E[\underline{X}])^{T} \underline{a}^{T}\right.\right.\)
\[
=\underline{a} E\left[\underline{X X}^{T}-\underline{X} E[\underline{X}]^{T}-E[\underline{X}] \underline{X}^{T}\right.
\]
\(\left.+E[X] E[X]]^{T}\right] \underline{a}^{T}\)
\(=\underline{a}\left(E\left[X X^{T}\right]-E[\underline{X}] E\left[\underline{X}^{T}\right)^{T} \underline{a}^{T}\right.\)
\(=E\left[a X X^{T} \underline{a}^{T}\right]-E[a X] E\left[\underline{X}^{T} \underline{a}^{T}\right]\)
\(=\operatorname{Var}[a X]\)
c. \(\operatorname{Var}[\underline{Y}]=\left|\begin{array}{cccc}\sigma_{Y_{1}}^{2} & \sigma_{Y_{2} Y_{2}} & \cdots & \sigma_{Y_{1} Y_{p}} \\ & \sigma_{Y_{2}}^{2} & \cdots & \sigma_{Y_{2} Y_{p}} \\ \text { symmetry } & & \vdots \\ & & & \sigma_{Y_{p}}^{2}\end{array}\right|\)
\[
\begin{aligned}
\sigma_{r_{i}}^{2} & =\operatorname{Var}\left[Y_{i}\right]=\operatorname{Var}\left[a_{i} X\right]=a_{i} \operatorname{Var}[X] a_{i}^{T} \\
\sigma_{Y_{i}} Y_{j} & =E\left[Y_{i} Y_{j}\right]-E\left[Y_{i}\right] E\left[Y_{j}\right] \\
& =E\left[a_{i} X X^{T} \underline{a}_{j}^{T}\right]-E\left[a_{i} X\right] E\left[X^{T} \underline{a}_{j}^{T}\right] \\
& =\underline{a}_{i}\left(E\left[X X X^{T}\right]-E\left[\underline{\left.X] E\left[X^{T}\right]\right) a_{j}^{T}}\right.\right. \\
& =\underline{a}_{i} \operatorname{Var}[\underline{X}] \underline{a}_{j}^{T}
\end{aligned}
\]
\[
\begin{aligned}
\operatorname{Var}[\underline{Y}] & =\left[\left.\begin{array}{l}
\underline{a}_{1} \\
\underline{a}_{2} \\
\vdots \\
\underline{a}_{p}
\end{array} \right\rvert\, \operatorname{Var}[X]\left[a_{1}^{T}, \underline{a}_{2}^{T}, \ldots, \underline{a}_{p}^{T}\right]\right. \\
& =A \operatorname{Var}\left[X \left[\underline{A}^{T}\right.\right.
\end{aligned}
\]

\section*{Problem 2.5-1}
a. \(\bar{x}=\frac{1}{n} \sum_{i=1}^{n} x_{i}=610\).

The results differ very little, indicating the validity of approximation in equation \(2.3-1\). The repeat values probably result from measurement errors, indicating that a value of the random variable could not be determined more precisely than one \(\mu \mathrm{mho} / \mathrm{cm}\).
b. \(\bar{x}=3.36\)
\(s_{X}^{2}=1.43\)
\[
\begin{aligned}
\bar{t}_{g} & =10^{\bar{x}}=2,290 \mathrm{gal} / \mathrm{d} / \mathrm{ft} \\
d_{g} & =10^{\left(\bar{x} \pm s_{X}\right)}=\bar{t}_{g} \times 10^{ \pm s_{X}} \\
& =(146,35900) \mathrm{gal} / \mathrm{d} / \mathrm{ft} .
\end{aligned}
\]

The large dispersion strongly suggests that \(\bar{t}_{g}\) may not represent the true geometric mean, since there is significant scatter in \(T\) in the vicinity of \(\bar{t}_{g}\).

\section*{Problem 2.5-2}

Let \(X\) be the specific conductance random variable and \(Y\) be the dissolved solids random variable. Then
\[

\]

\section*{Problem 2.6-1}
\[
\begin{aligned}
& \bar{x}=3 \times 10^{-4} \quad s_{X}^{2}=6.6 \times 10^{-5} \\
& P\left(\bar{X}^{2} /\left(S_{X}^{2} / n\right) \leq F_{\alpha}(1, n-1)\right) \\
& \left.=P\left(-\sqrt{F_{\alpha}(1, n-1)} \leq \bar{X} /\left(S_{X} / \sqrt{n}\right) \leq \sqrt{F_{\alpha}(1, n-1}\right)\right) \\
& =P\left(-2.26 \leq \bar{X} /\left(S_{X} / \sqrt{n}\right) \leq 2.26\right)=0.95 \\
& \bar{x} /\left(s_{X} / \sqrt{n}\right)=0.12 .
\end{aligned}
\]

The value of the statistic is well within the interval ( \(-2.26<0.12<2.26\) ). This result is expected, if the titration experiment is valid, as 95 percent of all values of the statistic \(\bar{X} /\left(S_{X} / \sqrt{n}\right)\), calculated from repeated random sampling, would be expected to fall in this interval. If the value had fallen outside the interval, one should feel uneasy because this should occur only 5 percent of the time. One would then be obligated to question whether the assumption \(\mu_{X}=0\) inherent to the titration test is valid.

\section*{Problem 2.8-1}
a. \(\bar{x}=10 \quad s_{X}^{2}=0.08 \quad n=7\)
\[
\begin{aligned}
& P \left\lvert\, \frac{\left(\bar{X}-\mu_{X}\right)^{2}}{S_{X}^{2} / n} \leq F_{\alpha}(1,6)\right. \\
= & P\left(-\sqrt{F_{\alpha}(1,6)} \leq \frac{\bar{X}-\mu_{X}}{S_{X} / \sqrt{n}} \leq \sqrt{F_{\alpha}(1,6)}\right) .
\end{aligned}
\]

Therefore, an interval can be constructed from
\[
-\sqrt{F_{\alpha}(1,6)} \leq \frac{\bar{x}-\mu_{X}}{s_{X} / \sqrt{n}} \leq \sqrt{F_{\alpha}(1,6)}
\]
where \(F_{\alpha}(1,6)=5.99\). Thus
\[
-s_{X} \sqrt{\frac{F_{\alpha}(1,6)}{n}} \leq\left(\bar{x}-\mu_{X}\right) \leq s_{X} \sqrt{\frac{F_{\alpha}(1,6)}{n}}
\]
or
\[
\bar{x}+s_{X} \sqrt{\frac{F_{\alpha}(1,6)}{n}} \geq \mu_{X} \geq \bar{x}-s_{X} \sqrt{\frac{F_{\alpha}(1,6)}{n}}
\]
and
\(10.26 \geq \mu_{X} \geq 9.74\).
b. As 95 percent of all intervals so constructed will contain \(\mu_{X}\), there is a 0.95 probability that this interval contains \(\mu_{X}\).

\section*{Problem 2.9-1}
a. \(H_{0}: \mu_{X}=0\)
\(H_{1}: \mu_{X} \neq 0\).
b. From equation 2.6-21, the statement \(P\left(\right.\) reject \(H_{0} / H_{0}\) true \()=\alpha\) becomes
\[
P\left(\left.\frac{S_{1} / \sigma_{1}^{2}}{S_{2}^{2} / \sigma_{2}^{2}}>a \right\rvert\, \frac{\sigma_{1}^{2}}{\sigma_{2}^{2}}=1\right)=P\left(S_{1}^{2} / S_{2}^{2}>c\right)=0.05
\]
where \(S_{1}^{2} / S_{2}^{2}\) is an \(F(24,15)\) random variable. Because the critical region is defined by values of \(F(24,15)\) greater than or equal to \(c, c\) must be equal to \(F_{0.05}(24,15)\) (see equation 2.6-15). Thus, \(c=2.29\) and, because \(s_{1}^{2} / s_{2}^{2}=1.31\), we accept \(H_{0}\) at 0.05 significance level.
c. Hypothesis to be tested for rejection:
\[
H_{0}: \mu_{X}=9.8
\]

Alternate hypothesis:
\[
H_{1}: \mu_{X} \neq 9.8
\]

From equation 2.9-7,
\[
\left.\begin{array}{rl} 
& P\left(\frac{\bar{X}-\mu_{0}}{S_{X} / \sqrt{n}}<-\sqrt{F_{\alpha}\left(\nu_{1}, \nu_{2}\right)}\right) \\
+ & P\left(\frac{\bar{X}-\mu_{0}}{S_{X} / \sqrt{n}}>\sqrt{F_{\alpha}\left(\nu_{1}, \nu_{2}\right)}\right.
\end{array}\right)
\]

Thus, the critical value for the statistic \(\left(\bar{X}-\mu_{0}\right)^{2} /\left(S_{X}^{2} / n\right)\) is \(F_{\alpha}(1, n-1)\), which, at the 0.05 significance level, has a value of 5.99 . Because the statistic for the random sample in question takes on the value 3.50, we are forced to accept the possibility that \(\mu_{X}=0\) : We cannot safely reject \(H_{0}\) at the 0.05 significance level.

\section*{Problem 3.1-1}

Upon substitution of equation 4 into equation 3 , one obtains
\[
h=\beta_{1} X_{1}+\beta_{2} X_{2}+\beta_{3} X_{3}
\]

For \(\boldsymbol{n}\) observations
\[
\underline{Y}=\underline{X} \underline{\beta}+\underline{\epsilon}
\]
where \(\underline{Y}-\underline{h}=\underline{\epsilon}\),
\[
\underline{Y}=\left[\begin{array}{c}
Y_{1} \\
\vdots \\
Y_{n_{s}}
\end{array}\right] \underline{h}=\left[\begin{array}{r}
h_{1} \\
\vdots \\
h_{n_{s}}
\end{array}\right] \quad \underline{\beta}=\left[\begin{array}{r}
\beta_{1} \\
\vdots \\
\beta_{3}
\end{array}\right] \quad \underline{\epsilon}=\left[\begin{array}{r}
\epsilon_{1} \\
\vdots \\
\epsilon_{n_{s}}
\end{array}\right]
\]
and
\(X=\left[\begin{array}{lll}X_{11} & X_{12} & X_{13} \\ \vdots & \vdots & \vdots \\ X_{n_{s} 1} & X_{n_{s} 2} & X_{n_{s} 3}\end{array}\right]\).
In the regression model:
\(\underline{Y}=\) observed dependent variable vector;
\(\underline{\underline{h}}=\) computed dependent variable vector;
\(s=\) independent variable (distance along stream tube);
\(X=\) sensitivities; and \(\bar{\beta}=\) parameters.
a. \(\quad S(\underline{b})=[\underline{Y}-X b]^{T}[\underline{Y}-\underline{X b}]\)
where
\[
\underline{b}=\left(\begin{array}{l}
b_{1} \\
b_{2} \\
b_{3}
\end{array}\right)
\]

That is,
\[
\begin{aligned}
S(\underline{b}) & =\left[Y_{1}-\left(b_{1} X_{11}+b_{2} X_{12}+b_{3} X_{13}\right)\right]^{2} \\
& +\left[Y_{2}-\left(b_{1} X_{21}+b_{2} X_{22}+b_{3} X_{23}\right)\right]^{2} \\
& +\cdots \\
& +\left[Y_{n_{s}}-\left(b_{1} X_{n_{s} 1}+b_{2} X_{n_{s} 2}+b_{3} X_{n_{s} 3}\right)\right]^{2}
\end{aligned}
\]
b. \(\quad S(\underline{b})=[\underline{Y}-\underline{X} b]^{T} \omega[Y-\underline{X} b]\)
where \(\underline{\omega}=\underline{V}^{-1}\) and
\[
\underline{Y}=\left|\begin{array}{llc}
\sigma_{1}{ }^{2} / \sigma^{2} & & \underline{0}{ }^{2} /{ }^{2} \\
& \sigma_{2} / \sigma^{2} & \\
\underline{0} & & \sigma_{n_{s}}^{2} / \sigma^{2}
\end{array}\right|
\]
so that
\[
\underline{V}^{-1}=\left|\begin{array}{ccc}
\sigma_{1}^{-2} & \sigma_{2}^{-2} & \\
\\
0 & & \ddots \\
\underline{0} & & \sigma_{n_{s}}^{-2}
\end{array}\right| \sigma^{2} .
\]

That is,
\[
\begin{aligned}
S(\underline{b}) & =\sigma^{2}\left\{\left[\frac{1}{\sigma_{1}}\left(Y_{1}-\left(b_{1} X_{11}+b_{2} X_{12}+b_{3} X_{13}\right)\right)\right]^{2}\right. \\
& +\left[\frac{1}{\sigma_{2}}\left(Y_{2}-\left(b_{1} X_{21}+b_{2} X_{22}+b_{3} X_{23}\right)\right)\right]^{2} \\
& +\ldots \\
& \left.+\left[\frac{1}{\sigma_{n_{s}}}\left(Y_{n_{s}}-\left(b_{1} X_{n_{s} 1}+b_{2} X_{n_{s} 2}+b_{3} X_{n_{s} 3}\right)\right)\right]^{2}\right\} .
\end{aligned}
\]
c. \(S(\underline{b})=[\underline{Y}-\underline{X} b]^{T} \underline{\omega}[\underline{Y}-\underline{X} b]\)
where
\[
\underline{Y}-\underline{X} \underline{b}=\left|\begin{array}{ccc}
Y_{1}-\left(b_{1} X_{11}\right. & +b_{2} X_{12} & \left.+b_{3} X_{13}\right) \\
Y_{n_{3}}-\left(b_{1} X_{n_{3} 1}+b_{2} X_{n_{3} 2}+b_{3} X_{n_{3}}\right) \\
h_{b} & -b_{2}
\end{array}\right|
\]
and \(\underline{\omega}=\underline{V}^{-1}\), so that
\[
\underline{V}=\left[\begin{array}{ccc}
1 & & \\
\ddots & \underline{0} \\
\underline{0} & & \sigma_{h_{b}}^{2} / \sigma^{2}
\end{array}\right]
\]
and
\[
\underline{V}^{-1}=\left[\begin{array}{lll}
1 & \ddots & \underline{0} \\
0 & 1 & \sigma^{2} / \sigma_{h_{b}}^{2}
\end{array}\right]
\]

That is,
\[
\begin{aligned}
S(\underline{b}) & =\left[Y_{1}-\left(b_{1} X_{11}+b_{2} X_{12}+b_{3} X_{13}\right)\right]^{2} \\
& +\ldots \\
& +\left[Y_{n_{s}}-\left(b_{1} X_{n_{s} 1}+b_{2} X_{n_{s} 2}+b_{3} X_{n_{3}}\right)\right]^{2} \\
& +\sigma^{2}\left[h_{b}-b_{2}\right]^{2} / \sigma_{h_{b}}^{2} .
\end{aligned}
\]

\section*{Problem 3.2-1}
a. Note that if \(\underline{\omega}\) is diagonal
\[
\begin{aligned}
\underline{X}^{T} \underline{\omega} X & =\left\{\sum_{i=1}^{n_{s}} X_{i k} \omega_{i} X_{i j}+\delta_{k 2} \delta_{2 j} \sigma^{2} / \sigma_{h_{b}}^{2}\right\} \\
& =\left\{\sum_{i=1}^{n_{s}+1} X_{i k} \omega_{i} X_{i j}\right\}
\end{aligned}
\]
and
\[
\begin{aligned}
\underline{X}^{T} \underline{\omega} Y & =\left\{\sum_{i=1}^{n_{s}} X_{i k} \omega_{i} Y_{i+\delta_{k 2} h_{b} \sigma^{2} / \sigma_{h_{b}}^{2}}\right\} \\
& =\left\{\begin{array}{l}
\sum_{i=1}^{n_{s}+1} \\
i k \omega_{i} Y_{i}
\end{array}\right\}
\end{aligned}
\]
where
\[
\begin{aligned}
& \delta_{i j}=\left\{\begin{array}{llll}
1 & i=j \\
0 & i \neq j
\end{array}\right. \\
& \underline{X}^{T}=\left[\begin{array}{lllll}
X_{11} & X_{21} & & X_{n_{1} 1} & 0 \\
X_{12} & X_{22} & \ldots & X_{n_{s} 2} & 1 \\
X_{13} & X_{23} & & X_{n_{s} 3} & 0
\end{array}\right] \\
&=\left[\begin{array}{llll}
X_{11} & X_{21} & X_{n_{s} 1} & X_{n_{s}+1,1} \\
X_{12} & X_{22} & \ldots & X_{n_{s} 2} \\
X_{13} & X_{n_{s}+1,2} & X_{n_{s} 3} & X_{n_{s}+1,3}
\end{array}\right] \\
& \underline{\underline{\omega}}=\left[\begin{array}{lll}
1 & \ddots & \underline{0} \\
\underline{0} & & \sigma^{2} / \sigma_{h_{b}}^{2}
\end{array}\right] \\
&=\left[\begin{array}{llll}
\omega_{1} & & \underline{0} \\
& \ddots & & \\
\underline{0} & & \omega_{n_{s}} & \omega_{n_{s}+1}
\end{array}\right]
\end{aligned}
\]
and
\[
\left.\underline{Y}=\left|\begin{array}{l}
Y_{1} \\
Y_{2} \\
\ldots \\
Y_{n_{s}} \\
h_{b}
\end{array}\right|=\left\lvert\, \begin{array}{l}
Y_{1} \\
Y_{2} \\
\ldots \\
Y_{n_{s}} \\
Y_{n_{s+1}}
\end{array}\right.\right]
\]

Thus, the normal equations -
\[
\underline{X}^{T} \underline{\underline{\omega} X} \hat{\underline{b}}=\underline{X}^{T} \underline{\underline{\omega}}
\]
can be written
\[
\sum_{j=1}^{p} \sum_{i=1}^{n_{s}+1} X_{i k} \omega_{i} X_{i j} \hat{b}_{j}=\sum_{i=1}^{n_{s}+1} X_{i k} \omega_{i} Y_{i}, \quad k=1,2, \ldots, p
\]
or
\[
\sum_{i=1}^{n_{s}+1}\left(\sum_{j=1}^{p} X_{i k} \omega_{i} X_{i j} \hat{b}_{j}\right)=\sum_{i=1}^{n_{s}+1} X_{i k} \omega_{i} Y_{i}, k=1,2, \ldots, p
\]

As a comparison, the normal equations may be derived without using matrix techniques. Let \(n=n_{s}+1\) so that
\[
S(\underline{b})=\sum_{i=1}^{n} e_{i}^{2} \omega_{i}
\]

The model equations are given as
\[
Y_{i}=\sum_{j=1}^{p} X_{i j} b_{j}+e_{i}, \quad i=1,2, \ldots, n
\]

Therefore
\[
S(\underline{b})=\sum_{i=1}^{n}\left(Y_{i}-\sum_{j=1}^{p} X_{i j} b_{j}\right)^{2} \omega_{i}
\]
\[
A_{c}=\left[\begin{array}{lll}
\left(a_{22} a_{33}-a_{23}\right)^{2} & \left(a_{13} a_{23}-a_{12} a_{33}\right) & \left(a_{12} a_{23}-a_{13} a_{22}\right) \\
\left(a_{13} a_{23}-a_{12} a_{33}\right) & \left(a_{11} a_{33}-a_{13}\right)^{2} & \left(a_{12} a_{13}-a_{11} a_{23}\right) \\
\left(a_{12} a_{23}-a_{13} a_{22}\right) & \left(a_{12} a_{13}-a_{11} a_{23}\right) & \left(a_{11} a_{22}-a_{12}\right)^{2}
\end{array}\right]
\]

The determinant of \(\underline{A}\) is
\[
\begin{aligned}
|\underline{A}| & =a_{11}\left(a_{22} a_{33}-a_{23}^{2}\right)+a_{12}\left(a_{13} a_{23}-a_{12} a_{33}\right) \\
& +a_{13}\left(a_{12} a_{23}-a_{13} a_{22}\right) .
\end{aligned}
\]

The inverse of \(\underline{A}\) is
\[
\underline{A}^{-1}=\underline{A}_{c}^{T}|\underline{A}|
\]

Let \(\underline{A}=\underline{X}^{T} \underline{\omega} \boldsymbol{X}\) where
\[
\underline{\omega}=\left[\begin{array}{llc}
1 & & \underline{0} \\
\ddots & 1 & \underline{0} \\
\underline{0} & & \sigma^{2} / \sigma_{h_{b}}^{2}
\end{array}\right] .
\]

Thus by adding the prior information to the information in table 3, there results:
\begin{tabular}{lcc}
\hline & Data Set 1 & Data Set 2 \\
\hline\(\sum_{j} X_{j 2}^{2} \omega_{j}\) & 3.53161157 & 3.12700831 \\
\(\sum_{j} X_{j 2} \omega_{j} Y_{j}\) & 127.0992273 & 122.5185789 \\
\hline
\end{tabular}
and, for data set 1 ,
\(\underline{X}^{T} \underline{\omega} \underline{\underline{x}}=\left[\begin{array}{llr}3.3250 & 1.6750 & 418,750 \\ & 3.53161157 & 418,750 \\ \text { symmetric } & 83,340,625,000\end{array}\right]\)
or, for data set 2,
\(\underline{X}^{T} \underline{\omega} X=\left[\begin{array}{llr}2.8500 & 1.6500 & 412,500 \\ & 3.12700831 & 412,500 \\ \text { symmetric } & 8.3325 \times 10^{10}\end{array}\right]\).
Then, for data set 1 ,
\[
\underline{A}_{c}=\left[\begin{array}{lll}
1.18975153 \times 10^{11} & 3.575601563 \times 10^{10} & -777,456.0949 \\
\text { symmetric } & 1.017560156 \times 10^{11} & -690,937.5 \\
& & 8.93698347
\end{array}\right]
\]
and
\[
\left|\underline{X}^{T} \underline{\omega X}\right|=1.299239702 \times 10^{11} .
\]

Therefore
\[
\left[\underline{X}^{T} \underline{\omega X}\right]^{-1}=\left[\begin{array}{llr}
0.9157290438 & 0.2752072275 & -5.983931169 \times 10^{-6} \\
\text { symmetric } & 0.7831966299 & -5.318014058 \times 10^{-6} \\
& & 6.878625596 \times 10^{-11}
\end{array}\right] .
\]

Similarly, for data set 2
\[
\underline{A}_{c}=\left[\begin{array}{lll}
9.040171743 \times 10^{10} & 3.267 \times 10^{10} & -609,265.9279 \\
\text { symmetric } & 6.732 \times 10^{10} & -495,000 \\
& & 6.189473684
\end{array}\right]
\]
and
\[
\left|\underline{X}^{T} \underline{\omega X}\right|=6.022819942 \times 10^{10} .
\]

\section*{Therefore}
\[
\left[\underline{X}^{T} \underline{\omega X}^{-1}=\left[\begin{array}{lll}
1.500986553 & 0.5424369368 & -1.011595787 \times 10^{-5} \\
& 1.117748839 & -8.218741466 \times 10^{-6} \\
\text { symmetric } & & 1.027670384 \times 10^{-10}
\end{array}\right] .\right.
\]
c. By adding the prior information to the information in table 3, we obtain for \(\underline{X}^{T} \underline{\underline{\omega}}\) :
\begin{tabular}{lll}
\hline & \multicolumn{1}{c}{ Data Set 1 } & \multicolumn{1}{c}{ Data Set 2 } \\
\hline\(\sum_{j} X_{j 1} \omega_{j} Y_{j}\) & 192.18350 & 168.2030 \\
\(\sum_{j} X_{j 2} \omega_{j} Y_{j}\) & 127.0992273 & 122.5185789 \\
\(\sum_{j} X_{j 3} \omega_{j} Y_{j}\) & \(26,879,687.5\) & \(26,583,550\) \\
\hline
\end{tabular}

By evaluating \(\left(\underline{\underline{X}}^{T} \underline{\underline{\omega}}\right)^{-1} \underline{X}^{T} \underline{\omega} Y\), estimates \(\underline{b}\) are obtained as:
\begin{tabular}{lll}
\hline & \multicolumn{1}{c}{ Data Set 1 } & \multicolumn{1}{c}{ Data Sot 2 } \\
\hline\(\hat{b}_{1}\) & 50.12043881 & 50.01097198 \\
\(\hat{b}_{2}\) & 9.487418691 & 9.701194703 \\
\(\hat{b}_{3}\) & 0.00002302475114 & 0.00002342971729 \\
\hline
\end{tabular}

\section*{Problem 3.3-1}

The answers to parts \(a, b\), and \(c\) are found in the section at the end of the problem where aids in debugging the computer code are given. The authors' computer code and output are:

DIMENSION T(10), S(10),F(10), \(\mathrm{Z}(2,10), \mathrm{C}(2,2), \mathrm{G}(2), \mathrm{D}(2)\)
10 FORMAT (I5,6F10.0,I5)
20 FORMAT ( \(1 \mathrm{HI}, 35 \mathrm{HNO}\). OF OBSERVATIONS (N) ------- \(=\), I 7

\$/1H , 35HDISTANCE FROM WELL CENTER (R) -- = , G11.5
\$/1H ,35HINITIAL TRANSMISSIVITY (TO) ---- = , G11.5
\$/1H ,35HINITIAL STORAGE COEFFICIENT (SO) \(=\), G11.5
\$/1H ,35HDAMPING PARAMETER (AP) -------- \(=\), Gll. 5
\$/1H , 35HCLOSURE CRITERION (ER) --------- \(=\), G11.5
\(\$ / 1 \mathrm{H}, 35 \mathrm{HMAXIMUM}\) NO. OF ITERATIONS (ITMX) \(=, \mathrm{I} 7\) )
30 FORMAT (8F10.0)
32 FORMAT (10hOTIMES (T))
34 FORMAT (23HOOBSERVED DRAWDOWNS (S))
36 FORMAT ( 15 HOITERATION NO. , I4
\$/36H CURRENT ESTIMATES OF PARAMETERS (B))
38 FORMAT ( 28 HOSOLUTION FAILED TO CONVERGE)
40 FORMAT ( \(1 \mathrm{H}, 10(\mathrm{GIJ} .5,2 \mathrm{X})\) ))
42 FORMAT (19HOSOLUTION CONVERGED)
44 FORMAT ( 29 HOFINAL COMPUTED DRAWDOWNS (F))
46 FORMAT (16HORESIDUALS (F-S))
48 FORMAT (18HOERROR VARIANCE \(=\), G11.5)
50 FORMAT (12HOVAR(T) \(=, \operatorname{Gll} .5 / 12 \mathrm{H} \operatorname{COV}(\mathrm{T}, \mathrm{S})=, \mathrm{Gll} .5\)
\(\$ / 12 \mathrm{H} \operatorname{VAR}(\mathrm{S})=, \mathrm{G} 1.5)\)
52 FORMAT (26HOSCALED SENSITIVITIES (Z):/5H TO T)
54 FORMAT ( 5 H TO S)
READ AND PRINT INPUT DATA
\(\operatorname{READ}(5,10) \mathrm{N}, \mathrm{Q}, \mathrm{R}, \mathrm{TO}, \mathrm{SO}, \mathrm{AP}, \mathrm{ER}, \operatorname{ITMX}\)
\(\operatorname{WRITE}(6,20) \mathrm{N}, \mathrm{Q}, \mathrm{R}, \mathrm{T} 0, \mathrm{SO}, \mathrm{AP}, \operatorname{ER}, \operatorname{ITMX}\)
\(\operatorname{READ}(5,30)(\mathrm{T}(\mathrm{I}), \mathrm{I}=1, \mathrm{~N})\)
WRITE \((6,32)\)
\(\operatorname{WRITE}(6,40)(T(I), \mathrm{J}=1, \mathrm{~N})\)
\(\operatorname{READ}(5,30)(S(I), I=1, N)\)
\(\operatorname{WRITE}(6,34)\)
\(\operatorname{WRITE}(6,40)(S(I), I=1, N)\)
\(\mathrm{FU}=\mathrm{R} * \mathrm{R} / 4\).
\(\mathrm{FW}=\mathrm{Q} / 12.5664\)
DMAX=ER+1.
DO \(140 \mathrm{KNT}=1\), ITMX
UTMP \(=F U *\) S0/T0
WTMP \(=F W /\) T0
DO \(80 \mathrm{I}=1, \mathrm{~N}\)
COMPUTE NEW DRAWDOWN (F(I))
COMPUTE U AND W(U) FIRST
\(\mathrm{U}=\mathrm{UTMP} / \mathrm{T}(\mathrm{I})\)
\(\mathrm{WU}=\mathrm{W}(\mathrm{U})\)
THEN \(\mathrm{F}(\mathrm{I})\)
F(I) \(=\) WTMP*WU
COMPUTE SCALED SENSITIVITIES (Z(I,J))
TMP \(=\operatorname{EXP}(-\mathrm{U})\)
\(Z(1, I)=W T M P *(T M P-W U)\)
\(80 \mathrm{Z}(2, \mathrm{I})=-\) WTMP*TMP
CHECK FOR CONVERGENCE
IF(DMAX.LT.ER) GO TO 150
ASSEMBLE COEFFICIENT MATRIX (C(I,J)) AND
GRADIENT VECTOR (G(J))
DO \(86 \mathrm{~J}=1,2\)
D0 \(84 \mathrm{I}=\mathrm{J}, 2\)
\(\mathrm{C}(\mathrm{I}, \mathrm{J})=0\).
\(84 \mathrm{C}(\mathrm{J}, \mathrm{I})=0\).
\(86 \mathrm{G}(\mathrm{J})=0\).
DO \(110 \mathrm{~K}=1, \mathrm{~N}\)
TMP=S(K)-F(K)
DO \(100 \mathrm{~J}=1,2\)
DO \(90 \mathrm{I}=\mathrm{J}, 2\)
\(C(I, J)=Z(I, K) * Z(J, K)+C(I, J)\)
```

    90 C(J,I)=C(I,J)
    100 G(J)=2(J,K)*TMP+G(J)
    110 CONTINUE
        INVERT COEFFICIENT MATRIX
        DET=C(1,1)*C(2,2)-C(1,2)*C(2,1)
        TMP=C(1,1)
        C(1,1)=C(2,2)/DET
        C(2,2)=TMP/DET
        C(1,2)=-C(1,2)/DET
        C(2,1)=C(1,2)
            COMPUTE PARAMETER DISPLACEMENTS (D(I)), MAX. DISPLACEMENT
            (DMAX), AND PARAMETERS (TO AND SO)
        DMAX=0.
        DO 130 J=1,2
        D(J)=C(1,J)*G(1)+C(2,J)*G(2)
        TMP=ABS(D(J))
        IF(TMP.GT.DMAX) DMAX=TMP
    130 CONTINUE
        T0=T0*(1.+AP*D(1))
        SO=SO*(1.+AP*D(2))
            PRINT PARAMETERS
            WRITE(6,36) KNT
            WRITE(6,40) T0,S0
    140 CONTINUE
        WRITE (6,38)
        GO TO 160
    150 WRITE (6,42)
        PRINT DRAWDOWNS, RESIDUALS (F(I)-S(I)), AND SCALED
    SENSITIVITIES
    160 WRITE(6,44)
        WRITE(6,40) (F(I),I=1,N)
        DO 170 I=1,N
    170 F(I)=F(I)-S(I)
        WRITE(6,46)
        WRITE (6,40) (F(I), I=1,N)
        WRITE (6,52)
        WRITE(6,40) ( Z(1,I),I=1,N)
        WRITE (6,54)
        WRITE(6,40) (Z(2,I),I=1,N)
    C COVARIANCE MATRIX FOR PARAMETERS
VAR=0.
DO 180 I=1,N
180 VAR=VAR+F(I)*F(I)
VAR=VAR/(N-2.)
WRITE (6,48) VAR
C(1,1)=T0*C(1,1)*T0*VAR
C(1,2)=T0*C(1,2)*S0*VAR
C(2,2)=S0*C(2,2)*S0*VAR
WRITE(6,50) C(1,1),C(1,2),C(2,2)
STOP
END
FUNCTION W(X)
COMPUTE THE WELL FUNCTION OF X
H=0.
IF(X.GT.10.) Go то 20
W=-0.577216-ALOG(X)+X
TERM=X
DO 10 J=2,36
R.J=J
TERM=-TERM*X/RJ
TMP=TERM/RJ
W=W+TMP
IF(ABS(TMP).LT.1.E-7) GO TO 20
10 CONTINUE
20 RETURN
END

```
C
C
C
\begin{tabular}{|c|c|c|c|c|}
\hline \multicolumn{5}{|l|}{1NO. OF OBSERVATIONS ( N ) ------- \(=7\)} \\
\hline \multicolumn{5}{|l|}{PUMPING RATE (Q) -------------- = 1.1600} \\
\hline \multicolumn{5}{|l|}{DISTANCE FROM WELL CENTER (R) -- = 175.00} \\
\hline \multicolumn{5}{|l|}{INITIAL TRANSMISSIVITY (TO) ---- = . 10000} \\
\hline \multicolumn{5}{|l|}{INITIAL STORAGE COEFFICIENT \((S 0)=.50000 \mathrm{E}-03\)} \\
\hline \multicolumn{5}{|l|}{DAMPING PARAMETER (AP) --------- = 1.0000} \\
\hline \multicolumn{5}{|l|}{CLOSURE CRITERION (ER) --------- \(=10000 \mathrm{E}-01\)} \\
\hline MAXIMUM NO. OF ITERATIONS (ITMX) = & \((\) ITMX \()=10\) & & & \\
\hline \multicolumn{5}{|l|}{OTIMES (T)} \\
\hline \(480.00 \quad 1020.0 \quad 1500.0\) & 1500.0 2040.0 & 2700.0 & 3720.0 & 4920.0 \\
\hline \multicolumn{5}{|l|}{OOBSERVED DRAWDOWNS (S)} \\
\hline \(1.7100 \quad 2.2300\) 2.5400 & 2.54002 .7700 & 3.0400 & 3.2500 & 3.5600 \\
\hline \multicolumn{5}{|l|}{OITERATION NO. 1} \\
\hline \multicolumn{5}{|l|}{CURRENT ESTIMATES OF PARAMETERS (B)} \\
\hline . 11188 . 54748 E -03 & & . & & \\
\hline \multicolumn{5}{|l|}{OITERATION NO. 2} \\
\hline \multicolumn{5}{|l|}{CURRENT ESTIMATES OF PARAMETERS (B)} \\
\hline \multicolumn{5}{|l|}{. 11347 . \(55219 \mathrm{E}-03\)} \\
\hline \multicolumn{5}{|l|}{OITERATION NO. 3} \\
\hline \multicolumn{5}{|l|}{CURRENT ESTIMATES OF PARAMETERS (B)} \\
\hline \multicolumn{5}{|l|}{. 11349 . 55221E-03} \\
\hline \multicolumn{5}{|l|}{OSOLUTTON CONVERGED} \\
\hline \multicolumn{5}{|l|}{OFINAL COMPUTED DRAWDOWNS (F)} \\
\hline 1.6715 2.2521 2.5564 & \(2.5564 \quad 2.8012\) & 3.0256 & 3.2832 & 3.5086 \\
\hline \multicolumn{5}{|l|}{ORESIDUALS (F-S)} \\
\hline -.38538E-01 .22079E-01 .16407E-01 & .16407E-01 .31217E-01 & -.14393E-01 & . 3321 3E-01 & -. \(51356 \mathrm{E}-01\) \\
\hline \multicolumn{5}{|l|}{OSCALED SENSITIVITIES (Z):} \\
\hline T0 T & & & & \\
\hline -.91882 -1.4679 -1.7630 & \(-1.7630-2.0026\) & -2.2234 & -2.4779 & -2.7014 \\
\hline \multicolumn{5}{|l|}{T0 S} \\
\hline -. 75264 -.78421 -.79343 & -. 79343 -. 79866 & -. 80223 & -. 80527 & -.80724 \\
\hline \multicolumn{5}{|l|}{OERROR VARIANCE \(=.14328 \mathrm{E}-02\)} \\
\hline \multicolumn{5}{|l|}{OVAR \((T)=.95030 \mathrm{E}-05\)} \\
\hline \multicolumn{5}{|l|}{\(\operatorname{COV}(\mathrm{T}, \mathrm{S})=-.11369 \mathrm{E}-06\)} \\
\hline \(\operatorname{VAR}(S)=.14595 E-08\) & & & & \\
\hline
\end{tabular}

\section*{Problem 3.3-2}
b.


c.
\(\underline{J}_{1}=\left|\begin{array}{l}h_{2}+h_{4}-2 h_{1} \\ h_{1}+h_{5}-2 h_{2} \\ 0 \\ h_{1}+2 h_{5}-4 h_{4}+h_{B 1} \\ h_{2}+2 h_{4}+h_{8}-4 h_{5} \\ 0 \\ 0 \\ h_{5}-2 h_{8}+h_{B 1} \\ 0\end{array}\right|\)

\(J_{3}=\left[\begin{array}{l}1 / 2 a^{2} \\ 1 / 2 a^{2} \\ 0 \\ a^{2} \\ a^{2} \\ 0 \\ 0 \\ 1 / 2 a^{2} \\ 0\end{array}\right] \quad \underline{J}_{4}=\left[\begin{array}{l}0 \\ 0 \\ 0 \\ 0 \\ 0 \\ a \\ 0 \\ 0 \\ a\end{array}\right]\)
d.
1. Let \(r=0\).
2. Compute \(D_{r}\) and \(q_{r}\).
3. \(h_{r}=D_{r}^{-1} q_{r}\)
4. Obtain \(f\left(\xi, \xi_{\xi}\right)\) from \(\underline{h}_{r}\) by deleting node 7 .
5.
\[
\left.\begin{array}{l}
\left(\frac{\partial \underline{h}}{\partial b_{j}}\right)_{r}=\underline{D}_{r}^{-1}\left[\left(\frac{\partial q}{\partial b_{j}}\right)_{r}-\left(\frac{\partial \underline{D}}{\partial b_{j}}\right)_{r} \underline{h}_{r}\right] \\
\left.\left.\begin{array}{l}
\text { Obtain } \underline{X}_{j}^{r} \text { from } \\
\text { deleting node } 7 .
\end{array}\right|_{j=1,2, \ldots, p} ^{\partial b_{j}}\right)_{r}^{\text {by }}
\end{array}\right)_{j=1}
\]
6. Compute \(\underline{X}_{T}^{T} \underline{\omega} X_{r}\) and \(\underline{X}_{r}^{T} \underline{\omega}\left(\underline{Y}-f\left(\xi, \xi_{r}\right)\right)\).
7. Define \(\underline{C}_{r}=\left\{C_{i i}^{r}\right\}=\left\{\left(X_{\underline{T}}^{T} X_{r}\right)_{i i}^{-1 / 2}\right\}\).
8. Compute \(S_{-}^{T} \underline{\omega}_{r}\) and \(\underline{S}_{r}^{T} \underline{\omega}\left(\underline{Y}-f\left(\xi, \underline{b_{r}}\right)\right)\) where \(\underline{S}_{r}=X_{r} C_{r}\).
9. \(\underline{\delta}_{r+1}=\left(\underline{S}_{-}^{T} \underline{\omega} \underline{S}_{r}+\mu \underline{I}\right)^{-1} S_{r}^{T} \underline{\underline{\omega}}\left(\underline{Y}-f\left(\xi_{2} \underline{b}_{r}\right)\right)\)
10. \(\underline{d}_{r+1}=\underset{\sim}{C} \underline{\delta}_{r+1}\).
11. \(\underline{b}_{r+1}=\rho d_{r+1}+b_{r}\).
12. If \(\left|d_{i}^{r+1} / c\right|>\epsilon\) (where \(c=b_{i}^{r}\) for \(b_{i}^{r} \neq 0\), and \(c=1\) for \(b_{i}^{r}=0\) ) for any \(i=1,2, \ldots, p\), then increment \(r\) by one and return to 2 . If not then:
13. \(\underline{h}_{\text {final }}=\underline{D}_{r+1}^{-1} q_{r+1}\).

Values of \(\mu\) and \(\rho\) can be computed at each iteration by using algorithms defined by equations \(3.3-28\) through \(3.3-30\) if desired.

\section*{Problem 4.2-1}

\section*{Data Set 1}

Nodes in columns 2 through 11 have observations. Nodes in columns 1 and 12 form specified head boundaries. Spacing: Cell row 1, 1 ft ; cell columns 1 and 11, 50 ft ; cell columns 2 through 10, 100 ft (figure 1).

\section*{Data Set 2.}

Nodes in columns 2 through 10 have observations. Nodes in columns 1 and 11 form specified head boundaries. Spacing: Cell row \(1,1 \mathrm{ft}\); cell columns 1 through \(10,100 \mathrm{ft}\) (figure 2). The input data to the regression program for the two data sets are shown in figures 3 and 4.

If both \(T\) and \(W\) were estimated, the problem would be singular because the only unique parameter is \(W / T\).


Figure 1


Figure 2
LAKE OHPUPU, DATA SET 1.
\begin{tabular}{|c|c|c|c|c|c|c|c|c|c|c|c|c|}
\hline 12 & 2 & 1 & 10 & 1 & 3 & 0 & 02 & 2 & 1 & 1 & & \\
\hline & 50 & & 0 & & 0 & & 0 & . 25 & & & & \\
\hline DX & 1 & 0 & & & & & & & & & & \\
\hline 1 & 11 & 1 & 1 & & 1 & 1 & & & & & & \\
\hline & 50 & & 100 & & 100 & & 100 & 100 & & 100 & 100 & 100 \\
\hline & 100 & & 100 & & 50 & & & & & & & \\
\hline DY & 1 & 0 & & & & & & & & & & \\
\hline 1 & 1 & 1 & 1 & & 1 & & & & & & & \\
\hline CX & 1 & 0 & & & & & & & & & & \\
\hline 1 & 11 & 1 & 1 & & 1 & & & & & & & \\
\hline CY & 1 & 0 & & & & & & & & & & \\
\hline 1 & 11 & 1 & 1 & & 1 & & & & & & & \\
\hline VL & 1 & 0 & & & & & & & & & & \\
\hline 1 & 11 & 1 & 1 & & 0 & & & & & & & \\
\hline HR & 1 & 0 & & & & & & & & & & \\
\hline 1 & 12 & 1 & 2 & & 0 & & & & & & & \\
\hline QR & 1 & 0 & & & & & & & & & & \\
\hline 1 & 11 & 1 & 1 & & 1 & & & & & & & \\
\hline HC & 1 & 0 & & & & & & & & & & \\
\hline 2 & 11 & 1 & 2 & & 11 & & & & & & & \\
\hline IZN & 1 & 0 & & & & & & & & & & \\
\hline 1 & 11 & 1 & 1 & 1 & & & & & & & & \\
\hline 1 & 2 & 1 & & 50 & & . 5 & 48.33 & & 1 & & & \\
\hline 2 & 3 & 1 & & 150 & & . 5 & 45.76 & & 1 & & & \\
\hline 3 & 4 & 1 & & 250 & & . 5 & 42.08 & & 1 & & & \\
\hline 4 & 5 & 1 & & 350 & & . 5 & 38.34 & & 1 & & & \\
\hline 5 & 6 & 1 & & 450 & & . 5 & 35.30 & & 1 & & & \\
\hline 6 & 7 & 1 & & 550 & & . 5 & 31.00 & & 1 & & & \\
\hline 7 & 8 & 1 & & 650 & & . 5 & 25.85 & & 1 & & & \\
\hline 8 & 9 & 1 & & 750 & & . 5 & 21.76 & & 1 & & & \\
\hline 9 & 10 & 1 & & 850 & & . 5 & 16.11 & & 1 & & & \\
\hline 10 & 11 & 1 & & 950 & & . 5 & 12.48 & & 1 & & & \\
\hline 1 & 0 & 0 & 0 & 1 & & & & & & & & \\
\hline 1 & & 0 & & & & & & & & & & \\
\hline 1 & & 1 & & 1 & & 0 & . 001 & & & & & \\
\hline IN & 2 & 0 & & & & & & & & & & \\
\hline 1 & 1 & 1 & 2 & -1 & & & & & & & & \\
\hline 12 & 12 & 1 & 2 & -1 & & & & & & & & \\
\hline 1 & 2 & 2 & 2 & & 0 & & 0 & & & & & \\
\hline 1 & 1 & & 49 & & & & & & & & & \\
\hline 1 & 2 & & 49 & & & & & & & & & \\
\hline 2 & 2 & 3 & 3 & & 1.1 & & 1.1 & & & & & \\
\hline 12 & 1 & & 11 & & & & & & & & & \\
\hline 12 & 2 & & 11 & & & & & & & & & \\
\hline
\end{tabular}

LAKE OHPUPU, DATA SET 2.


Figure 4

\section*{Problem 4.2-2}

Because prior information is available, all parameters can be estimated. Data points are located in areas of relatively high sensitivity for
all parameters. More data points in areas of highest sensitivity might improve results for parameters having low sensitivity. The input data to the regression program are shown in figure 1.



Figure 1

\section*{Problem 5.4-1}
a. Let \(\underline{b}_{0}=\underline{0}\). Then
\[
s^{2}=\frac{\underline{Y}^{T} \underline{\underline{\omega}} \underline{\underline{Y}}-\hat{\underline{b}}^{T} \underline{X}^{T} \underline{\underline{\omega}} \underline{Y}}{n-p}
\]


\[
\begin{aligned}
& \hat{b}=\left[\begin{array}{c}
50.12043881 \\
9.487418691 \\
0.00002302475114
\end{array}\right] \text { or } \\
& 50.01097198 \\
& 9.701194703 \\
& 0.00002342971729 \\
& \begin{array}{l}
\underline{X}^{T} \underline{\omega} Y=\left[\begin{array}{l}
192.18350 \\
127.0992273 \\
26,879,687.5
\end{array}\right] \quad \text { or } \quad\left[\begin{array}{l}
168.2030 \\
122.5185789 \\
26,583,550
\end{array}\right] \\
\underline{Y}^{T} \underline{\omega} Y=11,459.5411 \quad \text { or } \quad 10,225.4391
\end{array} \\
& \underline{b}^{\underline{b^{\prime}}} \underline{\underline{T}}^{T} \underline{\underline{Y}}=11,457.06305 \text { or } \quad 10,223.41717 \\
& s^{2}=0.30975625 \quad \text { or } \quad 0.2888471429 \\
& s / \Delta Y_{s}=0.5565574993 /(48.33-12.48)=0.0155 \text { or } \\
& 0.5374450138 /(47.13-14.96)=0.0167
\end{aligned}
\]

The fit is fairly good.
b.
\(\left(\underline{X}^{T} \underline{\omega X}^{-1} s^{2}=\right.\)
\(0.30975625\left[\begin{array}{lrr}0.9157290438 & 0.2752072275 & -5.983931169 \times 10^{-6} \\ \text { symmetric } & 0.7831966299 & -5.318014058 \times 10^{-6} \\ & 6.878625596 \times 10^{-11}\end{array}\right]\)
\[
=\left[\begin{array}{rr}
0.2836527946 & 0.08524715876-1.853560079 \times 10^{-6} \\
\text { symmetric } & 0.2426000511-1.647288092 \times 10^{-6} \\
& 2.130697270 \times 10^{-11}
\end{array}\right]
\]
or
\(0.2888471429\left[\begin{array}{llr}1.500986553 & 0.5424369368 & -1.011595787 \times 10^{-5} \\ & 1.117748839 & -8.218741466 \times 10^{-6} \\ \text { symmetric } & & 1.027670384 \times 10^{-10}\end{array}\right]\)
\[
=\left[\begin{array}{rr}
0.43355567744 & 0.1566813594 \\
& 0.3 .921965528 \times 10^{-6} \\
\text { symmetric } & 2.968396543 \times 10^{-11}
\end{array}\right]
\]

Because the standard errors (the square roots of the diagonal elements) are small, the parameters are well determined.

or


The problem is fairly well conditioned.

\section*{Problem 5.5-1}

The plots of \(\underline{d}\) and \(g\) (figures 1,2 ) suggest that correlation manifests itself in two ways: The variability of \(g\) from set to set is smaller than the variability of \(d\) from set to set, and the plots of \(g\) do not appear to be linear on normal probability paper. The plot of \(\hat{e}\) (figure 3 ) does not appear to differ very much from either the plots of \(\underline{d}\) or \(g\), although the plot of \(\hat{e}\) does show the same nonlinear trend as displayed by the plots of \(g\). Therefore, one may say that the plot of \(\frac{\hat{e}}{}\) does not appear to differ significantly from the plot of a \(N\left(\underline{0},(\underline{I}-\underline{R}) s^{2}\right)\) random variable. Furthermore, this distribution of residuals suggests that the Theis model is adequate to describe the observed-drawdown data set. Input data for the residuals analysis program are shown in figure 4.

\section*{Problem 5.5-2}
a. \(\quad s^{2}=0.98677\)
\(R_{y}=0.99964\)
\(s / \Delta Y_{s}=0.99336297 / 175.18=0.00567\)
b. The problem is not well conditioned. Entries \((1,8),(1,9),(8,9)\), and \((12,13)\) (and their symmetric counterparts) of the scaled least squares matrix have absolute values greater than 0.9 , although these large values yielded large (absolute value \(>0.9\) ) parameter correlations only for entries \((1,8)\) and \((12,13)\). The correlation between
parameters 12 and 13 is exceptionally large so that these parameters are behaving as one. Because parameters 12 and 13 are the transmissivity \(T_{3}\) and recharge \(W_{3}\) in zone 3, respectively, the physical interpretation is that the ratio \(W_{3} / T_{3}\) is much more unique than either \(W_{3}\) or \(T_{3}\). The only poorly determined parameter is the flow across boundary zone \(2, q_{B 2}\). Be cause of the large value of transmissivity in the aquifer zone adjacent to this boundary flow zone, the gradient near the boundary is low. Hence, the flow across the boundary is estimated poorly. Well discharges \(Q_{1}\) and \(Q_{2}\) and transmissivity \(T_{2}\) are estimated quite well. The drawdown cones provide large gradients that serve to determine \(Q_{1}, Q_{2}\), and \(T_{2}\) precisely.
c. Effects of correlation within \(g\) are not large. They may be exhibited as a slight steepening of the curves for \(g\) (figure 2) compared to those for \(\underline{d}\) (figure 1) on the normal probability plots. The plot of \(\underline{\hat{i}}\) (figure 3) is very similar to those of \(g\), which suggests that the distribution of \(\hat{\hat{u}}\) does not differ significantly from a \(\left.N(\underline{0}, \underline{I}-\underline{R}) s^{2}\right)\) distribution. Note that \(\hat{\hat{e}}\) could not have been used instead of \(\underline{\hat{u}}\) to make the comparison because the weight matrix \(\omega\) is not equal to \(\underline{\underline{I}}\). The plot of \(\hat{u}_{j}\) versus \(\vec{f}_{j}\) (figure 4) shows no pattern. However, the plot of \(\hat{e}_{j}\) versus Cartesian coordinate (figure 5) shows a group of negative residuals in the upper center of the area. This sign pattern was inherited from the sign pattern of the original errors \(\underline{\epsilon}\) that were generated (recall that this exercise is based on a hypothetical problem), and the original errors are ran\(\operatorname{dom} N(0,1)\) deviates. Hence, the sign pattern occurred entirely by chance. The lesson is that apparently nonrandom patterns can, and often do, develop by chance, and the analyst must learn to distinguish true problem areas from apparent ones.

The input data for the residuals analysis program are shown in figure 6.
Ordered residual distribution: (The calculated values are prior information residuals \(\hat{u}_{j}=\hat{e}_{j} \omega_{j}^{1 / 2}\), where \(\left.\omega_{j}^{1 / 2} s\left(\operatorname{Var}\left(\epsilon_{p j}\right)\right)^{1 / 2}.\right)\)
\begin{tabular}{|c|c|c|c|}
\hline No. & \(\hat{u}_{j}\) & No. & \(\hat{u}^{\boldsymbol{u}}\) \\
\hline 1 & -2.3900 & 21 & -0.013826 \\
\hline 2 & -1.4247 & 22 & \((-97000+97000) 5.1204 \times 10^{-4}=\approx 0\) \\
\hline 3 & -1.3248 & 23 & 0.034151 \\
\hline 4 & -1.2986 & 24 & 0.036783 \\
\hline 5 & -1.1685 & 25 & \((-50961+51000) 9.7389 \times 10^{-4}=0.037982\) \\
\hline 6 & -0.95063 & 26 & (0.080716-0.08) \(124.17=0.088906\) \\
\hline 7 & \((0.00031149-0.0004) 8278.0=-0.73269\) & 27 & \((5.4730-5.4) 1.8396=0.13429\) \\
\hline 8 & -0.72189 & 28 & 0.22868 \\
\hline 9 & \[
-0.70275
\] & 29 & 0.24557 \\
\hline 10 & \((0.00013516-0.00017) 19478=-0.67861\) & 30 & 0.33619 \\
\hline 11 & -0.50751 & 31 & 0.40574 \\
\hline 12 & -0.43737 & 32 & 0.41665 \\
\hline 13 & -0.43737 & 33 & 0.64095 \\
\hline 14 & -0.35588 & 34 & 0.65517 \\
\hline 15 & -0.33702 & 35 & (5.1211-4.8)2.0695 \(=0.66452\) \\
\hline 16 & -0.26802 & 36 & \((487.89-420) 0.011826=0.80287\) \\
\hline 17 & \((10.198-10.4) 0.95516=-0.19294\) & 37 & 0.90021 \\
\hline 18 & -0.18186 & 38 & 0.98513 \\
\hline 19 & -0.087697 & 39 & 1.1503 \\
\hline 20 & -0.080833 & 40
41 & \[
\begin{aligned}
& 1.3703 \\
& 1.6158
\end{aligned}
\] \\
\hline & & 41 & 1.6158 \\
\hline
\end{tabular}


Figure 1


Figure 2


Figure 3
\begin{tabular}{rrrrrrrr}
2 & 7 & 0 & 5 & 1 & .0014328 & & \\
\hline
\end{tabular}

Figure 4


Figure 1


Figure 2


Figure 3

Figure 4


Figure 5
\begin{tabular}{rrrrrrrrrrrr}
14 & 32 & 9 & 4 & 437 & .98677 & & & & & & \\
& 1 & & & & & & & & & & \\
3 & 5 & 6 & 7 & 9 & 10 & 13 & 14 & & & & \\
1940 & & 1020 & & 1.04 & & .48 & & .54 & .00012 & 84 & .000051 \\
.008 & & & & & & & & & & &
\end{tabular}

Figure 6

\section*{Problem 5.6-1}
a. Compute
\(\underline{H}\left(\underline{X}^{T} \underline{\omega X}\right)^{-1} \underline{H}^{T}\)
where \(\underline{H}=\left[\begin{array}{lll}0 & 0 & 1\end{array}\right]\). Then,
\(\underline{H}\left(\underline{X}^{T} \underline{\omega X}^{-1} \underline{H}^{T}=\widehat{\operatorname{Var}}\left(\hat{b}_{3}\right) / s^{2}\right.\).
By using equation 5.6-11,
\[
\begin{aligned}
& \begin{aligned}
W & =\frac{\left(\tilde{\beta}_{3}-\hat{b}_{3}\right)\left(\operatorname{Var}\left(\hat{b}_{3}\right) / s^{2}\right)^{-1}\left(\tilde{\beta}_{3}-\hat{b}_{3}\right)}{s^{2}} \\
& =\frac{\left(\tilde{\beta}_{3}-\hat{b}_{3}\right)^{2}}{\left.\hat{\operatorname{Var}( } \hat{b}_{3}\right)}
\end{aligned} \\
& \begin{aligned}
H_{0} & : \beta_{3}=\frac{0.0003}{10}=3 \times 10^{-5} \\
H_{1} & : \beta_{3} \neq 3 \times 10^{-5} \\
w & =\frac{\left(3 \times 10^{-5}-2.3025 \times 10^{-5}\right)^{2}}{2.1307 \times 10^{-11}} \\
& =2.273 \text { for data set } 1 .
\end{aligned}
\end{aligned}
\]
\(F_{0.05}(1,8)=5.318 .^{\cdot} . H_{0}\) accepted
Values of the ratio of recharge to transmissivity, \(\beta_{3}=W / T\), computed by the two methods are not significantly different.
\[
w=\frac{\left(3 \times 10^{-5}-2.3430 \times 10^{-5}\right)^{2}}{2.9684 \times 10^{-11}}
\]
\(=1.454\) for data set 2
\(F_{0.05}(1,7)=5.591 .{ }^{\cdot} H_{0}\) accepted.
The Maxey-Eakin estimate of \(W / T\) could be used as prior information in the regression model, but an estimate of \(\operatorname{Var}(W / T)\) would also be needed.
b. This test is the same as the one in a except that
\[
\begin{aligned}
& H_{0}: \beta_{3}=0 H_{1}: \beta_{3} \neq 0 \\
& w=\frac{\left(0-2.3025 \times 10^{-5}\right)^{2}}{2.1307 \times 10^{-11}} \\
&=24.88 \text { for data set } 1
\end{aligned}
\]
\(\therefore H_{0}\) rejected.
\[
\begin{aligned}
w & =\frac{\left(0-2.3430 \times 10^{-5}\right)^{2}}{2.9684 \times 10^{-11}} \\
& =18.49 \text { for data set } 2 \\
\therefore . & H_{0} \text { rejected. }
\end{aligned}
\]
\(W / T\) is significantly different from zero. Thus, recharge is a significant variable in the regression equation.
c. \(\quad \tilde{\beta}_{3}=\hat{b}_{3} \pm \sqrt{q F_{\alpha}(q, n-p)} s_{b 3}\)
\[
\begin{aligned}
& =2.3025 \times 10^{-5} \pm \sqrt{5.318} \times 4.6159 \times 10^{-6} \\
& =2.3025 \times 10^{-5} \\
& \pm 1.0645 \times 10^{-5} \text { for data set } 1 . \\
\tilde{\beta}_{3} & =2.3430 \times 10^{-5} \pm \sqrt{5.591} \times 5.4483 \times 10^{-6} \\
& =2.3430 \times 10^{-5} \\
& \pm 1.2883 \times 10^{-5} \text { for data set } 2 .
\end{aligned}
\]

\section*{Problem 5.6-2}

For \(T\) extreme:
\[
\begin{gathered}
T=\hat{T} \pm \frac{\sqrt{2 F_{0.05}(2,5)}}{s_{b T}} \widehat{\mathrm{Var}}(T) \\
=0.11349 \pm \frac{\sqrt{2 \times 5.7861}}{0.0030827}\left(0.95030 \times 10^{-5}\right) \\
=0.11349 \pm 0.010487 \\
S=\hat{S} \pm \frac{\sqrt{2 F_{0.05}(2,5)}}{s_{b T}} \widehat{\operatorname{Cov}}(T, S) \\
=0.00055221 \pm \frac{\sqrt{2 \times 5.7861}}{0.0030827}\left(-0.11369 \times 10^{-6}\right) \\
=0.00055221 \mp 0.00012546 .
\end{gathered}
\]

For \(S\) extreme:
\[
\begin{aligned}
& T=\hat{T} \pm \frac{\sqrt{2 F_{0.05}(2,5)}}{s_{b S}} \widehat{\operatorname{Cov}(T, S)} \\
= & 0.11349 \pm \frac{\sqrt{2 \times 5.7861}}{3.8203 \times 10^{-5}}\left(-0.11369 \times 10^{-6}\right) \\
= & 0.11349 \mp 0.010124
\end{aligned}
\]
\[
S=\hat{S} \pm \frac{\sqrt{2 F_{0.05}(2,5)}}{s_{b S}} \widehat{\mathrm{Var}}(S)
\]
\[
=0.00055221 \pm \frac{\sqrt{2 \times 5.7861}}{3.8203 \times 10^{-5}}\left(0.14595 \times 10^{-8}\right)
\]
\[
=0.00055221 \pm 0.00012996 \text {. }
\]

\section*{Problem 5.6-3}
\[
\underline{\tilde{b}}=\underline{\hat{b}} \pm \frac{\sqrt{q F_{\alpha}(q, n-p)}}{s_{b i}} V_{b i}
\]
where \(\underline{V}_{b}=\left(\underline{X}^{T} \underline{\omega X}\right)^{-1} s^{2}\). Let \(\left(\underline{X}^{T} \underline{\omega X}\right)^{-1} s^{2}=A\) \(=\left\{A_{i j}\right\}\). Then the above equation can be written in algebraic form for calculations as
\[
\tilde{b}_{j}=\hat{b}_{j} \pm \sqrt{q F_{\alpha}(q, n-p)} \frac{A_{i j}}{\sqrt{A_{i i}}} .
\]

Calculations:
\[
\sqrt{q F_{\alpha}(q, n-p)}=\sqrt{2 F_{0.05}(2,27)}=\sqrt{2 \times 3.359}=2.592
\]

For \(T_{3}=\beta_{12}\) extreme:
\[
\begin{aligned}
& \tilde{b}_{1}=80.978 \pm 9.2811 \\
& \tilde{b}_{2}=935.15 \pm 208.42 \\
& \tilde{b}_{3}=-97000 \mp 85.314 \\
& \tilde{b}_{4}=-50961 \mp 29.377 \\
& \tilde{b}_{5}=10.198 \pm 0.11520 \\
& \tilde{b}_{6}=5.1211 \pm 1.3758 \times 10^{-2} \\
& \tilde{b}_{7}=5.4730 \pm 1.2758 \times 10^{-2} \\
& \tilde{b}_{8}=65.754 \pm 7.2341 \\
& \tilde{b}_{9}=3.1149 \times 10^{-4} \pm 3.0988 \times 10^{-5}
\end{aligned}
\]
\[
\begin{aligned}
& \tilde{b}_{10}=487.89 \pm 2.5660 \\
& \tilde{b}_{11}=-1.3995 \times 10^{-4} \mp 4.9215 \times 10^{-5} \\
& \tilde{b}_{12}=13.288 \pm 9.3998 \\
& \tilde{b}_{13}=1.3516 \times 10^{-4} \pm 9.3069 \times 10^{-5} \\
& \tilde{b}_{14}=8.0716 \times 10^{-2} \pm 3.5370 \times 10^{-5}
\end{aligned}
\]

For \(q_{B 1}=\beta_{1}\) extreme:
\[
\begin{aligned}
\tilde{b}_{1} & =80.978 \pm 43.911 \\
\vec{b}_{2} & =935.15 \pm 489.14 \\
\vec{b}_{3} & =-97000 \mp 222.36 \\
b_{4} & =-50961 \mp 73.632 \\
\tilde{b}_{5} & =10.198 \pm 0.27770 \\
\tilde{b}_{6} & =5.1211 \pm 3.2747 \times 10^{-2} \\
\tilde{b}_{7} & =5.4730 \mp 3.2200 \times 10^{-3} \\
\tilde{b}_{8} & =65.754 \pm 25.550 \\
\tilde{b}_{9} & =3.1149 \times 10^{-4} \pm 5.4093 \times 10^{-5} \\
\tilde{b}_{10} & =487.89 \pm 6.2288 \\
\tilde{b}_{11} & =-1.3995 \times 10^{-4} \mp 1.1701 \times 10^{-4} \\
\tilde{b}_{12} & =13.288 \pm 1.9832 \\
\tilde{b}_{13} & =1.3516 \times 10^{-4} \pm 2.0377 \times 10^{-5} \\
\tilde{b}_{14} & =8.0716 \times 10^{-2} \pm 1.5001 \times 10^{-4} .
\end{aligned}
\]

\section*{Problem 5.7-1}
a. Let \(\underline{X}=\left\{X_{i j}\right\},\left(\underline{X}^{T} \underline{\omega X}\right)^{-1}=\left\{A_{i j}\right\}\). Then
\[
\begin{gathered}
\underline{X}\left(\underline{X}^{T} \underline{\omega X}\right)^{-1}=\left\{\sum_{k=1}^{p} X_{i k} A_{k \ell}\right\}=\left\{C_{i \ell}\right\} \\
\underline{X}\left(\underline{X}^{T} \underline{\omega X}^{-1} \underline{X}^{T}=\left\{\sum_{\ell=1}^{p} C_{i k} X_{j \ell}\right\}=\right. \\
\left\{\sum_{\ell=1}^{p} \sum_{k=1}^{p} X_{i k} A_{k \ell} X_{j \ell}\right\}=\left\{\sum_{\ell=1}^{p} X_{j \ell} \sum_{k=1}^{p} X_{i k} A_{k \ell}\right\} .
\end{gathered}
\]

If \(i=j\), then the entry is
\(\sum_{\ell=1}^{p} \sum_{k=1}^{p} X_{i k} A_{k \ell} X_{i \ell}=\sum_{\ell=1}^{p} X_{i \ell} \sum_{k=1}^{p} X_{i k} A_{k \ell}\).
Compute for \(\boldsymbol{i = j = 1}\) for data set 1 .
\(\underline{X}_{1}=\left[\begin{array}{lll}0.95 & 0.05 & 23,750\end{array}\right]\).
\[
\left(\underline{X}^{T} \underline{\omega}\right)^{-1}=\left[\begin{array}{rr}
0.91572904380 .2752072275 & -5.983931169 \times 10^{-6} \\
\text { symmetric } & 0.7831966299 \\
-5.318014058 \times 10^{-6} \\
& 6.878625596 \times 10^{-11}
\end{array}\right]
\]
\(\underline{X}_{1}\left(\underline{X}^{T} \underline{\omega X}^{-1} \underline{X}_{1}^{T}=0.95(0.95 \times 0.9157290438\right.\)
\(+0.05 \times 0.2752072275\)
\(\left.+23,750 \times\left(-5.983931169 \times 10^{-6}\right)\right)\)
\(+0.0510 .95 \times 0.2752072275\)
\(+0.05 \times 0.7831966299\)
\(\left.+23,750 \times\left(-5.318014058 \times 10^{-6}\right)\right)\)
\(+23,750 \times\left(0.95 \times\left(-5.983931169 \times 10^{-6}\right)\right.\)
\(+0.05 \times\left(-5.318014058 \times 10^{-6}\right)\)
\(\left.\left.+23,750 \times 6.878625596 \times 10^{-11}\right)\right)\)
\(=0.6106927103\).
\[
\begin{aligned}
& \underline{X}\left(\underline{X}^{T} \underline{\omega X}^{-1} \underline{X}^{T}=\underline{X C C^{-1}}\left(\underline{X}^{T} \underline{\omega X}\right)^{-1}\left(\underline{C C^{-1}}\right)^{T} \underline{X}^{T}\right. \\
& =\underline{X C C}^{T} \underline{X}^{T} \underline{\omega X C}^{-1} \underline{C}^{T} \underline{\underline{X}}^{T} \\
& \left.=\underline{X C})(\underline{X C})^{T} \underline{\omega(\underline{X} C}\right)^{-1}(\underline{X C})^{T} \\
& =\underline{S}\left(\underline{S}^{T} \omega \bar{S}^{-1} \underline{S}^{T}\right. \text {. }
\end{aligned}
\]
b. \(\widehat{\operatorname{Var}}\left(\hat{f}_{j}\right)=\underline{X}_{j} \widehat{\operatorname{Var}(\underline{b}} \underline{X}_{j}^{T}\)

Let \(j=1\), data set 1 .
\[
\begin{aligned}
\widehat{\operatorname{Var}}\left(\hat{f}_{1}\right) & =\underline{X}_{1}\left(\underline{X}^{T} \underline{\omega X}^{-1} s^{2} \underline{X}_{1}^{T}\right. \\
& =0.6106927103 \times 0.30975625 \\
& =0.1891658838
\end{aligned}
\]
c. \(\quad f_{\beta j}=\hat{f}_{j} \pm \sqrt{p F_{\alpha}(p, n-p)} s_{y j}\) \(F_{0.05}(3,8)=4.066\)
\(\hat{f}_{1}=0.95 \hat{b}_{1}+0.05 \hat{b}_{2}+23,750 \hat{b}_{3}\)
\(=0.95(50.1204)+0.05(9.48742)\)
\begin{tabular}{|c|c|c|c|c|c|c|}
\hline 22 & 0 & \multicolumn{2}{|l|}{2.0014328} & \multirow[b]{3}{*}{3.0256} & \multirow[b]{3}{*}{3.2832} & \multirow[b]{3}{*}{3.5086} \\
\hline . 11349 & . 00055221 & & & & & \\
\hline 1.6715 & 2.2521 & 2.5564 & 2.8012 & & & \\
\hline 1 & 1 & 1 & 1 & 1 & 1 & 1 \\
\hline -8.09604 & -1362.96 & & & & & \\
\hline -12.9342 & -1420.13 & & & & & \\
\hline -15.5344 & -1436.83 & & & & & \\
\hline -17.6456 & -1446.30 & & & & & \\
\hline -19.5912 & -1452.76 & & & & & \\
\hline -21.8336 & -1458.27 & & & & & \\
\hline -23.8030 & -1461.84 & & & & & \\
\hline . 10300 & . 00067767 & & & & & \\
\hline 1.5947 & 2.2223 & 2.5541 & 2.8218 & 3.0676 & 3.3503 & 3.5979 \\
\hline . 12361 & . 00042225 & & & & & \\
\hline 1.7821 & 2.3239 & 2.6058 & 2.8320 & 3.0390 & 3.2763 & 3.4839 \\
\hline
\end{tabular}

Because \(F_{0.05}(2,5)=5.7861,0.09 / F=0.0156\) and \(1 / F=0.173\), so that the model is almost roughly linear.

The input data for the modified Beale's measure program are shown in figure 1.

\section*{Problem 6.2-2}

From the computer output, the modified Beale's measure is 0.29808 . Based on \(F_{0.05}(2,27)\) \(=3.359\), the modified Beale's measure indicates that the model is at the point of being highly
nonlinear. Hence, linear theory can be applied to use the \(W\) statistic based on \(q=2\) only as a very rough approximation.
The input data to the regression code, modified to compute the modified Beale's measure, are shown in figure 1.

CLASS PROBLEM
TWO-DIMENSIONAL FLOW
SEVERAL ZONES


Figure 1
\begin{tabular}{|c|c|c|c|c|c|c|c|}
\hline 16 & 11 & 10 & & 9400 & 8400 & -89.36 & 1 \\
\hline 17 & 7 & 11 & & 6000 & 9400 & 6.68 & 1 \\
\hline 18 & 13 & 11 & & 11400 & 9400 & -1.5.32 & 1 \\
\hline 19 & 3 & 12 & & 2000 & 10400 & 16.88 & 1 \\
\hline 20 & 5 & 12 & & 4000 & 10400 & 1.5 .87 & 1 \\
\hline 21 & 9 & 12 & & 7400 & 10400 & 4.48 & 1 \\
\hline 22 & 11 & 12 & & 9400 & 10400 & -1.8.34 & 1 \\
\hline 23 & 13 & 13 & & 11400 & 11400 & -2.47 & 1 \\
\hline 24 & 14 & 13 & & 13400 & 11400 & 8.10 & 1 \\
\hline 25 & 3 & 14 & & 2000 & 12400 & 54.12 & 1 \\
\hline 26 & 5 & 14 & & 4000 & 12400 & 38.27 & 1 \\
\hline 27 & 10 & 14 & & 8400 & 12400 & . 053 & 1 \\
\hline 28 & 12 & 14 & & 10400 & 12400 & -2.92 & 1 \\
\hline 29 & 7 & 15 & & 6000 & 13400 & 8.30 & 1 \\
\hline 30 & 14 & 15 & & 12400 & 13400 & 4.54 & 1 \\
\hline 31 & 2 & 15 & & 1000 & 14400 & 85.82 & 1 \\
\hline 32 & 11 & 15 & & 9400 & 14400 & 2.26 & 1 \\
\hline 1 & 8 & 8 & 0 & 9 & & & \\
\hline 2 & 10 & 10 & 0 & 11 & & & \\
\hline 3 & 12 & 12 & 0 & 13 & & & \\
\hline 4 & 10 & 10 & 14 & 0 & & & \\
\hline 8 & & 0 & & & & & \\
\hline 9 & & 012 & & & & & \\
\hline 10 & & 84 & & & & & \\
\hline 11 & & 0 & & & & & \\
\hline 12 & & 0 & & & & & \\
\hline 13 & & 051 & & & & & \\
\hline 14 & & 008 & & & & & \\
\hline 1 & & 754 & & 5.754 & 0 & . 00031149 & \\
\hline 2 & & . 89 & & 87.89 & & . 00013995 & \\
\hline 3 & & 288 & & 3.288 & 0 & . 00013516 & \\
\hline 4 & & . 89 & & 87.89 & . 080716 & 0 & \\
\hline 1 & 16 & 7 & 16 & 1 & 80.978 & 0 & . 01 \\
\hline 7 & 16 & 15 & 16 & 2 & 935.15 & 0 & . 00056 \\
\hline 11 & 10 & 11 & 10 & 3 & -97000 & 1940 & 1 \\
\hline 12 & 8 & 12 & 8 & 4 & -50961 & 1020 & 1 \\
\hline IN & 1 & 0 & & & & & \\
\hline 15 & 15 & 5 & 16 & -1 & 0 & & \\
\hline 1 & 10 & 5 & 6 & & 1.04 & . 48 & \\
\hline 15 & 16 & & 198 & & & & \\
\hline 15 & 15 & & & & & & \\
\hline 15 & 14 & & & & & & \\
\hline 15 & 13 & & & & & & \\
\hline 15 & 12 & & & & & & \\
\hline 15 & 11 & & & & & & \\
\hline 15 & 10 & & & & & & \\
\hline 15 & 9 & & & & & & \\
\hline 15 & 8 & & & & & & \\
\hline 15 & 7 & & 211 & & \multirow[t]{2}{*}{} & & \\
\hline 2 & 2 & 6 & 7 & & & & \\
\hline 15 & 6 & & 211 & & \multirow{3}{*}{. 48} & . 54 & \\
\hline 15 & 5 & & 730 & & & & \\
\hline 4 & & & & & & & \\
\hline
\end{tabular}

Figure 1-Continued
190.259
\(2 \quad 1143.6\)
\(3-97085\)
4 - 50990
\(5 \quad 10.313\)
\(6 \quad 5.1349\)
\(7 \quad 5.4858\)
\(8 \quad 72.988\)
3.4248E-4
490.46
11-1.8916E-4
\(12 \quad 22.688\)
13 2.2823E-4
8.0751E-2 71.697
\(\begin{array}{ll}2 & 726.73 \\ 3 & -96915\end{array}\)
3 -96915
\(4 \quad-50932\)
\(5 \quad 10.083\)
\(6 \quad 5.1073\)
75.4602 58.520
2.8050E-4
\(10 \quad 485.32\)
11-9.0735E-5
\(12 \quad 3.8882\)
13 4.2091E-5
14 8.0681E-2
\(1 \quad 124.97\)
21424.3
\(3-97222\)
\(4 \quad-51035\)
\(5 \quad 10.476\)
\(6 \quad 5.1538\)
\(7 \quad 5.4698\)
\(8 \quad 91.304\)
3.6558E-4
\(10 \quad 494.12\)
11-2.5696E-4
\(12 \quad 15.271\)
\(131.5554 \mathrm{E}-4\)
14 8.0866E-2
\(1 \quad 36.987\)
2446.01
\(3 \quad-96778\)
4 - 50887
\(5 \quad 9.9203\)
\(6 \quad 5.0884\)
\(7 \quad 5.4762\)
\(8 \quad 40.204\)
\(2.5740 \mathrm{E}-4\)
\(10 \quad 481.66\)
11-2.2940E-5
\(12 \quad 11.305\)
\(131.1478 \mathrm{E}-4\)
14 8.0566E-2
2
. 98677

\section*{Problem 6.3-1}

\section*{Compute}
\[
\begin{aligned}
\hat{\gamma}= & \left(\underline{Y}_{p}-\underline{X}_{p} \underline{b}^{*}\right)^{T}\left[s^{2} \underline{X}_{p}\left(\underline{X}_{s}^{T} \underline{V}_{s}^{-1} \underline{X}_{s}\right)^{-1} \underline{X}_{p}^{T}+\underline{U}\right]^{-1} \\
& \cdot\left(\underline{Y}_{p}-\underline{X}_{p} \underline{\hat{b}}^{*}\right) .
\end{aligned}
\]

For data set 1
\[
\underline{\hat{b}}^{*}=\left|\begin{array}{l}
50.018 \\
9.1954 \\
2.5008 \times 10^{-5}
\end{array}\right|
\]

For data set 2
\[
\hat{b}^{*}=\left|\begin{array}{l}
50.055 \\
9.7914 \\
2.2766 \times 10^{-5}
\end{array}\right| .
\]

For data set 1
\[
\begin{aligned}
& \left(\underline{X}_{s}^{T} \underline{V}_{s}^{-1} \underline{X}_{s}\right)^{-1} s^{2}= \\
& \left.\begin{array}{rr}
2.5550 \times 10^{-1} 8.9780 \times 10^{-2} & -1.7349 \times 10^{-6} \\
2.5550 \times 10^{-1} & -1.7349 \times 10^{-6} \\
\text { symmetric } & 2.0715 \times 10^{-11}
\end{array}\right] \\
& s^{2} \underline{X}_{p}\left(\underline{X}_{s}^{T} \underline{V}_{s}^{-1} \underline{X}_{s}\right)^{-1} \underline{X}_{p}^{T}=1 \times 2.5550 \times 10^{-1} \times 1 \\
& = \\
& =0.25550 \\
& \underline{U}=(1.1)^{2}
\end{aligned}=1.218 .
\]
\[
\begin{aligned}
& {\left[s^{2} \underline{X}_{p}\left(\underline{X}_{s}^{T} \underline{V}_{s}^{-1} \underline{X}_{s}\right)^{-1} \underline{X}_{p}^{T}+\underline{U}^{-1}\right.} \\
& =0.68236 \\
& \underline{Y}_{p}-\underline{X}_{p} \underline{\underline{b}^{*}}=11-9.1954=1.8046 \\
& \hat{\gamma}=1.8046 \times 0.68236 \times 1.8046=2.222 \\
& \chi_{0.05}^{2}(1)=3.841 \quad \therefore H_{0} \text { accepted }
\end{aligned}
\]

Prior and pure regression estimates of \(\beta_{2}=h_{b}\) are in agreement.

For data set 2
\[
\left(\underline{X}_{s}^{T} \underline{V}_{s}^{-1} \underline{X}_{s}\right)^{-1} s^{2}=
\]
\(\left[\begin{array}{rrr}5.4121 \times 10^{-1} & 2.6265 \times 10^{-1} & -3.9795 \times 10^{-6} \\ & 5.4121 \times 10^{-1} & -3.9795 \times 10^{-6} \\ \text { symmetric } & & 4.3413 \times 10^{-11}\end{array}\right]\)
\[
\underline{U}=(0.95)^{2}=0.9025
\]
\[
s^{2} \underline{X}_{p}\left(\underline{X}_{s} \underline{V}_{s}^{-1} \underline{X}_{s}\right)^{-1} \underline{X}_{p}^{T}+\underline{U}=0.54121+0.9025
\]
\[
=1.44371
\]
\[
\left[s^{2} \underline{X}_{p}\left(\underline{X}_{s} \underline{V}_{s}^{-1} \underline{X}_{s}\right)^{-1} \underline{X}_{p}^{T}+\underline{U}^{-1}\right.
\]
\[
=0.69266
\]
\[
\underline{Y}_{p}-\underline{X}_{p} \underline{\underline{b}}^{*}=9.5-9.7914=-0.2914
\]
\[
\hat{\gamma}=-0.2914 \times 0.69266 \times(-0.2914)
\]
\[
=0.5882
\]
\(\therefore H_{0}\) accepted.```


[^0]:    ${ }^{1}{ }_{1}$ metric ton $=1$ megagram

[^1]:    ${ }^{1}$ Based on $\log _{10} T$, table 2.2-2

[^2]:    ${ }^{1}$ Throughout this text singly underlined symbols represent vectors and doubly underlined symbols represent matrices.

[^3]:    ${ }^{1}$ Use of the trade names in this report is for identification purposes only and does not constitute endorsement by the U.S. Geological Survey.

[^4]:    ${ }^{1}$ These dimensions are approximate. The exact sizes required are calculated and printed in subroutine ORDER. $N_{m}$ is $N_{x}+3$ or $N_{y}+3$, whichever is smaller.

