# University of Cape Town Department of Mathematical Statistics 

## COLLINEARITY AND CONSEQUENCES FOR ESTIMATION:

a STUDY and SIMOLATION

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A thesis prepared under the supervision of

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in fulfilment of the requirements for the degree of Master of Science in Mathematical Statistics

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## CONTENTS

## ACKNOVLEDGEMENTS

## INTRODUCTION

## 1 THE LINEAR REGRESSION MODEL

1.1 The Model 1-01
1.2 Ordinary Least Squares estimation 1 1-02
1.3 Singular Value Decomposition 1-03
1.4 Distributions 1-06
1.4.1 Univariate normal $1-06$
1.4.2 Multivariate normal $1-06$
1.4.3 Central $\chi^{2} \quad 1-07$
1.4.4 Central F 1-07
1.4.5 Central t 1-07
1.4.6 Non- central $\chi^{2} \quad 1-08$
1.4.7 Non- central F $1-08$
1.4.8 Non-central t 1-09
1.5 Variance Inflation Factors 1-09
1.6 Variable Diagnostics 1-10
1.6.1 Analysis of variance (ANOVA) $1-10$
1.6.2 Subset selection of regressor variables 1-13
1.6.2.1 All possible regressions 1-13
1.6.2.2 Stepwise regression 1-17
1.7 Case Diagnostics 1-19
1.7.1 Outliers 1-19
1.7.2 Influence $1-20$
1.8 Bias and Jackknifing 1-22
1.8.1 Biased estimation $1-22$
1.8.2 Jackknifing 1-23
1.9 Vector and Matrix Norms, and Decomposition 1-27
1.9.1 Vector norms 1-27
1.9.2 Matrix norms 1-28
1.9.3 Decomposition ..... 1-30
1.9.3.1 SVD ..... 1-30
1.9.3.2 QR decomposition ..... 1-31
1.10 Subset Selection using Singular Value Decomposition ..... 1-33
1.11 Probability limit ..... 1-38
2 COLLINEARITY
2.1 Definitions ..... 2- 02
2.2 Detecting Collinearity ..... 2-11
2.2.1 Nature and sensivity of estimates ..... 2-11
2.2.2 Correlation matrix of scaled regressors ..... 2-12
2.2.3 Determinant of $X^{\prime} X$ ..... 2-12
2.2.4 Departure from orthogonality ..... 2-13
2.2.5 Smallest singular value ..... 2-14
2.2.6 Condition number and condition index ..... 2-15
2.2.7 Regression coefficient variance decomposition ..... 2-17
2.2.8 Mixed condition index ..... 2-19
2.2.9 Variance inflation factors ..... 2-20
2.2.9.1 VIF as variance ratio ..... 2-23
2.2.9.2 Use of VIF's ..... 2-25
2.2.10 Signal-to-noise tests ..... 2-26
2.2.11 Prior information use ..... 2-29
2.2.12 Detection methods and Mason's definition ..... 2-31
2.3 Collinearity in Practice ..... 2-33
2.3.1 Sources and origins ..... 2-33
2.3.2 Effects of collinearity ..... 2-34
2.3.2.1 Geometric interpretation of collinearity ..... 2-35
2.3.2.2 Collinearity of two regressor variables ..... 2-36
2.3.2.3 Inflation of variance ..... 2-38
2.3.2.4 Unexpected coefficient values and signs ..... 2-38
2.3.2.5 Unstable regression coefficients ..... 2-40
2.3.2.6 Linear combinations of regression variables ..... 2-40
2.4 Centering and Standardization of the $X$ matrix ..... 2-41
2.5 Perturbation ..... 2-44
2.5.1 Perturbation bounds ..... 2-44
2.5.2 Computational accuracy ..... 2-45
2.6 Detecting and Handling Collinearity ..... 2-49
2.7 Summary ..... 2-50
3 PRINCIPAL COIPONENTS REGRESSION MODEL
3.1 Introduction ..... 3-01
3.2 The PC estimator ..... 3- 03
3.3 Properties of $\hat{\beta}_{\mathrm{PC}}$ ..... 3- 04
3.4 Eliminating PC's ..... 3- 05
3.4.1 Theoretical or Structural Norms ..... 3-06
3.4.1.1 Classical F test ..... 3-06
3.4.1.2 Strong $\mathbb{H S E}$ criterion ..... 3- 07
3.4.1.3 First weak MSE criterion ..... 3- 08
3.4.1.4 Confidence interval norm ..... 3- 09
3.4.2 Predictive norms ..... 3- 09
3.4.2.1 Weak predictive MSE criterion ..... 3-10
3.4.2.2 Squared bias of prediction ..... 3-10
3.5 Comments and critique ..... 3-11
3.6 Summary ..... 3-13
4 RIDGE REGRESSION
4.1 Introduction ..... 4- 01
4.2 The RR estimator ..... 4-01
4.3 Properties of $\hat{\beta}_{\mathrm{R}}$ ..... 4- 01
4.4 Estimation of the ridge parameter k ..... 4- 05
4.5 General comments ..... 4-06
4.5.1 Bayesian view ..... 4-06
4.5.2 Critique of ridge regression ..... 4- 07
4.6 Generalized Ridge Regression Estimators ..... 4- 08
4.6.1 K non-stochastic ..... 4- 08
4.6.1.1 The GRR estimator ..... 4- 08
4.6.1.2 Properties of $\hat{\delta}_{\mathrm{K}}$ ..... 4- 09
4.6.2 K stochastic ..... 4-10
4.6.2.1 The GRR estimator ..... 4-10
4.6.2.2 Properties of $\left[\hat{\delta}_{\hat{K}}\right]_{\mathrm{i}}$ ..... 4-12
4.6.3 Choice of k in $\mathrm{K}=\mathrm{kI}$ ..... 4-13
4.7 Jackknifed Ridge Estimator ..... 4-15
4.7.1 The JR estimator ..... 4-15
4.7.2 Properties of $\hat{\delta}_{\text {JW }}$ (non-stochastic K) ..... 4-17
4.7.3 Additional comments on $\hat{\delta}_{\text {JW }}$ (non-stochastic K) ..... 4-19
4.7.4 Additional comments on $\hat{\delta}_{\text {JW }}$ (stochastic K) ..... 4-21
4.8 Summary ..... 4-24
5 Fractional principals in estiuation
5.1 Shrunken Estimators ..... 5- 01
5.1.1 The SH estimator ..... 5-01
5.1.2 Properties (d fixed) ..... 5-01
5.1.3 Properties (d stochastic) ..... 5- 04
5.2 Fractional Principal Component (FPC) regression ..... 5- 04
5.2.1 FPC estimators ..... 5-05
5.2.2 Relationship to other estimators ..... 5-05
5.2.3 Properties of FPC estimators ..... 5-06
5.2.4 0ptimal values for fractions ..... 5- 08
5.2.5 0ptimal FPC estimators ..... 5- 09
5.3 Summary ..... 5-11
6 LATENT ROOT REGRESSION ESTIMATION
6.1 The LRR estimator ..... 6- 01
6.2 Properties and Problems ..... 6-07
6.3 Elimination of independent variables ..... 6-08
6.4 Summary ..... 6- 09
7 TOTAL LEAST SQQURES
7.1 Introduction ..... 7-01
7.2 TLS Technique ..... 7-01
7.2.1 Assumptions and notation ..... 7-01
7.2.2 Geometric view of TLS ..... 7-04
7.2.3 Definition ..... 7-05
7.2.4 Generalized definition ..... 7-07
7.2.5 Properties of TLS estimator ..... 7-08
7.3 Use of TLS in the errors-in- variables model ..... 7-10
7.3.1 The Model ..... 7-10
7.3.2 Estimation by Least Squares ..... 7-12
7.3.3 Estimation by Total Least Squares ..... 7-16
7.4 Subset Selection based on TLS for Prediction ..... 7-17
7.5 Summary ..... 7-19
8 COMPARISON OF ESTTMATORS
8.1 Basic Comparisons ..... 8-01
8.2 Judgement of estimators ..... 8-05
8.2.1 Unbiasedness ..... 8-05
8.2.2 MSE criteria ..... 8-05
8.3 Comparison of FPC estimators ..... 8-06
8.4 Summary ..... 8-09
9 INFLDENCE AND COLLINEARITY
9.1 Influential points ..... 9- 01
9.2 Detecting collinearity-influential points ..... 9-02
9.2.1 Leverage ..... 9- 02
9.2.2 Outlier sum of squares ..... 9- 04
9.2.3 Andrews-Pregibon statistic ..... 9- 06
9.2.4 Cook's statistic ..... 9- 08
9.2.5 DFFITS ..... 9-10
9.2.6 Variance inflation factor ..... 9-12
9.2.7 Condition index ..... 9-16
9.2.7.1 Row deletion and condition indices ..... 9-17
9.2.7.2 Collinearity-influential points ..... 9-20
9.3 Influential Variables ..... 9-22
9.4 Further Remarks on Influence ..... 9-26
9.4.1 Weighted Least Squares ..... 9-26
9.4.2 Influence Measures in Ridge Regression ..... 9-26
9.5 Summary ..... 9-27
10 SIY ULATION STUDY
10.1 Introduction ..... 10-01
10.2 Data ..... 10-01
10.3 Estimation Programs ..... 10-03
10.4 Simulation results ..... 10-04
10.5 Discussion of the results ..... 10-11
10.6 Recommendations and further research ..... 10-12
10.7 Summary ..... 10-12
REFERENCES
APPENDIX A Useful formulae and derivations
APPENDIX B Program
APPENDIX C Notation
BIBLIOGRAPHY

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## INTRODUCTION

Collinearity is a phenomenon that occurs in linear statistical models. It cannot be described in simple terms as present or absent but rather in terms of degree and consequences. Consequences include on the one hand the issue of computational accuracy, and on the other hand, the statistically inherent instability of the estimators. The statistical and mathematical theory and background required to understand collinearity is presented in Chapter 1. Chapter 2 discusses collinearity per se and in Chapters 3 to 6 we draw from the literature to present various biased statistical estimation procedures and their properties. These procedures include principal components, ridge, generalized ridge, shrunken, fractional and latent root regression methods. In Chapter 7 some attempt is made to incorporate the errors-in-variables model in the discussion of collinearity, in as much as it also admits a perturbation framework. Chapter 8 consists of a summary of all the proposed estimators. In Chapter 9 the question of influence is discussed. Chapter 10 presents the results of a simulation study.

Some original comments on some estimation techniques are presented in the appropriate chapters.

An alphabetical listing of notation is presented in Appendix $C$ for the convenience of the reader interested only in elements of the study.

An extensive bibliography has been completed from a literature search on the MATHSCI database and from additional sources found by the author.

In the published literature there is no study of the relative efficiencies of estimators as comprehensive as the study presented here. Lee (1986), in an unpublished doctoral thesis, apparently examined more biased estimators. This thesis borrows partly from the work of Chalton (1990) and explores the relative efficiencies of the estimators across different parameter-vector orientations, different common error variance sizes, and different collinearity severities within a simple and convenient collinearity framework. Some interesting and anomalous properties of estimators emerge.

Sources for sections of the simulation and estimation programs are acknowledged in Chapter 10, but the main program and subroutines written for this study are listed in Appendix B. Annotations highlight key changes in routines.

Avenues for further research are sketched in Chapters 4, 9 and 10.

## Chapter 1

## THE LINEAR REGRESSION MODEL

### 1.1 The Model

The linear regression model is given by

$$
\begin{equation*}
Y=X \beta+\epsilon \tag{1.1}
\end{equation*}
$$

where $Y$ is a nx1 observed response vector, $\epsilon$ is a nx1 vector of uncorrelated random error variables with expectation $E(\epsilon)=0$, and variance matrix $\operatorname{Var}(\epsilon)=V(\epsilon)=\sigma^{2} I$, $\beta$ is a px1 vector of regression coefficients that must be estimated, and
$X$ is a nxp matrix of fixed regressors or independent variables, whose rank is p (we will assume that $\mathrm{n}>\mathrm{p}$ ).

We will not always assume that the $X$ matrix has been standardized. If there is a constant present in the regression model we will assume that it is represented in the $X$ matrix as a column of ones. If we want the $X$ matrix to be scaled so that the product matrix $X^{\prime} X$ is in correlation form, that will be stated explicity.

By centering we imply that the mean of each regressor column is subtracted from the relevant column. By standardizing $X$ we mean that $X$ has been scaled so that the length of each column of $X$ is unity (eg. $X_{i}{ }^{\prime} X_{i}=1$, where $X_{i}$ denotes the i-th column of $X$ ). When the columns have been scaled (centered and standardized) the product matrix $X^{\prime} X$ is in correlation form. In correlation form each of the elements of $X^{\prime} X$ will lie between -1 and +1 . For example let $p=2$, then the correlation matrix of $X^{\prime} X$ is

$$
\left[\begin{array}{ll}
1 & \mathrm{r}_{12} \\
\mathrm{r}_{12} & 1
\end{array}\right]
$$

where $\mathrm{r}_{12}$ is the observed coefficient of correlation between the variables
represented by the first two columns of $X$, and

$$
\begin{aligned}
r_{12} & =\frac{\sum\left(X_{i 1}-\bar{x}_{1}\right)\left(X_{i 2}-\bar{x}_{2}\right)}{\left[\Sigma\left(X_{i 1}-\bar{x}_{1}\right)^{2}\right]^{\frac{1}{2}}\left[\Sigma\left(X_{i 2}-\bar{x}_{2}\right)^{2}\right]^{\frac{1}{2}}} \\
& =\frac{\sum X_{i_{1}} X_{i 2}-n \bar{x}_{1} \bar{x}_{2}}{\left[\left(\sum X_{i 1}^{2}-n \bar{X}_{1}^{2}\right)\left(\sum X_{i 2}^{2}-n \bar{X}_{2}^{2}\right)\right]^{\frac{1}{2}}}
\end{aligned}
$$

where $\bar{x}_{j}$ is the mean of the $j$-th column, $X_{i j}$ is the $i$-th element of the j -th column and the summation is from $\mathrm{i}=1(1) \mathrm{n}$.

Consequences of data centering for collinearity diagnosis are presented in Chapter 2.

### 1.2 Ordinary Least Square estimation

If $\hat{\beta}$ is the ordinary least square estimator (0LSE) of $\beta$ in (1.1), minimizing $(Y-X \beta)^{\prime}(Y-X \beta)$ over all $\beta$, then

$$
\begin{equation*}
\hat{\beta}=\left(X^{\prime} X\right)^{-1} X^{\prime} Y \tag{1.2.1}
\end{equation*}
$$

and the minimum sum of squares of residuals is

$$
\begin{equation*}
\operatorname{RSS}=\left(\hat{\epsilon}^{\prime} \hat{\epsilon}\right)=(\mathrm{Y}-\mathrm{X} \hat{\beta})^{\prime}(\mathrm{Y}-\mathrm{X} \hat{\beta})=\operatorname{SSE}(\hat{\beta}) \tag{1.2.2}
\end{equation*}
$$

Properties:

1. $E(\hat{\beta})=\beta \quad$ (unbiased)
2. $V(\hat{\beta})=\sigma^{2}\left(X^{\prime} X\right)^{-1}$
3. $\mathrm{MSE}=\sigma^{2}\left(\mathrm{X}^{\prime} \mathrm{X}\right)^{-1}$
4. $\quad$ TMSE $=\sigma^{2} \operatorname{tr}\left[\left(X^{\prime} X\right)^{-1}\right]$
5. $\hat{\beta}$ is the best linear unbaised estimator (BLUE) of $\beta$
6. Let $L_{1}=$ Euclidean distance from $\hat{\beta}$ to $\beta$ then

$$
\begin{aligned}
\mathrm{L}_{1}^{2} & =(\hat{\beta}-\beta)^{\prime}(\hat{\beta}-\beta) \\
\mathrm{E}\left(\mathrm{~L}_{1}^{2}\right) & =\sigma^{2} \operatorname{tr}\left(\mathrm{X}^{\prime} \mathrm{X}\right)^{-1} \\
\mathrm{E}\left(\hat{\beta}^{\prime} \hat{\beta}\right) & =\beta^{\prime} \beta+\sigma^{2} \operatorname{tr}\left(\mathrm{X}^{\prime} \mathrm{X}\right)^{-1}
\end{aligned}
$$

7. If $\epsilon$ is distributed normally then

$$
V\left(L_{1}^{2}\right)=2 \sigma^{4} \operatorname{tr}\left[\left(X^{\prime} X\right)^{-2}\right] \quad(\text { from } A .1)
$$

8. If the eigenvalues of $X^{\prime} X$ are denoted by

$$
\begin{aligned}
\lambda_{\max } & =\lambda_{1} \geq \lambda_{2} \geq \ldots \ldots \ldots \geq \lambda_{\mathrm{p}}=\lambda_{\mathrm{min}}>0, \text { then } \\
\mathrm{E}\left(\mathrm{~L}_{1}^{2}\right) & =\sigma^{2} \operatorname{tr}\left[\left(\mathrm{X}^{\prime} \mathrm{X}\right)^{-1}\right] \\
& =\sigma^{2} \Sigma 1 / \lambda_{\mathrm{i}} \quad \text { (from A.2) and } \\
V\left(\mathrm{~L}_{1}^{2}\right) & =2 \sigma^{4} \Sigma 1 / \lambda_{\mathrm{i}}^{2}
\end{aligned}
$$

9. If we assume that the random error terms $\epsilon_{\mathrm{i}}$ are Gaussian (normal) in distribution, the maximum likelihood estimators (MLE's) and the BLUE's coincide (Gauss-Markov Theorem, Searle (1971, p87))

### 1.3 Singular Value Decomposition

The singular value decomposition (SVD) of a matrix is discussed in textbooks such as Stewart (1973, p318), Golub and Van Loan (1983, p16) and Lawson and Hanson (1974, Chapter 4). These discussions can be summarized as follows:

Let $X$ be an nxp matrix of rank $k$. Then there is an nxn orthogonal matrix $U$, an pxp orthogonal matrix $V$, and an nxp matrix $\Delta$ such that

$$
\mathrm{U}^{\prime} \mathrm{XV}=\Delta, \quad \mathrm{X}=\mathrm{U} \Delta \mathrm{~V}^{\prime}
$$

$$
\text { where } \begin{aligned}
\Delta: \operatorname{nxp} & =\left[\begin{array}{cc}
\operatorname{Da} & 0 \\
0 & 0
\end{array}\right] \\
\mathrm{D}_{\alpha} & =\operatorname{Diag}\left(\sqrt{\lambda_{1}}, \sqrt{\lambda_{2}} \ldots \sqrt{\lambda_{\mathrm{k}}}\right)
\end{aligned}
$$

and $\quad \sqrt{\lambda}_{1} \geq \sqrt{\lambda_{2}} \geq \ldots \sqrt{\lambda_{k}}>0$
$\sqrt{\lambda_{i}}$ is called the $i$-th singular values of $X$, and $\lambda_{i}$ is the $i$-th characteristic value (eigenvalue) of $X^{\prime} X$. The vector columns of $U$ are left singular vectors of $X$ and the columns of $V$ are right singular vectors of $X$.

Unless the converse is stated explicitly, we will assume that $X$ has full column rank, i.e. the rank of $X$ is $p$. Observe that we may also write the SVD (also called the basic structure of $X$ ) as

$$
\begin{align*}
\mathrm{X} & =\mathrm{U} \Delta V^{\prime} \quad(\mathrm{X}: \mathrm{nxp}, \mathrm{U}: \mathrm{nxp}, \Delta: \mathrm{pxp}, \mathrm{~V}: \mathrm{pxp})  \tag{1.3.1}\\
\Delta & =\mathrm{D}_{a} \\
\mathrm{~V} & =\left[\mathrm{v}_{1} \ldots \mathrm{v}_{\mathrm{p}}\right], \quad \mathrm{v}_{\mathrm{i}}: \mathrm{px} 1 \\
\mathrm{U} & =\left[\mathrm{u}_{1} \ldots \mathrm{u}_{\mathrm{p}}\right], \quad \mathrm{u}_{\mathrm{i}}: \mathrm{nx} 1 \\
\mathrm{~V}^{\prime} \mathrm{V} & =\mathrm{U}^{\prime} \mathrm{U}=\mathrm{I}_{\mathrm{p}}
\end{align*}
$$

Note that $V^{\prime} V=I_{p}$ but that $U^{\prime} \neq I_{n}$. Using the SVD of $X$ the following equations will be useful in the subsequent chapters.:
1.

$$
\begin{aligned}
X & =\left[u_{1} \ldots u_{p}\right] \operatorname{Diag}\left(\sqrt{\lambda_{1}}, \sqrt{\lambda_{2}} \ldots \sqrt{\lambda_{k}}\right)\left[v_{1} \ldots v_{p}\right]^{\prime} \\
& =\sum_{i=1}^{p} \sqrt{\lambda_{i}} u_{i} v_{i}^{\prime}
\end{aligned}
$$

2. 

$$
\begin{align*}
X^{\prime} X & =V \Delta U^{\prime} U \Delta V^{\prime} \\
& =V \Delta^{2} V^{\prime} \\
& =\sum \lambda_{i} v_{i} v_{i}^{\prime} \tag{1.3.3}
\end{align*}
$$

3. 

$$
\begin{align*}
\left(X^{\prime} X\right)^{-1} & =V \Delta^{-2} V^{\prime} \\
& =\sum_{i=1}^{p} v_{i} v_{i}^{\prime} / \lambda_{i} \tag{1.3.4}
\end{align*}
$$

4. 

$$
\begin{align*}
\hat{\beta} & =\left(X^{\prime} X\right)^{-1} X^{\prime} Y \\
& =V \Delta^{-2} V^{\prime} V \Delta U^{\prime} Y \\
& =V \Delta^{-1} U^{\prime} Y  \tag{1.3.5}\\
& =\sum_{i=1}^{p} V_{i} u_{i} Y / \sqrt{\lambda_{i}}
\end{align*}
$$

Let $c_{i}=u_{i}{ }^{\prime} Y \sqrt{\lambda_{i}}$ then

$$
\begin{equation*}
\hat{\beta}=\sum_{i=1}^{p} \mathrm{v}_{\mathrm{i}} \mathrm{c}_{\mathrm{i}} / \lambda_{\mathrm{i}} \tag{1.3.6}
\end{equation*}
$$

5. $V(\hat{\beta})=\sigma^{2} \sum_{i=1}^{p} v_{i} V_{i}{ }^{\prime} / \lambda_{i} \quad$ (from 1.3.4)

The SVD of the augmented matrix [ X Y] is:

$$
\left[\begin{array}{ll}
\mathrm{X} & \mathrm{Y} \tag{1.3.8}
\end{array}\right]=\tilde{\mathrm{U}} \tilde{\Delta} \tilde{\mathrm{~V}}^{\prime}
$$

with

$$
\begin{align*}
\tilde{U} & =\left[\tilde{u}_{1} \ldots \tilde{u}_{p+1}\right], \tilde{u}_{i}: n \times 1 \\
\tilde{V} & =\left[\tilde{v}_{1} \ldots \tilde{\mathrm{v}}_{\mathrm{p}+1}\right], \tilde{v}_{\mathrm{i}}:(\mathrm{p}+1) \times 1 \\
\tilde{\Delta} & =\operatorname{Diag}\left[w_{1}, w_{2}, \ldots w_{\mathrm{p}+1}\right] \\
w_{1} & \geq w_{2} \geq \ldots \ldots \geq w_{\mathrm{p}+1} \\
\tilde{U}^{\prime} \tilde{\mathrm{U}} & =\tilde{V}^{\prime} \tilde{\mathrm{V}}=\mathrm{I}_{\mathrm{p}+1} \tag{1.3.9}
\end{align*}
$$

The following notation will be used in discussing the augmented matrix:

1. Let $\tilde{v}_{i},{ }_{j}$ be the $j$-th component of the i-th right singular vector $\tilde{v}_{i}$.
2. $\tilde{\mathrm{v}}_{\mathrm{i}}^{0}$ is the p -dimensional vector containing the first p components of the i- th right singular vector $\tilde{v}_{i}$ of [X Y] of dimension $p+1$, thus

$$
\begin{equation*}
\tilde{v}_{i}=\left[\tilde{v}_{i}^{0} \tilde{v}_{i}, p+1\right]^{\prime} \tag{1.3.10}
\end{equation*}
$$

### 1.4 Distributions

### 1.4.1 Univariate normal

When the random variable $X$ has a normal (Gaussian) distribution with mean $\mu$ and variance $\sigma^{2}$, we will write $X \sim N\left(\mu, \sigma^{2}\right)$. The density function of $X$, for $-\infty<x<+\infty$, is

$$
\begin{equation*}
\mathrm{f}(\mathrm{x})=\frac{1}{\sigma \sqrt{2 \pi}} \exp \left[-\frac{1}{2}(\mathrm{x}-\mu)^{2} / \sigma^{2}\right] \tag{1.4.1}
\end{equation*}
$$

### 1.4.2 Multivariate normal

When the random variables in $X^{\prime}=\left[\begin{array}{llll}X_{1} & X_{2} & \ldots X_{n}\end{array}\right]$ have a joint multivariate normal distribution with vector of means $\boldsymbol{\mu}$ and positive-definite variance- covariance matrix $V$ we write $X \sim N(\mu, \mathbf{V})$. The density function of $X$ is then

$$
\begin{equation*}
f\left(X_{1}, X_{2}, \ldots, X_{n}\right)=\frac{\exp \left[-\frac{1}{2}(X-\mu)^{\prime} \mathbf{V}^{-1}(X-\mu)\right]}{(2 \pi)^{\frac{n}{2}}|\mathbf{V}|^{\frac{1}{2}}} \tag{1.4.2}
\end{equation*}
$$

When $E\left(X_{i}\right)=\mu$ for all i then $\mu=\mu 1$ and if the $X_{i}$ 's are mutually independent, all with the same variance $\sigma^{2}$, then $\mathbf{V}=\sigma^{2} I$ and we write $X \sim N\left(\mu 1, \sigma^{2} I\right)$.

### 1.4.3 Central $\chi^{2}$

When $X \sim N(0, I)$ then $U=\sum^{n} \quad X_{i}^{2}$ has the central $\chi^{2}$-distribution with $n$ degrees of freedom. The density function is

$$
\begin{equation*}
f(u)=\frac{u^{\frac{n-2}{2}} \exp (-u / 2)}{(2)^{\frac{n}{2}} \Gamma(n / 2)} \text { for } u>0 \tag{1.4.3}
\end{equation*}
$$

### 1.4.4 Central F

Two independent variables each having $\chi^{2}$-distributions form the basis of the F-distribution. Thus if
$\mathrm{U}_{1} \sim \chi_{\mathrm{n}_{1}}^{2}$ and $\mathrm{U}_{2} \sim \chi_{\mathrm{n}_{2}}^{2}$ then $\mathrm{V}=\frac{\mathrm{U}_{1} / \mathrm{n}_{1}}{\mathrm{U}_{2} / \mathrm{n}_{2}} \sim \mathrm{~F}\left(\mathrm{n}_{1}, \mathrm{n}_{2}\right)$, the F -distribution with $n_{1}$ and $n_{2}$ degrees of freedom. The density function is

$$
\begin{equation*}
f(v)=\frac{\Gamma\left(\frac{1}{2}\left(n_{1}+n_{2}\right)\right) n_{1}{ }^{\frac{1}{2} n_{1}} n_{2}{ }^{\frac{1}{2} n_{2}} v^{\frac{1}{2} n_{1}-1}}{\Gamma\left(\frac{1}{2} n_{1}\right) \Gamma\left(\frac{1}{2} n_{2}\right)\left(n_{2}+n_{1} v\right)^{\frac{1}{2}\left(n_{1}+n_{2}\right)}} \quad \text { for } v>0 \tag{1.4.4}
\end{equation*}
$$

### 1.4.5 Central t

The ratio of a standard normal variable to the root of an independent variable that has a $\chi^{2}$-distribution is the basis of Student's t-distribution. Thus when $\mathrm{X} \sim \mathrm{N}(0,1)$ and U is $\mathrm{U} \sim \chi_{\mathrm{I}}^{2}$, independent of X , then

$$
\mathrm{Z}=\mathrm{X} /(\mathrm{U} / \mathrm{n})^{\frac{1}{2}} \sim \mathrm{t}_{\mathrm{n}},
$$

the t -distribution with n degrees of freedom. Its density function is

$$
\begin{equation*}
f(z)=\frac{\Gamma\left(\frac{1}{2}(n+1)\right)}{\sqrt{n \pi} \Gamma\left(\frac{1}{2} n\right)}\left[1+\frac{z^{2}}{n}\right], \text { for }-\infty<z<\infty \tag{1.4.5}
\end{equation*}
$$

### 1.4.6 Non-central $\chi^{2}$

When $X \sim N(\mu, I)$ and $U=\sum_{i=1}^{n} X_{i}^{2}$, the resulting distribution of $u$ is the non- central $\chi^{2}$ with $n$ degrees of freedom and non-centrality parameter $\gamma$,

$$
\begin{equation*}
\gamma=\mu^{\prime} \boldsymbol{\mu} / 2 \tag{1.4.6}
\end{equation*}
$$

Reference to the distribution is by means of the symbol $\chi^{2}(\mathrm{n}, \gamma)$. The density function of the non-central $\chi^{2}$-distribution $\chi^{2}(n, \gamma)$ is

$$
\begin{equation*}
f(u)=\exp (-\gamma) \sum_{k=0}^{\infty} \frac{\gamma^{k} u^{\frac{1}{2}(n+2 k-1)} \exp \left(-\frac{1}{2} u\right)}{k!(2)^{\frac{1}{2} n+k} \quad \Gamma\left(\frac{1}{2} n+k\right)}, \text { for } u>0 \tag{1.4.7}
\end{equation*}
$$

Some texts prefer to regard $\mu^{\prime} \mu$ as the non-centrality parameter with corresponding adjustments to the form of the density function.

### 1.4.7 Non-central F

If $U_{1}$ and $U_{2}$ are independent and

$$
\begin{aligned}
\mathrm{U}_{1} & \sim \chi^{2}\left(\mathrm{n}_{1}, \gamma\right) \text { and } \mathrm{U}_{2} \sim \chi_{\mathrm{n}_{2}}^{2}\left(\text { or } \chi^{2}\left(\mathrm{n}_{2}, 0\right)\right) \\
\text { then } \quad \mathrm{V} & =\frac{\mathrm{U}_{1} / \mathrm{n}_{1}}{\mathrm{U}_{2} / \mathrm{n}_{2}} \text { is distributed as } \mathrm{F}\left(\mathrm{n}_{1}, \mathrm{n}_{2} ; \gamma\right),
\end{aligned}
$$

the non-central $F$-distribution with $n_{1}$ and $n_{2}$ degrees of freedom and
non- centrality parameter $\gamma$. Its density function is

$$
f(v)=\sum_{k=0}^{\infty} \exp (-\gamma) \frac{\gamma^{k} \Gamma\left(\frac{1}{2}\left(n_{1}+n_{2}+2 k\right)\right)}{k!\quad n_{1}{ }^{\frac{1}{2} n_{1}+k} n_{2}^{\frac{1}{2} n_{2}} v^{\frac{1}{2} n_{1}+k-1}},
$$

Here too some texts prefer to regard $\mu^{\prime} \mu$ as the non-centrality parameter.

### 1.4.8 Non-central t

If $X \sim N(\mu, 1)$ and if $U \sim \chi_{n}^{2}$, independently of $X$, then $T=X /(U / n)^{\frac{1}{2}}$ has the non-central t-distribution, $\mathrm{t}(\mathrm{n}, \mu)$, with n degrees of freedom with the noncentrality parameter $\mu$. The density function is

$$
\begin{equation*}
f(t)=\frac{n^{\frac{n}{2}}}{\Gamma\left(\frac{1}{2} n\right)} \frac{e^{-\frac{1}{2} \mu^{2}}}{\left(n+t^{2}\right)^{\frac{1}{2}(n+1)}} \sum_{k=0}^{\infty} \frac{\Gamma\left(\frac{1}{2}(n+k+1)\right) \mu^{k} 2^{\frac{k}{2}} t^{k}}{k!\left(n+t^{2}\right)^{\frac{k}{2}}} \tag{1.4.8}
\end{equation*}
$$

### 1.5 Variance Inflation Factors

The variance inflation factors (VIF's) were first defined by Marquardt (1970) as the diagonal elements in the inverse of the correlation matrix of $X^{\prime} X$. Thus the $i$ - th variance inflation factor ( $\mathrm{VIF}_{\mathrm{i}}$ ) is:

$$
\begin{equation*}
\operatorname{VIF}_{i}=\frac{\left(X^{\prime} X\right)_{i i}^{-1}}{\left(X_{i}^{\prime} X_{i}\right)^{-1}} \tag{1.5.1}
\end{equation*}
$$

where $\left(X^{\prime} X\right)_{i i}^{-1}$ denotes the i-th diagonal elements of $\left(X^{\prime} X\right)^{-1}$ and $X_{i}$ is the i-th column of $X$. Note that the columns of $X$ are not necessarily scaled or centered.

Discussion of VIF's and their use in diagnosing collinearity is presented in Chapter 2.

### 1.6 Variable Diagnostics

### 1.6.1 Analysis of variance (ANOVA)

Suppose $Y$ is modelled as a linear function of the regressors, with an intercept term. Then $\hat{\epsilon}_{\mathrm{i}}$ (the residual term for the i - th observation) is

$$
\hat{\epsilon}_{i}=Y_{i}-\hat{Y}_{i}
$$

and $\hat{\epsilon}$ will be a $n x 1$ column vector of all the $n$ residual values.

The deviation $Y_{i}-\bar{Y}$ (a quantity measuring the variation of the observations $Y_{i}$ and around their mean) can be decomposed as follows

$$
\underbrace{\mathrm{Y}_{\mathrm{i}}-\bar{Y}}_{\mathrm{I}}=\underbrace{\hat{Y}_{i}-\bar{Y}}_{\mathrm{II}}+\underbrace{\mathrm{Y}_{i}-\hat{Y}_{i}}_{\mathrm{III}}
$$

where $I$ is the total deviation, II is the deviation of fitted 0LS regression value around the overall mean and III is the deviation of the observed value from the regression line. The figure below (Neter and Wasserman (1974)) shows this decomposition for one of the observations.


The sums of these squared deviations satisfy the same additive relationship, due to the mixed terms of the expression having zero sum (see, for example, Neter and Wasserman (1974), for a proof).

$$
\begin{aligned}
& \quad \sum_{i=1}^{n}\left(Y_{i}-\bar{Y}\right)^{2}=\sum_{i=1}^{n}\left(\hat{Y}_{i}-\bar{Y}\right)^{2}+\sum_{i=1}^{n}\left(Y_{i}-\hat{Y}_{i}\right)^{2} \\
& \text { or } \quad \operatorname{SST0}=S S R \quad+\quad \text { SSE }
\end{aligned}
$$

where SSTO is the total sum of squares (corrected for the mean) with n-1 degrees of freedom (df), SSR is the regression sum of squares with p-1 degrees of freedom ( $p$ - 1 independent regressor variables) and SSE denotes the error sum of squares with $n$ - $p$ degrees of freedom ( $p$ parameters are fitted).

In matrix notation and for any value of $p$ the sums of squares are

$$
\begin{align*}
\mathrm{SSTO} & =\mathrm{Y}^{\prime} \mathrm{Y}-\mathrm{n} \overline{\mathrm{Y}}^{2}  \tag{1.6.1}\\
\mathrm{SSR} & =\hat{\beta}^{\prime} X^{\prime} \mathrm{Y}-\mathrm{n} \overline{\mathrm{Y}}^{2}  \tag{1.6.2}\\
\mathrm{SSE} & =\mathrm{Y}^{\prime} \mathrm{Y}-\hat{\beta}^{\prime} X^{\prime} \mathrm{Y} \tag{1.6.3}
\end{align*}
$$

A sum of squares divided by its degrees of freedom is called a mean square (HS). The breakdown of the total sum of squares and associated degrees of freedom are displayed in the form of an analysis of variance table (ANOVA table).

ANOVA Table

| Source | SS | $d f$ | HS |
| :--- | :--- | :--- | :--- |
| Regression | $S S R=\hat{\beta}^{\prime} X^{\prime} Y-n \bar{Y}^{2}$ | $\mathrm{p}-1$ | $\mathrm{MSR}=\mathrm{SSR} /(\mathrm{p}-1)$ |
| Error | $\mathrm{SSE}=\mathrm{Y}^{\prime} Y-\hat{\beta}^{\prime} X^{\prime} Y$ | $\mathrm{n}-\mathrm{p}$ | $\mathrm{MSE}=\mathrm{SSE} /(\mathrm{n}-\mathrm{p})$ |
| Total | $\mathrm{SSTO}=\mathrm{Y}^{\prime} Y-\mathrm{n} \bar{Y}^{2}$ | $\mathrm{n}-1$ |  |

Sometimes the random variable SSE will also be specified as $\operatorname{SSE}\left(X_{1}, X_{2}, \ldots, X_{p}\right)$, where the bracket denotes the subset of the independent
regressor variables that are included in the model. Where this notation is not explicitly used, the set of regressor variables in the model will be clear in the context.

The use of the term MSE here (for the scalar random variable: mean square error) is not to be confused with the non-stochastic matrix MSE, a matrix of expectations corresponding to the (matrix) sum of the variance and bias matrices of a multivariate parameter estimator.

The coefficient of multiple determination is denoted by $R^{2}$ and is defined as

$$
\begin{align*}
& \mathrm{R}^{2}=\frac{\mathrm{SSR}}{\mathrm{SST0}}=1-\frac{\mathrm{SSE}}{\mathrm{SST0}}  \tag{1.6.4}\\
& 0 \leq \mathrm{R}^{2} \leq 1
\end{align*}
$$

$\mathrm{R}^{2}$ measures the proportionate reduction in the variation of Y achieved by the introduction of the entire set of $X$ variables considered in the model. Sometimes for clarity $\mathrm{R}^{2}$ is denoted by $\mathrm{R}_{\mathrm{y} \varphi}^{2}$, where $\varphi$ denotes the set of independent variables that are included in the model (i.e. for $X: n \times p$ including a column of ones to fit an intercept, the $\mathrm{R}^{2}$ of the full model is $\left.R_{y \varphi}^{2}=R_{y 12 \ldots p-1}^{2}\right)$.

The coefficient of multiple correlation $R$ is the positive square root of $\mathrm{R}^{2}$

$$
\begin{equation*}
\mathrm{R}=\sqrt{\mathrm{R}^{2}} \tag{1.6.5}
\end{equation*}
$$

In the case of simple regression $(p=2), R$ is the absolute value of coefficient of correlation $\left|r_{i j}\right|$ where $i$ and $j$ denote the dependent response variable $Y$ and a single regressor variable $X$. For $p \geq 2$, the value of $R$ is the (simple) correlation coefficient between the observed and estimated Y-values, and is consequently always positive.

The coefficient of partial determination is defined as

$$
\begin{equation*}
R_{y i, \phi}^{2}=\frac{\operatorname{SSE}\left(\left\{X_{\phi}\right\}\right)-\operatorname{SSE}\left(X_{i},\left\{X_{\phi}\right\}\right)}{\operatorname{SSE}\left(\left\{X_{\phi}\right\}\right)} \tag{1.6.6}
\end{equation*}
$$

where $\phi$ denotes the set of regressor $X$ variables already in the model prior to fitting $X_{i}$ : For example when $p=4$ and we want to find $R_{y t .234}^{2}$ then $\phi=234$ and $\left\{X_{\phi}\right\}=X_{2}, X_{3}, X_{4}$. Thus

$$
\mathrm{R}_{\mathrm{y} 1.234}^{2}=\frac{\operatorname{SSE}\left(\mathrm{X}_{2}, \mathrm{X}_{3}, \mathrm{X}_{4}\right)-\operatorname{SSE}\left(\mathrm{X}_{1}, \mathrm{X}_{2}, \mathrm{X}_{3}, \mathrm{X}_{4}\right)}{\operatorname{SSE}\left(\mathrm{X}_{2}, \mathrm{X}_{3}, \mathrm{X}_{4}\right)}
$$

The coefficient of partial determination measures the marginal contribution of a regressor variable $X_{i}$, given that other specified regressors are already included in the model.

### 1.6.2 Subset selection of regressor variables

Although $p$ regressor variables are available, not all of them may be necessary for an adequate fit of model to the data. After the functional form of each regressor variable is obtained (i.e. $X_{i}^{2}, \log \left(X_{j}\right), X_{i} X_{j}$, and so on), we seek a 'best' subset of regressor variables. This 'best' subset is not necessarily unique but may be one of a unique set of 'best' subsets.

To find a subset there are basically two strategies, all possible regressions and stepwise regression (which we take to include the special cases of forward selection and backward elimination.)

### 1.6.2.1 All possible regressions

In the all possible regressions search procedure, all possible regression equations are computed and selection of a 'best' equation is performed under some criterion ( $\mathrm{R}^{2}$, Adjusted $\mathrm{R}^{2}$, MSE and $\mathrm{C}_{\mathrm{p}}$ ). If there are $(\mathrm{p}-1)=\mathrm{k}$ independent variables and one intercept term there will be $2^{\mathrm{k}}$ possible
equations. For example if $p=3$ (constant, $X_{1}, X_{2}$ ) the following $2^{2}$ models are possible:
$\mathrm{E}(\mathrm{Y})=\beta_{0} ; \mathrm{E}(\mathrm{Y})=\beta_{0}+\beta_{1} \mathrm{X}_{1} ; \mathrm{E}(\mathrm{Y})=\beta_{0}+\beta_{2} \mathrm{X}_{2} ; \mathrm{E}(\mathrm{Y})=\beta_{0}+\beta_{1} \mathrm{X}_{1}+\beta_{2} \mathrm{X}_{2} ;$ where the meaning and the values of the coefficients $\beta_{0}, \beta_{1}, \beta_{2}$ is different in each model.

## (i) $\mathbf{R}^{2}$ Criterion

The coefficient of multiple determination $R^{2}$ defined in (1.6.4), is computed for each $2^{\mathrm{k}}$ equations. $\mathrm{R}^{2}$ will be a maximum when all p regressor variables are included in the equation. We therefore want to find a minimal subset for which $R^{2}$ has stabilized close to its maximum (i.e. when including another variable in the model, the increase in $R^{2}$ is very small).

## (ii) Adjusted $\mathbf{R}^{\mathbf{2}}$

Adding more independent variables to the model can only increase $\mathrm{R}^{2}$ and never reduce it. A modified measure that recognizes the number of independent variables is introduced. The adjusted coefficient of multiple determination, denoted $R_{a}^{2}$, is defined as:

$$
\mathrm{R}_{\mathrm{a}}^{2}=1-\frac{\mathrm{n}-1}{\mathrm{n}-\mathrm{p}} \frac{\mathrm{SSE}}{\operatorname{SST} 0}=1-\frac{\mathrm{n}-1}{\mathrm{n}-\mathrm{p}}\left(1-\mathrm{R}^{2}\right)
$$

One then computes $R_{a}^{2}$ for each equation and seeks a set (or more than one set) of independent variables which maximizes $R_{a}^{2}$.

## (iii) MSE Criterion

One may compute the MSE for each model equation and seek a set (or more than one set) of independent variables which minimizes MSE. The $\mathrm{R}^{2}$ criterion
does not take into account the number of parameters in the model, whereas the MSE criterion does take that number into account (MSE = SSE/(n-p)).

## (iv) $\mathrm{C}_{\mathrm{p}}$ Criterion

The $\mathrm{C}_{\mathrm{p}}$ criterion, proposed by Mallows (1964), is based on the 'total squared error'. Define the quantity $\Gamma_{p}$

$$
\Gamma_{p}=\left[\sum_{i=1}^{n}\left(\nu_{i}-\eta_{i}\right)^{2}+\sum_{i=1}^{n} \operatorname{Var}\left(\hat{Y}_{i}\right)\right] / \sigma^{2}
$$

where $\nu_{\mathrm{i}}=\nu\left(\mathrm{X}_{1 \mathrm{i}}, \mathrm{X}_{2 \mathrm{i}}, \ldots\right)$ is the expected value from true equation for the conditional expectation of ( $\mathrm{Y}_{\mathrm{i}} \mid \mathrm{X}_{1 \mathrm{i}} \ldots \mathrm{X}_{\mathrm{pi}}$ ),

$$
\begin{aligned}
\eta_{\mathrm{i}} & =\beta_{0}+\sum_{\mathrm{j}=1}^{\mathrm{k}} \beta_{\mathrm{j}} \mathrm{X}_{\mathrm{ij}}, \quad \text { is expected value from fitted equation } \\
\nu_{\mathrm{i}}-\eta_{\mathrm{i}} & =\text { bias at the } \mathrm{i} \text { - th data point } \\
\mathrm{p} & =\mathrm{k}+1 \text { when } \beta_{0} \text { is present } \\
& =\mathrm{k} \text { when } \beta_{0} \text { is absent } \\
\sum_{\mathrm{i}=1}^{\mathrm{n}\left(\nu_{\mathrm{i}}-\eta_{\mathrm{i}}\right)^{2}} & =\text { the sum of squared bias }\left(\mathrm{SSB}_{\mathrm{p}}\right)
\end{aligned}
$$

Now the residual sum of squares (denoted by $\operatorname{SSE}_{\mathrm{p}}$ ), from a fitted equation involving the $p$ estimated coefficients, has the expectation:

$$
\begin{equation*}
E\left(\operatorname{SSE}_{\mathrm{p}}\right)=\mathrm{SSB}_{\mathrm{p}}+(\mathrm{n}-\mathrm{p}) \sigma^{2} \tag{1.6.7}
\end{equation*}
$$

Denote the i-th row of $X$ by $x_{i}^{\prime}$, thus

$$
\begin{align*}
X^{\prime} & =\left[\mathrm{x}_{1}, \mathrm{x}_{2}, \ldots, \mathrm{x}_{\mathrm{n}}\right] \text { and }  \tag{1.6.8}\\
\mathrm{X}^{\prime} \mathrm{X} & =\sum_{\mathrm{i}=1}^{\mathrm{n}} \mathrm{x}_{\mathrm{i}} \mathrm{x}_{\mathrm{i}}^{\prime} \tag{1.6.9}
\end{align*}
$$

$$
\begin{align*}
\operatorname{Var}\left(\hat{Y}_{\mathrm{i}}\right) & =\operatorname{variance} \text { of the fitted value } \hat{Y}_{\mathrm{i}} \\
& =\operatorname{Var}\left(\mathrm{x}_{\mathrm{i}}^{\prime} \hat{\beta}\right) \\
& =\sigma^{2} \mathrm{x}_{\mathrm{i}}^{\prime}\left(\mathrm{X}^{\prime} \mathrm{X}\right)^{-1} \mathbf{x}_{\mathrm{i}} \tag{1.6.10}
\end{align*}
$$

$$
\begin{align*}
\sum_{\mathrm{i}=1}^{\mathrm{n}} \operatorname{Var}\left(\hat{\mathrm{Y}}_{\mathrm{i}}\right) & =\sigma^{2} \sum_{\mathrm{i}=1}^{\mathrm{n}} \mathrm{x}_{\mathrm{i}}^{\prime}\left(\mathrm{X}^{\prime} \mathrm{X}\right)^{-1} \mathrm{x}_{\mathrm{i}} \\
& =\sigma^{2} \operatorname{tr}\left\{\sum_{\mathrm{i}=1}^{\mathrm{n}} \mathrm{x}_{\mathrm{i}}^{\prime}\left(\mathrm{X}^{\prime} X\right)^{-1} \mathrm{x}_{\mathrm{i}}\right\} \\
& =\sigma^{2} \operatorname{tr}\left\{\mathrm{I}_{\mathrm{p}}\right\} \\
& =\sigma^{2} \mathrm{p} \tag{1.6.11}
\end{align*}
$$

thus $\quad \Gamma_{p}=E\left(\operatorname{SSE}_{\mathrm{p}}\right) / \sigma^{2}-(\mathrm{n}-\mathrm{p})+\mathrm{p}$.

If $\sigma^{2}$ is estimated by $\hat{\sigma}^{2}$ (after $p$ parameters are fitted), an estimator of $\Gamma_{p}$, denoted by $C_{p}$, is:

$$
\begin{equation*}
C_{p}=\operatorname{SSE}_{\mathrm{p}} /\left(\hat{\sigma}^{2}\right)-(\mathrm{n}-2 \mathrm{p}) \tag{1.6.12}
\end{equation*}
$$

When there is no bias in the regression equation

$$
\begin{aligned}
\mathrm{E}\left[\mathrm{C}_{\mathrm{p}} / \nu_{\mathrm{i}} \equiv \eta_{\mathrm{i}}\right] & =(\mathrm{n}-\mathrm{p}) \sigma^{2} / \sigma^{2}-(\mathrm{n}-2 \mathrm{p}) \\
& =\mathrm{p}
\end{aligned}
$$

Thus, when the $C_{p}$ values for all possible regressions are plotted against $p$, those regressions with little bias will tend to cluster near the line $C_{p}=p$, while those for equations with substantial bias will fall above this line. With this criterion we identify the sets of independent variables that lead to smallest $C_{p}$ and we would prefer those sets that have little bias (i.e. those near the line).

All possible regressions require much computation (in some contexts this is a disadvantage, but for p large a t-directed search can be performed). For a comprehensive discussion on all possible regression see Daniel and Vood (1980).

### 1.6.2.2 Stepwise regression

Some practitioners prefer stepwise regression because this technique requires less computation than all-possible subsets regression. This search method computes a sequence of regression equations. At each step an independent variable is added or deleted. The common criterion for adding (or deleting) some regressor variable examines the effect of that variable which produces the greatest reduction (or smallest increase) in the error sums of squares, at each step. Under stepwise regression we can distinguish basically three procedures (i) forward selection, (ii) backward elimination procedure and (iii) forward selection with a view back.

## (i) Forward Selection Procedure

In the forward selection procedure, the emphasis is on finding the best single predictor, then the best two predictors (which include the best single predictor) then the best three predictors (which include the best two predictors, and in turn the best single predictor), and so forth. The procedure as outlined by Graybill (1976) is as follows:

1. Compute all correlation coefficients (or $\mathrm{R}_{\mathrm{yi}}^{2}$ for $\mathrm{i}=1,2, \ldots$ p) between Y and $X_{1}, X_{2}, \ldots, X_{p}$, that is compute $\mathrm{r}_{\mathrm{y} 1}^{2}, \mathrm{r}_{\mathrm{y} 2}^{2}, \ldots, \mathrm{r}_{\mathrm{yp}}^{2}$. Choose the largest, suppose it is $r_{y_{1}}^{2}$; then $X_{1}$ is the best single predictor of $Y$.
2. Compute all squared multiple correlation coefficients of $Y$ with all pairs of independent variables involving $X_{1}$, that is compute $R_{y 12}^{2}, R_{y 13}^{2}$, $R_{y 14}^{2}, \ldots, R_{y 1 p}^{2}$, and select the largest. Suppose it is $R_{y 12}^{2}$ then $X_{1}$ and $X_{2}$ are the best two predictors of $Y$ which include the best single predictor $X_{1}$.
3. Compute all multiple correlation coefficients (or their squares) of $Y$ with all sets of three variables that include ( $X_{1}$ and $X_{2}$, that is compute $\mathrm{R}_{\mathrm{y} 123}^{2}, \mathrm{R}_{\mathrm{y} 124}^{2}, \ldots, \mathrm{R}_{\mathrm{y} 12 \mathrm{p}}^{2}$ and select the largest.

At every step in the forward selection procedure we want to determine if the addition of one more variable, will 'appreciably' improve the estimation of Y. If we found that a new variable will improve the resulting estimator of $Y$ we include it and continue, but otherwise the forward selection procedure is terminated because a 'best' subset has been found. Estimation is improved if the corresponding estimate of error variance is sufficiently less than the current estimate.

Another way to formulate this strategy is as follows: We ask if regressors $1,2, \ldots, q, q+1$ yield a better estimate of $Y$ than do regressors $1,2, \ldots, q$ (where $1,2, \ldots, q$ have been determined as above). Thus we examine

$$
\begin{equation*}
\mathrm{H}_{0}: \rho_{\mathrm{y} 12 \ldots \mathrm{q}+1}^{2}=\rho_{\mathrm{y} 12 \ldots \mathrm{q}}^{2} \tag{1.6.13}
\end{equation*}
$$

which is true if and only if

$$
\begin{equation*}
\mathrm{H}_{0}: \rho_{\mathrm{yq}+1}^{2} \cdot 12 \ldots \mathrm{q}=0 \tag{1.6.14}
\end{equation*}
$$

Compute the test statistic $W$, where

$$
W=\frac{(n-q-2) R_{y q+1}^{2} \cdot 12 \ldots q}{1-R_{y q+1}^{2} \cdot 12 \ldots q}
$$

where $\mathrm{R}^{2}$ is the estimated sample estimate of the population value $\rho^{2}$.

If $H_{o}$ is true, $W$ is distributed as $F(1, n-q-2)$. The hypotheses $H_{0}$ is rejected (for a size a test) if and only if $w$ the computed value of $W$ satisfies $w \geq F(a: 1, n-q-2)$, the critical value of the $F(1, n-q-2)$ distribution. So the forward selection procedure is terminated at the step where $H_{0}$ is rejected. In some contexts $a$ is chosen to be quite large, or
almost equivalently, the tabulated $F$-value criterion is replaced by a suitable constant (eg. $\mathrm{F}_{\mathrm{in}}=2.00$ by default in BMDP)

## (ii) Backward Elimination Procedure.

This search procedure is the opposite of forward selection. One starts with the full model and then the less important regressors are eliminated one at a time. The basic steps in the procedure are given in Draper and Smith (1981):

1. A regression equation containing all variables is computed.
2. The partial $F$-test value is calculated for every variable treated as though it were the last variable to enter the regression equation.
3. The lowest partial $F$-test value, $F_{L}$, is compared with a preselected significance level $F_{0}$. Then if $F_{L}<F_{0}$ we remove the variable $X_{L}$ from the equation, and recompute the regression equation without $X_{L}$, then re-enter step 2 again. If $F_{L}>F_{0}$, we adopt the regression equation.

## (iii) Forward Selection with a View Back

This method works just like the forward selection with the difference that at each step one looks back at the independent variables already in the model, examines them and decides if one of them should be dropped.

### 1.7 Case Diagnostics

Outliers are observed values that do not fit the model. Influential cases are observations which can markedly effect the estimation process. Their influence arises from their relationships with the other observations. It is possible for a particular case to be an outlier and to be influential.

### 1.7.1 Outliers

Generally speaking since the true errors are not observable the analyst has
to rely on the estimated error terms. In some situations however the estimated error terms are substantially effected by the influential cases. It is therefore advisable to examine the estimated error terms along with corresponding measures of influence.

### 1.7.2 Influence

For 0LS, the vector of ordinary residuals, $\hat{\epsilon}$ is given by

$$
\begin{align*}
\hat{\epsilon} & =Y-X \hat{\beta} \\
& =Y-X\left(X^{\prime} X\right)^{-1} X^{\prime} Y \\
& =[I-H] Y \tag{1.7.1}
\end{align*}
$$

where $H=X\left(X^{\prime} X\right)^{-1} X^{\prime}$ and $\hat{Y}$ is the vector of fitted values.

The matrix $H$ is called the Hat matrix, because it maps $Y$ into $\hat{Y}=H Y$ (Hoaglin and Velsh (1978)). The matrix $H$ is symmetric ( $H^{\prime}=H$ ), idempotent $(H H=H)$, and a projection matrix (into the column space of $X$ ). The diagonal elements of the Hat matrix, whose role as a diagnostic measure will be discussed in Chapter 9 are:

$$
\begin{equation*}
h_{i}=h_{i i}=X_{i}^{\prime}\left(X^{\prime} X\right)^{-1} X_{i} \tag{1.7.2}
\end{equation*}
$$

where $X_{i}$ ' is the i-th row of $X$. The diagonal elements are known as the leverage values.

Several transformations of the ordinary residuals have been proposed for use in diagnostic procedures (e.g. see Cook and Weisberg (1982, p17)). The most important are the standardized residuals and the studentized residuals. The standardized residual (also called the studentized residual (Cook and Weisberg (1982)) is defined as

$$
\begin{equation*}
r_{i}=\frac{\hat{\epsilon}_{\mathrm{i}}}{\hat{\sigma} \sqrt{1-h_{\mathrm{i} i}}}, \quad i=1,2, \ldots, n \tag{1.7.3}
\end{equation*}
$$

where $\hat{\sigma}$ is the residual mean square. It does not strictly follow a t-distribution because $\hat{\epsilon}_{\mathrm{i}}$ and $\hat{\sigma}$ are not independent. When $\sigma$ is estimated by $\hat{\sigma}(i)$, the estimated error variance when the i-th row of $X$ and $Y$ have been deleted, the result is a studentized residual

$$
\begin{equation*}
\mathrm{t}_{\mathrm{i}}=\frac{\hat{\epsilon}_{\mathrm{i}}}{\hat{\sigma}(\mathrm{i}) \sqrt{1-h_{\mathrm{i}}}}, \quad \mathrm{i}=1,2, \ldots, \mathrm{n} \tag{1.7.4}
\end{equation*}
$$

which is distributed as Student's t with n - $\mathrm{p}-1$ degrees of freedom. A simple formula for $\hat{\sigma}(i)$ (Belsley et al. (1980, p14)) uses

$$
\begin{equation*}
(\mathrm{n}-\mathrm{p}-1)\{\hat{\sigma}(\mathrm{i})\}^{2}=(\mathrm{n}-\mathrm{p}) \hat{\sigma}^{2}-\frac{\hat{\epsilon}_{\mathrm{i}}^{2}}{1-\mathrm{h}_{\mathrm{i} i}} \tag{1.7.5}
\end{equation*}
$$

The measure DFFITS (Belsley et al. (1980, p15)), is the standardized change in the fitted value of a case when it is deleted, is given for the i-th case by

$$
\begin{equation*}
\operatorname{DFFITS}_{i}=\left[\frac{h_{i \mathrm{i}}}{1-\mathrm{h}_{\mathrm{i} i}}\right]^{\frac{1}{2}} \frac{\hat{\epsilon}_{\mathrm{i}}}{\hat{\sigma}(\mathrm{i}) \sqrt{1-h_{\mathrm{i} i}}} \tag{1.7.6}
\end{equation*}
$$

Cook's (squared) distance (Cook (1977)) of an estimator $\tilde{\beta}$ from the 0LS $\hat{\beta}$ is defined as

$$
\begin{equation*}
\mathrm{C}=\mathrm{D}^{2}=(\hat{\beta}-\tilde{\beta})^{\prime} X^{\prime} X(\hat{\beta}-\tilde{\beta}) /\left(\mathrm{p} \hat{\sigma}^{2}\right) \tag{1.7.7}
\end{equation*}
$$

The distance is regarded as large when $D^{2}>F(1-a, p, n-p)$, where $\mathrm{F}(1-\alpha, \mathrm{p}, \mathrm{n}-\mathrm{p})$ is the $1-\alpha$ probability point of the central F - distribution with $p$ and $n-p$ degrees of freedom. It is widely known that $D^{2}$ does not follow an

F-distribution, but is an effective measure of relative change in estimated coefficients.

### 1.8 Bias and Jackknifing

### 1.8.1 Biased estimation

Least square estimators (LSE's or 0LSE's) are the best linear unbiased estimators (BLUE's) of the elements of the parameter vector $\beta$. Amongst linear unbiased estimators the LSE's have the smallest variances. In the presence of collinearity one or more of these variances can be inflated to such an extent that the corresponding estimators become unacceptable. The 'fly in the ointment' with the least squares criterion is its requirement of unbiasedness (Marquardt and Snee (1975)). A major reduction in variance can be obtained as a result of allowing a little bias. If one looks beyond the class of unbiased estimators, it is possible to find some biased estimators with smaller variances than the variances of the LSE's. Some of these biased estimators will perform better than LSE's in the presence of collinearity, in the sense of reduced mean square error (MSE).

Variance and bias in an estimator $\hat{\beta}_{\mathrm{i}}$

LSE


BIASED
ESTIMATOR

MSE may be used to assess the performance of regression estimators. In the regression model (1.1)

$$
\mathrm{Y}=\mathrm{X} \beta+\epsilon,
$$

if $\tilde{\beta}$ is an estimator of $\beta$, the MSE of $\tilde{\beta}$ is defined as

$$
\begin{align*}
\operatorname{MSE}(\tilde{\beta}) & =\mathrm{E}\left[(\tilde{\beta}-\beta)(\tilde{\beta}-\beta)^{\prime}\right] \\
& =\mathrm{V}(\tilde{\beta})+\mathrm{bb}^{\prime} \tag{1.8.1}
\end{align*}
$$

where $\mathrm{b}=\mathrm{E}(\tilde{\beta})-\beta$ is the bias vector.

The total mean squared error (TMSE) of $\tilde{\beta}$ is defined as

$$
\begin{align*}
\operatorname{TMSE}(\tilde{\beta}) & =\operatorname{tr}[\operatorname{MSE}(\tilde{\beta})] \\
& =\Sigma \operatorname{var}\left(\tilde{\beta}_{\mathrm{i}}\right)+\Sigma \mathrm{b}_{\mathrm{i}}^{2} \tag{1.8.2}
\end{align*}
$$

### 1.8.2 Jackknifing

The jackknife technique was introduced by quenouille (1956) and Tukey (1958). The jackknife is a general method for reducing the bias in an estimator and for obtaining a measure of the variance of the resulting estimator by sample reuse.

Let $X=\left[x_{1} \ldots x_{n}\right]^{\prime}$. The subscript -i with any matrix will mean that the i-th row has been deleted, i.e. with $X_{-i}$ we mean the $X$ matrix with its i-th row deleted. In a vector $Y$ the subscript $i$ indicates the $i$-th element of the vector (i.e. $Y_{i}$ ) but the subscript -i, indicates the subvector of $Y$ remaining after the i-th element has been deleted.

Define

$$
\begin{align*}
\hat{\epsilon}_{\mathrm{i}} & =Y_{i}-x_{i}^{\prime} \hat{\beta}  \tag{1.8.3}\\
\hat{\epsilon} & =Y-X \hat{\beta} \\
& =\left[\hat{\epsilon}_{1}, \ldots, \hat{\epsilon}_{\mathrm{n}}\right]^{\prime}  \tag{1.8.4}\\
\mathrm{h}_{\mathrm{i}} & =\mathrm{x}_{\mathrm{i}}^{\prime}\left(X^{\prime} X\right)^{-1} \mathrm{x}_{\mathrm{i}} \tag{1.8.5}
\end{align*}
$$

The least square estimator obtained by deleting the i - th row ( $\mathrm{x}_{\mathrm{i}}^{\prime}, \mathrm{Y}_{\mathrm{i}}$ ) of the data is:

$$
\begin{align*}
\hat{\beta}_{-i} & =\left(X_{-i}^{\prime} X_{-i}\right)^{-1} X_{-i}^{\prime} Y_{-i}  \tag{1.8.6}\\
& =\left[X^{\prime} X-x_{i} x_{i}^{\prime}\right]^{-1}\left[X^{\prime} Y-x_{i} Y_{i}\right] \\
& =\left[\left(X^{\prime} X\right)^{-1}+\left(X^{\prime} X\right)^{-1} x_{i}\left(I-x_{i}^{\prime}\left(X^{\prime} X\right)^{-1} x_{i}\right)^{-1} x_{i}^{\prime}\left(X^{\prime} X\right)^{-1}\right]\left[X^{\prime} Y-x_{i} Y_{i}\right] \\
& =\hat{\beta}-\left(X^{\prime} X\right)^{-1} x_{i}\left[Y_{i}-Y_{i} h_{i}-x_{i}^{\prime} \hat{\beta}+h_{i} Y_{i}\right]\left(1-h_{i}\right)^{-1} \\
& =\hat{\beta}-\left(X^{\prime} X\right)^{-1} x_{i}\left[Y_{i}-x_{i}^{\prime} \hat{\beta}\right]\left(1-h_{i}\right)^{-1} \\
& =\hat{\beta}-\left(X^{\prime} X\right)^{-1} x_{i}\left[\hat{\epsilon}_{i}\right]\left(1-h_{i}\right)^{-1} \tag{1.8.7}
\end{align*}
$$

for $i=1,2, \ldots, n$.

This equation illustrates the effect of an influential point ( $h_{i}$ close to 1) on the coefficients. Under the assumption that $Y_{i}$ can be modelled simultaneously with $Y_{-i}$, the scalar $\hat{\epsilon}_{\mathrm{i}} /\left(1-\mathrm{h}_{\mathrm{i}}\right)$ has zero expectation but large variance $\sigma^{2} /\left(1-h_{i}\right)$, and an outlying $\mathrm{x}_{\mathrm{i}}$ in the row space of $\mathrm{X}_{-\mathrm{i}}$ will tend to have a large influence on the choice of estimates. However if in fact the extrapolation from $Y_{-i}=X_{-i} \beta+\epsilon_{-i}$ to suggest values for $X_{i}{ }^{\prime} \beta$ is not justified, using the full data set will lead to $\hat{\beta}$ values that are generally sufficiently different from $\hat{\beta}_{-i}$ as to be misleading, and in particular, biased (for $\beta$ in the model for the reduced data set).

In what follows, we assume that the same $\beta$ is operative in the full and reduced data sets. Define the pseudovalue for $\mathrm{i}=1,2, \ldots, \mathrm{n}$ as

$$
\begin{align*}
\mathrm{P}_{\mathrm{i}} & =\mathrm{n} \hat{\beta}-(\mathrm{n}-1) \hat{\beta}_{-\mathrm{i}}  \tag{1.8.8}\\
& =\hat{\beta}+(\mathrm{n}-1)\left(\mathrm{X}^{\prime} \mathrm{X}\right)^{-1} \mathrm{x}_{\mathrm{i}}\left[\hat{\epsilon}_{\mathrm{i}}\right] /\left(1-\mathrm{h}_{\mathrm{i}}\right)
\end{align*}
$$

then the jackknifed estimator is given by

$$
\begin{align*}
\hat{\beta}_{\mathrm{J}} & =\mathrm{n}^{-1} \sum \mathrm{P}_{\mathrm{i}} \\
& =\hat{\beta}+(\mathrm{n}-1) \mathrm{n}^{-1}\left(X^{\prime} X\right)^{-1} \sum \mathrm{x}_{\mathrm{i}} \hat{\epsilon}_{\mathrm{i}} /\left(1-\mathrm{h}_{\mathrm{i}}\right)  \tag{1.8.9}\\
E\left(\hat{\beta}_{\mathrm{J}}\right) & =\beta+(\mathrm{n}-1) \mathrm{n}^{-1}\left(X^{\prime} X\right)^{-1} \sum \mathrm{x}_{\mathrm{i}} E\left(\hat{\epsilon}_{\mathrm{i}}\right) /\left(1-h_{\mathrm{i}}\right) \\
& =\beta \quad\left(E\left(\hat{\epsilon}_{\mathrm{i}}\right)=0\right) \tag{1.8.10}
\end{align*}
$$

Since $\operatorname{Var}(\hat{\beta})=\sigma^{2}\left(X^{\prime} X\right)^{-1}$ and $\operatorname{Var}(\hat{\epsilon})=\sigma^{2}\left[I-X\left(X^{\prime} X\right)^{-1} X^{\prime}\right]$, and because $\hat{\beta}$ and $\hat{\epsilon}$ are uncorrelated, we have

$$
\begin{align*}
\operatorname{Var}\left(\hat{\beta}_{\mathrm{J}}\right)= & \sigma^{2}\left(\mathrm{X}^{\prime} \mathrm{X}\right)^{-1}+\sigma^{2}\left((\mathrm{n}-1) \mathrm{n}^{-1}\right)^{2}\left(\mathrm{X}^{\prime} \mathrm{X}\right)^{-1} \\
& \times\left[\sum \mathrm{x}_{\mathrm{i}}\left[\mathrm{I}-\mathrm{x}_{\mathrm{i}}^{\prime}\left(\mathrm{X}^{\prime} \mathrm{X}\right)^{-1} \mathrm{x}_{\mathrm{i}}\right] \mathrm{x}_{\mathrm{i}}^{\prime} /\left(1-\mathrm{h}_{\mathrm{i}}\right)^{2}\right]\left(\mathrm{X}^{\prime} \mathrm{X}\right)^{-1} \tag{1.8.11}
\end{align*}
$$

The jackknife distribution-free estimate of variance for the parameter estimator ( $\hat{\beta}_{\mathrm{J}}$ ) is defined as

$$
\begin{align*}
\mathrm{V}_{\mathrm{J}} & =[\mathrm{n}(\mathrm{n}-1)]^{-1} \Sigma\left(\mathrm{P}_{\mathrm{i}}-\hat{\beta}_{\mathrm{J}}\right)\left(\mathrm{P}_{\mathrm{i}}-\hat{\beta}_{\mathrm{J}}\right)^{\prime}  \tag{1.8.12}\\
& =\operatorname{Var}\left(\hat{\beta}_{\mathrm{J}}\right)
\end{align*}
$$

This sample moment estimator (on the $n$ pseudo-values) may be used to estimate both $\operatorname{Var}\left(\hat{\beta}_{\mathrm{J}}\right)$ and $\operatorname{Var}(\hat{\beta})$ (Hinkley (1977)). He also pointed out the following shortcomings of the method, namely

1. $\hat{\beta}_{\mathrm{J}}$ is different from the original estimator $(\hat{\beta})$, is unbiased for $\beta$ but has in general a larger variance than the LSE (Gauss Markov).
2. $V_{J}$ is in general, biased for estimating $\operatorname{Var}\left(\hat{\beta}_{J}\right)$ or $\operatorname{Var}(\hat{\beta})$.

These problems stem from the balanced nature of the ordinary jackknife, which neglects the unbalanced nature of the regression data. Hinkley (1977) proposed a weighted modification. The weighted pseudo-value

$$
\begin{align*}
\mathrm{Q}_{\mathrm{i}} & =\hat{\beta}+\mathrm{n}\left(1-\mathrm{h}_{\mathrm{i}}\right)\left(\hat{\beta}-\hat{\beta}_{-\mathrm{i}}\right) \\
& =\hat{\beta}+\mathrm{n}\left(1-\mathrm{h}_{\mathrm{i}}\right)\left(\hat{\beta}-\hat{\beta}-\left(\mathrm{X}^{\prime} \mathrm{X}\right)^{-1} \mathrm{x}_{\mathrm{i}}\left[\hat{\epsilon}_{\mathrm{i}}\right]\left(1-\mathrm{h}_{\mathrm{i}}\right)^{-1}\right) \\
& =\hat{\beta}+\mathrm{n}\left(\left(\mathrm{X}^{\prime} \mathrm{X}\right)^{-1} \mathrm{X}_{\mathrm{i}} \hat{\epsilon}_{\mathrm{i}}\right) \tag{1.8.13}
\end{align*}
$$

The weighted jackknife estimator (denoted by $\hat{\beta}_{\mathrm{JW}}$ ) is

$$
\begin{align*}
\hat{\beta}_{J W} & =n^{-1} \sum_{i=1}^{n} Q_{i}  \tag{1.8.14}\\
& =\hat{\beta}
\end{align*}
$$

and the variance estimator

$$
\begin{align*}
V_{J W} & =[n(n-p)]^{-1} \Sigma\left(Q_{i}-\hat{\beta}_{J W}\right)\left(Q_{i}-\hat{\beta}_{J W}\right)^{\prime} \\
& =[n(n-p)]^{-1} \Sigma\left[\hat{\beta}^{+n}\left(\left(X^{\prime} X\right)^{-1} x_{i} \hat{\epsilon}_{i}\right)-\hat{\beta}\right]\left[\hat{\beta}^{n}+n\left(\left(X^{\prime} X\right)^{-1} x_{i} \hat{\epsilon}_{i}\right)-\hat{\beta}\right]^{\prime} \\
& =n(n-p)^{-1}\left(X^{\prime} X\right)^{-1}\left(\Sigma \hat{\epsilon}_{i}^{2} x_{i} x_{i}^{\prime}\right)\left(X^{\prime} X\right)^{-1} \tag{1.8.15}
\end{align*}
$$

$V_{J W}$ will be biased in unbalanced cases but is robust against error variance heterogeneity. (Lemma 2, Appendix of Hinkley (1977))

The above description of the jackknife only takes into account the deletion of one single row at a time. Therefore it is called the delete-one jackknife method. Wu (1986) proposed a class of weighted modifications allowing for the deletion of an arbitrary number of observations.

### 1.9 Vector and Matrix Norms, and Decompositions

### 1.9.1 Vector norms

A vector norm (or simply a norm) on $\mathrm{R}^{\mathrm{n}}$ is a function $\nu: \mathrm{R}^{\mathrm{n}} \rightarrow \mathrm{R}$ that satisfies the following conditions (Stewart (1973)):

$$
\begin{align*}
& \text { 1. } \mathrm{x} \neq 0 \Rightarrow \nu(\mathrm{x})>0 \\
& \text { 2. } \quad \nu(a \mathrm{x})=|a| \nu(\mathrm{x})  \tag{1.9.1}\\
& 3 . \\
& 3(\mathrm{x}+\mathrm{y}) \leq \nu(\mathrm{x})+\nu(\mathrm{y})
\end{align*}
$$

The conditions $1,2,3$, are also termed definiteness, homogeneity, and triangle inequality conditions.

Three norms on $\mathrm{R}^{\mathrm{n}}$ that are frequently used in analyzing matrix processes, are the $1-, 2^{-}$, and $\infty$ - norms.

The 1-norm of a vector $y$, is defined as

$$
\begin{equation*}
\|y\|_{1}=\sum_{i=1}^{n}\left|y_{i}\right|, \tag{1.9.2}
\end{equation*}
$$

where $y_{i}$ is the $i$-th element of the vector $y: n \times 1$.

The 2-norm of a vector $y$, is defined as

$$
\begin{equation*}
\|y\|_{2}=\sqrt{y^{\prime} y} \tag{1.9.3}
\end{equation*}
$$

The 2-norm is sometimes called the Euclidean norm of a vector y.

The $\infty$-norm of a vector $y$, is defined as

$$
\begin{equation*}
\|y\|_{\infty}=\max \left\{\left|y_{i}\right|: i=1,2, \ldots, n\right\} \tag{1.9.4}
\end{equation*}
$$

and is sometimes called the maximum norm (max-norm) or the Chebyshev norm.

The norms defined in (1.9.2), (1.9.3) and (1.9.4) are special cases of the Hölder norms or vector p-norms defined by

$$
\begin{equation*}
\|y\|_{p}=\sqrt[p]{\sum_{i=1}^{p}\left|y_{i}\right|^{p}}, 1 \leq p<\infty \tag{1.9.5}
\end{equation*}
$$

$\left(\|y\|_{\infty}\right.$ is $\lim \left(\|y\|_{p}\right)$ as $\left.p \rightarrow \infty\right)$

### 1.9.2 Matrix norms

A function $\quad \nu: \mathrm{R}^{\mathrm{mxn}} \rightarrow \mathrm{R}$ is a matrix norm on $\mathrm{R}^{\mathrm{mxn}}$ if

1. $\mathrm{A} \neq 0 \Rightarrow \nu(\mathrm{~A})>0, A \in \mathrm{R}^{\mathrm{mxn}}$,
2. $\nu(a \mathrm{~A})=|a| \nu(\mathrm{A}), \mathrm{A} \in \mathrm{R}^{\mathrm{mxn}}, a \in \mathrm{R}$,
3. $\nu(\mathrm{A}+\mathrm{B}) \leq \nu(\mathrm{A})+\nu(\mathrm{B}), \mathrm{A}, \mathrm{B} \in \mathrm{R}^{\mathrm{mxn}}$
4. $\nu(\mathrm{AB}) \leq \nu(\mathrm{A}) \nu(\mathrm{B})$.

Condition (4) is known as the submultiplicative or consistency condition. If a function statisfies (1)-(3) and not necessarily (4), it is called a generalized matrix norm.

The Frobenius norm of a matrix A is defined as

$$
\begin{equation*}
\|A\|_{F}=\sqrt{\sum_{\mathrm{i}} \mathrm{a}_{\mathrm{i} j \mathrm{j}}^{2}} \tag{1.9.7}
\end{equation*}
$$

The Frobenius norm can also be shown to satisfy

$$
\begin{equation*}
\|A\|_{F}^{2}=\operatorname{tr}\left[A^{\prime} A\right] \tag{1.9.8}
\end{equation*}
$$

This norm (1.9.8) is sometimes called the Euclidean matrix norm, $1_{2}$ norm, the Schur norm, or the Hilbert-Schmidt norm.

A unitarily invariant matrix norm is a norm that satisfies

$$
\begin{equation*}
\left\|U^{\prime} X V\right\|=\|X\| \tag{1.9.10}
\end{equation*}
$$

for all unitary matrices $U$ and $V$. (Although the symbols $U$ and $V$ are used in the SVD to indicate unique matrices, here we wish the identity to hold for all other conformable unitary matrices, as well as those $U$ and $V$ of the SVD).

The matrix $p$ - norm of a matrix is defined from vector $p$ - norms as

$$
\begin{equation*}
\|A\|_{p}=\sup _{x \neq 0} \frac{\|A x\|_{p}}{\|x\|_{p}} \tag{1.9.11}
\end{equation*}
$$

where $p \in(1,2, \infty)$

$$
\begin{equation*}
\text { e.g. } \quad\|A\|_{2}=\sup _{x \neq 0} \frac{\|A x\|_{2}}{\|x\|_{2}} \tag{1.9.12}
\end{equation*}
$$

The maximum column-sum matrix norm $\|\cdot\|_{1}$ of A is defined as

$$
\begin{equation*}
\|A\|_{1}=\max _{i=1}^{n}\left|a_{i j}\right|, 1 \leq j \leq n \tag{1.9.13}
\end{equation*}
$$

The maximum row-sum matrix norm $\|\cdot\|_{\infty}$ of A is defined as

$$
\begin{equation*}
\|A\|_{\infty}=\max \sum_{j=1}^{n}\left|a_{i j}\right|, 1 \leq i \leq n \tag{1.9.14}
\end{equation*}
$$

The spectral norm $\|\cdot\|_{2}$ of $A$ is defined as

$$
\begin{equation*}
\|A\|_{2}=\max \left\{\sqrt{\lambda}: \lambda \text { is an eigenvalue of } A^{\prime} A\right\} \tag{1.9.15}
\end{equation*}
$$

### 1.9.3 Decomposition

### 1.9.3.1 SVD

If the SVD of X is given by (1.3.1) and

$$
\sqrt{\lambda_{1}} \geq \sqrt{\lambda_{2}} \geq \ldots \geq \sqrt{\lambda_{\mathrm{r}}}>\sqrt{\lambda_{\mathrm{r}+1}}=\ldots \sqrt{\lambda_{\mathrm{p}}}=0
$$

then

$$
\begin{align*}
r(X) & =r \leq p  \tag{1.9.16}\\
N(X) & =\operatorname{span}\left\{v_{r+1}, \ldots, v_{p}\right\}  \tag{1.9.17}\\
R(X) & =\operatorname{span}\left\{u_{1}, \ldots, u_{r}\right\} \tag{1.9.18}
\end{align*}
$$

where $r(X)$ is the rank of $X, N(X)$ is the null space of $X$, and $R(X)$ is the range of $X$.

Then the Frobenius norm of $X$ can be written as

$$
\begin{equation*}
\|\mathrm{X}\|_{\mathrm{F}}^{2}=\lambda_{1}+\lambda_{2}+\cdots+\lambda_{\mathrm{r}} \tag{1.9.19}
\end{equation*}
$$

and the matrix 2-norm of X is

$$
\begin{equation*}
\|X\|_{2}=\sqrt{\lambda_{1}} \tag{1.9.20}
\end{equation*}
$$

Some authors call (1.9.20) the spectral norm and define it as

$$
\|A\|_{2}=\max \|A x\|_{2}, \text { for }\|x\|=1
$$

Stewart (1987), omits the subscript 2. Proofs of these properties can be found Golub and Van Loan (1983, Chapter 2) or Horn and Johnson (1987).

### 1.9.3.2 QR decomposition

The following decomposition of a matrix is known as the $Q \mathrm{R}$ decomposition:

Let $X: n x p$ and $Y: n x 1$ be given and suppose that an orthogonal matrix $\ell: n x n$ exists and is computable, with the property that

$$
Q^{\prime} X=R=\left[\begin{array}{l}
R_{1}: \operatorname{pxp}  \tag{1.9.21}\\
0:(n-p) \times p
\end{array}\right]
$$

is upper triangular.

If

$$
Q^{\prime} Y=\left[\begin{array}{l}
c: p \times 1  \tag{1.9.22}\\
d:(n-p) \times 1
\end{array}\right]
$$

then $\quad\|\mathrm{X} \beta-\mathrm{Y}\|_{2}^{2}=\left\|\mathrm{Q}^{\prime} \mathrm{X} \beta-\mathrm{Q}^{\prime} \mathrm{Y}\right\|_{2}^{2} \quad($ from 1.9.10 $)$

$$
\begin{align*}
& =\left\|\left[\begin{array}{l}
\mathrm{R}_{1} \\
0
\end{array}\right] \beta-\left[\begin{array}{l}
\mathrm{c} \\
\mathrm{~d}
\end{array}\right]\right\|_{2}^{2} \\
& =\left\|\mathrm{R}_{1} \beta-\mathrm{c}\right\|_{2}^{2}+\|\mathrm{d}\|_{2}^{2} \tag{1.9.23}
\end{align*}
$$

for any $\beta \in \mathrm{R}^{\mathrm{p}}$.

If $r(X)=p(i . e . X$ has full rank), then the 0LSE $\hat{\beta}$ may be obtained from the upper triangular system $\mathrm{R}_{1} \hat{\beta}=\mathrm{c}$, and the minimum sum of squares satisfies $\|X \hat{\beta}-Y\|_{2}^{2}=\|\mathrm{d}\|_{2}^{2}$.

If $X$ is rank deficient $(r(X)<p)$ then at least one diagonal entry in $R$ is zero and the $Q R$ factorization does not necessarily produce an orthonormal basis for $R(A)$. Therefore the $Q R$ factorization must be modified to produce an orthonormal basis for the $X$ range. This modified algorithm is known as QR with Column Pivoting:

Let II be a suitable (pxp) permutation matrix used to interchange the columns of $X$ so that the independent columns are moved to initial column positions. Then

$$
\text { XII }=Q R \quad \text { where } \underset{\mathrm{nxp}}{\mathrm{R}}=\left[\begin{array}{cc}
\mathrm{R}_{11} & \mathrm{R}_{12}  \tag{1.9.24}\\
0 & 0 \\
\mathrm{r} & \mathrm{p}-\mathrm{r}
\end{array}\right] \begin{gathered}
\mathrm{n}-\mathrm{r}
\end{gathered}
$$

where $\operatorname{rank}(X)=r \leq p, R_{11}$ is upper triangular and non-singular. Thus

$$
\begin{align*}
\|X \beta-Y\|_{2}^{2} & =\left\|\left(Q^{\prime} X I I\right)\left(\Pi^{\prime} \beta\right)-Q^{\prime} Y\right\|_{2}^{2} \\
& =\left\|\left[\begin{array}{ll}
\mathrm{R}_{11} & \mathrm{R}_{12} \\
0 & 0
\end{array}\right]\left[\begin{array}{l}
\mathrm{Z}_{1} \\
\mathrm{Z}_{2}
\end{array}\right]-\left[\begin{array}{c}
\mathrm{c} \\
\mathrm{~d}
\end{array}\right]\right\|_{2}^{2} \\
& =\left\|\mathrm{R}_{11} \mathrm{Z}_{1}-\left(\mathrm{c}-\mathrm{R}_{12} \mathrm{Z}_{2}\right)\right\|_{2}^{2}+\|\mathrm{d}\|_{2}^{2} \tag{1.9.25}
\end{align*}
$$

where

$$
\Pi^{\prime} \beta=\left[\begin{array}{l}
\mathrm{Z}_{1}  \tag{1.9.26}\\
\mathrm{Z}_{2}
\end{array}\right] \begin{gathered}
\mathrm{r} \\
\mathrm{n}-\mathrm{r}
\end{gathered}
$$

and $Q^{\prime} Y$ is defined in (1.9.22). Thus, if $\|X \beta-Y\|_{2}^{2}$ is minimized then

$$
\Pi^{\prime} \beta=\left[\begin{array}{c}
\mathrm{R}_{11}^{-1}\left(\mathrm{c}-\mathrm{R}_{12} \mathrm{Z}_{2}\right)  \tag{1.9.27}\\
\mathrm{Z}_{2}
\end{array}\right]_{\mathrm{n}-\mathrm{r}}^{\mathrm{r}}
$$

If $Z_{2}$ is set to zero then we obtain the unique solution of smallest norm

$$
\Pi^{\prime} \beta_{\mathrm{basic}}=\left[\begin{array}{c}
\mathrm{R}_{11}^{-1} \mathrm{c}  \tag{1.9.28}\\
0
\end{array}\right] \underset{\mathrm{n}-\mathrm{r}}{\mathrm{r}}
$$

An algorithm for QR with column pivoting can be found on p165 of Golub and Van Loan (1983). Lawson and Hanson (1974) describe the above method on pp78-82 and refered to it as QR with column interchange strategy.

Golub and Van Loan (1983) show by examples that QR with column pivoting is not entirely reliable for detecting rank deficiency but that it works well in practice.

### 1.10 Subset Selection using Singular Value Decomposition.

In OLS $\hat{\beta}$ is the minimum norm solution, i.e. $\|X \beta-Y\|_{2}^{2}$ will be a minimum when $\beta$ is estimated by $\hat{\beta}=\sum_{i=1}^{\mathrm{p}} \mathrm{v}_{\mathrm{i}} \mathrm{u}_{\mathrm{i}}{ }^{\prime} \mathrm{Y} / \sqrt{\lambda_{\mathrm{i}}}$ (from 1.3.4). In the case of rank deficiency $\hat{\beta}$ can be approximated by

$$
\begin{equation*}
\hat{\beta}_{\hat{r}}=\sum_{i=1}^{\hat{r}} v_{i} u_{i}^{\prime} Y / \sqrt{\lambda_{i}} \tag{1.10.1}
\end{equation*}
$$

where $\hat{\mathrm{r}}$ is an estimate of the $\operatorname{rank}(\mathrm{X})=\mathrm{r}$, thus $\hat{\beta}_{\hat{\mathrm{r}}}$ minimizes $\left\|\mathrm{X}_{\hat{\mathrm{r}}} \beta-\mathrm{Y}\right\|_{2}$ and

$$
\begin{align*}
X_{\hat{r}} & =\sum_{i=1}^{\hat{r}} \sqrt{\lambda_{i}} u_{i} v_{i}^{\prime} \quad \text { (by using (1.3.2)) }  \tag{1.10.2}\\
& =U_{\hat{r}} \Delta_{\hat{r}} V_{\hat{r}}{ }^{\prime}
\end{align*}
$$

where $\mathrm{U}_{\hat{\mathrm{r}}}=\left[\mathrm{u}_{1}, \ldots, \mathrm{u}_{\hat{r}}\right], \mathrm{V}_{\hat{\mathrm{r}}}=\left[\mathrm{v}_{1}, \ldots, \mathrm{v}_{\hat{\mathrm{r}}}\right]$ and $\Delta_{\hat{\mathrm{r}}}=\operatorname{diag}\left(\sqrt{\lambda_{1}}, \ldots, \sqrt{\lambda_{\hat{r}}}\right) . \quad\left(\mathrm{X}_{\hat{\mathrm{r}}}\right.$ is the matrix that is the closest to $X$ that has rank $r$ ). Furthermore the residuals due to this method will be denoted by $\epsilon_{\hat{r}}$, and

$$
\begin{align*}
\hat{\epsilon}_{\hat{r}} & =Y-X \hat{\beta}_{\hat{r}} \\
& =Y-U_{\hat{r}} U_{\hat{r}^{\prime}} Y \\
& =\left(I-U_{\hat{r}} U_{\hat{r}^{\prime}}^{\prime}\right) Y \tag{1.10.4}
\end{align*}
$$

A subset selection procedure that is based on the SVD has been proposed by Golub, Klema and Stewart (1976). It is also described in Golub and Van Loan (1983, pp414-419). Their method proceeds as follows:

1. Compute the SVD of $X(1.3 .1)$ and use it to determine the rank of $X$. Denote this rank by $\hat{r}, \hat{r} \leq p \leq n$.
2. Calculate a permutation matrix $P$ such that the columns of the matrix $X_{1}: n x \hat{r}$ in $X P=\left[\begin{array}{ll}X_{1} & X_{2}\end{array}\right]$ are 'sufficiently independent'.
3. Predict $Y$ with $X b$ where $b=P\left[\begin{array}{ll}z^{\prime} & 0\end{array}\right]^{\prime}$ and $z: \hat{r} x 1$ minimizes $\left\|X_{1} z-Y\right\|_{2}^{2}$. Denote the residuals due to this method by

$$
\begin{align*}
\hat{\epsilon}_{Z} & =Y-X b \\
& =Y-X_{1} z \\
& =Y-X_{1}\left(X_{1}^{\prime} X_{1}\right)^{-1} X_{1} Y \\
& =\left(I-B_{1} B_{1}^{\prime}\right) Y \tag{1.10.5}
\end{align*}
$$

where $B_{1}=X_{1}\left(X_{1}^{\prime} X_{1}\right)^{-\frac{1}{2}}$.
The key step is 2 , since

$$
\min \left\|X_{1} Z-Y\right\|_{2}=\|X b-Y\|_{2} \geq \min \|X \beta-Y\|_{2}
$$

We want some bounds on the smallest singular value $\sqrt{\lambda_{\hat{r}}\left(X_{1}\right)}$ of $X_{1}$. For clarity the singular values of any matrix, say $A$, will be written here as $\sqrt{\lambda_{i}(\Lambda)}$. Golub and Van Loan (1983, pp416-418) state and prove the following theorems:

Theorem 1.10.1

Let the SVD of $X$ be given by (1.3.1) and define the matrix $X_{1}: n x \hat{r}$, for $\hat{r} \leq \operatorname{rank}(X)$, by

$$
\mathrm{XP}=\left[\begin{array}{ll}
\mathrm{X}_{1} & \mathrm{X}_{2}
\end{array}\right]
$$

where $P$ is a pxp permutation matrix. If

$$
P^{\prime} V=\left[\begin{array}{cc}
\ddot{\mathrm{V}}_{11} & \ddot{\mathrm{~V}}_{12} \\
\ddot{\mathrm{~V}}_{21} & \ddot{\mathrm{~V}}_{22}
\end{array}\right] \begin{gathered}
\hat{\mathrm{r}} \\
\hat{\mathrm{r}}-\hat{\mathrm{r}} \\
\hat{\mathrm{p}}-\hat{\mathrm{r}}
\end{gathered}
$$

and $\ddot{\mathrm{V}}_{11}$ is nonsingular, then

$$
\begin{equation*}
\frac{\sqrt{\lambda_{\hat{\mathrm{r}}}(\mathrm{X})}}{\left\|\ddot{\mathrm{V}}_{1}^{-1}\right\|_{2}} \leq \sqrt{\lambda_{\hat{\mathrm{r}}}\left(\mathrm{X}_{1}\right)} \leq \sqrt{\lambda_{\hat{\mathrm{r}}}(\mathrm{X})} \tag{1.10.7}
\end{equation*}
$$

Thus the permutation $P$ must be chosen in such away that $\ddot{V}_{11}$ is as well conditioned as possible. A solution to this problem would be QR with column-pivoting factorization of the $\hat{r} x p$ matrix $\left[\begin{array}{lll}V_{11}^{\prime} & V_{21}^{\prime}\end{array}\right]$ where $V$ is partitioned as follows:

$$
\begin{gather*}
V=\left[\begin{array}{cc}
V_{11} & V_{12} \\
V_{21} & V_{22}
\end{array}\right] \begin{array}{c}
\hat{r} \\
p-\hat{r} \\
\hat{r}
\end{array} \mathrm{p}^{\hat{r}} \tag{1.10.8}
\end{gather*}
$$

Thus

$$
Q^{\prime}\left[\begin{array}{ll}
\ddot{\mathrm{V}}_{11}^{\prime} & \ddot{\mathrm{V}}_{21}^{\prime}
\end{array}\right]=Q^{\prime}\left[\begin{array}{ll}
\mathrm{V}_{11}^{\prime} & \mathrm{V}_{21}^{\prime}
\end{array}\right] P=\left[\begin{array}{ll}
\mathbf{R}_{11} & \mathbf{R}_{12} \tag{1.10.9}
\end{array}\right]
$$

where $\mathbb{Q}$ is orthogonal, $P$ is an pxp permutation matrix and $R_{11}$ is upper triangular. Then from (1.10.9)

$$
\left[\begin{array}{ll}
\mathrm{V}_{11}^{\prime} & \mathrm{V}_{21}^{\prime}
\end{array}\right] \mathrm{P}=\mathrm{Q}\left[\begin{array}{ll}
\mathrm{R}_{11} & \mathrm{R}_{12}
\end{array}\right]
$$

Thus

$$
\begin{align*}
P^{\prime}\left[\begin{array}{l}
V_{11} \\
V_{21}
\end{array}\right] & =\left[\begin{array}{ll}
R_{11}^{\prime} \\
R_{21}^{\prime}
\end{array}\right] Q^{\prime} \\
{\left[\begin{array}{ll}
\ddot{V}_{11} \\
\ddot{V}_{21}
\end{array}\right] } & =\left[\begin{array}{ll}
R_{11}^{\prime} & Q^{\prime} \\
R_{21}^{\prime} & Q^{\prime}
\end{array}\right] \quad(\text { from }(1.10 .6)) \tag{1.10.10}
\end{align*}
$$

and

$$
\left\|\ddot{\mathrm{V}}_{11}^{-1}\right\|=\left\|\left(\mathrm{R}_{11}^{\prime} Q^{\prime}\right)^{-1}\right\|=\left\|\mathrm{R}_{11}^{-1}\right\|
$$

Vith column pivoting $R_{11}$ will be well conditioned, which in turn will produce a well conditioned $\ddot{\mathrm{V}}_{11}$.

To assess the above method of subset selection, compare the residuals from this subset procedure $\left(\hat{\epsilon}_{\mathrm{z}}\right.$, defined in (1.10.5)) with those residuals from the nearest rank- $\hat{\mathbf{r}} \operatorname{LS}\left(\hat{\epsilon}_{\hat{\mathrm{r}}}\right.$, defined in (1.10.4)). Define $\ddot{\mathrm{V}}_{11}$ as the leading $\hat{\mathrm{r}}$ - by- $\hat{\mathbf{r}}$ submatrix of $\mathrm{P}^{\prime} \mathrm{V}$, then

$$
\begin{equation*}
\left\|\hat{\epsilon}_{\hat{\mathrm{r}}}-\hat{\epsilon}_{\mathrm{Z}}\right\|_{2} \leq \frac{\left(\lambda_{\hat{\mathrm{r}}+1}(X)\right)^{\frac{1}{2}}}{\left(\lambda_{\hat{\mathrm{r}}}(X)\right)^{\frac{1}{2}}}\left\|\ddot{\mathrm{~V}}_{11}^{-1}\right\|_{2}\|\mathrm{Y}\|_{2} \tag{1.10.11}
\end{equation*}
$$

The norm $\left\|\hat{\epsilon}_{\hat{r}}-\hat{\epsilon}_{\mathrm{z}}\right\|_{2}$ can also be written as

$$
\begin{align*}
\left\|\hat{\epsilon}_{\hat{r}}-\hat{\epsilon}_{Z}\right\|_{2} & =\left\|Y-X \hat{\beta}_{\hat{r}}-\left(Y-X_{1} Z\right)\right\|_{2} \quad \text { (from 1.10.4 and 1.10.5) } \\
& =\left\|X_{1} Z-\sum_{i=1}^{\hat{r}} u_{i}\left(u_{i}^{\prime} Y\right)\right\|_{2} \quad \text { (from (1.10.1)) } \tag{1.10.12}
\end{align*}
$$

This sheds light on how well $X_{1} Z$ can predict the stable component of $Y$, i.e. ${ }^{\mathrm{r}} \Sigma\left(\mathrm{u}_{\mathrm{i}}{ }^{\prime} \mathrm{Y}\right)$. Any attempt to approximate $\stackrel{\mathrm{n}}{\mathrm{n}}_{\mathrm{E}}^{\mathrm{\Sigma}}\left(\mathrm{u}_{\mathrm{i}}{ }^{\prime} \mathrm{Y}\right)$, where $\mathrm{s}=\mathrm{r}+1$, can lead to $\mathrm{i}=1$

$$
\mathrm{i}=\mathrm{s}
$$

a large norm solution. Moreover (1.10.11) says that if $\lambda_{\tilde{r}+1} \ll \lambda_{\tilde{r}}$ then any reasonably independent subset of columns produces essentially the same-sized residual. On the other hand, if there is no well-defined gap in the singular values then the determination of $\hat{r}$ becomes difficult and the entire subset selection problem more complicated.

These observations together with theorem 1.10 .1 form the basis of an algorithm proposed by Golub, Klema and Stewart (1976):

Algorithm SX-0LS (SX-0LS comes from subset selection on $X$ where the method of estimation is 0LS).

Given $X \in R^{\mathrm{nxp}}, \mathrm{Y} \in \mathrm{R}^{\mathrm{n}}$ and a method for computing an integer $\hat{\mathrm{r}}$ that approximates rank X , then the following algorithm computes a permutation P and a vector $z \in \hat{R^{r}}$ such that the first $\hat{r}$ columns of $X P$ are independent and such that $\left\|\mathrm{XP}\left[\begin{array}{ll}\mathrm{Z}^{\prime} & 0\end{array}\right]^{\prime}-\mathrm{Y}\right\|_{2}$ is minimized.

Compute the SVD of $X$ (1.3.1) and use it to determine the rank of $X$. Denote this rank by $\hat{\mathrm{r}}$, where $\hat{\mathrm{r}} \leq \mathrm{r}(\mathrm{X}) \leq \mathrm{p} \leq \mathrm{n}$ and partition V as in (1.10.8).

Step one:

$$
\begin{gathered}
V=\left[\begin{array}{ll}
V_{11} & V_{12} \\
V_{21} & V_{22}
\end{array}\right] \hat{\mathbf{r}} \mathrm{p}-\hat{\mathrm{r}} \\
\hat{\hat{r}}
\end{gathered}
$$

Use QR with column pivoting (as described in §1.9) to compute

Step two:

$$
Q^{\prime}\left[\begin{array}{ll}
\mathrm{V}_{11}^{\prime} & \mathrm{V}_{21}^{\prime}
\end{array}\right] \mathrm{P}=\left[\begin{array}{ll}
\mathrm{R}_{11} & \mathrm{R}_{12}
\end{array}\right]
$$

and set

$$
\mathrm{XP}=\left[\begin{array}{ll}
\mathrm{X}_{1} & \mathrm{X}_{2}
\end{array}\right]
$$

Step three: Determine $z \in \hat{R^{r}}$ such that $\left\|Y-X_{1} z\right\|_{2}=$ min.

When $\hat{r}>p / 2$ it would be more economical to compute $P$ by applying the $Q R$
 $\left\|\ddot{\mathrm{V}}_{1}^{-1}\right\|_{2}=\left\|\ddot{\mathrm{V}}_{2}^{-1}\right\|_{2}$ (for a proof see Golub and Yan Loan (1983, Theorem 2.4.1). Again it should be noted that these methods are intended to be in the realm of numerical mathematics and do not appeal to any statistical properties per se.

### 1.11 Probability limit

$T$ is the probability limit ( plim ) of the statistic $t_{n}$, derived from a random sample of $n$ observations, if, for any $\epsilon>0$, the probability of $\left|t_{n}-T\right|<\epsilon$, approaches the limit probability 1 as $n \rightarrow \infty$.

## Chapter 2

## COLLINEARITY

One of the assumptions of the linear regression model (1.1), is that the fixed matrix $X$ of independent variables is a full column-rank matrix. Violation of this assumption leads to problems referred to as collinearity. This phenomenon of collinearity and near-collinearity was first described by Ragnar Frisch (1934) and he warned that in ignoring this structure within the independent regressor variables, one runs the risk of determining a regression equation that is absurd. Frisch believed that 'a substantial part of the regression and correlation analyses which have been made on economic data in recent years is nonsense'.

Some authors refer to collinearity and its analysis as multicollinearity (Gunst (1983)), conditioning (Belsley and 0ldford (1986)), confluence analysis (Frisch (1934)), ill-conditioning (Belsley and 01dford (1986)), singularity (Stewart (1987)) and non-orthogonality (Farrar and Glauber (1967)). Ve favour 'collinearity' as collinearity will always involve two or more vectors, and the prefix 'multi' is unnecessary (Kalman (1984)). To some degree all the foregoing terms can be regarded as synonymous. Collinearity can not be described in simple terms as being present or absent. Rather, what is important is the degree of collinearity and what effect this degree can have on the regression model. For the statistician, near-collinearities inflate the variances of regression coefficients and magnify the effects of error in the regression response variable. For the numerical analyst, collinearities combine with rounding errors to introduce inaccuracies in computations. A great deal of time and effort have been devoted to issues related to collinearity (see attached bibliography) and still the subject has abounds with paradoxes, ambiguities and open questions.

In this chapter some light is shed on this problem (although in the opinion of the writer collinearity is not so much a problem as an inherent part of the data set and model). We discuss definitions of collinearity in §2.1, ways of detecting it in §2.2, the effect of collinearity in §2.3, the
issue of centering in §2.4, and the theory of perturbation in §2.5. A summary of approaches to collinearity is presented in §2.6.

### 2.1 Definitions

Several definitions of collinearity appear in the literature, some more descriptive of symptoms than rigorous.

Collinearity is often viewed as a situation in which two or more predictor variables are highly correlated. This statement is inadequate because correlation is a statistical property of random variables, and the regressor variables need not be stochastic, since they could represent preselected variable values in a designed experiment. When there are no preselected variable values the 'correlation' may be merely a characteristic of a particular data set in a particular time period, say, and might not be expected to occur in other data sets of the same type at other times.

Another view is that collinearities are due to 'weak' or 'deficient' data (Farrar and Glauber (1967)). 'Deficient' implies an aberration in the datacollection process; however, collinearities among predictor variables are sometimes an inherent property of the phenomenon under study, in which case 'deficient' data would actually be a misnomer.

In the foregoing statements, there are connotations which imply more about the predictor variables than simply the existence of a collinearity. A formal definition of collinearity should imply nothing more about the predictor variables than the existence of the collinearity. The following definition is due to Johnston (1963), Silvey (1969), Mason et al. (1975) and others, and is in terms of the linear dependence of a set of column vectors, $X_{j}$, of the matrix $X$.

Definition 2.1.1

Vectors $X_{1}, X_{2}, \ldots, X_{p}$ are linearly dependent if there exist non-zero constants $c_{1}, c_{2}, \ldots, c_{p}$ such that

$$
\begin{equation*}
\sum_{j=1}^{p} c_{j} X_{j}=0 \tag{2.1.1}
\end{equation*}
$$

When the relationship in (2.1.1) is exact for a subset of the columns of $X$, exact collinearity exists. When (2.1.1) is only approximately true, nearcollinearity is said to exist.

Gunst (1983) refined definition (2.1.1) as follows:

Definition 2.1.2

Let a multiple linear regression model be defined as in (1.1). If for some specified $\eta \geq 0$ there exists a vector $c^{\prime}=\left[\begin{array}{lll}c_{1} & c_{2} \ldots c_{p}\end{array}\right]$, not all of whose elements are zero, such that

$$
\begin{equation*}
\sum_{\mathrm{j}=1}^{\mathrm{p}} \mathrm{c}_{\mathrm{j}} \mathrm{X}_{\mathrm{j}}=0 \text { with }\|0\| \leq \eta\|\mathrm{c}\| \tag{2.1.2}
\end{equation*}
$$

then a collinearity is said to exist among the predictor variables in X .

For definition (2.1.2) to have a practical meaning $\eta$ must be chosen suitably small. No fixed value of $\eta$ is suitable for all regression analyses since the predictor variables are not scale invariant and $\eta$ can be expected to depend on both $n$ and $p$ (the selection of $\eta$ would be discussed in §2.2.12).

A very loose definition of collinearity in (1.1) is:

Definition 2.1.3

Collinearity exists among the predictor variables if the columns of $X$ are not mutually orthogonal (the matrix $X$ is said to have orthogonal regressors
when it is such that $X^{\prime} X$ is diagonal.). The strength of the collinearity is then determined by some measure of how close the columns of predictor variables are to being linearly dependent.

In definition 2.1 .2 the strength of the collinearity can be gauged by the magnitude of $\|O\| /\|c\|$. In addition, the loose notion of a collinearity being simply equivalent to the non-orthogonality of $X$ (i.e. collinearity exists when the correlation matrix $X^{\prime} X$ is not diagonal) is replaced by equation (2.1.2). Clearly collinearity under (2.1.2) implies non-orthogonality, but not conversely.

Smith and Campbell (1980) maintain that high correlations and other measures of collinearity such as in definitions 2.1 .2 and 2.1 .3 are inadequate since the columns of predictor variables in $X$ can always be transformed so they are mutually orthogonal; i.e., there exists a non-singular matrix $T$ such that $A=X T$ and $A^{\prime} A=I$. In doing so, they imply that model (1.1) is equivalent to the following model:

$$
\begin{align*}
\mathrm{Y} & =(\mathrm{XT})\left(\mathrm{T}^{-1} \beta\right)+\epsilon  \tag{2.1.3}\\
& =A a+\epsilon
\end{align*}
$$

They claim that 'the parameters $\beta$ and $\alpha$ are uniquely related by $\beta=\mathrm{T} \alpha$ and their estimates should be also. It should make no diffference whether $\beta$ is estimated explicitly or implicitly from $\beta=A a$.'

Discussants of the Smith and Campbell paper point out that this type of reasoning ignores an important purpose of the regression analysis. One is usually interested in assessing the influence of the original predictor variables, not the transformed ones, on the response. Their arguments are summarized by Gunst (1983) as follows:
(i) the transformed predictor variables are linear combinations of variables which are not unit- or scale-free,
(ii) in a Bayesian or a sampling-theory framework it is usually easier to obtain prior information on the original model parameters than on arbitrary linear combinations of 'apples and oranges' (see Thisted (1980) section 2.3)
(iii) if one's interest is in specific parameters or parametric functions one must eventually transform back to the original predictor variables, and (iv) if the collinearity is severe enough to affect the computational accuracy of inversion of $X$ it will also affect the computational accuracy of calculations to find and to invert $T$, in which case the $A$ and $a$ are not as well-observed as a superficial reading of (2.1.3) might imply.

The assertion that the collinearity problem is equivalent to a problem of high variances on transformed predictor variables can lead to uncritical, 'black box' mathematical manipulations which obscure the very purpose of a regression analysis: accurate and precise parameter estimation within a chosen and convenient model. Collinearity is an inherent aspect of the analysis and can not be transformed away.

Belsley and 0ldford (1986, p104) give the following definition of illconditioning:

Suppose we have a system of continuous equations

$$
\begin{equation*}
\gamma=\mathrm{f}(\varphi) \tag{2.1.4}
\end{equation*}
$$

where $\gamma, \varphi$, and $\mathrm{f}($.$) are vectors and/or matrices. Then (2.1.4) can describe$ an estimator, a stochastic model, or, in general, any system of interest in which elements $\gamma$ are assumed to be dependent upon elements $\varphi$. Suppose an additive perturbation $\delta \varphi$ in $\varphi$ results in a perturbation in $\gamma$ equal to $\delta \gamma=f(\varphi+\delta \varphi)-f(\varphi)$. A function $g(\delta \varphi) \equiv f(\varphi+\delta \varphi)-f(\varphi)$ may be defined which maps the elements of $\delta \gamma$ of a given domain $\Omega$ to elements $\delta \gamma$ in the corresponding range set 1 . That is,

$$
\mathrm{g}: \delta \varphi \rightarrow \delta \gamma
$$

or

$$
g: \Omega \rightarrow \Lambda
$$

Let $\Lambda^{*}$ consists of all those perturbations $\delta \gamma$ which are considered, $a$ priori, to be reasonable given the set $\Omega$. Concern arises when,
corresponding to some $\delta \varphi$ in $\Omega$, there exists $\delta \gamma$ in $\Lambda$ which is not in $\Lambda^{*}$; i.e. small perturbations in $\delta \varphi$ result in perturbations $\delta \gamma$ which are not 'reasonably small' e.g. small changes in X give rise to large changes in $\left(X^{\prime} X\right)^{-1} X^{\prime} Y$. These considerations lead to the following terms and definitions.

Conditioning analysis: The specification of the conditioning triple $K=\left\{\mathrm{f}, \Omega, \Lambda^{*}\right\}$ followed by a determination of whether $\gamma=\mathrm{f}(\varphi)$ is illconditioned.

Ill-conditioning: Given $K$ and its implied $\Lambda$, if $\Lambda \nsubseteq \Lambda^{*}$ then $\gamma$ is said to be ill-conditioned with respect to $\varphi$ (or $\Omega$ ). Equivalently, one can call the system f ill-conditioned. (The concept of a condition number, denoted by $K(X)$, as a measure of ill-conditioning will be defined and discussed in §2.2.6.)

The condition triple $K$ completely specifies the conditioning analysis. Therefore clear and explicit specification of $K$ is essential. The elements of K must be 'contextually or structurally interpretable', meaning that values assigned to the triplet must be interpretable through one's a priori knowledge of that situation (e.g. in economics it is usually a priori knowledge that allows one to assume some regressor variables may have an error margin of, say, one percent). The values assigned to the triple will depend on the kind of conditioning analysis. The authors distinguish between three kinds of conditioning for the linear model of (1.1):
(i) Data conditioning:

Let $\gamma$ and $\varphi$ be defined by

$$
\begin{equation*}
\gamma=\mathrm{f}(\varphi) \equiv \mathrm{Z} \varphi \tag{2.1.5}
\end{equation*}
$$

(Here we use $\gamma, \mathrm{Z} \varphi$ for generality, in the case of 0 LS their equivalent would be $E(Y)$ or $Y$ and $X \beta$ ).

Consider perturbations

$$
\begin{equation*}
\delta \gamma=\mathrm{g}(\delta \varphi) \equiv \mathrm{Z} \delta \varphi \tag{2.1.6}
\end{equation*}
$$

with domain set

$$
\begin{equation*}
n=\left\{\delta \varphi:\|\delta \varphi\| /\|\varphi\|=m_{1}\right\} \tag{2.1.7}
\end{equation*}
$$

and acceptable response set

$$
\begin{equation*}
\Lambda^{*}=\left\{\delta \gamma:\|\delta \gamma\| /\|\gamma\|>\mathrm{m}_{2}\right\} \tag{2.1.8}
\end{equation*}
$$

That is, perturbations $\delta \varphi$ of fixed relative size $m_{1}$ are required to result in perturbations $\delta \gamma$ whose length is not less than $m_{2}$. (In the linear model context, these conditions imply that changes in $\beta$ should give rise to changes in the data $Y$ of $E(Y)$, of some minimal size.) If this cannot be the case, then the data of $Z$ are ill-conditioned with respect to $\Omega$ (2.1.7). Thus, Z is ill-conditioned if $\|\delta \varphi\| /\|\varphi\|=\mathrm{m}_{1}$ and $\|\delta \gamma\| /\|\gamma\|<\mathrm{m}_{2}$ which implies

$$
\|\delta \varphi\|=\mathrm{m}_{1}\|\varphi\| \text { and }\|\delta \gamma\|<\mathrm{m}_{2}\|\gamma\|
$$

thus,

$$
\begin{equation*}
\|\delta \gamma\| /\|\delta \varphi\|<\left(m_{2}\|\gamma\|\right) /\left(m_{1}\|\varphi\|\right) \tag{2.1.9}
\end{equation*}
$$

This inequality can be directly associated to the definition of collinearity given by Gunst in (2.1.2). From (2.1.2) the data set $Z$ is ill-conditioned if

$$
\|\theta\| /\|c\| \leq \eta
$$

If we take

$$
\begin{equation*}
\eta=\left(\mathrm{m}_{2}\|\gamma\|\right) /\left(\mathrm{m}_{1}\|\varphi\|\right) \tag{2.1.10}
\end{equation*}
$$

and replace $\mathrm{Zc}=\theta$ by $\mathrm{Z} \delta \mathrm{c}=\delta 0$ then we find the data set collinear according to Gunst's definition.
(ii) Estimator conditioning

In 0LS the estimator to be examined is

$$
\begin{equation*}
\hat{\beta}=\left(X^{\prime} X\right)^{-1} X^{\prime} Y=X \dagger Y, \quad\left(X \dagger=\left(X^{\prime} X\right)^{-1} X^{\prime}\right) \tag{2.1.11}
\end{equation*}
$$

In the notation of (2.1.4) $\gamma=\hat{\beta}$, and $\mathrm{f}(\varphi)=\mathrm{f}(\mathrm{X}, \mathrm{Y})$. We consider three types of perturbation under this conditioning:

Firstly: If only $X$ is perturbed, the triplet of the condition analysis is:

$$
\begin{align*}
& \Omega=\left\{\delta \varphi \equiv[\delta \mathrm{X}, \delta \mathrm{Y}]:\|\delta \mathrm{X}\| /\|\mathrm{X}\| \leq \mathrm{m}_{1},\|\delta \mathrm{Y}\|=0\right\}  \tag{2.1.12}\\
& \Lambda^{*}=\left\{\delta \hat{\beta}:\|\delta \hat{\beta}\| /\|\hat{\beta}\| \leq \mathrm{m}_{2}\right\} \tag{2.1.13}
\end{align*}
$$

and

$$
\begin{equation*}
\delta \hat{\beta}=\mathrm{g}(\delta \mathrm{X})=(\mathrm{X}+\delta \mathrm{X}) \dagger \mathrm{Y}-\mathrm{X} \dagger \mathrm{Y} \tag{2.1.14}
\end{equation*}
$$

Belsley and 01dford (1986, p111) suggest choices of 0.01 for $m_{1}$ and $m_{2}$ approximately $20 \mathrm{~m}_{1}$. If the range $\Lambda$ of (2.1.14) based on the domain $\Omega$ of (2.1.12) contains any element not in $\Lambda^{*}$, the DLS estimate is ill-conditioned with respect to $\Omega$. Thus if small relative changes in the $X$ matrix can produce large relative changes in the estimate, the estimate is said to be ill-conditioned.

To determine if an OLS estimate is ill-conditioned we must calcualate

$$
\begin{equation*}
\sup \|\delta \hat{\beta}\| /\|\hat{\beta}\| \leq \mathrm{m}_{2} \text { for } \delta \mathrm{X} \in \Omega \tag{2.1.15}
\end{equation*}
$$

To evaluate (2.1.15) is difficult but Belsley and Oldford (1986) show that (2.1.15) is bounded from above by

$$
\begin{equation*}
m_{1} K(X) R^{-1}\left[2+\left(1-R^{2}\right)^{1 / 2} K(X)\right] \tag{2.1.16}
\end{equation*}
$$

where $K(X)$ is the condition number defined in (2.2.5), and $R$ is the uncentered multiple correlation coefficient defined in Chapter one. As a rough guide the quantity $2 m_{1} K(X) R^{-1}$ is compared to $m_{2}$, and if it is much larger than $\mathrm{m}_{2}$, then $\hat{\beta}$ is declared ill-conditioned. Note that $K(X)$ is an important multiplicative factor: for example, 1 percent relative change in $X$ produces approximately $K(X)$ percent change in $\|\delta \hat{\beta}\| /\|\hat{\beta}\|$.

Secondly : If only $Y$ is perturbed, the triplet of the condition analysis is:

$$
\begin{align*}
& \Omega=\left\{\delta \varphi \equiv[\delta \mathrm{X}, \delta \mathrm{Y}]:\|\delta \mathrm{Y}\| /\|\mathrm{Y}\| \leq \mathrm{m}_{1},\|\delta \mathrm{X}\|=0\right\}  \tag{2.1.17}\\
& \Lambda^{*}=\left\{\delta \hat{\beta}:\|\delta \hat{\beta}\| /\|\hat{\beta}\| \leq \mathrm{m}_{2}\right\}  \tag{2.1.18}\\
& \delta \hat{\beta}=\mathrm{g}(\delta \mathrm{Y})=\mathrm{X} \dagger \delta \mathrm{Y} \tag{2.1.19}
\end{align*}
$$

The relative bound is again the quantity $2 m_{1} K(X) R^{-1}$.

Thirdly: Instead of perturbing $Y$ around its observed value as above we perturb it against its theoretical or expected value, $X \beta$, thus $\mathrm{Y}+\delta \mathrm{Y}=\mathrm{X} \beta+\epsilon$. The conditioning triplet would then be:

$$
\begin{align*}
\Omega & =\left\{\delta \mathrm{Y}:\|\delta \mathrm{Y}\| /\|\mathrm{X} \beta\| \leq \mathrm{m}_{1}\right\}  \tag{2.1.20}\\
\Lambda^{*} & =\left\{\delta \hat{\beta}:\|\delta \hat{\beta}\| /\|\beta\| \leq \mathrm{m}_{2}\right\} \tag{2.1.21}
\end{align*}
$$

and

$$
\begin{equation*}
\delta \hat{\beta}=\mathrm{X} \dagger(\mathrm{Y}+\delta \mathrm{Y})-\mathrm{X} \dagger \mathrm{Y}=\hat{\beta}-\beta \tag{2.1.22}
\end{equation*}
$$

The relative bound for $m_{2}$ is $2 \mathrm{~K}(\mathrm{X})\|\epsilon\| /\|\mathrm{X} \beta\| \leq 2 \mathrm{~m}_{1} \mathrm{~K}(\mathrm{X})$.

Comment: 1. To distinguish this basis for a conditioning analysis from the others, it is called stochastically-based conditioning analysis.
2. Since $\|\delta \hat{\beta}\|^{2}=(\hat{\beta}-\beta)^{\prime}(\hat{\beta}-\beta)$, in determining whether $\hat{\beta}$ is ill-conditioned one also determines whether the maximal squared error of the resulting OLSE is less than some number $\mathrm{m}_{2}\|\beta\|$.
3. $\mathrm{m}_{1}^{-1}$ is chosen to be the minimum 'signal-to-noise' ratio expected to be encountered in model (1.1).
(iii) Criterion conditioning

Parameters $\beta$ are often estimated by minimizing some criterion function of the data and parameters. Let the criterion function be denoted by $Q(X, Y, \beta)$. In OLS, the criterion to be minimized is

$$
\begin{equation*}
Q(X, Y, \beta)=(Y-X \beta)^{\prime}(Y-X \beta) \tag{2.1.23}
\end{equation*}
$$

It is desirable that large changes $\delta \hat{\beta} \equiv \hat{\beta}-\beta$ from $\beta$ should be detectable by the selected criterion function. The condition triplet is given by

$$
\begin{align*}
\Omega & =\left\{\delta \beta:\|\delta \hat{\beta}\|=m_{1}>0\right\}  \tag{2.1.24}\\
\Lambda^{*} & =\left\{\delta Q:\|\delta Q\| /(\inf \|\delta Q\|) \leq m_{2}\right\} \tag{2.1.25}
\end{align*}
$$

and

$$
\begin{equation*}
\delta \underline{q}=\mathrm{g}(\delta \hat{\beta})=2 \delta \hat{\beta}^{\prime}\left(\mathrm{X}^{\prime} \mathrm{X}\right) \delta \hat{\beta} \tag{2.1.26}
\end{equation*}
$$

If $\|\hat{\beta}\| \neq 0$, then $\|\delta \hat{\beta}\|$ in (2.1.24) can be replaced by $\|\delta \hat{\beta}\| /\|\hat{\beta}\|=m_{1}$. The criterion (2.1.23) is ill-conditioned with respect to $\Omega$ if $K^{2}(X)$ is greater then $\mathrm{m}_{2}$.

In (2.1.25) the worst possible effect is used in defining $\Lambda^{*}$, that is the worst effect denotes the boundary of acceptability. However, should perturbations $\delta \hat{\beta}$ in $\Omega$ produce $\delta Q$ 's which are not in $\Lambda^{*}$, then there exist $\hat{\beta}^{*}$ that differ substantially from $\hat{\beta}$ but which are relatively indistinguishable from $\hat{\beta}$ by criterion $Q($.$) . If \inf \|\delta Q\|=0$, for $\delta \hat{\beta} \epsilon \Omega$, then $\beta$ is said to be inestimable with respect to the criterion. The concept of inestimability will be discussed further in $\S 2.3$ under the effects of collinearity.

In conclusion we note how all these conditioning criteria are connected to the condition number. The upper bounds given in the above discussion will be clearer after reading $\S 2.5$ on perturbation theory.

### 2.2 Detecting Collinearity

The distinction between defining and detecting collinearity is very thin. Some authors use ways of detecting collinearity as an implicit method of defining collinearity. For example, (i) if some or one singular value of $X$ is small, then $X$ is collinear; (ii) if the determinant of $X^{\prime} X$ approaches zero, then $X$ is near-singular; (iii) large VIF's or large condition numbers are taken as indicators of collinearity. The main issue is that the user of 0LS must be aware of what is happening in the space of the regressor variables.

Many collinearity measures have been described and we discuss both historical and recent examples.

### 2.2.1 Nature and sensitivity of estimates

An unexpected sign in estimated coefficient of a model, or a low t-statistic, corresponding to a variable which for other than statistical reasons is viewed as an important regressor, or sensitivity of results to deletion of one or more rows or columns, will generally alert an informed reader to data or model inadequacy. These phenomena are sometimes cited as evidence of collinearity, or their occurrence is ascribed to collinearity.

However none of the phenomena are necessary for collinearity to exist, and they cannot therefore become definitive criteria for detecting collinearity, but rather should be grouped under its possible effects.

### 2.2.2 Correlation matrix of scaled regressors

Examining the correlation matrix, $X^{\prime} X$, of the scaled regressors, is a commonly employed procedure because $X^{\prime} X$ is a standard output of most regression packages. High values of the off-diagonal elements of $X^{\prime} X$, i.e. $\left(X^{\prime} X\right)_{i j}$, can be an indication of collinearity, between the two regressors $X_{i}$ and $X_{j}$. This method can detect only pairwise collinearity and it is possible for three or more variates to be collinear while no two of the variates taken alone are highly correlated. Thus, the method may be helpful but not conclusive.

### 2.2.3 Determinant of $\mathrm{X}^{\prime} \mathrm{X}$

When the $X$ matrix is standardized, $0 \leq \operatorname{det}\left(X^{\prime} X\right) \leq 1$, and for $X$ with some columns exactly collinear, $X^{\prime} X$ is singular and $\operatorname{det}\left(X^{\prime} X\right)=0$. $0 n$ the other hand when $\operatorname{det}\left(X^{\prime} X\right)=1$ the columns of $X$ are mutually orthogonal. These facts lead to the notion that in all other cases some degree of collinearity exists and becomes most severe as the determinant approaches zero.

Stewart (1987) disapproves of the use of the determinant to detect collinearities and describes its use as dependent upon the 'unhappy notion that $\operatorname{det}\left(X^{\prime} X\right)$ bears a close relation to near collinearity'. The determinant is excessively sensitive to scaling, for example, $\operatorname{det}\left(k X^{\prime} X\right)=k^{p} \operatorname{det}\left(X^{\prime} X\right)$, A small determinant, may imply little or nothing about the invertibilty of a matrix. For example, the matrix $\mathrm{kI}_{\mathrm{n}}$, whose determinant is $\mathrm{k}^{\mathrm{n}}$ and can be made arbitrarily small, has a simple inverse $\left[k I_{n}\right]^{-1}=k^{-1} I_{n}$, for $k \neq 0$.

### 2.2.4 Departure from orthogonality

Farrar and Glauber (1967) defined collinearity in terms of departures from orthogonality, because they claim that 'orthogonality lends itself easily to formulation as a statistical hypothesis'. The authors then present a three-stage test for the presence, location, and pattern of collinearity. The stages they describe (p104) are:

1) Compute $\operatorname{det}\left(X^{\prime} X\right)$, and test for the severity and presence of collinearity. The authors assume $X^{\prime} X$ is in correlation form and transform $\operatorname{det}\left(X^{\prime} X\right)$ into an approximate Chi Square statistic (see their equation (3)). They claim that 'a meaningful scale is provided against which departure from orthogonality, and hence the gradient between singularity and orthogonality, can be calibrated.'
2) In the second stage, the value of the statistic

$$
\begin{equation*}
\left(\left(X^{\prime} X\right)_{i i}^{-1}-1\right)(n-p) /(p-1) \tag{2.2.1}
\end{equation*}
$$

where $\left(X^{\prime} X\right)_{i i}^{-1}$ is the $i$ - th diagonal element of $\left(X^{\prime} X\right)^{-1}$, is computed and if the underlying distribution of the $i$-th column elements is normal, the statistic has the F -distribution with $\mathrm{p}-1$ and n - p degrees of freedom. (For a complete discussion of the second stage see p102)
3) In the third stage, the authors define the coefficient of partial correlation (denoted by $\mathrm{R}_{\mathrm{ij}}$ 。) between the i - th and j -th columns of X , while all other columns are held constant, as

$$
\begin{equation*}
R_{i j}=-\left(X^{\prime} X\right)_{i j}^{-1}\left(\left(X^{\prime} X\right)_{i j}^{-1}\left(X^{\prime} X\right)_{j j}^{-1}\right)^{-\frac{1}{2}} \tag{2.2.2}
\end{equation*}
$$

A transformation of $\mathrm{R}_{\mathrm{ij}}$. has a known t-distribution, and this property, according to Farrar and Glauber, permits one to assess the significance of the computed $\mathrm{R}_{\mathrm{ij}}$.

The test proposed by Farrar and Glauber (1967) is rejected by several authors, e.g. Kumar (1975) and Belsley et al. (1980). The rejections are mostly based on the fact that Farrar and Glauber (1967) assume that the X matrix to be stochastic, and that the rows of the $X$ matrix are independently distributed. $0^{\prime} H a g e n$ and McCabe (1975) question the validity of 'statistical' interpretation of a measure of collinearity, and Belsley et al. (1980) added that 'there are no distributional implications from the linear regression model for specific null hypotheses (such as orthogonality) on the nature of the data matrix $X$, against which tests can be made.'

### 2.2.5 Smallest singular value

The technique of examining the smallest singular value is very popular amongst numerical analysts and complete discussions can be found in Lawson and Hanson (1974), Stewart (1973,1987) and Wilkinson (1965). Amongst statisticians this idea was suggested by Kendall (1957) and Silvey (1969). Basically one computes the SVD of $X$ and examines the singular values, and collinearity is indicated by the presence of one or more 'small' singular values. Equivalently some authors use the eigenvalues of $X^{\prime} X$ which are just the squares of the singular values of $X$. We note that some authors mean by 'small' a singular value near zero, whereas others mean one or more singular values are 'small' in relation to others.

Stewart (1987) expresses the smallest singular value of X in terms of the Euclidean norm \|.\| as

$$
\begin{equation*}
\inf (\mathrm{X}) \stackrel{\operatorname{def}}{=} \min \|X \nu\|, \text { for }\|\nu\|=1 \tag{2.2.3}
\end{equation*}
$$

Its square is the smallest eigenvalue of the crossproduct matrix $X^{\prime} X$. The justification is the following result due to Eckart and Young (1936), and generalized by Mirsky (1960):
$\inf (X)$ is the spectral norm of the smallest matrix $E$ such that
$X+E$ is exactly collinear.

Thus $\inf (X)$ measures the absolute distance of $X$ from collinearity. The fact that $\inf (X)$ is an absolute measure makes it difficult to interpret in the absence of information about the size of $X$, or putting it differently, it is difficult to know what is meant by the term 'small' which clearly implies a basis of comparison. This question leads to our next measure of collinearity, the condition number.

### 2.2.6 Condition number and condition index

The condition number, $K(X)$, of a matrix $X$ is defined in terms of the (matrix) spectral norm $\|\cdot\|$ as

$$
\begin{equation*}
K(X)=\|X\|\|X \dagger\| \tag{2.2.5}
\end{equation*}
$$

where $\mathrm{X}^{\dagger}$ is the pseudoinverse of X , i.e.

$$
\begin{equation*}
X^{\dagger}=\left(X^{\prime} X\right)^{-1} X^{\prime} \tag{2.2.6}
\end{equation*}
$$

From the $S V D, X=U \Delta V^{\prime}$, the $S V D$ of $X^{\dagger}$ is $V \Delta \dagger^{\prime}$, where $\Delta^{\dagger}$ is the pseudoinverse of $\Delta$, specifically a diagonal matrix with non-zero entries taken as the reciprocals of nonzero diagonal elements of $\Delta$. Hence the singular values of $X^{\dagger}$ are the reciprocals of those of $X$, whence $\inf (X)=\| X \dagger^{-1}$ and it follows that

$$
\begin{equation*}
K^{-1}(X)=\frac{\inf (X)}{\|X\|} \tag{2.2.7}
\end{equation*}
$$

Thus $K^{-1}$ is just $\inf (X)$ scaled by the norm of $X$. In terms of the condition number, the Eckart-Young-Mirsky theorem (2.2.4) reads as follows:

The smallest matrix $E$ for which $X+E$ is collinear satisfies

$$
\begin{equation*}
\|E\| /\|X\|=K^{-1}(X) \tag{2.2.8}
\end{equation*}
$$

$K^{-1}(X)$ gives a lower bound on the relative distance to collinearity.

The degree of ill-conditioning under the measure $K(X)$ depends on how small the minimum singular value is relative to the maximum singular value. From (2.2.7) the condition number is always greater than one. Thus, using the definitions of a spectral norm, $K(X)$ can be written as:

$$
\begin{equation*}
K(X)=\sqrt{\lambda_{1}} / \sqrt{\lambda_{\mathrm{p}}} \geq 1 \tag{2.2.9}
\end{equation*}
$$

The lower bound of $K(X)$ will be reached when all the columns of $X$ are orthonormal. The condition number provides a measure of the sensitivity of the solution to the normal equations, to small changes in X or Y . Useful perturbation theory which has been cast in terms of the condition number, will be discussed in $\S 2.5$.

The condition number is extended to provide a set of the condition indices, and the i-th condition index is defined as

$$
\begin{equation*}
\eta_{\mathrm{i}}=\sqrt{\lambda_{1}} / \sqrt{\lambda_{\mathrm{i}}} \quad i=1, \ldots, \mathrm{p} \tag{2.2.10}
\end{equation*}
$$

The maximum singular value $\left(\sqrt{\lambda_{1}}\right)$ is used as a 'yardstick' against which smallness can be measured. The largest condition index is the condition number and the indices can be ordered as

$$
\mathrm{K}(\mathrm{X})=\eta_{\mathrm{p}} \geq \eta_{\mathrm{p}-1} \geq \ldots \geq \eta_{2} \geq \eta_{1}=1
$$

Any 'large' condition indices can indicate the presence of dependencies. The term 'large' in connection with condition indexes can not be fixed at a definite number. Empirically Belsley et al. (1980) suggest that condition indices around 5-10 indicate weak dependencies that may be starting to affect the regression estimates. Condition indices of 30 to 100 indicate moderate to strong dependencies, and indices larger than 100 indicate serious collinearity problems. Detecting collinearity by condition indices alone is not enough. Some other insight is necessary to find which columns of $X$ are involved in the collinearity.

We may note that the condition number has its defects:
(i) it is too crude for statistical applications, as it uses matrix norms to distil a large amount of information into a single number and other distillations such as VIF's may be better. Belsley (1987), claims the 'full set of condition indices and variance-decomposition proportions (see §2.2.7) must be compared to other measures', and not the condition indices alone.
(ii) The condition number has its own scaling problem which will be discussed in §2.4.

To bring this numerical measure to bear on statistical questions as to what is the effect of the collinearity, we need methods such as the technique of Variance Decomposition due to Belsley et al. (1980).

### 2.2.7 Regression coefficient variance decomposition

Using the SVD, $\mathrm{X}=\mathrm{U} \Delta V^{\prime}$, the variance-covariance matrix of $\hat{\beta}$ is

$$
\operatorname{Var}(\hat{\beta})=\sigma^{2} \sum_{i=1}^{\mathrm{p}} \mathrm{v}_{\mathrm{i}} \mathrm{v}_{\mathrm{i}}{ }^{\prime} / \lambda_{\mathrm{i}}
$$

and for the $j$-th component of $\hat{\beta}$

$$
\begin{equation*}
\operatorname{Var}\left(\hat{\beta}_{\mathrm{j}}\right)=\sigma_{\mathrm{i}}^{2} \sum_{\mathrm{i}=1}^{\mathrm{p}} \mathrm{v}_{\mathrm{ij}}^{2} / \lambda_{\mathrm{i}} \tag{2.2.11}
\end{equation*}
$$

where $v_{i j}$ is the $j$-th element of the $i$-th eigenvector.

Note that (2.2.11) decomposes $\operatorname{Var}\left(\hat{\beta}_{\mathrm{j}}\right)$ into a sum of components. Those components associated with near dependencies (small $\lambda_{i}$ ), will be large relative to the other components. This fact suggests that an unusually high proportion of the variance associated with estimators of two or more coefficents (at least two columns of $X$ are involved in the dependency), and
provides evidence that the corresponding near dependency is causing problems. Not all regression coefficients need be affected. If the j-th variable is not significantly involved in the near-singularity, its coefficient in the i-th eigenvector, $v_{i j}$, will be near zero and its regression coefficient will remain stable even in the presence of the collinearity.

Define the $\mathrm{i}, \mathrm{j}$-th variance-decomposition proportion of the variance of the $j$-th regression coefficent associated with the i-th component as

$$
\begin{equation*}
\phi_{i j}=v_{i j}^{2} / \lambda_{i} \tag{2.2.12}
\end{equation*}
$$

and

$$
\begin{equation*}
\phi_{\text {TOTAL }}=\sum_{i=1}^{p} v_{\mathrm{i} j}^{2} / \lambda_{\mathrm{i}} \tag{2.2.13}
\end{equation*}
$$

Then the variance- decomposition proportions are

$$
\begin{equation*}
\pi_{\mathrm{ij}}=\phi_{\mathrm{ij}} / \phi_{\mathrm{TOTAL}} \tag{2.2.14}
\end{equation*}
$$

Setting up a summary table of the $\pi_{i j}$, patterns of high variance-decomposition proportions will be clear. Such a table will have the following form:

| Associated | Proportions of |  |  |  |
| :--- | :---: | :---: | :---: | :---: |
| Singular values | $\operatorname{Var}\left(\hat{\beta}_{1}\right)$ | $\operatorname{Var}\left(\hat{\beta}_{2}\right)$ | $\ldots$ | $\operatorname{Var}\left(\hat{\beta}_{\mathrm{p}}\right)$ |
| $\sqrt{\lambda_{1}}$ | $\pi_{11}$ | $\pi_{12}$ | $\cdots$ | $\pi_{1 \mathrm{p}}$ |
| $\sqrt{\lambda_{2}}$ | $\pi_{21}$ | $\pi_{22}$ | $\ldots$ | $\pi_{2 \mathrm{p}}$ |
| $\vdots$ | $\vdots$ | $\vdots$ |  | $\vdots$ |
| $\sqrt{\lambda_{\mathrm{p}}}$ | $\pi_{\mathrm{p} 1}$ | $\pi_{\mathrm{p} 2}$ | $\cdots$ | $\pi_{\mathrm{pp}}$ |

A suggested diagnostic procedure for collinearity would then:

1. identify high condition indices (say, greater than 30), associated with
2. high variance-decomposition proportions (say, greater than 0.5) for two or more variances of estimated regression coefficients.

Once the variates involved in each near-dependency have been identified by their high variance-decomposition proportions, the near-dependency itself can be examined, e.g. by regressing one of the variates involved on the others.

In the case of multiple near-singularities, the variance decomposition table may be dominated by the linear transformation with the smallest eigenvalue so that the effect of other near-singularities may not be apparent. To overcome this problem, Rawlings (1988) suggested that the variance contribution of the other linear transformations should then be found by rescaling each column of the table so that the proportions add to one after each removal of a dominating linear transformation.

### 2.2.8 Mixed condition index

Another measure of collinearity suggested by Thisted (1980) involves the ratios of the squares of the eigenvalues and is defined as

$$
\begin{equation*}
\mathrm{mci}=\sum_{i=1}^{\mathrm{p}} \lambda_{\mathrm{p}}^{2} / \lambda_{\mathrm{i}}^{2} \tag{2.2.15}
\end{equation*}
$$

Values of mci near unity indicate high collinearity; values greater than 2.0 indicate little or no collinearity. In this collinearity measure it is the relative size of the smallest eigenvalue to all the others, not its absolute size or its size relative to the largest eigenvalue, that is important.

### 2.2.9 Variance inflation factors

VIF's were defined by Marquardt (1970) as the diagonal elements in the inverse of the correlation matrix of $X^{\prime} X$. Thus the $i$ - th variance inflation factor ( $\mathrm{VIF}_{\mathrm{i}}$ ) is:

$$
\begin{equation*}
\operatorname{VIF}_{i}=\frac{\left(X^{\prime} X\right)_{i}^{-1}}{\left(X_{i}^{\prime} X_{i}\right)^{-1}} \geq 1 \tag{2.2.16}
\end{equation*}
$$

where $\left(X^{\prime} X\right)_{i}^{-1}$ denotes the $i$ - th diagonal elements of $\left(X^{\prime} X\right)^{-1}$ and $X_{i}$ is the i-th column of $X$, not necessarily scaled (centered and standardized). When $X$ is scaled ( $X^{\prime} X$ is in correlation form) then the $V I F_{i}$ can be written as

$$
\begin{equation*}
\mathrm{VIF}_{\mathrm{i}}=\frac{1}{1-\mathrm{R}_{\mathrm{i}}^{2}} \tag{2.2.17}
\end{equation*}
$$

where $R_{i}^{2}$ is the coefficient of determination from the regression of $X_{i}$ on the other independent variables. If there are near-singularities involving $X_{i}$ and the other independent variables, $R_{i}^{2} \simeq 1$ and $V_{i f}$ will be large. If $X_{i}$ is orthogonal to the other independent variables, $R_{i}^{2} \simeq 0$ and $V I F_{i} \simeq 1$. Indications of serious collinearity will be some $V I F_{i}>10$ (if one exists, there are likely be two or more VIF's > 10).

Stewart (1987) demonstrated that the VIF's can be used to detect collinearity in the model (1.1) and called the square root of VIF $\mathrm{V}_{\mathrm{i}}$ the i-th collinearity index, denoted by $\kappa_{\mathrm{i}}$ :

$$
\begin{equation*}
\kappa_{\mathrm{i}}=\left(\operatorname{VIF}_{\mathrm{i}}\right)^{\frac{1}{2}} \tag{2.2.18}
\end{equation*}
$$

or

$$
\mathrm{VIF}_{\mathrm{i}}=\kappa_{\mathrm{i}}^{2}
$$

Stewart also derived and proved the following properties of $\kappa_{\mathrm{i}}$ :
(1) If $X$ has unit column scaling then

$$
\begin{equation*}
\max \left(\kappa_{\mathrm{i}}\right) \leq[\inf (\mathrm{X})]^{-1} \leq \sqrt{\mathrm{p}} \max \left(\kappa_{\mathrm{i}}\right) \tag{2.2.19}
\end{equation*}
$$

and

$$
\max \left(\kappa_{\mathrm{i}}\right) \leq K(\mathrm{X}) \leq \sqrt{\mathrm{p}} \max \left(\kappa_{\mathrm{i}}\right)
$$

(2) The smallest pertubation $\delta \mathrm{X}_{\mathrm{j}}$ in $\mathrm{X}_{\mathrm{j}}$ that will make X exactly collinear satisfies

$$
\begin{equation*}
\frac{\left\|\delta \mathrm{X}_{\mathrm{j}}\right\|}{\left\|\mathrm{X}_{\mathrm{j}}\right\|}=\kappa_{\mathrm{j}}^{-1} \tag{2.2.20}
\end{equation*}
$$

(For a proof see Golub, Hoffman and Stewart (1987)). This statement is simplified to the following rule of thumb: 'if $\kappa_{j}^{-1} \cong 10^{t}$ then perturbations in the $t$-th digits of the components of $X_{j}$ can make the problem collinear. Another way of saying the same thing is that one should be troubled about a model if the number of digits in $\kappa_{\mathrm{j}}$ is not less than the number of accurate digits in the components of $X_{j}$,
(3) For $\mathrm{j}=1,2, \ldots, \mathrm{p}$ and $\mathrm{i} \neq \mathrm{j}$

$$
\begin{equation*}
\max \left(\kappa_{\mathrm{i}}\right) \geq \sqrt{1+\left(\kappa_{\mathrm{j}}^{2}-1\right) /(\mathrm{p}-1)^{2}} \tag{2.2.21}
\end{equation*}
$$

Thus when there is one large coefficient there will be others. Therefore because collinearity is a 'group phenomenon' the naive use of the condition indices in selecting a variable to delete from an unsatisfactory model, is clearly incorrect.

However the condition indices can be used to assess the ill effects of nearcollinearity on regression coefficients. Define the importance of the $j$-th variable

$$
\begin{equation*}
\mathrm{IMP}_{\mathrm{j}}=2 \kappa_{\mathrm{j}} \hat{\sigma} /\|\mathrm{Y}\| \tag{2.2.22}
\end{equation*}
$$

where $\hat{\sigma}$ is the usual estimator of $\sigma$. Choose levels of importance, $\iota_{j}$, (between 0 and 1), above which the $X_{j}$ would be considered important, e.g. if $I M P_{j}>\iota_{j}$ reject the model. For example, a model which produce a $I M P_{j}$ of 0.5 is unsatisfactory, since $X_{j}$ accounts for $50 \%$ of the response and is judged insignificant.

Stewart (1987) also suggests the following regression diagnostics to assess the effects of errors in the variables:

Let the errors in the $j$-th column of $X$ have mean $\mu_{j}$ and variance $\sigma_{j}$, and set

$$
\begin{array}{rlrl}
\mathrm{e}_{\mathrm{j}} & =\left(\mu_{\mathrm{j}}^{2}+\sigma_{\mathrm{j}}^{2}\right)^{1 / 2}, & \text { if there is no constant term, } \\
& =\sigma_{\mathrm{j}} & & \text { otherwise }
\end{array}
$$

Compute

$$
\begin{equation*}
\tau_{\mathrm{j}}=\left[(\mathrm{n}-\mathrm{p})^{1 / 2} \kappa_{\mathrm{j}} \mathrm{e}_{\mathrm{j}}\right] /\left\|\mathrm{X}_{\mathrm{j}}\right\| \tag{2.2.23}
\end{equation*}
$$

If $\tau_{j}>1 / 3$, reject the model. The errors are so influential that the diagnostic procedure cannot be trusted. For more measures to judge the errors see pp75-77 of his paper.

Schall and Dunne (1987b) generalize the VIF's so that they can be associated with more than one parameter, and more generally, with arbitrary sets of linear functions of $\beta$. This generalization, and the distinction of marginal and partial VIF's result in a complete set of collinearity diagnostics. Thus, the generalized set of variance inflation factors can be used to perform the three main tasks of collinearity diagnostics: (i) they measure variance inflation in the parameter estimates (ii) they can be used as collinearity indices and (iii) they can detect the nature and the sources of collinearity.

### 2.2.9.1 VIF as variance ratio

Under model (1.1) the variance of $\hat{\beta}_{\mathrm{i}}$ is given by

$$
V\left(\hat{\beta}_{\mathrm{i}}\right)=\sigma^{2}\left(\mathrm{X}^{\prime} X\right)_{\mathrm{i}}^{-1}
$$

and to avoid confusion $V\left(\hat{\beta}_{\mathrm{i}}\right)$ will be written as $V(X)$ meaning variance of $\hat{\beta}_{i}$ under the full model (1.1). Context will allow us to dispense with the subscript.

Clearly $V\left(\hat{\beta}_{i}\right)$ can also be written as

$$
\begin{equation*}
\mathrm{V}\left(\hat{\beta}_{\mathrm{i}}\right)=\frac{\sigma^{2}}{\left(\mathrm{X}_{\mathrm{i}}^{\prime} \mathrm{X}_{\mathrm{i}}\right)^{-1}} \mathrm{VIF}_{\mathrm{i}} \tag{2.2.24}
\end{equation*}
$$

showing that the variance of the $i$-th regression coefficient is directly proportional the $\operatorname{VIF}_{\mathrm{i}}$ : If $\operatorname{VIF}_{\mathrm{i}}$ is large the $\mathrm{V}\left(\hat{\beta}_{\mathrm{i}}\right)$ will be inflated accordingly.

Consider the model

$$
\begin{equation*}
\mathrm{Y}-\mathrm{X}_{(\mathrm{i})} \beta_{(\mathrm{i})}=\mathrm{X}_{\mathrm{i}} \beta_{\mathrm{i}}+\epsilon \tag{2.2.25}
\end{equation*}
$$

where $X_{(i)}$ and $\beta_{(i)}$ are obtained by dropping the $i$ - th column from $X$ and the i-th component from $\beta$, and where the parameters $\beta_{(i)}$ are assumed to be known. Thus the left hand side of (2.2.25) is observable. Under (2.2.25) the variance of $\beta_{i}$, denoted by $V\left(X_{i}\right)$ is given by

$$
\begin{equation*}
V\left(X_{i}\right)=\sigma^{2}\left(X_{i}^{\prime} X_{i}\right)^{-1} \tag{2.2.26}
\end{equation*}
$$

The i-th variance inflation factor can be written as the ratio

$$
\begin{equation*}
\operatorname{VIF}_{i}=\frac{\left(X^{\prime} X\right)_{i}^{-1}}{\left(X_{i}^{\prime} X_{i}\right)^{-1}}=\frac{\sigma^{2}\left(X^{\prime} X\right)_{i}^{-1}}{\sigma^{2}\left(X_{i}^{\prime} X_{i}\right)^{-1}}=\frac{V(X)}{V\left(X_{i}\right)} \tag{2.2.27}
\end{equation*}
$$

Thus, the VIF $_{\mathrm{i}}$ can be described as the loss of information on $\beta_{\mathrm{i}}$ due to having the covariates $X_{(i)}$ and the unknown parameters $\beta_{(i)}$ in the model (1.1) as compared to the model (2.2.25). Expression of VIF's as variance ratios can be further factorized (Schall and Dunne (1987b), p4) as follows:

Let $X_{(i)}$ be partitioned as

$$
X_{(i)}=\left[\begin{array}{ll}
X_{1} & X_{2} \tag{2.2.28}
\end{array}\right]
$$

and denote the variance inflation $\left(\mathrm{VIF}_{\mathrm{i}}\right)$ due to the covariates $\mathrm{X}_{(\mathrm{i})}$ as

$$
\operatorname{VIF}_{\mathrm{i}}=\operatorname{VIF}_{\mathrm{i}}\left(\mathrm{X}_{(\mathrm{i})}\right)=\operatorname{VIF}_{\mathrm{i}}\left(\left[\begin{array}{ll}
\mathrm{X}_{1} & \mathrm{X}_{2} \tag{2.2.29}
\end{array}\right]\right)
$$

Thus

$$
\begin{align*}
\operatorname{VIF}_{\mathrm{i}} & =\operatorname{VIF}_{\mathrm{i}}\left(\left[\mathrm{X}_{1} \mathrm{X}_{2}\right]\right) \\
& =\frac{V\left(\left[\mathrm{X}_{1} \mathrm{X}_{2} X_{i}\right]\right)}{V\left(X_{i}\right)} \\
& =\frac{V\left(\left[X_{1} X_{i}\right]\right)}{V\left(X_{i}\right)} \cdot \frac{V\left(\left[X_{1} X_{2} X_{i}\right]\right)}{V\left(\left[X_{1} X_{i}\right]\right)} \\
& =\operatorname{VIF}_{i}\left(X_{1}\right) \cdot \operatorname{VIF}_{i}\left(X_{2} \mid X_{1}\right) \tag{2.2.30}
\end{align*}
$$

$\operatorname{VIF}_{\mathrm{i}}\left(\mathrm{X}_{2} \mid \mathrm{X}_{1}\right)$ is a partial variance inflation factor (PVIF), where as $\operatorname{VIF}_{\mathrm{i}}\left(\mathrm{X}_{1}\right)$ is a marginal variance inflation factor (MVIF). The factorization provides a refinement in the collinearity diagnostics. If (say) $\mathrm{VIF}_{\mathrm{i}}\left(\mathrm{X}_{1}\right)$ is large but $\operatorname{VIF}_{\mathrm{i}}\left(\mathrm{X}_{2} \mid \mathrm{X}_{1}\right)$ is small (close to one), then collinearity exists between $X_{(i)}$ and $X_{1}$, and that given $X_{1}$ in the model, having $X_{2}$ in the model does not further inflate the variance of $\hat{\beta}_{i}$.

The authors then generalize the above factorization with respect to single components to a factorization for parameter subsets. Let I $\mathcal{C}\{1, \ldots p\}$
denote the index set specifying the parameter subset $\beta_{I}$, and $X$ be partitioned as $\mathrm{X}=\left[\mathrm{X}_{(\mathrm{I})} \mathrm{X}_{\mathrm{I}}\right]$. Then model (2.2.25) can be written as

$$
\begin{equation*}
Y-X_{(I)} \beta_{(I)}=X_{I} \beta_{I}+\epsilon \tag{2.2.31}
\end{equation*}
$$

The generalized variance inflation factor (GVIF) VIF $_{I}$ is defined as

$$
\begin{equation*}
\operatorname{VIF}_{I}=\frac{\left|\left(X^{\prime} X\right)_{I I}^{-1}\right|}{\left|\left(X_{I}^{\prime} X_{I}\right)^{-1}\right|} \tag{2.2.32}
\end{equation*}
$$

where $\left(X^{\prime} X\right)_{I I}^{-1}$ denotes the submatrix of $\left(X^{\prime} X\right)^{-1}$ formed by the elements which fall into the rows and columns indexed by $I$.

### 2.2.9.2 Jse of YIF's

Firstly: Identify disconnected subsets of variables

By disconnected subsets we mean that collinearity may exist within a subset but not between subsets. This identification is performed by partitioning of $X=\left[X_{(I)} X_{I}\right]$, and the subsets $X_{(I)}$ and $X_{I}$ are said to be disconnected if the GVIF is small. This procedure is repeated to find disconnected subsets of variables in both $X_{(I)}$ and $X_{I}$. If $s \leq p$ disconnected subsets are found then the number $c \leq s$ are subsets containing more than one variable is the 'number of collinearities'.

Secondly: Factorize within each of the c subsets

Insight into the nature of the collinearities can be found by looking at the c subsets individually. Each of them is partitioned as described above. For example assume one disconnected collinear subset is [ $\left.\begin{array}{llll}X_{1} & X_{2} & X_{i}\end{array}\right]$ then VIF ${ }_{i}$ can be partitioned as in (2.2.30) as

$$
\begin{aligned}
\operatorname{VIF}_{\mathrm{i}} & =\operatorname{VIF}_{\mathrm{i}}\left(\mathrm{X}_{1}\right) \cdot \operatorname{VIF}_{\mathrm{i}}\left(\mathrm{X}_{2} \mid \mathrm{X}_{1}\right) \\
& =\operatorname{VIF}_{\mathrm{i}}\left(\mathrm{X}_{2}\right) \cdot \operatorname{VIF}_{\mathrm{i}}\left(\mathrm{X}_{1} \mid \mathrm{X}_{2}\right)
\end{aligned}
$$

If $\operatorname{VIF}_{\mathrm{i}}\left(\mathrm{X}_{1}\right)$ is large, then $\mathrm{X}_{\mathrm{i}}$ and $\mathrm{X}_{1}$ are collinear; similarly, if $\operatorname{VIF}_{\mathrm{i}}\left(\mathrm{X}_{2}\right)$ is large then $X_{i}$ and $X_{2}$ are collinear. If both MVIF's are small, then the PVIF's $\operatorname{VIF}_{\mathrm{i}}\left(\mathrm{X}_{2} \mid \mathrm{X}_{1}\right)$ and $\operatorname{VIF}_{\mathrm{i}}\left(\mathrm{X}_{1} \mid \mathrm{X}_{2}\right)$ must be large, and the collinearity involves $X_{i}$ and at least one variable from both $X_{1}$ and $X_{2}$.

Computations of all the GVIF's can be obtained easily from the correlation matrix of each step in an all subsets regression problem. In the case where $p$ is too large for all subsets regression a subset of all possible GVIF's can be obtained from stepwise regression procedure.

It is recommended by several authors (Stewart (1987), Schall and Dunne (1987b)) that the $X$ matrix should be mean centered before calculating the VIF's (see also §2.4)

### 2.2.10 Signal-to-noise tests

Belsley (1982) suggested that the presence of harmful collinearity and other forms of weak data could be assessed through a test for signal-to-noise. In this test the size of the parameter variance (noise) is assessed relative to the magnitude of the parameter (signal). This test is then combined with other collinearity diagnostics (condition number, variance-decomposition) to provide a test for the presence of harmful collinearity and/or short data.

The general test is constructed as follow:

Assume model (1.1) and the following partitioning of the model

$$
\begin{equation*}
\mathrm{Y}=\mathrm{X}_{\mathrm{a}} \beta_{\mathrm{a}}+\mathrm{X}_{\mathrm{b}} \beta_{\mathrm{b}}+\epsilon \tag{2.2.33}
\end{equation*}
$$

where $X_{a}: \operatorname{nxp}_{\mathrm{a}}, \beta_{\mathrm{a}}: \mathrm{p}_{\mathrm{a}} \times 1, \mathrm{X}_{\mathrm{b}}: \operatorname{nxp}_{\mathrm{b}}, \beta_{\mathrm{b}}: \mathrm{p}_{\mathrm{b}} \mathrm{x} 1$ and $\mathrm{p}_{\mathrm{a}}+\mathrm{p}_{\mathrm{b}}=\mathrm{p}$. Then the marginal distribution of the estimate of $\beta_{\mathrm{b}}$ is

$$
\begin{equation*}
\hat{\beta}_{\mathrm{b}} \sim \mathrm{~N}\left(\beta_{\mathrm{b}}, \mathrm{~V}\left(\hat{\beta}_{\mathrm{b}} \mid \mathrm{X}\right)\right) \tag{2.2.34}
\end{equation*}
$$

where $V\left(\hat{\beta}_{b} \mid X\right)$ is the variance- covariance matrix of $\hat{\beta}_{b}$ conditional on $X$ and is

$$
\begin{equation*}
\mathrm{V}\left(\hat{\beta}_{\mathrm{b}} \mid \mathrm{X}\right)=\sigma^{2}\left(\mathrm{X}_{\mathrm{b}}^{\prime} \mathrm{M}_{\mathrm{a}} \mathrm{X}_{\mathrm{b}}\right)^{-1} \tag{2.2.35}
\end{equation*}
$$

where $M_{a}=I-X_{a}\left(X_{a}^{\prime} X_{a}\right)^{-1} X_{a}^{\prime}$ and $\left(X_{a}^{\prime} X_{a}\right)^{-1}$ is assumed to exist.

Let $\beta_{\mathrm{b}}^{*}$ be any $\mathrm{p}_{\mathrm{b}}$ vector. Then the signal-to-noise of the 0LSE $\hat{\beta}_{\mathrm{b}}$ relative to $\beta_{\mathrm{b}}^{*}$ is defined as

$$
\begin{equation*}
\tau^{2}=\left(\beta_{\mathrm{b}}-\beta_{\mathrm{b}}^{*}\right)^{\prime}\left[\mathrm{V}\left(\hat{\beta}_{\mathrm{b}} \mid \mathrm{X}\right)\right]^{-1}\left(\beta_{\mathrm{b}}-\beta_{\mathrm{b}}^{*}\right) \tag{2.2.36}
\end{equation*}
$$

We note that $(2.2 .36)$ reduces to $\tau=\beta_{\mathrm{b}} /\left[\mathrm{V}\left(\hat{\beta}_{\mathrm{b}} \mid \mathrm{X}\right)\right]^{\frac{1}{2}}$, when $\beta_{\mathrm{b}}^{*}=0, \mathrm{p}_{\mathrm{b}}=1$ and $\beta_{\mathrm{b}}$ is simply the b - th element of $\beta$. The inverse of $\tau$ is often called the coefficient of variation, but $\tau$ itself is the non-centrality parameter of the non-central t-distribution. The term signal-to-noise ratio refers to the magnitude of $\tau$, but the authors use the same name for $\tau^{2}$.

Belsley (1982) shows that

$$
\phi^{2}=\left(\hat{\beta}_{\mathrm{b}}-\beta_{\mathrm{b}}^{*}\right)^{\prime}\left(\mathrm{X}_{\mathrm{b}}^{\prime} \mathrm{H}_{\mathrm{a}} \mathrm{X}_{\mathrm{b}}\right)\left(\hat{\beta}_{\mathrm{b}}-\beta_{\mathrm{b}}^{*}\right) /\left\{\mathrm{p}_{\mathrm{b}} \hat{\sigma}^{2}\right\} \sim \mathrm{F}\left(\mathrm{p}_{\mathrm{b}}, \mathrm{n}-\mathrm{p} ; \tau^{2}\right)
$$

a non-central $F$ with $p_{b}$ and $n-p$ degrees of freedom and non-centrality parameter $\tau^{2}$. Hence, under $H_{0}: \tau^{2}=\tau_{*}^{2}$, we have

$$
\phi^{2} \sim \mathrm{~F}\left(\mathrm{p}_{\mathrm{b}}, \mathrm{n}-\mathrm{p} ; \tau_{*}^{2}\right)
$$

A practical drawback of this test is that it requires knowledge of $\beta_{\mathrm{b}}$ and $V\left(\hat{\beta}_{\mathrm{b}} \mid X\right)$ to stipulate directly a value for $\tau_{\text {* }}^{2}$. The authors propose a practical definition for an 'adequate level' for the signal to noise ratio that does not require $\beta_{\mathrm{b}}$ and $\mathrm{V}\left(\hat{\beta}_{\mathrm{b}} \mid \mathrm{X}\right)$. This measure is an increasing
function of a single, selectable parameter $\gamma \in[0,1)$, and can be made stringent ( $\gamma$ chosen near unity) or relaxed ( $\gamma$ chosen small).

The (1-a) critical F-values for the test of

$$
\begin{aligned}
& \mathrm{H}_{0}: \tau^{2}=\tau_{*}^{2} \text { against } \\
& \mathrm{H}_{\mathrm{a}}: \tau^{2}>\tau_{*}^{2}
\end{aligned}
$$

are tabulated in Belsley (1982, appendix A), where $\gamma$ is chosen by the experimenter.

The authors then set up a strategy to determine harmful collinearity and short data, using results from the sequence of (1) the collinearity diagnostics (condition indices and varance-decomposition) of Belsley et al. (1980) followed by (2) a test for adequate signal-to-noise for the $\hat{\beta}_{b}$ (usually the signals are calculated for individual $\hat{\beta}_{\mathrm{i}}$ 's). The four possible outcomes of these tests are given below:

Inadequacy of
Signal-to-noise

| Collinearity present |  |  |
| :---: | :---: | :---: |
|  | No | Yes |
| No | 1 | 2 |
| Yes | 3 | 4 |

In situations 1 to 4 we may say:
(1) everything seems acceptable
(2) collinearity is present, but has not resulted in inadequate signal-tonoise. This situation augurs well for the use of the model for prediction purposes, particularly, but not necessarily only, if the collinear relations continue into the prediction period.
(3) 'the length of $X_{b}$ is short, e.g. short data', data problems exist, but collinearity is not the culprit.
(4) structural estimation may be adversely affected, but prediction is not necessarily affected if the collinear relations extend into the prediction period.

### 2.2.11 Prior information use

Swamy et al. (1985) argue that any method of detecting collinearity which does not use prior information about $\beta$ may not be successful in diagnosing the presence of real collinearity. The measure of collinearity should take into account the dependence of $\hat{\beta}$ on $X^{\prime} X$ as well as the prior and sample information about $\beta$ as a primary characteristic.

They considered two types of prior information. A non-Bayesian form may investigate $W \beta$ lying on or within the ellipsoid $(\beta-\bar{\beta})^{\prime} W^{\prime} \Delta_{1}^{-1} W(\beta-\bar{\beta})=r^{2}$ with known $W, \bar{\beta}, \Delta_{1}$ and r. To incorporate the non-Bayesian ellipsoid of prior information the method of constrained least squares is proposed. The constrained LSE of $\beta$ (denoted by $\hat{\beta}_{c}$ ) subject to $(\beta-\bar{\beta})^{\prime} W^{\prime} \Delta_{1}^{-1} W(\beta-\bar{\beta})=r^{2}$ is

$$
\begin{equation*}
\hat{\beta}_{c}=\left(X^{\prime} X+\sigma^{2} \mu W^{\prime} \Delta_{1}^{-1} W\right)^{-1}\left(X^{\prime} Y+\sigma^{2} \mu W^{\prime} \Delta_{1}^{-1} W \bar{\beta}\right) \tag{2.2.37}
\end{equation*}
$$

where $\mu$ is chosen such that $\left(\hat{\beta}_{c}-\bar{\beta}\right)^{\prime} W^{\prime} \Delta_{1}^{-1} W\left(\hat{\beta}_{c}-\bar{\beta}\right)=r^{2}$ (see Swamy et al. (1985), p406) for remarks on $\mu$ ).

Swamy et al. (1985) then propose a specific biased version of the estimator in (2.2.37)

$$
\begin{equation*}
\hat{\beta}_{\mathrm{s}}=\left(X^{\prime} X+s^{2} \hat{\mu} W^{\prime} \Delta_{1}^{-1} W\right)^{-1} X^{\prime} Y \tag{2.2.38}
\end{equation*}
$$

where $\mathrm{s}^{2}=(\mathrm{Y}-\mathrm{X} \hat{\beta})^{\prime}(\mathrm{Y}-\mathrm{X} \hat{\beta}) /(\mathrm{n}-\mathrm{p})$ and $\hat{\mu}$ is the value of $\mu$ selected
according to an empirical rule of Thurman et al. (1984). Usually, $\hat{\mu}$ determined by this rule will be a finite, positive quantity (see Swamy et al.(1985, p407)). (Although the authors do not state it explicity, they imply that in (2.2.38) the $\bar{\beta}$ is 0 ).

The modified coefficient of multiple determination is defined as

$$
\begin{equation*}
\mathrm{R}_{\mathrm{s}}^{2}=\left(\hat{\beta}_{\mathrm{s}}^{\prime} \mathrm{X}^{\prime} \mathrm{X}_{\mathrm{s}}+2 \hat{\beta}_{\mathrm{s}}^{\prime} \mathrm{X}^{\prime} \hat{\epsilon}_{\mathrm{s}}\right) /\left(\mathrm{Y}^{\prime} \mathrm{Y}\right) \tag{2.2.39}
\end{equation*}
$$

where $\hat{\epsilon}_{\mathrm{s}}=\mathrm{Y}-\mathrm{X} \hat{\beta}_{\mathrm{s}}$. The 'truncated model' is

$$
\begin{equation*}
Y=X(h)^{\beta}(h)+\epsilon_{h} \tag{2.2.40}
\end{equation*}
$$

in which the $h$-th independent variable is not included but the other $p-1$ variables are included. Then the biased estimates of the truncated model are

$$
\begin{equation*}
\left(\hat{\beta}_{\mathrm{s}}\right)_{-h}=\left[X_{(h)^{\prime}}^{X}(h)+s^{2} \hat{\mu}\left(W^{\prime} \Delta_{1}^{-1} W\right)_{-h},-h\right]^{-1} X_{(h)}^{\prime} Y \tag{2.2.41}
\end{equation*}
$$

where $X_{(h)}$ is obtained by deleting the $h$-th column from $X$, and $\left(W^{\prime} \Delta_{1}^{-1} W\right)_{-h},-h$ is obtained by deleting the $h$ - th row and the $h$-th column from $\left(W^{\prime} \Delta_{1}^{-1} W\right)$. The corresponding coefficient of determination is obtained from

$$
\begin{equation*}
\left[1-\left(R_{\mathrm{s}}\right)_{-\mathrm{h}}^{2}\right]=\left(\hat{\epsilon}_{\mathrm{s}}\right)_{-\mathrm{h}}^{\prime}\left(\hat{\epsilon}_{\mathrm{s}}\right)_{-\mathrm{h}} \tag{2.2.42}
\end{equation*}
$$

The modified incremental contribution of the $h$ - th variable is then

$$
R_{s}^{2}-\left(R_{s}\right)_{-h}^{2}
$$

Based on this increment their measure of collinearity (denoted by $\tilde{m}$ ) is defined as

$$
\begin{equation*}
\tilde{m}=R_{s}^{2}-\sum_{h=1}^{p}\left(R_{s}^{2}-\left(R_{s}\right)_{-h}^{2}\right) \tag{2.2.43}
\end{equation*}
$$

The authors give a lengthy discussion of the bounds of $\tilde{m}$ (see p410-411) and we only give the computationally simpler bounds for $\tilde{m}$ here, namely

$$
\begin{align*}
& m_{L}=R_{s}^{2}-p\left[\max \left(R_{s}^{2}-\left(R_{s}\right)_{-h}^{2}\right)\right], \text { for } 1 \leq h \leq p  \tag{2.2.44}\\
& m_{U}=R_{s}^{2}-p\left[\min \left(R_{s}^{2}-\left(R_{s}\right)_{-h}^{2}\right)\right], \text { for } 1 \leq h \leq p \tag{2.2.45}
\end{align*}
$$

The value of $\tilde{m}$ is interpreted as follows:
(1) If $\tilde{\mathrm{m}} \cong 0$ then collinearity is absent.
(2) If $X^{\prime} X$ is non-diagonal $|\tilde{m}|$ is closer to the value of $m_{L}$ or $m_{U}$ than to zero and then serious collinearity exists.

A Bayesian view may assume $V \beta$ distributed a priori with mean $W \bar{\beta}$ and covariance matrix $V \Delta W^{\prime}$, where $W$ is a known rectangular matrix with full row rank. Here $\bar{\beta}$ and $\Delta$ may be unknown. When the Bayesian prior distribution of $W \beta$ is available the authors derive the posterior mean of $\beta$ as

$$
\begin{equation*}
\hat{\beta}_{\mathrm{G}}=\mathrm{E} \beta+\operatorname{cov}(\beta)\left[\sigma^{2}\left(\mathrm{X}^{\prime} \mathrm{X}\right)^{-1}+\operatorname{cov}(\beta)\right]^{-1}(\hat{\beta}-\mathrm{E} \beta) \tag{2.2.46}
\end{equation*}
$$

where $E \beta$ and $\operatorname{cov}(\beta)$ are respectively the prior mean and the prior covariance matrix of $\beta$. For a derivation of (2.2.46) see Thurman et al. (1984). Kashyap et al. (1984) give an operational version of (2.2.46) which is nearly minimax.

### 2.2.12 Detection methods and Mason's definition

In $\S 2.1$, Mason's definition of collinearity is

$$
\begin{equation*}
\sum_{j=1}^{p} c_{j} X_{j}=0 \quad \text { with }\|\theta\| \leq \eta\|c\| \tag{2.1.2}
\end{equation*}
$$

In this whole section we will assume $X$ is scaled so that $X^{\prime} X$ is in correlation form. Then the above definition (2.1.2) is related to the measures of detection.

Correlation coefficients: Let $r_{i j}$ denote the correlation coefficient between $X_{i}$ and $X_{j}$, as $\left|r_{i j}\right| \rightarrow 1$ the vectors $X_{i}$ and $X_{j}$ approach linear dependence:

$$
\begin{equation*}
x_{i}-\operatorname{sign}\left(r_{i j}\right) x_{j}=0 \text { with }\|0\|=\sqrt{2}\left(1-\left|r_{i j}\right|\right)^{\frac{1}{2}} \tag{2.2.47}
\end{equation*}
$$

From definition (2.1.2), if $\left|r_{i j}\right| \geq 1-\eta^{2}$ then a pairwise collinearity exists between $X_{i}$ and $X_{j}$. The choice of $\eta$ will depend on the goals of the research project.

Variance inflation factors: Consider the definition of VIF $_{\mathrm{i}}$

$$
\begin{equation*}
\operatorname{VIF}_{\mathrm{i}}=\frac{1}{1-\mathrm{R}_{\mathrm{i}}^{2}} \tag{2.2.17}
\end{equation*}
$$

where $R_{i}^{2}$ is the coefficient of determination from the regression of $X_{i}$ on the other independent variables. Denote by $X_{(i)}$ the $X$ matrix with the i-th column deleted and let $b_{i}=\left(X_{(i)} X_{(i)}\right)^{-1} X_{(i)} X_{i}$ be the $0 L S$ regression coefficient when $X_{i}$ is regressed on the remaining $p-1$ columns of $X$ (i.e. on $\left.X_{(i)}\right)$. The corresponding 'residual vector' is then

$$
\begin{equation*}
X_{i}-X_{(i)} b_{i}=0 \quad \text { with }\|\theta\|=\left(1-R_{i}^{2}\right)^{\frac{1}{2}} \tag{2.2.48}
\end{equation*}
$$

Thus if $\left(1-R_{i}^{2}\right)^{\frac{1}{2}}\left(=\left(\text { VIF }_{\mathrm{i}}\right)^{-\frac{1}{2}}\right) \leq \eta\left(1+\mathrm{b}_{\mathrm{i}}{ }^{\prime} \mathrm{b}_{\mathrm{i}}\right)^{\frac{1}{2}}$, equation (2.2.48) defines a collinearity among the predictor variable.

The smallest singular values: From (1.3.2) the SVD of X is

$$
X=\sum_{i=1}^{p} \sqrt{\lambda_{i}} u_{i} v_{i}^{\prime}
$$

so that

$$
\begin{equation*}
\mathrm{X}_{\mathrm{j}}=\sqrt{\lambda_{\mathrm{j}} \mathrm{u}_{\mathrm{j}} \quad \mathrm{j}=1,2, \ldots, \mathrm{p} .} \tag{2.2.49}
\end{equation*}
$$

and

$$
X_{v_{j}}=\sum_{i=1}^{p} X_{i} v_{j i}
$$

where $v_{j i}$ is the $i$-th component of the $j$-th singular vector $v_{j}$. If $\sqrt{\lambda_{j}}$ is suitably small, equation (2.2.49) defines a collinearity since

$$
\begin{equation*}
\sum_{i=1}^{p} X_{i} v_{j i}=\theta \text { with }\|\theta\|=\sqrt{\lambda_{j}} \tag{2.2.50}
\end{equation*}
$$

satisfies definition (2.1.2) when $\sqrt{\lambda_{j}} \leq \eta$. The authors suggested a cutoff value for $\eta$ is about 0.3 (i.e. $\lambda_{\mathrm{j}} \leq \eta^{2}=0.1$ ).

The condition number: If a condition number of 30 is selected as a cutoff value for collinearity, then $\eta=\sqrt{\lambda_{1}} / 30$ is the appropriate cutoff for $\sqrt{\lambda_{\mathrm{p}}}$. Belsley and 0ldford's results in the section on conditioning analysis, $\S 2.1$, can be related to Mason's approach.

### 2.3 Collinearity in Practice

In practice collinearity becomes harmful when estimation or hypothesis testing is influenced more by the relationship between the regressor variables than by the relationship between the response and the regressor variables. Such an influence can result in poor parameter estimates and restrictions on the applicability and generality of the model in use.

### 2.3.1 Sources and origins

Collinearity or near-singularity may arise in several ways (for detailed discussions see Mason, Gunst and Webster (1975) and Rawlings (1988):

1. An over-defined model is one in which there are more regressor variables than observations. This type of model arises frequently in medical research where many elements of information are recorded on each individual in a study.
2. An in-built mathematical constraint in variables that forces them to add to a constant will generate a collinearity. Generating new variables as transformation of other variables can produce a collinearity among the set of variables involved e.g. ratios or powers of variables frequently may be nearly collinear with the original variables.
3. Component variables of a system may show near linear dependencies because of biological or physical constraints of the system (e.g. various measures of size of an organism will show dependencies). Such correlation structures are properties of the system and can be expected to be present in all observations obtained from the system. Gunst (1983) refered to this type of collinearity as 'population-inherent collinearities'
4. Inadequate sampling occurs when the experimenter unknowingly samples only from a subspace of the space of the regressor variables. Collinearities due to sampling deficiencies are a property of the particular data set which has been collected and would not be expected to occur in data sets arising from alternative sampling.
5. Poor experimental design may give rise to collinearities. If possible the levels of the experimental factors are generally chosen in such a way so that the different treatment factors are statistically orthogonal to each other.
6. Outliers can induce artificial collinearites among the predictor variables, and will be discussed in Chapter 8.

Identifying the origin of collinearity is not always possible but it is important to illustrate likely sources in each instance.

### 2.3.2 Effects of collinearity

The impact of collinearity on least squares methodology is very serious if primary interest is in the regression coefficients or if the purpose is to identify 'important' variables in the estimation process. The solution is very unstable, i.e. small changes (random noise or rounding effects ) in the $Y$ or $X$, can cause drastic changes in the estimates of the regression coefficients (e.g. change in sign), and the variances of the regression
coefficients for the regressor variables involved in the near-singularity, become very large.

In discussing effects of collinearities, the notation of Chapter 1 will be used; for convenience we repeat the following:

The variance covariance matrix of the 0LS estimator is given by

$$
\begin{equation*}
V(\hat{\beta})=\sigma^{2} \sum_{i}^{p} \mathrm{v}_{\mathrm{i}} \mathrm{v}_{\mathrm{i}}^{\prime} / \lambda_{\mathrm{i}} \tag{2.3.1}
\end{equation*}
$$

Because the bias is zero the

$$
\begin{equation*}
\operatorname{MSE}(\hat{\beta})=\sigma^{2} \sum_{\mathrm{i}=1}^{\mathrm{p}} \mathrm{v}_{\mathrm{i}} \mathrm{v}_{\mathrm{i}} / \lambda_{\mathrm{i}} \tag{2.3.2}
\end{equation*}
$$

and the TMSE (which is the same as $E\left(L_{1}^{2}\right)$ in section 1.2) is then

$$
\begin{equation*}
\operatorname{TMSE}(\hat{\beta})=\sigma^{2} \sum_{i=1}^{\mathrm{p}} 1 / \lambda_{\mathrm{i}}>\sigma^{2} / \lambda_{\mathrm{p}} \tag{2.3.3}
\end{equation*}
$$

### 2.3.2.1 Geometric interpretation of collinearity

The impact of collinearity can be illustrated geometrically. The figure on p2- 37 and the interpretation below are from Rawlings (1988, p163).

The plane is the $X$-space and the heavy dot, with the shaded area around it, represents $E(Y)$. The shaded area represents the distribution of projections $\hat{Y}$ of $Y$, onto the $X$-space which one might obtain from repeated samplings of the dependent variable. Panels (a) and (b) represent the case where $X_{1}$ and $X_{2}$ are orthogonal. Panels (c) and (d) represent the case where $X_{1}$ and $X_{2}$ are nearly collinear: the angle between the vectors is small. The position of $X_{2}$ relative to $E(Y)$ remains the same in all cases; the position of $X_{1}$ has been shifted to create the collinearity.


### 2.3.2.2. Collinearity of two regressor variables

Assume $X^{\prime} X$ is in correlation form, then

$$
\left(X^{\prime} X\right)^{-1}=\left[\begin{array}{ll}
1 /\left(1-r_{12}^{2}\right) & -\mathrm{r}_{12} /\left(1-\mathrm{r}_{12}^{2}\right)  \tag{2.3.4}\\
-\mathrm{r}_{12} /\left(1-\mathrm{r}_{12}^{2}\right) & 1 /\left(1-\mathrm{r}_{12}^{2}\right)
\end{array}\right]
$$

where $\mathrm{r}_{12}$ is the sample correlation coefficient between $\mathrm{X}_{1}$ and $\mathrm{X}_{2}$. When the correlation between $X_{1}$ and $X_{2}$ increases ( $r_{12}^{2} \rightarrow 1$ ) and applying L'Hospital's rule

$$
\begin{equation*}
\operatorname{Var}\left(\hat{\beta}_{\mathrm{j}}\right) \rightarrow \infty \text { as } \mathrm{r}_{12}^{2} \rightarrow 1 \tag{2.3.5}
\end{equation*}
$$

and

$$
\operatorname{Cov}\left(\hat{\beta}_{1}, \hat{\beta}_{2}\right) \rightarrow \pm \infty \text {, depending on whether } r_{12} \rightarrow \pm 1
$$

A strong pairwise linear relationship between $X_{1}$ and $X_{2}$ results in very large variances and covariances for the estimators of the regression
coefficients. Sastry (1970) was the first to show these limits. His equation (7) is missing a square on the denominator term, ( $1-\mathrm{r}_{12}^{2}$ ) and thus the limit he derived in his equation (8) is incorrect.

For the two regressor variables the estimator of the regression coefficients are

$$
\begin{align*}
\hat{\beta} & =\left(X^{\prime} X\right)^{-1} X^{\prime} Y=\left[\begin{array}{l}
1 /\left(1-r_{12}^{2}\right) X_{1}^{\prime} Y-r_{12} /\left(1-r_{12}^{2}\right) X_{2}^{\prime} Y \\
-\mathrm{r}_{12} /\left(1-\mathrm{r}_{12}^{2}\right) X_{1}^{\prime} Y+1 /\left(1-\mathrm{r}_{12}^{2}\right) X_{2}^{\prime} Y
\end{array}\right] \\
& =\left[\begin{array}{l}
\left(\mathrm{X}_{1}^{\prime} \mathrm{Y}-\mathrm{r}_{12} \mathrm{X}_{2}^{\prime} \mathrm{Y}\right) /\left(1-\mathrm{r}_{12}^{2}\right) \\
\left(-\mathrm{r}_{12} \mathrm{X}_{1}^{\prime} \mathrm{Y}+\mathrm{X}_{2}^{\prime} \mathrm{Y}\right) /\left(1-\mathrm{r}_{12}^{2}\right)
\end{array}\right] \tag{2.3.6}
\end{align*}
$$

Mason et al. (1975) (by using the results in Sastry (1970)) show that as $\mathrm{r}_{12} \rightarrow 1$, and assuming that $X_{1}^{\prime} Y$ and $X_{2}^{\prime} Y$ becomes identical (the equality is justified since if $X_{1}$ and $X_{2}$ have perfect correlation between them, they would each have the same correlation with $Y$ ). If $\mathrm{r}_{12} \rightarrow-1$, the correlation between $X_{1}$ and $Y$ and between $X_{2}$ and $Y$ would be identical but opposite in sign, $\hat{\beta}_{1} \rightarrow X_{1}^{\prime} Y / 2$ and $\hat{\beta}_{2} \rightarrow-X_{1}^{\prime} Y / 2$. The limit operation forces $\hat{\beta}_{1}$ and $\hat{\beta}_{2}$ to become equal but opposite in sign regardless of the true parameter values $\beta_{1}$ and $\beta_{2}$.

The t-ratio for the partial regression coefficient of the first prediction limit for this model (given the second predictor) is

$$
t^{2}=\frac{r_{1 y}^{2}\left(1-r_{12}\right)^{2}(n-3)}{1+2 r_{1 y}^{2} r_{12}-2 r_{1 y}^{2}-r_{12}^{2}}
$$

where $r_{1 y}$ is the sample correlation coefficient between $X_{1}$ and $Y, r_{2 y}$ is the sample correlation coefficient between $X_{2}$ and $Y$, and we let $r_{1 y}=r_{2 y}$, when $r_{12}$ approaches 1 . The limit of the t-ratio for $\left|r_{1 y}\right| \neq 0$ is (Crocker (1971)):

$$
\begin{aligned}
\lim _{\mathrm{r}_{12} \rightarrow 1} \mathrm{t}^{2} & =\lim _{\mathrm{r}_{12} \rightarrow 1} \frac{\mathrm{r}_{1 \mathrm{y}}^{2}\left(1-\mathrm{r}_{12}\right)^{2}(\mathrm{n}-3)}{\left(1-\mathrm{r}_{12}\right)\left(1+\mathrm{r}_{12}-2 \mathrm{r}_{1 \mathrm{y}}^{2}\right)} \\
& =\lim _{\mathrm{r}_{12} \rightarrow 1} \frac{\mathrm{r}_{1 \mathrm{y}}^{2}\left(1-\mathrm{r}_{12}\right)(\mathrm{n}-3)}{\left(1+\mathrm{r}_{12}-2 \mathrm{r}_{1 \mathrm{y}}^{2}\right)} \\
& =\frac{0}{2\left(1-\mathrm{r}_{1 \mathrm{y}}^{2}\right)}=0
\end{aligned}
$$

### 2.3.2.3 Inflation of variance

In the presence of near collinearity $\lambda_{\mathrm{p}} \rightarrow 0$ so that the $\operatorname{var}(\hat{\beta})$ is inflated and $\operatorname{TMSE}(\hat{\beta}) \rightarrow \infty$. From (2.2.24) the individual variance of the $i$ - th element of $\hat{\beta}$ is

$$
\begin{equation*}
\mathrm{V}\left(\hat{\beta}_{\mathrm{i}}\right)=\frac{\sigma^{2}}{\left(\mathrm{X}_{\mathrm{i}} X_{\mathrm{i}}\right)^{-1}} \operatorname{VIF}_{\mathrm{i}} \quad \text { for } \mathrm{i}=1,2, \ldots \mathrm{p} \tag{2.3.7}
\end{equation*}
$$

Thus the variance of the estimator of the i-th regression coefficient is directly proportional to VIF $_{\mathrm{i}}$ 's. If VIF $_{\mathrm{i}}$ is large (indicating collinearity) then the variance will be inflated as well. In the case of a near orthogonal design VIF $\approx 1$, and there is no effect on the variance.

Inflation of the variance will also mean that the null hypothesis $H_{0}: \beta_{i}=0$ will be more likely to be accepted. For a detail discussion on parametric inference see Gunst (1983).

### 2.3.2.4 Unexpected coefficient values and signs

Collinearities can result in $\hat{\beta}_{i}$ to 'have the wrong sign' (Farrar and Glauber (1967)), and magnitudes of values that disagree with well-established theory of previous empirical studies. The notion 'wrong sign' can only be welldefined in a Bayesian framework where the 'correct' sign can be assumed to be known from a prior distribution. Mullet (1976) pointed out that the
'wrong sign' need not be the result of collinearity. 0ther possible explanations for an incorrect sign are: (i) limited range of regressor variable values, (ii) model misspecification, and (iii) computational error. To these we may add (iv) outliers in the response variable and (v) influential cases.

We illustrate the effect of collinearity by considering the OLSE which, from (1.3.5), can be written as

$$
\begin{align*}
\hat{\beta} & =\sum_{i=1}^{p} v_{i} u_{i} \prime Y / \sqrt{\lambda_{i}} \\
& =\sum_{i=1}^{p} v_{i} c_{i} / \lambda_{i} \text { where } c_{i}=u_{i}^{\prime} Y \sqrt{\lambda_{i}} \tag{2.3.8}
\end{align*}
$$

Assume in the SVD of $X$ that the eigenvalues are ordered, e.g.

$$
\lambda_{1} \geq \lambda_{2} \geq \ldots \geq \lambda_{\mathrm{p}-1} \geq \lambda_{\mathrm{p}}
$$

Suppose that $\lambda_{p}$ is much smaller than $\lambda_{p-1}$ (implying a single very strong collinearity), so much smaller that the summation in (2.3.8) is completely dominated by it, e.g.

$$
\hat{\beta} \simeq v_{p} c_{p} / \lambda_{p}
$$

Then we may infer

$$
\begin{equation*}
\hat{\beta}_{\mathrm{j}} \simeq \mathrm{v}_{\mathrm{pj}} \mathrm{c}_{\mathrm{p}} / \lambda_{\mathrm{p}} \tag{2.3.9}
\end{equation*}
$$

Gunst and Mason (1980) claim two characteristics of expression (2.3.9) when a strong collinearity occurs in $X$ and $\mathrm{v}_{\mathrm{pj}} \neq 0$ :
(i) the estimates tend to be large in magnitude due to the multiplier $\lambda_{\mathrm{p}}^{-1}$, unless $\lambda_{p}^{-1}$ is complemented by a small value of $c_{p}$ or $v_{p j}$
(ii) the signs of the estimates tend to be determined more by the collinearity associated with $v_{p}$ than by relationship of the predictor variables with the response: i.e. if $c_{1}>0$, the sign of $\hat{\beta}_{j}$ is the same as that of $v_{p j}$; if $c_{1}<0$, the sign of $\hat{\beta}_{j}$ is opposite that of $v_{p j}$.

The second of these claims may be something of an overstatement, because the SVD admits $\left(-u_{i},-v_{i}\right)$ in place of $\left(u_{i}, v_{i}\right)$, and the relationship of the predicted values and the predictor variables with the response is certainly implicit in $c_{i}=u_{i}^{\prime} Y$, for each value of $i$, even for $i=p$. It may be better to say that the inherent impression associated with $\lambda_{p}$ being as small as suggested, implies that the notion of the sign of the regression coefficients is also imprecise, in that small stochastic variation in $Y$ may result in substantial changes in coefficient estimates, including some sufficiently large as to give rise to apparent sign changes.

### 2.3.2.5 Unstable regression coefficients

In the presence of collinearity a small perturbation in $X$ or $Y$ can result in a relatively unstable regression coefficient $\hat{\beta}_{\mathrm{i}}$. In $\S 2.2$, we have given some insight into situations when perturbation in $X$ will cause harmful collinearity and further results on perturbation will be given in $\S 2.5$

### 2.3.2.6 Linear combinations of regression variables

Poor precision in the estimation of individual parameters does not imply that the estimated model is a poor predictor. Although some individual parameters may be estimated poorly, the $Y$ value may be predicted adequately as the whole vector of $\hat{\beta}_{\mathrm{i}}$ 's is used. When collinearity persists into the prediction area the collinearity is not harmful. Use of the model outside the defined field (extrapolation) will result in poor prediction.

### 2.4. Centering and Standardization of the $X$ matrix

In the literature several conflicting views appear on the question of whether data in the $X$-matrix should be mean-centered before collinearity is assessed. Belsley (1984) contrasts with authors like Stewart (1987), Schall and Dunne (1987b), Gunst (1983), Marquardt (1980) and Marquardt and Snee (1975) who advocate mean centering. There is less argument on the question on whether $X$ should be standardized although the question of 'how the standardizing must be done could be vague'. Stewart (1987) pointed out that any combination of three elements could be standardized: the matrix $X$, the vector $\beta$ (its elements should be close together), or the matrix E (defined in (2.2.4).

The standardizing of $X$ is accomplished by dividing the elements of each column vector by the square root of the sum of squares of the elements, so that the length of each vector, (the root sum of squares of each column) is unity. Standardizing ensures that the measurement of the $X$ variables is uniform (e.g. some columns of the regressor variables may be measured in inches while others could be measured in centimeters) and in fact unit free. Marquardt and Snee (1975) recommended that in some contexts estimates from standardized variates could provide readier parameter interpretability.

Standardizing is essential before eigenanalysis is used for purposes of detecting collinearity, to prevent the eigenanalysis from being dominated by one or two of the independent variables. Independent variables in their original units of measure would contribute unequally to the total sum of squares and, hence, to the eigenvalues.

In §2.2.5 we pointed out that the condition number has its own scaling problem. Stewart (1987) shows that if we partition

$$
\begin{equation*}
X=\left[X_{(p)} X_{p}\right], \tag{2.4.1}
\end{equation*}
$$

where $X_{(p)}$ is all the columns of $X$ except the $p$-th column, and write

$$
\mathrm{X}_{\alpha}=\left[\mathrm{X}_{(\mathrm{p})}{ }^{\left.\alpha \mathrm{X}_{\mathrm{p}}\right] .}\right.
$$

When $\alpha$ approaches zero then

$$
\lim \left\|\mathrm{X}_{\alpha}\right\|=\left\|\mathrm{X}_{(\mathrm{p})}\right\|
$$

and

$$
\left\|\mathrm{X}_{\alpha}^{\dagger}\right\| \cong a^{-1}\left\|\mathrm{X}_{\mathrm{p}}^{\dagger}\right\|
$$

where $X_{p}^{\dagger}$ is the $p$-th row of $X^{\dagger}$. The rest of the matrix becomes so small in comparison to the last row that the norm is only a function of the last row. It follows then that

$$
\begin{equation*}
\mathrm{K}\left(\mathrm{X}_{\alpha}\right)=\alpha^{-1}\left\|\mathrm{X}_{(\mathrm{p})}\right\|\left\|\mathrm{X}_{\mathrm{p}}^{\dagger}\right\| \rightarrow \infty \tag{2.4.2}
\end{equation*}
$$

Thus by scaling down any column of $X$, the condition number can be made arbitrary large and cause 'artificial ill-conditioning'. Therefore it is recommended that before computing the condition number, the columns should be standardized to have unit column length (Belsley et al. (1980, appendix 3B and §3.3))

Centering makes all independent variables orthogonal to the intercept column and hence removes any collinearity that involves the intercept (see the discussion later in this section on collinearity involving the intercept term). 'Nonessential collinearity' (Marquardt and Snee (1975)) is thus removed. Centering is recommended in order to eliminate collinearities which are due to the origins of the predictor variables and it can often provide computational benefits when small storage or low precision prevail.

The effect of centering on VIF's is discussed by Schall and Dunne (1987b). Let $X$ from a model that includes $a$ constant term be partitioned $\mathrm{X}=\left[\begin{array}{ll}1 \mathrm{X} & (\mathrm{i}) \\ X_{\mathrm{i}}\end{array}\right]$, as in (2.2.28), and $\beta=\left[\beta_{0} \beta_{(\mathrm{i})}^{\prime} \beta_{\mathrm{i}}\right]^{\prime}$ where $\beta_{0}$ is the intercept term. Then similarly to (2.2.30)

$$
\begin{equation*}
\operatorname{VIF}_{\mathrm{i}}=\operatorname{VIF}_{\mathrm{i}}(1) \cdot \operatorname{VIF}_{\mathrm{i}}\left(\mathrm{X}_{(\mathrm{i})} \mid 1\right) \tag{2.4.3}
\end{equation*}
$$

The partial variance inflation factor $\operatorname{VIF}_{\mathrm{i}}\left(\mathrm{X}_{(\mathrm{i})} \mid \mathbf{1}\right)$ is the variance inflation factor for the parameter $\beta_{i}$ which would have been obtained from the mean

### 2.5 Perturbation

The numerical data that constitute $X$ and $Y$ have only a limited number of accurate digits and after those digits the data are completely uncertain and hence arbitrary. Situations should be avoided where a 'small' change in the data produces 'large' changes in the solution. Results relating to perturbation of the pseudoinverse $\left(X \dagger=\left(X^{\prime} X\right)^{-1} X^{\prime}\right)$ or the solution of the 0LS problem ( $\hat{\beta}$ ) have been given by a number of authors. For treatments of this perturbation problem or special cases of the problem, see Stewart (1977 and 1973), Lawson and Hanson (1974), and Wedin (1969, 1973). The theory is so well described in textbooks (e.g. Lawson and Hanson (1974)) that only a few results will be given here.

### 2.5.1 Perturbation bounds

The notation for this section is as follows:

Perturbations in $\mathrm{X}, \mathrm{Y}$ and $\hat{\beta}$ will be denoted by $\delta \mathrm{X}, \delta \mathrm{Y}$, and $\delta \hat{\beta}$ respectively. (In previous sections we sometimes use the notation E for $\delta X$, (adopted from Stewart (1987)). The condition number will as previously be denoted by $K(X), \hat{\epsilon}$ will be the residual vector. Let $\hat{\beta}=X \dagger Y$ and $\tilde{\mathrm{b}}=\tilde{\mathrm{X}} \dagger \mathrm{Y}$, where $\tilde{\mathrm{X}}$ is the perturbed regression matrix $\tilde{\mathrm{X}}=\mathrm{X}+\mathrm{E}$ or $\mathrm{X}+\delta \mathrm{X}$.

An important inequality, used in $\S 2.1$ to derive upper bounds for the different conditioning analysis, is from Golub and Van Loan (1983, 6.1-10, p141):

$$
\begin{equation*}
\|\delta \hat{\beta}\| /\|\hat{\beta}\| \leq \mathrm{K}(\mathrm{X}) \mathrm{R}^{-1}\left[2+\left(1-\mathrm{R}^{2}\right)^{1 / 2} \mathrm{~K}(\mathrm{X})\right] \nu+0\left(\nu^{2}\right) \tag{2.5.1}
\end{equation*}
$$

where $\nu=\max \left(\|\delta \mathrm{Y}\| /\|\mathrm{Y}\|,\|\delta \mathrm{X}\| /\|\mathrm{X}\|<[\mathrm{K}(\mathrm{X})]^{-1}\right)$, and R is the uncentered multiple correlation coefficient of $Y$ regressed on $X$. Both $X$ and $X+\delta X$ are assumed to be of full rank. We note that the condition number plays a vital role in the inequality. A bound for (2.1.15) was established by choosing
$\nu=\|\delta X\| /\|X\|=m_{1}$. In Chapter 9 of Lawson and Hanson (1974) (2.5.1) is generalized to different ranks of $X(e . g \operatorname{rank}(X)=k=p<n$, $\mathrm{n}=\mathrm{p}=\mathrm{k}=\operatorname{rank}(\mathrm{X}), \quad \mathrm{p}>\mathrm{n}=\mathrm{k}=\operatorname{rank}(\mathrm{X}))$.

A second inequality is due to Stewart (1987, equation 3.4)

$$
\begin{equation*}
\frac{\|\tilde{\mathrm{B}}-\hat{\beta}\|}{\|\hat{\beta}\|} \leq \mathrm{K}(\mathrm{X}) \frac{\|\mathrm{E}\|}{\|\mathrm{X}\|}+[\mathrm{K}(\mathrm{X})] \frac{\|\mathrm{E}\|}{\|\mathrm{X}\|} \frac{\|\hat{\epsilon}\|}{\|\mathrm{X}\|\|\hat{\beta}\|}+0\left(\|\mathrm{E}\|^{2}\right) \tag{2.5.2}
\end{equation*}
$$

The derivation of (2.5.2) can be achieved by using equation (3.24) of Stewart (1977) and then applying the triangular and submultiplicative inequalities for matrix norms (appearing in §1.9).

A bound like (2.5.2) is unnecessarily pessimistic, due to the repeated applications of the triangular and submultiplicative inequalities, and the fact that each application represents another backing off from sharpness. Whereas numerical analysts are not concerned about the bound (2.5.2) because their errors originate from rounding on a digital computer and can be made very small by using different routines (for a discussion of various routines to calculate $\hat{\beta}$, see Lawson and Hanson (1974)), the statistician on the other hand must deal with measurement errors in recording data, and here the lack of sharpness hurts.

### 2.5.2 Computational accuracy

There exists concern about the numerical accuracy of common computer programs calculating 0LS solutions. One of the first to study this phenomenon was Longley (1967). He found that different regression programs resulted in very different solutions, including differences in sign and first significant digit. More recently this phenomenon was studied by several authors including Beaton, Rubin and Barone (1976), Lesage and Simon (1988), Simon and Lesage (1988), and comprehensively by Randa11 and Rayner (1987).

In their study Beaton, Rubin and Barone (1976) took the Longley data set and added small perturbations to it, before running it on various different regression programs. The results exhibited unstable regression coefficients which could have been as a result of the errors in the input data itself (the data was analyzed by different investigators) or problems with the model (collinearity). They suggested that if one chooses a regression program it should be one that uses the Gram-Schmidt or Givens method. If possible matrix inversion should be avoided. The disadvantage with GramSchmidt is it is too expensive for a general program as it requires a pass over the data for each independent variable.

In their study on numerical accuracy Randall and Rayner (1987) rate the Cholesky algorithm (described in Chapter 7 of Graybill (1976)) as the best method of performing the least square calculations. It has the following desirable properties: it is compact, whether in recording calculator results or in computer storage requirements, is computationally economical, and is computationally stable. The basis of this method is the factorization of

$$
X^{\prime} X=L L^{\prime}
$$

where $L$ is a unique positive lower triangular matrix of order p. Randall and Rayner (1987) also show how the Cholesky algorithm is related to the QR factorisation. They concentrate their study on the Cholesky rather than the $Q R$ as 'its efficiency on the score of economy of computing operation; Chambers (1977) ranks the Cholesky ahead of any $Q \mathrm{R}$ method in this respect.'

A measure of the accuracy of a calculated coefficient b, is Vampler's (1970) 'count of the number of correct significant digits' in $b$, and termed the Wampler accuracy (VA) of $b$. Let $\tilde{b}$ denote the value of the coefficient as computed free from round-off errors, then
(i) $W A=-\log _{10}(|b-\tilde{b}| /|\tilde{b}|)$, if $|b-\tilde{b}| \neq 0$ and $\tilde{b} \neq 0$,
(ii) $W A=-\log _{10}(|b-\tilde{b}|)$, if $|b-\tilde{b}| \neq 0$ and $\tilde{b}=0$,
(iii) $W A=$ the approximate number of decimal digits (d) with which the machine computes, if $b-\tilde{b}=0$.

By studying different data sets and implementing the Cholesky method Randall and Rayner (1987) found:

1. Standardizing $X$ (scaling the columns to unity), has no effect on accuracy. They report 'scaling per se can have no appreciable effect on accuracy with the Cholesky method', and can be associated with serious losses of accuracy (see their discussion on pp21-27).
2. Provided double precision is used, the Cholesky method will satisfy all normal requirements of accuracy without the need for centering (unless centering is appropriate for other reasons), even in cases of extreme ill-conditioning. (The average gain in accuracy due to centering may be between 2-3 digits)
3. If at all possible double precision arithmetic should be used at all times.
4. The accuracy of the Cholesky method can be improved further by the double Cholesky (discussed on their pp36-42).

Lesage and Simon (1988) find that $Q R$ decomposition provides greater accuracy than the Cholesky decomposition for the Wampler benchmark data set. Furthermore in their data set standardizing of $X$ does not have a substantial impact on accuracy, and they show that although centering does improve the accuracy of Cholesky, this improvement is not due to centering itself. Rather the improvement occurs because Cholesky with centering is a hybrid algorithm mixing Cholesky with the more accurate QR decomposition. In summary, they do not recommend standardizing (it appears to cause a slight decline in accuracy), but support that centering appears useful for the Cholesky algorithm, whereas in the $Q R$ algorithm, centering is irrelevant.

Simon and Lesage (1988) study the impact of ill-conditioning on numerical accuracy arising from two types of collinearity: (i) a collinear relation involving the intercept column and each of the independent variables. The 'slope' columns become simultaneously and individually collinear with the
intercept and with each other, (ii) a near linear relationship exists among each of the independent variables but does not involve the intercept term.

They investigated the numerical accuracy of different regression software packages by using an artificial data set called 'the modified Wampler Benchmark (the original Wampler benchmark was described by Vampler (1980), modified by Lesage and Simon (1988) and again modified in this study to take into account collinearity involving the intercept.)

The Modified Wampler benchmark data matrix $X(n x(n-1))$ and dependent variable vector $Y$ ( $n \times 1$ ) are shown in expression (2.5.3)

$$
\mathbf{X}=\left[\begin{array}{ccccc}
1 & 1+\gamma & 1+\gamma & \ldots & 1+\gamma  \tag{2.5.3}\\
1 & \gamma+\epsilon & \gamma & \ldots & \gamma \\
1 & \gamma & \gamma+\epsilon & \ldots & \gamma \\
\vdots & : & : & : & : \\
1 & \gamma & \gamma & \ldots & \gamma+\epsilon \\
1 & \gamma & \gamma & \ldots & \gamma
\end{array}\right] \quad \mathbf{Y}=\left[\begin{array}{c}
(\mathrm{n}-1)+(\mathrm{n}-2) \gamma+\epsilon \\
(\mathrm{n}-2) \gamma+\epsilon \\
(\mathrm{n}-2) \gamma+\epsilon \\
\vdots \\
(\mathrm{n}-2) \gamma+\epsilon \\
(\mathrm{n}-1)+(\mathrm{n}-2) \gamma-\epsilon
\end{array}\right]
$$

where the parameter labelled $\gamma$ controls the severity of ill-conditioning of the first type and $\epsilon$ controls the severity of the second type. $\gamma$ and $\epsilon$ control the severity of both types independently. As the parameter $\epsilon$ is decreased towards zero, the collinearities among all of the last ( $\mathrm{n}-2$ ) columns increases. If $\gamma$ increases the last ( $n-2$ ) columns of the matrix $X$ become more nearly collinear with the intercept column.

Their results can be summarized as follows: Both types of collinearity can adversely affect the numerical accuracy of regression estimator calculation. Integer-valued data were found to react quite differently from non-integer data during the computational process of centering, and the computational benefits associated with centering were found to be much greater for integer-valued data than for non-integer data. Centering the data matrix mitigates but does not prevent accuracy problems. Accuracy problems are not confined to the intercept estimate, but extend to all coefficients. The numerical accuracy of the regression algorithm coefficient estimates are
measured in the manner of Wampler (1980) as:

$$
\begin{equation*}
a c c_{\mathrm{j}}=-\log _{10}\left(\left|\beta_{\mathrm{j}}-\hat{\beta}_{\mathrm{j}}\right| / \beta_{\mathrm{j}}\right) \tag{2.5.4}
\end{equation*}
$$

where $\beta_{\mathrm{j}}$ are the true value and $\hat{\beta}_{\mathrm{j}}$ the estimated value.

If one looks critically at this study there are various important limitations: the Vampler benchmark data matrix $X$, has approximately rank 1 , thus very severe collinearity exists, almost the worst case possible. For the artificial data in (2.5.3) the fitted $Y$ values and the residuals are

$$
\hat{\mathrm{Y}}=\left[\begin{array}{c}
(\mathrm{n}-1)+(\mathrm{n}-2) \gamma  \tag{2.5.5}\\
1+(\mathrm{n}-2) \gamma+\epsilon \\
1+(\mathrm{n}-2) \gamma+\epsilon \\
1+(\mathrm{n}-2) \gamma+\epsilon \\
1+(\mathrm{n}-2) \gamma
\end{array}\right] \quad \hat{\epsilon}=\left[\begin{array}{c}
\epsilon \\
-1 \\
-1 \\
\vdots \\
-1 \\
(\mathrm{n}-2)-\epsilon
\end{array}\right]
$$

which indicate that case 1 in influential, case $n$ is an outlier and the errors in $Y$ are highly correlated.

Furthermore the standard errors of the regression coefficients are very high (calculated for various values of $\gamma$ and $\epsilon$ to be between 78 to 7936, with the possibility of even larger values). If one is working with such high values of SE's of estimators, one may wonder what is the point of centering at all, and of reporting the accuracy of the estimators.

### 2.6 Detecting and handling collinearity

The first step in successfully coping with collinearity is an understanding of the nature and effects of collinearities and an ability to determine when they occur in a data set (Gunst (1983)).

In $\S 2.2$ various ways of detecting collinearities have been discussed. It is of utmost importance that collinearity should be detected. Any method can be used, but it may even be advisable to use several of them. If one were to use for instance VIF's, it is also good practise to look at other methods
(e.g condition number, condition indices, variance-decomposition, and small eigenvalues), to get a multi-faceted insight into the problem. What is important here is the identification of collinearity and not which particular method one uses to detect it. The user may have to calculate some of the measures as most regression computer programs are not designed to warn automatically of the presence of near-collinearities.

Once collinearity is identified no easy remedy is at hand. Any remedy will depend on the objective of the model fitting exercise. If the objective of the study is prediction, collinearity will cause no harmful effects if the collinearity proceeds into the prediction area and no serious extrapolation is made within or outside the row-space of $X$. When primary interest is in estimation of the regression coefficients, other alternatives should be considered. One is augmentation of the data in the directions of the collinearities, e.g. obtain new data or additional data such that the row-space is expanded to remove the near-singularity. Unfortunately this is frequently impractical or impossible.

Subset selection of variables to remove the collinearity should be applied with great care, as this approach may result in removing some of the important regression variables. Hoerl et al. (1986) recommend against subset selection as a general strategy to combat collinearity. In the face of severe collinearity one of the best alternatives is to use those biased estimators that are not so severely effected by collinearity. An array of biased estimators will be described elsewhere in this dissertation. Choosing one of them will depend on the circumstances of the problem, and estimates may perform differently in different situations.

### 2.7 Summary

In this chapter we defined collinearity. We discussed ways of detecting collinearity and the effect of collinearity on regression estimates. The issue of centering and the concepts of perturbation were also introduced. Finally a summary of approaches to collinearity was presented.

## Chapter 3

## PRINCIPAL COYPONENTS REGRESSION MODEL

In this chapter we discuss Principal Components Regression (PCR). We present in $\S 3.1$ a general overview of PCR, in $\S 3.2$ the estimator, in $\S 3.3$ some properties of PCR estimators, in $\S 3.4$ methods of elimination of PC's and in $\S 3.5$ comments and critique of PCR.

### 3.1 Introduction

Let $X^{\prime} X$ be in correlation form and consider model (1.1), and the SVD of $X$ (1.3):

Then

$$
\begin{aligned}
& \mathrm{Y}=\mathrm{X} \beta+\epsilon \quad \text { and } \quad \mathrm{X}=\mathrm{U} \Delta V^{\prime} . \\
& \mathrm{Y}=\mathrm{XV} V^{\prime} \beta+\epsilon \quad \text { where }\left(\mathrm{VV}^{\prime}=\mathrm{I}_{\mathrm{p}}\right)
\end{aligned}
$$

Let $\quad Z=X V, \quad \delta=V^{\prime} \beta$ then $Y=Z \delta+\epsilon$
$\mathrm{Z}=\left[\begin{array}{llll}\mathrm{z}_{1} & z_{2} & \ldots \mathrm{z}_{\mathrm{p}}\end{array}\right]$ is the nxp matrix of principal components. $\mathrm{z}_{\mathrm{i}}$ is the $i$ - th principal component.

$$
\begin{align*}
Z^{\prime} Z & =V^{\prime} X^{\prime} X V \\
& \left.=V^{\prime} V \Delta^{2} V^{\prime} V \quad \text { (from }(1.3 .3)\right) \\
& =\Delta^{2} \tag{3.1.2}
\end{align*}
$$

Using (3.1.1) the LS estimator of $\delta$ is:

$$
\begin{aligned}
\hat{\delta} & =\left(Z^{\prime} Z\right)^{-1} Z^{\prime} Y \\
& =\Delta^{-2} Z^{\prime} Y \quad(\text { from }(3.1 .2)) \\
& =\Delta^{-2} V^{\prime} X^{\prime} Y \\
& =\Delta^{-2} V^{\prime} X^{\prime} X \hat{\beta} \quad\left(X^{\prime} Y=X^{\prime} X \hat{\beta}\right)
\end{aligned}
$$

$$
\begin{align*}
& =\Delta^{-2} V^{\prime} V \Delta^{2} V^{\prime} \hat{\beta} \quad(\text { from 1.3.3) } \\
& =V^{\prime} \hat{\beta} \tag{3.1.3}
\end{align*}
$$

A principal components estimator of $\beta$ is obtained by deleting one or more of the principal components $\left(z_{j}\right)$ and corresponding parameter $\left(\delta_{j}\right)$ and then making a transformation back into the original parameter space.

Equivalently one constructs new parameters ( $\delta_{\mathbf{i}}$ ) and chooses a subset of them. No variables in $X$ are removed, but some PC's are dropped as if the corresponding $\delta_{i}$ were zero. The final model equation is as complex as the one formed using LS (a criticism of PC regression by Wetherill (1986)). This approach may seem feasible rather than meaningful, but in the presence of collinearity it may remove collinearities and is claimed to give an insight into the original objective, namely coefficient estimation.

Let $Z=\left[Z_{1}: Z_{2}\right]$ where $Z_{2}$ contains the $P C$ 's $z_{i}$ that must be retained and $Z_{2}$ the $z_{j}$ to be deleted. Let us assume that $r$ PC's are deleted; $Z_{2}$ will then contain $r$ columns and $Z_{1}$ p-r columns. Then (3.1.1) may be rewritten as

$$
\begin{aligned}
\mathrm{Y} & =\left[\begin{array}{ll}
\mathrm{Z}_{1} & \mathrm{Z}_{2}
\end{array}\right]\left[\begin{array}{l}
\delta_{1} \\
\delta_{2}
\end{array}\right]+\epsilon \quad\left(\mathrm{Z}_{1}: \mathrm{nx}(\mathrm{p}-\mathrm{r}) ; \mathrm{Z}_{2}: \mathrm{nxr}\right) \\
& =\mathrm{Z}_{1} \delta_{1}+\mathrm{Z}_{2} \delta_{2}+\epsilon
\end{aligned}
$$

Let $Z_{1}=X V_{1} \quad$ (the PC's to be retained and $V=\left[\begin{array}{ll}V_{1} & V_{2}\end{array}\right]$ )

Then to remove the $z_{j}$ 's (all those contained in $Z_{2}$, which is orthogonal to $Z_{1}$ ) one sets $\delta_{2}=0$, which implies that $Z_{2} \delta_{2}=0$. Hence we obtain the principal component estimators (PCE) of $\delta_{1}$ :

$$
\begin{aligned}
\hat{\delta}_{1} & =\left(Z_{1}^{\prime} Z_{1}\right)^{-1} Z_{1}^{\prime} Y \\
E\left[\hat{\delta}_{1}\right] & =\Delta^{-2} Z_{1}^{\prime} Y=\delta_{1} \quad\left(\text { unbiased for } \delta_{1}\right)
\end{aligned}
$$

$$
\mathrm{V}\left[\hat{\delta}_{1}\right]=\sigma^{2}\left(Z_{1}^{\prime} Z_{1}\right)^{-1}
$$

In (3.1.1) we use a transformation to reparameterize the model:

$$
\begin{aligned}
\delta & =V^{\prime} \beta \\
{\left[\begin{array}{l}
\delta_{1} \\
\delta_{2}
\end{array}\right] } & =\left[\begin{array}{l}
V_{1}^{\prime} \\
V_{2}^{\prime}
\end{array}\right]
\end{aligned}
$$

PC regression can be viewed as the use of a restricted least squares estimator (Hill et al.(1977) and Johnson et al.(1973)). The restrictions arise from the assumption that $\delta_{2}=0$, which implies that $V_{2}^{\prime} \beta=0$.

Three advantages of seeing the PCE's as equivalent to Restricted Least Squares (RLS) (in which the restrictions are applied to the noisy collinearities), are claimed by Hill et al. (1977) namely:
(i) Data reduction techniques generally impose restrictions upon the associated parameter space.
(ii) Explicit recognition of the restrictions permits an evaluation of their theoretical implications.
(iii) RLS formulation is a convenient vehicle for the statistical investimation of the partitioning of $Z$.

We comment in $\S 3.5$ on these claims.

### 3.2 The PC estimator

Let $\hat{\beta}_{p c}$ denote the PCE of $\beta$. Using (1.3.5) with $\sum_{i=p-r+1}^{p} v_{i} c_{i} / \lambda_{i}$ set to 0 by the r restrictions $V_{2}^{\prime} \beta=0$, the $\operatorname{PCE}\left(\hat{\beta}_{\mathrm{pc}}\right)$ of $\beta$ can be written as

$$
\begin{equation*}
\hat{\beta}_{p c}=\sum_{i=1}^{p-r} v_{i} c_{i} / \lambda_{i} \tag{3.2.1}
\end{equation*}
$$

### 3.3 Properties of $\hat{\beta}_{\mathrm{pc}}$

Some of the properties of PCR discussed in the introduction are restated here for a convenient comparison with ridge regression (Chapter 4).

1. Relationship to OLSE

$$
\begin{align*}
\hat{\beta}_{\mathrm{pc}} & =\hat{\beta}-\sum_{i=\mathrm{p}-\mathrm{r}+1}^{\mathrm{p}} \mathrm{v}_{\mathrm{i}} \mathrm{c}_{\mathrm{i}} / \lambda_{\mathrm{i}}  \tag{3.3.1}\\
& =\hat{\beta}-\mathrm{V}_{2} \mathrm{~V}_{2} \hat{\beta}^{\hat{\beta}} \tag{3.3.2}
\end{align*}
$$

2. Expectation

$$
\begin{aligned}
\mathrm{E}\left(\hat{\beta}_{\mathrm{pc}}\right) & \left.=\left[\mathrm{I}_{\mathrm{p}}-\mathrm{V}_{2} \mathrm{~V}_{2}^{\prime}\right]\right] \beta \\
& =W \beta
\end{aligned}
$$

where $W=\left[I_{p}-V_{2} V_{2}^{\prime}\right]=W^{\prime}=V_{1} V_{1}^{\prime} \quad\left(\right.$ from (1.3.1) and $\left.I_{p}=V_{1} V_{1}^{\prime}+V_{2} V_{2}^{\prime}\right)$.

Thus the bias vector is

$$
\begin{align*}
\mathrm{b} & =W \beta-\beta \\
& =(W-I) \beta \\
& =-V_{2} V_{2}^{\prime} \beta \tag{3.3.3}
\end{align*}
$$

3. Variance

$$
\begin{align*}
V\left(\hat{\beta}_{\mathrm{pc}}\right) & =\mathrm{V}(\mathrm{~W} \hat{\beta}) \\
& =W V(\hat{\beta}) W \quad\left(W^{\prime}=W\right) \\
& =\sigma^{2} W\left(X^{\prime} X\right)^{-1} W \quad(\text { from } \S 1.2) \\
& =\sigma^{2} V_{1} V_{1}^{\prime}\left(V \Delta^{-2} V^{\prime}\right) V_{1} V_{1}^{\prime} \\
& =\sigma^{2} V_{1} V_{1}^{\prime}\left(V_{1} \Delta_{1}^{-2} V_{1}^{\prime}+V_{2} \Delta^{-2} V_{2}^{\prime}\right) V_{1} V_{i}^{\prime} \\
& =\sigma^{2}\left[V_{1} \Delta_{1}^{-2} V_{1}^{\prime}\right] \quad\left(V_{1} \text { and } V_{2} \text { are orthogonal }\right) \\
& =\sigma^{2} \sum_{i=1}^{p-r} V_{i} V_{i}^{\prime} / \lambda_{i} \tag{3.3.4}
\end{align*}
$$

where

$$
\Delta=\left[\begin{array}{ll}
\Delta_{1} 0  \tag{3.3.5}\\
0 & \Delta_{2}
\end{array}\right] \begin{aligned}
& \Delta_{1}:(\mathrm{p}-\mathrm{r}) \mathrm{x}(\mathrm{p}-\mathrm{r}) \\
& \Delta_{2}: \mathrm{rxr}
\end{aligned}
$$

4. 

$$
\begin{align*}
\operatorname{MSE}\left(\hat{\beta}_{\mathrm{pc}}\right) & =\mathrm{E}\left[\left(\hat{\beta}_{\mathrm{pc}}-\beta\right)\left(\hat{\beta}_{\mathrm{pc}}-\beta\right)^{\prime}\right] \\
& =\operatorname{Var}\left(\hat{\beta}_{\mathrm{pc}}\right)+\mathrm{bb} \\
& =\sigma^{2} \sum_{\mathrm{i}=1}^{\mathrm{p}-\mathrm{r}} \mathrm{v}_{\mathrm{i}} \mathrm{v}_{\mathrm{i}}^{\prime} / \lambda_{\mathrm{i}}+\mathrm{V}_{2} \mathrm{~V}_{2}^{\prime} \beta \beta^{\prime} \mathrm{V}_{2} \mathrm{~V}_{2}^{\prime} \tag{3.3.6}
\end{align*}
$$

(from (3.3.4) and (3.3.3))

$$
\text { 5. } \quad \begin{align*}
\operatorname{TMSE}\left(\hat{\beta}_{\mathrm{pc}}\right) & =\operatorname{tr}\left(\mathrm{WV}(\hat{\beta}) \mathrm{W}+[\mathrm{I}-\mathrm{W}] \beta \beta^{\prime}[\mathrm{I}-\mathrm{W}]\right) \\
& =\operatorname{tr}\left(W^{\prime} \mathrm{W} \sigma^{2}\left(X^{\prime} X\right)^{-1}\right)+\operatorname{tr}\left(\beta^{\prime}\left(-\mathrm{V}_{2} \mathrm{~V}_{2}^{\prime}\right)\left(-\mathrm{V}_{2} \mathrm{~V}_{2}^{\prime}\right) \beta\right) \\
& =\sigma^{2} \operatorname{tr}\left(\left[\mathrm{~V}_{1} \Delta_{1}^{-2} \mathrm{~V}_{1}^{\prime}\right]+\beta^{\prime}\left(\mathrm{V}_{2} \mathrm{~V}_{2}^{\prime}\right) \beta\right. \\
& =\sigma^{2} \sum_{i=1}^{-\mathrm{F}} 1 / \lambda_{\mathrm{i}}+\beta^{\prime}\left(\mathrm{V}_{2} \mathrm{~V}_{2}^{\prime}\right) \beta \quad \text { (from A.2) } \tag{3.3.7}
\end{align*}
$$

### 3.4 Eliminating PC's

To address the collinearity problem a logical choice might be to include in $\mathrm{Z}_{2}$ those components corresponding to 'small' eigenvalues of $\mathrm{X}^{\prime} \mathrm{X}$ (Massy (1965)). Inclusion will ensure that the variance of $\hat{\beta}_{\mathrm{pc}}$ will be small. Simultaneously the bias of PCE will increase so that the MSE and TMSE can be large and predictions of the response poor. Components with small eigenvalues may be highly correlated with the dependent variable. Massy (1965) therefore recommended that the selection rule should be one of balancing the preservation of the sample variation in $X^{\prime} X$ against the correlation of the PC's with the response.

The choice of which PC's to eliminate will depend on the purpose of the regression. We will distinguish between two purposes. The linear regression model may be estimated to test some theoretical or structural hypothesis (§3.4.1), or its sole use may be that of a prediction equation (§3.4.2).

### 3.4.1 Theoretical or Structural Norms

PCE's are RLS estimators and can be expressed in terms of structural norms and tests of implicit or explicit restrictions. The term norm here is not exactly the same as the term norm defined in $\S 1.9$ although that idea contributes to this section: most of the test statistics defined in this section involve eigenvalues and traces and the test statistics are ratios of norms. Four plausible tests are (Hill et al. (1977)):

### 3.4.1.1 Classical F Test

To exclude $r$ PC's the following restriction matrix ( R ) is appropriate:

$$
\underset{r \times p}{R}=\left[\begin{array}{cc}
0 & 0 \\
r \times p-r & I_{r}
\end{array}\right]
$$

To evaluate the hypothesis that $R \delta=0$ (setting $Z_{2} \delta_{2}=0$ ), the F statistic can be used. Under the null hypothesis the statistic $V$ is:

$$
\begin{equation*}
V=\left[\frac{(n-p+r) \hat{\sigma}_{R}^{2}-(n-p) \hat{\sigma}^{2}}{(n-p) \hat{\sigma}^{2}}\right]\left[\frac{n-p}{r}\right] \tag{3.4.1}
\end{equation*}
$$

where $\hat{\sigma}_{\mathrm{R}}^{2}$ is the maximum likelihood estimator (MLE) of $\sigma^{2}$ in the restricted model, and $\hat{\sigma}^{2}$ is the MLE of $\sigma^{2}$ in the full model. Now $V$ has a central F-distribution with $r$ and $n-p$ degrees of freedom, when $R \delta=0$, and when $\mathrm{R} \delta \neq 0$ we have

$$
V \sim F(r, n-p ; \gamma)
$$

where

$$
\begin{align*}
\gamma & =\frac{(R \delta)^{\prime}\left[R\left(Z^{\prime} Z\right)^{-1} \mathrm{R}^{\prime}\right]^{-1} \mathrm{R} \delta}{2 \sigma^{2}}  \tag{3.4.2}\\
& \left.=\frac{\delta_{2}^{\prime}\left[\begin{array}{ll}
0 & \left.\left.\mathrm{I}_{\mathrm{r}}\right]\left[\begin{array}{ll}
\Delta_{1} & 0 \\
0 & \Delta_{2}
\end{array}\right]\left[\begin{array}{l}
0 \\
\mathrm{I}_{\mathrm{r}}
\end{array}\right]\right]^{-1} \\
2 \sigma^{2} & \delta_{2}
\end{array}\right.}{\left\{\begin{array}{l}
\text { from (3.1.2) } \\
\text { and (3.3.5) }
\end{array}\right.}\right\}
\end{align*}
$$

$$
=\frac{\sum_{i=p-r+1}^{p} \delta_{i}^{2} \lambda_{i}}{2 \sigma^{2}}
$$

To evaluate the hypothesis that $\mathrm{R} \delta=0$ is to test that $\gamma=0$.
$H_{0}$ will be accepted ( $\gamma \doteqdot 0$ ) when:
(i) $\sigma^{2}$ is very large.
(ii) $\delta_{i}^{2} \lambda_{i}$ is small for every i
(iii) a combination of (i) and (ii) applies

In (ii) when $\lambda_{i}$ is small but $\delta_{i}^{2}$ is large enough, the $F$ test will tend to reject the hypothesis.

### 3.4.1.2 Strong MSE criterion

We say PCE is better in the strong MSE sense if

$$
\begin{aligned}
\operatorname{MSE}\left(\hat{\beta}_{\mathrm{PC}}\right) & \leq \operatorname{MSE}(\hat{\beta}) \\
\text { i.e. } \operatorname{MSE}\left(\mathrm{h}^{\prime} \hat{\beta}_{\mathrm{PC}}\right) & \left.\leq \operatorname{MSE}\left(\mathrm{h}^{\prime} \hat{\beta}\right) \quad \text { (for every } \mathrm{h} \neq 0, \mathrm{~h}: \mathrm{px} 1\right)
\end{aligned}
$$

This inequality holds if and only if for $\gamma$ in (3.4.2) we have

$$
\gamma \leq \frac{1}{2}
$$

Under the null hypothesis that PCE is better than LSE in the strong MSE sense

$$
V-F(r, n-p ; \gamma), \quad \gamma \leq \frac{1}{2}
$$

This norm is less stringent than the classical F test $(\gamma=0)$.

### 3.4.1.3 First weak YSE criterion

A weaker criterion than the above is to say PCE is better than the OLSE if
or

$$
\begin{aligned}
\operatorname{tr}\left(\operatorname{MSE}\left(\hat{\beta}_{\mathrm{PC}}\right)\right) & \leq \operatorname{tr}(\operatorname{MSE}(\hat{\beta})) \\
\operatorname{TMSE}\left(\hat{\beta}_{\mathrm{PC}}\right) & \leq \operatorname{TMSE}(\hat{\beta})
\end{aligned}
$$

where

$$
\begin{aligned}
\operatorname{tr}\left(\operatorname{MSE}\left(\hat{\beta}_{\mathrm{PC}}\right)\right. & =\left(\sigma^{2} \sum_{\mathrm{i}=1}^{\mathrm{p}-\mathrm{r}} 1 / \lambda_{\mathrm{i}}+\left(\beta^{\prime}\left(\mathrm{V}_{2} \mathrm{~V}_{2}^{\prime}\right) \beta\right)\right) \quad(\text { from }(3.3 .5)) \\
& =\sigma^{2} \sum_{\mathrm{i}=1}^{\mathrm{p}-\mathrm{r}} 1 / \lambda_{\mathrm{i}}+\sum_{\mathrm{i}=\mathrm{p}-\mathrm{r}+1}^{\mathrm{p}} \delta_{\mathrm{i}}^{2} \\
\operatorname{tr}(\operatorname{MSE}(\hat{\beta})) & =\sigma^{2} \sum_{\mathrm{i}=1}^{\mathrm{p}} 1 / \lambda_{\mathrm{i}}
\end{aligned}
$$

and

$$
\begin{align*}
\operatorname{tr}\left(\operatorname{MSE}\left(\hat{\beta}_{\mathrm{PC}}\right)\right)-\operatorname{tr}(\operatorname{MSE}(\hat{\beta})) & =\sum_{\mathrm{i}=\mathrm{p}-\mathrm{r}+1}^{\mathrm{p}} \delta_{\mathrm{i}}^{2}-\sigma_{\mathrm{i}=\mathrm{p}-\mathrm{r}+1}^{\mathrm{p}} 1 / \lambda_{\mathrm{i}} \\
& =\sum_{\mathrm{i}=\mathrm{p}-\mathrm{r}+1}^{\mathrm{p}}\left(\delta_{\mathrm{i}}^{2} \lambda_{\mathrm{i}}-\sigma^{2}\right) / \lambda_{\mathrm{i}} \tag{3.4.3}
\end{align*}
$$

In this case $\mathrm{V} \sim \mathrm{F}(\mathrm{r}, \mathrm{n}-\mathrm{p} ; \gamma)$ where the non-centrality parameter $\gamma$ is given by

$$
\begin{equation*}
\gamma \leq \frac{1}{2} d_{L} \operatorname{tr}\left[\left(X^{\prime} X\right)^{-1} R^{\prime}\left[R\left(X^{\prime} X\right)^{-1} R^{\prime}\right]^{-1} R\left(X^{\prime} X\right)^{-1}\right] \tag{3.4.4}
\end{equation*}
$$

and $d_{L}$ is the largest eigenvalue of the expression under the trace operator. This test was proposed by Vallace (1972) and refined by Yancey, Judge and Bock (1973).

For $\hat{\beta}_{\mathrm{PC}}$ to be 'better' than LSE (3.4.3) must be negative. We note that (3.4.3) would be very sensitive to collinearity, due to the division by $\lambda_{i}$ : the smaller any $\lambda_{i}$ the greater the potential advantage of PCE over OLSE under this criterion.

### 3.4.1.4 Confidence interval norm

A norm proposed by Cheng and Iglarch (1976) generates a highly specialized test, and will only be reported here:

Assume $\beta_{\mathrm{i}}$ is of primary interest ( $\beta_{\mathrm{i}}$ is any element of $\beta,\left(\hat{\beta}_{\mathrm{PC}}\right)_{\mathrm{i}}$ is the corresponding element of $\beta_{\mathrm{PC}}$ ). If

$$
\operatorname{Pr}\left(\left|\hat{\beta}_{\mathrm{i}}-\beta_{\mathrm{i}}\right|<\eta\right)<\operatorname{Pr}\left(\left|\left(\hat{\beta}_{\mathrm{PC}}\right)_{\mathrm{i}}-\beta_{\mathrm{i}}\right|<\eta\right)
$$

then $\hat{\beta}_{\mathrm{PC}}$ is judged superior to $\hat{\beta}$. The value of $\eta$ is determined by the investigator ( $\eta$ is the measure of how 'close' the investigator wants the estimated $\beta$ to its true value (see Hill et al. (1977) eq. 13)).

### 3.4.2 Predictive norms

When the purpose of the regression is one of forecasting, a norm involving $\hat{Y}$ $(\hat{Y}=X \hat{\beta})$ is implied. The total mean square error of prediction is (Allen (1974) eq. 6):

$$
\begin{align*}
\operatorname{TMSE}\left(\mathrm{Y}_{\mathrm{f}}\right) & =\mathrm{E}\left(\hat{\mathrm{Y}}-\mathrm{Y}_{\mathrm{f}}\right)^{\prime}\left(\hat{\mathrm{Y}}-\mathrm{Y}_{\mathrm{f}}\right) \\
& =\mathrm{n} \sigma^{2}+\underset{\mathrm{i}=1}{\mathrm{n}} \mathrm{~V}\left(\hat{\mathrm{Y}}_{\mathrm{i}}\right)+\underset{\mathrm{i}=1}{\mathrm{n}}\left[\mathrm{E}\left(\mathrm{Y}_{\mathrm{i}}-\hat{\mathrm{Y}}_{\mathrm{i}}\right)\right]^{2} \\
& =\mathrm{n} \sigma^{2}+\quad \mathrm{E}\left\{(\hat{\mathrm{Y}}-\mathrm{X} \beta)^{\prime}(\hat{\mathrm{Y}}-\mathrm{X} \beta)\right\} \\
& =\mathrm{n} \sigma^{2}+\operatorname{TMSE}(\hat{\mathrm{Y}}) \tag{3.4.5}
\end{align*}
$$

Thus (3.4.5) depends upon the unknown population parameters $\beta$ and $\sigma^{2}$. Allen (1974) suggested using $\hat{\beta}$ and $\hat{\sigma}^{2}$. The disadvantage of this approach is that this measure can not discriminate among various prospective estimators. Two predictive norms that do not have this disadvantage are discussed in §3.4.2.1 and §3.4.2.2.

### 3.4.2.1 Veak predictive MSE criterion

Consider two alternative estimators, $\mathrm{X} \hat{\beta}_{\mathrm{PC}}$ and $\mathrm{X} \hat{\beta}$, of $\hat{\mathrm{Y}}$. The PCE is a better estimator if

$$
\mathrm{E}\left[\left(\mathrm{X} \beta-\mathrm{X} \hat{\beta}_{\mathrm{PC}}\right)^{\prime}\left(\mathrm{X} \beta-\mathrm{X} \hat{\beta}_{\mathrm{PC}}\right)\right] \leq \mathrm{E}\left[(\mathrm{X} \beta-\mathrm{X} \hat{\beta})^{\prime}(\mathrm{X} \beta-\mathrm{X} \hat{\beta})\right]
$$

The $F$ test can again be applied using $V$ as in (3.3.1). $\hat{Y}_{P C}$ will be better than $\hat{Y}$ if for $\gamma$ as in (3.4.2) we have

$$
\begin{aligned}
& \gamma \leq \frac{r}{2} \\
& V-F(r, p-r ; \gamma)
\end{aligned}
$$

Note that the degrees of freedom change from n-p to p-r. Critical values for this test are in Goodnight and Vallace (1972).

### 3.4.2.2 Squared bias of prediction

The term $\quad{ }_{\Sigma}^{n}\left[E\left(Y_{i}-\hat{Y}_{i}\right)\right]^{2}$ is the total squared bias of the predictions, and appears $i_{i}$ in (3.4.5). It is also called the squared bias of prediction. If a norm based on this criterion is considered then it is possible to construct a multiple comparison test for all general hypotheses for the model in (1.1). This norm does not involve $\sigma^{2}$ and it only considers the squared bias. This procedure and other references to it can be found in Hill et al.(1977).

### 3.5 Comments and critique

Some of the above criteria are quite overwhelming. One simple way to eliminate PC's is to use restricted least squares. For more details on the theory and tests of hypotheses involving such restrictions, any textbook on linear models can be consulted. (eg. Searle (1971)).

Rawlings (1988) states that a good working rule is to eliminate those PC's that satisfy two conditions:
(i) they cause serious variance inflation because of small eigenvalues
(ii) the corresponding estimated regression coefficients ( $\left.\hat{\beta}_{\mathrm{PC}}\right)_{\mathrm{i}}$ are not significantly different from zero.

Lott (1973) described the Strong MSE criterion (§3.4.1.2) and shows in a Monte Carlo study that under this criterion PCE performed better than OLSE.

For more examples on the PC method see Hill et al. (1977), Wetherill (1986) and Jolliffe $(1972,1973,1982)$. In Jolliffe (1982) four examples are given to demonstrate that it is not always the PC's with small eigenvalues that are eliminated.

The justification that a MSE criterion can underpin the PCR approach is technically correct, but is only adequate when the actual $V_{2}^{\prime} \beta$ parameter values are sufficiently small. Specifically we require

$$
\begin{aligned}
\operatorname{MSE}\left(\hat{\beta}_{\mathrm{PC}}\right) & =\sigma^{2} \sum_{\mathrm{i}=1}^{\mathrm{p}-\mathrm{r}} \mathrm{v}_{\mathrm{i}} \mathrm{v}_{\mathrm{i}}^{\prime} / \lambda_{\mathrm{i}}+\mathrm{V}_{2} \mathrm{~V}_{2}^{\prime} \beta \beta^{\prime} \mathrm{V}_{2} \mathrm{~V}_{2}^{\prime} \\
& \leq \operatorname{var}(\hat{\beta}) \\
\operatorname{Var}(\hat{\beta}) & =\sigma^{2} \sum_{\mathrm{i}=1}^{\mathrm{p}-\mathrm{r}} \mathrm{v}_{\mathrm{i}} \mathrm{v}_{\mathrm{i}}{ }^{\prime} / \lambda_{\mathrm{i}}+\sigma^{2} \sum_{\mathrm{p}-\mathrm{r}+1}^{\mathrm{p}} \mathrm{v}_{\mathrm{i}} \mathrm{v}_{\mathrm{i}}^{\prime} / \lambda_{\mathrm{i}}
\end{aligned}
$$

and the data may not admit belief that $V_{2}^{\prime} \beta$ is in fact small enough. These comments underpin suggestions of Rawlings (1988) previously mentioned.

PC Regression initially amounts to choosing that part of the parameter space spanned by $V_{2}^{\prime} \beta$, about which one has virtually no intrinsic information in the data since $\operatorname{Var}\left(V_{2}^{\prime} \hat{\beta}\right)=\operatorname{diag}\left(\lambda_{\mathrm{p}-\mathrm{r}+1}^{-1}, \ldots, \lambda_{\mathrm{p}}^{-1}\right)$ is very large. One then claims to have perfect information that $V_{2}^{\prime} \beta=V_{2}^{\prime} \beta_{\mathrm{PC}}=0$ in the sense that $\operatorname{Var}\left(V_{2} \hat{\beta}_{\mathrm{PC}}\right)=0$. This additional perfect information does not affect that part of the parameter space spanned by $V_{1}{ }^{\prime} \beta$, which is intrinsically independent of any information whatsoever on $V_{2}{ }^{\prime} \beta$. In consequence looking at $V_{1} \hat{\beta}_{\mathrm{PC}}$ will tell us nothing different from what is already known from $V_{1}{ }^{\prime} \hat{\beta}$. Any alleged advantages in estimation from $P C R$ must result from the apparently reduced variances of estimators of $A X \beta$ where

$$
A X=A X V V^{\prime}=B V^{\prime}=\left[B_{1}: B_{2}\right]\left[\begin{array}{l}
V_{1}^{\prime} \\
Y_{2}^{\prime}
\end{array}\right],
$$

since it follows directly from (3.3.5) that

$$
\begin{aligned}
\operatorname{Var}(\mathrm{AX} \hat{\beta}) & =\operatorname{Var}\left(\mathrm{B}_{1} \mathrm{~V}_{1}^{\prime} \hat{\beta}\right)+\operatorname{Var}\left(\mathrm{B}_{2} \mathrm{~V}_{2}^{\prime} \hat{\beta}\right) \\
& =\operatorname{Var}\left(\mathrm{B}_{1} \mathrm{~V}_{1}^{\prime} \hat{\beta}_{\mathrm{PC}}\right)+\operatorname{Var}\left(\mathrm{B}_{2} \mathrm{~V}_{2}^{\prime} \hat{\beta}\right) \\
& >\operatorname{Var}\left(\mathrm{B}_{1} \mathrm{~V}_{1}^{\prime} \hat{\beta}_{\mathrm{PC}}\right)+\operatorname{Var}\left(\mathrm{B}_{2} \mathrm{~V}_{2}^{\prime} \hat{\beta}_{\mathrm{PC}}\right) \quad\left(\mathrm{B}_{2} \neq 0\right)
\end{aligned}
$$

Nonetheless such advantages are artificial, and if the restrictions are not justifiable a priori, it makes more sense to simply isolate the region in the parameter space about which further information data could be usefully applied. Equivalently: it should suffice to specify that $V_{2}^{\prime} \beta$ is not well approximated by the $X \beta$ of the model.

The point is that a restriction such as $V_{2}^{\prime} \beta=V_{2}^{\prime} \hat{\beta}$ imposed on the model will admit $V^{\prime} \hat{\beta}_{\mathrm{PC}}=V^{\prime} \hat{\beta}$ and will reduce the variance of the estimator from $r \quad \mathrm{p}$-r $\sigma^{2} \Sigma 1 / \lambda_{i}$ to $\sigma^{2} \Sigma 1 / \lambda_{i}$ without introducing bias. Such an approach is no $i=1$
more artificial than the ordinary $P C$ restriction $V_{2}^{\prime} \beta=0$. It is also rank-reducing in whatever sense $P C$ estimation is rank reducing.

### 3.6 Summary

In this chapter we presented a general overview of PCR. We defined the estimator and its properties. We discussed methods of eliminating PC's and gave some critique of PCR.

## Chapter 4

## RIDGE REGRESSION

### 4.1 Introduction

Ridge regression (RR) was first proposed by Hoerl and Kennard (1970a and 1970b). Ridge is an estimation procedure for the model (1.1) based on adding a small positive constant to the diagonal elements of $X^{\prime} X$ (which we assume in this section will be in correlation form). Instead of inverting $X^{\prime} X$ which is ill-conditioned in the presense of collinearity, one inverts $X^{\prime} X+k I$, where $k$ is chosen in such a way that the estimators of $\beta$ become stable.

### 4.2 The RR estimator

The ridge estimator is defined as the solution to

$$
\begin{equation*}
\left(X^{\prime} X+k I\right)^{-1} \hat{\beta}_{R}=X^{\prime} Y, \quad k \geq 0 \tag{4.2.1}
\end{equation*}
$$

### 4.3 Properties of $\hat{\beta}_{\mathbf{R}}$

1. Relationship to 0LSE

$$
\begin{align*}
\hat{\beta}_{R} & =W X^{\prime} Y \quad \text { where } \quad W=\left(X^{\prime} X+k I\right)^{-1}=W^{\prime} \\
& =W X^{\prime} X \hat{\beta} \\
& =Z \hat{\beta} \quad \text { where } \quad Z=W X^{\prime} X=I-k W=Z^{\prime} \tag{4.3.1}
\end{align*}
$$

Note

1. If $\mathrm{k}=0$, then $\hat{\beta}_{\mathrm{R}}=\hat{\beta}$.
2. If $k \rightarrow \infty$, then $Z$ approaches 0 , and $\hat{\beta}_{R} \rightarrow 0$.
3. Expectation

$$
\begin{align*}
\mathrm{E}\left(\hat{\beta}_{\mathrm{R}}\right) & =\mathrm{E}\left(W X^{\prime} X \hat{\beta}\right) \\
& =W X^{\prime} \mathrm{X} \mathrm{E}(\hat{\beta}) \\
& =W X^{\prime} X \beta \\
& =\mathrm{Z} \beta \\
& =\beta-\mathrm{kV} \beta \tag{4.3.2}
\end{align*}
$$

Thus, $\hat{\beta}_{\mathrm{R}}$ is biased for $\beta$ and we denote the bias of $\hat{\beta}_{\mathrm{R}}$ by $\mathrm{b}=-\mathrm{k} W \beta$

3 Variance

$$
\begin{align*}
V\left(\hat{\beta}_{R}\right) & =V\left(W X^{\prime} X \hat{\beta}\right) \\
& =W X^{\prime} X V(\hat{\beta}) X^{\prime} X W \\
& =\sigma^{2} W X^{\prime} X W \\
& =\sigma^{2} Z\left(X^{\prime} X\right)^{-1} Z^{\prime} \tag{4.3.3}
\end{align*}
$$

4. The mean squared error (MSE) of $\hat{\beta}_{\mathrm{R}}$ satisfies

$$
\begin{align*}
\operatorname{MSE}\left(\hat{\beta}_{\mathrm{R}}\right) & =\mathrm{E}\left[\left(\hat{\beta}_{\mathrm{R}}-\beta\right)\left(\hat{\beta}_{\mathrm{R}}-\beta\right)^{\prime}\right] \\
& =\mathrm{V}\left(\hat{\beta}_{\mathrm{R}}\right)+\mathrm{bb}^{\prime} \\
& =\sigma^{2} Z\left(\mathrm{X}^{\prime} \mathrm{X}\right)^{-1} \mathrm{Z}^{\prime}+\mathrm{k}^{2} \mathrm{~W} \beta \beta^{\prime} \mathrm{W}^{\prime} \quad(\text { from (4.3.2) and (4.3.3)) } \\
& =\sigma^{2} \mathrm{Z}\left(\mathrm{X}^{\prime} \mathrm{X}\right)^{-1} \mathrm{Z}^{\prime}+[\mathrm{Z}-\mathrm{I}] \beta \beta^{\prime}[\mathrm{Z}-\mathrm{I}]^{\prime} \tag{4.3.4}
\end{align*}
$$

5. The total mean squared error (TMSE) of $\hat{\beta}_{\mathrm{R}}$ satisfies.

$$
\begin{aligned}
\operatorname{TMSE}\left(\hat{\beta}_{\mathrm{R}}\right) & =\operatorname{tr}\left(\operatorname{MSE}\left(\hat{\beta}_{\mathrm{R}}\right)\right) \\
& =\operatorname{tr}\left(\mathrm{V}\left(\hat{\beta}_{\mathrm{R}}\right)\right)+\beta^{\prime}(Z-I)^{\prime}(Z-I) \beta
\end{aligned}
$$

Let

$$
\begin{aligned}
\gamma_{1}(\mathrm{k}) & =\operatorname{tr}\left[\mathrm{V}\left(\hat{\beta}_{\mathrm{R}}\right)\right] \\
& =\sigma^{2} \operatorname{tr}\left[\mathrm{Z}^{\prime} Z\left(X^{\prime} \mathrm{X}\right)^{-1}\right] \\
& =\sigma^{2} \operatorname{tr}\left[\mathrm{~V} \Delta^{2}\left(\Delta^{2}+\mathrm{kI}\right)^{-1} \mathrm{~V}^{\prime} \mathrm{V}\left(\Delta^{2}+\mathrm{kI}\right)^{-1} \Delta^{2} \mathrm{~V}^{\prime} \mathrm{V} \Delta^{-2} \mathrm{~V}^{\prime}\right] \quad(\text { from (1.3.1)) } \\
& =\sigma^{2} \operatorname{tr}\left[\Delta^{2}\left(\Delta^{2}+\mathrm{kI}\right)^{-2}\right] \\
& =\sigma^{2} \sum_{\mathrm{i}} \lambda_{\mathrm{i}}\left(\lambda_{\mathrm{i}}+\mathrm{k}\right)^{-2}
\end{aligned}
$$

Also let

$$
\begin{aligned}
\gamma_{2}(\mathrm{k}) & =\beta^{\prime}(\mathrm{Z}-\mathrm{I})^{\prime}(\mathrm{Z}-\mathrm{I}) \beta \\
& =\mathrm{k}^{2} \beta^{\prime} W^{\prime} W \beta
\end{aligned}
$$

thus

$$
\begin{equation*}
\operatorname{THSE}\left(\hat{\beta}_{\mathrm{R}}\right)=\sigma_{\mathrm{i}}^{2} \underset{\mathrm{i}}{\Sigma} \lambda_{\mathrm{i}}\left(\lambda_{\mathrm{i}}+\mathrm{k}\right)^{-2}+\mathrm{k}^{2} \beta^{\prime} W^{\prime} W \beta \tag{4.3.5}
\end{equation*}
$$

Comments:
(i) $\gamma_{1}(\mathrm{k})$ is the total variance of the ridge estimators.
(ii) $\gamma_{2}^{1}(k)$ will be considered as the measure of total squared bias that is introduced when $\hat{\beta}_{\mathrm{R}}$ is used rather than $\hat{\beta}$
(iii) $\gamma_{1}(k)$ is a continuous monotonically decreasing function of $k$.
(iv) $\gamma_{2}^{1}(k)$ is a continuous monotonically increasing function of $k$, as can be easily seen from

$$
\begin{aligned}
\gamma_{2}(\mathrm{k}) & =\mathrm{k}^{2} \beta^{\prime}\left[\mathrm{X}^{\prime} \mathrm{X}+\mathrm{kI}\right]^{-2} \beta \\
& =\mathrm{k}^{2} \beta^{\prime} \mathrm{V}\left(\Delta^{2}+\mathrm{kI}\right)^{-2} V^{\prime} \beta \\
& =k^{2} \Sigma \delta_{\mathrm{i}}^{2} /\left(\lambda_{\mathrm{i}}+k\right)^{2} \text { for } \delta=V^{\prime} \beta
\end{aligned}
$$

(v) The squared bias $\gamma_{2}(k)$ approaches $\beta^{\prime} \beta$ as an upper limit, as can be easily seen from

$$
\begin{aligned}
\gamma_{2}(\mathrm{k}) & =\mathrm{k}^{2} \Sigma \delta_{\mathrm{i}}^{2} /\left(\lambda_{\mathrm{i}}+\mathrm{k}\right)^{2} \\
& =\Sigma \delta_{\mathrm{i}}^{2} /\left(\lambda_{\mathrm{i}} / \mathrm{k}+1\right)^{2}, \text { thus }
\end{aligned}
$$

$$
\begin{aligned}
\lim _{k \rightarrow \infty} \gamma_{2}(k) & =\delta^{\prime} \delta \\
& =\beta^{\prime} \beta
\end{aligned}
$$

(vi) Hoerl and Kennard (1970b) show that there always exists a value of $k$ such that $\operatorname{TMSE}\left(\hat{\beta}_{R}\right)<\operatorname{TMSE}(\hat{\beta})$. But this value of $k$ is in general bounded by functions of the unknown parameter values $\beta_{i}$ (i.e $\left(\mathrm{k}<\sigma^{2} / \delta_{\max }^{2}\right)$ ). An optimal choice of $\mathrm{k}=\mathrm{p} \sigma^{2} / \beta^{\prime} \beta$ is operationalized and discussed in §4.6.3.
(vii) The so-called Admissibility Condition of Mayer and Willke (1973) states: A class of estimators $E$ will be called (mean square) admissible if for every problem there is an estimator e in E such that $G(e)<G(\hat{\beta})=\operatorname{tr}(V(\hat{\beta})) \quad$ (where the symbol $G$ of Mayer and Willke is the THSE in our terms).
6. The residual sum of squares for the ridge estimator is

$$
\begin{align*}
\operatorname{SSE}\left(\hat{\beta}_{\mathrm{R}}\right) & =\left(Y-X \hat{\beta}_{\mathrm{R}}\right)^{\prime}\left(Y-X \hat{\beta}_{\mathrm{R}}\right) \\
& =Y^{\prime} Y-Y^{\prime} X \hat{\beta}_{\mathrm{R}}-\hat{\beta}_{\mathrm{R}}^{\prime} X^{\prime} Y+\hat{\beta}_{\mathrm{R}} X^{\prime} X \hat{\beta}_{\mathrm{R}} \\
& =Y^{\prime} Y-\hat{\beta}_{\mathrm{R}}^{\prime} X^{\prime} Y-\left[\mathrm{Y}^{\prime} X-\hat{\beta}_{\mathrm{R}} X^{\prime} X\right] \hat{\beta}_{\mathrm{R}} \\
& =Y^{\prime} Y-\hat{\beta}_{\mathrm{R}} X^{\prime} Y-\left[\hat{\beta}_{\mathrm{R}}^{\prime} \mathrm{W}^{-1}-\hat{\beta}_{\mathrm{R}} X^{\prime} X\right] \hat{\beta}_{\mathrm{R}} \\
& =Y^{\prime} Y-\hat{\beta}_{\mathrm{R}} X^{\prime} Y-\hat{\beta}_{\mathrm{R}}^{\prime}\left[\mathrm{W}^{-1}-X^{\prime} X\right] \hat{\beta}_{\mathrm{R}} \\
& =Y^{\prime} Y-\hat{\beta}_{\mathrm{R}} X^{\prime} \mathrm{Y}-\mathrm{k} \hat{\beta}_{\mathrm{R}}^{\prime} \hat{\beta}_{\mathrm{R}} \\
& >\operatorname{SSE}(\hat{\beta}) \tag{4.3.6}
\end{align*}
$$

7. The squared length of $\hat{\beta}_{\mathrm{R}}$ is less than the squared length of $\hat{\beta}$ for all $\mathrm{k}>0$ (Hoerl and Kennard (1970a, equation 2.8).

$$
\begin{align*}
\text { Proof } \quad \hat{\beta}_{R}^{\prime} \hat{\beta}_{R} & =\hat{\beta}^{\prime} Z^{\prime} Z \hat{\beta} \quad(\text { from }(4.3 .1)) \\
& =\hat{\beta}^{\prime} V \Delta^{2} V^{\prime} V\left(\Delta^{2}+k I\right)^{-1} V^{\prime} V\left(\Delta^{2}+k I\right)^{-1} V^{\prime} V \Delta^{2} V^{\prime} \hat{\beta} \\
& (\text { from }(4.3 .1) \text { and }(1.3 .1)) \\
& =\hat{\beta}^{\prime} \Delta^{4}\left(\Delta^{2}+k I\right)^{-2} \hat{\beta}  \tag{4.3.7}\\
& <\hat{\beta}^{\prime} \hat{\beta} \quad \text { for } \mathrm{k}>0,
\end{align*}
$$

Since $\lambda_{i}^{2}\left(\lambda_{i}+k\right)^{-2}<1$ for all i.
8. Marquardt (1970) and Allen (1974) stated that the ridge estimator is formally equivalent to an LSE for the extended data that is obtained when the actual data are supplemented by a fictitious set of data points taken from an orthogonal experiment $H_{k}\left(H_{k}\right.$ is orthogonal up to a scale factor $\sqrt{\mathrm{k}}$ ). The response is taken to be zero for each of these data points.

Let the augmented model be

$$
\left[\begin{array}{l}
\mathrm{Y}  \tag{4.3.8}\\
0
\end{array}\right]=\left[\begin{array}{l}
\mathrm{X} \\
\mathrm{H}_{\mathrm{k}}
\end{array}\right] \beta+\left[\begin{array}{l}
\epsilon_{1} \\
\epsilon_{2}
\end{array}\right]
$$

If $X^{*}$ is the augmented matrix, we obtain

$$
\begin{aligned}
\hat{\beta} & =\left(X^{*} X^{*}\right)^{-1} X^{*} Y^{\prime} \\
& =\left(X^{\prime} X+H_{k}^{\prime} H_{k}\right)^{-1} X^{\prime} Y
\end{aligned}
$$

since $H_{k}$ is orthogonal up to a scale factor, i.e. $H_{k}{ }^{\prime} H_{k}=k I_{p}$. It can be seen that the 0LSE of the augmented matrix is equivalent to a ridge estimator. We claim that the notion of pseudo-observations at zero for the response variable implicity suggests shrinkage on the usual estimates of $\beta$ (towards zero).

### 4.4 Estimation of the ridge parameter $k$

One way of choosing $k$ involves the use of the 'ridge trace'. The ridge trace is formed by plotting $\hat{\beta}_{\mathrm{R}}$ against k , as k varies through the interval $[0,1]$. Marquardt (1970) stated that values of $k$ typically are $10^{-4}<k<1$. A value of $k$ is chosen at the point where the estimates of all $\beta_{i}$ have stabilized. If the estimates of some of the regression coefficients do not stabilize, the corresponding variables could be good candidates for deleting (dropping) from the regression equation.

Wetherill et al. (1986) suggested that the VIF can also be used by choosing k such that the largest VIF's lie between one and ten. This criterion is Marquardt's (1970) rule of thumb.

Usually there is a range of values of $k$ which give equivalent results from a practical point of view. One problem of selecting an estimate of $k$ (viewed as a parameter), via the ridge trace, is that the estimator of $k$ (as a function of $y$ ) is technically a random variable. This fact complicates the theory of confidence limits and hypothesis tests, due to introduction of bias and dependence. Because of the introduced bias, the MSE and the THSE of $\hat{\beta}_{\mathrm{R}}$ are dependent upon the true unknown coefficient vector $\beta$ (Marquardt and Snee (1975)). In $\S 4.6 .2$ the fact that $k$ is dependent on the data is investigated further.

For more comments on the use of the ridge trace and the choice of $k$, the reader is referred to Hoerl and Kennard (1970b, p65), Galpin (1978, p38-58) and section $\S 4.6 .3$ of this chapter. Examples in the articles to which references have been made, illustrate the effect of varying k. The ridge regression procedure is summarized in recipe form by Rawlings (1988, p340). An alternative method to the ridge trace for estimating $k$ is also discussed by Rawlings.

### 4.5. General comments

### 4.5.1 Bayesian view

The ridge estimator can also be derived using Bayesian theory. From a Bayesian viewpoint a prior distribution is defined for the parameters ( $\beta, \sigma^{2}$ ) and this 'prior knowledge' is then incorporated with the model (1.1) to find the conditional and posterior distributions for the data. (For a general discussion on Bayes estimation in the linear model see Lindley and Smith (1972)).

In the linear regression model (1.1) if we assume that

$$
\text { (i) } \quad \epsilon \sim \mathrm{N}\left(0, \sigma^{2} \mathrm{I}\right)
$$

and we further assume a priori knowledge on $\beta$ :

$$
\text { (ii) } \beta \sim \mathrm{N}(a, \mathrm{Z}) \quad(a, \mathrm{Z} \text { are known })
$$

then the posterior distribution of $\beta$ given $Y$, i.e. of the conditional random variable ( $\beta \mid \mathrm{Y}$ ), can be found using Bayes Theorem (Hsiang (1976)):

$$
\begin{equation*}
(\beta \mid Y) \sim N\left(\left(Z^{-1}+1 / \sigma^{2} X^{\prime} X\right)^{-1}\left(Z^{-1} a+1 / \sigma^{2}\right) X^{\prime} Y ;\left(Z^{-1}+1 / \sigma^{2} X^{\prime} X\right)^{-1}\right) \tag{4.5.1}
\end{equation*}
$$

If $a=0$ and $Z=\left(\sigma^{2} / k\right) I$ then the posterior mean of $\beta$ becomes

$$
\left(X^{\prime} X+k I\right)^{-1} X^{\prime} Y
$$

the ridge estimator and with posterior variance $\sigma^{2} \mathrm{~W}$. Thus a ridge estimator can be considered as the posterior mean based on a normal prior parameter distribution with mean 0 and variance matrix $\left(\sigma^{2} / k\right) I$. Observe that this variance matrix implies a connection between the scale of the residuals (through $\sigma^{2}$ ) and the degree of the collinearity (through $k$ ). This posterior mean estimator is equivalent to the LSE after augmenting the data $Y=X \beta+\epsilon$ with the statistically independent pseudo-observations $a=k^{\frac{1}{2}} \beta+\epsilon_{*}$, which incorporates the belief that $\mathrm{k}^{\frac{1}{2}} \beta=a-\epsilon_{*}$. In effect the estimator $\hat{\beta}$ is shrunk towards the value $\mathrm{k}^{\frac{1}{2}} a$. The Bayesian argument does not depend on X being in a scaled form.

### 4.5.2 Critique of ridge regression

Smith and Campbell (1980) present the following objections to the use of ridge estimators

1. The standardising and scaling of the $X^{\prime} X$ matrix to correlation form, is generally presented before the assumption of a unique parameter k.
2. The prior information that is implicit in ridge regression is inadequately justified because of the assumptions that
(i) parameters have zero means ( $a=0$ )
(ii) parameters have zero covariances
(iii) parameters have identical variances
'Ridge regression's weakness is the use of a loose representation of a priori beliefs and the relevance at times of ad hoc pseudoinformation'.
3. Ridge users may obtain estimates that have preassigned values. 'If one blithely manipulates the data with no regard for the appropriateness of the implicit ridge prior distributions, the ridge estimates may literally be anything.,

### 4.6 Generalized Ridge Regression Estimators

In this section the generalized ridge regression estimator (GRRE) using a diagonal matrix $K$ in place of $k I$ is discussed and we distinguish between the non-stochastic and the stochastic ridge parameter. We discuss in §4.6.1 the GRRE when $K$ is non-stochastic and in $\S 4.6 .2$ the GRRE when $K$ is estimated and therefore stochastic. In the latter case the choice of $K$ usually depends on the data and this complicates the sampling distributions of GRRE

### 4.6.1 I non- stochastic

### 4.6.1.1 The GRR estimator

By using the model (3.1.1) to obtain the PC reparameterization we have, $\mathrm{Y}=\mathrm{Z} \delta+\epsilon$. Hoerl and Kennard (1970b) defined the general ridge regression
estimator (GRRE) of $\delta$ as

$$
\begin{equation*}
\hat{\delta}_{K}=\left[Z^{\prime} Z+K\right]^{-1} Z^{\prime} Y \tag{4.6.1}
\end{equation*}
$$

where $K=\operatorname{diag}\left[k_{1}, k_{2}, \cdots, k_{p}\right], \quad k_{i} \geq 0 \quad \forall i$, and $K$ is known.

### 4.6.1.2 Properties of $\hat{\delta}_{\mathrm{K}}$

1. Relationship to 0LSE

By manipulation, the GRRE of $\delta$ (denoted by $\hat{\delta}_{\mathrm{K}}$ ) can also be written as:

$$
\begin{align*}
\hat{\delta}_{K} & =\left[\Delta^{2}+K\right]^{-1} \Delta^{2} \hat{\delta} \quad(\text { from (3.1.2) and (3.1.3)) } \\
& =\left[I+\Delta^{-2} K\right]^{-1} \hat{\delta} \\
& =\left[I-\left[\Delta^{2}+K\right]^{-1} K\right] \hat{\delta} \tag{4.6.2}
\end{align*}
$$

and the GRRE of $\beta=V \delta$ is

$$
\begin{align*}
\hat{\beta}_{K} & =V \hat{\delta}_{G R} \\
& =V\left[I-\left[\Delta^{2}+K\right]^{-1} K\right] \hat{\delta} \\
& =V\left[I-\left[\Delta^{2}+K\right]^{-1} K\right] V^{\prime} \hat{\beta} \tag{4.6.3}
\end{align*}
$$

When $K=k I$ the GRRE of $\beta$ is the ordinary ridge estimator $\left(\hat{\beta}_{\mathbf{R}}\right)$.
2. Expectation

$$
\begin{align*}
E\left(\hat{\delta}_{K}\right) & =\left[I-\left[\Delta^{2}+K\right]^{-1} K\right] E(\hat{\delta}) \\
& =\left[I-\left[\Delta^{2}+K\right]^{-1} K\right] \delta \tag{4.6.4}
\end{align*}
$$

Thus the bias of $\hat{\delta}_{K}$ is

$$
\begin{equation*}
\mathrm{b}=-\left[\Delta^{2}+\mathrm{K}\right]^{-1} \mathrm{~K} \delta \tag{4.6.5}
\end{equation*}
$$

3. Variance

$$
\begin{align*}
\operatorname{Var}\left(\hat{\delta}_{K}\right) & =\operatorname{Var}\left(\left[\Delta^{2}+K\right]^{-1} \Delta^{2} \hat{\delta}\right) \\
& =\left[\Delta^{2}+\mathbf{K}\right]^{-1} \Delta^{2} \operatorname{Var}(\hat{\delta}) \Delta^{2}\left[\Delta^{2}+K\right]^{-1} \\
& =\sigma^{2}\left[\Delta^{2}+K\right]^{-1} \Delta^{2} \Delta^{-2} \Delta^{2}\left[\Delta^{2}+K\right]^{-1} \\
& =\sigma^{2}\left[\Delta^{2}+K\right]^{-1} \Delta^{2}\left[\Delta^{2}+K\right]^{-1} \tag{4.6.6}
\end{align*}
$$

4. Then by using (4.6.6) and (4.6.5) the MSE of $\hat{\delta}_{\mathrm{K}}$ is:

$$
\begin{align*}
\operatorname{MSE}\left(\hat{\delta}_{K}\right) & =\mathrm{E}\left(\hat{\delta}_{\mathbf{K}}-\delta\right)\left(\hat{\delta}_{\mathrm{K}}-\delta\right)^{\prime} \\
& =\operatorname{Var}\left(\hat{\delta}_{\mathrm{K}}\right)+\mathrm{bb} \mathbf{b}^{\prime} \\
& =\sigma^{2}\left[\Delta^{2}+K\right]^{-2} \Delta^{2}+\left[\Delta^{2}+\mathrm{K}\right]^{-1} \mathrm{~K} \delta \delta^{\prime} \mathrm{K}\left[\Delta^{2}+\mathrm{K}\right]^{-1} \tag{4.6.7}
\end{align*}
$$

5. The total mean square error (TMSE) of $\hat{\delta}_{\mathrm{K}}$ is

$$
\begin{align*}
\operatorname{TMSE}\left(\hat{\delta}_{\mathrm{K}}\right) & =\operatorname{tr}\left(\operatorname{MSE}\left(\hat{\delta}_{\mathrm{K}}\right)\right) \\
& =\sum_{\mathrm{i}=1}^{\mathrm{p}}\left\{\sigma^{2}\left(\lambda_{\mathrm{i}}+\mathrm{k}_{\mathrm{i}}\right)^{-2} \lambda_{\mathrm{i}}+\left[\lambda_{\mathrm{i}}+\mathrm{k}_{\mathrm{i}}\right]^{-2} \mathrm{k}_{\mathrm{i}}^{2} \delta_{\mathrm{i}}^{2}\right\} \\
& =\sum_{\mathrm{i}=1}^{\mathrm{p}} \frac{\sigma^{2} \lambda_{\mathrm{i}}+\mathrm{k}_{\mathrm{i}}^{2} \delta_{\mathrm{i}}^{2}}{\left(\lambda_{\mathrm{i}}+\mathrm{k}_{\mathrm{i}}\right)^{2}} \tag{4.6.8}
\end{align*}
$$

### 4.6.2 K stochastic

### 4.6.2.1 The GRR estimator

In the case of $K$ a stochastic variable, the GRRE of $\delta$ is

$$
\begin{align*}
\hat{\delta}_{\hat{K}} & =\left[\Delta^{2}+\hat{\mathbf{K}}\right]^{-1} \Delta^{2} \hat{\delta} \\
& =\left[I-\left[\Delta^{2}+\hat{K}\right]^{-1} \hat{\mathbf{K}}\right] \hat{\delta} \tag{4.6.9}
\end{align*}
$$

using (4.6.2), where $\hat{K}$ is an estimator of $K$, i.e $\hat{K}=\operatorname{diag}\left[\hat{\mathrm{k}}_{1}, \ldots, \hat{\mathrm{k}}_{\mathrm{p}}\right]$

Then the $\mathbf{i}$ - th element of $\hat{\delta}_{\mathbf{K}}$ is

$$
\begin{equation*}
\left[\hat{\delta}_{\hat{K}}\right]_{i}=\frac{\lambda_{i}}{\lambda_{i}+\hat{k}_{i}} \hat{\delta}_{i} \tag{4.6.10}
\end{equation*}
$$

which shows that $\left[\hat{\delta}_{\hat{K}}\right]_{i}$ is a shrinkage estimator of $\delta_{i}$ where the shrinkage factor ( $\partial_{\mathrm{i}}$ ) is:

$$
\begin{equation*}
\partial_{i}=\frac{\lambda_{i}}{\lambda_{i}+\hat{k}_{i}} \tag{4.6.11}
\end{equation*}
$$

Minimizing (4.6.6) term-by-term to find an optimum value for $k_{i}$ yields

$$
\begin{equation*}
k_{i}(o p t)=\frac{\sigma^{2}}{\delta_{i}^{2}} \quad(i=1,2, \ldots, p) \tag{4.6.12}
\end{equation*}
$$

Hoerl and Kennard (1970b) then estimate $\mathbf{k}_{\mathrm{i}}$ by using the LS estimates of $\sigma^{2}$ and $\delta_{i}$, thus

$$
\begin{equation*}
\hat{\mathrm{k}}_{\mathrm{i}}=\frac{\hat{\sigma}^{2}}{\hat{\delta}_{\mathrm{i}}^{2}} \tag{4.6.13}
\end{equation*}
$$

Inserting these values into (4.6.11) gives

$$
\begin{aligned}
\partial_{i} & =\frac{\lambda_{i}}{\lambda_{i}+\hat{\sigma}^{2} / \hat{\delta}_{i}^{2}} \\
& =\frac{F_{i}}{F_{i}+1}
\end{aligned}
$$

where $F_{i}=\lambda_{i} \hat{\delta}_{\mathrm{i}}^{2} / \hat{\sigma}^{2}$, which is the same as the F -ratio for testing the hypothesis $H_{0}: \delta_{i}=0$. Highly significant values of parameters estimates $\hat{\delta}_{i}$ will be scarcely shrunk at all, but relatively low values of the estimates $\hat{\delta}_{i}$ will be shrunken substantially.

### 4.6.2.2 Properties of $\left[\hat{\delta}_{\hat{\mathbf{R}}}\right]_{\mathrm{i}}$

Let $\tau_{\mathrm{i}}^{2}=\lambda_{\mathrm{i}} \delta_{\mathrm{i}}^{2} / \sigma^{2}$ (i.e. a non-centrality parameter of the F -distributions associated with $\delta_{i}$ ) and $\nu=(n-p)$, then the first and second moments of $\left[\hat{\delta}_{\hat{K}}\right]_{\mathbf{i}}$ are given by (Dwivedi, Srivastava and Hall (1980)):

$$
\begin{align*}
& \mathrm{E}\left(\left[\hat{\delta}_{\hat{\mathbf{K}}}\right]_{\mathrm{i}}\right)=\delta_{\mathrm{i}} \exp \left(-\tau_{\mathbf{i}}^{2} / 2\right) \sum_{\mathrm{a}=0}^{\infty} \sum_{\mathrm{j}=0}^{\infty}\left(\frac{\nu-1}{\nu}\right)^{\mathrm{a}} \times \frac{\left(\tau_{\mathrm{i}}^{2} / 2\right)^{\mathrm{j}}}{\Gamma(\mathrm{j}+1)} \\
& \times \frac{\Gamma(\mathrm{a}+(\nu / 2)) \Gamma(\mathrm{j}+(\nu+3) / 2) \Gamma(\mathrm{j}+(5 / 2))}{\Gamma(\nu / 2) \Gamma(\mathrm{a}+\mathrm{j}+(\nu+5) / 2) \Gamma(\mathrm{j}+(3 / 2))} \tag{4.6.14}
\end{align*}
$$

$$
\begin{align*}
& \mathbb{E}\left(\left[\hat{\delta}_{\hat{\mathbf{K}}}\right]_{\mathrm{i}}^{2}\right)=\delta_{\mathrm{i}}^{2} \exp \left(-\tau_{\mathrm{i}}^{2} / 2\right) \sum_{\mathrm{a}=0}^{\infty} \sum_{\mathrm{j}=0}^{\infty}(\mathrm{a}+1) \times\left(\frac{\nu-1}{\nu}\right)^{\mathrm{a}} \\
& \times \frac{\left(\tau_{\mathrm{i}}^{2} / 2\right)^{\mathrm{j}-1}}{\Gamma(\mathrm{j}+1)} \times \frac{\Gamma(\mathrm{a}+(\nu / 2)) \Gamma(\mathrm{j}+(\nu+3) / 2) \Gamma(\mathrm{j}+(7 / 2))}{\Gamma(\nu / 2) \Gamma(\mathrm{a}+\mathrm{j}+(\nu+7) / 2) \Gamma(\mathrm{j}+(1 / 2))} \tag{4.6.15}
\end{align*}
$$

Using (4.6.14) and (4.6.15) the total squared bias and total mean square error of $\left[\hat{\delta}_{\hat{K}}\right]_{\mathrm{i}}$ can be computed. These authors then define the relative bias (RB) and relative mean square error (RHSE) of $\left[\hat{\delta}_{\hat{K}}\right]_{i}$ as:

$$
\begin{gather*}
\operatorname{RB}\left(\left[\hat{\delta}_{\hat{\mathbf{K}}}\right]_{\mathrm{i}}\right)=\mathrm{E}\left[\left(\left[\hat{\delta}_{\mathbf{K}}\right]_{\mathrm{i}}-[\delta]_{\mathrm{i}}\right) /[\delta]_{\mathrm{i}}\right]  \tag{4.6.16}\\
\operatorname{RMSE}\left(\left[\hat{\delta}_{\hat{\mathbf{K}}}\right]_{\mathrm{i}}\right)=\mathrm{E}\left[\left(\left[\hat{\delta}_{\hat{\mathrm{K}}}\right]_{\mathrm{i}}-[\delta]_{\mathrm{i}}\right) /[\delta]_{\mathrm{i}}\right]^{2} \tag{4.6.17}
\end{gather*}
$$

which are functions of $\nu$ and $\tau_{i}$ only. The efficiency of the LSE of $\delta_{i}$
relative to this ridge estimator is

$$
\begin{align*}
\eta_{\mathrm{i}} & =\left(\operatorname{MSE}\left(\left[\hat{\delta}_{\hat{\mathrm{K}}}\right]_{\mathrm{i}}\right) / \operatorname{Var}[\hat{\delta}]_{\mathrm{i}}\right) \times 100  \tag{4.6.18}\\
& =100 \tau_{\mathrm{i}}^{2} \operatorname{RMSE}\left(\left[\hat{\delta}_{\hat{\mathrm{K}}}\right]_{\mathrm{i}}\right)
\end{align*}
$$

By selecting different values for $\tau_{\mathrm{i}}^{2}$ (between 0.02 and 50 ) and $\nu$ (between 1 and 100) the authors concluded:

1. The ridge estimator ( $\left[\hat{\delta}_{\hat{K}}\right]_{\mathrm{i}}$ ) is biased in a direction opposite to the sign of the coefficient (i.e. shrunk toward zero).
2. The RB is a decreasing function of $\tau_{\mathrm{i}}^{2}$ and and increasing function of $\nu$ (the degrees of freedom for error).
3. RMSE decreases as $\tau_{\mathrm{i}}^{2}$ increases.
4. Changes in RMSE are more rapid for small values of $\tau_{i}^{2}$ and $\nu$.
5. The ridge estimator is more efficient than the LSE as long as $\tau_{\mathrm{i}}^{2}<2$. (We suggest that to operationalise this criterion one estimates $\tau_{i}^{2}$ by substitution of the sample estimators. Further research is planned to investigate behaviour of estimators chosen on this basis.)

### 4.6.3 Choice of $k$ in $K=k I$

Hoerl, Kennard and Baldwin (1975) give an algorithm for the selection of $k$ (in ordinary ridge estimation). They distinguish between two cases: First if $X^{\prime} X=I$ then a minimum TMSE is obtained if $k=p \sigma^{2} / \beta^{\prime} \beta$. Secondly, in the general form (as discussed in §4.6.1 and §4.6.2), a minimum MSE is obtained when $\mathrm{k}_{\mathrm{i}}=\sigma^{2} / \delta_{\mathrm{i}}^{2}$ as in (4.6.12). These individual $\mathrm{k}_{\mathrm{i}}$ are combined to form a single value for $k$. Large values of $\delta_{i}$ will blow up the ordinary mean of $k_{i}$, therefore they (HKB) suggested the harmonic mean ( $k_{h}$ ) of the
$k_{i}$ 's. The two methods give the same value of $k$ since

$$
\begin{aligned}
1 /\left(\mathrm{k}_{\mathrm{h}}\right) & =\frac{1}{\mathrm{p}} \sum_{\mathrm{i}=1}^{\mathrm{p}}\left(1 / \mathrm{k}_{\mathrm{i}}\right) \\
& =\frac{1}{\mathrm{p}} \sum_{\mathrm{i}=1}^{\mathrm{p}}\left(\delta_{\mathrm{i}}^{2} / \sigma^{2}\right) \\
& =\left(1 / \mathrm{p} \sigma^{2}\right) \delta^{\prime} \delta \\
& =\left(1 / \mathrm{p} \sigma^{2}\right) \beta^{\prime} \mathrm{V} V^{\prime} \beta \\
\therefore \mathrm{k}_{\mathrm{h}} & =\mathrm{p} \sigma^{2} / \beta^{\prime} \beta
\end{aligned}
$$

and

$$
\begin{equation*}
\hat{\mathrm{k}}_{\mathrm{h}}=\mathrm{p} \hat{\sigma}^{2} / \hat{\beta}^{\prime} \hat{\beta} \tag{4.6.19}
\end{equation*}
$$

Substituting $\hat{k}$ from (4.6.19) in the TMSE (4.6.8) give the estimated minimum TMSE as

$$
\begin{aligned}
\operatorname{TMSE}\left(\hat{\delta}_{K}\right) & =\sum_{i=1}^{p} \frac{\hat{\sigma}^{2} \lambda_{i}+\left(p \hat{\sigma}^{2} / \hat{\delta}^{\prime} \hat{\delta}\right)^{2} \hat{\delta}_{i}^{2}}{\left(\lambda_{i}+\left(p \hat{\sigma}^{2} / \hat{\delta}^{\prime} \hat{\delta}\right)\right)^{2}} \\
& =\sum_{i=1}^{p} \frac{\hat{\sigma}^{2} \lambda_{i}\left(\hat{\delta}^{\prime} \hat{\delta}\right)^{2}+\left(p \hat{\sigma}^{2}\right)^{2} \hat{\delta}_{i}^{2}}{\left(\hat{\delta}^{\prime} \hat{\delta} \lambda_{i}+p \hat{\sigma}^{2}\right)^{2}}
\end{aligned}
$$

By using simulation on three data sets they show that the algorithm (discussed in the previous paragraph for selecting $k$ ) has the following properties (p111):

1. The use of the ridge estimator with biasing parameter $\hat{\mathrm{k}}_{\mathrm{h}}=\mathrm{p} \hat{\sigma}^{2} / \hat{\beta}^{\prime} \hat{\beta}$ has a probability greater than 0.5 of producing estimates with a smaller MSE than that of OLSE.
2. The probability of a smaller TMSE (than that of LSE) using $\hat{\mathrm{k}}_{\mathrm{h}}$ increases as $p$ (the number of independent variables) increases.
3. The probability of a smaller THSE using $\hat{\mathrm{k}}_{\mathrm{h}}$ increases as X ' X becomes less well- conditioned.
4. The probability of a smaller TMSE using $\hat{\mathrm{k}}_{\mathrm{h}}$ increases as the noise (measured by $\hat{\sigma}^{2}$ ) increases.

They concluded that every ridge trace should have the point corresponding to $\mathrm{k}=\hat{\mathrm{k}}_{\mathrm{h}}$ computed $\operatorname{explicitly}$.

### 4.7 Jackknifed Ridge Estimator

One disadvantage of the ridge estimator is that it may have a serious bias. To reduce this bias Singh, Chaubey and Dwivedi (1986) use the jackknife procedure (as described in §1.8.2 for OLSE) on the generalized ridge estimator to construct a jackknifed ridge estimator (JRE). Some authors (e.g. Nomura (1988)) call this JRE the almost unbiased generalized ridge regression (AUGRR) estimator. This estimator is similar in form to the ridge estimator and has the same asymptotic properties. For the moment $K$ is fixed (non-stochastic).

### 4.7.1 The JR estimator

Let $Y_{-i}$ and $Z_{-i}^{\prime}$ denote respectively the vector $Y$ with its i-th coordinate deleted and the matrix $Z^{\prime}$ with its i-th column deleted, and let $\left(\hat{\delta}_{K}\right)_{-i}$ be given by (4.6.1) with $Z^{\prime}$ and $Y$ replaced by $Z_{-i}^{\prime}$ and $Y_{-i}$ respectively. Clearly $Z_{-i}$ is the matrix $Z$ with its i-th row deleted.

$$
\begin{equation*}
\left(\hat{\delta}_{\mathrm{K}}\right)_{-\mathrm{i}}=\left[\mathrm{Z}_{-\mathrm{i}}^{\prime} \mathrm{Z}_{-\mathrm{i}}+\mathrm{K}\right]^{-1} \mathrm{Z}_{-\mathrm{i}}^{\prime} \mathrm{Y}_{-\mathrm{i}} \tag{4.7.1}
\end{equation*}
$$

$$
\text { Note that } \begin{align*}
Z^{\prime} Z & =\left[\begin{array}{ll}
Z_{-i}^{\prime} & z_{i}
\end{array}\right]\left[\begin{array}{c}
Z_{-i} \\
z_{i}^{i}
\end{array}\right] \\
& =\left[\begin{array}{ll}
Z_{-i}^{\prime} & Z_{-i}+Z_{i} z_{i}^{\prime}
\end{array}\right] \tag{4.7.2}
\end{align*}
$$

$$
\text { and } \begin{align*}
Z^{\prime} Y & =\left[\begin{array}{ll}
Z_{-i}^{\prime} & Z_{i}
\end{array}\right]\left[\begin{array}{l}
Y_{-i} \\
y_{i}
\end{array}\right] \\
& =\left[\begin{array}{ll}
Z_{-i}^{\prime} Y_{-i} & +z_{i} y_{i}
\end{array}\right] \tag{4.7.3}
\end{align*}
$$

Define $A=Z^{\prime} Z+K, \quad W_{i}=z_{i}^{\prime} A^{-i} z_{i}$ (corresponding to $h_{i}$ in the $0 L S$ estimation) and $u_{i}=y_{i}-z_{i}^{\prime} \hat{\delta}_{k}$ (corresponding to the 0LS $\hat{\epsilon}_{i}$ ). Then (4.7.1) can be be written as

$$
\begin{align*}
\left(\hat{\delta}_{K}\right)_{-i} & =\left[Z^{\prime} Z-z_{i} z_{i}^{\prime}+K\right]^{-1}\left[Z^{\prime} Y-z_{i} y_{i}\right] \\
& =\left[A-z_{i} z_{i}\right]^{-1}\left[Z^{\prime} Y-z_{i} y_{i}\right] \\
& =\left[A^{-1}+A^{-1} z_{i}\left(1-z_{i} A^{-1} z_{i}\right)^{-1} z_{i} A^{-1}\right]\left[Z^{\prime} Y-z_{i} y_{i}\right] \\
& =A^{-1} Z^{\prime} Y-A^{-1} z_{i} y_{i}+A^{-1} z_{i}\left(1-W_{i}\right)^{-1} z_{i} A^{-1}\left[Z^{\prime} Y-z_{i} y_{i}\right] \\
& =\left(\hat{\delta}_{K}\right)-A^{-1} z_{i}\left(1-W_{i}\right)^{-1}\left[y_{i}\left(1-H_{i}\right)-z_{i} A^{-1}\left[Z^{\prime} Y-z_{i} y_{i}\right]\right] \\
& (\operatorname{from}(4.6 .1)) \\
& =\left(\hat{\delta}_{K}\right)-A^{-1} z_{i}\left(1-W_{i}\right)^{-1}\left[y_{i}-y_{i} H_{i}-z_{i}^{\prime} A^{-1} Z^{\prime} Y+W_{i} y_{i}\right] \\
& =\left(\hat{\delta}_{K}\right)-A^{-1} z_{i}\left(1-W_{i}\right)^{-1}\left[y_{i}-z_{i}^{\prime}\left(\hat{\delta}_{K}\right)\right] \\
& =\hat{\delta}_{K}-\frac{A^{-1} z_{i}}{\left(1-W_{i}\right.}\left[u_{i}\right] \tag{4.7.4}
\end{align*}
$$

In the same manner as (1.8.8) the pseudovalues are defined as

$$
\begin{equation*}
P_{i}=n \hat{\delta}_{K}-(n-1)\left(\hat{\delta}_{K}\right)_{-i} \tag{4.7.5}
\end{equation*}
$$

and the jackknife ridge estimator (JRE) is

$$
\begin{align*}
\hat{\delta}_{J} & =\frac{1}{n} \Sigma P_{i} \\
& =\frac{1}{n} \Sigma\left(n \hat{\delta}_{K}-(n-1)\left(\hat{\delta}_{K}\right)_{-i}\right) \\
& \left.=\hat{\delta}_{K}+\frac{n-1}{n} A^{-1} \sum_{i=1}^{n} \frac{z_{i} u_{i}}{\left(1-w_{i}\right.}\right) \tag{4.7.6}
\end{align*}
$$

Singh, Chaubey and Drivedi (1986) reject these pseudo-values ( $\mathrm{P}_{\mathrm{i}}$ ), as they are defined symmetrically with respect to the observations, while the
regression models are often unbalanced with the lack of balance reflecting through the 'squared distance' $w_{i}$. Further, since the variance of the differences $\left(\hat{\delta}_{K}-\left(\hat{\delta}_{K}\right)_{-i}\right)=A^{-1} \frac{z_{i} u_{i}}{\left(1-W_{i}\right)}$ is an increasing function of $w_{i}$, they define the pseudo-values in the manner of (1.8.13) as

$$
\begin{equation*}
Q_{i}=\hat{\delta}_{K}+n\left(1-W_{i}\right)\left(\hat{\delta}_{K}-\left(\hat{\delta}_{K}\right)_{-i}\right) \tag{4.7.7}
\end{equation*}
$$

Then the weighted jackknife estimator (denoted by $\hat{\delta}_{J W}$ ) is

$$
\begin{align*}
\hat{\delta}_{J W} & =\frac{1}{n} \Sigma Q_{i} \\
& =\hat{\delta}_{K}+A^{-1} \Sigma z_{i} u_{i} \quad \text { (using (4.7.4) } \\
& =\hat{\delta}_{K}+A^{-1} \Sigma z_{i}\left(y_{i}-z_{i}^{\prime} \hat{\delta}_{K}\right) \\
& =\hat{\delta}_{K}+A^{-1}\left[\Sigma z_{i} y_{i}-\Sigma z_{L_{i}} z_{i}^{\prime} \hat{\delta}_{K}\right] \\
& =\hat{\delta}_{K}+A^{-1}\left[Z^{\prime} Y-Z^{\prime} Z \hat{\delta}_{K}\right] \\
& =\hat{\delta}_{K}+A^{-1}\left[A \hat{\delta}_{K}-Z^{\prime} Z \hat{\delta}_{K}\right] \quad\left(\hat{\delta}_{K}=A^{-1} Z^{\prime} Y\right) \\
& =\hat{\delta}_{K}+A^{-1} K \hat{\delta}_{K} \quad\left(A=Z^{\prime} Z+K\right) \tag{4.7.8}
\end{align*}
$$

When $\mathrm{k}=\mathrm{k}_{1}=\ldots=\mathrm{k}_{\mathrm{p}}$ the AUGRR estimator is called the almost unbiased ordinary ridge regression (AUORR) estimator (see (4.7.21) for an estimate of $k$ ).

### 4.7.2 Properties of $\hat{\delta}_{\mathrm{JW}}$ (non- stochastic K)

1. Relationship to OLSE

The weighted JRE can also be written as

$$
\begin{align*}
\hat{\delta}_{J W} & =\left[I+A^{-1} K\right]\left[I-A^{-1} K\right] \hat{\delta} \quad(\text { from }(4.6 .2)) \\
& =\left[I-\left[A^{-1} K\right]^{2}\right] \hat{\delta} \tag{4.7.9}
\end{align*}
$$

and the i-th element of $\hat{\delta}_{J W}$ is:

$$
\begin{equation*}
\left[\hat{\delta}_{\mathrm{JW}}\right]_{\mathrm{i}}=\left[1-\frac{\mathrm{k}_{\mathrm{i}}^{2}}{\left(\lambda_{\mathrm{i}}+\mathrm{k}_{\mathrm{i}}\right)^{2}}\right] \hat{\delta}_{\mathrm{i}} \tag{4.7.10}
\end{equation*}
$$

## 2. Expectation

$$
\begin{align*}
E\left(\hat{\delta}_{J W}\right) & =\left[I-\left[A^{-1} K\right]^{2}\right] \delta \\
\operatorname{Bias}\left(\hat{\delta}_{J W}\right) & =-\left[A^{-1} K\right]^{2} \delta \tag{4.7.11}
\end{align*}
$$

and the bias of the i-th element is:

$$
\begin{equation*}
\operatorname{Bias}\left[\left(\hat{\delta}_{\mathrm{JW}}\right)_{\mathrm{i}}\right]=-\left[\frac{\mathrm{k}_{\mathrm{i}}^{2}}{\left(\lambda_{\mathrm{i}}+\mathrm{k}_{\mathrm{i}}\right)^{2}}\right] \delta_{\mathrm{i}} \tag{4.7.12}
\end{equation*}
$$

Note that the JRE $\left(\hat{\delta}_{\mathbf{J W}}\right)$ has smaller bias than that of the GRRE, where bias of the i-th element $\left[\left(\hat{\delta}_{K}\right)\right]_{i}$ is $-\left[\frac{k_{i}}{\left(\lambda_{i}+k_{i}\right)}\right] \delta_{i}$, thus

$$
\begin{equation*}
\left|\operatorname{Bias}\left[\left(\hat{\delta}_{\mathrm{K}}\right)\right]_{\mathrm{i}}\right|-\left|\operatorname{Bias}\left[\left(\hat{\delta}_{\mathrm{JW}}\right)_{\mathrm{i}}\right]\right|=\left[\frac{\lambda_{\mathrm{i}} \mathrm{k}_{\mathrm{i}}}{\left(\lambda_{\mathrm{i}}+\mathrm{k}_{\mathrm{i}}\right)^{2}}\right]\left|\delta_{\mathrm{i}}\right|>0 \tag{4.7.13}
\end{equation*}
$$

3. Variance

$$
\begin{aligned}
\operatorname{Var}\left(\left[\left(\hat{\delta}_{J W}\right)\right]\right. & =\left[I-\left[A^{-1} K\right]^{2}\right] \operatorname{Var}(\hat{\delta})\left[I-\left[A^{-1} K\right]^{2}\right] \\
& \left.=\left[I-\left[A^{-1} K\right]^{2}\right] \sigma^{2} \Delta^{-2}\left[I-\left[A^{-1} K\right]^{2}\right]\right)
\end{aligned}
$$

and the variance of the $i$ - th component is:

$$
\begin{align*}
\operatorname{Var}\left[\left(\hat{\delta}_{\mathrm{JW}}\right)_{\mathrm{i}}\right] & =\left[1-\frac{\mathrm{k}_{\mathrm{i}}^{2}}{\left(\lambda_{\mathrm{i}}+\mathrm{k}_{\mathrm{i}}\right)^{2}}\right]^{2} \frac{\sigma^{2}}{\lambda_{\mathrm{i}}} \\
& =\sigma^{2}\left[\frac{\lambda_{\mathrm{i}}\left(\lambda_{\mathrm{i}}+2 \mathrm{k}_{\mathrm{i}}\right)^{2}}{\left(\lambda_{\mathrm{i}}+\mathrm{k}_{\mathrm{i}}\right)^{4}}\right] \tag{4.7.14}
\end{align*}
$$

4. By using (4.7.14) and (4.2.12) the MSE of the components of $\hat{\delta}_{J W}$ is

$$
\begin{align*}
\operatorname{MSE}\left[\left(\hat{\delta}_{\mathrm{JW}}\right)_{\mathrm{i}}\right] & =\operatorname{Var}\left[\left(\hat{\delta}_{\mathrm{JW}}\right)_{\mathrm{i}}\right]+\operatorname{bias}^{2} \\
& =\sigma^{2}\left[\frac{\lambda_{\mathrm{i}}\left(\lambda_{\mathrm{i}}+2 \mathrm{k}_{\mathrm{i}}\right)^{2}}{\left(\lambda_{\mathrm{i}}+\mathrm{k}_{\mathrm{i}}\right)^{4}}\right]+\left[\frac{\mathrm{k}_{\mathrm{i}}^{4}}{\left(\lambda_{\mathrm{i}}+\mathrm{k}_{\mathrm{i}}\right)^{4}}\right] \delta_{\mathrm{i}}^{2} \tag{4.7.15}
\end{align*}
$$

5. The total mean square error (TMSE) of $\hat{\delta}_{J W}$ is obtained from (4.7.15) and the definition of TMSE (1.8.2):

$$
\begin{equation*}
\operatorname{TMSE}\left(\hat{\delta}_{J W}\right)=\sum_{i=1}^{p}\left\{\sigma^{2} \frac{\lambda_{i}\left(\lambda_{i}+2 k_{i}\right)^{2}}{\left(\lambda_{i}+k_{i}\right)^{4}}+\frac{k_{i}^{4}}{\left(\lambda_{i}+k_{i}\right)^{4}} \delta_{i}^{2}\right\} \tag{4.7.16}
\end{equation*}
$$

### 4.7.3 Additional comments on ( $\hat{\delta}_{\mathrm{JW}}$ ) (non-stochastic K )

In comparing component-wise the JRE (or AUGRR estimator) with the GRRE and the OLSE in terms of the MSE, the following theorems were developed by Nomura (1988). The comparison of the various estimators will be discussed further in Chapter 8.

From (4.6.7) and (4.6.15) we have

$$
\begin{aligned}
\operatorname{MSE}\left(\hat{\delta}_{\mathrm{i}}\right) & =\sigma^{2} / \lambda_{\mathrm{i}}, \quad \mathrm{i}=1, \ldots, \mathrm{p} \\
\operatorname{MSE}\left[\left(\hat{\delta}_{\mathrm{K}}\right)_{\mathrm{i}}\right] & =\left[\sigma^{2} \lambda_{\mathrm{i}}+\delta_{\mathrm{i}}^{2} \mathrm{k}_{\mathrm{i}}^{2}\right] /\left(\lambda_{\mathrm{i}}+\mathrm{k}_{\mathrm{i}}\right)^{2} \\
\operatorname{MSE}\left[\left(\hat{\delta}_{\mathrm{JW}}\right)_{\mathrm{i}}\right] & =\left\{\sigma^{2} \lambda_{\mathrm{i}}^{3}+4 \sigma^{2} \lambda_{\mathrm{i}}^{2} \mathrm{k}_{\mathrm{i}}+4 \sigma^{2} \lambda_{\mathrm{i}} \mathrm{k}_{\mathrm{i}}^{2}+\delta_{\mathrm{i}}^{2} \mathrm{k}_{\mathrm{i}}^{4}\right\} /\left(\lambda_{\mathrm{i}}+\mathrm{k}_{\mathrm{i}}\right)^{4}
\end{aligned}
$$

## Theorem 4.7.3.1

1. $\operatorname{MSE}\left[\left(\hat{\delta}_{\mathrm{JW}}\right)_{\mathrm{i}}\right]>\operatorname{MSE}\left[\left(\hat{\delta}_{\mathrm{K}}\right)_{\mathrm{i}}\right]$ for $0<\mathrm{k}_{\mathrm{i}}<\mathrm{K} 1_{\mathrm{i}}$,
2. $\operatorname{MSE}\left[\left(\hat{\delta}_{\mathrm{JW}}\right)_{\mathrm{i}}\right]<\operatorname{MSE}\left[\left(\hat{\delta}_{\mathrm{K}}\right)_{\mathrm{i}}\right]$ for $\mathrm{K}_{\mathrm{i}}<\mathrm{k}_{\mathrm{i}}<\infty$,
where $\mathrm{K} 1_{\mathrm{i}}=\left[3 \sigma^{2}-\lambda_{\mathrm{i}} \delta_{\mathrm{i}}^{2}+\left\{\left(3 \sigma^{2}+\lambda_{\mathrm{i}} \delta_{\mathrm{i}}^{2}\right)^{2}+4 \sigma^{2} \lambda_{\mathrm{i}} \delta_{\mathrm{i}}^{2}\right\}^{1 / 2}\right] / 4 \delta_{\mathrm{i}}^{2}>0$.
Proof is given in Nomura (1988, p732).

## Theorem 4.7.3.2

1. If $\lambda_{i} \delta_{i}^{2}-\sigma^{2} \leq 0$, then

$$
\operatorname{MSE}\left[\left(\hat{\delta}_{\mathrm{J} W}\right)_{\mathrm{i}}\right]<\operatorname{MSE}\left[(\hat{\delta})_{\mathrm{i}}\right] \text { for } 0<\mathrm{k}_{\mathrm{i}}<\omega
$$

2. If $\lambda_{\mathrm{i}} \delta_{\mathrm{i}}^{2}-\sigma^{2} \leq 0$, then there exists a strictly positive constant $K 2_{\mathrm{i}}=\left[2 \sigma^{2} \lambda_{\mathrm{i}}+\left(2 \sigma^{4} \lambda_{\mathrm{i}}^{2}+2 \sigma^{2} \lambda_{\mathrm{i}}^{3} \delta_{\mathrm{i}}^{2}\right)^{1 / 2}\right] /\left(\lambda_{\mathrm{i}} \delta_{\mathrm{i}}^{2}-\sigma^{2}>0\right.$ such that $\operatorname{MSE}\left[\left(\hat{\delta}_{J W}\right)_{\mathrm{i}}\right]<\operatorname{MSE}\left[(\hat{\delta})_{\mathrm{i}}\right]$ for $0<\mathrm{k}_{\mathrm{i}}<\mathrm{K} 2_{\mathrm{i}}$ and $\operatorname{MSE}\left[\left(\hat{\delta}_{\mathrm{JW}}\right)_{\mathrm{i}}\right]>\operatorname{MSE}\left[(\hat{\delta})_{\mathrm{i}}\right]$ for $\mathrm{K} 2_{\mathrm{i}}<\mathrm{k}_{\mathrm{i}}<\infty$.

Proof is given in Nomura (1988, p733)

Furthermore, differentiating the $\operatorname{MSE}\left[\left(\hat{\delta}_{J W}\right)_{i}\right]$ with respect to $k_{i}$ yields:

$$
\begin{equation*}
\frac{\operatorname{dMSE}\left[\left(\hat{\delta}_{\mathrm{JW}}\right)_{\mathrm{i}}\right]}{d \mathrm{k}_{\mathrm{i}}}=\frac{4 \lambda_{\mathrm{i}} \mathrm{k}_{\mathrm{i}}\left(\delta_{\mathrm{i}}^{2} \mathrm{k}_{\mathrm{i}}^{2}-2 \sigma^{2} \mathrm{k}_{\mathrm{i}}-\sigma^{2} \lambda_{\mathrm{i}}\right)}{\left(\lambda_{\mathrm{i}}+\mathrm{k}_{\mathrm{i}}\right)^{5}} \tag{4.7.17}
\end{equation*}
$$

thus the optimal value of $k_{i}$ for which $\operatorname{MSE}\left[\left(\hat{\delta}_{J W}\right)_{i}\right]$ is a minimum is:

$$
\begin{equation*}
\mathrm{k}_{\mathrm{i}}(\mathrm{opt})=\frac{\sigma^{2}+\sqrt{\sigma^{4}+\sigma^{2} \lambda_{\mathrm{i}} \delta_{\mathrm{i}}^{2}}}{\delta_{\mathrm{i}}^{2}} \tag{4.7.18}
\end{equation*}
$$

If $\lambda_{i} \delta_{i}^{2}-\sigma^{2}>0$ then

$$
\begin{equation*}
\mathrm{K}_{1}{ }_{\mathrm{i}}<\mathrm{k}_{\mathrm{i}}(\mathrm{opt})<\mathrm{K} 2_{\mathrm{i}} \tag{4.7.19}
\end{equation*}
$$

We observe that this is a condition that depends on the data source not the data and is therefore not operationable.

In the case of AUORR defined below (4.7.8), a single value of $k$ can be formed by combining all the $\mathrm{k}_{\mathrm{i}}$ through the harmonic mean (compare this with §4.6.3). Thus

$$
\begin{equation*}
\mathrm{k}_{\mathrm{h}}=\mathrm{p} \sigma^{2} /\left(\sum_{\mathrm{i}=1}^{\mathrm{p}}\left[\delta_{\mathrm{i}}^{2} /\left\{1+\sqrt{1+\lambda_{\mathrm{i}}\left(\delta_{\mathrm{i}}^{2} / \sigma^{2}\right)}\right\}\right]\right) \tag{4.7.20}
\end{equation*}
$$

### 4.7.4 Additional comments on ( $\hat{\delta}_{\mathrm{JW}}$ ) (stochastic K)

In equation (4.7.20) an estimate of $\mathrm{k}_{\mathrm{h}}$ can be obtained from estimating $\delta$ and $\sigma$ by their LSE's

$$
\begin{equation*}
\hat{\mathrm{k}}_{\mathrm{h}}=\mathrm{p} \hat{\sigma}^{2} /\left(\mathrm{\Sigma}^{\mathrm{L}}\left[\hat{\delta}_{\mathrm{i}}^{2} /\left\{1+\sqrt{1+\lambda_{\mathrm{i}}\left(\hat{\delta}_{\mathrm{i}}^{2} / \hat{\sigma}^{2}\right)}\right\}\right]\right) \tag{4.7.21}
\end{equation*}
$$

where $\hat{\sigma}^{2}=(Y-Z \hat{\delta})^{\prime}(Y-Z \hat{\delta}) /(n-p)$.

Nomura (1988) performed a simulation study on six estimators namely:

1. OLSE
2. AUORR estimator using the operational $\hat{\mathrm{k}}_{\mathrm{h}}$ of (4.7.21)
3. Ordinary RR estimator using the Hoerl, Kennard and Baldwin ridge parameter (4.6.19)
4. AUORR estimator using the operational $\hat{\mathrm{k}}_{\mathrm{h}}$ of (4.6.19)
5. Ordinary RR estimator using the $\hat{k}$ as suggested by Lawless and Wang (1976) $\hat{k}_{L W}=p \hat{\sigma}^{2} / \sum_{i=1}^{p} \lambda_{i} \hat{\delta}_{i}^{2}$
6. AUORR estimator using the $\hat{\mathrm{k}}_{\mathrm{LW}}$.

For a description of the simulation, the data and the results see pp735-742 of his article. He reported the following findings for different estimators referenced 1 to 6 as above.
(1) The absolute value of the bias for 2 and 4 is smaller than that of 3 and 5. This result coincides with (4.7.13) for $k$ non- stochastic.
(2) The estimators 2 and 4 are more efficient in terms of the MSE and predictive MSE (defined as the $E(\tilde{\delta}-\delta)^{\prime} X^{\prime} X(\tilde{\delta}-\delta)$ ) than 3 and 4 when the signal to noise ratio is relatively large, and more efficient than 1 (0LSE) when the signal to noise ratio is relatively small.
(3) The estimator 2 performs very well against 4.

We note that these findings are applicable to the data set and models used by Nomura and are thus only suggestive of what one might expect.

The AUGRR discussed so far was due to Singh, Chaubey and Dwivedi (1986) and was not operationable. Ohtani (1986) considered the operational AUGRR estimator based on the idea of Kadiyala (1984). Following Kadiyala (1984) the almost unbiased generalized ridge estimator (non-operationable) for the i-th element of $\delta$ is:

$$
\begin{equation*}
\left[\delta_{0}\right]_{\mathrm{i}}=\left[\delta_{\mathrm{K}}\right]_{\mathrm{i}}+\mathrm{k}_{\mathrm{i}}\left[\delta_{\mathrm{K}}\right]_{\mathrm{i}} /\left(\lambda_{\mathrm{i}}+\mathrm{k}_{\mathrm{i}}\right) \tag{4.7.22}
\end{equation*}
$$

In the above expression if $\mathrm{k}_{\mathrm{i}}$ and $\left[\delta_{\mathrm{K}}\right]_{\mathrm{i}}$ are replaced by $\hat{\mathrm{k}}_{\mathrm{i}}$ and $\left[\hat{\delta}_{\mathrm{K}}\right]_{\mathrm{i}}$ we obtain the following operational AUGRR estimator:

$$
\begin{align*}
{\left[\hat{\delta}_{\hat{0}}\right]_{\mathrm{i}} } & =\left(1+\hat{\mathbf{k}}_{\mathrm{i}} /\left(\lambda_{\mathrm{i}}+\hat{\mathbf{k}}_{\mathrm{i}}\right)\right)\left[\hat{\delta}_{\mathrm{K}}\right]_{\mathrm{i}} \\
& =\left(\left(\lambda_{\mathrm{i}}+2 \hat{\mathrm{k}}_{\mathrm{i}}\right) /\left(\lambda_{\mathrm{i}}+\hat{\mathrm{k}}_{\mathrm{i}}\right)\right)\left[\lambda_{\mathrm{i}} \hat{\delta}_{\mathrm{i}} /\left(\lambda_{\mathrm{i}}+\hat{\mathrm{k}}_{\mathrm{i}}\right)\right] \quad(\text { from }(4.6 .10)) \\
& =\frac{\left(\lambda_{\mathrm{i}} \hat{\delta}_{\mathrm{i}}^{2}+2 \hat{\sigma}^{2}\right) \lambda_{\mathrm{i}} \hat{\delta}_{\mathrm{i}}^{3}}{\left(\lambda_{\mathrm{i}} \hat{\delta}_{\mathrm{i}}^{2}+\hat{\sigma}^{2}\right)^{2}} \quad(\text { from }(4.6 .13)) \tag{4.7.23}
\end{align*}
$$

Let $\tau_{\mathrm{i}}^{2}=\lambda_{\mathrm{i}} \delta_{\mathrm{i}}^{2} / \sigma^{2}$ and $\nu=\mathrm{n}-\mathrm{p}$, then the first and second moments of $\left[\hat{\delta}_{\hat{0}}\right]_{\mathrm{i}}$ are

$$
\begin{align*}
\mathrm{E}\left(\left[\hat{\delta}_{\hat{0}}\right]_{\mathrm{i}}\right)=[ & \delta_{\mathrm{i}} \\
& \exp \left(-\tau_{\mathrm{i}}^{2} / 2\right) /\{\sqrt{\pi} \Gamma(\nu / 2) \\
& \times \sum_{\mathrm{g}=0}^{\infty}\left[2^{\mathrm{g}+1}\left(\tau_{\mathrm{i}}^{2}\right)^{\mathrm{g}} \Gamma(\mathrm{~g}+(\nu+3) / 2) /(2 \mathrm{~g}+1)!\right. \\
& \times \int_{0}^{1}\left[(\mathrm{t}+2(1-\mathrm{t}) / \nu) /(\mathrm{t}+(1-\mathrm{t}) / \nu)^{2}\right]  \tag{4.7.24}\\
& \times \mathrm{t}^{\mathrm{g}+3 / 2}(1-\mathrm{t})^{\nu / 2-1} \mathrm{dt}
\end{align*}
$$

$$
\begin{align*}
\mathrm{E}\left(\left[\hat{\delta}_{\hat{0}}^{\perp}\right]_{\mathrm{i}}^{2}\right)=[ & \delta_{\mathrm{i}}^{2} \exp \left(-\tau_{\mathrm{i}}^{2} / 2\right) /\{\sqrt{\pi} \Gamma(\nu / 2) \\
& \times \sum_{\mathrm{g}=0}^{\infty}\left[2^{\mathrm{g}+1}\left(\tau_{\mathrm{i}}^{2}\right)^{\mathrm{g}-1} \Gamma(\mathrm{~g}+(\nu+3) / 2) /(2 \mathrm{~g})!\right] \\
& \times \int_{0}^{1}\left[(\mathrm{t}+2(1-\mathrm{t}) / \nu)^{2} /(\mathrm{t}+(1-\mathrm{t}) / \nu)^{4}\right] \\
& \times \mathrm{t}^{\mathrm{g}+5} / 2(1-\mathrm{t})^{\nu / 2-1} \mathrm{dt} \tag{4.7.25}
\end{align*}
$$

Then the relative bias (RB) and relative mean square error (RMSE) of $\left[\hat{\delta}_{\hat{0}}\right]_{\mathrm{i}}^{2}$ as defined in (4.6.16) and (4.6.17) are

$$
\begin{gather*}
\operatorname{RB}\left(\left[\hat{\delta}_{0}\right]_{\mathrm{i}}^{2}\right)=\mathrm{E}\left[\left(\left[\hat{\delta}_{0}\right]_{\mathrm{i}}-[\hat{\delta}]_{\mathrm{i}}\right) /[\hat{\delta}]_{\mathrm{i}}\right]  \tag{4.7.26}\\
\operatorname{RMSE}\left(\left[\hat{\delta}_{0}\right]_{\mathrm{i}}^{2}\right)=\mathrm{E}\left[\left(\left[\hat{\delta}_{0}\right]_{\mathrm{i}}-[\hat{\delta}]_{\mathrm{i}}\right)^{2} /[\hat{\delta}]_{\mathrm{i}}^{2}\right] \tag{4.7.27}
\end{gather*}
$$

which are functions of $\tau_{i}^{2}$ and $\nu$ only. Ohtani (1986) then defined the relative efficiencies of the AUGRR estimator and the GRRE to the OLSE as:

$$
\begin{align*}
& \operatorname{REFF}\left(\left[\hat{\delta}_{\hat{0}}\right]_{\mathrm{i}}\right)=\operatorname{MSE}\left([\hat{\delta}]_{\mathrm{i}}\right) / \operatorname{MSE}\left(\left[\hat{\delta}_{\hat{0}}\right]_{\mathrm{i}}\right)  \tag{4.7.28}\\
& \operatorname{REFF}\left(\left[\hat{\delta}_{\hat{\mathrm{K}}}\right]_{\mathrm{i}}\right)=\operatorname{MSE}\left([\hat{\delta}]_{\mathrm{i}}\right) / \operatorname{MSE}\left(\left[\hat{\delta}_{\hat{\mathrm{K}}}\right]_{\mathrm{i}}\right) \tag{4.7.29}
\end{align*}
$$

In (4.6.18) efficiency of OLSE to GRRE was defined as the inverse of (4.7.29) and presented as a percentage.

By selecting different values for $\tau_{\mathrm{i}}^{2}$ (between 0.25 and 0.40 ) and $\nu$ (between 10 and 60) the author concluded:

1. The reduction in bias is greater for AUGRR estimator than for GRRE but the relative efficiency of the AUGRR estimator is less than the GRRE.
2. The AUGRR estimator is more efficient than the OLSE as long as the value of $\tau_{\mathrm{i}}^{2}<1.5$ (For GRRE the corresponding value was 2 ).

Ohtani (1986) concluded that the AUGRR estimation is rather inferior to the GRRE.

### 4.8 Summary

In this chapter we introduced ridge, generalized ridge and the jackknifed ridge estimators. The estimators were defined, and their properties discussed. Methods of estimating $k$ of K were introduced.

## Chapter 5

## fractional principals in estimation

In this chapter the principle of fractional estimation using the augmented matrix and in terms of shrinkage estimators, will be discussed.

### 5.1 Shrunken Estimators

### 5.1.1 The SII estimator

Mayer and Villke (1973) discuss the use of several shrunken (SH) estimators of the form:

$$
\begin{equation*}
\hat{\beta}_{\mathrm{SH}}=\mathrm{d} \hat{\beta} \tag{5.1.1}
\end{equation*}
$$

where $0 \leq d \leq 1$ is a deterministically (fixed) or stochastically defined constant. By using (1.3.6) this estimator can be written as

$$
\begin{equation*}
\hat{\beta}_{S H}=d \sum_{i=1}^{p} v_{i} c_{i} / \lambda_{i} \tag{5.1.2}
\end{equation*}
$$

Thus the effect of $d / \lambda_{i}$ is to reduce the magnitude of the estimates.

### 5.1.2 Properties (d fixed)

1. Relationship to OLSE

From the definition it is clear that

$$
\begin{equation*}
\hat{\beta}_{\mathrm{SH}}=\mathrm{d} \hat{\beta} \tag{5.1.1}
\end{equation*}
$$

2. Expectation

$$
\begin{align*}
\mathrm{E}\left[\hat{\beta}_{\mathrm{SH}}\right] & =\mathrm{d} \beta  \tag{5.1.3}\\
& =\beta-(1-\mathrm{d}) \beta
\end{align*}
$$

Thus the estimator is biased, and the bias is

$$
\begin{equation*}
b=-(1-d) \beta \tag{5.1.4}
\end{equation*}
$$

3. Variance

$$
\begin{aligned}
\operatorname{Var}\left(\hat{\beta}_{\mathrm{SH}}\right) & =\mathrm{d}^{2} \operatorname{Var}(\hat{\beta}) \\
& =\mathrm{d}^{2} \sigma^{2}\left(\mathrm{X}^{\prime} \mathrm{X}\right)^{-1} \\
& =\sigma^{2} \mathrm{~d}^{2} \sum_{\mathrm{i}=1}^{\mathrm{p}} \mathrm{v}_{\mathrm{i}} \mathrm{v}_{\mathrm{i}}^{\prime} / \lambda_{\mathrm{i}}
\end{aligned}
$$

4. 

$$
\begin{align*}
\operatorname{MSE}\left[\hat{\beta}_{\mathrm{SH}}\right] & =\operatorname{Var}\left(\hat{\beta}_{\mathrm{SH}}\right)+\mathrm{bb}^{\prime} \\
& =\sigma^{2} \mathrm{~d}^{2}\left(\mathrm{X}^{\prime} \mathrm{X}\right)^{-1}+(1-\mathrm{d})^{2} \beta \beta^{\prime} \tag{5.1.7}
\end{align*}
$$

5. $\quad \operatorname{TMSE}\left[\hat{\beta}_{\mathrm{SH}}\right]=\operatorname{tr}\left(\operatorname{Var}\left(\hat{\beta}_{\mathrm{SH}}\right)+\mathrm{bb}^{\prime}\right)$

$$
\begin{align*}
& =\sigma^{2} \mathrm{~d}^{2} \sum_{i=1}^{p} 1 / \lambda_{i}+(1-\mathrm{d})^{2} \operatorname{tr}\left(\beta \beta^{\prime}\right) \\
& =\sigma^{2} \mathrm{~d}^{2} \sum_{i=1}^{p} 1 / \lambda_{i}+(1-\mathrm{d})^{2} \beta^{\prime} \beta
\end{align*}
$$

The SHE with $0 \leq d \leq 1$ and $d$ non-stochastic guarantees smaller variances than 0LS for the parameter estimators. The bias, however could offset the reduction in variances. Furthermore the $S H E$ is (total mean square) admissible (see comment 7 in the ridge section), as proved by Mayer and Willke (1973):

$$
\begin{align*}
\operatorname{TMSE}\left[\hat{\beta}_{\mathrm{SH}}\right] & <\operatorname{tr}[\operatorname{Var}(\hat{\beta})] \text { if and only if } \\
\mathrm{d} & >\frac{\beta^{\prime} \beta-\operatorname{tr}(\operatorname{Var}(\hat{\beta}))}{\beta^{\prime} \beta+\operatorname{tr}(\operatorname{Var}(\hat{\beta}))} \tag{5.1.9}
\end{align*}
$$

Since the RHS in (5.1.9) is less than 1 , there does exist a number $\mathrm{d}<1$ such that $\operatorname{TMSE}\left[\hat{\beta}_{\mathrm{SH}}\right]<\operatorname{tr}[\operatorname{Var}(\hat{\beta})]$. In fact if we minimize (5.1.8) with respect to $d$ we get:

$$
\frac{\delta \operatorname{TMSE}\left(\hat{\beta}_{\mathrm{SH}}\right)}{\delta \mathrm{d}}=\sigma^{2} 2 \mathrm{~d} \sum_{i=1}^{\mathrm{p}} 1 / \lambda_{\mathrm{i}}-2(1-\mathrm{d}) \beta^{\prime} \beta=0
$$

thus

$$
\begin{align*}
\sigma^{2} \mathrm{~d} \sum_{\mathrm{i}=1}^{\mathrm{p}} 1 / \lambda_{\mathrm{i}}+\mathrm{d} \beta^{\prime} \beta & =\beta^{\prime} \beta \\
\mathrm{d} & =\frac{\beta^{\prime} \beta}{\beta^{\prime} \beta+\operatorname{tr}(\operatorname{Var}(\hat{\beta}))}
\end{align*}
$$

The $d$ in (5.1.10) minimize (5.1.8), and the minimum TMSE of the SHE is

$$
\begin{align*}
\min \left\{\operatorname{TMSE}\left[\hat{\beta}_{\mathrm{SH}}\right]\right\} & =\mathrm{d}^{2} \operatorname{tr}\left(\operatorname{Var}(\hat{\beta})+\left(1-2 \mathrm{~d}+\mathrm{d}^{2}\right) \beta^{\prime} \beta\right. \\
& =\mathrm{d}^{2}\left(\beta^{\prime} \beta+\operatorname{tr}(\operatorname{Var}(\hat{\beta}))+\beta^{\prime} \beta-2 \mathrm{~d} \beta^{\prime} \beta\right. \\
& =\frac{\left(\beta^{\prime} \beta\right)^{2}}{\beta^{\prime} \beta+\operatorname{tr}(\operatorname{Var}(\hat{\beta}))}+\beta^{\prime} \beta-2 \frac{\left(\beta^{\prime} \beta\right)^{2}}{\beta^{\prime} \beta+\operatorname{tr}(\operatorname{Var}(\hat{\beta}))} \\
& =\frac{\beta^{\prime} \beta(\operatorname{tr}(\operatorname{Var}(\hat{\beta}))}{\beta^{\prime} \beta+\operatorname{tr}(\operatorname{Var}(\hat{\beta}))} \tag{5.1.11}
\end{align*}
$$

### 5.1.3 Properties (d stochastic)

The only known SHE of the form given in (5.1.1) with $d$ stochastic, and possessing any optimal properties, is the estimator due to James and Stein (1961).

Providing $\mathrm{p} \geq 3$ and $\mathrm{X}^{\prime} \mathrm{X}=\mathrm{I}$ the SHE is given by (5.1.1) with

$$
\mathrm{d}=\max \left\{0,\left(1-\mathrm{cv} / \hat{\beta}^{\prime} \hat{\beta}\right)\right\}
$$

where $0<c<2(p-2) /(v+2)$ and $v$ is the degrees of freedom for the residual sums of squares in OLS. The drawback of this SHE is the requirement that $p \geq 3$ and $X^{\prime} X=I$. Particularly the latter requirement eliminates most practical regression problems and completely eliminates the possible presence of collinearity.

Sclove (1968) modified the James-Stein estimator by suggesting shrinkage of just a subset of the components of $\hat{\beta}$. If the shrinkage is applied to those $\mathrm{p}-\mathrm{r}$ (where r is the rank of X ) components with smallest eigenvalues, his estimator (denoted be $\hat{\beta}_{\text {SSH }}$ ) may be written as,

$$
\hat{\beta}_{\mathrm{SSH}}=\left[\begin{array}{cc}
\mathrm{I}_{\mathrm{r}} & 0  \tag{5.1.12}\\
0 & \mathrm{dI} \\
\mathrm{p}-\mathrm{r}
\end{array}\right] \hat{\beta}
$$

for $0 \leq \mathrm{d} \leq 1$. Note that for $\mathrm{d}=0$ this estimator reduces to the PCE.

### 5.2 Fractional Principal Component (FPC) regression

In this section we firstly show that most of the biased estimators so far discussed, can be considered as FPC estimators. Then we obtain two new biased estimators due to Lee and Birch (1988).

### 5.2.1 FPC estimators

Consider the model defined in (3.1.1):

$$
\begin{align*}
Y & =Z \delta+\epsilon  \tag{5.2.1}\\
\text { with } \quad &  \tag{5.2.2}\\
\text { and } \quad Z^{\prime} Z & =\Delta^{2} \quad(\text { from }(3.1 .2)) \\
&  \tag{5.2.3}\\
& \\
& \\
& =\left(Z^{\prime} Z\right)^{-1} Z^{\prime} Y \quad(\text { from }(3.1 .3)) \\
&
\end{align*}
$$

Then a typical member of the class of FPC regression estimators is defined as:

$$
\begin{equation*}
\hat{\delta}_{\mathrm{FPC}}=\mathrm{F} \hat{\delta} \tag{5.2.4}
\end{equation*}
$$

where $\mathrm{F}=\operatorname{Diag}\left(\mathrm{f}_{1}, \mathrm{f}_{2}, \ldots, \mathrm{f}_{\mathrm{p}}\right)$ and $0 \leq \mathrm{f}_{\mathrm{j}} \leq 1, \mathrm{j}=1, \ldots, \mathrm{p}$. In addition the FPC estimator of $\hat{\beta}$ is

$$
\begin{equation*}
\hat{\beta}_{\mathrm{FPC}}=\mathrm{V} \hat{\delta}_{\mathrm{FPC}}=\mathrm{VF} \hat{\delta}=\mathrm{VFV}^{\prime} \hat{\beta} \tag{5.2.5}
\end{equation*}
$$

The diagonal matrix $F$ is termed the fraction matrix and the diagonal elements, $f_{j}$, are called the fractions. (This FPC estimator is one of a special class of linear transforms of $\hat{\delta}$ introduced by Mayer and Villke (1973)).

### 5.2.2 Relationship to other estimators

## 1. Least squares:

$$
\hat{\delta}_{\mathrm{FPC}}=\hat{\delta} \text { with } \mathrm{F}=\mathrm{I}
$$

## 2. Ridge

The ridge estimator of $\delta$ is

$$
\begin{aligned}
\hat{\delta}_{R} & =\left(\Delta^{2}+\mathrm{kI}\right)^{-1} \Delta^{2} \hat{\delta} \quad \text { (from (4.3.2) } \\
& =\mathrm{F} \hat{\delta}
\end{aligned}
$$

with $\mathrm{F}=\left(\Delta^{2}+\mathrm{kI}\right)^{-1} \Delta^{2}$, where $\mathrm{k}>0$

## 3. Generalized ridge

The i-th element of the generalized ridge estimator is

$$
\left[\hat{\delta}_{\mathrm{K}}\right]_{\mathrm{i}}=\lambda_{\mathrm{i}}\left(\lambda_{\mathrm{i}}+\mathrm{k}_{\mathrm{i}}\right)^{-1} \hat{\delta}_{\mathrm{i}}
$$

Thus, take $\mathrm{f}_{\mathrm{i}}=\lambda_{\mathrm{i}}\left(\lambda_{\mathrm{i}}+\mathrm{k}_{\mathrm{i}}\right)^{-1}, \mathrm{k}_{\mathrm{i}} \geq 0$.

## 4. Principal components

The PCE applies 1 to the first $\mathrm{p}-\mathrm{r}$ components (those components retained in the estimator) and 0 to the rest (those components that the estimator deletes).

## 5. Shrinkage

The concept of the FPC estimator is closely related to that of SHE's (5.1.1). Both estimators shrink the length of the LSE vector of the parameters toward the origin. For SHE the fractions are constant ( $f_{i}=d$ for all i) but FPC takes the individual PC's and shrinks each of them. Thus some PC's receive greater emphasis in estimating $\delta$ than others.

### 5.2.3 Properties of FPC estimators

1. Expectation

$$
\begin{align*}
\mathrm{E}\left(\hat{\delta}_{\mathrm{FPC}}\right) & =\mathrm{F} \delta \\
& =\delta-[\mathrm{I}-\mathrm{F}] \delta \tag{5.2.6}
\end{align*}
$$

Thus, the bias is

$$
\begin{equation*}
\mathrm{b}=-[\mathrm{I}-\mathrm{F}] \delta \tag{5.2.7}
\end{equation*}
$$

2. Variance

$$
\begin{align*}
\operatorname{Var}\left(\hat{\delta}_{\mathrm{FPC}}\right) & =\operatorname{FVar}(\hat{\delta}) \mathrm{F} \\
& =\sigma^{2} \mathrm{~F}^{-2} \mathrm{~F} \tag{5.2.8}
\end{align*}
$$

3. $\quad \operatorname{MSE}\left(\hat{\delta}_{\mathrm{FPC}}\right)=\operatorname{Var}\left(\hat{\delta}_{\mathrm{FPC}}\right)+\mathrm{bb}^{\prime}$

$$
\begin{equation*}
=\sigma^{2} \mathrm{~F} \Delta^{-2} \mathrm{~F}+[\mathrm{I}-\mathrm{F}] \delta \delta^{\prime}[\mathrm{I}-\mathrm{F}] \tag{5.2.9}
\end{equation*}
$$

Thus FPCE is an improvement over OLS in terms of MSE if $\operatorname{MSE}(\hat{\delta})-\operatorname{MSE}\left(\hat{\delta}_{\mathrm{FPC}}\right)=\sigma^{2} \Delta^{-2}\left(\mathrm{I}-\mathrm{F}^{2}\right)-(\mathrm{I}-\mathrm{F}) \delta \delta^{\prime}(\mathrm{I}-\mathrm{F})=\mathrm{S}$ (say), is a positive semi-definite matrix. $S$ will be positive semi-definite if $y^{\prime} S y \geq 0$ for any vector $y: n x 1$. Thus, where $\left(I-F^{2}\right)^{-}$is the generalized inverse of $\left(I-F^{2}\right)$, we have

$$
y^{\prime} \Delta^{2} S\left(I-F^{2}\right)^{-} y=y^{\prime}\left[\sigma^{2} I-\Delta^{2}(I-F) \delta \delta^{\prime}(I-F)\left(I-F^{2}\right)^{-}\right] y \quad \text { for all } y
$$

Consequently, the necessary and sufficient condition for $S$ to be positive semi-definite is

$$
\begin{equation*}
\sigma^{2} \geq \delta^{\prime} \Delta^{2}(\mathrm{I}-\mathrm{F})^{2}\left(\mathrm{I}-\mathrm{F}^{2}\right)^{-} \delta \tag{5.2.10}
\end{equation*}
$$

Thus the positive semi-definitenes of $S$ will depend on the vector $\delta$, the degree of collinearity ( $\Delta$ ), and $\sigma$.
4. $\operatorname{TMSE}\left(\hat{\delta}_{\mathrm{FPC}}\right)=\operatorname{tr}\left[\operatorname{MSE}\left(\left(\hat{\delta}_{\mathrm{FPC}}\right)\right]\right.$

$$
\begin{equation*}
=\sigma^{2} \sum_{i=1}^{p} f_{i}^{2} / \lambda_{i}+\sum_{i=1}^{p}\left(1-f_{i}\right)^{2} \delta_{i}^{2} \tag{5.2.11}
\end{equation*}
$$

We note that:
(i) If $\mathrm{f}_{\mathrm{i}} \rightarrow 0$, then the $\mathbf{i}$-th diagonal element of $\operatorname{Var}\left(\hat{\delta}_{\mathrm{FPC}}\right) \rightarrow 0$ and the $(\text { bias })^{2} \rightarrow \delta_{\mathrm{i}}^{2}$.
(ii) If $\mathrm{f}_{\mathrm{i}} \rightarrow 1$, then the i - th diagonal element of $\operatorname{Var}\left(\hat{\delta}_{\mathrm{FPC}}\right) \rightarrow[\operatorname{Var}(\hat{\delta})]_{\mathrm{i}}$ and the $(\text { bias })^{2} \rightarrow 0$.

### 5.2.4 Optimal values for fractions

One set of optimal values of the fractions is obtained by minimizing $\operatorname{TMSE}\left(\hat{\delta}_{\mathrm{FPC}}\right)$ with respect to the $\mathrm{f}_{\mathrm{j}}$ 's. Thus

$$
\frac{\operatorname{dTMSE}\left(\hat{\delta}_{\mathrm{FPC}}\right)}{d f_{\mathrm{j}}}=2 \sigma^{2} \mathrm{f}_{\mathrm{j}} / \lambda_{\mathrm{j}}-2 \delta_{\mathrm{j}}\left(1-\mathrm{f}_{\mathrm{j}}\right)
$$

Let the $j$-th optimal value be denoted by $f_{j}^{0}$, thus

$$
\begin{equation*}
f_{j}^{0}=\delta_{j}^{2} \lambda_{j}\left(\sigma^{2}+\delta_{j}^{2} \lambda_{j}\right)^{-1} \tag{5.2.12}
\end{equation*}
$$

and the min TMSE will be

$$
\begin{align*}
& \min \left\{\operatorname{TMSE}\left(\hat{\delta}_{\mathrm{FPC}}\right)\right\}=\sigma^{2} \sum_{\mathrm{i}=1}^{\mathrm{p}} \delta_{\mathrm{j}}^{4} \lambda_{\mathrm{j}}^{2}\left(\sigma^{2}+\delta_{\mathrm{j}}^{2} \lambda_{\mathrm{j}}\right)^{-2} / \lambda_{\mathrm{i}}+\sum_{\mathrm{i}=1}^{\mathrm{p}} \sigma^{4}\left(\sigma^{2}+\delta_{\mathrm{j}}^{2} \lambda_{\mathrm{j}}\right)^{-2} \delta_{\mathrm{j}}^{2} \\
& =\sum_{\substack{\mathrm{i}=1 \\
\mathrm{p}}}^{\mathrm{p}}\left(\sigma^{2} / \lambda_{\mathrm{i}}\right)\left[\delta_{\mathrm{j}}^{2} \lambda_{\mathrm{i}}\right]\left(\sigma^{2}+\delta_{\mathrm{j}}^{2} \lambda_{\mathrm{j}}\right)^{-1} \\
& =\sum_{i=1}^{p}\left(\sigma^{2} / \lambda_{i}\right) f_{i}^{0}  \tag{5.2.13}\\
& =\sum_{i=1}^{i=1}\left(\sigma^{2} / \lambda_{i}\right)\left[1-\left(1-f_{i}^{0}\right)\right] \\
& =\operatorname{TMSE}(\hat{\delta})-\sum_{i=1}^{p}\left(\sigma^{2} / \lambda_{i}\right)\left(1-f_{i}^{0}\right) \tag{5.2.14}
\end{align*}
$$

The second term in the above equation (i.e. $\left.\quad \Sigma^{p}\left(\sigma^{2} / \lambda_{i}\right)\left(1-f_{i}^{0}\right)\right)$ is the maximum reduction in THSE when the optimal FPC is used instead of 0LS. This quantity will always be positive and will increase as the degree of collinearity increases (because of the term $\lambda_{i}^{-1}$ ).

The theoretically optimal fraction matrix cannot be used in practice since the $f_{j}^{0}$, defined in (5.2.12), contain the unknown parameters $\delta_{j}$ and $\sigma^{2}$. Therefore Lee and Birch (1988) suggested two new biased estimators both using the optimal FPC but with different approaches in estimating $f_{j}^{0}$.

### 5.2.5 Optimal FPC estimators

In §4.6.1 (generalized ridge) it was shown that an optimal value for $\mathrm{k}_{\mathrm{i}}$ is obtained when $\mathrm{k}_{\mathrm{i}}=\sigma^{2} / \delta_{\mathrm{i}}^{2}$. Hoerl and Kennard (1970) estimated $\mathrm{k}_{\mathrm{i}}$ by estimating $\sigma^{2}$ and $\delta_{\mathrm{i}}^{2}$ as their 0LS estimates. They used an iterative method to obtain a new value for $\hat{\mathbf{k}}_{\mathrm{i}}=\hat{\sigma}^{2} / \hat{\delta}_{\mathrm{K}}^{2}$ and continued until the estimate stabilised. The same method is applied to form the iterative (optimal) FPC estimator (abbreviated as FPCI). For the FPCI estimator the iterative scheme of the optimal fraction is:

$$
\begin{equation*}
f_{j, G R}(t+1)=\frac{\lambda_{j}}{\lambda_{j}+S^{2} /\left[\hat{\delta}_{K}(t)\right]_{j}^{2}}, \quad t=0,1,2, \ldots \tag{5.2.15}
\end{equation*}
$$

where t denotes the iteration number, $\mathrm{s}^{2}$ is the 0LS estimate of $\sigma^{2}$ and $\left[\hat{\delta}_{K}(\mathrm{t})\right]_{\mathrm{j}}$ is the generalized ridge estimate of $\delta_{j}$ at the t - th iteration with $\left[\hat{\delta}_{\mathrm{K}}(0)\right]_{\mathrm{j}}=\hat{\delta}_{\mathrm{j}}$. The iteration continues until there is stability achieved in the length of the generalized ridge estimator $\left(\left[\hat{\delta}_{\mathrm{K}}(\mathrm{t})\right]_{\mathrm{j}}\right)$.

In the presence of collinearity the starting values of 0LS may be severely perturbed and therefore it may be more beneficial to use a biased estimator as initial value. If one considers a PCE as the initial value, then the
iterative scheme (5.2.15) can be used to compute the fractions, where the starting value for $\left[\hat{\delta}_{K}(0)\right]_{j}$ would be $\left[\hat{\delta}_{P C}\right]_{j}$ and $S^{2}$ is replaced with $\mathrm{s}^{2}(\mathrm{t})=(\mathrm{Y}-\mathrm{Za}(\mathrm{t}))^{\prime}(\mathrm{Y}-\mathrm{Za}(\mathrm{t})) /(\mathrm{n}-\mathrm{p})$. Here $\mathrm{a}(\mathrm{t})$ is the estimate of $\delta$ at the t-th iteration, $\mathrm{a}(0)=\hat{\delta}_{\mathrm{PC}}$. Fractions obtained in this way will be denoted by $f_{j}, P C(t+1)$.

The final resulting estimator, denoted by $\hat{\delta}_{\text {FPCI }}$, with the limiting fraction $\operatorname{matrix} \mathrm{F}_{\mathrm{PCI}}$, is

$$
\begin{equation*}
\hat{\delta}_{\mathrm{FPCI}}=\mathrm{F}_{\mathrm{PCI}} \hat{\delta} \tag{5.2.16}
\end{equation*}
$$

where $\mathrm{F}_{\mathrm{PCI}}=\operatorname{Diag}\left(\mathrm{f}_{1}^{*}, \mathrm{PC}, \ldots, \mathrm{f}_{\mathrm{p}, \mathrm{PC}}^{*}\right)$, and $\mathrm{f}_{\mathrm{j}, \mathrm{PC}}^{*}=\lim _{\mathrm{t}}\left[\mathrm{f}_{\mathrm{j}}, \mathrm{PC}(\mathrm{t}+1)\right]$, for $j=1, \ldots, p$. The FPCI estimator is thus formed from the combined concepts of the PC estimator and the iterative generalized ridge estimator. It is interesting to note that as early as Marquardt (1970), the possibility is suggested of a combined estimator i.e. $\quad \hat{\delta}_{c}=\left(Z_{1}^{\prime} Z_{1}+k I\right)^{-1} Z_{1}^{\prime} Y$ where $k$ is chosen to deflate the effects of the remaining near-zero eigenvalues left over in $Z_{1}$.

The second biased estimator due to Lee and Birch (1988) is based on the iterative ridge estimator concept. In this scheme the fraction in (5.2.15) becomes

$$
\begin{equation*}
\mathrm{f}_{\mathrm{j}, \mathrm{R}}(\mathrm{t}+1)=\frac{\lambda_{\mathrm{j}}}{\lambda_{\mathrm{j}}+\mathrm{s}^{2} /\left[\hat{\delta}_{\mathrm{R}}(\mathrm{t})^{\prime} \hat{\delta}_{\mathrm{R}}(\mathrm{t}) / \mathrm{p}\right]}, \mathrm{t}=0,1,2, \ldots \tag{5.2.17}
\end{equation*}
$$

where $\hat{\delta}_{R}(\mathrm{t})$ is the ridge estimate of $\delta$ at the t - th iteration with $\hat{\delta}_{\mathrm{R}}(0)=\hat{\delta}$. Just as in the FPCI estimator, the authors replaced the OLSE by the PCE $\left(\hat{\delta}_{R}(0)=\hat{\delta}_{P C}\right)$ and the $s^{2}$ by $s^{2}(t)$, and denoted the fractions by $f_{j}, P C V(t+1)$. The resulting estimator, denoted by $\hat{\delta}_{\text {FPCV }}$, is defined as:

$$
\begin{equation*}
\hat{\delta}_{\mathrm{FPCV}}=\mathrm{F}_{\mathrm{PCV}} \hat{\delta} \tag{5.2.18}
\end{equation*}
$$

where $F_{P C V}=\operatorname{Diag}\left(f_{1}^{*}, P C V, \ldots, f_{p, P C V}^{*}\right)$, and $f_{j}^{*}, P C V=\lim _{t}\left[f_{j, P C V}(t+1)\right]$, for $j=1, \ldots, p$.

The authors have observed that a 1-step version of both (5.2.15) and (5.2.17) exhibited already improved estimation properties over other biased estimators in the data sets they studied. More comments on these estimators and their performances against other estimators will appear in Chapter 8 and Chapter 10 .

### 5.3 Summary

In this chapter we introduced the shrinkage estimators as well as various fractional principal component regression estimators. The estimators were defined, and their properties given and discussed.

## Chapter 6

## LATENT ROOT REGRESSION ESTIMATION

Latent root regression (LRR)is similar to PC regression with the difference that it operates on the augmented matrix of explanatory variables and response in correlation form rather than the matrix of regressors alone. The method was originally proposed by Vebster et al. (1974) and by Hawkins (1973).

### 6.1 The LRR estimator

Define the centered and standardized vector $\mathrm{Y}^{*}$ ( $\mathrm{n} \times 1$ ) as

$$
\begin{equation*}
Y^{*}=(Y-\bar{y} 1) / s_{y} \tag{6.1.1}
\end{equation*}
$$

where $\bar{y}=\sum_{i=1}^{n} Y_{i} / n \quad$ and $s_{y}^{2}=\sum_{i=1}^{n}\left(Y_{i}-\bar{y}\right)^{2}$
(Instead of just saying $Y$ is centered and standardized we introduce this notation as it is useful in developing the estimator).

Assume that X is centered and standardized and form the augmented matrix $\left[\mathrm{X} \mathrm{Y}^{*}\right]$, so that $\left[\mathrm{X} \mathrm{Y}^{*}\right]^{\prime}\left[\mathrm{X} \mathrm{Y}{ }^{*}\right]$ will be in correlation form. Then by using the SVD of an augmented matrix (1.3.8),

$$
\text { with } \begin{aligned}
& {\left[\mathrm{XY} \mathrm{Y}^{*}\right] }=\tilde{\mathrm{U}} \tilde{\Delta} \tilde{\mathrm{~V}}^{\prime} \\
& \tilde{\mathrm{U}}=\left[\tilde{\mathrm{u}}_{1} \ldots \tilde{\mathrm{u}}_{\mathrm{p}+1}\right], \tilde{\mathrm{u}}_{\mathrm{i}}: \mathrm{nx} 1 \\
& \tilde{\mathrm{~V}}=\left[\tilde{\mathrm{v}}_{1} \ldots \tilde{\mathrm{v}}_{\mathrm{p}+1}\right], \tilde{\mathrm{v}}_{\mathrm{i}}:(\mathrm{p}+1) \times 1 \\
& \tilde{\Delta}=\operatorname{diag}\left[\omega_{1}, \omega_{2}, \ldots \omega_{\mathrm{p}+1}\right] \\
& \omega_{1} \geq \omega_{2} \geq \ldots \ldots \omega_{\mathrm{p}+1} \\
& \tilde{\mathrm{U}}^{\prime} \tilde{\mathrm{U}}=\tilde{\mathrm{V}}^{\prime} \tilde{\mathrm{V}}=\mathrm{I}_{\mathrm{p}+1}
\end{aligned}
$$

and the following notation defined in Chapter 1:

1. Let $\tilde{v}_{i}, j$ be the $j$-th component of the i-th right singular vector $\tilde{v}_{i}$.
2. $\tilde{v}_{i}^{0}$ is the $p$-dimensional vector containing the first $p$ components of the i-th right singular vector $\tilde{v}_{i}$ of $\left[\mathrm{X} \mathrm{Y}^{*}\right]$ of dimension $\mathrm{p}+1$, thus

$$
\tilde{v}_{i}=\left[\tilde{v}_{i}^{0} \tilde{v}_{i}, p+1\right]^{\prime} .
$$

Then by using the definition of an eigenvalue, one can write

$$
\left[\begin{array}{ll}
\mathrm{X} & \left.\mathrm{Y}^{*}\right] \omega_{\mathrm{i}}^{2}=[\mathrm{X}
\end{array} \mathrm{Y}^{*}\right]\left[\begin{array}{ll}
\tilde{\mathrm{v}}_{\mathrm{i}}^{\prime} & \tilde{\mathrm{v}}_{\mathrm{i}}, \mathrm{p}+1
\end{array}\right]^{\prime} .
$$

If $\omega_{\mathrm{i}}{ }^{2} \approx 0$ (indicating a collinearity) then

$$
\begin{align*}
{\left[\mathrm{X} \mathrm{Y}^{*}\right]\left[\tilde{\mathrm{v}}_{\mathrm{i}}^{0} \tilde{\mathrm{v}}_{\mathrm{i}, \mathrm{p}+1}\right]^{\prime} } & \approx 0 \quad(0: \mathrm{nx} 1) \text { or } \\
X \tilde{\mathrm{v}}_{\mathrm{i}}^{0}+\mathrm{Y}^{*} \tilde{\mathrm{v}}_{\mathrm{i}, \mathrm{p}+1} & \approx 0 \tag{6.1.3}
\end{align*}
$$

If $\tilde{v}_{i, p+1}$ (the weight of $Y^{*}$ ) is not trivially close to zero, the response variable $\left(Y^{*}\right)$ is involved in the collinearity. The key difference between detecting collinearities in $X^{\prime} X$ and those in the augmented correlation matrix is that $Y$ can affect the collinearities in $\left[X Y^{*}\right]^{\prime}\left[X Y^{*}\right]$, and any such collinearity has predictive value since (6.1.3) implies

$$
\begin{align*}
& \hat{Y}^{*} \approx-\left(\tilde{v}_{i}, p+1\right)^{-1} X \tilde{v}_{i}^{0}  \tag{6.1.4}\\
& \hat{Y}=\bar{y} 1-s_{y}\left(\tilde{v}_{i}, p+1\right)^{-1} X \tilde{v}_{i}^{0} \quad(\text { from }(6.1 .1)) \tag{6.1.5}
\end{align*}
$$

Definition 6.1: When $w_{i} \approx 0$ and $\left(\tilde{v}_{i, p+1}\right) \approx 0$ the multicollinearity is referred to as a non-predictive collinearity, i.e., a collinearity among the
predictor variables that is of little value in predicting the response variable.

Gunst and Mason (1980) suggest values of $\omega_{i} \leq 0.1$ and ( $\left.\tilde{v}_{i}, p^{+}+1\right) \leq 0.1$ should be investigated for non-predictive collinearities)

Providing all $\left(\tilde{v}_{i},{ }^{p}+1\right) \neq 0$, one can define $\mathrm{p}+1$ prediction equations like (6.1.5) (one for each eigenvector). Let the i-th prediction equation be denoted by $\hat{\mathbf{Y}}_{\mathrm{i}}$, thus, for $\mathrm{i}=1,2, \ldots, \mathrm{p}+1$

$$
\begin{equation*}
\hat{\mathbf{Y}}_{i}=\overline{\mathrm{y}} 1-\mathrm{s}_{\mathrm{y}}\left(\tilde{\mathrm{v}}_{\mathrm{i}}, \mathrm{p}+1\right)^{-1} \mathrm{X} \tilde{\mathrm{v}}_{\mathrm{i}}^{0} . \tag{6.1.6}
\end{equation*}
$$

Linear combinations of the predictors in (6.1.6) will be used to obtain estimates of the parameters of the model. Consider the following arbitrary linear combination of the predictors:

$$
\hat{\mathbf{Y}}=\sum_{i=1}^{p+1} a_{i}\left(\tilde{v}_{i}, p^{+1}\right) \hat{\mathbf{Y}}_{i} .
$$

A value for $\mathrm{a}_{\mathrm{i}}$ will be obtained in (6.1.15). Imposing the restriction

$$
\begin{align*}
& \sum_{i=1}^{p+1} a_{i}\left(\tilde{v}_{i}, p+1\right)=1  \tag{6.1.7}\\
& \hat{\mathbf{Y}}=\sum_{\mathrm{i}=1}^{\mathrm{p}+1} \mathrm{a}_{\mathrm{i}}\left(\tilde{\mathrm{v}}_{\mathrm{i}}, \mathrm{p}+1\right)\left\{\overline{\mathrm{y}} \mathbf{1}-\mathrm{s}_{\mathrm{y}}\left(\tilde{\mathrm{v}}_{\mathrm{i}}, \mathrm{p}^{+1}\right)^{-1} \mathrm{X} \tilde{\mathrm{v}}_{\mathrm{i}}^{0}\right\} \\
& =\bar{y} 1-s_{y} X \sum_{i=1}^{p+1} a_{i} \tilde{v}_{i}^{0}  \tag{6.1.8}\\
& \text { yields }
\end{align*}
$$

The residual sum of squares using this predictor is

$$
\begin{aligned}
(Y-\hat{Y})^{\prime}(Y-\hat{Y}) & =\left(Y-\bar{y} 1+s_{y} X \sum_{i=1}^{p+1} a_{i} \tilde{v}_{i}^{0}\right)^{\prime}\left(Y-\bar{y} 1+s_{y} X \sum_{i=1}^{p+1} a_{i} \tilde{v}_{i}^{0}\right) \\
& =s_{y}^{2}\left[\left(Y^{*}+X \sum_{i=1}^{p+1} a_{i} \tilde{v}_{i}^{0}\right)^{\prime}\left(Y^{*}+X \sum_{i=1}^{p+1} a_{i} \tilde{v}_{i}^{0}\right)\right. \\
& \left.=s_{y}^{2}\left\{\left[X Y^{*}\right]\left[\sum_{i=1}^{p+1} a_{i} \tilde{v}_{i}^{0}\right]^{\prime}\right]^{\prime}\right\}^{\prime}\left\{\left[X Y^{*}\right]\left[\sum_{i=1}^{p+1} a_{i} \tilde{v}_{i}^{0} 1\right]^{\prime}\right\}
\end{aligned}
$$

and $\left.\underset{i=1}{p+1} \sum_{i} \tilde{v}_{i}^{0} 1\right]^{\prime}=\left[\sum_{i=1}^{p+1} a_{i} \tilde{v}_{i}^{0}, \sum_{i=1}^{p+1} a_{i}\left(\tilde{v}_{i}, p+1\right)\right]^{\prime} \quad$ (from (6.1.7))

$$
\begin{align*}
& =\left[\sum_{i=1}^{p+1} a_{i}\left[\tilde{v}_{i}^{\prime} \tilde{v}_{i}, p^{+1}\right]\right]^{\prime} \\
& =\left[\sum_{i=1}^{p+1} a_{i} \tilde{v}_{i}^{\prime}\right]^{\prime} \quad \text { (from 1.3.10) } \\
& =\left[a^{\prime} \tilde{V}^{\prime}\right]^{\prime} \\
& =\tilde{V} a \tag{6.1.9}
\end{align*}
$$

where $a^{\prime}=\left[a_{1}, \ldots, a_{p+1}\right]$. Thus the RSS is

$$
\begin{align*}
& =\mathrm{s}_{\mathrm{y}}^{2} \mathrm{a}^{\prime} \tilde{\Delta}^{2} \mathrm{a} \quad(\text { from (1.3.8)) }  \tag{6.1.10}\\
& =\mathrm{s}_{\mathrm{y}}^{2} \sum_{\mathrm{i}=1}^{\mathrm{p}+1} \mathrm{a}_{\mathrm{i}}^{2} \omega_{\mathrm{i}}^{2} \tag{6.1.11}
\end{align*}
$$

To find the value of a for which this residual sum of squares will be a minimum, one has to minimize $f(a)$ where

$$
f(a)=s_{y}^{2} \sum_{i=1}^{p+1} a_{i}^{2} w_{i}^{2}-2 u_{0}\left(\sum_{i=1}^{p+1} a_{i}\left(\tilde{v}_{i}, p^{+1}\right)-1\right)
$$

and $2 u_{0}$ is a Lagrangian multiplier. Now, for $j=1, \ldots, p+1$

$$
\begin{equation*}
\frac{d f(a)}{d a_{j}}=2 s_{y}^{2} a_{j} \omega_{i}^{2}-2 u_{0}\left(\tilde{v}_{j}, p+1\right)=0 \tag{6.1.12}
\end{equation*}
$$

so

$$
\begin{equation*}
a_{j}=u_{o}\left(\tilde{v}_{j}, p+1\right) /\left(s_{y}^{2} \omega_{j}^{2}\right) \tag{6.1.13}
\end{equation*}
$$

and from the restriction in (6.1.7)

$$
\begin{align*}
& \sum_{j=1}^{p+1} a_{j}\left(\tilde{v}_{j}, p+1\right)=u_{o} \sum_{j=1}^{p+1}\left(\tilde{v}_{j}, p+1\right)^{2} /\left(s_{y}^{2} w_{j}^{2}\right) \\
& u_{o}=1 /\left\{\sum_{j=1}^{p+1}\left(\tilde{v}_{j}, \mathrm{p}+1\right)^{2} /\left(\mathrm{s}_{\mathrm{y}}^{2} \omega_{\mathrm{j}}^{2}\right)\right\} \tag{6.1.14}
\end{align*}
$$

Substituting this value in (6.1.13) for $j=1, \ldots, p+1$ gives

$$
\begin{align*}
a_{j} & =\left[1 /\left\{\sum_{j=1}^{p+1}\left(\tilde{v}_{j}, p+1\right)^{2} /\left(s_{y}^{2} \omega_{j}^{2}\right)\right\}\right]\left(\tilde{v}_{j}, p+1\right) /\left(s_{y}^{2} \omega_{j}^{2}\right) \\
& =\left(\tilde{v}_{j}, p+1\right) \omega_{j}^{-2}\left\{\underset{j=1}{p+1}\left(\tilde{v}_{j}, p+1\right)^{2} / \omega_{j}^{2}\right\}^{-1} \tag{6.1.15}
\end{align*}
$$

Then if we put this value in (6.1.8):

$$
\begin{align*}
\hat{\mathbf{Y}} & \left.=\bar{y} 1-\operatorname{s}_{y} \sum_{\substack{p+1}}^{\sum_{i=1}^{p+1}} \tilde{v}_{i}, p+1\right) \omega_{i}^{-2}\left\{\sum_{j=1}^{p+1}\left(\tilde{v}_{j}, p+1\right)^{2} / \omega_{j}^{2}\right\}^{-1} \tilde{v}_{i}^{0} \\
& =\bar{y} 1+X \sum_{i} f_{i} \tilde{v}_{i}^{0} \tag{6.1.16}
\end{align*}
$$

where

$$
\begin{equation*}
f_{i}=-s_{y}\left(\tilde{v}_{i}, p+1\right) \omega_{i}^{-2}\left\{\sum_{j=1}^{p+1}\left(\tilde{v}_{j}, p+1\right)^{2} / \omega_{j}^{2}\right\}^{-1} \tag{6.1.17}
\end{equation*}
$$

Suppose that there are s non-predictive near-singularities, i.e.

$$
\omega_{p+1-s}=\omega_{p+1-s-1}=\ldots=\omega_{p+1-1}=\omega_{p+1} \approx 0
$$

and

$$
\tilde{\mathrm{v}}_{\mathrm{p}+1-\mathrm{s}, \mathrm{p}+1}=\ldots=\tilde{\mathrm{v}}_{\mathrm{p}+1-1}, \mathrm{p}+1=\tilde{\mathrm{v}}_{\mathrm{p}+1}, \mathrm{p}+1 \approx 0
$$

Then

$$
\begin{equation*}
\hat{\mathbf{Y}}=\overline{\mathrm{y}} 1+X \sum_{\mathrm{i}=1}^{\mathrm{p}+1-\mathrm{s}} \mathrm{f}_{\mathrm{i}} \tilde{\mathrm{v}}_{\mathrm{i}}^{0} \tag{6.1.18}
\end{equation*}
$$

Then the latent root estimator ( $\hat{\beta}_{\mathrm{LR}}$ of $\beta$ ) is

$$
\begin{equation*}
\hat{\beta}_{L R}=\sum_{i=1}^{p+1-s} f_{i} \tilde{v}_{i}^{0} \tag{6.1.19}
\end{equation*}
$$

and $f_{i}$ is as defined in (6.1.17) for $i=1,2, \ldots, p+1-s$

$$
f_{i}=-s_{y}\left(\tilde{v}_{i}, p+1\right) \omega_{i}^{-2}\left\{\sum_{j=1}^{p+1}\left(\tilde{v}_{j}, p+1\right)^{2} / \omega_{j}^{2}\right\}^{-1}
$$

Some authors (i.e. Gunst and Mason (1980) define $\hat{\beta}_{L R}=\sum_{i=1}^{p+1} f_{i} \tilde{v}_{i}^{0}$ where $f_{i}=0$ if $\omega_{i} \approx 0$ and $\tilde{v}_{p+1, p+1} \approx 0 ;$ otherwise $f_{i}$ will have the value defined in (6.1.17).

If $\tilde{V}=\left[\tilde{\mathrm{v}}_{1} \ldots \tilde{\mathrm{v}}_{\mathrm{p}+1}\right], \quad \tilde{\mathrm{v}}_{\mathrm{i}}:(\mathrm{p}+1) \mathrm{x} 1$ is partitioned as $\tilde{V}=\left[\begin{array}{ll}\tilde{\mathrm{V}}_{1} & \tilde{\mathrm{~V}}_{2}\end{array}\right]$ where $\tilde{\mathrm{V}}_{1}=\left[\tilde{\mathrm{v}}_{1} \ldots \tilde{\mathrm{v}}_{\mathrm{p}+1-\mathrm{s}}\right], \quad \tilde{\mathrm{V}}_{2}=\left[\tilde{\mathrm{v}}_{\mathrm{p}+2-\mathrm{s}} \ldots \tilde{\mathrm{v}}_{\mathrm{p}+1}\right]$, then by defining the vector $f=\left[f_{A}^{\prime} f_{B}^{\prime}\right]^{\prime}$, and $f_{A}=\left[f_{1}, f_{2}, \ldots f_{p+1-s}\right]^{\prime}, f_{B}=\left[f_{p+2-s} \ldots f_{p+1}\right]^{\prime}$, LRR can be viewed as the use of a restricted least squares estimator where the restriction is $\tilde{V}_{2} f_{B}=0$

When $\mathrm{s}=0$ we have $\hat{\beta}_{\mathrm{Lr}}=\hat{\beta}$, since both minimize SSE without imposing any restriction on the estimators and the RSS will be given by (6.1.11) and (6.1.15) as

$$
\begin{align*}
\text { RSS } & =s_{y}^{2} \sum_{j=1}^{p+1} a_{j}^{2} w_{j}^{2} \\
& =s_{y}^{2} \sum_{\substack{p+1 \\
p+1}}\left(\tilde{v}_{j}, p+1\right)^{2} w_{j}^{-4}\left\{\sum_{j=1}^{p+1}\left(\tilde{v}_{j}, p+1\right)^{2} / w_{j}^{2}\right\}^{-2} w_{j}^{2} \\
& =s_{y}^{2}\left\{\sum_{j=1}^{2}\left(\tilde{v}_{j}, p+1\right)^{2} / w_{j}^{2}\right\}^{-1} \\
& =u_{o} \tag{6.1.20}
\end{align*}
$$

When $\mathrm{s} \neq 0$ then

$$
\begin{align*}
\text { RSS } & =s_{y}^{2} \sum_{j=1}^{\mathrm{p}+1-\mathrm{s}} \mathrm{a}_{\mathrm{j}}^{2} \omega_{\mathrm{j}}^{2} \\
& =\mathrm{s}_{\mathrm{y}}^{2}\left\{\sum_{\mathrm{j}=1}^{\mathrm{p}+\mathrm{s}-\mathrm{s}}\left(\tilde{v}_{\mathrm{j}}, \mathrm{p}+1\right)^{2} / \omega_{\mathrm{j}}^{2}\right\}^{-1} \tag{6.1.21}
\end{align*}
$$

Comment:

1. $\sum_{j=1}^{p+1-s}\left(\tilde{v}_{j}, p+1\right)^{2} / \omega_{j}^{2} \leq \sum_{j=1}^{p+1}\left(\tilde{v}_{j}, p+1\right)^{2} / \omega_{j}^{2} \quad$ so that RSS in (6.1.21) will be larger ${ }^{j=1}$ than RSS in $(6.1 .20)^{j=1}$ (the full model).
2. If one is removing non-predictive near-singularities (6.1.21) will be approximately equivalent to (6.1.20) as the $s$ non-predictive collinearities that have been removed contribute very little to $\sum_{\mathrm{j}=1}^{\mathrm{p}+1}\left(\tilde{v}_{\mathrm{j}}, \mathrm{p}+1\right)^{2} / \omega_{\mathrm{j}}^{2}$. Thus the $\mathrm{j}=1$
estimated $\hat{\sigma}^{2}$ using (6.1.20) and (6.1.21) will be reasonably close. When the number of non-predictive collinearities (s) increases, $\hat{\sigma}^{2}$ for the full model will be larger then that of (6.1.21) as the degrees of freedom for the latter will be $n-p-1+s$ (pooling the error associated with the $s$ nonpredictive collinearities with the error sums of squares of the restricted model (6.1.21)).

### 6.2 Properties and Problems

1. Expectation

$$
\begin{align*}
E\left(\hat{\beta}_{L R}\right) & =E\left[\sum_{i=1}^{p+1} f_{i} \tilde{v}_{i}^{0}-\sum_{i=t}^{p+1} f_{i} \tilde{v}_{i}^{0}\right] \text { where } t=(p+1-s)+1 \\
& =\beta-E\left[\sum_{i=t}^{p+1} f_{i} \tilde{v}_{i}^{0}\right] \tag{6.2.1}
\end{align*}
$$

2. The LRRE has no known MSE function and so can not be directly compared with other estimators. Gunst et al. (1976) report that in many cases the LRRE behaves similarly to the PCE so that

$$
\begin{equation*}
\operatorname{TMSE}\left(\hat{\beta}_{\mathrm{LR}}\right) \approx \operatorname{TMSE}\left(\hat{\beta}_{\mathrm{PC}}\right) \tag{6.2.2}
\end{equation*}
$$

Thus LRRE yields a biased estimator, but one which apparently greatly reduces the TMSE of the estimator. The theoretical conditions necessary for (6.2.2) to hold are unknown at the present time.
3. No criterion of the smallness of the eigenvalues and last element of the corresponding eigenvector that admits deletion, is as yet adequately defined. The decision will depend on the number of independent variables, the sample size ( n ) and the degree of ill-conditioning. Webster et al. (1974) suggested in their example ( $n=12$ and six regressors) that suitably small would mean $\omega_{i}^{2} \leq 0.05$ and $\tilde{\mathrm{v}}_{\mathrm{i}}, \mathrm{p}^{+1} \leq 0.10$.
4. Though Vebster et al. (1974) suggested that accuracy of computer algorithms for determining eigenvectors and eigenvalues under various degrees of ill-conditioning should be investigated, this may not be a relevant problem today.

Note that LRRE is in some sense simply a computational device like TLS with a justification but no complete statistical development.

### 6.3. Elimination of independent variables

Webster et al. (1974) as well as Gunst et al. (1976) discussed and compared OLS and LRRE on their performances in eliminating independent variables via the stepwise technique. Both authors preferred the backward elimination procedure (see Chapter 1) and Webster et al. (1974) proposed an F-statistic for LRR similar to the OLS F-statistic. They also suggested that perhaps several regressor variables could be deleted at the first stage of LRR.

Thus in a backward elimination procedure, they applied the following variable selection rule:

0LS: Determine the minimum $F$-value for each of the $p$ variables in the model (1.1); if the minimum $F<1$, delete that particular regressor from the predictor.

LRR: Determine each of the F -values for the p variables in the model (6.1.18); for each $F<1$, delete the corresponding regressor from the predictor.

The advantage of LRR over OLS is that by using LRR one firstly eliminates the non-predictive collinearities (unmasks them) while in OLS they remain masked. The true influences of the independent variables on the dependent variable are more clearly represented. In many cases several independent variables may be eliminated at the first stage when the computations are the easiest.

### 6.4 Summary

In this chapter we introduced latent root regression estimation. Ve discussed the estimator, its properties, related problems and methods of eliminating independent variables.

## Chapter 7

## TOTAL LEAST SQUARES

Total least squares (TLS) is a method of fitting that is appropriate when there are errors in both the observation vector $Y(n x 1)$ and in the independent variable matrix $X$ ( $n x p$ ). Errors in X-variables models contrast with the classical regression model where the data matrix $X$ is assumed to be error free.

### 7.1 Introduction

The problem of parameter estimation when all the variables involve error, has a long history. For an overview of the historical development see Van Huffel and Vandewalle (1985). Recently the linear errors-in-variables (all observations are coming from some unknown true values plus measurement errors) has been treated by the method of total least squares (TLS) based on the singular value decomposition (SVD) as proposed by Golub and Van Loan (1980) and further developed by Van Huffel, Vandewalle and Staar (1984).

In $\S 7.2$ the TLS technique is discussed, the use of TLS in the errors-invariables model in $\S 7.3$ and subset selection based on TLS for prediction, in §7.4. The comparison of TLS estimator with the estimators discussed in Chapters 3 to 6 as well as the 0LSE will be discussed in Chapter 8.

### 7.2 TLS Technique

### 7.2.1 Assumptions and notation

1. The Total Least Square Estimate (TLSE) of $\beta$ will be denoted by $\hat{\beta}_{\text {TLS }}$.
2. The SVD of the augmented matrix [ X Y] is given in (1.3.8):

$$
\left[\begin{array}{ll}
\mathrm{X} & \mathrm{Y}
\end{array}\right]=\tilde{\mathrm{U}} \tilde{\mathrm{U}} \tilde{\mathrm{~V}}^{\prime}
$$

$$
\text { with } \begin{align*}
\tilde{\mathrm{U}} & =\left[\tilde{\mathrm{u}}_{1} \ldots \tilde{\mathrm{u}}_{\mathrm{p}+1}\right], \tilde{\mathrm{u}}_{\mathrm{i}}: \mathrm{nx} 1 \\
\tilde{\mathrm{~V}} & =\left[\tilde{\mathrm{v}}_{1} \ldots \tilde{\mathrm{v}}_{\mathrm{p}+1}\right], \tilde{\mathrm{v}}_{\mathrm{i}}:(\mathrm{p}+1) \times 1 \\
\tilde{\mathrm{U}}^{\prime} \tilde{\mathrm{U}} & =\tilde{\mathrm{V}}^{\prime} \tilde{\mathrm{V}}=\mathrm{I}_{\mathrm{p}+1} \\
\tilde{\Delta} & =\operatorname{diag}\left[\omega_{1}, \omega_{2}, \ldots, \omega_{\mathrm{p}+1}\right] \\
\omega_{1} & \geq w_{2} \geq \ldots \ldots \geq w_{\mathrm{p}+1}
\end{align*}
$$

For simplicity, assume $w_{1}>w_{2}>\ldots \ldots>w_{p+1}$.
3. $\tilde{v}_{i}, j$ is the $j$ th component of the $i$ th right singular vector $\tilde{v}_{i}$ the i- th column of $\tilde{V}$
4. $\tilde{v}_{i}^{0}$ is the p-dimensional vector containing the first $p$ components of the i-th right singular vector $\tilde{v}_{i}$ of [X Y], so we have

$$
\tilde{v}_{i}=\left[\begin{array}{ll}
\tilde{v}_{i}^{0} & \tilde{v}_{i}, \mathrm{p}+1 \tag{7.2.2}
\end{array}\right]^{\prime}
$$

5. $[\mathrm{X} \hat{\mathrm{Y}}]$ is the LS approximation of $[\mathrm{X} Y]$ with $\hat{\mathrm{Y}}$ the orthogonal projection of $Y$ onto the column space $R(X)$ of $X$. [ $\tilde{X} \tilde{Y}]$ is the TLS approximation of $\left[\begin{array}{ll}\mathrm{X} & \mathrm{Y}\end{array}\right]$.
6. The following theorem is due to Van Huffel and Vandewalle (1987) and is useful in the generalization of TLS:

If $1 \leq j \leq p+1$ then
(a) $\tilde{\mathrm{v}}_{\mathrm{j}, \mathrm{p}^{+1}}=0 \Rightarrow \mathrm{Y}_{\perp} \tilde{u}_{\mathrm{j}}$ and $\tilde{\mathrm{Y}}_{\mathrm{A}} \tilde{\mathrm{u}}_{\mathrm{j}}$
(when $\tilde{v}_{j, p+1}=0$, then $\tilde{u}_{j} \omega_{j} \tilde{v}_{j}^{\prime}=\omega_{j}\left[\tilde{u}_{j} \tilde{v}_{j}^{0} 0\right]$ )
(b) The eigenvalues of $X^{\prime} X$ interlace with those of $[\mathrm{X} \mathrm{Y}]^{\prime}\left[\begin{array}{l}\mathrm{X} \\ \mathrm{Y}\end{array}\right]$ as follows:

$$
\left.\omega_{1}^{2} \geq \lambda_{1} \geq w_{2}^{2} \geq \lambda_{2} \geq \cdots \geq \omega_{\mathrm{p}}^{2} \geq \lambda_{\mathrm{p}} \geq \omega_{\mathrm{p}+1}^{2} \quad \text { (fromA. } 3\right)
$$

and if we assume that the singular values of $X$ are disjoint. Then

$$
\begin{align*}
& \tilde{\mathrm{v}}_{\mathrm{j}, \mathrm{p}+1}=0 \Rightarrow \omega_{\mathrm{j}}^{2}=\lambda_{\mathrm{i}}, \tilde{\mathrm{u}}_{\mathrm{j}}= \pm \mathrm{u}_{\mathrm{i}}, \text { and } \tilde{\mathrm{v}}_{\mathrm{j}}=\left[ \pm \mathrm{v}_{\mathrm{i}}{ }^{\prime} 0\right]^{\prime} \\
& \quad \text { with i } \in\{\mathrm{j}-1, \mathrm{j}\}, 1 \leq \mathrm{i} \leq \mathrm{p},  \tag{7.2.4}\\
& \tilde{\mathrm{v}}_{\mathrm{j}, \mathrm{p}+1}=0 \Rightarrow \hat{\mathrm{Y}}_{\perp} \tilde{\mathrm{u}}_{\mathrm{j}} \tag{7.2.5}
\end{align*}
$$

7. The correction vector $\epsilon=\mathrm{Y}-\hat{\mathrm{Y}}$ is the 0LS approximation error and the correction matrix $[\Delta \tilde{X} \Delta \tilde{Y}]=[\mathrm{X}-\tilde{\mathrm{X}}, \mathrm{Y}-\tilde{Y}]$ the TLS approximation error.
8. By manipulation the pseudo-equation $Y=X \beta$ can be written as (Golub and Van Loan (1980)) :

$$
\begin{align*}
\mathrm{X} \beta-\mathrm{Y} & =0 \\
{[\mathrm{X}} & \mathrm{Y}]\left[\begin{array}{r}
\beta \\
-1
\end{array}\right]
\end{align*}=0
$$

9. The following result of Eckart and Young (1936) is used in the development of the TLS technique:

Let $r(X)=k \leq p \leq n$ and let $\hat{X}$ be an approximation to $X$ that satisfies

$$
\begin{align*}
r(\hat{X}) & \leq k  \tag{7.2.7}\\
\|\hat{X}-X\| & =\inf \|(\bar{X} \quad-X)\|  \tag{7.2.8}\\
& r(\bar{X}) \leq k
\end{align*}
$$

where $\|\cdot\|$ is a unitarily invariant matrix norm as defined in (1.9.10).

Set

$$
{\sqrt{\lambda_{k+1}}}=\sqrt{\lambda_{k}+2}=\ldots=\sqrt{\lambda_{p}}=0
$$

then by using the $\operatorname{SVD}(1.3 .1), \hat{X}$ can be written as

$$
\begin{equation*}
\hat{\mathrm{X}}=\mathrm{UD} \hat{D}_{a} V^{\prime} \tag{7.2.9}
\end{equation*}
$$

where

$$
\hat{\mathrm{D}}_{a}=\operatorname{Diag}\left[\sqrt{\lambda_{1}}, \sqrt{\lambda_{2}}, \ldots \sqrt{\lambda_{\mathrm{k}}}, 0, \ldots 0\right]
$$

In the case of the Frobenius norm, Eckart and Young (1936) showed that the $\hat{X}$ of (7.2.9) satisfies (7.2.7) and (7.2.8). For a generalization of the Eckart and Young approximation see Golub, Hoffman and Stewart (1987).

### 7.2.2 Geometric view of TLS

The figures below are from Van Huffel and Vandewalle (1985):

(a) The LS solution is obtained by projecting $Y$ orthogonally onto $R(X)$ and solving $\mathrm{X} \hat{\beta}=\mathrm{Y}$.
(b) The TLS solution is obtained by approximating the columns $X_{i}$ of $X$ and $Y$ by $X_{i}$ and $\hat{Y}$ until $\hat{Y}$ is in the space $R(\hat{X})$ generated by the columns $\hat{X}_{i}$ and solving $\hat{X} \beta=\hat{\mathrm{Y}}$.

### 7.2.3 Definition

Given an over-determined set of $n$ linear equations $X \beta=Y$ in $p$ unknowns, the total least-squares (TLS) solution is the minimum-norm solution $\hat{\beta}_{\text {TLS }}$ of the set of $n$ linear equations

$$
\begin{equation*}
\tilde{\mathrm{x}} \hat{\mathrm{P}}_{\mathrm{TLS}}=\tilde{\mathrm{Y}} \tag{7.2.10}
\end{equation*}
$$

where $\tilde{X}$ and $\tilde{Y}$ are determined such that

$$
\begin{align*}
\tilde{Y} & \in \mathbb{R}(\tilde{X})  \tag{7.2.11}\\
\|\Delta X \Delta Y\|_{F} & =\|[X Y]-[\tilde{X} \tilde{Y}]\|_{F} \quad \text { is minimal } \tag{7.2.12}
\end{align*}
$$

Using the Eckart-Young principle, equations (7.2.11) and (7.2.12) will be satisfied by making $\omega_{p+1}$ in (7.2.1) zero.

The TLS approximation of $\tilde{\Delta}$ defined in (7.2.1) can then itself be approximated by

$$
\hat{\tilde{\Delta}}=\operatorname{diag}\left[\omega_{1}, \omega_{2}, \ldots, \omega_{p}, 0\right]
$$

and

$$
[\tilde{\mathrm{X}} \tilde{\mathrm{Y}}]=\tilde{\mathrm{U}} \tilde{\mathrm{D}} \tilde{\mathrm{~V}}^{\prime}
$$

The lower rank approximation for (7.2.6) is then

$$
\left[\begin{array}{ll}
\tilde{\mathrm{X}} & \tilde{\mathrm{Y}}
\end{array}\right]\left[\begin{array}{ll}
\hat{\beta}_{\mathrm{TLS}}^{\prime} & -1]^{\prime}=0 \quad(7.2 .13)
\end{array}\right.
$$

The TLS solution is obtained by scaling the last column $\tilde{v}_{p+1}$ of $\tilde{V}$ so that its last component $\tilde{v}_{p+1, p+1}=-1$. It is easy to verify that

$$
\begin{equation*}
\left[\hat{\beta}_{\mathrm{TLS}}^{\prime}-1\right]^{\prime}=-\tilde{\mathrm{v}}_{\mathrm{p}+1} / \tilde{\mathrm{v}}_{\mathrm{p}+1}, \mathrm{p}+1 \tag{7.2.14}
\end{equation*}
$$

by showing

$$
\begin{aligned}
& {\left[\begin{array}{ll}
\tilde{X} & \tilde{Y}
\end{array}\right]\left[-\tilde{\mathrm{v}}_{\mathrm{p}+1} / \tilde{\mathrm{v}}_{\mathrm{p}+1}, \mathrm{p}+1\right]=\tilde{\mathrm{U}} \tilde{\Delta} \tilde{V}^{\prime}\left[-\tilde{\mathrm{v}}_{\mathrm{p}+1} / \tilde{\mathrm{v}}_{\mathrm{p}+1, \mathrm{p}+1}\right]} \\
& =\tilde{U} \tilde{\Delta}\left[\tilde{v}_{1} \ldots \tilde{v}_{p+1}\right]^{\prime}\left[-\tilde{v}_{p+1} / \tilde{v}_{p+1}, p+1\right] \\
& =\tilde{U} \tilde{\Delta}\left[0,0, \ldots, 0,-1 / \tilde{\mathrm{v}}_{\mathrm{p}+1}, \mathrm{p}+1\right]^{\prime} \\
& \text { (the columns of } \tilde{V} \text { are orthogonal) } \\
& \left.=0: \operatorname{nx} 1 \quad\left(\tilde{\Delta}\left[0,0, \ldots, 0,-1 / \tilde{\mathrm{v}}_{\mathrm{p}+1}, \mathrm{p}+1\right]^{\prime}=0\right)\right)
\end{aligned}
$$

If we transpose (7.2.14) and multiply each side by $[\mathrm{X} Y]^{\prime}[\mathrm{X} Y$, we obtain

$$
\left.\begin{array}{rl}
{\left[\begin{array}{ll}
X^{\prime} X & X^{\prime} Y \\
\mathrm{Y}^{\prime} \mathrm{X} & \mathrm{Y}^{\prime} \mathrm{Y}
\end{array}\right]\left[\begin{array}{c}
\hat{\beta}_{\mathrm{TLS}}
\end{array}\right]} & =\tilde{V}^{2} \tilde{\mathrm{~V}}^{\prime}\left[-\tilde{\mathrm{v}}_{\mathrm{p}+1} / \tilde{\mathrm{v}}_{\mathrm{p}+1}, \mathrm{p}+1\right.
\end{array}\right] \quad(\text { by using }(7.2 .1))
$$

Using the SVD (1.3.1) of X this estimator can be written as

$$
\begin{align*}
\hat{\beta}_{\mathrm{TLS}} & =\left[V \Delta^{2} V^{\prime}-\omega_{\mathrm{p}_{+1} \mathrm{I}}^{2} I\right]^{-1} V \Delta U^{\prime} Y \\
& =V\left[\Delta^{2}-\omega_{p_{+1}}^{2} I\right]^{-1} V^{\prime} V \Delta U^{\prime} Y \\
& =\sum_{j=1}^{p} v_{j}\left(\lambda_{j}-\omega_{p_{+1}}^{2}\right)^{-1} \sqrt{\lambda_{j}} u_{j}^{\prime} Y \\
& =\sum_{j=1}\left(\lambda_{j}-\omega_{\mathrm{p}+1}^{2}\right)^{-1} \sqrt{\lambda_{j}}\left(u_{j}^{\prime} Y\right) v_{j} \tag{7.2.16}
\end{align*}
$$

where $u_{j}^{\prime} Y$ is a scalar. Equation (7.2.16) is applicable when $\omega_{p}>\omega_{p+1}$ and $\tilde{\mathrm{v}}_{\mathrm{p}+1}, \mathrm{p}+1 \neq 0$. If $\tilde{\mathrm{v}}_{\mathrm{p}+1}, \mathrm{p}+1=0$, (7.2.14) become infinite and Golub and Van Loan's (1980) algorithm for TLS fails to compute a finite TLS solution. Consequently such problems can not be solved with TLS. However Van Huffel and Vandewalle (1987) claimed that the TLS approximation [X Y] can still be determined by making the next smallest singular value in (7.2.1) zero.

Note that the authors definition of non-predictive collinearities is: $\tilde{v}_{j},{ }_{p+1}(j=k, \ldots, p+1)$ are exactly zero and $w_{j}$ is suitably small; in contrast to the definition 6.1 (below (6.1.5) of the Chapter 6) where $\tilde{v}_{j},{ }_{p+1}(j=k, \ldots, p+1) \approx 0$.

### 7.2.4 Generalized definition

In generalizing of TLS to cases where $\omega_{p+1}$ coincides with other singular values $\omega_{\mathrm{k}-1}>\omega_{\mathrm{k}}=\ldots \omega_{\mathrm{p}}=\omega_{\mathrm{p}+1}$, and some or all of the $\tilde{\mathrm{v}}_{\mathrm{j}, \mathrm{p}+1}(\mathrm{j}=\mathrm{k}, \ldots, \mathrm{p}+1)$ are zero, Van Huffel and Vandewalle (1987) give two definitions (their equations 2.8 and 2.9):

Definition one: If $\omega_{\mathrm{p}+1}$ coincides with other singular values $\omega_{\mathrm{k}-1}>w_{\mathrm{k}}=\ldots \omega_{\mathrm{p}}=\omega_{\mathrm{p}+1}$, then $\left[\hat{\beta}_{\mathrm{TLS}}^{\prime}-1\right]^{\prime}$ is a linear combination of the corresponding right singular vectors $\tilde{v}_{k} \ldots \tilde{\mathrm{v}}_{\mathrm{p}+1}$ such that the solution $\hat{\beta}_{\text {TLS }}$ has minimum norm.

If not all components $\tilde{v}_{j},_{p+1}(j=k, \ldots, p+1)$ are zero the solution is

$$
\begin{equation*}
\hat{\beta}_{T L S}=\sum_{j=k}^{p+1} \varrho_{j} \tilde{v}_{j}^{0} \quad \text { with } \varrho_{j}=-\frac{\tilde{v}_{j}, p^{+1}}{\sum_{j=k}^{p+1} \tilde{v}_{j}^{2},_{p+1}} \tag{7.2.17}
\end{equation*}
$$

When $\mathrm{k}=\mathrm{p}+1,(7.2 .17)$ reduces to

$$
\begin{aligned}
\hat{\beta}_{\mathrm{TLS}} & =\varrho_{\mathrm{p}+1} \tilde{\mathrm{v}}_{\mathrm{p}+1}^{0} \quad \text { with } \varrho_{\mathrm{p}+1}=-\frac{\tilde{\mathrm{v}}_{\mathrm{p}+1}, \mathrm{p}+1}{} \\
& =-\tilde{\mathrm{v}}_{\mathrm{p}+1}^{0} / \tilde{v}_{\mathrm{p}+1}, \mathrm{p}+\mathrm{p}+1
\end{aligned}
$$

When $r\left(\left[\begin{array}{ll}X & Y\end{array}\right]\right)=p+1-r(r$ singular values are near zero) we will assume $w_{\mathrm{p}+1-\mathrm{r}}>\psi_{\mathrm{p}+1-\mathrm{r}+1}=\ldots w_{\mathrm{p}+1-1}=\psi_{\mathrm{p}+1}, \quad$ and $\mathrm{k}=(\mathrm{p}+1)-(\mathrm{r}-1)=\mathrm{p}+2-\mathrm{r}$.

If, however, all r collinearities which appear in $[\mathrm{X} Y$ ] are nonpredictive, then TLS approximates [X Y] with [ $\tilde{\mathrm{X}} \tilde{\mathrm{Y}}]$ by making the next larger singular value $w_{\mathrm{p}+1-\mathrm{r}}$ zero. Then by using equation (7.2.3) to (7.2 5) the TLS estimator is given by

Definition tho: If $r>0, \tilde{v}_{j}{ }_{p+1}(j=p+2-r, \ldots, p+1)$ are exactly zero (hence by using (7.2.3) to (7.2.5) we know that $Y$ is orthogonal to $\tilde{u}_{j}$ ) and with $w_{p-r}>w_{p+1-r}$, we have:

$$
\begin{equation*}
\hat{\beta}_{T L S}=\sum_{j=1}^{p-r}\left(\lambda_{j}-\omega_{p+1-r}^{2}\right)^{-1} \sqrt{\lambda_{j}}\left(u_{j}^{\prime} Y\right) v_{j} \tag{7.2.18}
\end{equation*}
$$

with $r$ the number of $r$ non- predictive collinearities.

From a computational point of view definition 6.1 would be more practical, and in applying (7.2.18) the $\tilde{v}_{j},{ }_{p+1}(j=p+2-r, \ldots, p+1)$ are exactly zero only when known a priori or assumed and set to 0 .

### 7.2.5 Properties of TLS estimator

For simplicity we work with (7.2.18):

$$
\begin{aligned}
\hat{\beta}_{T L S} & =\sum_{\substack{j=1 \\
p-r}}^{p-r}\left(\lambda_{j}-w_{p+1-r}^{2}\right)^{-1} \sqrt{\lambda_{j}}\left(u_{j}^{\prime} Y\right) v_{j} \\
& =\sum_{j=1}^{\sum}\left(\lambda_{j}-w_{p+1-r}^{2}\right)^{-1} c_{j} v_{j}, \text { where } c_{j}=\sqrt{\lambda_{j}}\left(u_{j}^{\prime} Y\right)
\end{aligned}
$$

The following properties are reported by Van Huffel and Vandewalle (1985, p15):

$$
\begin{equation*}
\mathrm{E}\left(\hat{\beta}_{\mathrm{TLS}}\right)=\beta+\omega_{\mathrm{p}+1-\mathrm{r}}^{2} \sum_{\mathrm{j}=1}^{\mathrm{p}-\mathrm{r}}\left(\lambda_{\mathrm{j}}-\omega_{\mathrm{p}+1-\mathrm{r}}^{2}\right)^{-1}\left(\mathrm{v}_{\mathrm{j}}^{\prime} \beta\right) \mathrm{v}_{\mathrm{j}}-\sum_{\mathrm{j}=\mathrm{t}}^{\mathrm{p}}\left(\mathrm{v}_{\mathrm{j}}^{\prime} \beta\right) \mathrm{v}_{\mathrm{j}} \tag{7.2.19}
\end{equation*}
$$

where $\mathrm{t}=\mathrm{p}-\mathrm{r}+1$

$$
\begin{align*}
\operatorname{Var}\left(\hat{\beta}_{\mathrm{TLS}}\right)= & \sigma^{2} \sum_{j=1}^{\mathrm{p}-\mathrm{r}} \lambda_{\mathrm{j}}\left(\lambda_{\mathrm{j}}-w_{\mathrm{p}+1-\mathrm{r}}^{2}\right)^{-2} \mathrm{v}_{\mathrm{j}} \mathrm{v}_{\mathrm{j}}^{\prime}  \tag{7.2.20}\\
\operatorname{MSE}\left(\hat{\beta}_{\mathrm{TLS}}\right)= & \sigma^{2} \sum_{j=1}^{\mathrm{p}-\mathrm{r}} \lambda_{\mathrm{j}}\left(\lambda_{\mathrm{j}}-\omega_{\mathrm{p}+1-\mathrm{r}}^{2}\right)^{-2} \mathrm{v}_{\mathrm{j}} \mathrm{v}_{\mathrm{j}}^{\prime}+ \\
& \quad \omega_{\mathrm{p}+1-\mathrm{r}}^{4} \sum_{\mathrm{j}=1}^{\mathrm{p}-\mathrm{r}}\left(\lambda_{\mathrm{j}}-\omega_{\mathrm{p}+1-\mathrm{r}}^{2}\right)^{-2}\left(\mathrm{v}_{\mathrm{j}}^{\prime} \beta\right)^{2}-\sum_{j=\mathrm{t}}^{\mathrm{p}}\left(\mathrm{v}_{\mathrm{j}}^{\prime} \beta\right)^{2} \tag{7.2.21}
\end{align*}
$$

The above properties ( $(7.2 .19)$, (7.2.20) and (7.2.21)) were given without a proof. Van Huffel (pers. comm, 1990) confirms that her proofs assume the following:
(i) $\omega_{\mathrm{p}+1-\mathrm{r}}^{2}$ is fixed and thus the expectation of the TLSE is the same as that of the RRE with $k=-\omega_{\mathrm{p}_{+1}-\mathrm{r}}^{2}$.
(ii) Furthermore $\sum^{\mathrm{p}}\left(\lambda_{\mathrm{j}}-\omega_{\mathrm{p}+1-\mathrm{r}}^{2}\right)^{-1} \sqrt{\lambda_{j}}\left(\mathrm{u}_{\mathrm{j}}^{\prime} \mathrm{Y}\right) \mathrm{v}_{\mathrm{j}}=0$, where $\mathrm{t}=\mathrm{p}+1-\mathrm{r}$, on the grounds that TLS and $\stackrel{\mathrm{j}}{\mathrm{t}}$ (

Statisticians may reject the expectation of the TLSE (7.2.19) and what follows firstly as $\psi_{\mathrm{p}+1-\mathrm{r}}^{2}$ is not a constant but a function of Y , and therefore random. Furthermore the quantity $\sum_{j=t}^{p}\left(\lambda_{j}-\omega_{p_{+1}-r}^{2}\right)^{-1} \sqrt{\lambda_{j}}\left(u_{j}^{\prime} Y\right) v_{j}$, is actually infinite (division by 0 ).

Van Huffel (pers. comm, 1990) comments 'my main motivation for the expressions ( 7.2 .19 ), (7.2.20) and (7.2.21)) is to give the reader the
feeling that TLS is not appropriate for estimation of parameters in models with error-free predictor variables. TLS is not devised as a biased regression estimator like $R R$ (although close similarities between the expressions of both estimators exist) but is only appropriate for parameter estimation in models with errors in all variables.'

Ve may further observe that from a computational point of view if $X^{\prime} X$ is ill-conditioned $\left[X^{\prime} X-\omega_{p+1}^{2}\right]^{-1}$ will be worse-conditioned. TLS is therefore an alternative estimation procedure to handle the errors-in-variables phenomenon rather than collinearity. As a device it approximates $X$ with $\tilde{X}$ generally of lower rank and imposes restrictions (7.2.10) in place of the normal equations. If there are in fact no errors in the variables then (7.2.15) implies that $\hat{\beta}_{\text {TLS }}$ would be biased and have a larger variance than the OLSE $\hat{\beta}$.

### 7.3 Use of TLS in the errors- in-variables model

### 7.3.1 The Model

The general errors-in-variables model is (Van Huffel and Vandewalle (1985)):

$$
\begin{align*}
Y_{0} & =X_{0} \beta_{0}+\epsilon \\
X & =X_{0}+X_{e}  \tag{7.3.1}\\
Y & =Y_{0}+Y_{e}
\end{align*}
$$

where
$\beta_{0}: \mathrm{px} 1$ is the unobservable vector of parameters to be estimated
$X$ : nxp is the observed matrix on the $p$ independent variables
Y : nx1 is the observed response vector
$Y_{0}, X_{0}$ contain the true but unobservable variables, and
$X_{e}$ and $Y_{e}$ are their measurement or observation errors (unobserved).

In addition to the errors in the variables ( $\mathrm{X}_{\mathrm{e}}, \mathrm{Y}_{\mathrm{e}}$ ) there is an error vector $\epsilon$ in the equation. In the classical or pure error-in-variables model $\epsilon=0$.

It is assumed that the error variances of $\left[X_{e}\right]_{i j}$ and $\left[Y_{e}\right]_{i}$ and their covariances $\mathrm{E}\left(\left[\mathrm{X}_{\mathrm{e}}\right]_{\mathrm{ij}},\left[\mathrm{Y}_{\mathrm{e}}\right]_{\mathrm{i}}\right)$ are known while $\beta_{0}$ and the variance $\sigma_{\epsilon}^{2}$, of the equation, are unknown.

Define the $n x(p+1)$ error matrix as

$$
\Delta_{\mathrm{e}}=\left[\begin{array}{ll}
\mathrm{X}_{\mathrm{e}} & \mathrm{Y}_{\mathrm{e}}+\epsilon \tag{7.3.2}
\end{array}\right]
$$

Denote the $i$ - th row of $\Delta_{e}$ by $\left[\Delta_{\mathrm{e}}\right]_{i}^{\prime}$ and partition $X_{e}$ by its rows and denote each row by a subscript. Thus $X_{e}=\left[\left(X_{e}\right)_{1},\left(X_{e}\right)_{2}, \ldots,\left(X_{e}\right)_{n}\right]^{\prime}$ and $\left(X_{e}\right)_{i}^{\prime}$ is a 1 xp row vector. This implies that $\left[\Delta_{\mathrm{e}}\right]_{\mathrm{i}}^{\prime}=\left[\left(\mathrm{X}_{\mathrm{e}}\right)_{\mathrm{i}}^{\prime} \quad\left[\left(\mathrm{Y}_{\mathrm{e}}\right)_{\mathrm{i}}+\epsilon_{\mathrm{i}}\right]\right]$, where $\left(Y_{e}\right)_{i}+\epsilon_{i}$ is the $i$-th element of the sum of $Y_{e}$ and $\epsilon$. Assume the following
(i) $\Delta_{\mathrm{e}}$ and X are stochastically independent random matrices
(ii) row vectors of $\Delta_{\mathrm{e}}$ are stochastically independent and identically distributed
(iii) $E\left(\Delta_{e}\right)=0$

Define the covariance matrix of the error variables as

$$
\begin{align*}
& \Sigma_{\Delta_{\mathrm{e}}}=\mathrm{E}\left[\left[\Delta_{\mathrm{e}}\right]_{\mathrm{i}}\left[\Delta_{\mathrm{e}}\right]_{\mathrm{i}}^{\prime}\right] \quad(:(\mathrm{p}+1) \mathrm{x}(\mathrm{p}+1) \text { matrix })  \tag{7.3.3}\\
& =E\left[\begin{array}{ll}
\left(X_{e}\right)_{i}\left(X_{e}\right)_{i}^{\prime} & \left(X_{e}\right)_{i}\left[\left(Y_{e}\right)_{i}+\epsilon_{i}\right] \\
{\left[\left(Y_{e}\right)_{i}+\epsilon_{i}\right]\left(X_{e}\right)_{i}^{\prime}} & {\left[\left(Y_{e}\right)_{i}+\epsilon_{i}\right]\left[\left(Y_{e}\right)_{i}+\epsilon_{i}\right]}
\end{array}\right]
\end{align*}
$$

Denote
$E\left[\left(X_{e}\right)_{i}\left(X_{e}\right)_{i}^{\prime}\right]=\Sigma_{X_{e}}(\operatorname{pxp}), E\left[\left(X_{e}\right)_{i}\left[\left(\mathrm{Y}_{\mathrm{e}}\right)_{\mathrm{i}}+\epsilon_{\mathrm{i}}\right]\right]=\sigma_{\mathrm{Y}_{\mathrm{e}} \mathrm{X}_{\mathrm{e}}}(\mathrm{px} 1), \quad \mathrm{E}\left[\epsilon_{\mathrm{i}}^{2}\right]=\sigma_{\epsilon}^{2}(1 \mathrm{x} 1)$
and $E\left[\left(Y_{e}\right)_{i}^{2}\right]=\sigma_{Y_{e}}^{2}(1 \times 1)$ and from the assumptions:

$$
\Sigma_{\Delta_{e}}=\left[\begin{array}{cc}
\Sigma_{\mathrm{x}_{\mathrm{e}}} & \sigma_{\mathrm{Y}_{\mathrm{e}} \mathrm{X}_{\mathrm{e}}}  \tag{7.3.4}\\
\sigma_{\mathrm{Y}_{\mathrm{e}} \mathrm{X}_{\mathrm{e}}}^{\prime} & \sigma_{\mathrm{Y}_{\mathrm{e}}}^{2}+\sigma_{\epsilon}^{2}
\end{array}\right]
$$

The term $\sigma_{Y_{e} X_{e}}$ depends on the error variables only and is assumed to be known, and $\left[\Sigma_{\Delta_{e}}\right]_{n+1},{ }_{n+1}$ depends on $\sigma_{\epsilon}^{2}$ (unknown).

The model (7.3.1) can be manipulated to give

$$
\begin{align*}
\mathrm{Y}_{0} & =\mathrm{X}_{0} \beta_{0}+\epsilon \\
\mathrm{Y}-\mathrm{Y}_{\mathrm{e}} & =\left(\mathrm{X}-\mathrm{X}_{\mathrm{e}}\right) \beta_{0}+\epsilon \\
\mathrm{Y} & =\mathrm{X} \beta_{0}-\mathrm{X}_{\mathrm{e}} \beta_{0}+\epsilon+\mathrm{Y}_{\mathrm{e}} \\
& =\mathrm{X} \beta_{0}+0 \tag{7.3.5}
\end{align*}
$$

where

$$
\begin{align*}
\theta & =\left(\mathrm{Y}_{\mathrm{e}}+\epsilon\right)-\mathrm{X}_{\mathrm{e}} \beta_{0} \\
& =-\Delta_{\mathrm{e}}\left[\beta_{0}^{\prime}-1\right]^{\prime} \quad \text { (from 7.3.2) } \tag{7.3.6}
\end{align*}
$$

This is not a simple regression situation: $X_{e}$ is a random matrix and it is correlated with the error term 0 in contrast to the regression situation defined in (1.1).

### 7.3.2 Estimation by Least Squares

The 0LS estimation of $\beta_{0}$ (denoted by $\hat{\beta}_{0}$ ) in (7.3.5) will yield an inconsistent estimator for $\beta_{0}$ :

$$
\begin{align*}
\hat{\beta}_{0} & =\left(X^{\prime} X\right)^{-1} X^{\prime} Y \\
& =\beta_{0}+\left(X^{\prime} X\right)^{-1} X^{\prime} 0 \\
& =\beta_{0}+\left(\left(X^{\prime} X\right) / n\right)^{-1}\left(X^{\prime} \theta / n\right) \tag{7.3.7}
\end{align*}
$$

Schneeweiss (1976) shows that the following probability limits (plim) (defined in §1.11) exist and we will denote them by:

$$
\begin{align*}
\operatorname{plim}\left(n^{-1} X_{0} X_{0}\right) & =\mathbb{M}_{X_{0}} \\
\operatorname{plim}\left(n^{-1} \Delta_{e}^{\prime} \Delta_{e}\right) & =\mathbb{H}_{\Delta_{e}} \\
\operatorname{plim}\left(n^{-1} X_{0} \Delta_{e}\right) & =0 \\
\operatorname{plim}\left(n^{-1} X^{\prime} X\right) & =\mathbb{M}_{X_{0}}+\Sigma_{X_{e}}=\mathbb{M}_{X}  \tag{7.3.8}\\
\operatorname{plim}\left(n^{-1} X^{\prime} \Delta_{e}\right) & =\left[\Sigma_{X_{e}} \sigma_{Y_{e} X_{e}}\right]: \operatorname{px}(p+1) \\
\operatorname{plim}\left(n^{-1} X^{\prime} \theta\right) & =\left[\Sigma_{X_{e}} \sigma_{Y_{e} X_{e}}\right]\left[\beta_{0}^{\prime}-1\right]^{\prime} \\
\operatorname{plim}\left(n^{-1} \theta^{\prime} \theta\right) & =\beta_{0}^{\prime} \Sigma_{X_{e}} \beta_{0}-2 \sigma_{Y_{\mathrm{e}}} X_{e} \beta_{0}+\left(\sigma_{\epsilon}^{2}+\sigma_{Y_{e}}^{2}\right)=\sigma_{\theta}^{2}
\end{align*}
$$

Then the plim of $\hat{\beta}_{0}$ is:

$$
\begin{align*}
\operatorname{plim}\left(\hat{\beta}_{0}\right) & =\beta_{0}+\operatorname{plim}\left\{\left(\left(\mathrm{X}^{\prime} \mathrm{X}\right) / \mathrm{n}\right)^{-1}\left(\mathrm{X}^{\prime} 0 / \mathrm{n}\right)\right\} \\
& =\beta_{0}+\mathbb{M}_{\mathrm{X}}^{-1}\left[\Sigma_{\mathrm{X}_{\mathrm{e}}} \beta_{0}-\sigma_{\mathrm{Y}_{\mathrm{e}} \mathrm{X}_{\mathrm{e}}}\right] \quad(\text { from }(7.3 .8)) \tag{7.3.9}
\end{align*}
$$

Also the 0LS estimator of $\sigma_{\theta}^{2}$ is inconsistent:

$$
\begin{aligned}
(\mathrm{n}-\mathrm{p}) \hat{\sigma}_{\theta}^{2} & =\mathrm{Y}^{\prime}\left(\mathrm{I}-\mathrm{X}\left(\mathrm{X}^{\prime} \mathrm{X}\right)^{-1} \mathrm{X}^{\prime}\right) Y \\
& =\theta^{\prime}\left(\mathrm{I}-\mathrm{X}\left(\mathrm{X}^{\prime} X\right)^{-1} X^{\prime}\right) \theta
\end{aligned}
$$

Then the plim of $\hat{\sigma}_{\theta}^{2}$ is:

$$
\begin{align*}
\operatorname{plim}\left(\hat{\sigma}_{\theta}^{2}\right) & \left.=\sigma_{\theta}^{2}-\operatorname{plim}\left(\mathrm{n}(\mathrm{n}-\mathrm{p})^{-1}\left(0^{\prime} \mathrm{X} / \mathrm{n}\right)\left(\mathrm{X}^{\prime} \mathrm{X} / \mathrm{n}\right)^{-1}\left(\mathrm{X}^{\prime} 0\right] / \mathrm{n}\right)\right) \\
& =\sigma_{\theta}^{2}-\left[\Sigma_{\mathrm{X}_{\mathrm{e}}} \beta_{0}-\sigma_{\mathrm{Y}_{\mathrm{e}} \mathrm{X}_{\mathrm{e}}}\right]^{\prime} \mathrm{M}_{\mathrm{X}}^{-1}\left[\Sigma_{\mathrm{X}_{\mathrm{e}}} \beta_{0}-\sigma_{\mathrm{Y}_{\mathrm{e}} \mathrm{X}}\right] \tag{7.3.10}
\end{align*}
$$

From (7.3.9) and (7.3.10) it is clear that $\hat{\beta}_{0}$ is asymptotically biased and $\hat{\sigma}_{\theta}^{2}$ will always be asymptotically underestimated. When $\Sigma_{X_{e}}$ and $\sigma_{Y_{e} X_{e}}$ are known, this bias can be removed by the corrected least square (CLS) estimator (Schneeweiss (1976) equation 3.1), defined as:

$$
\begin{equation*}
\hat{\beta}_{C L S}=\left(X^{\prime} X-n \Sigma_{X_{e}}\right)^{-1}\left(X^{\prime} Y-n \sigma_{Y_{e} X_{e}}\right) \tag{7.3.11}
\end{equation*}
$$

Then

$$
\begin{aligned}
& \hat{\beta}_{\mathrm{CLS}}-\beta_{0}=\left(\mathrm{X}^{\prime} \mathrm{X} / \mathrm{n}-\Sigma_{\mathrm{X}_{\mathrm{e}}}\right)^{-1}\left(\mathrm{X}^{\prime} \mathrm{Y} / \mathrm{n}-\sigma_{\mathrm{Y}_{\mathrm{e}} \mathrm{X}_{\mathrm{e}}}-\left(\mathrm{X}^{\prime} \mathrm{X} / \mathrm{n}-\mathrm{\Sigma}_{\mathrm{X}_{\mathrm{e}}}\right) \beta_{0}\right) \\
& =\left(X^{\prime} \mathrm{X} / \mathrm{n}-\Sigma_{\mathrm{X}_{\mathrm{e}}}\right)^{-1}\left(\mathrm{X}^{\prime} \mathrm{X} \beta_{0} / \mathrm{n}+\mathrm{X}^{\prime} \theta / \mathrm{n}-\sigma_{\mathrm{Y}_{\mathrm{e}} \mathrm{X}_{\mathrm{e}}} \mathrm{X}^{\prime} \mathrm{X} \beta_{0} / \mathrm{n}+\Sigma_{\mathrm{X}_{\mathrm{e}}} \beta_{0}\right) \\
& \text { (from (7.3.5)) } \\
& =\left(X^{\prime} \mathrm{X} / \mathrm{n}-\Sigma_{\mathrm{X}_{\mathrm{e}}}\right)^{-1}\left(\mathrm{X}^{\prime}\left(-\Lambda_{\mathrm{e}}\left[\beta_{0}^{\prime}-1\right]^{\prime}\right) / \mathrm{n}+\Sigma_{\mathrm{X}_{\mathrm{e}}} \beta_{0}-\sigma_{\mathrm{Y}_{\mathrm{e}} \mathrm{X}_{\mathrm{e}}}\right) \\
& \text { (from (7.3.6)) } \\
& =\left(\mathrm{X}^{\prime} \mathrm{X} / \mathrm{n}-\Sigma_{\mathrm{X}_{\mathrm{e}}}\right)^{-1}\left(-\mathrm{X}^{\prime} \Delta_{\mathrm{e}}\left[\beta_{0}^{\prime}-1\right]^{\prime} / \mathrm{n}+\left[\Sigma_{\mathrm{x}_{\mathrm{e}}} \sigma_{\mathrm{Y}_{\mathrm{e}} \mathrm{X}_{\mathrm{e}}}\right]\left[\beta_{0}^{\prime}-1\right]^{\prime}\right) \\
& =\left(X^{\prime} X / n-\Sigma_{X_{e}}\right)^{-1}\left(\left(-X^{\prime} \Delta_{e} / n+\left[\Sigma_{X_{e}} \sigma_{Y_{e} X_{e}}\right]\right)\left[\beta_{0}^{\prime}-1\right]^{\prime}\right)
\end{aligned}
$$

and hence

$$
\begin{aligned}
\operatorname{plim}\left(\hat{\beta}_{\mathrm{CLS}}-\beta_{0}\right) & =\mathbb{Y}_{\mathrm{X}_{0}}^{-1}\left(\left(-\left[\Sigma_{\mathrm{x}_{\mathrm{e}}} \sigma_{\mathrm{Y}_{\mathrm{e}} \mathrm{X}_{\mathrm{e}}}\right]+\left[\begin{array}{ll}
\Sigma_{\mathrm{x}_{\mathrm{e}}} & \left.\left.\left.\sigma_{\mathrm{Y}_{\mathrm{e}}}\right]\right)\left[\beta_{0}^{\prime}-1\right]^{\prime}\right) \\
(\mathrm{from}(7.3 .8))
\end{array}\right.\right.\right. \\
& =0
\end{aligned}
$$

Thus $\hat{\beta}_{\text {cLS }}$ is a consistent estimator of $\beta_{0}$ (i.e. is asymptotically unbiased). A consistent estimator of $\sigma_{\theta}^{2}$ is

$$
\hat{\sigma}_{\mathrm{CLS}}^{2}=\hat{\theta}^{\prime} \hat{\theta} / \mathrm{n},
$$

where

$$
\begin{aligned}
\hat{\theta} & =\mathrm{Y}-\mathrm{X} \hat{\beta}_{\text {CLS }} \\
& \left.=\mathrm{X} \beta_{0}+0-X \hat{\beta}_{C L S} \quad \text { (by using }(7.3 .5)\right) \\
& =\mathrm{X}\left(\beta_{0}-\hat{\beta}_{\mathrm{CLS}}\right)+\theta
\end{aligned}
$$

From

$$
\begin{aligned}
\hat{\sigma}_{\mathrm{CLS}}^{2} & =\hat{\theta}^{\prime} \hat{\theta} / \mathrm{n} \\
& =\left[\left(\beta_{0}-\hat{\beta}_{\mathrm{CLS}}\right)^{\prime} \mathrm{X}^{\prime}+\theta^{\prime}\right]\left[\mathrm{X}\left(\beta_{0}-\hat{\beta}_{\mathrm{CLS}}\right)+\theta\right] / \mathrm{n} \\
& =\left[\left(\beta_{0}-\hat{\beta}_{\mathrm{CLS}}\right)^{\prime} \mathrm{X}^{\prime} \mathrm{X}\left(\beta_{0}-\hat{\beta}_{\mathrm{CLS}}\right)+2 \theta^{\prime} \mathrm{X}\left(\beta_{0}-\hat{\beta}_{\mathrm{CLS}}\right)+\theta^{\prime} \theta\right] / \mathrm{n}
\end{aligned}
$$

we have

$$
\begin{align*}
\operatorname{plim}\left(\hat{\sigma}_{\mathrm{CLS}}^{2}\right) & =\sigma_{\theta}^{2} \quad(\operatorname{from}(7.3 .8)) \\
& =\beta_{0}^{\prime} \Sigma_{\mathrm{X}_{\mathrm{e}}} \beta_{0}-2 \sigma_{\mathrm{Y}_{\mathrm{e}} \mathrm{X}_{\mathrm{e}}}^{\prime} \beta_{0}+\left(\sigma_{\epsilon}^{2}+\sigma_{\mathrm{Y}_{\mathrm{e}}}^{2}\right) \tag{7.3.12}
\end{align*}
$$

Hence,

$$
\left(\sigma_{\epsilon}^{2}+\sigma_{\mathrm{Y}_{\mathrm{e}}}^{2}\right)=\sigma_{\theta}^{2}-\beta_{0}^{\prime} \Sigma_{\mathrm{X}_{\mathrm{e}}} \beta_{0}+2 \sigma_{\mathrm{Y}_{\mathrm{e}} \mathrm{X}_{\mathrm{e}}^{\prime}} \beta_{0}
$$

The CLS estimate of ( $\sigma_{\epsilon}^{2}+\sigma_{\mathrm{Y}_{\mathrm{e}}}^{2}$ ) is:

$$
\begin{equation*}
\left(\hat{\sigma}_{\epsilon}^{2}+\hat{\sigma}_{\mathrm{Y}_{\mathrm{e}}}^{2}\right)_{\mathrm{CLS}}=\hat{\sigma}_{\mathrm{CLS}}^{2}-\hat{\beta}_{\mathrm{CLS}}^{\prime} \Sigma_{\mathrm{X}_{\mathrm{e}}} \hat{\beta}_{\mathrm{CLS}}+2 \sigma_{\mathrm{Y}_{\mathrm{e}}}^{\prime} \hat{\beta}_{\mathrm{CLS}} \tag{7.3.13}
\end{equation*}
$$

and Schneeweiss (1976) shows that this estimator is consistent. He also shows that

$$
\left(\hat{\sigma}_{\epsilon}^{2}\right)_{\mathrm{CLS}}=\left(\hat{\sigma}_{\mathrm{CLS}}^{2}-\hat{\beta}_{\mathrm{CLS}}^{\prime} \Sigma_{\mathrm{x}_{\mathrm{e}}} \hat{\beta}_{\mathrm{CLS}}+2 \sigma_{\mathrm{Y}_{\mathrm{e}} \mathrm{X}_{\mathrm{e}}} \hat{\beta}_{\mathrm{CLS}}\right)-\left(\hat{\sigma}_{\mathrm{Y}}^{2}\right)_{\mathrm{CLS}}
$$

estimates $\left(\sigma_{\epsilon}^{2}\right)$ consistently under the assumptions that $E\left(Y_{e}\right)=E(\epsilon)=0$, $\mathrm{E}\left(\mathrm{Y}_{\mathrm{e}} \epsilon^{\prime}\right)=0, \mathrm{E}\left(\mathrm{Y}_{\mathrm{e}} \mathrm{Y}_{\mathrm{e}}{ }^{\prime}\right)=\sigma_{\mathrm{Y}_{\mathrm{e}}}^{2} \mathrm{I}, \mathrm{E}\left(\epsilon \epsilon^{\prime}\right)=\sigma_{\epsilon}^{2} \mathrm{I}$, and $\sigma_{\mathrm{Y}_{\mathrm{e}}}^{2}$ is known.

An alternative formula for (7.3.13) is

$$
\begin{equation*}
\left(\hat{\sigma}_{\epsilon}^{2}+\hat{\sigma}_{Y_{\mathrm{e}}}^{2}\right)_{\mathrm{CLS}}=\mathrm{Y}^{\prime} \mathrm{Y} / \mathrm{n}-\left(\mathrm{Y}^{\prime} \mathrm{X} / \mathrm{n}-\sigma_{\mathrm{Y}_{\mathrm{e}}}{ }_{\mathrm{X}}\right) \hat{\beta}_{\mathrm{CLS}} \tag{7.3.14}
\end{equation*}
$$

When $\epsilon, \mathrm{Y}_{\mathrm{e}}, \mathrm{X}_{\mathrm{e}}$ and $\mathrm{X}_{0}$ are normally distributed, with the rows of $\Delta_{\mathrm{e}}$ and $\mathrm{X}_{0}$ are identically and independently distributed and under the assumptions already stated, then all these estimators are maximum likelihood estimators. (Johnston (1963)).

### 7.3.3 Estimation by Total Least Squares

The TLS estimator of $\beta$, defined in (7.2.15), is:

$$
\hat{\beta}_{\mathrm{TLS}}=\left[\mathrm{X}^{\prime} \mathrm{X}-\omega_{\mathrm{p}+1}^{2} \mathrm{I}\right]^{-1} \mathrm{X}^{\prime} \mathrm{Y}
$$

and the CLS estimator of $\beta$, defined in (7.3.11), is:

$$
\hat{\beta}_{C L S}=\left(X^{\prime} X-n \Sigma_{X_{e}}\right)^{-1}\left(X^{\prime} Y-n \sigma_{Y_{e}} X_{e}\right)
$$

Hence, TLS and CLS will yield the same consistent estimator if
(i) $\sigma_{\mathrm{Y}_{\mathrm{e}} \mathrm{x}_{\mathrm{e}}}=0$, (zero correlation between the measurement errors), and
(ii) $\Sigma_{x_{e}}=\sigma_{v}^{2} I_{p}$, and $\omega_{p_{+}}^{2} / n$ is a consistent estimator of $\sigma_{v}^{2}$, i.e. all error variables $\left(X_{e}\right)_{i}$ are stochastically independent and have equal variance $\sigma_{v}^{2}$, consistently estimated by $w_{p_{+1}}^{2} / n$.

Van Huffel and Vandewalle (1985, Theorem 4.3-1) prove that $\sigma_{v}^{2}$ is consistently estimated by $\omega_{\mathrm{p}+1}^{2} / \mathrm{n}$ if $\left(\sigma_{\epsilon}^{2}+\sigma_{\mathrm{Y}_{\mathrm{e}}}^{2}\right)=\sigma_{v}^{2}$. They use (7.3.14) to show that $\left(\sigma_{\epsilon}^{2}+\sigma_{\mathrm{Y}_{\mathrm{e}}}^{2}\right)$ can be consistently estimated by

$$
\begin{equation*}
\left(\sigma_{\mathrm{E}}^{2}+\sigma_{\mathrm{Y}_{\mathrm{e}}}^{2}\right)_{\mathrm{TLS}}=\omega_{\mathrm{p}_{+1}}^{2} / \mathrm{n} \tag{7.3.15}
\end{equation*}
$$

Their theorem also holds in case of $s$ non-predictive collinearities. The variance of the equation error $\sigma_{\epsilon}^{2}$ is always unknown, so a consistent
estimator is only possible is $\sigma_{\epsilon}^{2}=0$. Hence, TLS estimates consistently the parameters of a pure errors-in-variables model with known covariance matrix of the error variables. This condition implies that (7.3.4) will be known totally (all error variables are stochastically independent with equal variance $\sigma_{v}^{2}$ ):

$$
\Sigma_{\Delta_{\mathrm{e}}}=\left[\begin{array}{ll}
\sigma_{v}^{2} \mathrm{I}_{\mathrm{p}} & 0  \tag{7.3.16}\\
0 & \sigma_{v}^{2}
\end{array}\right]
$$

Since $\Sigma_{\mathrm{x}_{\mathrm{e}}}$ is known it is always possible to obtain independence and constancy of variance by transforming the data. This means that TLS always yields a consistent estimate of the parameters of a classical errors-invariables model provided the appropriate scaling of the data has been performed.

Furthermore, $\Sigma_{x_{e}}$ must only be known, up to a factor of proportionality.

The assumption of known error variance means that if the data is not scaled exactly or if the equation error is not zero, the TLS estimate may not be consistent $\left(\mathrm{E}\left(\sigma_{0}^{2}\right) \neq \omega_{\mathrm{p}+1}^{2} / n\right)$. Ketellapper (1983) investigated the impact of the violation of this assumption and found that TLS is then also preferable to 0LS, except when the error variance is highly overestimated.

For a discussion of the asymptotic properties of TLS in the error-invariable models see Van Huffel and Vandewalle (1985) and Schneeweiss (1976).

### 7.4 Subset Selection based on TLS for Prediction

In $\S 1.10$ subset selection via $Q R$ factorization with column pivoting was discussed, and this technique can be extended to TLS. Van Huffel and Vandewalle (1985) distinguish between three methods:

1. Algorithm SX-TLS (subset selection on $X$ where the method of estimation is TLS)

This method is a variant of SX-0LS (§1.10). The third step of Algorithm SX-OLS was to compute a OLS solution for the subset equation $X_{1} Z=Y$. If some of the variables are perturbed then a higher accuracy of the solution $Z$ and the predicted response $Y$ can be obtained by using TLS. Thus Algorithm SX-TLS is the same as SX-0LS except that the third step is replaced by:
built upe $Z \in R^{\hat{\mathrm{r}}}$ such that $\left\|\left[\hat{X}_{1}\right]_{\mathrm{e}} \hat{\mathrm{Y}}_{\mathrm{e}}\right\|_{\mathrm{F}}^{2}$ is minimized and $\hat{\mathrm{Y}}_{\mathrm{e}} \in \mathrm{R}\left(\hat{X}_{1}\right)$.
2. Algorithm SXY-TLS (subset selection on $\left[\begin{array}{ll}\mathrm{X} & \mathrm{Y}\end{array}\right]$ where the method of estimation is TLS)

Given $X: n x p, Y: n x 1$ and a method of computing an integer $r$ that approximates the rank of $[\mathrm{X} Y$ ], the following algorithm computes a permutation P and a vector z: $\hat{r} x 1$ such that the first $\hat{r}$ columns $X_{1}$ of $X P$ are independent and such that $\left\|\left[\hat{X}_{1}\right]_{\mathrm{e}} \hat{\mathrm{Y}}_{\mathrm{e}}\right\|_{\mathrm{F}}^{2}$ is minimized and $\hat{\mathrm{Y}} \in \mathbb{R}\left(\hat{\mathrm{X}}_{1}\right)$.

First step: Compute the SVD of $[\mathrm{X} \mathrm{Y}](1.3 .8)$ and determine $\hat{\mathrm{r}} \leq \operatorname{rank}[\mathrm{X} \mathrm{Y}]$, partition $\tilde{V}$

$$
\begin{gather*}
\tilde{\mathrm{V}}=\left[\begin{array}{cc}
\tilde{\mathrm{V}}_{11} & \tilde{\mathrm{~V}}_{12} \\
\tilde{\mathrm{~V}}_{21} & \tilde{\mathrm{~V}}_{22} \\
\tilde{\mathrm{~V}}_{1 \mathrm{y}} & \tilde{\mathrm{~V}}_{2 \mathrm{y}}
\end{array}\right]  \tag{7.4.1}\\
\underset{\mathrm{r}}{\hat{\mathrm{r}}} \mathrm{p}-\hat{\mathrm{r}}+\mathrm{\hat{r}} \\
1
\end{gather*}
$$

Second step: use QR with column piving to compute

$$
\begin{array}{r}
Q^{\prime}\left[\begin{array}{ll}
\tilde{\mathrm{V}}_{11}^{\prime} & \tilde{\mathrm{V}}_{21}^{\prime}
\end{array}\right] \mathrm{P}=\mathrm{Q}^{\prime}\left[\begin{array}{ll}
\ddot{\mathrm{V}}_{11}^{\prime} & \ddot{\mathrm{V}}_{21}^{\prime}
\end{array}\right]=\left[\begin{array}{ll}
\mathrm{R}_{11} & \mathrm{R}_{12}
\end{array}\right] \text { and } \operatorname{set} \mathrm{XP}=\begin{array}{ll}
{\left[\begin{array}{ll}
\mathrm{X}_{1} & \mathrm{X}_{2}
\end{array}\right]} \\
\hat{\mathrm{r}} & \mathrm{p}-\hat{\mathrm{r}}
\end{array}
\end{array}
$$

Third step: determine the TLS solution $z: \hat{r} x 1$ of $X_{1} Z=Y . \quad R_{11}$ is nonsingular and $\left\|\ddot{\mathrm{V}}_{1}^{-1}\right\|_{2}=\left\|\left[\mathrm{R}_{1}^{-1}\right]\right\|_{2}$.

The SXY-TLS algorihtm tends to maximize the $r$ - th singular value of $\left[X_{1} \quad Y\right]$. Indeed from (1.10.7) we have

$$
\begin{equation*}
\frac{w_{\hat{r}}}{\left\|\ddot{\mathrm{~V}}_{1}^{-1}\right\|_{2}} \leq \omega_{\hat{r}}\left(\left[X_{1} Y\right]\right) \leq \omega_{\hat{r}} \tag{7.4.2}
\end{equation*}
$$

The extra information provided by the variable $Y$ in this algorithm might $\hat{r}$
make the stable components of $Y$ predicted by SXY- TLS, i.e. $\Sigma \tilde{u}_{i}{ }^{\prime} Y$, superior $i=1$
to the stable components predicted by SX-LS $\Sigma \mathrm{u}_{\mathrm{i}}{ }^{\prime} \mathrm{Y}$, when errors are introduced into all variables of the model under consideration.
3. Algorithm SXY-VTLS (subset selection on $\left[\begin{array}{ll}X & Y\end{array}\right]$ where the method of estimation is TLS, with a variant)

It follows step one and two of SXY-TLS and then instead of computing step three, $z$ is obtained directly from step two as

$$
z=R_{i 1}^{-1} Q^{\prime} V_{1 y}
$$

(Here the authors assumed that the ( $\mathrm{n}-\hat{\mathrm{r}}+1$ ) smallest singular values of [X Y] are non-zero only due to perturbations.)

### 7.5 Summary

In this chapter TLS was introduced. The estimator was defined and its properties and use of for errors-in-variables models were discussed, along with subset selection based on TLS. It was observed that TLS is not an appropriate tool for dealing with collinearity.

## Chapter 8

## COMPARISON OF ESTIMATORS

In the previous chapters regression estimation techniques (0LS, PC, RR, GRR, LRR, FPC, SHE and TLS) have been defined, their properties discussed and some comparisons or remarks have been made. In this chapter we represent some of these remarks and properties in a unified manner in order to compare estimators and give some general direction on finding 'best' estimators. The estimators can be divided into two groups: The first group includes all those that operate on the regressor matrix $X$, namely OLS, PC, RR, GRR (and its variations), SH and FPC (and its variations): the second group consist of those that operate on the augmented matrix [X Y], namely LRR and TLS.

### 8.1 Basic Comparisons

In this section we present tabular summaries of the expectation and expected square error properties of the estimators. The tables 8.1 and 8.2 are separated largely for convenience of presentation, as the properties to which they refer are essentially interrelated. Table 8.3 presents some choices (mostly those used in Chapter 10) of $k, K$, and $d$.

Table 8.1 Expectation properties of estimators and references

| Type | Chapter | Definition | Bias | Variance |
| :---: | :---: | :---: | :---: | :---: |
| 0LS | 1 | $\hat{\beta}=\sum_{i=1}^{p} v_{i} c_{i} / \lambda_{i},$ | 0 | $\sigma^{2} \sum_{i=1}^{\mathrm{p}} \mathrm{v}_{\mathrm{i}} \mathrm{v}_{\mathrm{i}}^{\prime} / \lambda_{\mathrm{i}}$ |
|  | 3 | $\hat{\delta}=\left(Z^{\prime} Z\right)^{-1} Z^{\prime} Y$ | 0 | $\sigma_{i=1}^{2} \sum_{i=1} 1 / \lambda_{i}$ |
|  |  | $=V^{\prime} \hat{\beta}$ |  |  |
| PC | 3 | $\hat{\beta}_{P C}=\sum_{i=1}^{p-r} v_{i} c_{i} / \lambda_{i}$, | $-\mathrm{V}_{2} \mathrm{~V}_{2}{ }^{\prime} \beta$ | $\sigma_{i=1}^{\sigma^{2}-\mathrm{r}} \mathrm{v}_{\mathrm{i}} \mathrm{v}_{\mathrm{i}}^{\prime} / \lambda_{\mathrm{i}}$ |
| RR | 4.2 | $\hat{\beta}_{\mathrm{R}}=\sum_{i=1}^{p}\left(\lambda_{i}+\mathrm{k}\right)^{-1} \mathrm{c}_{\mathrm{i}} \mathrm{v}_{\mathrm{i}}$ | $-\mathrm{k}\left(\mathrm{X}^{\prime} \mathrm{X}+\mathrm{kI}\right)^{-1} \beta$ | $\sigma^{2} \mathrm{~V}\left[\Delta^{2}+\mathrm{kI}\right]^{-2} \Delta^{2} \mathrm{~V}^{\prime}$ |
| GRR | 4.6 | $\hat{\delta}_{\mathrm{K}}=\left(\Delta^{2}+\mathrm{K}\right)^{-1} \Delta^{2} \hat{\delta}$ | $-\left(\Delta^{2}+K\right)^{-1} \mathrm{~K} \delta$ | $\sigma^{2}\left(\Delta^{2}+K\right)^{-2} \Delta^{2}$ |
| AUORR | 4.7 | $\begin{aligned} & \hat{\delta}_{J W}=\left[\mathrm{I}-\left[\mathrm{kA}^{-1}\right]^{2}\right] \hat{\delta} \\ & \text { where } A=\mathrm{Z}^{\prime} \mathrm{Z}+\mathrm{K} \end{aligned}$ | $-\left[k A^{-1}\right]^{2} \delta$ | $\sigma^{2}\left[\mathrm{I}-\left[\mathrm{kA}{ }^{-1}\right]^{2}\right]^{2} \Delta^{-2}$ |
| AUGRR | 4.7 | $\left[\hat{\delta}_{0}\right]_{\mathrm{i}}=\frac{\left(\lambda_{\mathrm{i}} \hat{\delta}_{\mathrm{i}}^{2}+2 \hat{\sigma}^{2}\right) \lambda_{\mathrm{i}} \hat{\delta}_{\mathrm{i}}^{3}}{\left(\lambda_{\mathrm{i}} \hat{\delta}_{\mathrm{i}}^{2}+\hat{\sigma}^{2}\right)^{2}}$ | see (4.7.24) | and (4.7.25) |
| SH | 5 | $\hat{\beta}_{S H}=\mathrm{d} \sum_{i=1}^{\mathrm{p}} \mathrm{v}_{\mathrm{i}} \mathrm{c}_{\mathrm{i}} / \lambda_{\mathrm{i}}$ | - $(1-\mathrm{d}) \beta$ | $\sigma^{2} \mathrm{~d}^{2} \sum_{\mathrm{i}=1}^{\mathrm{p}} \mathrm{v}_{\mathrm{i}} \mathrm{v}_{\mathrm{i}}^{\prime} / \lambda_{\mathrm{i}}$ |
| FPC | 5 | $\hat{\delta}_{\text {FPC }}=\mathrm{F} \hat{\delta}$ | $-[\mathrm{I}-\mathrm{F}] \delta$ | $\sigma^{2} \mathrm{FD}^{-2} \mathrm{~F}$ |
|  | 5 | $\hat{\beta}_{\text {FPC }}=V^{\prime} V^{\prime} \hat{\beta}$ | $-\left[\mathbf{I}-\mathrm{VFV}^{\prime}\right] \beta$ | $\sigma^{2} \mathrm{VF} \Delta^{-2} \mathrm{FV}^{\prime}$ |
| LRR | 6 | $\hat{\beta}_{\text {LR }}=\sum_{i=1}^{p+1-s} \mathrm{f}_{\mathrm{i}} \tilde{\mathrm{V}}_{\mathrm{i}}^{0}$ | $-E\left[\sum_{i=t}^{p+1} f_{i} \tilde{v}_{i}^{0}\right]$ | - |
| TLS | 7 | $\hat{\beta}_{\mathrm{TLS}}=\sum_{\mathrm{j}=1}^{\mathrm{p}}\left(\lambda_{\mathrm{j}}-\omega_{\mathrm{p}+1}^{2}\right)^{-1} \mathrm{c}_{\mathrm{j}} \mathrm{v}_{\mathrm{j}}$ | - | - |

Table 8.2 Mean square errors of estimators

| Type | MSE | TMSE |
| :---: | :---: | :---: |
| 0LS | $\sigma^{2} \sum_{i=1}^{p} v_{i} v_{i}^{\prime} / \lambda_{i}$ | $\sigma^{2} \sum_{i=1}^{p} 1 / \lambda_{i}$ |
| PC | ${\sigma^{2}}_{\sum_{i=1}^{p-r} v_{i} v_{i}^{\prime} / \lambda_{i}+V_{2} V_{2}^{\prime} \beta \beta^{\prime} V_{2} V_{2}^{\prime}, ~}^{\prime}$ | $\sigma_{i=1}^{\mathrm{a}-\mathrm{r}} 1 / \lambda_{\mathrm{i}}+\beta^{\prime} V_{2} V_{2}^{\prime} \beta$ |
| RR | $\sigma^{2} \mathrm{~V}\left[\Delta^{2}+\mathrm{kI}\right]^{-2} \Delta^{2} \mathrm{~V}^{\prime}+\mathrm{k}^{2}\left(\mathrm{X}^{\prime} \mathrm{X}+\mathrm{kI}\right)^{-1} \beta \beta^{\prime}\left(\mathrm{X}^{\prime} \mathrm{X}+\mathrm{kI}\right)^{-1}$ | $\sum_{i=1}^{p} \frac{\sigma^{2} \lambda_{i}+k^{2} \delta_{i}^{2}}{\left(\lambda_{i}+k\right)^{2}}$ |
| GRR | $\sigma^{2}\left(\Delta^{2}+K\right)^{-2} \Delta^{2}+\left(\Delta^{2}+K\right)^{-1} \mathrm{~K} \delta \delta^{\prime} \mathrm{K}\left(\Delta^{2}+\mathrm{K}\right)^{-1}$ | $\sum_{\mathrm{i}=1}^{\mathrm{p}} \frac{\sigma^{2} \lambda_{\mathrm{i}}+\mathrm{k}_{\mathrm{i}}^{2} \delta_{\mathrm{i}}^{2}}{\left(\lambda_{\mathrm{i}}+\mathrm{k}_{\mathrm{i}}\right)^{2}}$ |
| AUORR | $\sigma^{2}\left[\mathrm{I}-\left[\mathrm{kA}^{-1}\right]^{2}\right]^{2} \Delta^{-2}+\left[\mathrm{kA}^{-1}\right]^{2} \delta \delta^{\prime}\left[\mathrm{kA}^{-1}\right]^{2}$ | $\sum_{i=1}^{p}\left\{\sigma^{2} \frac{\lambda_{i}\left(\lambda_{\mathrm{i}}+2 \mathrm{k}_{\mathrm{i}}\right)^{2}}{\left(\lambda_{\mathrm{i}}+\mathrm{k}_{\mathrm{i}}\right)^{4}}+\frac{\mathrm{k}_{\mathrm{i}}^{4} \delta_{\mathrm{i}}^{2}}{\left(\lambda_{\mathrm{i}}+\mathrm{k}_{\mathrm{i}}\right)^{4}}\right\}$ |
| AUGRR | see (4.7.24) and (4.7.25) |  |
| SH | $\sigma^{2} \mathrm{~d}^{2} \sum_{\mathrm{i}=1}^{\mathrm{p}} \mathrm{v}_{\mathrm{i}} \mathrm{v}_{\mathrm{i}}^{\prime} / \lambda_{\mathrm{i}}+(1-\mathrm{d})^{2} \beta \beta^{\prime}$ | $\sigma^{2} \mathrm{~d}^{2} \sum_{i=1}^{p} 1 / \lambda_{\mathrm{i}}+(1-\mathrm{d})^{2} \beta^{\prime} \beta$ |
| FPC | $\sigma^{2} \mathrm{~F} \Delta^{-2} \mathrm{~F}+[\mathrm{I}-\mathrm{F}] \delta \delta^{\prime}[\mathrm{I}-\mathrm{F}]$ | $\sigma^{2} \sum_{i=1}^{p} f_{i}^{2} / \lambda_{i}+\sum_{i=1}^{p}\left(1-f_{i}\right)^{2} \delta_{i}^{2}$ |
| LRR | $\approx \operatorname{MSE}\left(\hat{\beta}_{\mathrm{PC}}\right)$ | $\approx \operatorname{TMSE}\left(\hat{\beta}_{\mathrm{PC}}\right)$ |
| TLS | unknown | unknown |

Table 8.3 Choices of $\mathrm{k}, \mathrm{K}$, and d
type estimated value use for suggested by
$\mathrm{k}(\mathrm{HK}) \quad \mathrm{p} \hat{\sigma}^{2} / \hat{\beta}^{\prime} \hat{\beta}$
RRE Hoerl, Kennard and Baldwin
$k$ (LW) $\quad \mathrm{p} \hat{\sigma}^{2} / \sum_{i=1}^{p} \lambda_{i} \hat{\delta}_{i}^{2}$
RRE Lawless and Vang
$\mathrm{K}($ HKB $) \quad \hat{\mathbf{k}}_{\mathrm{i}}=\hat{\sigma}^{2} / \delta_{\mathrm{i}}^{2}$
GRRE Koerl, Kennard and Baldwin
$\mathrm{K}(\mathrm{T}) \quad \hat{\mathrm{k}}_{\mathrm{i}}=\lambda_{\mathrm{i}} /\left(\mathrm{F}_{\mathrm{i}}+1\right), \mathrm{F}_{\mathrm{i}}=\lambda_{\mathrm{i}} \hat{\delta}_{\mathrm{i}}^{2} / \hat{\sigma}^{2} \quad$ GRRE $\quad$ Troskie (1990)
$\mathrm{k} \quad \mathrm{p} \hat{\sigma}^{2} /\left(\sum_{\mathrm{i}=1}^{\mathrm{p}}\left[\hat{\delta}_{\mathrm{i}}^{2} /\left\{1+\sqrt{1+\lambda_{\mathrm{i}}\left(\hat{\delta}_{\mathrm{i}}^{2} / \hat{\sigma}^{2}\right)}\right\} \quad\right.\right.$ AUORR $\quad$ Nomura
d $\quad \hat{\beta}^{\prime} \hat{\beta} /\left\{\operatorname{tr}\left([\operatorname{var}(\hat{\beta})]+\hat{\beta}^{\prime} \hat{\beta}\right\} \quad\right.$ SHE *

* an estimated value that minimizes the TLSE of the SHE (5.1.10)


### 8.2 Judgement of estimators

### 8.2.1 Unbiasedness

In the class of unbiased estimators, the OLSE is the best linear estimator (BLUE) in the sense of minimum variance. In the presense of collinearity, the variance of OLSE can be inflated (due to small $\lambda_{i}$ 's) so that other biased estimators will be more suitable under a changed criterion, e.g. mimimum mean square error. If the collinearity between some of the regressors is however consistently continued in the prediction area, the effect of this collinearity on predictions will be less serious. If $\sigma^{2}$ is sufficiently small, $\beta$ can be estimated by OLS with good accuracy even if strong collinearities exist in $X$. Thus the choice of whether or not to use the OLSE should be based on the magnitudes of the $\lambda_{i}$ and unknown $\sigma^{2}$.

### 8.2.2. MSE criteria

Consider two competing estimators $b_{1}$ and $b_{2}$. If the matrix difference

$$
\begin{equation*}
S=\operatorname{MSE}\left(b_{2}\right)-\operatorname{MSE}\left(b_{1}\right) \tag{8.2.1}
\end{equation*}
$$

is positive semi-definite (psd), then $b_{1}$ is to be preferred to $b_{2}$. $S$ will be psd if $w^{\prime} S_{w} \geq 0$ for any non-zero vector w:nx1. The TMSE equivalent is that $b_{1}$ is prefered to $b_{2}$ whenever

$$
\begin{equation*}
\operatorname{TMSE}\left(w^{\prime} b_{2}\right) \geq \operatorname{TMSE}\left(w^{\prime} b_{1}\right), \text { for every vector } w \tag{8.2.2}
\end{equation*}
$$

This MSE criterion is the so called 'strong MSE criterion', and a weaker criterion for $b_{1}$ to be preferred to $b_{2}$ is that

$$
\begin{equation*}
\operatorname{TMSE}\left(\mathrm{b}_{2}\right) \geq \operatorname{TMSE}\left(\mathrm{b}_{1}\right) \tag{8.2.3}
\end{equation*}
$$

Looking at these criteria it is worthwhile to point out:
(i) $\operatorname{TMSE}\left(\mathrm{b}_{\mathrm{i}}\right)$ is the Euclidian distance between $\mathrm{b}_{\mathrm{i}}$ and $\beta$. One therefore
seeks an estimator that minimizes this norm.
(ii) The criteria (8.2.1) to (8.2.3) were defined in PC estimation as ways to determine which PC's to eliminate.
(iii) Although only some criteria are explicity stated here, there is a whole range available. For instance, all those criteria applied in PC to eliminate PC's can be generalized. A detailed discussion of criteria appears in Yinod and Ullah (1981, Chapter 2).
(iv) Some authors perform comparisons based on the relative efficiency (RE) of each estimator to OLSE (i.e. see section 4.7 where this RE concept was used to compare the various GRRE's). There appears to be no statistical analysis of these efficiency ratios, literally only comparisons of the values (i.e. $2<3$, or 2.9 is slightly better than 3 , and so on)
(v) Empirical comparisons of estimators reveals that no one estimator is clearly superior to the others. The conditions for superiority depend on the degree of collinearity, the orientation of $\delta$, and the value of $\sigma^{2}$. These factors should always be considered when choosing an estimator. Although some rough guidelines can be given, the 'best' estimator of any problem will be unique to that particular problem and no 'recipe' is available.
(vi) Simulation studies have been performed to compare the different estimators under various degrees of collinearity, changing $\sigma^{2}$ values, and different orientations of the $\delta$ 's. Some of these results have been reported in the appropriate foregoing chapters, and a study of some FPC estimators will be reported in the next section.

### 8.3 Comparison of FPC estimators

In Chapter 5 it was shown that most estimators (0LS, RR, GRR, PC) can be considered as FPC estimators. All the FPC estimators give an improvement over OLS in terms of the MSE criteria if

$$
\begin{equation*}
\sigma^{2} \geq \delta^{\prime} \Delta^{2}(\mathrm{I}-\mathrm{F})^{2}\left(\mathrm{I}-\mathrm{F}^{2}\right)^{-} \delta \tag{5.2.10}
\end{equation*}
$$

Thus the improvement over 0LS will depend on the vector $\delta$, the degree of collinearity ( $\Delta$ ), and $\sigma^{2}$.

In (5.2.12) the optimal values for the fractions were obtained as

$$
\mathrm{f}_{\mathrm{j}}^{0}=\delta_{\mathrm{j}}^{2} \lambda_{\mathrm{j}}\left(\sigma^{2}+\delta_{\mathrm{j}}^{2} \lambda_{\mathrm{j}}\right)^{-1}
$$

and the min TMSE will be

$$
\min \left\{\operatorname{TMSE}\left(\hat{\delta}_{\mathrm{FPC}}\right)\right\}=\operatorname{TMSE}(\hat{\delta})-\sum_{\mathrm{i}=1}^{\mathrm{p}}\left(\sigma^{2} / \lambda_{\mathrm{i}}\right)\left(1-\mathrm{f}_{\mathrm{i}}^{0}\right) \quad(\text { from }(5.2 .14))
$$

By using this optimal value two new estimators (FPCI and FPCV) were defined in (5.2.16) and (5.2.18).

To evaluate different FPC estimators under a variety of conditions Lee and Birch (1988) set up the following artificial data sets:
$\mathrm{p}=4: \mathrm{n}=20$; eigenvalues $: 2.96072,1.02801,0.0112$, and 0.00012 ; condition number $=157.07762$.

Two orientations of $\delta$ were considered:

$$
\begin{aligned}
& \text { I: } \delta=\left[\begin{array}{llll}
385.112 & 91.9509 & 333.138 & -198.061
\end{array}\right] \\
& \text { II: } \delta=\left[\begin{array}{llll}
-51.4399 & -86.844 & -28.844 & 542.225
\end{array}\right]
\end{aligned}
$$

Two values of $\sigma$ were chosen as 5 and 10 .

They then compared 0LS, RR (where $k$ is obtained by an iterative method starting with OLS), GRR ( $K$ is obtained by an iterative method as in (5.2.15) starting with 0LS), two PC estimators (called PC(1) and PC(2) to denote the number of deleted eigenvalues and vectors) and the two new estimators (FPCI and FPCV with iteration initial value $\operatorname{PC}(1)$ ), by means of a Monte Carlo
study with the number of repetitions set at 50 . The performances of the estimators across the four combinations and the 50 repetitions was summarized by computing the standardized empirical mean square error defined as

$$
\begin{equation*}
\operatorname{SEMSE}(b)=\sum_{j=1}^{4} \sum_{i=1}^{50} \frac{\left(\mathrm{~b}_{\mathrm{j}} \mathrm{i}^{-}-\beta_{\mathrm{j}}\right)^{2}}{50 \sigma_{\mathrm{j}}^{2}, \mathrm{LS}} \tag{8.3.1}
\end{equation*}
$$

where $b_{j i}$ is the $j$-th element of $b$, an estimator of $\beta$, in the $i$-th Monte Carlo repetition and $\sigma_{j}^{2}, L S$ is the theoretical variance of $b_{j}, L S, j=1, \ldots, 4$. Lee (1986) shows that the $\operatorname{SEMSE}(\hat{\beta})$ is independent of both the orientation of the parameter vector and the value of $\sigma^{2}$, but does depend on the structure of the regressor variables. The use of the SEMSE(b) facilitates comparisons of estimators both within a given combination and across the combinations. The results of this Monte Carlo study are displayed in Table II of Lee and Birch (1988), where the evaluations are based on the relative efficiency of each estimator compared to $\hat{\beta}$. Their results and suggestions can be summarized as follows:

1. While $P C(1)$ is far superior to LS in terms of SEMSE, just the opposite is true for $\mathrm{PC}(2)$ estimator. This suggests that the user of PC regression should exercise caution in determining which principal components should be deleted. This finding is analogous to the test and ideas described in Chapter 3. It is interesting that the authors find for orientation II and $\sigma=5$, a PCE with only the third component deleted more efficient than LS, $\mathrm{PC}(1)$ and $\mathrm{PC}(2)$. They do not however give any explanation as to why this phenomenon occurred and it may constitute another reason why one should always be on the alert, when deleting PC's.
2. The relative behaviour of the RR, GRR, FPCI and FPCV estimators is dependent upon the orientation of the parameter vector. The success of the FPCI and FPCV depends on the proper choice of the PC estimator used as the initial values. Once this choice is made the authors suggest that either FPCI or FPCV could be computed, and for their simulation each gives practically the same results.
3. As $\sigma^{2}$ increases, the relative efficiency of estimators increased with respect to 0LS.
4. The FPCV estimator is always more efficient than the RR estimator, but not always more efficient than the GRR estimator. The FPCI estimator is not always more efficient than the RR and GRR estimator.
5. In conclusion: among the biased estimators, no single estimator dominates the others, in terms of TMSE, across all conditions.

### 8.4 Summary

In this chapter the various estimators were summarized. Judgement and comparisons of estimators were discussed and the idea of a class of FPC estimators was introduced.

## Chapter 9

## INFLUENCE AND COLLINEARITY

In this chapter the influence of observations as well as variables on the collinearity of the model is investigated. We introduce in $\S 9.1$ the concept of influential points, ways of determining collinearity-influential points in $\S 9.2$, influential variables in $\S 9.3$ and in $\S 9.4$ the concept of influence in ridge regression estimators and weighted least squares. In this chapter we will assumme that $X$ in (1.1) is a full column rank matrix.

### 9.1 Influential points

Traditionally, collinearity has been associated with the columns and the column space of $X$. However, as already pointed out in Chapter 2, the collinearity structure of the data can be strongly affected by a few observations (Belsley, et al. (1980), Mason and Gunst (1985), Draper and John (1981)). The term influential is used to describe an observation whose inclusion in a data set substantially changes regression coefficent estimates, predicted responses, or the results of inferential procedures (Mason and Gunst (1985)). Not all outliers are necessarily collinearityinfluence points and vice versa.

A formal definition for a collinearity-influential point or case is given by Walker (1989):

Let $\mathrm{x}_{\mathrm{i}}{ }^{\prime}$ be the i - th row of the X matrix and let $\eta_{\mathrm{k}}$ and $\eta_{\mathrm{k}}(-\mathrm{i})$ be the $k$ - th condition indices computed with and without $x_{i}{ }^{\prime}$, respectively. The i-th point ( $\mathrm{x}_{\mathrm{i}}{ }^{\prime}$ ) is a collinearity-influential point if, for a predetermined value $\delta$,

$$
\begin{equation*}
\left|\eta_{\mathrm{k}}-\eta_{\mathrm{k}}(-\mathrm{i})\right|>\delta \eta_{\mathrm{k}} \text { for } \mathrm{k}=2,3, \ldots, \mathrm{p} \tag{9.1.1}
\end{equation*}
$$

This definition focusses upon the potential of the case to influence the estimation and inference procedures. But the actual influence depends on the $Y$-value observed for that case. Influential points in the sense of

Walker can be classified into two broad categories: those that mask collinearity and those that induce collinearity, illustrated in the following figures from Walker (1989)


Masking collinearity

$$
\left(\eta_{\mathrm{k}} \ll \eta_{\mathrm{k}}(-\mathrm{i})\right)
$$



Inducing collinearity $\left(\eta_{\mathrm{k}} \gg \eta_{\mathrm{k}}(-\mathrm{i})\right)$

### 9.2 Detecting collinearity-influential points

The first step in finding collinearity-influential points is to detect outliers. Various methods of detecting outliers are available in the literature. Common methods include graphical representation of the residuals and standardized residuals versus the individual fitted $Y$ and observed $X$ variables, and normal probability plots (see for instance Cook and Weisberg (1982) and Daniel and Wood (1980)). High leverage, as defined below, is also an indication that a point is potentially influential. Once a observation is potentially influential it is necessary to see what actual influence it has on the estimator and inference. Various measures of this influence will be discussed.

### 9.2.1 Leverage

From $\S 1.7$ the Hat matrix is

$$
H=X\left(X^{\prime} X\right)^{-1} X^{\prime}
$$

It maps $Y$ into $\hat{Y}$

$$
\hat{Y}=H Y
$$

and the diagonal elements of $H$ are known as the leverage values:

$$
\begin{equation*}
h_{i j}=x_{i}^{\prime}\left(X^{\prime} X\right)^{-1} x_{i} \tag{1.7.2}
\end{equation*}
$$

The following properties of $h_{i i}$ are pertinent here:

$$
\begin{array}{ll}
\text { (i) } h_{i \mathrm{i}}=\mathrm{h}_{\mathrm{i} i}^{2}+\sum_{\mathrm{j} \neq \mathrm{i}} \mathrm{~h}_{\mathrm{ij}}^{2} & \text { (from H idempotent) } \\
\text { (ii) } 0 \leq h_{\mathrm{i} i} \leq 1 & \text { (H is a projection matrix) } \\
\text { (iii) } \sum_{i=1}^{\mathrm{n}} \mathrm{~h}_{\mathrm{i} i}^{2}=\mathrm{p} & \text { (X is full column rank) } \tag{9.2.3}
\end{array}
$$

When $h_{i i}$ is large, the $i$-th case is called a high-leverage point. Hoaglin and Welsh (1978) suggested that a point has high leverage if $h_{i i}>2 p / n$. Mason and Gunst (1985) use $\mathrm{h}_{\mathrm{i}}>2(\mathrm{p}+1) / \mathrm{n}$. The influence of the response value $Y_{i}$ on the fitted value $\hat{Y}_{i}$ is reflected in the corresponding leverage. For the two extreme cases, ( $\mathrm{h}_{\mathrm{i} \mathrm{i}}=0$ or 1 ) we have:

If $h_{i \mathrm{i}}=0$, then $\underset{\mathrm{j} \neq \mathrm{i}}{\sum} \mathrm{h}_{\mathrm{ij}}^{2}=0 \quad$ (from (9.2.1)) and thus $\mathrm{h}_{\mathrm{ij}}=0$ and $\hat{Y}_{i}=0$. Thus $\hat{Y}_{i}$ must be fixed at zero by design - it is not effected by any $Y_{i}$.

If $h_{i \mathrm{i}}=1$, then $\underset{\mathrm{j} \neq \mathrm{i}}{\sum} h_{\mathrm{i} j}^{2}=0 \quad\left(\right.$ from (9.2.1)) and thus $h_{i j}=0$ and $\hat{Y}_{\mathrm{i}}=\mathrm{Y}_{\mathrm{i}}$, implying that $\hat{\epsilon}_{\mathrm{i}}=0$. The model fits the data value exactly.

### 9.2.2 Outlier sum of squares

Once outliers or influential observations are detected model (1.1) can be partitioned as

$$
\left[\begin{array}{l}
\mathrm{Y}_{\mathrm{a}}  \tag{9.2.4}\\
\mathrm{Y}_{\mathrm{b}}
\end{array}\right]=\left[\begin{array}{l}
\mathrm{X}_{\mathrm{a}} \\
\mathrm{X}_{\mathrm{b}}
\end{array}\right] \beta+\left[\begin{array}{c}
\epsilon_{\mathrm{a}} \\
\epsilon_{\mathrm{b}}
\end{array}\right]
$$

where by some rearrangement of rows the $k$ observations that are influential or outlying are contained in $\left[\mathrm{Y}_{\mathrm{b}} \mathrm{X}_{\mathrm{b}}\right]$. Thus $\mathrm{Y}_{\mathrm{a}}$ and $\epsilon_{\mathrm{a}}$ are $(\mathrm{n}-\mathrm{k}) \mathrm{x} 1$ vectors, $\mathrm{Y}_{\mathrm{b}}$ and $\epsilon_{\mathrm{b}}$ are kx1 vectors, $\mathrm{X}_{\mathrm{a}}$ is a ( $\mathrm{n}-\mathrm{k}$ ) xp matrix and $\mathrm{X}_{\mathrm{b}}$ is a kxp matrix. The residuals $\hat{\epsilon}=[\mathrm{I}-\mathrm{H}] \mathrm{Y}$ under this model can also be partitioned as

$$
\left[\begin{array}{c}
\hat{\epsilon}_{a}  \tag{9.2.5}\\
\hat{\epsilon}_{\mathrm{b}}
\end{array}\right]=\left[\begin{array}{rr}
\mathrm{I}_{\mathrm{n}-\mathrm{k}^{-}}-\mathrm{X}_{\mathrm{a}}\left(\mathrm{X}^{\prime} \mathrm{X}\right)^{-1} \mathrm{X}_{\mathrm{a}}^{\prime} & -\mathrm{X}_{\mathrm{a}}\left(X^{\prime} X\right)^{-1} \mathrm{X}_{\mathrm{b}}^{\prime} \\
-X_{b}\left(X^{\prime} X\right)^{-1} X_{a}^{\prime} & I_{k^{-}} X_{b}\left(X^{\prime} X\right)^{-1} X_{b}^{\prime}
\end{array}\right]\left[\begin{array}{c}
\mathrm{Y}_{\mathrm{a}} \\
\mathrm{Y}_{\mathrm{b}}
\end{array}\right]
$$

Deleting the suspected outliers gives a model with $E\left[Y_{a}\right]=X_{a} \beta$. Draper and John (1981) suggested the following alternative model:

$$
\left[\begin{array}{l}
\mathrm{Y}_{\mathrm{a}}  \tag{9.2.6}\\
\mathrm{Y}_{\mathrm{b}}
\end{array}\right]=\left[\begin{array}{ll}
\mathrm{X}_{\mathrm{a}} & 0 \\
\mathrm{X}_{\mathrm{b}} & \mathrm{I}_{\mathrm{k}}
\end{array}\right]\left[\begin{array}{c}
\beta \\
a
\end{array}\right]+\left[\begin{array}{c}
0_{\mathrm{a}} \\
\theta_{\mathrm{b}}
\end{array}\right]
$$

where $\alpha$ is a kx1 vector of additional parameters. The estimates of $\beta$ and $a$ under this model are:

$$
\begin{align*}
& \mathrm{Y}_{\mathrm{a}}=\mathrm{X}_{\mathrm{a}} \tilde{\beta}, \text { thus } \tilde{\beta}=\left(\mathrm{X}_{\mathrm{a}}^{\prime} \mathrm{X}_{\mathrm{a}}\right)^{-1} \mathrm{X}_{\mathrm{a}}^{\prime} \mathrm{Y}_{\mathrm{a}}  \tag{9.2.7}\\
& \mathrm{Y}_{\mathrm{b}}=\mathrm{X}_{\mathrm{b}} \tilde{\beta}+\tilde{a}
\end{align*}
$$

Thus, by using (9.2.5)

$$
\begin{aligned}
\tilde{a} & =Y_{b}-X_{b} \tilde{\beta} \\
& =\left[I-X_{b}\left(X^{\prime} X\right)^{-1} X_{b}^{\prime}\right]^{-1}\left[\hat{\epsilon}_{\mathrm{b}}+X_{b}\left(X^{\prime} X\right)^{-1} X_{a}^{\prime} Y_{a}\right]-X_{b}\left(X_{a}^{\prime} X_{a}\right)^{-1} X_{a}^{\prime} Y_{a} \\
& =\left[I-X_{b}\left(X^{\prime} X\right)^{-1} X_{b}^{\prime}\right]^{-1}\left[\hat{\epsilon}_{b}+X_{b}\left(X^{\prime} X\right)^{-1} X_{a}^{\prime} Y_{a}-\left[I-X_{b}\left(X^{\prime} X\right)^{-1} X_{b}^{\prime}\right] X_{b}\left(X_{a}^{\prime} X_{a}\right)^{-1} X_{a}^{\prime} Y_{a}\right]
\end{aligned}
$$

and

$$
\begin{aligned}
X_{b}\left(X^{\prime} X\right)^{-1} & X_{a}^{\prime} Y_{a}-\left[I-X_{b}\left(X^{\prime} X\right)^{-1} X_{b}^{\prime}\right] X_{b}\left(X_{a}^{\prime} X_{a}\right)^{-1} X_{a}^{\prime} Y_{a} \\
& =X_{b}\left(X^{\prime} X\right)^{-1} X_{a}^{\prime} Y_{a}-X_{b}\left(X_{a}^{\prime} X_{a}\right)^{-1} X_{a}^{\prime} Y_{a}+X_{b}\left(X^{\prime} X\right)^{-1} X_{b}^{\prime} X_{b}\left(X_{a}^{\prime} X_{a}\right)^{-1} X_{a}^{\prime} Y_{a} \\
& =X_{b}\left(X^{\prime} X\right)^{-1}\left[I-\left(X^{\prime} X\right)\left(X_{a}^{\prime} X_{a}\right)^{-1}+X_{b}^{\prime} X_{b}\left(X_{a}^{\prime} X_{a}\right)^{-1}\right] X_{a}^{\prime} Y_{a} \\
& =X_{b}\left(X^{\prime} X\right)^{-1}\left[I-\left\{\left(X_{a}^{\prime} X_{a}+X_{b}^{\prime} X_{b}\right)-X_{b}^{\prime} X_{b}\right\}\left(X_{a}^{\prime} X_{a}\right)^{-1}\right] X_{a}^{\prime} Y_{a} \\
& =X_{b}\left(X^{\prime} X\right)^{-1}[0] X_{a}^{\prime} Y_{a} \\
& =0
\end{aligned}
$$

we have

$$
\begin{align*}
\tilde{a} & =\left[\mathrm{I}-\mathrm{X}_{\mathrm{b}}\left(\mathrm{X}^{\prime} \mathrm{X}\right)^{-1} \mathrm{X}_{\mathrm{b}}^{\prime}\right]^{-1}\left[\hat{\epsilon}_{\mathrm{b}}+0\right] \\
& =\left[\mathrm{I}-\mathrm{X}_{\mathrm{b}}\left(\mathrm{X}^{\prime} \mathrm{X}\right)^{-1} \mathrm{X}_{\mathrm{b}}^{\prime}\right]^{-1} \hat{\epsilon}_{\mathrm{b}} \tag{9.2.8}
\end{align*}
$$

The adjusted observations $\mathrm{Y}_{\mathrm{b}}$ with expectations $\mathrm{X}_{\mathrm{b}} \beta$ can be estimated by $\mathrm{Y}_{\mathrm{b}}$ - $\tilde{a}$, then model (9.2.4) becomes

$$
\left[\begin{array}{l}
\mathrm{Y}_{\mathrm{a}}  \tag{9.2.9}\\
\mathrm{Y}_{\mathrm{b}}-\tilde{a}
\end{array}\right]=\left[\begin{array}{l}
\mathrm{X}_{\mathrm{a}} \\
\mathrm{X}_{\mathrm{b}}
\end{array}\right] \beta+\left[\begin{array}{l}
\theta_{\mathrm{a}} \\
\theta_{\mathrm{b}}
\end{array}\right]
$$

thus,

$$
\begin{align*}
\hat{\beta} & =\left(X^{\prime} X\right)^{-1} X^{\prime} Y \\
Y_{b}-\tilde{a} & =Y_{b}-\left[I-X_{b}\left(X^{\prime} X\right)^{-1} X_{b}^{\prime}\right]^{-1} \hat{\epsilon}_{b} \quad(\text { from } 9.2 .8) \\
& =\left[I-X_{b}\left(X^{\prime} X^{-1} X_{b}^{\prime}\right]^{-1} X_{b}\left(X^{\prime} X\right)^{-1} X_{a}^{\prime} Y_{a} \quad(\text { from } 9.2 .5)\right. \\
\hat{\theta}_{b} & =\left(Y_{b}-\tilde{a}\right)-X_{b} \hat{\beta} \\
& =\left[I-X_{b}\left(X^{\prime} X\right)^{-1} X_{b}^{\prime}\right]^{-1} X_{b}\left(X^{\prime} X\right)^{-1} X_{a}^{\prime} Y_{a}-X_{b}\left(X^{\prime} X\right)^{-1} X_{a}^{\prime} Y_{a} \\
& \quad-X_{b}\left(X^{\prime} X\right)^{-1} X_{b}^{\prime}\left[I-X_{b}\left(X^{\prime} X\right)^{-1} X_{b}^{\prime}\right]^{-1} X_{b}\left(X^{\prime} X\right)^{-1} X_{a}^{\prime} Y_{a} \\
& =\left[I-\left[I-X_{b}\left(X^{\prime} X\right)^{-1} X_{b}^{\prime}\right]-X_{b}\left(X^{\prime} X\right)^{-1} X_{b}^{\prime}\right]\left[I-X_{b}\left(X^{\prime} X\right)^{-1} X_{b}^{\prime}\right]^{-1} X_{b}\left(X^{\prime} X\right)^{-1} X_{a}^{\prime} Y_{a} \\
& =0 \tag{9.2.10}
\end{align*}
$$

$$
\begin{align*}
\hat{\theta}_{\mathrm{a}} & =Y_{a}-X_{a} \hat{\beta} \\
& =Y_{a}-X_{a}\left(X^{\prime} X\right)^{-1} X_{a}^{\prime} Y_{a}-X_{a}\left(X^{\prime} X\right)^{-1} X_{b}^{\prime}\left[I-X_{b}\left(X^{\prime} X\right)^{-1} X_{b}^{\prime}\right]^{-1} X_{b}\left(X^{\prime} X\right)^{-1} X_{a}^{\prime} Y_{a} \\
& =\left[I-X_{a}\left(X^{\prime} X\right)^{-1} X_{a}^{\prime}-X_{a}\left(X^{\prime} X\right)^{-1} X_{b}^{\prime}\left[I-X_{b}\left(X^{\prime} X\right)^{-1} X_{b}^{\prime}\right]^{-1} X_{b}\left(X^{\prime} X\right)^{-1} X_{a}^{\prime}\right] Y_{a} \tag{9.2.11}
\end{align*}
$$

The extra sum of squares due to fitting $a$ in model (9.2.6), as compared with model (9.2.4) is

$$
\begin{equation*}
Q_{\mathrm{k}}=\hat{\epsilon}_{\mathrm{b}}^{\prime}\left[\mathrm{I}-\mathrm{X}_{\mathrm{b}}\left(\mathrm{X}^{\prime} \mathrm{X}\right)^{-1} \mathrm{X}_{\mathrm{b}}^{\prime}\right]^{-1} \hat{\epsilon}_{\mathrm{b}} \tag{9.2.12}
\end{equation*}
$$

Gentleman and Wilk (1975), called $Q_{k}$ the outlier sum of squares associated with the observations $Y_{b}$. Then the F-statistic, associated with the hypothesis $H_{0}: a=0$ is:

$$
\begin{align*}
F & =\frac{\hat{\epsilon}_{b}^{\prime}\left[I-X_{b}\left(X^{\prime} X\right)^{-1} X_{b}^{\prime}\right]^{-1} \hat{\epsilon}_{b}}{\hat{\epsilon}^{\prime} \hat{\epsilon}^{-}-\hat{\epsilon}_{b}^{\prime}\left[I-X_{b}\left(X^{\prime} X\right)^{-1} X_{b}^{\prime}\right]^{-1} \hat{\epsilon}_{b}} \cdot \frac{n-p-k}{k} \\
& =\frac{Q_{k}}{\operatorname{RSS}-Q_{k}} \cdot \frac{n-p-k}{k} \tag{9.2.13}
\end{align*}
$$

### 9.2.3 Andrews-Pregibon statistic

Andrews and Pregibon (1978) consider the following two augmented matrices:

$$
\begin{aligned}
& X_{1}^{*}=\left[\begin{array}{ll}
X & Y
\end{array}\right] \quad(\text { from model }(1.1)) \\
& X_{2}^{*}=\left[\begin{array}{lll}
X & D & Y
\end{array}\right] \quad(\text { from model }(9.2 .6))
\end{aligned}
$$

where $D$ is the matrix with which the $X$ matrix was augmented in (9.2.6)

$$
D=\left[\begin{array}{l}
0 \\
I_{k}
\end{array}\right]
$$

The Andrews- Pregibon (AP) statistic is then defined as:

$$
\begin{equation*}
A P_{i j}^{k} \ldots=\left|X_{2}^{*} X_{2}^{*}\right| /\left|X_{1}^{*} X_{1}^{*}\right| \tag{9.2.14}
\end{equation*}
$$

where ij.. denote the $k$ subscripts that estimate $Y_{b}$ (the outliers). For example ${A P_{34}}_{2}$ means the Andrews-Pregibon statistic for 2 suspected outliers cases 3 and 4. Then from the result of Draper and John (1981, p25) we have

$$
\begin{align*}
& \left|X_{2}^{*} X_{2}^{*}\right|=\left|X^{\prime} X\right| \cdot \operatorname{RSS} \cdot \mid\left(I-X_{b}\left(X^{\prime} X\right)^{-1} X_{b}^{\prime} \mid \cdot\left(1-Q_{\mathrm{k}} / \operatorname{RSS}\right)\right.  \tag{9.2.15}\\
& \left|X_{1}^{*} X_{1}^{*}\right|=\left|X^{\prime} X\right| \cdot \operatorname{RSS}
\end{align*}
$$

where RSS is the residual sums of squares from the full model (1.1). Then (9.2.14) can be written as the dimensionless quantity

$$
\begin{equation*}
A P_{i j}^{k} \ldots=\mid\left(I-X_{b}\left(X^{\prime} X\right)^{-1} X_{b}^{\prime} \mid \cdot\left(1-Q_{k} / \text { RSS }\right)\right. \tag{9.2.16}
\end{equation*}
$$

and Andrews and Pregibon (1978) regard the quantity ( $1-A P_{i j}^{k} .$. ) 'as the proportion of the volume generated by $X_{1}^{*}$ attributable to the $k$ observation ( $\mathrm{ij} . ..)^{\prime}$. Hence an interpretation of $A P_{i j}^{k} \ldots$ is that small values of $A P_{i j}^{k}$.. are associated with influential observations.

The AP statistic in (9.2.16) consists of two factors:
(i) $\mid\left(\mathrm{I}-\mathrm{X}_{\mathrm{b}}\left(\mathrm{X}^{\prime} \mathrm{X}\right)^{-1} \mathrm{X}_{\mathrm{b}}^{\prime} \mid\right.$ only involves the regressor variables, and $\mathrm{X}_{\mathrm{b}}(\mathrm{XX})^{-1} \mathrm{X}_{\mathrm{b}}$, is the leverage of Hoaglin and Welsh (1978). For $\mathrm{k}=1$ the leverage will be $x_{i}^{\prime}\left(X^{\prime} X\right)^{-1} x_{i}=h_{i \mathrm{i}}$, and thus $\mid\left(I-X_{b}\left(X^{\prime} X\right)^{-1} X_{b}^{\prime} \mid=1-h_{i \mathrm{i}}\right.$. Small values of $\mid\left(I-X_{b}\left(X^{\prime} X\right)^{-1} X_{b} \mid\right.$ correspond to large $h_{i i}$, revealing high-leverage points. This factor reflects potential influence of the cases.
(ii) $\left(1-\mathbf{Q}_{\mathrm{k}} /\right.$ RSS $)=\left(\right.$ RSS $\left.-\mathbb{Q}_{\mathrm{k}}\right) /$ RSS is a decreasing function of the F-statistic of (9.2.13). Small values of (RSS - $Q_{k}$ )/RSS indicate high influence (large sum of squares for outliers). This factor reflects actual influence of the observed responses on the fitting.

### 9.2.4 Cook's statistic

In (1.7.7) Cook's distance was defined as

$$
\begin{equation*}
\mathrm{C}=\mathrm{D}^{2}=(\hat{\beta}-\tilde{\beta})^{\prime} \mathrm{X}^{\prime} \mathrm{X}(\hat{\beta}-\tilde{\beta}) / \mathrm{p} \hat{\sigma}^{2} \tag{1.7.7}
\end{equation*}
$$

To determine the degree of influence the i-th data point has on the estimate, $\tilde{\beta}$ will be replaced by $\hat{\beta}_{-i}$, where $\hat{\beta}_{-i}$ is the least square estimate of $\beta$ with the $i$ - th point deleted. Thus for $i=1,2, \ldots, n$

$$
\begin{equation*}
C_{i}=\left(\hat{\beta}-\hat{\beta}_{-\mathrm{i}}\right)^{\prime} \mathrm{X}^{\prime} X\left(\hat{\beta}-\hat{\beta}_{-\mathrm{i}}\right) / \mathrm{p} \hat{\sigma}^{2} \tag{9.2.17}
\end{equation*}
$$

Computed values of $C_{i}$ can be compared to the $F(p, n-p)$ distribution. For example if $C_{i}$ equals the 0.50 value of the corresponding $F$-distribution, then deletion of the $i$-th case moves the estimate of $\beta$ to the edge of a $50 \%$ confidence ellipsoid relative to $\hat{\beta}$. Cook (1977) suggests that for an uncomplicated analysis one would like each $\hat{\beta}_{-i}$ to stay well within a $10 \%$ confidence region. The comparison of $C_{i}$ to an $F$ is only used for converting $C_{i}$ to a familiar scale, and $C_{i}$ is not distributed as $F$ (Cook and Veisberg (1982)).

The statistic $C_{i}$ can be simplified as follows:

$$
\begin{align*}
\hat{\beta}-\hat{\beta}_{-i} & =\hat{\beta}-\dot{\beta}+\left(X^{\prime} X\right)^{-1} \mathrm{x}_{\mathrm{i}} \hat{\epsilon}_{\mathrm{i}} /\left(1-\mathrm{h}_{\mathrm{ii}}\right) \quad(\text { from }(1.8 .7)) \\
& =\left(X^{\prime} X\right)^{-1} \mathrm{x}_{\mathrm{i}} \hat{\epsilon}_{\mathrm{i}} /\left(1-\mathrm{h}_{\mathrm{i} \mathrm{i}}\right) \tag{9.2.18}
\end{align*}
$$

thus

$$
\begin{align*}
\mathrm{C}_{\mathrm{i}} & =\mathrm{x}_{\mathrm{i}}^{\prime}\left(X^{\prime} X\right)^{-1} \mathrm{x}_{\mathrm{i}}\left(\hat{\epsilon}_{\mathrm{i}}\left(1-\mathrm{h}_{\mathrm{i} i}\right)^{-1}\right)^{2} / \mathrm{p} \hat{\sigma}^{2} \\
& =\mathrm{p}^{-1}\left[\hat{\epsilon}_{\mathrm{i}} /\left\{\hat{\sigma}\left(1-\mathrm{h}_{\mathrm{i} i}\right)^{\frac{1}{2}}\right\}\right]^{2}\left[\mathrm{~h}_{\mathrm{i} i} /\left(1-\mathrm{h}_{\mathrm{i} i}\right)\right] \\
& =\mathrm{p}^{-1} \mathrm{r}_{\mathrm{i}}^{2}\left[\mathrm{~h}_{\mathrm{i} i} /\left(1-\mathrm{h}_{\mathrm{i} i}\right)\right] \tag{9.2.19}
\end{align*}
$$

Here $r_{i}=\hat{\epsilon}_{\mathrm{i}} /\left\{\hat{\sigma}\left(1-\mathrm{h}_{\mathrm{if}}\right)^{\frac{1}{2}}\right\}$ is the standardized residual as defined in (1.7.3) and it is a measure of the degree to which the $i$-th observation can be considered as an outlier from the assumed model. We note that $r_{i}$ is a monotonic function of $Q_{1}$ (9.2.12). The ratio $h_{i i} /\left(1-h_{i i}\right)$ measures the relative sensitivity of the estimate $\hat{\beta}$, to potential outlying values, so that large values of this ratio indicate the associated point has heavy weight in the determination of $\hat{\beta}$.

The (squared) distance $C_{i}$ can be extended to contexts in which more than one case is an outlier ( $k>1$, in the outlier model (9.2.6)). Thus $\tilde{\beta}$ in (1.7.7) is computed from model (9.2.9):

$$
\begin{align*}
\tilde{\beta} & =\left(X^{\prime} X\right)^{-1} X^{\prime} Y \\
& =\left(X^{\prime} X\right)^{-1}\left[X_{a}^{\prime} X_{b}^{\prime}\right]\left[Y_{a}^{\prime}\left(Y_{b}-\tilde{a}\right)^{\prime}\right]^{\prime} \\
& =\left(X^{\prime} X\right)^{-1}\left[X_{a}^{\prime} Y_{a}+X_{b}^{\prime} Y_{b}-X_{b}^{\prime} \tilde{a}\right] \\
& =\left(X^{\prime} X\right)^{-1}\left[X^{\prime} Y-X_{b}^{\prime} \tilde{a}\right] \\
& =\hat{\beta}-\left(X^{\prime} X\right)^{-1} X_{b}^{\prime} \tilde{a} \tag{9.2.20}
\end{align*}
$$

thus

$$
\begin{equation*}
\hat{\beta}-\tilde{\beta}=\left(X^{\prime} X\right)^{-1} X_{b}^{\prime} \tilde{a} \tag{9.2.21}
\end{equation*}
$$

The distance $C_{i}$ is now indicated by $C_{i j} \ldots$, where the subscripts ij.. denote the cases contained in $\left[\mathrm{Y}_{\mathrm{b}} \mathrm{X}_{\mathrm{b}}\right.$ ], thus

$$
\begin{align*}
C_{i j} \ldots & =\frac{\tilde{a}^{\prime} X_{b}\left(X^{\prime} X\right)^{-1}\left(X^{\prime} X\right)\left(X^{\prime} X\right)^{-1} X_{b}^{\prime} \tilde{a}}{p \hat{\sigma}^{2}} \\
& =\frac{Q_{k}}{p \hat{\sigma}^{2}} \frac{\tilde{a}^{\prime} X_{b}\left(X^{\prime} X\right)^{-1} X_{b}^{\prime} \tilde{a}+\tilde{a}^{\prime} \tilde{a}-\tilde{a}^{\prime} \tilde{a}}{Q_{k}} \\
& =\frac{Q_{k}}{p \hat{\sigma}^{2}}\left[\frac{\tilde{a}^{\prime} \tilde{a}}{Q_{k}}-\frac{\tilde{a}^{\prime}\left[I-X_{b}\left(X^{\prime} X\right)^{-1} X_{b}^{\prime}\right] \tilde{a}}{Q_{k}}\right] \\
& =\frac{Q_{k}}{p \hat{\sigma}^{2}}\left[\frac{\tilde{a}^{\prime} \tilde{a}}{Q_{k}}-\frac{\hat{\epsilon}_{b}^{\prime}\left[I-X_{b}\left(X^{\prime} X\right)^{-1} X_{b}^{\prime}\right]^{-1} \hat{\epsilon}_{b}}{Q_{k}}\right] \quad \text { (from (9.2.8)) } \\
& =\frac{Q_{k}}{p \hat{\sigma}^{2}}\left[\frac{\tilde{a}^{\prime} \tilde{a}}{Q_{k}}-1\right] \quad \text { (from (9.2.12) } \tag{9.2.22}
\end{align*}
$$

### 9.2.5 DFFITS

The change of fit on forecasting when an observation is deleted, is defined (Belsley et al. (1980)) as

$$
\begin{align*}
\operatorname{DFFIT}_{\mathrm{i}} & =\hat{\mathrm{Y}}_{\mathrm{i}}-[\hat{\mathrm{Y}}(-\mathrm{i})]_{\mathrm{i}} \\
& =\mathrm{x}_{\mathrm{i}}^{\prime}\left[\hat{\beta}-\hat{\beta}_{-\mathrm{i}}\right] \\
& =\mathrm{h}_{\mathrm{i} \mathrm{i}} \hat{\epsilon}_{\mathrm{i}} /\left(1-\mathrm{h}_{\mathrm{ii}}\right) \quad(\text { from }(9.2 .18)) \tag{9.2.23}
\end{align*}
$$

where $[\hat{Y}(-i)]_{i}$ is the $i$-th element of $[\hat{Y}(-i)]$, and $[\hat{Y}(-i)]$ is the estimated Y obtained by using $\hat{\beta}_{-\mathrm{i}}$. For scaling purposes (9.2.23) is divided by
$\hat{\sigma}(i)\left\{h_{i j}\right\}^{\frac{1}{2}}$, where $\hat{\sigma}(i)$ is the estimated error variance when the $i$-th row of $X$ and $Y$ have been deleted (see (1.7.5)). The measure DFFITS (Belsley et al. ( 1980 , p15)) is the standardized change in the fitted value of a case when it is deleted, and is given for the i-th case by

$$
\begin{equation*}
\operatorname{DFFITS}_{i}=\left[\frac{h_{i i}}{1-h_{i_{i}}}\right]^{\frac{1}{2}} \frac{\hat{\epsilon}_{i}}{\hat{\sigma}(i) \sqrt{1-h_{i i}}} \tag{1.7.6}
\end{equation*}
$$

The authors warned against the misuse of deleting high-influence data points solely to effect a desired change in a particular estimated coefficient, or t-value. Once a high-influence point is identified it should only be deleted if shown to be in error, or if it has the effect of inducing collinearity, which, can be an undesirable property of the model.

The authors suggested a size-adjusted cut-off value for DFFITS $_{i}$ as $2 \sqrt{\mathrm{p} / \mathrm{n}}$, taking into account the sample size ( $n$ ) as well as the number of variables in the model. If the DFFITS $_{i}$ are divided into distinct groups and if a noticable gap appears, caution should be exercised in deleting observations.

The statistic DFFITS can be extended to include more than one data point. If the data points to be deleted are indicated by I, where I c $\{1,2, \ldots, n\}$ then

$$
\begin{equation*}
\operatorname{MDFFITS}_{\mathrm{I}}=\left[\hat{\beta}-\hat{\beta}_{-\mathrm{I}}\right]^{\prime} X_{-\mathrm{I}}^{\prime} X_{-\mathrm{I}}\left[\hat{\beta}-\hat{\beta}_{-\mathrm{I}}\right] \tag{9.2.24}
\end{equation*}
$$

where $\hat{\beta}_{-I}$ is the 0LS estimator with the $\{I\}$ rows deleted and $X_{-I}$ is the $X$ matrix with the $\{I\}$ rows deleted.

### 9.2.6 Variance inflation factor

Schall and Dunne (1987b) suggested the following statistic for the detection of collinearity-influential points:

$$
\begin{equation*}
R_{i j}=\frac{V I F_{i}^{-j}}{V I F_{i}} \tag{9.2.25}
\end{equation*}
$$

where $\operatorname{VIF}_{\mathrm{i}}^{-j}$ denotes the i -th variance inflation factor obtained after deletion of the $j$-th observation from model (1.1). When $R_{i j} \approx 1$ for all $i$ and $j$, there are no collinearity influential points; if $R_{i j} \ll 1, x_{j}$ is a point that induces collinearity; and if $\mathrm{R}_{\mathrm{ij}} \gg 1$, $\mathrm{x}_{\mathrm{j}}$ hides (masks) collinearity.

A straightforward generalization of (9.2.25) can be defined as

$$
\begin{equation*}
\mathrm{R}_{\mathrm{IJ}}=\frac{\mathrm{VIF}_{\mathrm{I}}^{-\mathrm{J}}}{\mathrm{VIF}_{\mathrm{I}}} \tag{9.2.26}
\end{equation*}
$$

where $\mathrm{I} \subset\{1, \ldots, \mathrm{p}\}$, $\mathrm{J} \subset\{1, \ldots, \mathrm{n}\}$ and $\operatorname{VIF}_{\mathrm{I}}^{-\mathrm{J}}$ is the variance inflation factor of the set of I variables and $J$ is the subset of data points that are deleted before calculating it.

The augmented model (9.2.6), where $I_{k}$ reduces to $u_{j}$, is

$$
\mathrm{Y}=\left[\begin{array}{ll}
\mathrm{X} & \mathrm{u}_{\mathrm{j}}
\end{array}\right]\left[\begin{array}{l}
\beta  \tag{9.2.27}\\
\lambda
\end{array}\right]+\epsilon
$$

Since here $u_{j}$ is the $j$-th unit vector in the space $R^{n}$ (not the $j$-th column of $U$ in the SVD of $X$ ), $V I F_{i}^{-j}$ can be obtained as the partial variance inflation factor $\operatorname{VIF}_{\mathrm{i}}\left(\mathrm{X}_{-\mathrm{i}} \mid \mathrm{u}_{\mathrm{j}}\right)$ :

$$
\begin{align*}
R_{i j} & =\frac{\operatorname{VIF}_{\mathrm{i}}^{-j}}{\operatorname{VIF}_{\mathrm{i}}} \\
& =\frac{\operatorname{VIF}_{\mathrm{i}}\left(X_{-i} \mid u_{j}\right)}{\operatorname{VIF}_{\mathrm{i}}\left(X_{-i}\right)} \\
& =\frac{\operatorname{VIF}_{\mathrm{i}}\left(\left[X_{-i} u_{j}\right]\right)}{\operatorname{VIF}_{\mathrm{i}}\left(X_{-i}\right) \cdot \operatorname{VIF}_{\mathrm{i}}\left(u_{j}\right)} \tag{9.2.28}
\end{align*}
$$

Let $e_{i}$ denote the residual vector of $X_{i}$ when this variable is regressed on $X_{-i}$, thus

$$
\begin{equation*}
\mathrm{e}_{\mathrm{i}}^{\prime} \mathrm{e}_{\mathrm{i}}=\mathrm{X}_{\mathrm{i}}{ }^{\prime}\left[\mathrm{I}-\mathrm{X}_{-\mathrm{i}}\left(\mathrm{X}_{-i}^{\prime} \mathrm{X}_{-\mathrm{i}}\right)^{-1} \mathrm{X}_{-\mathrm{i}}^{\prime}\right] X_{\mathrm{i}} \tag{9.2.29}
\end{equation*}
$$

and, similarly let $e_{i}^{-j}$ denote the residual vector of $X_{i}$ when the chosen regressors are $\left[\mathrm{X}_{-\mathrm{i}} \mathrm{u}_{\mathrm{j}}\right]$, thus

$$
\left(e_{i}^{-j}\right)^{\prime}\left(e_{i}^{-j}\right)=X_{i}^{\prime}\left\{I-\left[\begin{array}{ll}
X_{-i} & u_{j}
\end{array}\right]\left[\begin{array}{ll}
X_{-i}^{\prime} X_{-i} & X_{-i}^{\prime} u_{j}  \tag{9.2.30}\\
u_{j}^{\prime} X_{-i} & u_{j}^{\prime} u_{j}
\end{array}\right]^{-1}\left[\begin{array}{ll}
X_{-i} & u_{j}
\end{array}\right]^{\prime}\right\} X_{i}
$$

Note that the first diagonal element of $\left\{\left[\begin{array}{ll}X & u_{j}\end{array}\right]^{\prime}\left[\begin{array}{ll}X & u_{j}\end{array}\right]\right\}^{-1}$ is

$$
\begin{aligned}
& \left\{\left[\begin{array}{ll}
X & u_{j}
\end{array}\right]^{\prime}\left[\begin{array}{ll}
X & u_{j}
\end{array}\right]\right\}_{11}^{-1}=\left[\begin{array}{lll}
X_{i}{ }^{\prime} X_{i} & X_{i}{ }^{\prime} X_{-i} & X_{i}{ }^{\prime} u_{j} \\
X_{-i}^{\prime} X_{i} & X_{-i}^{\prime} X_{-i} & X_{-i}^{\prime} u_{j} \\
u_{j}{ }^{\prime} X_{i} & u_{j} X_{-i} & u_{j}{ }^{\prime} u_{j}
\end{array}\right]_{11}^{-1} \\
& =\left\{X_{i}{ }^{\prime} X_{i}-\left[X_{i}{ }^{\prime} X_{-i} X_{i}^{\prime} u_{j}\right]\left[\begin{array}{ll}
X_{-i}^{\prime} X_{-i} & X_{-i}^{\prime} u_{j} \\
u_{j}{ }^{\prime} X_{-i} & u_{j}^{\prime} u_{j}
\end{array}\right]^{-1}\left[\begin{array}{l}
X_{-i}^{\prime} X_{i} \\
u_{j}{ }^{\prime} X_{i}
\end{array}\right]^{-1}\right. \\
& =\left\{X_{i}^{\prime} X_{i}-X_{i}^{\prime}\left[\begin{array}{ll}
X_{-i} & u_{j}
\end{array}\right]\left[\begin{array}{ll}
X_{-i}^{\prime} X_{-i} & X_{-i}^{\prime} u_{j} \\
u_{j}^{\prime} X_{-i} & u_{j}^{\prime} u_{j}
\end{array}\right]^{-1}\left[\begin{array}{l}
X_{-i}^{\prime} \\
u_{j}{ }^{\prime}
\end{array}\right]^{X_{i}}\right\}^{-1}
\end{aligned}
$$

$$
\begin{align*}
& =\left\{X_{i}^{\prime}\left[I-\left[\begin{array}{ll}
X_{-i} & u_{j}
\end{array}\right]\left[\begin{array}{ll}
X_{-i}^{\prime} X_{-i} & X_{-i}^{\prime} u_{j} \\
u_{j}^{\prime} X_{-i} & u_{j}^{\prime} u_{j}
\end{array}\right]^{-1}\left[\begin{array}{l}
X_{-i}^{\prime} \\
u_{j}
\end{array}\right]\right] X_{i}\right\}^{-1} \\
& =\left\{\left(e_{i}^{-j}\right)^{\prime}\left(e_{i}^{-j}\right)\right\}^{-1} \tag{9.2.31}
\end{align*}
$$

The first diagonal element of $\left[X^{\prime} X\right]^{-1}$, is

$$
\begin{align*}
{\left[X^{\prime} X\right]_{1}^{-1} } & =\left[\begin{array}{ll}
X_{i}^{\prime} X_{i} & X_{i}^{\prime} X_{-i} \\
X_{-i}^{\prime} X_{i} & X_{-i}^{\prime} X_{-i}
\end{array}\right]_{11}^{-1} \\
& =\left\{X_{i}^{\prime} X_{i}-X_{i}^{\prime} X_{-i}\left(X_{-i}^{\prime} X_{-i}\right)^{-1} X_{-i}^{\prime} X_{i}\right\}^{-1} \\
& =\left\{X_{i}^{\prime}\left[I-X_{-i}\left(X_{-i}^{\prime} X_{-i}\right)^{-1} X_{-i}^{\prime}\right] X_{i}\right\}^{-1} \\
& =\left\{e_{i}^{\prime} e_{i}\right\}^{-1} \tag{9.2.32}
\end{align*}
$$

and the first diagonal element of $\left\{\left[\begin{array}{ll}X_{i} & u_{j}\end{array}\right]^{\prime}\left[\begin{array}{ll}X_{i} & u_{j}\end{array}\right]\right\}^{-1}$ is

$$
\begin{align*}
\left\{\left[\begin{array}{ll}
x_{i} & u_{j}
\end{array}\right]^{\prime}\left[\begin{array}{ll}
x_{i} & u_{j}
\end{array}\right]\right\}_{1 i}^{-1} & =\left[\begin{array}{ll}
x_{i}^{\prime} X_{i} & x_{i}^{\prime} u_{j} \\
u_{j}^{\prime} X_{i} & u_{j}^{\prime} u_{j}
\end{array}\right]_{11}^{-1} \\
& =\left[\begin{array}{ll}
x_{i}^{\prime} X_{i} & x_{i j} \\
x_{i j} & 1
\end{array}\right]_{11}^{-1} \\
& =1 /\left(x_{i}^{\prime} X_{i}-x_{i j}^{2}\right)\left[\begin{array}{cc}
1 & -x_{i j} \\
-x_{i . j} & X_{i}^{\prime} x_{i}
\end{array}\right]_{11} \\
& =1 /\left(x_{i}^{\prime} X_{i}-x_{i, j}^{2}\right) \tag{9.2.33}
\end{align*}
$$

Then

$$
\begin{aligned}
\operatorname{VIF}_{i}\left(X_{-i} u_{j}\right) & =\frac{V\left(X_{i} X_{-i} u_{j}\right)}{V\left(X_{i}\right)} \\
& =\frac{\left\{\left[\begin{array}{ll}
X & u_{j}
\end{array}\right]^{\prime}\left[\begin{array}{ll}
X & u_{j}
\end{array}\right]\right\}_{i i}^{-1}}{\left(X_{i}^{\prime} X_{i}\right)^{-1}}
\end{aligned}
$$

$$
\begin{align*}
& =\left\{\left(e_{i}^{-j}\right)^{\prime}\left(e_{i}^{-j}\right)\right\}^{-1} /\left(X_{i}^{\prime} X_{i}\right)^{-1} \quad(\text { from }(9.2 .31)) \\
& =\left(X_{i}^{\prime} X_{i}\right) /\left\{\left(e_{i}^{-j}\right)^{\prime}\left(e_{i}^{-j}\right)\right\}  \tag{9.2.34}\\
\operatorname{VIF}_{i}\left(X_{-i}\right) & =\frac{V\left(X_{i} X_{-i}\right)}{V\left(X_{i}\right)} \\
& =\frac{\left(X^{\prime} X\right)_{i}^{-1}}{\left(X_{i}^{\prime} X_{i}\right)^{-1}} \\
& =\left\{e_{i}^{\prime} e_{i}\right\}^{-1} /\left(X_{i}^{\prime} X_{i}\right)^{-1} \quad(\text { from }(9.2 .32)) \\
& =\left(X_{i}^{\prime} X_{i}\right) /\left\{e_{i}^{\prime} e_{i}\right\} \tag{9.2.35}
\end{align*}
$$

and

$$
\begin{align*}
\operatorname{VIF}_{i}\left(u_{j}\right) & =\frac{V\left(X_{i} u_{j}\right)}{V\left(X_{i}\right)} \\
& =\frac{\left\{\left[X_{i} u_{j}\right]^{\prime}\left[X_{i} u_{j}\right]\right\}_{11}^{-1}}{\left(X_{i}^{\prime} X_{i}\right)^{-1}} \\
& =\left(X_{i}^{\prime} X_{i}\right) /\left(X_{i} X_{i}-x_{i j}^{2}\right) \quad(\text { from }(9.2 .33)) \tag{9.2.36}
\end{align*}
$$

where $\mathrm{x}_{\mathrm{ij}}$ are the j -th element of $\mathrm{X}_{\mathrm{i}}$. Then by inserting (9.2.34), (9.2.35) and (9.2.36), into (9.2.28), $\mathrm{R}_{\mathrm{ij}}$ can be written as

$$
\begin{aligned}
R_{i j} & =\frac{\operatorname{VIF}_{i}\left(\left[X_{-i} u_{j}\right]\right)}{\operatorname{VIF}_{i}\left(X_{-i}\right) \cdot \operatorname{VIF}_{i}\left(u_{j}\right)} \\
& =\frac{\left(X_{i}^{\prime} X_{i}\right)}{\left(e_{i}^{-j}\right)^{\prime}\left(e_{i}^{-j}\right)} \frac{e_{i}^{\prime} e_{i}}{X_{i} X_{i}} \frac{X_{i}^{\prime} X_{i}-x_{i j}^{2}}{X_{i} X_{i}}
\end{aligned}
$$

$$
\begin{align*}
& =\frac{e_{i}^{\prime} e_{i}}{\left(e_{i}^{-j}\right)^{\prime}\left(e_{i}^{-j}\right)} \frac{X_{i} X_{i}-x_{i j}^{2}}{X_{i}^{\prime} X_{i}} \\
& =\frac{e_{i}^{\prime} e_{i}}{e_{i}^{\prime} e_{i}-e_{i j}^{2} /\left(1-\left[P_{-i}\right]_{j j}\right)} \frac{X_{i}^{\prime} X_{i}-x_{i j}^{2}}{X_{i}{ }^{\prime} X_{i}}  \tag{9.2.37}\\
& =R_{i j}^{1} \cdot R_{i j}^{2}
\end{align*}
$$

where $\left[P_{-i}\right]_{j j}=X_{-i}\left(X_{-i}^{\prime} X_{-i}\right)^{-1} X_{-i}^{\prime}$ is the projection matrix onto the column space of the variables $X_{-i}$, and $e_{i j}$ is the $j$-th elements of $e_{i}$. Schall and Dunne (1987b) claim that the quantity ( $\left.\mathrm{R}_{\mathrm{ij}}^{1}-1\right) /(\mathrm{n}-\mathrm{p})$ follows an $\mathrm{F}(1, \mathrm{n}-\mathrm{p} ; \gamma)$ distribution under the assumption that the $\mathrm{x}_{\mathrm{ij}}, \mathrm{j}=1, \ldots, \mathrm{n}$ are independent and identically distributed normal variates. The F-distribution is central $(\gamma=0)$ when $X_{-i}$ includes an intercept. The quantity $\left(\mathrm{R}_{\mathrm{ij}}^{2}-1\right) /(\mathrm{n}-1)$ follows an $\mathrm{F}(1, \mathrm{n}-1 ; \gamma)$-distribution which in general is non-central when $X_{i}$ is not centered. Thus the two factors in (9.2.37) can be calibrated against these distributions. A similar factorization of the multiple case (9.2.26) is possible.

### 9.2.7 Condition index

The condition indices of a matrix are one of the measures proposed in Chapter 2 to detect collinearity. It is important to determine whether an index is high because of collinearity-inducing points or low because of collinearty masking points. The effect on the condition indices when a single row is deleted were studied by Walker (1989) and Hadi (1988), and based on this knowledge Hadi (1988) proposed two measures to detect influential points, one of which was also proposed by Valker (1989).

### 9.2.7.1. Row deletion and condition indices

Let the condition indices of the matrix $X_{-i}$ be denoted by

$$
\begin{equation*}
\eta_{\mathrm{j}}(-\mathrm{i}) \tag{9.2.38}
\end{equation*}
$$

where the subscript $j$ denotes the $j$-th condition index and (-i) is an indicator of the deleted row. If one wants to assess the influence of the i-th row, the condition indices of $X_{-i}$ can be computed and compared to the condition indices of the full matrix. One drawback of this approach is that it requires the computation of the eigenvalues of $(n+1)$ matrices each of order k×k. Hadi (1988) derives an approximation to $\eta_{j}(-i)$ without actually computing the eigenvalues of $\mathrm{X}_{-\mathrm{i}} \mathrm{X}_{-\mathrm{i}}, \mathrm{i}=1,2, \ldots, \mathrm{n}$. His development consist of two special cases (case one and two) and then the generalized case (case 3 ).

Case one: A theorem due to Kempthorne (1985), when $X$ is an $n \times 2$ matrix, is summarized by (Hadi (1988), p146):

Theorem one: If $X$ is $n \times 2$ and the columns of $X$ are normalized to have length 1 , then the square of the condition index of $X_{-i}$ is

$$
\begin{equation*}
\eta_{\mathrm{p}}(-\mathrm{i})=\frac{1+\left(1-4 / \mathrm{S}_{\mathrm{i}}\right)^{\frac{1}{2}}}{1-\left(1-4 / \mathrm{S}_{\mathrm{i}}\right)^{\frac{1}{2}}} \tag{9.2.39}
\end{equation*}
$$

where

$$
\begin{equation*}
S_{i}=\frac{\left[\left(1-x_{i 1}^{2}\right)+\left(1-x_{i 2}^{2}\right)\right]^{2}}{\left(1-x_{i 1}^{2}\right)\left(1-x_{i 2}^{2}\right)-\left(\xi-x_{i 1} x_{i 2}\right)^{2}} \tag{9.2.40}
\end{equation*}
$$

and the i-th row of $X$ is partitioned as

$$
\mathrm{x}_{\mathrm{i}}^{\prime}=\left[\begin{array}{ll}
\mathrm{x}_{\mathrm{i} 1} & \mathrm{x}_{\mathrm{i} 2}
\end{array}\right]
$$

and

$$
\begin{equation*}
\xi=\sum_{i=1}^{n} x_{i 1} x_{i 2} \tag{9.2.41}
\end{equation*}
$$

Second case: The second special case, due to Dorsett (1982), involves the i-th row of $X$ lying in the direction of the $j$-th eigenvector, and is summarized by (Hadi (1988), p147):

Theorem two: Let $\mathrm{v}_{1}, \mathrm{v}_{2}, \ldots, \mathrm{v}_{\mathrm{p}}$ be the set of orthonormal eigenvectors associated with the eigenvalues $\lambda_{1} \geq \lambda_{2} \geq \ldots \geq \lambda_{p}$ of $X^{\prime} X$. If $x_{i}=a v_{j}$, then the (possibly unordered) eigenvalues of $X_{-i}^{\prime} X_{-i}$ are

$$
\begin{equation*}
\lambda_{1}, \lambda_{2}, \ldots, \lambda_{\mathrm{j}-1}, \lambda_{\mathrm{j}}-\mathrm{a}^{2}, \lambda_{\mathrm{j}+1}, \ldots, \lambda_{\mathrm{p}} \tag{9.2.42}
\end{equation*}
$$

Thus when $x_{i}$ lies on the direction of the $j$-th eigenvector, the deletion of $x_{i}$ will deflate the $j$-th eigenvalue. If the conditions of Theorem 2 hold, the following three conclusions can be made:
(i) If $\mathrm{j}=1$, then $\left[\eta_{\mathrm{p}}(-\mathrm{i})\right]^{2}=\max \left(\lambda_{1}-\mathrm{a}^{2}, \lambda_{2}\right) / \lambda_{\mathrm{p}}$, and the deletion of $\mathrm{x}_{\mathrm{i}}$ will decrease the condition index,
(ii) If $j=p$, then $\left[\eta_{p}(-i)\right]^{2}=\lambda_{1} /\left(\lambda_{p}-a^{2}\right)$, and the deletion of $x_{i}$ will increase the condition index, and
(iii) If $1<j<p$, then $\left[\eta_{p}(-i)\right]^{2}=\lambda_{1} / \min \left(\lambda_{p} ; \lambda_{j}-a^{2}\right)$, and the deletion of $x_{i}$ will have no effect on the condition index as long as $\lambda_{j}-\lambda_{p}>a^{2}$

A similar theorem is also given by Walker (1989) as Corollary 1 on p1681, but apparently has a slight error. The author assumed that $\max \left(\lambda_{1}-\mathrm{a}^{2}, \lambda_{2}\right)=\left(\lambda_{1}-\mathrm{a}^{2}\right)$.

Third case: Suppose there are no restrictions on $X$ : for example, the restrictions of theorem one and two are lifted, $X$ may or may not contain a constant column, or may or may not be normalized or standardized.

Let $v_{1}, v_{2}, \ldots, v_{p}$ be the set of orthonormal eigenvectors associated with the eigenvalues $\lambda_{1} \geq \lambda_{2} \geq \ldots \geq \lambda_{p}$ of $X^{\prime} X$. Let $\gamma_{i 1} \geq \gamma_{i 2} \geq \ldots \geq \gamma_{i p}$ be the eigenvalues of $X_{-i}^{\prime} X_{-i}$. Define $z_{i j}=x_{i}^{\prime} v_{j}$ and

$$
c_{i}=\sum_{j=1}^{p}\left(z_{i j} / \lambda_{j}\right)^{2}
$$

then
(i) if $\mathrm{j}=1$

$$
\begin{equation*}
\gamma_{i 1} \simeq \frac{\lambda_{1}^{2}-2 \lambda_{1} z_{i 1}^{2}+z_{i 1}^{2} x_{i}^{\prime} x_{i}}{\lambda_{1}-z_{i 1}^{2}} \tag{9.2.43}
\end{equation*}
$$

(ii) if $1<j<p$

$$
\begin{equation*}
\gamma_{\mathrm{i} \mathrm{j}} \simeq \lambda_{\mathrm{j}}-\mathrm{z}_{\mathrm{i} \mathrm{j}}^{2} \tag{9.2.44}
\end{equation*}
$$

(iii) if $\mathrm{j}=\mathrm{p}$

$$
\begin{equation*}
\gamma_{i p} \simeq \min \left(A_{j}, A_{p}\right) \tag{9.2.45}
\end{equation*}
$$

where

$$
A_{j} \simeq \min \left(\lambda_{j}-z_{i j}^{2}\right), \text { for } j \neq p
$$

and

$$
A_{p} \simeq \frac{\lambda_{p}\left(1-h_{i i}\right)\left[\lambda_{p}\left(1-h_{i i}\right)+z_{i p}^{2}\right]}{\lambda_{p}\left(1-h_{i i}\right)^{2}+2 z_{i p}^{2}\left(1-h_{i i}\right)+z_{i p}^{2} \lambda_{p} c_{i}}
$$

(iv) the condition index of $X_{-i}$ can be approximated by

$$
\begin{equation*}
\left[\eta_{\mathrm{p}}(-\mathrm{i})\right]^{2} \simeq \gamma_{\mathrm{i}_{1}} / \gamma_{\mathrm{ip}} \tag{9.2.46}
\end{equation*}
$$

The results in (i) and (iii) are proven in Hadi (1988, Theorem five). In simulation studies the author finds that the approximations of the condition number in (iv) are close to the actual computed condition indixes. If the condition of theorem two holds, then (iv) holds exactly. Although the author could not prove (ii) theoretically, he uses an empirical investigation to suggest it is a good approximation.

### 9.2.7.2 Collinearity-influential points

To assess the influence of the $i$ - th row on the condition index of X , Hadi (1988) suggested the following two measures:

Define

$$
\begin{equation*}
\delta_{\mathrm{i}}=\frac{\left[\eta_{\mathrm{p}}(-\mathrm{i})\right]-\left[\eta_{\mathrm{p}}\right]}{\left[\eta_{\mathrm{p}}\right]}, \quad \mathrm{i}=1,2, \ldots, \mathrm{n} \tag{9.2.47}
\end{equation*}
$$

where $\left[\eta_{\mathrm{p}}(-\mathrm{i})\right]$ is defined in (9.2.46) and $\eta_{\mathrm{p}}$ is the ordinary condition index, defined in Chapter 2. Thus $\delta_{i}$ is an approximation to the relative distance between the condition index of $X_{-i}$ and $X$. If $\delta_{i} \gg 1$ then the i-th case is a row/point that masks collinearity, and if $\delta_{i} \ll 1$ then the i-th case is a point that induces collinearity.

Although Hadi does not suggest it, it is possible that the approximation of all the collinearity indices could be derived from the result in (9.2.44), and (9.2.46) be generalized from the p - th index to all the collinearity indices. Once $\delta_{i}$ is identified as unequal to one, and the other collinearity indices examined and their $\delta_{i}$ 's computed, and say $\delta_{35}$ (3-rd case deleted, 5 -th condition index) is relatively large, is the $x_{35}$ point the outlier, or can one say anything about the 5 -th variable? The above questions could be a field for further investigation.

The second measure proposed by Hadi (1988) to diagnose collinearityinfluential points involves diagonal elements of the Hat matrix expressed a little differently. The $Z$ matrix, whose elements are $z_{i j}=x_{i}^{\prime} v_{j}$ can be
expressed as $Z=X V$ (this is the same $Z$ as defined in Chapter 3). Thus the Hat matrix can be expressed as

$$
\begin{align*}
H & =X\left(X^{\prime} X\right)^{-1} X^{\prime} \\
& \left.=Z V^{\prime}\left(V D^{2} V^{\prime}\right)^{-1} V Z \quad \text { (by using the SVD of } X^{\prime} X\right) \\
& =Z D^{-2} Z \\
& =W W^{\prime} \tag{9.2.48}
\end{align*}
$$

where $W=Z D^{-1}$, then $h_{i i}$ can be expressed as

$$
\begin{equation*}
h_{i j}=\sum_{j=1}^{p} z_{i j}^{2} / \lambda_{j}=\sum_{j=1}^{p} w_{i j}^{2} \tag{9.2.49}
\end{equation*}
$$

The diagnostic measures $\delta_{i}$ and $h_{i i}$ are supplemented by plots to give the analyst a comprehensive picture of the eigenstructure of the data. Collinearty-influential points can easily be seen on a graphical display of $\delta_{i}$, such as stem and leaf displays. Pairwise scatterplots of the columns of $Z$, or $W$ can be drawn, the quantities $W_{i j}^{2}$ are referred to as the leverage components (LC), and the scatterplots of $W$ as the LC plots.

Walker (1989) proposed a diagnostic very similar to Hadi's second diagnostic ( $\mathrm{h}_{\mathrm{i}}$, (9.2.49)). By using the SVD of X, Walker (1989) expressed the Hat matrix as

$$
\begin{align*}
H & =X\left(X^{\prime} X\right)^{-1} X^{\prime} \\
& =U D V^{\prime}\left(V D^{2} V^{\prime}\right)^{-1} V D U^{\prime} \\
& =U U^{\prime} \tag{9.2.50}
\end{align*}
$$

which is the same as (9.2.48), as $W=Z D^{-1}=X V D^{-1}=U D V^{\prime} V D^{-1}=U$. The diagonal element of H as

$$
\begin{equation*}
\mathrm{h}_{\mathrm{i} i}=\sum_{\mathrm{j}=1}^{\mathrm{p}} \mathrm{u}_{\mathrm{ij}}^{2} \tag{9.2.51}
\end{equation*}
$$

Velleman and Ypeelar (1980), call the $u_{i j}$ 's the orthogonal leverage components. Diagnostics proposed by Walker (1989) to detect collinearity influential points, involve the set of squared elements of the matrix U

$$
\begin{equation*}
u_{\mathrm{ij}}^{2} \quad i=1, \ldots, n, \quad j=1, \ldots, p \tag{9.2.52}
\end{equation*}
$$

Walker (1989) suggests a cut-off value for high leverage is taken as $\mathrm{p} / \mathrm{n}$. Once 'high' $u_{i j}^{2}$ 's, say $u_{k j}^{2}$, are identified, he fits the model again with the k -th case deleted and compares the condition indices.

### 9.3 Influential Variables

Extending the notion of collinearity-influential points, Schall and Dunne (1987a) also define collinearity-influential variables.

Consider the following augmented model:

$$
\mathrm{Y}=\left[\begin{array}{ll}
\mathrm{X} & \mathrm{~A}
\end{array}\right]\left[\begin{array}{l}
\beta  \tag{9.3.1}\\
\lambda
\end{array}\right]+\epsilon
$$

where $A: n x k$ is an arbitrary set of variables not necessarily dummy variables. Typically, A would consist of multiplicative interaction terms of the variables already in the model, or extra polynomial regression terms instead of the $p$ already fitted, or the covariates in ANACOVA.

The decision whether to include the variables A is usually based on the F-statistic associated with the variables. Even if this F-test indicates non-significance, A should be included if it has an effect on the estimation of the parameters already in the model. Thus, when deciding whether a particular variable should be included in the model, Schall and Dunne (1987a) advocate the additional use of Cook's distance

$$
\begin{equation*}
\mathrm{C}_{\mathrm{A}}=\frac{(\hat{\beta}-\tilde{\beta})^{\prime} X^{\prime} X(\hat{\beta}-\tilde{\beta})}{\hat{\sigma}^{2}} \frac{\mathrm{n}-\mathrm{p}}{\mathrm{p}} \tag{9.3.2}
\end{equation*}
$$

where $\hat{\beta}$ is the OLSE of model (1.1) and $\tilde{\beta}$ is the LSE of $\beta$ in the augmented model (9.3.1). If the statistic (9.3.2) indicates high influence of the variables $A$, it should be included in the model even if the $F$-test indicates non-significance. The variance of the estimates will increase, but large bias, due to the high influence of the variables, will be removed.

A second influence measure proposed by them is based on the AP-statistic, let

$$
\begin{aligned}
& X_{1}^{*}=\left[\begin{array}{ll}
X & Y
\end{array}\right] \\
& X_{2}^{*}=\left[\begin{array}{lll}
X & A^{*} & Y
\end{array}\right] \text { with } A^{*}=A\left(A^{\prime} A\right)^{-\frac{1}{2}}
\end{aligned}
$$

then the AP-statistic is defined as:

$$
\begin{equation*}
\mathrm{AP}_{\mathrm{A}}=\left|\mathrm{X}_{2}^{*} \mathrm{X}_{2}^{*}\right| /\left|\mathrm{X}_{1}^{*} \mathrm{X}_{1}^{*}\right| \tag{9.3.3}
\end{equation*}
$$

Alternative and computational forms of (9.3.2) and (9.3.3) involve factorization as in (9.2.16) and (9.2.22), yielding

$$
\begin{equation*}
\left.A P_{A}=(\operatorname{RSS}-S S A) / R S S\right) \cdot\left|A\left(I-X\left(X^{\prime} X\right)^{-1} X^{\prime}\right) A\right| /\left|A^{\prime} A\right| \tag{9.3.4}
\end{equation*}
$$

and

$$
\begin{equation*}
C_{A}=\frac{\operatorname{SSA}}{\operatorname{RSS}}\left[\frac{\tilde{\lambda}^{\prime} A^{\prime} A \tilde{\lambda}}{S S A}-1\right] \frac{n-p}{p} \tag{9.3.5}
\end{equation*}
$$

where SSA is the extra sum of squares due to fitting $A$ after $X$. SSA is similar to the $Q_{k}$ defined in (9.2.12), $S S A=\hat{\epsilon}^{\prime} A\left\{A^{\prime}\left[I-X\left(X^{\prime} X\right)^{-1} X^{\prime}\right] A\right\}^{-1} A^{\prime} \hat{\epsilon}$, and $\tilde{\lambda}$ is the LSE of $\lambda$ under the model (9.3.1), $\tilde{\lambda}=\left\{A^{\prime}\left[I-X\left(X^{\prime} X\right)^{-1} X^{\prime}\right] A\right\}^{-1} A^{\prime} \hat{\epsilon}$.

The first factor in (9.3.5) is a monotonic function of the F-statistic, as in (9.2.13), and is thus a measure of the statistical significance of the
variables. The first factor of (9.3.4) is a decreasing function of the F-statistic. The second factors of (9.3.4) and (9.3.5) are measures of the potential influence of the variables A .

When $k=1$ ( $\mathrm{A}: \mathrm{nx} 1)$ the second factor of (9.3.5) simplifies to

$$
\begin{align*}
{\left[\frac{\tilde{\lambda}^{\prime} A^{\prime} A \tilde{\lambda}}{S S A}-1\right] } & =\frac{\hat{\epsilon}^{\prime} A\left\{A^{\prime}\left[I-X\left(X^{\prime} X\right)^{-1} X^{\prime}\right] A\right\}^{-1} A^{\prime} A\left\{A^{\prime}\left[I-X\left(X^{\prime} X\right)^{-1} X^{\prime}\right] A\right\}^{-1} A^{\prime} \hat{\epsilon}}{\hat{\epsilon}^{\prime} A\left\{A^{\prime}\left[I-X\left(X^{\prime} X\right)^{-1} X^{\prime}\right] A\right\}^{-1} A^{\prime} \hat{\epsilon}}-1 \\
& =A^{\prime} A /\left\{A^{\prime}\left[I-X\left(X X^{\prime} X\right)^{-1} X^{\prime}\right] A\right\}-1 \tag{9.3.6}
\end{align*}
$$

The authors define the quantity ( $\left.1-\left\{A^{\prime}\left[I-X\left(X^{\prime} X\right)^{-1} X^{\prime}\right] A\right\} / A^{\prime} A\right)$ as the leverage of the variable $A$, extending the notion of the leverage of an observation described in §1.7.

The VIF associated with the variable A in the model (9.3.1) is given by

$$
\begin{aligned}
\operatorname{VIF}_{A}(X) & =\frac{V[A X]}{V[A]} \\
& =\frac{\left[\begin{array}{ll}
A^{\prime} A & A^{\prime} X \\
X^{\prime} A & X^{\prime} X
\end{array}\right]_{11}^{-1}}{\left[A^{\prime} A\right]_{11}^{-1}}
\end{aligned}
$$

If $A: n x 1$, then

$$
\begin{equation*}
\operatorname{VIF}_{A}(X)=\frac{\left[A^{\prime} A\right]}{A^{\prime}\left[I-X\left(X^{\prime} X\right)^{-1} X^{\prime}\right] A} \tag{9.3.7}
\end{equation*}
$$

Note that the second factor of (9.3.4) and (9.3.5), and the leverage, are monotonic functions of the variance inflation factor. Thus, they are also monotonic functions of the collinearity index associated with the variable A. Cook's distance and the AP-statistic therefore contain collinearity measures.

The measure of the influence of the variable $A$ on the variance inflation factor is given by

$$
\begin{align*}
& R_{i,(A)}=\frac{V I F_{i}\left(X_{-i} \mid A\right)}{\operatorname{VIF}_{i}\left(X_{-i}\right)} \\
&\left.=\frac{V\left[X_{i}\right.}{} X_{-i} A\right] V\left[X_{i}\right] \\
& V\left[X_{i} A\right]\left[X_{i}\right]  \tag{9.3.8}\\
& V\left[X_{i}\right]
\end{align*}
$$

The retention of variables with high influence seems to contradict a common method of handling collinearity in a regression model, namely dropping variables from the model. The statistics (9.3.4) and (9.3.5) take into account the bias of dropping a variable as well as the influence it would have on the estimates. In the extreme case when $A$ is orthogonal to $X$ its 'dropping' will have no influence in the estimates. However when A is exactly or nearly collinear to some of the variables, its effect is confounded with a linear combination of the variables. Because collinearity is a group phenomenon (Stewart (1987)) it is dangerous to drop variables from the model. Using influence statistics, offers at least in some cases a method to 'regularize' a collinear design (Stewart (1987)). It should be empasised that the statistic (9.3.2) is a diagnostic and not a variable selection criterion.

To summarize the following guidelines of Schall and Dunne (1987a) are useful in applying Cook's distance for variables:

Statistic (9.3.2) has two roles: Firstly it should be used as a new type of collinearity diagnostic: it indicates not only the potential for harmful effects of collinearity, but is also quantifies the actual influence of a variable. By using the decomposition of (9.3.5) the computation of (9.3.2) requires no extra computations when F -tests and VIF's are computed.

Secondly the statistic (9.3.2) can also be used for model-checking purposes. A small value of (9.3.2) is reassuring and in the case of a large value, corrective action of some sort must be considered, like mixed regression techniques (Belsley et al. (1980, Chapter 4)).

### 9.4 Further Remarks on Influence

We finally present research fields that were not covered in $\S 9.1$ to $\S 9.3$.

### 9.4.1 Veighted Least Squares

An influential observation can be set artificially to zero by using weighted least squares. In weighted least squares the modified normal equation is

$$
\begin{equation*}
X^{\prime} W X \tilde{\beta}=X^{\prime} W Y \tag{9.4.1}
\end{equation*}
$$

where $W=\operatorname{Diag}\left[1, \ldots, 1, \mathrm{w}_{\mathrm{i}}, 1, \ldots 1\right]$. Then the i -th observation can be 'deleted' or 'downweighted' by making $w_{i}$ arbitrary small. For reference to these techniques see Belsley et al. (1980), Cook and Weisberg (1982) and Cook (1986).

### 9.4.2 Influence measures in Ridge Regression

The notion of an influential observation in OLSE is extended to RRE by authors including Walker and Birch (1988), Lichtenstein and Velleman (1983) and Chalton (1990). Walker and Birch (1988) show that when RRE is used, the influence of each case changes is a function of the shrinkage parameter $k$. This change is mainly because of the behaviour of the residual as a function of the $k$ value:

$$
\begin{align*}
\hat{\epsilon}_{R} & =Y-X \hat{\beta}_{R} \\
& =Y-X\left(X^{\prime} X+k I\right)^{-1} X^{\prime} X \hat{\beta} \quad(\text { from } 4.3 .1) \tag{9.4.1}
\end{align*}
$$

where $\hat{\epsilon}_{\mathrm{R}}$ is the residual vector when the estimation is with the RRE.

Therefore when using RRE one could not rely on the influence measures obtained for 0LSE $(k=0)$. The authors suggested that 'once the value of $k$ is determined (by any method), influence measures should be computed for that $k$ '. They then defined influence measures to use with RRE, by extending $\mathrm{h}_{\mathrm{i} \mathrm{i}}$, DFFITS and Cook's distance to include RRE.

The concept of subset selection in OLSE is extended to biased estimators by Hoerl et al. (1986). Their results indicate that there is potential in using biased estimation to select subsets. They recommend that when applying RR, suspected superfluous variables may be deleted. However they warn that insignificance in an LS model does not necessarily imply that a variable is superfluous, particularly with collinear data.

### 9.5 Summary

This chapter introduced the concept of influential points and variables. Methods of determining collinearity-influential points were defined and the concept of influence in ridge regression was introduced.

## Chapter 10

## SIMULATION STUDY

### 10.1 Introduction

The purpose of this simulation study is to compare the performances of 13 different biased estimators, as well as OLSE on the simulated data discussed in $\S 10.2$. In $\S 10.3$ we discuss the program, giving a summary of the performance efficiencies of different estimators and tabulated results in $\S 10.4$, and some comments on the results in $\S 10.5$.

### 10.2 Data

The simulation study of this thesis follows that of McDonald and Galarneau (1975) and Wichern and Churchill (1978). The data sets were obtained from Chalton (1990) who generated them for a simulation study in his Ph.D thesis.

Chalton (1990), considers a five parameter model, with a sample size of 30 and the predictor variables generated from the following relationship:

For $\mathrm{j}=1,2,3$ and $\mathrm{i}=1,2, \ldots, 30$

$$
\begin{equation*}
X_{i j}=\left(1-a_{1}^{2}\right)^{\frac{1}{2}} Z_{i j}+a_{1} Z_{i 6} \tag{10.2.1}
\end{equation*}
$$

For $\mathrm{j}=4,5$ and $\mathrm{i}=1,2, \ldots, 30$

$$
\begin{equation*}
X_{i j}=\left(1-a_{2}^{2}\right)^{\frac{1}{2}} Z_{i j}+a_{2} Z_{i 6} \tag{10.2.2}
\end{equation*}
$$

where
(i) $\quad Z_{i j}$ are independent $N(0,1)$ variates generated by the SAS-function RANNOR. The seeds were not recorded by Chalton(1990), as the RANNOR function derives these from the time clock of the computer.
(ii) The parameters $a_{1}$ and $a_{2}$ determine the degree of collinearity between the predictor variables: $a_{1}^{2}$ is the theoretical correlation between any pair of the variables $\mathrm{X}_{1}, \mathrm{X}_{2}$ and $\mathrm{X}_{3}, a_{1} a_{2}$ is the theoretical correlation between the variables $X_{1}, X_{2}, X_{3}$ and $\mathrm{X}_{4}$ or $\mathrm{X}_{5}$, and $\alpha_{2}^{2}$ is the theoretical correlation between $\mathrm{X}_{4}$ and $X_{5}$.

Five different combinations of $\alpha_{1}^{2}$ and $a_{2}^{2}$ were considered, and two choices (orientations) of $\beta$, suggested by Newhouse and Oman (1971), namely the eigenvectors corresponding to the largest and smallest eigenvalues, denoted by $\beta_{\mathrm{L}}$ and $\beta_{\mathrm{S}}$. For these 10 combinations three different values of $\sigma$ where considered, namely $0.01,1.0$ and 5.0 .

Chalton (1990) generated the Y-vector as sets of 30 data points from the model (10.2.3) for each of the $(10 \times 3=30)$ combinations of orientations, $a_{1}^{2}$ and $\alpha_{2}^{2}$ values, and variance values.

For $\mathrm{i}=1,2, \ldots, 30$

$$
\begin{equation*}
\mathrm{Y}_{\mathrm{i}}=\beta_{0}+\beta_{1} \mathrm{X}_{\mathrm{i} 1}+\beta_{2} \mathrm{X}_{\mathrm{i} 2}+\beta_{3} \mathrm{X}_{\mathrm{i} 3}+\beta_{4} \mathrm{X}_{\mathrm{i} 4}+\beta_{5} \mathrm{X}_{\mathrm{i} 5}+\epsilon_{\mathrm{i}} \tag{10.2.3}
\end{equation*}
$$

where the $X_{i j}$ are unstandardized, $\beta_{0}$ is zero, and the $\epsilon_{\mathrm{i}}$ are independent $N\left(0, \sigma^{2}\right)$ variates. For the simulation study in this chapter, $\beta_{0}=10$ because we added 10 to the $Y$ 's generated by the author. For each combination of $X$, $\beta$ and $\sigma, 100$ replications of the $(30 \times 1) \epsilon$-vector were generated

The eigenvalues and condition numbers corresponding to the five different combinations of $a_{1}^{2}$ and $a_{2}^{2}$ are shown in Table 10.1. and the coefficients of $\beta_{\mathrm{L}}$ and $\beta_{\mathrm{S}}$ are shown in Table 10.2.

### 10.3 Estimation Programs

The program to obtain the different estimates is given in appendix B. It was written in Fortran 5 and ran on a PC. Double precision was used throughout although we have found in trial runs that it did not make much of a difference. The $X$ matrix was first standardized before any calculations were performed, then the SVD was computed. After any particular standardized estimate was obtained we transformed back to unstandardized parameter estimates before calculating the particular statistic of interest, as discussed later.

The SVD and the OLSE's were computed by using the subroutine SVDCMP and SVBKSB of Press et al. (1985). To obtain all the biased estimates, SVBKSB was modified for each particular estimation procedure.

One set of data (for one response vector $Y$ ) was run for comparative purposes by Dr. D. Chalton on an IBM4381 using SAS. The results for 12 of the 14 estimators were basically the same as those obtained with Fortran (agreeing to 6 and even 8 digits in estimating the 4 leading decimal digits of the $\beta^{\prime}$ s. For the two estimators FGRPC1 and FGRPC2 (to be explained below), we found that the estimates differ from the third and sometimes the second digit.

To avoid dividing by zero in these last two estimation procedures (FGRPC1 and FGRPC2) the $\mathrm{k}_{\mathrm{i}}$ 's were flagged as soon as the delta's ( $\delta$ 's) became smaller than $10^{-10}$ and in the subroutine that calculates the estimators, the delta's was then set equal to zero.

The criterion of observed TMSE was used as a basis for comparing estimators, and the comparisons were effected by means of ratios.

The program in appendix $B$ was tailored to find the summary statistics given in $\S 10.4$ and in the final runs only these statistics are printed. For those users wanting the estimates the obvious print lines should be added and the program should be modified to suit their needs.

### 10.4 Simulation Results

The definitions of the 14 estimators used are found in previous Chapters, especially the summary tables in Chapter 8. In Table 10.3 various abbreviations are explained. The FPC estimators require further clarification. In Chapter 5 the fractions for the FPCI and FPCV estimators (5.2.15) and (5.2.17) were described as iterative estimators. Because Lee and Birch (1988) observed that a 1-step version of both (5.2.15) and (5.2.17) exhibited already improved estimation properties over other biased estimators we considered only the following non-iterative FPC estimators:

FGRPC1: one step version of FPCI estimator, in (5.2.15) the $\left[\hat{\delta}_{K}(t)\right]_{j}^{2}$ is replaced by $\left[\hat{\delta}_{\mathrm{PC}}\right]_{\mathrm{j}}$, where the vector $\hat{\delta}$ consists of the $\hat{\delta}$ 's of the PC1 estimator, and where $s^{2}$ is the estimate of $\sigma^{2}$, using the PC1 estimator. Basically then FGRPC1 is the GRRE, where $\mathrm{k}_{\mathrm{i}}$ is estimated with the $\hat{\delta}^{\prime} \mathrm{s}$ and $\mathrm{s}^{2}$ of the PC1 estimators instead of the OLSE's as used by Hoerl, Kennard and Baldwin.

FGRPC2: the same as FGRPC1, but instead of using known PC1 estimates, we used the known PC2 estimates.

FRPC1 and FRPC2: similar to the ridge estimation procedure, but instead of using Hoerl, Kennard and Baldwin's $k$, we now estimate $k$ and $s^{2}$ using the PC1 and PC2 estimators. These procedures are one step versions of (5.2.17).

The performance of each estimator over the thirty combinations and the 100 repetitions was summarized by computing

$$
\begin{equation*}
\sum_{\mathrm{j}=1}^{5} \sum_{\mathrm{i}=1}^{100}\left(\tilde{\beta}_{\mathrm{ji}}-\beta_{\mathrm{j}}\right)^{2} \tag{10.4.1}
\end{equation*}
$$

where $\tilde{\beta}_{\mathrm{ji}}$ is the j -th element of $\tilde{\beta}$, the estimate of $\beta$ in the i -th repetition. The results of this simulation study are given in Tables 10.4, 10.5 , and 10.6 for $\sigma=0.01,1.0$ and 5.0 respectively. The evaluations are
based on the relative efficiency of each estimator compared to $\hat{\beta}$. Thus the tabulated relative efficiency values are

$$
\begin{equation*}
\sum_{j=1}^{5} \sum_{i=1}^{100}\left(\hat{\beta}_{j i}-\beta_{j}\right)^{2} / \sum_{j=1}^{5} \sum_{i=1}^{100}\left(\tilde{\beta}_{j i}-\beta_{j}\right)^{2} \tag{10.4.2}
\end{equation*}
$$

where $\tilde{\beta}$ is one of the estimators given in Table 10.3. Entries marked with a **, are very small values ( $\leq 0.0007$ ).

Table 10.1: Eigenvalues and condition numbers of $\mathbb{I}^{\prime} \mathbf{X}$ ( $\mathbb{X}$ is standardized).

| Correlations <br> $a_{1}^{2}: a_{2}^{2}$ | eigenvalues of $\mathrm{X}^{\prime} \mathrm{X}: \lambda_{\mathrm{i}}$ <br> (without $\beta_{0}$ ) | condition number <br> $\lambda_{1} / \lambda_{5}$ |
| :--- | :--- | :---: |
|  |  |  |
| $.99: .99$ | $(4.920,0.026,0.021,0.013,0.011)$ | 435 |
| $.99: .10$ | $(3.157,1.128,0.668,0.031,0.016)$ | 201 |
| $.90: .90$ | $(4.215,0.430,0.154,0.126,0.075)$ | 56 |
| $.90: .10$ | $(2.755,1.215,0.173,0.168,0.148)$ | 19 |
| $.70: .30$ | $(2.283,1.049,0.871,0.496,0.301)$ | 8 |

Table 10.2: $\beta$ used in generating $Y$

| Correlations$a_{1}^{2}: a_{2}^{2}$ | $\beta^{\prime}$ | Eigenvectors of $\mathrm{X}^{\prime} \mathrm{X}$ |  |  |  |  |
| :---: | :---: | :---: | :---: | :---: | :---: | :---: |
|  |  |  |  |  |  |
| . $99: .99$ | $\beta_{\mathrm{L}}^{\prime}$ | [ 0.4474 | 0.4473 |  | 0.4481 | 0.4470 | $0.4463]$ |
|  | $\beta_{\mathrm{S}}^{\prime}$ | [ 0.2846 | 0.4760 | -0.8302 | 0.0548 | $0.0163]$ |
| . $99: .10$ | $\beta_{\mathrm{L}}^{\prime}$ | [ 0.5534 | 0.5542 | 0.5510 | 0.1705 | $0.2323]$ |
|  | $\beta_{S}^{\prime}$ | [-0.7755 | 0.1675 | 0.6084 | 0.0190 | -0.0094] |
| . $90: .90$ | $\beta_{\mathrm{L}}^{\prime}$ | [ 0.4125 | 0.4547 | 0.4649 | 0.4383 | $0.4636]$ |
|  | $\beta_{\mathrm{S}}^{\prime}$ | [ 0.1821 | -0.3973 | 0.5673 | 0.3037 | -0.6284] |
| . $90: .10$ | $\beta_{L}^{\prime}$ | [ 0.5634 | 0.5489 | 0.5644 | 0.2354 | $0.0862]$ |
|  | $\beta_{\mathrm{S}}^{\prime}$ | [-0.0162 | 0.6908 | -0.7051 | 0.0655 | $0.1451]$ |
| . $70: .30$ | $\beta_{\mathrm{L}}^{\prime}$ | [ 0.5177 | 0.5611 | 0.5226 | 0.0286 | $0.3785]$ |
|  | $\beta_{\mathrm{S}}^{\prime}$ | [-0.6600 | 0.7049 | 0.0099 | 0.1956 | -0.1708] |

Table 10.3 Estimators and abbreviations

0LS Ordinary least squares estimator
RHK Ridge regression, k is estimated via Hoerl, Kennard and Baldwin
RLW Ridge regression, k is estimated via Lawless and Wang
GRHK Generalized ridge regression, $K$ is estimated via Hoerl, Kennard and Baldwin
GRCAS Generalized ridge regression, $K$ is estimated via Troskie
AUGRR almost unbiased generalized ridge regression
AUORR almost unbiased operational ridge regression
PC1 principal component regression, delete the smallest singular value
PC2 principal component regression, delete the two smallest singular values
FGRPC1 one step version of FPCI estimator, where the OLSE is replaced by the PC1 estimator.
FGRPC2 one step version of FPCI estimator, where the OLSE is replaced by the PC2 estimator.
FRPC1 one step version of FPCV estimator, where the OLSE is replaced by the PC1 estimator.
FRPC2 one step version of FPCV estimator, where the OLSE is replaced by the PC2 estimator.
SHE shrinkage estimator

Table 10.4 Relative efficiencies of biased estimators to 0LSE ( $\sigma=0.01$ )

Combinations of correlations and orientations of the $\beta$ 's

| $a_{1}^{2}: a_{2}^{2}$ | $99: 99$ |  | $99: 10$ |  | $90: 90$ |  | $90: 10$ |  | $70: 30$ |  |
| :--- | :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: |
| $\beta$ | $\beta_{\mathrm{L}}$ | $\beta_{\mathrm{S}}$ | $\beta_{\mathrm{L}}$ | $\beta_{\mathrm{S}}$ | $\beta_{\mathrm{L}}$ | $\beta_{\mathrm{S}}$ | $\beta_{\mathrm{L}}$ | $\beta_{\mathrm{S}}$ | $\beta_{\mathrm{L}}$ | $\beta_{\mathrm{S}}$ |
| OLS | 1.00 | 1.00 | 1.00 | 1.00 | 1.00 | 1.00 | 1.00 | 1.00 | 1.00 | 1.00 |
| RHK | 1.00 | 0.99 | 1.00 | 0.99 | 1.00 | 1.00 | 1.00 | 1.00 | 1.00 | 1.00 |
| RLW | 1.00 | 0.07 | 1.00 | 0.05 | 1.00 | 0.94 | 1.00 | 0.99 | 1.00 | 1.00 |
| GRHK | 1.35 | 1.33 | 1.37 | 1.03 | 1.05 | 0.99 | 0.84 | 0.93 | 0.97 | 0.96 |
| GRCAS | 1.32 | 1.25 | 1.33 | 1.05 | 1.04 | 1.01 | 0.89 | 0.97 | 0.97 | 1.00 |
| AUGRR | 1.11 | 1.14 | 1.10 | 1.00 | 1.04 | 0.98 | 0.93 | 0.95 | 1.00 | 0.96 |
| AUORR | 1.33 | 0.93 | 1.38 | 0.89 | 1.00 | 0.99 | 1.01 | 1.00 | 1.00 | 1.00 |
| PC1 | 0.78 | $* *$ | 2.49 | $* *$ | 0.03 | $* *$ | 0.24 | $* *$ | 0.04 | $* *$ |
| PC2 | 0.91 | $* *$ | 1.56 | $* *$ | 0.03 | $* *$ | 0.05 | $* *$ | 0.01 | $* *$ |
| FGRPC1 | 0.96 | $* *$ | 2.41 | $* *$ | 0.03 | $* *$ | 0.24 | $* *$ | 0.04 | $* *$ |
| FGRPC2 | 1.01 | $* *$ | 1.56 | $* *$ | 0.03 | $* *$ | 0.05 | $* *$ | 0.01 | $* *$ |
| FRPC1 | 1.00 | $* *$ | 1.00 | $* *$ | 1.00 | $* *$ | 1.00 | $* *$ | 1.00 | $* *$ |
| FRPC2 | 1.00 | $* *$ | 1.00 | $* *$ | 1.00 | $* *$ | 1.00 | $* *$ | 1.00 | $* *$ |
| SHE | 1.00 | 1.00 | 1.00 | 1.00 | 1.00 | 1.00 | 1.00 | 1.00 | 1.00 | 1.00 |

Table $\mathbf{1 0 . 5}$ Relative efficiencies of biased estimators to 0LSE ( $\sigma=1.0$ )

| $a_{1}^{2}: a_{2}^{2}$ | 99:99 |  | 99:10 |  | 90:90 |  | 90:10 |  | 70:30 |  |
| :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: |
| $\beta$ | $\beta_{\text {L }}$ | $\beta_{\mathrm{S}}$ | $\beta_{\text {L }}$ | $\beta_{\mathrm{S}}$ | $\beta_{\text {L }}$ | $\beta_{\text {S }}$ | $\beta_{\text {L }}$ | $\beta_{\text {S }}$ | $\beta_{\mathrm{L}}$ | $\beta_{\text {S }}$ |
| 0LS | 1.00 | 1.00 | 1.00 | 1.00 | 1.00 | 1.00 | 1.00 | 1.00 | 1.00 | 1.00 |
| RHK | 3.43 | 2.88 | 4.52 | 3.08 | 2.49 | 1.45 | 2.13 | 1.14 | 1.47 | 0.86 |
| RLW | 6.45 | 7.65 | 10.82 | 5.20 | 1.88 | 1.30 | 1.73 | 0.91 | 1.28 | 0.65 |
| GRHK | 2.23 | 1.98 | 2.16 | 1.76 | 2.03 | 1.39 | 2.12 | 1.26 | 1.94 | 1.33 |
| GRCAS | 1.75 | 1.65 | 1.72 | 1.54 | 1.66 | 1.35 | 1.70 | 1.25 | 1.61 | 1.25 |
| AUGRR | 1.39 | 1.32 | 1.36 | 1.23 | 1.31 | 1.12 | 1.34 | 1.08 | 1.30 | 1.16 |
| AUORR | 4.26 | 3.21 | 6.31 | 3.61 | 3.75 | 1.47 | 3.22 | 1.13 | 1.98 | 0.76 |
| PC1 | 1.59 | 1.41 | 2.90 | 1.79 | 1.79 | 0.67 | 1.78 | 0.47 | 1.48 | 0.24 |
| PC2 | 2.85 | 2.24 | 49.55 | 4.67 | 2.95 | 0.83 | 5.15 | 0.57 | 2.32 | 0.25 |
| FGRPC1 | 3.58 | 2.73 | 6.05 | 2.67 | 3.83 | 0.85 | 3.61 | 0.55 | 2.76 | 0.26 |
| FGRPC2 | 6.66 | 3.86 | 93.90 | 4.90 | 6.08 | 0.96 | 9.34 | 0.61 | 3.95 | 0.27 |
| FRPC1 | 5.05 | 4.21 | 10.36 | 4.26 | 3.46 | 1.24 | 2.55 | 0.79 | 1.57 | 0.37 |
| FRPC2 | 9.07 | 5.53 | 40.05 | 5.07 | 4.15 | 1.18 | 3.15 | 0.68 | 1.64 | 0.33 |
| SHE | 2.52 | 2.48 | 2.02 | 2.01 | 1.60 | 1.56 | 1.41 | 1.23 | 1.20 | 1.04 |

Table 10.6 Relative efficiencies of biased estimators to 0LSE ( $\sigma=5.0$ )

Combinations of correlations and orientations of the $\beta^{\prime}$ s

| $a_{1}^{2}: \alpha_{2}^{2}$ | $99: 99$ |  | $99: 10$ |  | $90: 90$ |  | $90: 10$ |  | $70: 30$ |  |
| :--- | ---: | ---: | ---: | ---: | ---: | ---: | :---: | ---: | :---: | :---: | :---: |
| $\beta$ | $\beta_{\mathrm{L}}$ | $\beta_{\mathrm{S}}$ | $\beta_{\mathrm{L}}$ | $\beta_{\mathrm{S}}$ | $\beta_{\mathrm{L}}$ | $\beta_{\mathrm{S}}$ | $\beta_{\mathrm{L}}$ | $\beta_{\mathrm{S}}$ | $\beta_{\mathrm{L}}$ | $\beta_{\mathrm{S}}$ |
| OLS | 1.00 | 1.00 | 1.00 | 1.00 | 1.00 | 1.00 | 1.00 | 1.00 | 1.00 | 1.00 |
| RHK | 3.40 | 3.81 | 4.49 | 5.59 | 3.64 | 3.66 | 3.35 | 3.57 | 2.59 | 2.43 |
| RLW | 360.75 | 162.19 | 79.16 | 61.83 | 12.68 | 14.59 | 6.61 | 7.16 | 2.99 | 3.05 |
| GRHK | 1.98 | 2.09 | 1.97 | 2.12 | 2.14 | 2.20 | 1.98 | 2.21 | 1.83 | 1.84 |
| GRCAS | 1.64 | 1.70 | 1.64 | 1.71 | 1.71 | 1.74 | 1.64 | 1.76 | 1.57 | 1.58 |
| AUGRR | 1.31 | 1.32 | 1.29 | 1.31 | 1.35 | 1.36 | 1.30 | 1.38 | 1.25 | 1.28 |
| AUORR | 4.14 | 4.56 | 5.95 | 7.40 | 4.21 | 4.09 | 3.91 | 3.94 | 2.95 | 2.62 |
| PC1 | 1.65 | 1.58 | 2.45 | 3.13 | 1.77 | 1.66 | 1.64 | 1.65 | 1.45 | 1.18 |
| PC2 | 3.41 | 3.19 | 50.07 | 35.37 | 3.32 | 3.30 | 5.18 | 3.75 | 2.50 | 1.97 |
| FGRPC1 | 3.41 | 3.38 | 4.47 | 6.40 | 3.83 | 3.56 | 3.12 | 3.34 | 2.59 | 2.00 |
| FGRPC2 | 7.46 | 7.39 | 85.81 | 56.71 | 7.18 | 6.91 | 9.03 | 6.24 | 4.28 | 3.10 |
| FRPC1 | 6.05 | 6.48 | 9.93 | 14.43 | 6.50 | 6.13 | 5.02 | 5.37 | 3.57 | 3.02 |
| FRPC2 | 14.81 | 16.22 | 128.40 | 75.83 | 12.37 | 10.80 | 12.17 | 8.86 | 5.33 | 4.36 |
| SHE | 2.70 | 2.99 | 2.39 | 2.65 | 2.87 | 2.99 | 2.57 | 2.94 | 2.40 | 2.48 |

### 10.5 Discussion of the results

Features of Tables 10.4-10.6 are:

1. From Table 10.4 we note that the OLSE performs satisfactorily when the collinearity is modest $(<56)$. As the collinearity increases a slight improvement over OLS was obtained by GRHK, GRCAS, AUGRR and AUORR (except for the second orientation of $\beta$ ). The PC1, PC2, FGRPC1, FGRPC2, FRPC1 and FRPC2 perform disastrously for the 2 -nd orientation of the $\beta^{\prime}$ s (the $\beta^{\prime}$ s corresponding to the smallest eigenvalue of $X^{\prime} X$ ). Only for the 99:10 combination did PC1, PC2, FGRPC1, FGRPC2, FRPC1 and FRPC2 perform clearly better than the OLS.
2. As $\sigma$ increases the relative efficiency of estimators increased with respect to the OLSE except for combinations 90:90, $90: 10$ and 70:30, under orientation 2 , in Table 10.5, where we again picked up the disastrous pattern of Table 10.4.
3. We observed the same pattern that Lee and Birch (1988) noted: The efficiencies of FPC estimators are proportional to the efficiencies of the PC1 and PC2 estimates. Furthermore when $\sigma \geq 1.0$ the fractional generalized ridge estimators (FGRPC1 and FGRPC2) outperformed the generalized ridge estimators. The fractional ridge estimators (FRPC1 and FRPC2) were also better relative to OLSE than the RHK estimators but for the RLW estimators a 'see-sab' situation is noted.
4. When $\sigma=5$ some of the biased estimators (e.g. RLW and FRPC2) appears to be exceptionally good. Further investigation with respect to the different orientations and combinations may explain this phenomenon.
5. The AUORR estimator generally performs better than RHK for $\sigma=1.0$ and 5.0. Contradicting Nomura (1988) we found that AUORR is not better than RLW for $\sigma=1$. When $\sigma=5$ RLV outperforms AUORR relative to OLS. RLW was always better (except for 3 orientations) than RHK relative to 0LS.
6. We found that the AUGRR was always inferior to the GRRE (both GRHK and GRCAS). This phenomenon was also found by Ohtani (1986).
7. The shrinkage estimators were always similar to or relatively better than OLSE, but rather inferior to other biased estimators.
8. In conclusion then: among the biased estimators, no single estimator appears to dominate the others, in terms of TMSE, across all conditions.

### 10.6 Recommendations and further research

According to Belsley et al. (1980) the above collinearity indices ( $\sqrt{435}$, $\sqrt{201}, \ldots, \sqrt{8})$, are generally weak with perhaps one being moderate. For further research we recommend that at least one higher collinearity index (say near 100) should be included. Though the results obtained in this simulation suggest the biased estimators are more efficient than 0LS, we seek to compare such estimators under a variety of combinations.

The extremely small values in Table 10.4 and 10.5 should be explored as well as the extremely high values in Table 10.6. Furthermore fixing the $\epsilon_{\mathrm{i}}$ vector across various $\sigma$ 's, the sampling distributions of $\beta_{i}$ for each $\sigma$, sampling distributions of $\sum_{j=1}^{5}\left(\hat{\beta}_{\mathrm{j}}-\beta_{\mathrm{j}}\right)^{2}$ for each $\sigma$ and the estimated covariance of $\hat{\beta}_{i}$ should be investigated.

In this simulation study we do not report $\beta_{0}$, as $\beta_{0}$ is simply estimated as the mean of the $Y$ vector in the case of the biased estimators, following a practice used elsewhere in the literature.

### 10.7 Summary

This chapter consisted of a simulation study and its general findings. The data, the program and the results were given and discussed, and areas for further research were described. Generally it appears that no biased class of estimators consistently outperforms other classes on the criterion of relative efficency. Specific conditions seem to be associated with the optimality of specific estimators, but it is not yet possible to define those conditions.

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## Appendix A

## USEFUL FORMOLAE AND DERIVATIONS

A. 1 If $\mathrm{x} \sim \mathrm{N}(\mu, \mathrm{V})$ then for A symmetric and conformable

$$
\begin{aligned}
& \mathrm{E}\left(\mathrm{x}^{\prime} \mathrm{Ax}\right)=\operatorname{tr}(\mathrm{AV})+\mu^{\prime} \mathrm{A} \mu \\
& \mathrm{~V}\left(\mathrm{x}^{\prime} \mathrm{Ax}\right)=2 \operatorname{tr}(\mathrm{AV})^{2}+4 \mu^{\prime} \mathrm{AVA} \mu
\end{aligned}
$$

(Searle (1971, pp55-57))
A. 2 Let $A$ be an $n x n$ matrix with eigenvalues $\lambda_{1}, \lambda_{2}, \ldots \lambda_{n}$, then

$$
\operatorname{tr}(\mathrm{A})=\Sigma \lambda_{\mathrm{i}} .
$$

If $A$ is symmetric and $\lambda_{i}>0 \quad \forall_{i}$
then $\operatorname{tr}\left(A^{-}\right)=\Sigma \lambda_{i}{ }^{-1}$
(Graybill (1969, pp223-225))
A. 3 Let $A$ be an nxn symmetric matrix with eigenvalues

$$
a_{1} \geq a_{2} \cdots \geq a_{\mathrm{n}}
$$

Let $k$ be an integer, $1 \leq k \leq n$. Let $B$ be the ( $n-1$ ) $x(n-1)$ symmetric matrix obtained by deleting the $k$-th row and column from $A$. Then the ordered eigenvalues $\beta_{\mathrm{i}}$ of B interlace with those of A as follows:

$$
a_{1} \geq \beta_{1} \geq a_{2} \geq \beta_{2} \geq \cdots \geq \beta_{\mathrm{n}-1} \geq a_{\mathrm{n}}
$$

(Lawson and Hanson (1974, p24) and the proof can be found in Wilkinson (1965, pp99-109)).
A. 4 Let $A=\left[a_{1} \ldots a_{p}\right]$ be $a$ column partitioning of $A: m x p$, and denote the i-th singular value of $A$ by $\sigma_{i}(A)$. If $A_{r}=\left[a_{1} \ldots a_{r}\right]$ then for $r=1, \ldots, p-1$
$\sigma_{1}\left(\mathrm{~A}_{\mathrm{r}+1}\right) \geq \sigma_{1}\left(\mathrm{~A}_{\mathrm{r}}\right) \geq \sigma_{2}\left(\mathrm{~A}_{\mathrm{r}+1}\right) \geq \ldots \geq \sigma_{\mathrm{r}}\left(\mathrm{A}_{\mathrm{r}+1}\right) \geq \sigma_{\mathrm{r}}\left(\mathrm{A}_{\mathrm{r}}\right) \geq \sigma_{\mathrm{r}}\left(\mathrm{A}_{\mathrm{r}+1}\right)$

Equivalently, adding a column to a matrix increases the largest singular value and diminishes the smallest (Golub and Van Loan (1983, p286).

## Appendix B

## PROGRAI

The program code is given in this appendix. Some of the comments and documentation were added using $\mathrm{T}^{3}$ to relate the programming to the development of the theory in Chapters 1 to 10 . The program-code is printed in bold and the comments unbold. Various abbreviations for the estimators were summarized in Table 10.3 of Chapter 10 . It is suggested that $\S 10.3$ should be read with this program.


C Declaration of variables

```
parameter (maxr=30,maxc=7) physical row/column dimension
character*14 FNAME input data file
character*14 OJTFNA output data file
double precision XSIM (30,105) simulated data matrix
double precision X(maxr,maxc) unstd X, without a column of ones
```

```
double precision TEMPX(maxr,maxc) unstd X, with a column of ones
double precision STDX(maxr,maxc) std X matrix (correlation form)
double precision Y(maxr) Y-vector
double precision U(maxr,maxc), W(maxc), V(maxc,maxc) SVD std X
double precision UU(maxr,maxc), WV(maxc), UV(maxc,maxc) SVD unstd X
double precision TV(maxc) }\mp@subsup{\lambda}{i}{}\mathrm{ 's, smallest set equal to 0 for PC
double precision DELTA(maxc) \hat{\delta}
double precision B(maxc) std \hat{\beta}
double precision BO(maxc) unstd \hat{\beta}
double precision AVE(maxc) vector - means of columns of X
double precision SS(maxc) vector - sums of squares of columns of X
double precision C(maxr) vector send to SVDSORT
double precision BB(maxc) Vector send to SVDSORT
double precision RES(maxr)
double precision KPC1(maxc),KPC2(maxc) \hat{K}, via PC1/PC2
double precision KHX(maxc) \hat{K}}\mathrm{ via HКВ
double precision KCAS(maxc) \hat{K}}\mathrm{ via Troskie
double precision sse, ssepc1, ssepc2 \hat{\sigma}}\mathrm{ via 0LS/PC1/PC2
double precision TEMPB(maxc)
double precision faclw \hat{k}}\mathrm{ via Lawless and Nang
double precision fachk \hat{k}}\mathrm{ via Hoerl, Kennard and Baldwin
double precision hkpc2, hkpc1 \hat{k}}\mathrm{ via PC2/PC1
double precision nomura \hat{k}}\mathrm{ via Nomura
double precision sinf flag, marking infinity
double precision meany mean of Y
double precision F F used in Troskie's method
double precision shrink shrinkage factor
double precision sum
double precision stemp 的'\hat{\beta}
double precision tmse trace of variance of OLSE
double precision tbeta(maxc) }
```

```
double precision sols(maxc), sauorr(maxc)
double precision srhk(maxc), srlw(maxc)
double precision sgrhk(maxc), sgrcas(maxc)
double precision saugrr(maxc), ssh(maxc)
double precision spc1(maxc), spc2(maxc)
```



```
    j=1
for i = 1,...,5
double precision sfgrp1(maxc), sfgrp2(maxc)
double precision sfrpc1(maxc), sfrpc2(maxc)
double precision tols, trhk
double precision rlw, tgrhk
double precision tgrcas, taugrr
double precision tpc1, tpc2
double precision tfgrp1, tfgrp2
double precision tfrpc1, tfrpc2
double precision tauorr, tssh
real tsigma }\sigma=0.1\mathrm{ or 1.0 or 5.0
integer N,P,X,1,M,R1,C1,R2,C2,R,rep
```

C Read in file names

```
write (*,*) 'Enter filename as name.dat, 1-14 characters'
```

read (*,85) FNAME
write (*,*) 'Enter output filename as name.out, 1-14 characters'
read (*,85) OUTPNA

C Set logical dimension of matrix, and define value near infinity
C Note: change here for new data sets, or change program to read in
C these dimensions.

$$
\begin{aligned}
& N=30 \\
& P=6 \\
& \operatorname{sinf}=1.0 d-10
\end{aligned}
$$

C Open files, output file exist, change here if either is a new file

```
open (15, file=FNANE, status='OLD')
open (16, file=00TFNA, status='OLD')
```

C Read true sigma and betas, and write to output file

```
write(16,85) OUTFNA
read (15,*) tsigma
write(16,*) 'This are the results of ', PNANE
write(16,*) 'sigma =', tsigma
read(15,*) (tbeta(i),i = 2,6)
write(16,*) 'True betas'
write(16,*) (tbeta(i),i = 2,6)
tbeta(1) = dble(10.0)
```

C Set trace sums for the 14 estimators equal to zero

```
do }10\textrm{i}=1,\textrm{p
    sols(i) = dble(0.0)
    srhk(i) = dble(0.0)
    srlw(i) = dble(0.0)
    sgrhk(i) = dble(0.0)
    sgrcas(i) = dble(0.0)
    saugrr(i) = dble(0.0)
    sauorr(i) = dble(0.0)
    spc1(i) = dble(0.0)
    spc2(i) = dble(0.0)
    sfgrp1(i) = dble(0.0)
    sfgrp2(i) = dble(0.0)
    sfrpc1(i) = dble(0.0)
    sfrpc2(i) = dble(0.0)
    ssh(i) = dble(0.0)
continue
```

C Read in simulated data matrix

$$
\begin{aligned}
& \text { do } 90 \mathrm{i}=1, \mathrm{~N} \\
& \quad \operatorname{read}\left(15,^{*}\right)(\operatorname{XSIM}(\mathrm{i}, 1), 1=1,105) \\
& \text { continue }
\end{aligned}
$$

90

C Idd a column of ones to the TEIPX matrix and copy the other columns to $\mathbb{X}$

```
do 92 i = 1,N
    TEMPX(i,1) = dble(1.0)
            do 91 j = 2,P
                TEIPX(i,j) = XSIM(i,(100 + j - 1))
                X(i,j-1) = TEIPX (i,j)
            continue
continue
```

c Make a copy of TEXPX into 0

CALL COPY (TEMPX, UU, $\operatorname{maxr}, \operatorname{maxc}, \mathrm{N}, \mathrm{P}$ )
c Compute the SVD, unsorted

CALL SVDCNP ( $\mathrm{UO}, \mathrm{N}, \mathrm{P}$, maxr, maxc, UV, UV)
write ( $16,{ }^{*}$ ) 'SVDCYP COMPUTED UNSTANDARDIZED MaTRIX'

C Sort the SVD

CALL SVDSORT( $\mathrm{UU}, \mathrm{UV}, \mathrm{UV}, \mathrm{N}, \mathrm{P}, \operatorname{maxr}, \max \mathrm{C}, \mathrm{C}, \mathrm{BB}$ )

C Standardize the matrix, and get back STDX
$\mathrm{IP}=\mathbf{P}$
$\mathbf{P}=\mathbf{P}-\mathbf{1}$
CALL STAND(X, $\max r, \max , N, P, \operatorname{STDX}$, AVE, SS)

C Copy STDI into U

$$
\text { CALL COPY(STDX, } \mathrm{V}, \operatorname{maxr}, \operatorname{maxc}, \mathrm{~N}, \mathrm{P})
$$

C Decompose matrix J

CALL SVDCIP ( $\mathbf{O}, \mathrm{N}, \mathrm{P}, \max , \max , \mathrm{V}, \mathrm{V})$

C Sort the SVD

```
CALL SVDSORT(J, V, V, N, P, maxr, maxc, C, BB)
write(16,*)'singular values of std \'
write (16,*) (V(j), j=1,P)
```

C Split the $Y$ and $\mathbb{X}$ part in XSIM

C SVD is computed, pick up all the Y's ( 100 Y -vectors), and calculate the
C corresponding estimators, and efficiency totals

```
do 500 rep \(=1,100\)
    meany \(=\) dble (0.0) \(\quad\) * calculate the mean of \(Y\)
        do \(94 \mathrm{i}=1, \mathrm{~N}\)
            Y(i) \(=\) ISIY (i,rep) + dble(10.0)
            meany = meany + Y(i)
94 continue
    meany = meany/dble(N)
```

C Compute the OLS estimator, using the SVD of the unstandardized $X$ matrix
CALL SVBKSB(UU, UV, UV, N, IP, $\operatorname{maxr}, \operatorname{maxc}, \mathrm{Y}, \mathrm{B})$
do $100 \mathrm{i}=1, \mathrm{IP}$
$\operatorname{sols}(i)=\operatorname{sols}(i)+(B(i)-\operatorname{tbeta}(i))^{* *} 2$
100
continue

C Compute the estimate of sigma (0LS), unstandardized $\mathbb{X}$

```
CALL SIGMA(TEYPX, maxr, maxc, N, IP, Y, B, RES, sse)
```

C Compute shrinkage estimator from unstandardized data using the $d$ in C equation (5.1.10)

$$
\text { tmse }=\operatorname{DBLE}(0.0)
$$

$$
\text { stemp }=\operatorname{DBLE}(0.0)
$$

do $110 \mathrm{i}=2, \mathrm{IP}$
stemp $=$ stemp $+B(i)^{* *} 2$
tmse $=$ tmse $+\operatorname{DBLE}(1) /\left(\mathrm{OV}(\mathrm{i})^{* *} 2\right)$
110 continue
shrink $=$ stemp $/\left(\left(\right.\right.$ sse ${ }^{*}$ tmse $)+$ stemp $)$
do $125 \mathrm{i}=1, \mathrm{IP}$
$\operatorname{ssh}(\mathrm{i})=\operatorname{ssh}(\mathrm{i})+\left(\left(\mathrm{B}(\mathrm{i})^{*} \text { shrink }\right)-\operatorname{tbeta}(\mathrm{i})\right)^{* *} 2$
125
continue

C Compute the OLS estimator - from standardized $\mathbb{X}$ matrix

CALL $\operatorname{SVBKSB}(\mathbf{U}, \mathrm{V}, \mathrm{V}, \mathrm{N}, \mathrm{P}, \operatorname{maxr}, \max \mathrm{c}, \mathrm{Y}, \mathrm{B})$

C Compute the DELTA'S - OLS

CALL VTBETA( $\mathrm{V}, \max \mathrm{C}, \mathrm{B}, \mathrm{P}, \mathrm{DELTA})$

C Compute the $K$ and $k$ values for RRE, GRRE, and the others

```
nomura = DBLE(0.0)
do 200 i=1,p
    KHK(i) = sse/(delta(i)**2)
    F = ((w(i)**2) * (DE1TA(i)**2))/sse
    KCAS(i) = (w(i)**2)/(DBLE(1) + F)
    sum = DSQRT(DBLE(1) + ((V(I)**2)*(DELTA (I)**2))/sse)
    nomura = nomura + (DELTA(I)**2)/(DBLE(1) + sum)
```

```
nomura = DBLE(P) * sse/nomura
```

C Compute the Hoerl, Kennard and Baldwin and the Lavless and Vang constant

```
fachk = DBLE(0.0)
faclv = DBLE(0.0)
do 220 i = 1,P
    fachk = fachk + B(i) ** 2
    faclu = faclw + (DELTA(i)**2) * (V(i)**2)
fachk = (DBLE(P) * sse)/fachk
faclw = (DBLE(P)* sse)/faclw
```

220 continue

C Compute the PC estimators, by setting the smallest singular value equal
C to 0 for PC1
do $250 \mathrm{i}=1,4$
$T W(i)=W(i)$

250 continue

$$
\operatorname{TY}(5)=\operatorname{DBLE}(0.0){ }^{*} \text { set smallest singular value to zero }
$$

C Compute PC1 estimator

$$
\text { CALL } \operatorname{SVBKSB}(\mathrm{D}, \mathrm{TV}, \mathrm{~V}, \mathrm{~N}, \mathrm{P}, \operatorname{maxr}, \operatorname{maxc}, \mathrm{Y}, \mathrm{~B})
$$

C Compute the estimate of sigma via PC1

```
CALL SIGMA(STDX, maxr, maxc, N, P, Y, B, RES, ssepc1)
ssepc1 = DBLE(n-p) * ssepc1 - DBLE(N) * (meany**2) correct for \overline{Y}
ssepc1 = ssepc1/(DBLE(N-P-1))
```

C Compute the ridge constant, $k$, using PC1 estimator

$$
\begin{aligned}
& \text { hkpc1 }=\operatorname{DBLE}(0.0) \\
& \text { do } 260 \mathrm{i}=1, \mathrm{P} \\
& \quad \mathrm{hkpc} 1=\mathrm{hkpc} 1+\mathrm{B}(\mathrm{i})^{* *} 2 \\
& \text { continue } \\
& \text { hkpc1 }=\operatorname{DBLE}(\mathrm{p})^{*} \operatorname{ssepc} 1 / \mathrm{hkpc} 1
\end{aligned}
$$

260

C Unstandardize the PCE's and compute efficiency

CALL BETA(SS, AVE, B, BJ, maxc, P, meany)
do $262 \mathbf{i}=1,(\mathbf{P}+1)$ here $\hat{\beta}_{0}-\beta$ was also computed, but not used $\operatorname{spc} 1(i)=\operatorname{spc} 1(i)+(B U(i)-\operatorname{tbeta}(i))^{* *} 2$
262
continue

C Compute the deltas ( $\mathrm{V}^{\prime} \mathrm{b}$ ) to obtain the constants for fractional
C rank estimators $\hat{\delta}=V^{\prime} \hat{\beta}_{\mathrm{pc}}$

CALL VtBETA( $\mathrm{V}, \operatorname{maxc}, \mathrm{B}, \mathrm{P}, \mathrm{dELTA})$

C Compute vector K , using PC1 estimator

$$
\begin{aligned}
& \text { do } \quad 280 \text { i }=1, \mathrm{P} \\
& \text { if }\left(\left(\operatorname{DELTA}(\mathrm{i})^{* *} 2\right) \cdot g e . \operatorname{sinf}\right) \text { then } \\
& \text { KPC1 }(\mathrm{i})=\operatorname{ssepc} 1 /\left(\operatorname{DELTA}(\mathrm{i})^{* *} 2\right) \\
& \text { else } \\
& \quad \operatorname{KPC} 1(\mathrm{i})=\operatorname{DBLE}(9999.9) \\
& \text { endif }
\end{aligned}
$$

280
continue

C Set the second smallest eigenvalue equal to zero

$$
\mathrm{TV}(4)=\operatorname{DBLE}(0.0)
$$

C Compute the PC2 estimator

$$
\operatorname{CALL} \operatorname{SVBKSB}(\mathrm{U}, \mathrm{TV}, \mathrm{~V}, \mathrm{~N}, \mathrm{P}, \max , \operatorname{maxc}, \mathrm{Y}, \mathrm{~B})
$$

C Compute the estimate of sigma using PC2 estimator

```
CALL SIGMA(STDX, maxr, maxc, N, P, Y, B, RES, ssepc2)
ssepc2 = DBLE(N-P)*ssepc2 - DBLE(N) * (meany**2)
ssepc2 = ssepc2/(DBLE(N-P-1))
```

C Compute k, using PC2

$$
\begin{aligned}
& \mathbf{h k p c} 2=\operatorname{DBLE}(0.0) \\
& \text { do } 290 \mathrm{i}=1, \mathrm{P} \\
& \quad \mathrm{hkpc} 2=\mathrm{hkpc} 2+\mathrm{B}(\mathrm{i})^{* *} 2 \\
& \text { continue } \\
& \text { hkpc2 }=\operatorname{DBLE}(\mathrm{P})^{*} \operatorname{ssepc} 2 / \mathrm{hkpc} 2
\end{aligned}
$$

290 continue

C Compute the V'b for PC2

CALL VTBETA(V, maxc, B, P, DELTA)

C Compute the vector K using PC 2

```
        do 330 i = 1,P
            if ((DELTA(i)**2).ge.sinf) then
            KPC2(i) = ssepc2/(DELTA(i)**2)
            else
            KPC2(i) = DBLE(9999.9)
            endif
```

C Onstandardize the PCE's and compute efficiency

$$
\begin{aligned}
& \text { CALL BETA(SS, AVE, B, BU, maxc, P, meany) } \\
& \text { do } 335 \mathrm{i}=1,(\mathrm{P}+1) \\
& \quad \operatorname{spc} 2(\mathrm{i})=\operatorname{spc} 2(\mathrm{i})+(\mathrm{BU}(\mathrm{i})-\operatorname{tbeta}(\mathrm{i}))^{* *} 2
\end{aligned}
$$

335 continue

C Compute the ridge estimators - using $k$ of Hoerl, Kennard and Baldwin

$$
\text { CALL RIDGE( } \mathrm{V}, \mathrm{~V}, \mathrm{~V}, \mathrm{~N}, \mathrm{P}, \max , \operatorname{maxc}, \mathrm{Y}, \mathrm{~B}, \mathrm{fach} \mathrm{k})
$$

C Unstandardize the RHKB estimators and compute efficiency

$$
\begin{aligned}
& \text { CALL BETA(SS, AVE,B,BD, MAXC, P, meany }) \\
& \text { do } 345 i=1,(\mathrm{P}+1) \\
& \quad \operatorname{srhk}(\mathrm{i})=\operatorname{srhk}(\mathrm{i})+(\mathrm{BD}(\mathrm{i})-\operatorname{tbeta}(\mathrm{i}))^{* * 2}
\end{aligned}
$$

345 continue

C Compute the ridge estimators - using $k$ of Lawless and Vang

CALL RIDGE( $\mathrm{D}, \mathrm{V}, \mathrm{V}, \mathrm{N}, \mathrm{P}, \operatorname{maxr}, \operatorname{maxc}, \mathrm{Y}, \mathrm{B}, \mathrm{faclw})$

C Unstandardize the ridge estimators and compute efficiency

CALL BETA (SS, AYE, B, BO, maxc, $P$, meany) do $355 \mathrm{i}=1,(\mathrm{P}+1)$ $\operatorname{srlw}(i)=\operatorname{srlw}(i)+(B D(i)-\operatorname{tbeta}(i))^{* *} 2$
355
continue

C Compute the GRRE via Hoerl, Kennard and Baldwin

CALL $\operatorname{GRRE}(\mathbf{J}, \mathrm{V}, \mathrm{V}, \mathrm{N}, \mathrm{P}, \operatorname{maxr}, \operatorname{maxc}, \mathrm{Y}, \mathrm{B}, \mathrm{Kh} \mathrm{k})$

C Unstandardize GRRE and compute efficiency

$$
\begin{aligned}
& \text { CALL BETA (SS, AVE, B, BU,maxc, P, meany }) \\
& \text { do } 365 \mathrm{i}=1,(\mathrm{P}+1) \\
& \quad \operatorname{sgrhk}(\mathrm{i})=\operatorname{sgrhk}(\mathrm{i})+(\mathrm{BU}(\mathrm{i})-\operatorname{tbeta}(\mathrm{i}))^{* *} 2 \\
& \text { continue }
\end{aligned}
$$

365

C Compute the GRRE via Troskie

CALL GRRE( $\mathbf{~}, \mathrm{V}, \mathrm{Y}, \mathrm{N}, \mathrm{P}, \max , \operatorname{maxc}, \mathrm{Y}, \mathrm{B}, \mathrm{Kcas})$

C Unstandardize GRRE and compute efficiency

CALL BETA (SS, AVE, B, BJ, maxc, P, meany) do $375 \mathrm{i}=1,(\mathrm{P}+1)$
$\operatorname{sgrcas}(i)=\operatorname{sgrcas}(i)+(B O(i)-\operatorname{tbeta}(i))^{* *} 2$
375 continue

C Compute the AUORR estimator, via nomura

CALL $\operatorname{AOORR}(\mathrm{O}, \mathrm{V}, \mathrm{V}, \mathrm{N}, \mathrm{P}, \operatorname{maxr}, \operatorname{maxc}, \mathrm{Y}, \mathrm{B}$, nomura)

C Unstandardize AUORR and compute efficiency

CALL BETA(SS, AVE, B, BJ, maxc, P, meany) do $385 \mathrm{i}=1,(\mathrm{P}+1)$
sauorr $(\mathrm{i})=\operatorname{sauorr}(\mathrm{i})+(\mathrm{BU}(\mathrm{i})-\operatorname{tbeta}(\mathrm{i}))^{* *} 2$
385 continue

C Compute the AUGRR estimator where estimate of $K$ is via Hoerl, Kennard and C Baldwin

CALL $\operatorname{AUGRR}(\mathbb{O}, \mathrm{V}, \mathrm{V}, \mathrm{N}, \mathrm{P}, \operatorname{maxr}, \max \mathrm{C}, \mathrm{Y}, \mathrm{B}, \mathrm{Kh} \mathbf{)}$

C Unstandardize AUGRR and compute efficiency

```
CALL BETA(SS, AVE, B, BU, maxc, P, meany)
do 395 i = 1,(P+1)
    saugrr(i) = saugrr(i) + (BO(i) - tbeta(i))**2
continue
```

C Compute FPC estimator, using GRR, where $\mathbb{K}$ is estimated by PC1


C Unstandardize the FPCI estimator, and compute efficiency

CALL BETA(SS, AVE, B, BJ, maxc, P, meany) do $400 \mathrm{i}=1,(\mathrm{P}+1)$
$\operatorname{sfgrp} 1(i)=\operatorname{sfgrp} 1(i)+(B O(i)-\operatorname{tbeta}(i))^{* *} 2$
400 continue

C Compute the FPC estimator, using GRR, where $K$ is estimated by PC2

CALL $\operatorname{GRRE}(\mathbf{O}, \mathrm{V}, \mathrm{V}, \mathrm{N}, \mathrm{P}, \operatorname{maxr}, \operatorname{maxc}, \mathrm{Y}, \mathrm{B}, \mathrm{Kpc} 2)$

C Unstandardize the FPCI estimator, and compute efficiency

CALL BETA(SS, AYE, B, BJ, maxc, P, meany) do $410 \mathbf{i}=1,(\mathrm{P}+1)$
$\operatorname{sfgrp2(i)}=\operatorname{sfgrp2(i)}+(B O(i)-\operatorname{tbeta}(i))^{* *} 2$
410 continue

C Compute the FPC estimator, using RR, where $k$ is estimated using PC1

CALL RIDGE( $\mathrm{D}, \mathrm{H}, \mathrm{V}, \mathrm{N}, \mathrm{P}, \max \mathrm{r}, \max \mathrm{C}, \mathrm{Y}, \mathrm{B}, \mathrm{hkpc} 1)$

C Unstandardize the FPCV estimator, and compute efficiency

```
CALL BETA(SS, AVE, B, BU, maxc, P, meany)
```

do $420 \mathrm{i}=1,(\mathrm{P}+1)$
$\operatorname{sfrpc} 1(i)=\operatorname{sfrpc} 1(i)+(B J(i)-\operatorname{tbeta}(i))^{* *} 2$

420 continue

C Compute the PPC estimator, using RR, where $k$ is estimated using PC2

> CALL RIDGE(J, V, V, N, P, maxr, maxc, Y, B, hkpc2)

C Unstandardize the FPCV estimator, and compute efficiency

```
CALL BETA(SS, AVE, B, BU, maxc, P, meany)
do 430 i = 1,(P+1)
    sfrpc2(i) = sfrpc2(i) + (BU(i) - tbeta(i))**2
                continue
```

500 continue

C Set sum of traces (over 5 regression coefficients) for the relevant
C estimators equal to zero

$$
\begin{aligned}
& \text { tols = DBLE (0.0) } \\
& \text { trhk }=\operatorname{DBLE}(0.0) \\
& \text { trlw }=\operatorname{DBLE}(0.0) \\
& \text { tgrhk }=\text { DBLE (0.0) } \\
& \text { tgrcas }=\operatorname{DBLE}(0.0) \\
& \text { taugrr }=\text { DBLE (0.0) } \\
& \text { tauorr }=\operatorname{DBLE}(0.0) \\
& \mathrm{tpc} 1=\operatorname{DBLE}(0.0) \\
& \mathrm{tpc} 2=\mathrm{DBLE}(0.0) \\
& \text { tfgrp1 }=\operatorname{DBLE}(0.0) \\
& \mathrm{tfgrp} 2=\operatorname{DBLE}(0.0) \\
& \operatorname{tfrpc} 1=\operatorname{DBLE}(0.0) \\
& \text { tfrpc2 }=\text { DBLE (0.0) } \\
& \text { tssh }=\operatorname{DBLE}(0.0)
\end{aligned}
$$

C Compute the efficiency totals for the different biased estimators

$$
\begin{aligned}
& \text { do } 550 \mathrm{j}=2,6 \\
& \text { tols }=\text { tols }+\operatorname{sols}(\mathrm{j}) \\
& \text { trhk }=\operatorname{trhk}+\operatorname{srhk}(\mathrm{j}) \\
& \text { trlw }=\mathbf{t r l w}+\operatorname{srlw}(\mathrm{j}) \\
& \text { tgrhk }=\text { tgrhk }+\operatorname{sgrhk}(j) \\
& \text { tgrcas }=\text { tgrcas }+\operatorname{sgrcas}(\mathrm{j}) \\
& \text { taugrr }=\text { taugrr }+\operatorname{saugrr}(\mathrm{j}) \\
& \text { tauorr }=\text { tauorr }+\operatorname{sauorr}(\mathrm{j}) \\
& \operatorname{tpc} 1=\operatorname{tpc} 1+\operatorname{spc} 1(\mathrm{j}) \\
& \text { tpc2 }=\mathrm{tpc} 2+\mathrm{spc} 2(\mathrm{j}) \\
& \operatorname{tfgrp} 1=\operatorname{tfgrp} 1+\operatorname{sfgrp} 1(j) \\
& \text { tfgrp2 }=\text { tfgrp2 }+\operatorname{sfgrp2(j)~} \\
& \operatorname{tfrpc} 1=\operatorname{tfrpc} 1+\operatorname{sfrpc}(j) \\
& \text { tfrpc2 }=\text { tfrpc2 }+\operatorname{sfrpc} 2(j) \\
& \text { tssh }=\mathbf{t s s h}+\mathbf{s s h}(\mathrm{j})
\end{aligned}
$$

550 continue

C Vrite the efficiency and the relative efficiency ratio to output file

```
        write(16,555) tols, tols/tols
555 format (' OLS', T10, 2F15.8)
        write(16,565) trhk, tols/trhk
565 format (' RHK', T10, 2F15.8)
write(16,575) trlw, tols/trlw
575 format (' RLV', T10, 2F15.8)
        write(16,585) tgrhk, tols/tgrhk
585 format ('GRHK', T10, 2P15.8)
        write(16,595) tgrcas, tols/tgrcas
595 format ('GRCAS', T10, 2F15.8)
        write(16,605) taugrr, tols/taugrr
        format ('AOGRR', T10, 2F15.8)
        write(16,615) tauorr, tols/tauorr
615 format ('AOORR', T10, 2F15.8)
        write(16,625) tpc1, tols/tpc1
```

```
625 format ('PC1',T10, 2P15.8)
    write(16,635) tpc2, tols/tpc2
635 format ('PC2', T10, 2F15.8)
    write(16,645) tfgrp1, tols/tfgrp1
645 format ('FGRPC1', T10, 2P15.8)
    write(16,655) tfgrp2, tols/tfgrp2
655 format ('FGRPC2', T10, 2F15.8)
    write(16,665) tfrpc1, tols/tfrpc1
665 format ('FRPC1', T10, 2P15.8)
    write(16,675) tfrpc2, tols/tfrpc2
675 format ('PRPC2', T10, 2P15.8)
    write(16,685) tssh, tols/tssh
685 format ('SHE', T10, 2P15.8)
    close (16)
    CLOSE (15)
    END * End of main program *
```

```
SUBROUTINE AUGRR (U, V, V, rows, col, maxr, maxc, Y, BETA, K)
```

C This subroutine computes the AUGRR estimator and is the same as SVBKSB
C except for a change in the root part. The comments under SYBKSB are also
$C$ relevant here. The vector $I$ is input value.

```
parameter (nmax = 100)
double precision O(maxr,maxc), V(maxc), V(maxc,maxc)
double precision Y(maxr), BETA(maxc)
double precision TMP(nmax), S
double precision K(maxc)
integer rows, col
do 12 j = 1,col
    S = DBLE(0.0)
    if (V(j).NE.DBLE(0.)) then
        do 11 i = 1,rows
                S = S + O(i,j) * Y(i)
                continue * note change to SVBKSB is done here*
                S = S * (V(j)* (DBLE (1.0)+K}(\mathbf{j})/(\mathbb{V}(\textrm{j})**2+K(j)))/(V(j)** 2+K(j)))
    endif
    TMP(j) = S
continue
do 14 J = 1,col
        S = DBLE(0.0)
        do 13 jj = 1,col
            S = S + V(j,jj) * TMP(jj)
        continue
BETA(j) = S
continue
return
end
```

```
SUBROUTINE AOORR (U, \(V\), \(V\), rows, col, maxr, maxc, Y, BETA, rc)
```

C This subroutine computes the AUORR estimator and is the same as SVBKSB
C except for a change in the root parts (in SVBKSB we change the $S$ part).
C The comments under SVBKSB are also relevant here. The ridge constant
C is denoted by rc.

```
parameter (nmax = 100)
double precision \(O(\max , \max ), V(\max c), V(\max c, \max )\)
double precision \(Y(\operatorname{maxr})\), BETA(maxc)
double precision TMP(nmax), S
double precision rc
integer rows, col
do \(12 \mathrm{j}=1, \mathrm{col}\)
    S = DBLE(0.0)
    if ( \(\mathbf{V}(\mathrm{j}) \cdot \mathrm{NE} . \mathrm{dble}(0.0)\) ) then
        do \(11 \mathrm{i}=1\), rows
            \(S=S+U(i, j) * Y(i)\)
        continue
                                    *changes to SVBKSB are done here*
            \(\mathbf{S}=S^{*}\left(\operatorname{DBLE}(1.0)-r c^{* *} 2 /\left(\left(\mathbf{V}(\mathbf{j})^{* *} 2+r c\right)^{* *} 2\right)\right) / \mathbf{V}(j)\)
    endif
\(\operatorname{TMP}(\mathrm{j})=\mathrm{S}\)
continue
do \(14 \mathrm{j}=1\), col
    \(\mathrm{S}=\mathrm{DBLE}(0.0)\)
    do \(13 \mathrm{jj}=1, \mathrm{col}\)
        \(S=S+V(j, j j) * T M P(j j)\)
    continue
    BETA \((\mathrm{j})=\mathrm{S}\)
continue
return
end
```

SUBRODTINE BETA (S, AVE, B, BU, maxc, row, meany)

C Aim: Unstandardize the betas
C Input to this subroutine:
C $\quad$ - vector containing the sum of squares of columns of $\mathbb{X}$
C AVE - vector containing the means of columns of $\mathbb{X}$,
C meany - mean of $Y$
C B - vector containing the standardized betas
C maxc - physical dimensions of S, AVE, B, BU
C row - logical dimension of S, AYE, B, BJ
C Return to main program BO - vector, of unstandardized betas

```
double precision S(maxc), AVE(maxc), B(maxc), BD(maxc)
double precision meany
integer row
BU(1) = DBLE(0.0)
do 100 i = 1,row
    BU(i+1) = B(i) * DBLE(1)/(DSqRT(S(I)))
    BO}(1)=\operatorname{BO}(1)+\operatorname{BU}(i+1) * AVE(I
100 continue
BO(1) = meany - BU(1)
return
end
```

SUBROUTINE COPY (A, B, maxr, maxc, ra, cA)

C Copies $\mathbf{A}(\mathrm{r} \Lambda, \mathrm{c} \Lambda$ ) into $B$ and loses previous B
double precision $A($ maxr, , maxc $)$
double precision $B(\operatorname{maxr}, \max c)$
integer ra, cA
do $100 \mathrm{i}=1, \mathrm{rA}$
do $20 \mathrm{j}=1, \mathrm{~cd}$
$B(i, j)=A(i, j)$
20
continue
100 continue
return
end

SUBROUTINE GRRE( $, V, V$, rows, col, maxr, $\operatorname{maxc}, Y$, BETA, $K)$

C This subroutine computes the generalized ridge regression estimator and
$C$ is the same as SVBKSB except for adding $\mathbf{k}_{\mathbf{i}}$ to the diagonal elements of
C $V$ when we divide. The comments under SVBKSB are also relevant here. K
$C$ is the ridge vector. If ( $k(i) . g e .1 \times 10^{* *} 10$ ), then DELTA ( $V^{\prime} B$ ) is set to 0

```
parameter ( \(n \max =100\) )
```



```
double precision \(Y(\operatorname{maxr})\), BETA (maxc)
double precision TMP(nmax), S, K(maxc)
integer rous, col
do \(12 \mathrm{j}=1\), col
    \(\mathrm{S}=\mathrm{DBLE}(0.0)\)
    If ( \(\mathbf{V}(\mathrm{j}) . \operatorname{NE} . \operatorname{DBLE}(0.0))\) then
        do \(11 \mathrm{i}=1\), rows
            \(S=S+W(i, j) * Y(i)\)
        continue
        \(\mathbf{S}=\mathbf{S}^{*}\left(\mathbf{V}(\mathbf{j}) /\left(\mathbf{V}(\mathbf{j})^{* *} \mathbf{2}+\mathbf{K}(\mathbf{j})\right)\right) \quad{ }^{*}\) changes made here*
        endif
        TIP \((\mathrm{j})=\mathrm{S}\)
```

    continue
    do \(15 \mathrm{~J}=1, \mathrm{col}\)
        \(\mathrm{S}=\mathrm{DBLE}(0.0)\)
        do \(14 \mathrm{jj}=1, \mathrm{col}\)
        \(\mathbf{S}=\mathbf{S}+\mathrm{V}(\mathrm{j}, \mathrm{j} \mathrm{j}){ }^{*} \operatorname{TMP}(\mathrm{j} j)\)
    continue
    \(\operatorname{BETA}(\mathrm{j})=\mathrm{S}\)
    continue
    return
    end
    
## SUBROUTINE RIDGE(J, $V$, $V$, ROVS, COL, MAXR, MAXC, Y, BETA, rc)

C This subroutine computes the ridge estimator and is the same as SVBKSB
$C$ except for adding $k$ to the diagonal elements of $V$ when we divide.
$C$ The comments under SVBKSB are also relevant here. The ridge constant
C is denoted by rc.

```
parameter (nmax = 100)
double precision \(V(\max r, \max ), \mathrm{V}(\operatorname{maxc}), \mathrm{V}(\max c, \max c)\)
double precision \(Y(\operatorname{maxr})\), BETA (maxc)
double precision TMP(nmax), S, rc
integer rows, col
do \(12 \mathrm{j}=1, \mathrm{col}\)
    S = DBLE (0.0)
    if ( \(\mathbf{V}(\mathrm{j})\).NE.DBLE(0.)) then
        do \(11 \mathrm{i}=1\), rows
                \(S=S+U(i, j) * Y(i)\)
        continue
        \(\mathbf{S}=\mathbf{S}^{*}\left(\mathbf{V}(\mathbf{j}) /\left(\mathbf{V}(\mathrm{j})^{* *} \mathbf{2}+\mathbf{r c}\right)\right) \quad{ }^{*}\) changes made here*
    endif
    \(\operatorname{TMP}(\mathrm{j})=\mathrm{S}\)
    continue
    do \(14 \mathrm{j}=1\), col
        \(\mathrm{S}=\mathrm{DBLE}(0.0)\)
        do \(13 \mathrm{jj}=1, \mathrm{col}\)
            \(\mathrm{S}=\mathrm{S}+\mathrm{V}(\mathrm{j}, \mathrm{jj}){ }^{*} \mathrm{TMP}(\mathrm{jj})\)
        continue
        \(\operatorname{BETA}(\mathrm{j})=\mathrm{S}\)
14 continue
return
end
```


## SUBROUTINE SIGMA (MAT, maxr, maxc, N, P, Y, B, RES, sse)

C This subroutine computes an estimate of sigma square, sse.
$C$ Mat is the $X$ matrix, $Y$ is the $Y$ vector, and $B$ is the vector of
$C$ regression coefficients. The physical dimensions are maxr and maxc and
$C$ the logical dimensions are $N$ and $P$.

```
double precision MAT(maxr,maxc), Y(maxr)
double precision RES(maxr), sse
double precision B(maxc), temp
integer N,P,i,j
sse = DBLE(0.0)
do 100 i = 1,N
    temp = DBLE(0.0)
    do 15 j = 1,P
        temp = temp + MAT(i,j) * B(j)
    continue
    RES(I) = Y(I) - temp
    sse = sse + (RES(I)**2)
sse = sse/DBLE(N-P)
return
end
```

15
100 continue

SUBROUTINE STAND (MAT, maxr, maxc, N, P, STDX, AVE, SS)

C This subroutine standardizes MAT (physical dimensions maxr, maxc:
C logical dimensions $N$ and $P$ ) by subtracting the mean from each column and
C then divides the column by its standard deviation.
C Note the matrix Mat should not contain a column of ones

```
double precision MAT(maxr, maxc), AVE(maxc)
double precision SS(maxc), STDX(maxr, maxc)
double precision temp
integer \(\mathrm{N}, \mathrm{P}, \mathrm{j}\), i
do \(6000 \mathrm{j}=1, \mathrm{P}\)
    \(\operatorname{AVE}(\mathrm{j})=\operatorname{DBLE}(0.0)\)
    do \(605 \mathrm{i}=1, \mathrm{~N}\)
                \(\operatorname{AVE}(\mathrm{j})=\operatorname{AVE}(\mathrm{j})+\operatorname{MT}(\mathrm{i}, \mathrm{j})\)
605 continue
    \(\operatorname{AVE}(\mathrm{j})=\operatorname{AVE}(\mathrm{j}) / \operatorname{DBLE}(\mathrm{N})\)
    do \(6010 \mathrm{j}=1, \mathrm{P}\)
        temp \(=\operatorname{DBLE}(0.0)\)
                do \(610 \mathrm{i}=1, \mathrm{~N}\)
                temp \(=\) temp \(+(\mathbf{M A T}(\mathbf{i}, \mathrm{j})-\operatorname{AVE}(\mathrm{j}))^{* *} 2\)
```

6000 continue
610
continue
SS(j) = temp
6010 continue
do $7000 \mathrm{i}=1, \mathrm{~N}$
do $6900 \mathrm{j}=1, \mathrm{P}$
$\operatorname{STDX}(\mathrm{i}, \mathrm{j})=(\operatorname{MAT}(\mathrm{i}, \mathrm{j})-\operatorname{AVE}(\mathrm{j})) / \operatorname{DSqRT}(\operatorname{SS}(\mathrm{j}))$
6900 continue
7000 continue
return
end

SUBROUTINE SVBKSB (U, V, V, M, N, MP, NP, B, X)

C The computation of the estimators is based on this subroutine. The
C program-code is from Press et al. (1985). The only change made to
C SVBKSB was a switch to double precision.
C
C AIM of SVBKSB: Solves $A X=B$ for a vector $\mathbb{X}$, where $\mathbb{A}$ is specified by
$C$ the arrays $U, V, V$ as returned by SVDCMP. $I$ and $N$ are the logical
C dimensions of $A$. $M P$ and NP are the physical dimensions of $A$. $B$ is the
C input right-hand side. $X$ is the output solution vector. No input
C quantities are destroyed, so the routine may be called sequentially with
C different $B$ 's

```
parameter ( \(\mathrm{nmac}=100\) ) Kaximum anticipated value of N
double precision \(\mathbf{~ ( M P , N P ) , ~ V ( N P ) , ~ V ( N P , N P ) ~}\)
double precision \(B(I P), X(N P)\), TMP(nmax), \(S\)
do \(12 \mathrm{j}=1, \mathrm{~N} \quad\) Calculate \(\mathrm{U}^{\prime} \mathrm{B}\)
    S = DBLE(0.)
    if \((\mathbf{V}(\mathbf{j})\).NE. \(\operatorname{DBLE}(0)\).\() then Nonzero result only if \mathbf{V}_{\mathrm{j}} \neq 0\)
        do \(11 \mathbf{i}=1\), \(\mathbf{I}\)
            \(S=S+V(i, j) * B(i)\)
```

        continue
        \(S=S / \mathbf{V}(\mathbf{j}) \quad \begin{aligned} & \text { division } \\ & \text { methods }\end{aligned} W_{j}\), changes made here for other
    endif
    \(\operatorname{TIP}(\mathbf{j})=S\)
    continue
do $14 \mathrm{j}=1, \mathrm{~N}$ Matrix multiply by V
$\mathrm{S}=\mathrm{DBLE}(0$.
do $13 \mathrm{jj}=1, \mathrm{~N}$
$S=S+V(j, j j) * T M P(j j)$
continue
$\mathbf{X}(\mathrm{j})=\mathbf{S}$
continue
return
end

SUBROUTINE SVDCMP (A, M, N, MP, NP, Y, Y)
$C$ This routine computes the SVD of the matrix $A$. The program-code is form
C Press et al. (1985) and is not given here. The only change made to
C SVDCIP was a switch to double precision. Remember to make a copy
C of A before calling SVDCIP as $U$ replaces $\mathbb{A}$ on output. Furthermore
$C$ the matrix $V$ is output as $V$ and not the transpose of $V$.

## SUBROUTINE SVDSORT(U, V, V, N, M, NP, MP, C, BB)

C This subroutine sorts the SVD from high to low and is borrowed from
C Troskie (1990).

```
double precision \(\mathbf{U}(\mathbf{N P}, \mathbf{I P}), \mathbf{V}(\mathbf{M P}), \mathrm{V}(\mathbf{M P}, \mathbf{I P})\)
double precision C(NP), BB(IP)
do \(90 \mathrm{k}=1, \mathbf{I}-1\)
    do \(100 \mathrm{j}=1, \mathbf{I}-\mathrm{K}\)
        if ( \(\mathbf{V}(\mathrm{j}) . \operatorname{LT} . W(\mathrm{j}+1))\) then
            HOLD \(=W(\mathrm{j})\)
            do \(110 \mathrm{l}=1, \mathrm{~N}\)
                \(\mathbf{C}(1)=\mathbf{U}(1, j)\)
            continue
            do \(1201=1, I\)
                \(\mathrm{BB}(1)=\mathrm{V}(1, \mathrm{j})\)
            continue
            \(\mathbf{V}(\mathrm{j})=\mathbf{V}(\mathrm{j}+1)\)
            do \(1301=1, N\)
                \(\mathrm{J}(1, \mathrm{j})=\mathbf{D}(1, \mathrm{j}+1)\)
            continue
            do \(1401=1\),I
                \(\mathbf{V}(1, j)=\mathbf{V}(1, j+1)\)
            continue
            \(W(j+1)=\) HOLD
            do \(1501=1, N\)
                \(\mathbf{O}(1, j+1)=C(1)\)
            continue
            do \(1601=1\),
                \(\mathrm{V}(1, \mathrm{j}+1)=\mathrm{BB}(1)\)
            continue
            endif
                            continue
return
end
```

90 continue

## SUbroutine vtbeta (V, maxc, B, P, delta)

C This subroutine computes the delta's, V'B, Here V, B are input matrices $C$ and the DELTAS are returned. Maxc is the physical dimensions, and $P$ the C logical dimensions of the matrices.

```
double precision \(V(\max c, \max c), B(\max c)\), DELTA(maxc)
double precision sum
integer maxc, \(P\)
do \(180 \mathrm{i}=1, \mathrm{P}\)
    sum \(=\operatorname{DBLE}(0.0)\)
    do \(175 \mathrm{j}=1, \mathrm{P}\)
        sum \(=\operatorname{sum}+V(j, i) * B(j)\)
    continue
    DELTA(i) = sum
180 continue
return
end
```


## Appendix C

## NOTATION

This dissertation is written in such a way that as the material develops, the relevant notation is described. Description in context is intended to obviate any confusion. However this summary of notation is also supplied for those readers interested only in particular chapters. The first part of the notation under $A$, represents some matrix operations on the matrix $A$.

A:nxp - matrix of order nxp
a - scalar or column- vector
$A^{\prime} \quad$ - transpose of $A$
$a_{i}$ - $\quad$ i-th element of $a$ when $a$ is $a$ vector
$A_{i}$ - i-th column of $A$ when $A$ is a matrix,
$a_{i j}$ - $i$-th row and $j$-th column element of the matrix $A$
$[A]_{i j}$ - i-th row and $j$-th column element of the matrix $A$
$A^{-}$- (Moore-Penrose) generalized inverse of A
$\mathrm{A}^{-1} \quad$ inverse of A
$|a|$ - absolute value of $a$
$|A|$ - determinant of $A$
$\lambda_{i}\left(A^{\prime} A\right)-\quad i-t h e i g e n v a l u e$ of $A^{\prime} A$
$\|a\|$ - vector norm of $a$
$\|a\|_{2}$ - Euclidean norm (or vector 2-norm) of a vector a
$\|A\|$ - matrix norm of $A$
$\|\mathrm{A}\|_{\mathrm{F}}$ - Frobenius norm of a matrix A
$\|\mathrm{A}\|_{\mathrm{p}} \quad-\quad$ matrix p - norm of A
$N(A)-\quad$ null space of $A$
$\operatorname{tr}[\mathrm{A}]$ - trace of A
$r(A)-\operatorname{rank}$ of $A$
$R(A)$ - range of $A$
$a^{\prime}=\left[a_{1}, \ldots, a_{p+1}\right]$, where $a_{i}$ is defined in Chapter 6
$a(0)=\hat{\delta}_{P C}$
$\mathrm{a}(\mathrm{t})$ - estimate of $\delta$ at the t - th iteration
$\operatorname{acc}_{j}-$ numerical accuracy of the $j$-th regression coefficient

AP - Andrews-Pregibon
$A P_{i j}^{k}-A n d r e w s-P r e g i b o n$ statistic for $k$ suspected outliers, cases $i, j, \ldots$
$A P_{A}$ - Andrews-Pregibon statistic for variables A
ANACOVA - analysis of covariance

ANOVA - analysis of variance

AUGRRE - almost unbiased generalized ridge regression estimator

```
AUORRE - almost unbiased operational ridge regression estimator
\beta:px1 - vector of regression coefficients that must be estimated
\beta=[[\begin{array}{ll}{\mp@subsup{\beta}{\textrm{a}}{\prime}}&{\mp@subsup{\beta}{\textrm{b}}{\prime}}\end{array}\mp@subsup{]}{}{\prime}
\betaa}:\mp@subsup{p}{a}{}\times
\beta
\mp@subsup{\beta}{b}{*}}:\mp@subsup{p}{b}{}\times1 - any pob vector
\beta
\beta
        specifying the subset
\betao - intercept term
\tilde{\beta}}\mathrm{ - any estimator of }\beta\mathrm{ , usually biased
\tilde{\beta}}\mathrm{ - LSE of }\beta\mathrm{ in the augmented model
\beta}- 0LSE of 
\delta\hat{\beta} - perturbation vector of \hat{\beta}
\hat{\beta}}\mp@subsup{\textrm{a}}{}{-
\mp@subsup{\beta}{b}{}
\mp@subsup{\hat{\beta}}{c}{}-constrained LSE of }
```

$\beta_{\text {basic }}$ - unique solution of smallest norm
$\hat{\beta}_{\text {LS }}$ - corrected least square estimator
$\hat{\beta}_{-\mathrm{i}}-$ OLSE of $\beta$ with the i -th case ( i - th row) deleted
$\hat{\beta}_{-1}-0 L S E$ of $\beta$ with the rows indexed by I deleted, I $\subset\{1,2, \ldots, \mathrm{n}\}$
$\beta_{(i)} \quad$ - subset of $\beta$ that corresponds to the $X_{(i)}$ part of $X$
$\hat{\beta}_{\text {FPC }}$ - fractional principal component estimator of $\beta$
$\hat{\beta}_{G}$ - posterior mean of $\beta$
$\hat{\beta}_{J}$ - jackknife estimator of $\beta$
$\hat{\beta}_{J W}$ - weighted jackknife estimator of $\beta$
$\hat{\beta}_{\mathrm{K}}$ - generalized ridge regression estimator of $\beta$
$\hat{\beta}_{\mathrm{LR}}$ - latent root estimator of $\beta$
$\hat{\beta}_{\mathrm{PC}}$ - principal component estimator of $\beta$
$\hat{\beta}_{\mathrm{R}}$ - ridge regression estimator of $\beta$
$\hat{\beta}_{\mathrm{SH}}$ - shrunken estimator of $\beta$
$\hat{\beta}_{\text {SSH }}$ - Sclove's modified shrunken estimator of $\beta$
$\hat{\beta}_{\text {ILS }}$ - total least square estimator of $\beta$
$\hat{\beta}_{\hat{\mathrm{r}}}$ - estimate of $\beta$, obtained when $X_{\hat{r}}$ is used instead of $X$
$\hat{\beta}_{\mathrm{s}}$ - biased estimator of $\beta$
$\left(\hat{\beta}_{\mathrm{s}}\right)_{-\mathrm{h}}$ - biased estimator of the truncated model
( $\beta \mid \mathrm{Y}$ ) - $\quad \beta$ given Y , Bayesian posterior random variable
b - bias vector
$\mathrm{b}_{1}, \mathrm{~b}_{2}$ - two competing estimators
b - estimator of $\beta$ in Chapter 8
$b_{j i}$ - $j$-th element of $b$, in the $i$-th Monte Carlo repetition
$b_{j, L S}-j$-th element of the OLSE
$\tilde{b}$ - OLSE for $\beta$, where the $X$ matrix is perturbed ( $\tilde{X}=X+\delta X)$

ह - value of a coefficient as computed free from round -off errors
$b_{i}$ - OLSE when $X_{i}$ is regressed on the remaining columns of $X$
$B_{1}=X_{1}\left(X_{1}^{\prime} X_{1}\right)^{-1 / 2}$
BLUE - best linear unbaised estimator
$c_{i}=u_{i}{ }^{\prime} Y \lambda_{i}$ element in describing the estimators of $\beta$
$c^{\prime}=\left[c_{1}, c_{2}, \ldots c_{p}\right]$, vector of constants
c - number of collinearities
$c_{i}=\sum_{j=1}^{p}\left(z_{i j} / \lambda_{j}\right)^{2}$
C - Cook's squared distance
$C_{i}$ - Cook's squared distance after i-th case deleted
$C_{i j} \ldots$ - Cook's squared distance when cases $i, j, \ldots$ are deleted
$C_{A}$ - Cook's squared distance when variables A are deleted
$C_{p}-$ Mallows ${ }^{\prime} C_{p}$

CLS - corrected least squares
$\operatorname{Cov}(\bullet, \bullet)-c o v a r i a n c e ~ b e t w e e n ~ t w o ~ r a n d o m ~ v a r i a b l e s ~$
d - deterministically (fixed) or stochasically defined constant
$\Delta: n x p$ - diagonal matrix of the singular values of $X$
$\Delta=\left[\begin{array}{ll}\Delta_{1} & \Delta_{2}\end{array}\right], \Delta_{1}:(\mathrm{p}-\mathrm{r}) \mathrm{x}(\mathrm{p}-\mathrm{r})$
$\Delta_{\mathrm{e}}: \mathrm{nx}(\mathrm{p}+1)$ - error matrix, in the errors-in-variables model
$\left[\Delta_{\mathrm{e}}\right]_{\mathrm{i}}$ - i-th row of the error matrix
$\Delta_{\hat{r}}: n x \hat{r}$ - diagonal matrix of the singular values of $X_{\hat{r}}$
df - degrees of freedom
$\operatorname{Diag}(\bullet)$ - diagonal matrix of elements of a vector argument
$\operatorname{diag}(\bullet)$ - vector of diagonal elements of a matrix argument
$D^{2}$ - Cook's squared distance
$D_{a}=\operatorname{Diag}\left(\sqrt{\lambda_{1}}, \sqrt{\lambda_{2}} \ldots \sqrt{\lambda_{k}}\right)$ - the non-zero singular values of $X$
$\hat{D}_{a}$ - diagonal matrix consisting of the singular values of $\hat{X}$
$\tilde{\Delta}$ - diagonal matrix of the singular values of the augmented matrix [X Y]
$\delta=V^{\prime} \beta$
$\delta_{i}-$ i-th element of $\delta$
$\delta_{\text {max }}-$ largest element of $\delta$
$\delta=\left[\begin{array}{ll}\delta_{1}^{\prime} & \delta_{2}^{\prime}\end{array}\right]^{\prime}, \delta_{1}: \mathrm{rx} 1$ and $\delta_{2}:(\mathrm{p}-\mathrm{r}) \mathrm{x} 1$
$\hat{\delta}$ - OLSE of $\delta$
$\hat{\delta}_{\mathrm{J}}-\quad$ jackknife ridge estimator of $\delta$
$\hat{\delta}_{J W}$ - weighted jackknife ridge estimator of $\delta$
$\left[\hat{\delta}_{J W}\right]_{i}-\quad$ i-th element of $\hat{\delta}_{J W}$
$\hat{\delta}_{\mathrm{K}}-\operatorname{GRRE} \delta, \mathrm{K}$ is non-stochastic
$\hat{\delta}_{\hat{\mathrm{K}}}$ - GRRE of $\delta, \mathrm{K}$ is stochastic
$\left[\hat{\delta}_{\hat{K}}\right]_{\mathrm{i}}$ or $\left(\hat{\delta}_{\hat{K}}\right)_{\mathrm{i}}$ i-th element of $\hat{\delta}_{\hat{K}}$
$\left[\hat{\delta}_{\hat{0}}\right]_{\mathrm{i}}$ - i-th element of operational AUGRRE
$\left[\hat{\delta}_{K}(\mathrm{t})\right]_{\mathrm{i}}-\operatorname{GRRE}$ of $\delta_{\mathrm{i}}$ at t - th iteration
$\left[\hat{\delta}_{K}(0)\right]_{i}=\hat{\delta}_{i}$
$\hat{\delta}_{\mathrm{R}}$ - ridge estimator of $\delta$
$\hat{\delta}_{\mathrm{R}}(\mathrm{t})$ - ridge estimate of $\delta$ at t - th iteration
$\hat{\delta}_{\mathrm{R}}(0)=\hat{\delta}$
$\hat{\delta}_{\text {FPC }}$ - fractional principal component estimator of $\delta$
$\hat{\delta}_{\mathrm{FPCI}}=\mathrm{F}_{\mathrm{PCI}} \hat{\delta}$
$\hat{\delta}_{\mathrm{FPCV}}=\mathrm{F}_{\mathrm{PCV}} \hat{\delta}$
$\partial_{\mathrm{i}}$ - shrinkage factor
$\delta \bullet$ - perturbation in •
$\mathrm{DFFIT}_{\mathrm{i}}$ - change of fit on forecasting after i-th observation is deleted DFFITS $_{i}$ - standardized change in fitted value of $i$ - th case after deletion

E - Mayer and Willke's class of estimators
$\epsilon: n x 1$ - nx1 vector of uncorrelated random error variables
$\epsilon=\left[\begin{array}{ll}\epsilon_{\mathrm{a}} & \epsilon_{\mathrm{b}}\end{array}\right]$
$\epsilon_{\mathrm{a}}:(\mathrm{n}-\mathrm{k}) \mathrm{x} 1$
$\epsilon_{\mathrm{b}}: \mathrm{kx} 1$
$\epsilon$ - scalar controlling collinearity in modified Wampler data set
$\hat{\epsilon}: n x 1-n x 1$ residual vector
$\hat{\epsilon}=\left[\begin{array}{ll}\hat{\epsilon}_{\mathrm{a}} & \hat{\epsilon}_{\mathrm{b}}\end{array}\right]$
$\hat{\epsilon}_{\mathrm{a}}:(\mathrm{n}-\mathrm{k}) \mathrm{x} 1$
$\hat{\epsilon}_{\mathrm{b}}: \mathrm{kx} 1$
$\hat{\epsilon}_{\mathrm{i}}$ - residual term for i-th observation
$\hat{\epsilon}_{-i}:(n-1) x 1-v e c t o r$ of residuals estimator after i-th case deleted $\hat{\epsilon}_{\hat{r}}-\quad$ residuals due to fitting $X_{\hat{r}}$
$\hat{\epsilon}_{R}$ - residual vector under RRE
$\hat{\epsilon}_{\mathrm{s}}-$ residuals due to using $\hat{\beta}_{\mathrm{s}}$
$\left(\hat{\epsilon}_{\mathrm{s}}\right)_{-\mathrm{h}}-$ residuals due to using $\left(\hat{\beta}_{\mathrm{s}}\right)_{-\mathrm{h}}$
$E($.$) - expectation of the scalar, vector or matrix argument$
$e_{j}$ - error in the $j$-th column of $X$
$e_{i}$ - residual vector of $X_{i}$ when regressed on $X_{-i}$
$e_{i}^{-j}$ - residual vector of $X_{i}$ when the regressors are $\left[\begin{array}{lll}X_{-i} & u_{j}\end{array}\right]$
$e_{i j}-j$-th element of $e_{i}$
$\eta_{\mathrm{i}}=\beta_{0}+\sum_{\mathrm{j}=1}^{\mathrm{k}} \beta_{\mathrm{j}} \mathrm{X}_{\mathrm{i} j}-$ expected value from fitted equation
$\eta_{\mathrm{i}}=\sqrt{\lambda_{1}} / \sqrt{\lambda_{\mathrm{i}}} \quad-\quad$ i-th condition index
$\eta_{\mathrm{k}}(-\mathrm{i})$ - k - th condition index computed without the i - th case
$\eta_{\mathrm{i}}$ - efficiency of the LSE of $\delta_{i}$ relative to the ridge estimator $\exp -\quad$ exponent
$f(x)$ - density function of random variable $X$
$f(X)$ - density function of random vector $X$
$f(\bullet)$ - function of a random variable, vector or a matrix
$\mathrm{F}\left(\mathrm{n}_{1}, \mathrm{n}_{2}: \gamma\right)$ - non-central F-distribution with $\mathrm{n}_{1}$ and $\mathrm{n}_{2}$ df and non- centrality parameter $\gamma$
$\mathrm{F}(a: 1, \mathrm{n}-\mathrm{q}-2) \quad$ - tabulated $100 \alpha$ percentile F - value for 1 and $\mathrm{n}-\mathrm{q}-2 \mathrm{df}$
$\mathrm{F}_{\mathrm{L}}$ - lowest partial F - test value
$F_{0}$ - preselected critical value
$F_{i n}$ - constant value chosen as criterion for inclusion of a variable
$F_{i}=\lambda_{i} \hat{\delta}_{i}^{2} / \sigma^{2}$
$F=\operatorname{Diag}\left[\mathrm{f}_{1}, \mathrm{f}_{2}, \ldots, \mathrm{f}_{\mathrm{p}}\right]$ - fraction matrix
$F_{P C I}=\operatorname{Diag}\left(f_{1}^{*},{ }_{P C}, \ldots, f_{p}^{*},{ }_{P C}\right)$

```
F
f
f
f
f
f
f
f
f
FPC - fractional principal component
FPCI - iterative fractional principal component estimator, via the GRR method
FPCV - iterative fractional principal component estimator, via the \(R R\) method
\(G(\bullet)\) - Hayer and Willkes's term for TMSE
\(\gamma\) - non-centrality parameter
\(\gamma\) - scalar controlling collinearity in (2.5.3)
\(\gamma_{1}(\mathrm{k})=\operatorname{tr}\left(\mathrm{V}\left(\hat{\beta}_{\mathrm{R}}\right)\right)-\) total variance of the ridge estimator
```

```
\(\gamma_{2}(\mathrm{k})=\beta^{\prime}(\mathrm{Z}-\mathrm{I})^{\prime}(\mathrm{Z}-\mathrm{I}) \beta-\) total squared bias of ridge estimator
\(\gamma_{i j}-j\)-th eigenvalue of \(X_{-i}^{\prime} X_{-i}\)
\(\Gamma(\bullet)\) - gamma function
\(\Gamma_{p}=\left[\sum_{i=1}^{n}\left(\nu_{i}-\eta_{i}\right)^{2}+\sum_{i=1}^{n} \operatorname{Var}\left(\hat{Y}_{i}\right)\right] / \sigma^{2}-\) Mallows' \(\Gamma_{p}\)
GRRE - generalized ridge regression estimator
GVIF - generalized variance inflation factor
GVIF \(_{\mathrm{i}}\) - \(\mathbf{i}\) - th generalized variance inflation factor
H:nxn - Hat matrix
\(H_{0}\) - null hypothesis
\(h_{i}=h_{i i}=x_{i}^{\prime}\left(X^{\prime} X\right)^{-1} x_{i}-\) leverage values
I - identity matrix
\(I_{p}-\quad(p x p)\) identity matrix
\(\operatorname{IMP} \mathrm{j}_{\mathrm{j}}\) - importance of the j -th variable
\(L_{j}\) - level of importance
\(\inf (X) \stackrel{\text { def }}{=} \min \|X \nu\|\), for \(\|\nu\|=1\)
\(\inf (X)\) - spectral norm of the smallest matrix \(E\) such that \(X+E\) is exactly
    collinear
```

JRE - jackknifed ridge estimator

```
k = p-1 - number of independent variables, excluding the column of one's
k - ridge constant
k - number of influential or outlier observations
K:pxp - diagonal matrix with ridge values }\mp@subsup{\textrm{k}}{\textrm{i}}{}\mathrm{ on the diagonal
k
k
\mp@subsup{\hat{k}}{h}{}}\mathrm{ - estimator of harmonic mean of }\mp@subsup{\textrm{k}}{\textrm{i}}{
\hat{K}}\mathrm{ - estimator of K
\mp@subsup{\hat{k}}{i}{} - i- th element of \hat{K}
\hat{k}
K}={F,\Omega,\mp@subsup{\Lambda}{}{*}}-condition tripl
K(X)}=\sqrt{}{\mp@subsup{\lambda}{1}{}}/\sqrt{}{\mp@subsup{\lambda}{\textrm{p}}{}}\mathrm{ - condition number of X (also indicated by }\mp@subsup{\eta}{\textrm{p}}{}\mathrm{ )
K1}\mp@subsup{1}{i}{\prime},K2\mp@subsup{2}{i}{}-\quad\mathrm{ boundary values of }\mp@subsup{\textrm{k}}{\textrm{i}}{
\kappa
L - unique positive lower triangular matrix of order p
L
\lambdai - i-th eigenvalue of X'X
```

```
\lambdamax - max eigenvalue of ( }\mp@subsup{X}{}{\prime}
    \lambdamin
    \lambda - LSE of }
    A - range set
    A* - range set
    LC - leverage components
    LRR - latent root regression
    LRRE - latent root regression estimator
    LSE - least square estimator
```



```
    \tilde{m} - Swamy's measure of collinearity
    m
    m
    mci = \sum \ \lambda ip
        i=1
    M}\mp@subsup{X}{\mp@subsup{X}{0}{}}{}=\operatorname{plim}(\mp@subsup{n}{}{-1}\mp@subsup{X}{0}{\prime}\mp@subsup{X}{0}{}
    M
```

$\mathbb{M}_{\Delta_{e}}=\operatorname{plim}\left(n^{-1} \Delta_{e}{ }^{\prime} \Delta_{e}\right)$

MDFFIT - statistic DFFITS extended to deletion of more than one data point

MLE - maximum likelihood estimator

MS - mean square

MSE - mean square error random variable in ANOVA

MSE - matrix mean square error (expectation)

MSR - mean square of regression
$\mu$ - expectation of random variable $X$
$\boldsymbol{\mu}$ - expectation of random vector $X$
$\hat{\mu}$ - estimate of $\mu$

WVIF - marginal variance inflation factor
n - scalar, usually the number of rows of Y or X
$\mathrm{N}\left(\mu, \sigma^{2}\right)$ - univariate normal distribution
$N\left(\mu, \sigma^{2} \mathrm{I}\right)$ - multivariate normal distribution with independence and homoscedasticity
$\nu_{i}=\nu\left(X_{1 i}, X_{2 i}, \ldots\right)-$ expected value from true equation for the conditional expectation of ( $\mathrm{Y}_{\mathrm{i}} \mid \mathrm{X}_{1 \mathrm{i}} \ldots \mathrm{X}_{\mathrm{pi}}$ )
$\nu=(\mathrm{n}-\mathrm{p}) \mathrm{df}$ for residual sum of squares

```
0 - vector or matrix whose elements are zero
OLSE - ordinary least square estimator
\omega
\Omega - domain
1 - column vector of ones
P - projection matrix
[P
p - number of independent variables
pa - number of independent variables in the subset }\mp@subsup{X}{a}{
\phiij}=\mp@subsup{v}{ij}{2}/\mp@subsup{\lambda}{i}{} - i,j-th variance-decomposition proportion of variance of
                                    j-th regression coefficient associated with i-th component
\Pi}\mp@subsup{|}{ij}{}=\mp@subsup{\phi}{ij}{}/\mp@subsup{\phi}{\mathrm{ TOTAL }}{}-\quad\mathrm{ variance decomposition proportions
\pi - pi
II:pxp - permutation matrix
P
PC - principal component
PC(i) - principal component estimator, with i components deleted
PCE - principal component estimator
```

```
PCR - principal component regression
plim - probability limit
PMSE - predicted mean square error
Pr(\bullet) - probability
psd - positive semi-definite
PVIF - partial variance inflation factor
Q:nxn - orthogonal-matrix in QR-decomposition of X
Q ( - extra sum of squares
Q ( - weighted pseudo-value
Q(X,Y,\beta) - criterion function to be minimized
\rho}\mp@subsup{}{}{2}\mathrm{ - population value estimated by R}\mp@subsup{}{}{2
R:nxp - upper triangular matrix in QR- decomposition of X
R
R:rxr = [0 I_ ] - restriction matrix
R
R
R
R i,(A) - measure of influence of variable A on variance inflation factor
```

R - coefficient of multiple correlation
$R^{2}$ - coefficient of multiple determination
$\mathrm{R}_{\mathrm{y} \varphi}^{2}$ - coefficient of multiple determination where $\varphi$ is set of independent variables included in model
$R_{y i . \phi}^{2}$ - coefficient of partial determination where $\phi$ denotes set of regressor $X$ variables already in model prior to fitting $X_{i}$.
$R_{i j}$. - coefficient of partial correlation between i-th and $j$-th columns of $X$, while all other columns are held constant
$R_{a}^{2}$ - adjusted coefficient of multiple determination
$R_{i}^{2}$ - coefficient of determination from regression of $X_{i}$ on other independent variables
$\mathrm{R}_{\mathrm{s}}^{2}$ - modified coefficient of multiple determination
$\left(R_{s}^{2}\right)_{-h}-$ modified coefficient of multiple determination after deletion of the h -th column
$R(X)$ - column space (range-space) of $X$
$\hat{\mathbf{r}}$ - estimate of $r(X)=r$
r - number of restrictions
$r_{i}-\quad$ standardized residual
$r_{12}$ - observed coefficient of correlation between the variables represented by the first two columns of $X$
$r_{i j}$ - coefficient of correlation between $i$ - th and $j$-th variables $r(X)$ - rank of $X$

RB - relative bias

RE - relative efficiency

RLS - restricted least squares

RMSE - relative mean square error

RR - ridge regression

RRE - ridge regression estimator

RSS - residual sum of squares
s - number of disconnected subsets
$s$ - number of non- predictive near-singularities

S - positive semi-definite matrix
$S_{i}$ - scalar quantity
$\sigma^{2}$ - population variance scalar
$\sigma_{j}^{2}, L S$ - theoretical variance of $b_{j}, L S$
$\sigma_{\mathrm{Y}_{\mathrm{e}}}^{2}: 1 \mathrm{x} 1=\mathrm{E}\left[\left(\mathrm{Y}_{\mathrm{e}}\right)_{\mathrm{i}}^{2}\right]$
$\sigma_{\mathrm{Y}_{\mathrm{e}} \mathrm{X}}: \mathrm{px} 1=\mathrm{E}\left[\left(\mathrm{X}_{\mathrm{e}}\right)_{\mathrm{i}}\left[\left(\mathrm{Y}_{\mathrm{e}}\right)_{\mathrm{i}}+\epsilon_{\mathrm{i}}\right]\right]$

```
\sigma\epsilon}\mp@subsup{\sigma}{\epsilon}{2}:1\textrm{x}1=\textrm{E}[\mp@subsup{\epsilon}{\textrm{i}}{2}
\sigma
\mp@subsup{\hat{\sigma}}{\mathrm{ CLS }}{2} - consistent estimator of 的
(\hat{\sigma}}\mp@subsup{\mp@code{\epsilon}}{2}{2}+\mp@subsup{\hat{\sigma}}{\mp@subsup{Y}{e}{}}{2}\mp@subsup{)}{\textrm{CLS}}{}-\mathrm{ consistent least square estimator of ( }\mp@subsup{\sigma}{\textrm{e}}{2}+\mp@subsup{\sigma}{\mp@subsup{\textrm{Y}}{\textrm{e}}{\prime}}{2}
\sigma
\mp@subsup{\hat{\sigma}}{}{2}}\mathrm{ - estimate of }\mp@subsup{\sigma}{}{2
```

$\hat{\sigma}(\mathrm{i})$ - estimated error variance when the i - th row of X and Y have been
deleted
$\Sigma_{X_{e}}: \operatorname{pxp}=E\left[\left(X_{e}\right)_{\mathrm{i}}\left(\mathrm{X}_{\mathrm{e}}\right)_{\mathrm{i}}{ }^{\prime}\right]$
$\Sigma_{\Delta_{e}}$ - covariance matrix of error variables of the errors-in-variables model
$s^{2}-\quad 0 L S$ estimate of SSE
$s^{2}$ - Swamy's estimate of $\sigma^{2}$
$s_{y}^{2}$ - sample variance of $Y$
SEMSE - standardized empirical mean square error
SH - shrunken
SHE - shrunken estimator

SS sum of squares

SSA - extra sum of squares due to fitting A after $X$
$\mathrm{SSB}_{\mathrm{p}}$ - sum of squared bias

SSE - error sum of squares or residual sums of squares
$\operatorname{SSE}\left(X_{1}, X_{2}, \ldots, X_{p}\right)-\operatorname{SSE}$ of all independent regressor variables ((...)) included in model

SSR - regression sum of squares

SST0 - total sum of squares (usually corrected for mean)

SX-0LS - subset selection on $X$ where method of estimation is 0LS

SX-TLS - subset selection on $X$ where method of esitmation is TLS

SXY- TLS - subset selection on [X Y] where method of estimation is TLS

SXY- VTLS - subset selection on [X Y] where method of estimation is TLS, with a variant

SVD - singular value decomposition of a matrix .
$\tau_{\mathrm{j}}=\left[(\mathrm{n}-\mathrm{p})^{1 / 2} \kappa_{\mathrm{j}} \mathrm{e}_{\mathrm{j}}\right] /\left\|\mathrm{X}_{\mathrm{j}}\right\|$
$\tau^{2}($ or $\tau)$ - signal- to-noise ratio
$\tau_{\mathrm{i}}^{2}$ - non- centrality parameter of F-distribution associated with $\delta_{i}$
$\tau_{*}^{2}-$ value of $\tau^{2}$ under null hypothesis

```
\tau
0:nx1 - residual vector for alternative model
0=[lll}\mp@subsup{0}{\textrm{a}}{
\hat{0},\mp@subsup{\hat{0}}{\textrm{a}}{},\mp@subsup{\hat{0}}{\textrm{b}}{}: estimates of 0, 有, 访
0 - sum of product terms for two columns of X
TLS - total least squares
TLSE - total least squares estimator
TMSE - total mean square error
TMSE( (Yf) - total mean square error of prediction
tr[\bullet] - trace of a matrix argument
tn
t}\mp@subsup{}{}{2}-t-rati
u - statistic
u
U:nxp - left singular vectors of X
U U
u
```

```
\(u_{i j}\) - \(j\)-th element of \(i\)-th left singular vector of \(X\), or so-called leverage components
\(u_{i}=y_{i}-z_{i}^{\prime} \hat{\delta}_{k}\)
\(u_{j}-j\)-th unit vector in the space \(R^{n}\)
\(\tilde{U}: n x(p+1)\) - left singular vectors of \([\mathrm{XY}]\)
\(\tilde{u}_{i}: n x 1\) - i-th left singular vector of \([\mathrm{XY}]\)
\(V_{J}\) - distribution-free estimate of variance for jackknife estimator
\(V_{J W}\) - distribution-free estimate of variance for weighted jackknife estimator
```

$\mathbf{V}$ - variance-covariance matrix, usually $\mathbf{V}=\sigma^{2} I$
$V: p x p$ - right singular vectors of $X$
$V_{\hat{r}}$ - right singular vectors of $X_{\hat{r}}$
$v_{i}: p x 1$ - i-th right singular vector of $X$
$v_{i j}-j$-th element of the $i$-th eigenvector
$\tilde{\mathrm{V}}:(\mathrm{p}+1) \mathrm{x}(\mathrm{p}+1)$ - right singular vectors of $[\mathrm{X} \mathrm{Y}]$
$\tilde{\mathrm{V}}_{\mathrm{i}}:(\mathrm{p}+1) \mathrm{x} 1$ - i-th right singular vector of $[\mathrm{XY}]$
$\tilde{v}_{i}, j-j$-th component of i-th right singular vector $\tilde{v}_{i}$

```
\mp@subsup{\tilde{v}}{i}{0}}0:px1 - vector containing the first p components of i-th right
    singular vector }\mp@subsup{\tilde{v}}{i}{}\mathrm{ of [X Y]
Var(.), V(.) - variance matrix
```



```
\varphi ~ - ~ r a n d o m ~ v a r i a b l e , ~ v e c t o r ~ o r ~ a ~ m a t r i x ~
VIF - variance inflation factor
VIF (X) - VIF associated with variable A in model (9.3.1)
VIF i or VIF (X (i) - i-th variance inflation factor
VIF ( }\mp@subsup{X}{1}{}|\mp@subsup{X}{2}{\prime}\mathrm{ ) - partial variance inflation factor of }\mp@subsup{X}{i}{}\mathrm{ in }\mp@subsup{X}{(i)}{(i)}=[\mp@subsup{X}{1}{}\mp@subsup{X}{2}{}
VIF}\mp@subsup{\textrm{i}}{(}{(X(i)
VIF
VIF - generalized variance inflation factor
VIF j
VIF -J - generalized variance inflation factor with cases indexed by J
    deleted
V = Diag[1,\ldots,1, Hi},1,\ldots,1] - matrix used in weighted LSE
Wi
WA - Wampler accuracy of b
```

$\left(W^{\prime} \Delta_{1}^{-1} W\right)_{-h},_{-h}-$ submatrix of $\left(W^{\prime} \Delta_{1}^{-1} W\right)$ after deleting $h-t h$ row and column
$\chi^{2}, \chi^{2}(\mathrm{n}, 0)$ or $\chi_{\mathrm{n}}^{2}$, central $\chi^{2}$ - distribution with n df
$\chi^{2}(\mathrm{n}, \gamma)$ - non-central $\chi^{2}$ - distribution with n df and non-centrality parameter $\gamma$

X:nxp - matrix of fixed regressors or explanatory variables.

X:nxp - observed matrix of $p$ explanatory variables
$X_{0}: n x p$ - true but unobservable variables
$X_{e}: n x p$ - measurement or observation errors (unobserved)
$\left(X_{e}\right)_{i}^{\prime}: 1 x p$ - i-th row of $X_{e}$
$\tilde{X}: \operatorname{nxp}-\quad$ perturbed $X$ matrix,$\tilde{X}=X+E$ or $X+\delta X$
$X^{\dagger}=\left(X^{\prime} X\right)^{-1} X^{\prime}-\quad$ pseudo- inverse or generalized inverse of $X$
$X_{\hat{r}}$ - rank $\hat{r}$ matrix approximation of $X$
$X_{i}: n x 1$ - i-th column of $X$
$\bar{x}_{j}$ - mean of $j$-th column
$X_{i j} \quad$ - $\quad$-th element of $j$-th column
$X_{(i)} \quad$ - matrix with i-th column deleted
$X_{I}$ - sub-matrix of $X$ containing columns indexed by $I$
$\left(X^{\prime} X\right)^{-1}{ }^{-1} \quad$ i-th diagonal elements of $\left(X^{\prime} X\right)^{-1}$ $\left(X^{\prime} X\right)_{\text {II }}$ - sub-matrix of $X$ formed by columns and rows indexed by I $x_{i}{ }^{\prime}: 1 \mathrm{xp}-\quad$ i-th row of $X$
$X_{\text {- }}-X$ matrix with $i$ - th row deleted
$\left\{X_{\phi}\right\}$ - set of regressor variables already in model
$\hat{X}$ - approximation to $X$ that satisfies $\operatorname{rank}(\hat{X}) \leq k=\operatorname{rank}(X)$
$\left[\begin{array}{ll}\mathrm{X} & \mathrm{Y}\end{array}\right]: \mathrm{nx}(\mathrm{p}+1)$ - augmented matrix
$[X \hat{Y}]$ - LS approximation of [X Y with orthogonal projection $\hat{Y}$ of $Y$ onto the column space $R(X)$ of $X$
$\left[\begin{array}{ll}\tilde{X} & \mathrm{Y}\end{array}\right]$ - TLS approximation of $\left[\begin{array}{ll}\mathrm{X} & \mathrm{Y}\end{array}\right]$
$[\Delta \tilde{X} \Delta \tilde{Y}]=[X-\tilde{X}, Y-\tilde{Y}]$ - TLS approximation error

Y:nx1 - vector of observed response
$Y_{i}$ - i-th element of $Y$
$[\hat{Y}(-i)]$ - estimated $Y$ obtained by using $\hat{\beta}_{-i}$
$[\hat{\mathrm{Y}}(-\mathrm{i})]_{\mathrm{i}}$ - i-th element of $[\hat{\mathrm{Y}}(-\mathrm{i})]$
$\overline{\mathrm{Y}}$ or $\overline{\mathrm{y}}$ - mean of Y
$Y^{*}: n x 1$ - standardized $Y$ values, corrected for mean, and scaled
$\hat{Y}: n x 1$ - fitted vector of $Y$
$\hat{Y}_{i} \quad$ - -th value of fitted responses
$Y_{-i}$ - sub-vector of $Y$ after deletion of i-th element
$Y_{0}$ :nx1 - true but unobservable variables
$\mathrm{Y}_{\mathrm{e}}: \mathrm{nx} 1$ - measurement or observation errors (unobserved)
Z:nxp - matrix of principal components ( $\mathrm{Z}=\mathrm{XY}$ )
||•\| - norm of a matrix argument

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