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# Likelihood Ratio Tests in Multivariate Linear Model 

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#### Abstract

The aim of this chapter is to review likelihood ratio test procedures in multivariate linear models, focusing on projection matrices. It is noted that the projection matrices to the spaces spanned by mean vectors in hypothesis and alternatives play an important role. Some basic properties are given for projection matrices. The models treated include multivariate regression model, discriminant analysis model, and growth curve model. The hypotheses treated involve a generalized linear hypothesis and no additional information hypothesis, in addition to a usual liner hypothesis. The test statistics are expressed in terms of both projection matrices and sums of squares and products matrices.


Keywords: algebraic approach, additional information hypothesis, generalized linear hypothesis, growth curve model, multivariate linear model, lambda distribution, likelihood ratio criterion (LRC), projection matrix

## 1. Introduction

In this chapter, we review statistical inference, especially likelihood ratio criterion (LRC) in multivariate linear model, focusing on matrix theory. Consider a multivariate linear model with $p$ response variables $y_{1}, \ldots, y_{p}$ and $k$ explanatory or dummy variables $x_{1}, \ldots, x_{k}$. Suppose that $y=\left(y_{1}, \ldots, y_{p}\right)^{\prime}$ and $x=\left(x_{1}, \ldots, x_{k}\right)^{\prime}$ are measured for $n$ subjects, and let the observation of the $i$ th subject be denoted by $y_{i}$ and $x_{i}$. Then, we have the observation matrices given by

$$
\begin{equation*}
\mathbf{Y}=\left(y_{1}, y_{2}, \ldots, \boldsymbol{y}_{n}\right)^{\prime}, \quad \mathrm{X}=\left(\boldsymbol{x}_{1}, \boldsymbol{x}_{2}, \ldots, \boldsymbol{x}_{n}\right)^{\prime} . \tag{1.1}
\end{equation*}
$$

It is assumed that $\boldsymbol{y}_{1}, \ldots, \boldsymbol{y}_{n}$ are independent and have the same covariance matrix $\boldsymbol{\Sigma}$. We express the mean of $\mathbf{Y}$ as follows:

$$
\begin{equation*}
\mathrm{E}(\mathbf{Y})=\boldsymbol{\eta}=\left(\boldsymbol{\eta}_{1}, \ldots, \boldsymbol{\eta}_{p}\right) . \tag{1.2}
\end{equation*}
$$

A multivariate linear model is defined by requiring that

$$
\begin{equation*}
\boldsymbol{\eta}_{i} \in \Omega \text { for all } i=1, \ldots, p, \tag{1.3}
\end{equation*}
$$

where $\Omega$ is a given subspace in the $n$ dimensional Euclid space $\mathrm{R}^{n}$. A typical $\Omega$ is given by

$$
\begin{equation*}
\Omega=\mathcal{R}[\mathbf{X}]=\left\{\boldsymbol{\eta}=\mathbf{X} \boldsymbol{\theta} ; \boldsymbol{\theta}=\left(\theta_{1}, \ldots, \theta_{k}\right)^{\prime},-\infty<\theta_{i}<\infty, i=1, \ldots, k\right\} . \tag{1.4}
\end{equation*}
$$

Here, $\mathcal{R}[\mathbf{X}]$ is the space spanned by the column vectors of $\mathbf{X}$. A general theory for statistical inference on the regression parameter $\Theta$ can be seen in texts on multivariate analysis, e.g., see [1-8]. In this chapter, we discuss with algebraic approach in multivariate linear model.

In Section 2, we consider a multivariate regression model in which $x_{i}^{\prime} s$ are explanatory variables and $\Omega=\mathcal{R}[\mathbf{X}]$. The maximum likelihood estimator (MLE)s and likelihood ratio criterion (LRC) for $\Theta_{2}=O$ are derived by using projection matrices. Here, $\Theta=\left(\Theta_{1} \Theta_{2}\right)$. The distribution of LRC is discussed by multivariate Cochran theorem. It is pointed out that projection matrices play an important role. In Section 3, we give a summary of projection matrices. In Section 4, we consider to test an additional information hypothesis of $y_{2}$ in the presence of $y_{1}$, where $y_{1}=\left(y_{1} \ldots, y_{q}\right)^{\prime}$ and $y_{2}=\left(y_{q+1} \ldots, y_{p}\right)^{\prime}$. In Section 5 , we consider testing problems in discriminant analysis. Section 6 deals with a generalized multivariate linear model which is also called the growth curve model. Some related problems are discussed in Section 7.

## 2. Multivariate regression model

In this section, we consider a multivariate regression model on $p$ response variables and $k$ explanatory variables denoted by $y=\left(y_{1}, \ldots, y_{p}\right)^{\prime}$ and $x=\left(x_{1}, \ldots, x_{k}\right)^{\prime}$, respectively. Suppose that we have the observation matrices given by (1.1). A multivariate regression model is given by

$$
\begin{equation*}
\mathbf{Y}=\mathbf{X} \Theta+\mathbf{E}, \tag{2.1}
\end{equation*}
$$

where $\Theta$ is a $k \times p$ unknown parameter matrix. It is assumed that the rows of the error matrix $\mathbf{E}$ are independently distributed as a $p$ variate normal distribution with mean zero and unknown covariance matrix $\boldsymbol{\Sigma}$, i.e., $\mathrm{N}_{\mathrm{p}}(0, \boldsymbol{\Sigma})$.

Let $L(\Theta, \Sigma)$ be the density function or the likelihood function. Then, we have

$$
-2 \log L(\boldsymbol{\Theta}, \boldsymbol{\Sigma})=n \log |\boldsymbol{\Sigma}|+\operatorname{tr} \boldsymbol{\Sigma}^{-1}(\mathbf{Y}-\mathbf{X} \Theta)^{\prime}(\mathbf{Y}-\mathbf{X} \Theta)+n p \log (2 \pi) .
$$

The maximum likelihood estimators (MLE) $\hat{\boldsymbol{\Theta}}$ and $\hat{\boldsymbol{\Sigma}}$ of $\boldsymbol{\Theta}$ and $\boldsymbol{\Sigma}$ are defined by the maximizers of $L(\Theta, \Sigma)$ or equivalently the minimizers of $-2 \log L(\Theta, \Sigma)$.

Theorem 2.1 Suppose that $\mathbf{Y}$ follows the multivariate regression model in (2.1). Then, the MLEs of $\Theta$ and $\Sigma$ are given as

$$
\begin{aligned}
& \hat{\boldsymbol{\Theta}}=\left(\mathbf{X}^{\prime} \mathbf{X}\right)^{-1} \mathbf{X} \mathbf{Y}, \\
& \hat{\boldsymbol{\Sigma}}=\frac{1}{n}(\mathbf{Y}-\mathbf{X} \hat{\boldsymbol{\Theta}})^{\prime}(\mathbf{Y}-\mathbf{X} \hat{\boldsymbol{\Theta}})=\frac{1}{n} \mathbf{Y}^{\prime}\left(\mathbf{I}_{n}-\mathbf{P}_{\mathbf{x}}\right) \mathbf{Y},
\end{aligned}
$$

where $\mathbf{P}_{\mathbf{x}}=\mathbf{X}\left(\mathbf{X}^{\prime} \mathbf{X}\right)^{-1} \mathbf{X}^{\prime}$. Further, it holds that

$$
-2 \log L(\hat{\Theta}, \hat{\Sigma})=n \log |\hat{\Sigma}|+n p\{\log (2 \pi)+1\} .
$$

Theorem 2.1 can be shown by a linear algebraic method, which is discussed in the next section. Note that $\mathbf{P}_{\mathbf{X}}$ is the projection matrix on the range space $\Omega=\mathcal{R}[\mathbf{X}]$. It is symmetric and idempotent, i.e.

$$
\mathbf{P}_{\mathrm{x}}^{\prime}=\mathbf{P}_{\mathrm{x}}, \quad \mathbf{P}_{\mathrm{x}}^{2}=\mathbf{P}_{\mathrm{x}} .
$$

Next, we consider to test the hypothesis

$$
\begin{equation*}
H: \mathrm{E}(\mathbf{Y})=\mathbf{X}_{1} \Theta_{1} \quad \Leftrightarrow \quad \Theta_{2}=\mathbf{O} \tag{2.2}
\end{equation*}
$$

against $K ; \boldsymbol{\Theta}_{2} \neq \mathrm{O}$, where $\mathbf{X}=\left(\begin{array}{ll}\mathbf{X}_{1} & \mathbf{X}_{2}\end{array}\right), \quad \mathbf{X}_{1} ; n \times j$ and $\boldsymbol{\Theta}=\left(\boldsymbol{\Theta}_{1}^{\prime} \quad \boldsymbol{\Theta}_{2}^{\prime}\right), \quad \boldsymbol{\Theta}_{1} ; j \times p$. The hypothesis means that the last $k-j$ dimensional variate $\boldsymbol{x}_{2}=\left(x_{j+1}, \ldots, x_{k}\right)^{\prime}$ has no additional information in the presence of the first $j$ variate $x_{1}=\left(x_{1}, \ldots, x_{j}\right)^{\prime}$. In general, the likelihood ratio criterion (LRC) is defined by

$$
\begin{equation*}
\lambda=\frac{\max _{H} L(\Theta, \Sigma)}{\max _{K} L(\Theta, \Sigma)} . \tag{2.3}
\end{equation*}
$$

Then we can express

$$
\begin{aligned}
-2 \log \lambda & =\min _{H}\{-2 \log L(\Theta, \Sigma)\}-\min _{K}\{-2 \log L(\Theta, \Sigma)\} \\
& =\min _{H}\left\{n \log |\Sigma|+\operatorname{tr}(\mathbf{Y}-\mathbf{X} \Theta)^{\prime}(\mathbf{Y}-\mathbf{X} \Theta)\right\} \\
& -\min _{K}\left\{n \log |\Sigma|+\operatorname{tr}(\mathbf{Y}-\mathbf{X} \Theta)^{\prime}(\mathbf{Y}-\mathbf{X} \Theta)\right\} .
\end{aligned}
$$

Using Theorem 2.1, we can expressed as

$$
\lambda^{2 / n} \equiv \Lambda=\frac{\left|n \hat{\Sigma}_{\Omega}\right|}{\left|n \hat{\Sigma}_{\omega}\right|} .
$$

Here, $\hat{\Sigma}_{\Omega}$ and $\hat{\Sigma}_{\omega}$ are the maximum likelihood estimators of $\Sigma$ under the model (2.1) or $K$ and $H$, respectively, which are given by

$$
\begin{align*}
n \hat{\mathbf{\Sigma}}_{\Omega} & =\left(\mathbf{Y}-\mathbf{X} \hat{\boldsymbol{\Theta}}_{\Omega}\right)^{\prime}\left(\mathbf{Y}-\mathbf{X} \hat{\Theta}_{\Omega}\right), \quad \hat{\boldsymbol{\Theta}}_{\Omega}=\left(\mathbf{X}^{\prime} \mathbf{X}\right)^{-1} \mathbf{X}^{\prime} \mathbf{Y} \\
& =\mathbf{Y}^{\prime}\left(\mathbf{I}_{n}-\mathbf{P}_{\Omega}\right) \mathbf{Y} \tag{2.4}
\end{align*}
$$

and

$$
\begin{align*}
n \hat{\boldsymbol{\Sigma}}_{\omega} & =\left(\mathbf{Y}-\mathbf{X}_{1} \hat{\Theta}_{1 \omega}\right)^{\prime}\left(\mathbf{Y}-\mathbf{X}_{1} \hat{\Theta}_{1 \omega}\right), \hat{\Theta}_{1 \omega}=\left(\mathbf{X}_{1}^{\prime} \mathbf{X}_{1}\right)^{-1} \mathbf{X}_{1}^{\prime} \mathbf{Y} \\
& =\mathbf{Y}^{\prime}\left(\mathbf{I}_{n}-\mathbf{P}_{\omega}\right) \mathbf{Y} \tag{2.5}
\end{align*}
$$

Summarizing these results, we have the following theorem.
Theorem 2.2Let $\lambda=\Lambda^{n / 2}$ be the LRC for testing $H$ in (2.2). Then, $\Lambda$ is expressed as

$$
\Lambda=\frac{\left|\mathbf{S}_{e}\right|}{\left|\mathbf{S}_{e}+\mathbf{S}_{h}\right|},
$$

where

$$
\begin{equation*}
\mathbf{S}_{e}=n \hat{\boldsymbol{\Sigma}}_{\Omega}, \quad \mathbf{S}_{h}=n \hat{\boldsymbol{\Sigma}}_{\omega}-n \hat{\boldsymbol{\Sigma}}_{\Omega}, \tag{2.7}
\end{equation*}
$$

and $\mathbf{S}_{\Omega}$ and $\mathbf{S}_{\omega}$ are given by (2.4) and (2.5), respectively.
The matrices $\mathbf{S}_{e}$ and $\mathbf{S}_{h}$ in the testing problem are called the sums of squares and products (SSP) matrices due to the error and the hypothesis, respectively. We consider the distribution of $\Lambda$. If a $p \times p$ random matrix $\mathbf{W}$ is expressed as

$$
\mathbf{W}=\sum_{j=1}^{n} \boldsymbol{z}_{j} \boldsymbol{z}_{j}{ }^{\prime}
$$

where $\mathbf{z}_{j} \sim \mathrm{~N}_{\mathrm{p}}\left(\boldsymbol{\mu}_{j}, \boldsymbol{\Sigma}\right)$ and $\boldsymbol{z}_{1}, \ldots, \boldsymbol{z}_{n}$ are independent, $\mathbf{W}$ is said to have a noncentral Wishart distribution with $n$ degrees of freedom, covariance matrix $\boldsymbol{\Sigma}$, and noncentrality matrix $\boldsymbol{\Delta}=\boldsymbol{\mu}_{1} \boldsymbol{\mu}_{1}^{\prime}+\cdots+\boldsymbol{\mu}_{n} \boldsymbol{\mu}_{n}^{\prime}$. We write that $\mathbf{W} \sim \mathrm{W}_{p}(n, \boldsymbol{\Sigma} ; \boldsymbol{\Delta})$. In the special case $\boldsymbol{\Delta}=\mathbf{O}, \mathbf{W}$ is said to have a Wishart distribution, denoted by $\mathbf{W} \sim \mathrm{W}_{p}(n, \boldsymbol{\Sigma})$.

Theorem 2.3 (multivariate Cochran theorem) Let $\mathbf{Y}=\left(\mathbf{y}_{1}, \ldots, \mathbf{y}_{n}\right)^{\prime}$, where $\mathbf{y}_{i} \sim \mathrm{~N}_{p}\left(\boldsymbol{\mu}_{i}, \boldsymbol{\Sigma}\right)$, $i=1, \ldots, n$ and $y_{1}, \ldots, \boldsymbol{y}_{n}$ are independent. Let $\mathbf{A}, \mathbf{A}_{1}$, and $\mathbf{A}_{2}$ be $n \times n$ symmetric matrices. Then:

1. $\mathbf{Y}^{\prime} \mathbf{A} \mathbf{Y} \sim \mathrm{W}_{p}(k, \boldsymbol{\Sigma} ; \boldsymbol{\Omega}) \Leftrightarrow \mathbf{A}^{2}=\mathbf{A}, \quad \operatorname{tr} \mathbf{A}=k, \quad \boldsymbol{\Omega}=\mathrm{E}(\mathbf{Y})^{\prime} \mathbf{A E}(\mathbf{Y})$.
2. $\mathbf{Y}^{\prime} \mathbf{A}_{1} \mathbf{Y}$ and $\mathbf{Y}^{\prime} \mathbf{A}_{2} \mathbf{Y}$ are independent $\Leftrightarrow \mathbf{A}_{1} \mathbf{A}_{2}=\mathbf{O}$.

For a proof of multivariate Cochran theorem, see, e.g. [3, 6-8]. Let $\mathbf{B}$ and $\mathbf{W}$ be independent random matrices following the Wishart distribution $\mathrm{W}_{p}(q, \Sigma)$ and $\mathrm{W}_{p}(n, \Sigma)$, respectively, with $n \geq p$. Then, the distribution of

$$
\Lambda=\frac{|\mathbf{W}|}{|\mathbf{B}+\mathbf{W}|}
$$

is said to be the $p$-dimensional Lambda distribution with $(q, n)$-degrees of freedom and is denoted by $\Lambda_{p}(q, n)$. For distributional results of $\Lambda_{p}(q, n)$, see $[1,3]$.

By using multivariate Cochran's theorem, we have the following distributional results:
Theorem 2.4Let $\mathbf{S}_{e}$ and $\mathbf{S}_{h}$ be the random matrices in (2.7). Let $\Lambda$ be the $\Lambda$-statistic defined by (2.6). Then,

1. $\mathbf{S}_{e} a n d \mathbf{S}_{h}$ are independently distributed as a Wishart distribution $W_{p}(n-k, \boldsymbol{\Sigma})$ and a noncentral Wishart distribution $\mathrm{W}_{p}(k-j, \boldsymbol{\Sigma} ; \boldsymbol{\Delta})$, respectively, where

$$
\begin{equation*}
\Delta=(\mathbf{X} \Theta)^{\prime}\left(\mathbf{P}_{\mathbf{x}}-\mathbf{P}_{\mathbf{x}_{1}}\right) \mathbf{X} \Theta . \tag{2.8}
\end{equation*}
$$

2. Under $H$, the statistic $\Lambda$ is distributed as a lambda distribution $\Lambda_{p}(k-j, n-k)$.

Proof. Note that $\mathbf{P}_{\Omega}=\mathbf{P}_{\mathbf{X}}=\mathbf{X}\left(\mathbf{X}^{\prime} \mathbf{X}\right)^{-1} \mathbf{X}^{\prime}, \mathbf{P}_{\omega}=\mathbf{P}_{\mathbf{X}_{1}}=\mathbf{X}_{1}\left(\mathbf{X}_{1}^{\prime} \mathbf{X}_{1}\right)^{-1} \mathbf{X},{ }^{\prime}$ and $\mathbf{P}_{\Omega} \mathbf{P}_{\omega}=\mathbf{P}_{\omega} \mathbf{P}_{\Omega}$. By multivariate Cochran's theorem the first result (1) follows by checking that

$$
\begin{aligned}
& \left(\mathbf{I}_{n}-\mathbf{P}_{\Omega}\right)^{2}=\left(\mathbf{I}_{n}-\mathbf{P}_{\Omega}\right),\left(\mathbf{P}_{\Omega}-\mathbf{P}_{\omega}\right)^{2}=\left(\mathbf{P}_{\Omega}-\mathbf{P}_{\omega}\right), \\
& \left(\mathbf{I}_{n}-\mathbf{P}_{\Omega}\right)\left(\mathbf{P}_{\Omega}-\mathbf{P}_{\omega}\right)=\mathbf{O} .
\end{aligned}
$$

The second result (2) follows by showing that $\Delta_{0}=\mathbf{O}$, where $\Delta_{0}$ is the $\Delta$ under $H$. This is seen that

$$
\Delta_{0}=\left(\mathbf{X}_{1} \Theta_{1}\right)^{\prime}\left(\mathbf{P}_{\Omega}-\mathbf{P}_{\omega}\right)\left(\mathbf{X}_{1} \Theta_{1}\right)=\mathbf{O}
$$

since $\mathbf{P}_{\Omega} \mathbf{X}_{1}=\mathbf{P}_{\omega} \mathbf{X}_{1}=\mathbf{X}_{1}$.
The matrices $\mathbf{S}_{e}$ and $\mathbf{S}_{h}$ in (2.7) are defined in terms of $n \times n$ matrices $\mathbf{P}_{\Omega}$ and $\mathbf{P}_{\omega}$. It is important to give expressions useful for their numerical computations. We have the following expressions:

$$
\mathbf{S}_{e}=\mathbf{Y}^{\prime} \mathbf{Y}-\mathbf{Y}^{\prime} \mathbf{X}\left(\mathbf{X}^{\prime} \mathbf{X}\right)^{-1} \mathbf{X}^{\prime} \mathbf{Y}, \quad \mathbf{S}_{h}=\mathbf{Y}^{\prime} \mathbf{X}\left(\mathbf{X}^{\prime} \mathbf{X}\right)^{-1} \mathbf{X}^{\prime} \mathbf{Y}-\mathbf{Y}^{\prime} \mathbf{X}_{1}\left(\mathbf{X}_{1}^{\prime} \mathbf{X}_{1}\right)^{-1} \mathbf{X}_{1}^{\prime} \mathbf{Y} .
$$

Suppose that $x_{1}$ is 1 for all subjects, i.e., $x_{1}$ is an intercept term. Then, we can express these in terms of the SSP matrix of $\left(y^{\prime}, x^{\prime}\right)^{\prime}$ defined by

$$
\mathbf{S}=\sum_{i=1}^{n}\binom{\boldsymbol{y}_{i}-\overline{\boldsymbol{y}}}{\boldsymbol{x}_{i}-\overline{\boldsymbol{x}}}\binom{\boldsymbol{y}_{i}-\overline{\boldsymbol{y}}}{\boldsymbol{x}_{i}-\overline{\boldsymbol{x}}}=\left(\begin{array}{ll}
\mathbf{S}_{y y} & \mathbf{S}_{y x}  \tag{2.9}\\
\mathbf{S}_{x y} & \mathbf{S}_{x x}
\end{array}\right),
$$

where $\overline{\mathbf{y}}$ and $\overline{\mathbf{x}}$ are the sample mean vectors. Along the partition of $\mathbf{x}=\left(\mathbf{x}_{1}^{\prime}, \mathbf{x}_{2}^{\prime}\right)^{\prime}$, we partition $\mathbf{S}$ as

$$
\mathbf{S}=\left(\begin{array}{lll}
\mathbf{S}_{y y} & \mathbf{S}_{y 1} & \mathbf{S}_{y 2}  \tag{2.10}\\
\mathbf{S}_{y 1} & \mathbf{S}_{11} & \mathbf{S}_{12} \\
\mathbf{S}_{y 2} & \mathbf{S}_{21} & \mathbf{S}_{22}
\end{array}\right) .
$$

Then,

$$
\begin{equation*}
\mathbf{S}_{e}=\mathbf{S}_{y y \cdot x}, \quad \mathbf{S}_{h}=\mathbf{S}_{y 2 \cdot 1} \mathbf{S}_{22 \cdot 1}^{-1} \mathbf{S}_{2 y \cdot 1} . \tag{2.11}
\end{equation*}
$$

Here, we use the notation $\mathbf{S}_{y y \cdot x}=\mathbf{S}_{y y}-\mathbf{S}_{y x} \mathbf{S}_{x x}^{-1} \mathbf{S}_{x y}, \mathbf{S}_{y 2 \cdot 1}=\mathbf{S}_{y 2}-\mathbf{S}_{y 1} \mathbf{S}_{11}^{-1} \mathbf{S}_{1 y}$, etc. These are derived in the next section by using projection matrices.

## 3. Idempotent matrices and max-mini problems

In the previous section, we have seen that idempotent matrices play an important role on statistical inference in multivariate regression model. In fact, letting $\mathrm{E}(\boldsymbol{Y})=\boldsymbol{\eta}=\left(\boldsymbol{\eta}_{1}, \ldots, \boldsymbol{\eta}_{p}\right)$ consider a model satisfying

$$
\begin{equation*}
\eta_{i} \in \Omega=\mathcal{R}[\mathbf{R}], \quad \text { for all } i=1, \ldots, p, \tag{3.1}
\end{equation*}
$$

Then the MLE of $\boldsymbol{\Theta}$ is $\hat{\boldsymbol{\Theta}}=\left(\mathbf{X}^{\prime} \mathbf{X}\right)^{-1} \mathbf{X} \mathbf{Y}$, and hence the MLE of $\eta$ is denoted by

$$
\hat{\eta}_{\Omega}=\mathbf{X} \hat{\Theta}=\mathbf{P}_{\Omega} \mathbf{Y} .
$$

Here, $\mathbf{P}_{\Omega}=\mathbf{X}\left(\mathbf{X}^{\prime} \mathbf{X}\right)^{-1} \mathbf{X}^{\prime}$. Further, the residual sums of squares and products (RSSP) matrix is expressed as

$$
\mathbf{S}_{\Omega}=\left(\mathbf{Y}-\hat{\eta}_{\Omega}\right)^{\prime}\left(\mathbf{Y}-\hat{\eta}_{\Omega}\right)=\mathbf{Y}^{\prime}\left(\mathbf{I}_{n}-\mathbf{P}_{\Omega}\right) \mathbf{Y} .
$$

Under the hypothesis (2.2), the spaces $\eta_{i}^{\prime}$ s belong are the same and are given by $\omega=\mathcal{R}\left[\mathbf{X}_{1}\right]$. Similarly, we have

$$
\begin{aligned}
\hat{\eta}_{\omega} & =\mathbf{X} \hat{\Theta}_{\omega}=\mathbf{P}_{\omega} \mathbf{Y}, \\
\mathbf{S}_{\omega} & =\left(\mathbf{Y}-\hat{\eta}_{\omega}\right)^{\prime}\left(\mathbf{Y}-\hat{\eta}_{\omega}\right)=\mathbf{Y}^{\prime}\left(\mathbf{I}_{n}-\mathbf{P}_{\omega}\right) \mathbf{Y},
\end{aligned}
$$

where $\widehat{\boldsymbol{\Theta}}_{\omega}=\left(\widehat{\boldsymbol{\Theta}}_{1 \omega}^{\prime} \mathbf{O}\right)^{\prime}$ and $\hat{\boldsymbol{\Theta}}_{1 \omega}=\left(\mathbf{X}_{1}^{\prime} \mathbf{X}_{1}^{\prime}\right)^{-1} \mathbf{X}_{1}^{\prime} \mathbf{Y}$. The LR criterion is based on the following decomposition of SSP matrices;

$$
\begin{aligned}
\mathbf{S}_{\omega}=\mathbf{Y}^{\prime}\left(\mathbf{I}_{n}-\mathbf{P}_{\omega}\right) \mathbf{Y} & =\mathbf{Y}^{\prime}\left(\mathbf{I}_{n}-\mathbf{P}_{\Omega}\right) \mathbf{Y}+\mathbf{Y}^{\prime}\left(\mathbf{P}_{\Omega}-\mathbf{P}_{\omega}\right) \mathbf{Y} \\
& =\mathbf{S}_{e}+\mathbf{S}_{h} .
\end{aligned}
$$

The degrees of freedom in the $\Lambda$ distribution $\Lambda_{p}\left(f_{k} f_{e}\right)$ are given by

$$
f_{e}=n-\operatorname{dim}[\Omega], \quad f_{h}=k-j=\operatorname{dim}[\Omega]-\operatorname{dim}[\omega] .
$$

In general, an $n \times n$ matrix $\mathbf{P}$ is called idempotent if $\mathbf{P}^{2}=\mathbf{P}$. A symmetric and idempotent matrix is called projection matrix. Let $\mathbf{R}^{n}$ be the $n$ dimensional Euclid space, and $\Omega$ be a subspace in $\mathbf{R}^{n}$. Then, any $n \times 1$ vector $\mathbf{y}$ can be uniquely decomposed into direct sum, i.e.,

$$
\begin{equation*}
\boldsymbol{y}=\boldsymbol{u}+\boldsymbol{v}, \quad \boldsymbol{u} \in \Omega, \quad \mathbf{v} \in \Omega^{\perp}, \tag{3.2}
\end{equation*}
$$

where $\Omega^{\perp}$ is the orthocomplement space. Using decomposition (3.2), consider a mapping

$$
\mathbf{P}_{\Omega}: y \rightarrow u, \quad \text { i.e. } \mathbf{P}_{\Omega} y=u
$$

The mapping is linear, and hence it is expressed as a matrix. In this case, $\boldsymbol{u}$ is called the orthogonal projection of $y$ into $\Omega$, and $\mathbf{P}_{\Omega}$ is also called the orthogonal projection matrix to $\Omega$. Then, we have the following basic properties:
(P1) $\mathbf{P}_{\Omega}$ is uniquely defined;
(P2) $\mathbf{I}_{n}-\mathbf{P}_{\Omega}$ is the projection matrix to $\Omega^{\perp}$;
(P3) $\mathbf{P}_{\Omega}$ is a symmetric idempotent matrix;
(P4) $\mathcal{R}\left[\mathbf{P}_{\Omega}\right]=\Omega$, and $\operatorname{dim}[\Omega]=\operatorname{tr} \mathbf{P}_{\Omega}$;
Let $\omega$ be a subset of $\Omega$. Then, we have the following properties:
(P5) $\mathbf{P}_{\Omega} \mathbf{P}_{\omega}=\mathbf{P}_{\omega} \mathbf{P}_{\Omega}=\mathbf{P}_{\omega}$.
(P6) $\mathbf{P}_{\Omega}-\mathbf{P}_{\omega}=\mathbf{P}_{\omega^{\perp} \cap \Omega}$, where $\omega^{\perp}$ is the orthocomplement space of $\omega$.
(P7) Let $\mathbf{B}$ be a $q \times n$ matrix, and let $\mathrm{N}(\mathbf{B})=\{\boldsymbol{y} ; \mathbf{B} y=\mathbf{0}\}$. If $\omega=\mathrm{N}[\mathbf{B}] \cap \Omega$, then $\omega^{\perp} \cap \Omega=R\left[\mathbf{P}_{\Omega} \mathbf{B}{ }^{\prime}\right]$.
For more details, see, e.g. [3, 7, 9, 10].
The MLEs and LRC in multivariate regression model are derived by using the following theorem.

## Theorem 3.1

1. Consider a function of $f(\boldsymbol{\Sigma})=\log |\Sigma|+\operatorname{tr} \boldsymbol{\Sigma}^{-1} \mathbf{S}$ of $p \times p$ positive definite matrix. Then, $f(\boldsymbol{\Sigma})$ takes uniquely the minimum at $\boldsymbol{\Sigma}=\mathbf{S}$, and the minimum value is given by

$$
\min _{\Sigma>0} f(\mathbf{\Sigma})=f(\mathbf{S})+p
$$

2. Let $\mathbf{Y}$ be an $n \times p$ known matrix and $\mathbf{X}$ an $n \times k$ known matrix of rank $k$. Consider a function of $p \times p$ positive definite matrix $\boldsymbol{\Sigma}$ and $k \times p$ matrix $\boldsymbol{\Theta}=\left(\theta_{i j}\right)$ given by

$$
g(\Theta, \Sigma)=m \log |\Sigma|+\operatorname{tr} \Sigma^{-1}(\mathbf{Y}-\mathbf{X} \Theta)^{\prime}(\mathbf{Y}-\mathbf{X} \Theta),
$$

where $m>0,-\infty<\theta_{i j}<\infty$, for $i=1, \ldots, k ; j=1, \ldots, p$. Then, $g(\boldsymbol{\Theta}, \boldsymbol{\Sigma})$ takes the minimum at

$$
\boldsymbol{\Theta}=\hat{\boldsymbol{\Theta}}=\left(\mathbf{X}^{\prime} \mathbf{X}\right)^{-1} \mathbf{X}^{\prime} \mathbf{Y}, \quad \Sigma=\hat{\boldsymbol{\Sigma}}=\frac{1}{m} \mathbf{Y}^{\prime}\left(\mathbf{I}_{n}-\mathbf{P}_{\mathbf{x}}\right) \mathbf{Y},
$$

and the minimum value is given by $m \log |\hat{\Sigma}|+m p$.
Proof. Let $\ell_{1}, \ldots, \ell_{p}$ be the characteristic roots of $\boldsymbol{\Sigma}^{-1} \mathbf{S}$. Note that the characteristic roots of $\boldsymbol{\Sigma}^{-1} \mathbf{S}$ and $\boldsymbol{\Sigma}^{-1 / 2} \mathbf{S} \Sigma^{-1 / 2}$ are the same. The latter matrix is positive definite, and hence we may assume $\ell_{1} \geq \cdots \geq \ell_{p}>0$. Then

$$
\begin{aligned}
f(\Sigma)-f(\mathbf{S}) & =\log \left|\Sigma \mathbf{S}^{-1}\right|+\operatorname{tr}\left(\Sigma^{-1} \mathbf{S}\right)-p \\
& =-\log \left|\Sigma^{-1} \mathbf{S}\right|+\operatorname{tr}\left(\Sigma^{-1} \mathbf{S}\right)-p \\
& =\sum_{i=1}^{p}\left(-\log \ell_{i}+\ell_{i}-1\right) \geq 0 .
\end{aligned}
$$

The last inequality follows from $x-1 \geq \log x(x>0)$. The equality holds if and only if $\ell_{1}=\cdots=\ell_{p}=1 \Leftrightarrow \boldsymbol{\Sigma}=\mathbf{S}$.

Next, we prove 2. we have

$$
\begin{aligned}
& \operatorname{tr} \Sigma^{-1}(\mathbf{Y}-\mathbf{X} \Theta)^{\prime}(\mathbf{Y}-\mathbf{X} \Theta) \\
& \quad=\operatorname{tr} \Sigma^{-1}(\mathbf{Y}-\mathbf{X} \hat{\Theta})^{\prime}(\mathbf{Y}-\mathbf{X} \hat{\Theta})+\operatorname{tr} \Sigma^{-1}\{\mathbf{X}(\hat{\Theta}-\Theta)\}^{\prime} \mathbf{X}(\hat{\Theta}-\Theta) \\
& \quad \geq \operatorname{tr} \Sigma^{-1} \mathbf{Y}^{\prime}\left(\mathbf{I}_{n}-\mathbf{P}_{\mathbf{X}}\right) \mathbf{Y} .
\end{aligned}
$$

The first equality follows from that $\mathbf{Y}-\mathbf{X} \boldsymbol{\Theta}=\mathbf{Y}-\mathbf{X} \widehat{\boldsymbol{\Theta}}+\mathbf{X}(\hat{\boldsymbol{\Theta}}-\boldsymbol{\Theta})$ and $(\mathbf{Y}-\mathbf{X} \hat{\boldsymbol{\Theta}})^{\prime} \mathbf{X}(\hat{\boldsymbol{\Theta}}-\boldsymbol{\Theta})=\mathbf{O}$. In the last step, the equality holds when $\Theta=\hat{\boldsymbol{\Theta}}$. The required result is obtained by noting that $\hat{\Theta}$ does not depend on $\Sigma$ and combining this result with the first result 1 .

Theorem 3.2Let $\mathbf{X}$ be an $n \times k$ matrix of rank $k$, and let $\Omega=\mathcal{R}[\mathbf{X}]$ which is defined also by the set $\{y$ : $\boldsymbol{y}=\mathbf{X} \boldsymbol{\theta}\}$, where $\boldsymbol{\theta}$ is a $k \times 1$ unknown parameter vector. Let $\mathbf{C}$ be a $c \times k$ matrix of rank $c$, and define $\omega$ by the set $\{\boldsymbol{y}: \boldsymbol{y}=\boldsymbol{X} \boldsymbol{\theta}, \mathbf{C} \boldsymbol{\theta}=\mathbf{0}\}$. Then,

1. $\mathbf{P}_{\Omega}=\mathbf{X}\left(\mathbf{X}^{\prime} \mathbf{X}\right)^{-1} \mathbf{X}^{\prime}$.
2. $\mathbf{P}_{\Omega}-\mathbf{P}_{\omega}=\mathbf{X}\left(\mathbf{X}^{\prime} \mathbf{X}\right)^{-1} \mathbf{C}^{\prime}\left\{\mathbf{C}\left(\mathbf{X}^{\prime} \mathbf{X}\right)^{-1} \mathbf{C}\right\}^{-1} \mathbf{C}\left(\mathbf{X}^{\prime} \mathbf{X}\right)^{-1} \mathbf{X}^{\prime}$.

Proof. 1 Let $\hat{\boldsymbol{y}}=\mathbf{X}\left(\mathbf{X}^{\prime} \mathbf{X}\right)^{-1} \mathbf{X}^{\prime}$ and consider a decomposition $y=\hat{y}+(y-\hat{y})$. Then, $\hat{y}^{\prime}(y-\hat{y})=\mathbf{0}$. Therefore, $\mathbf{P}_{\Omega} \boldsymbol{y}=\hat{\boldsymbol{y}}$ and hence $\mathbf{P}_{\Omega}=\mathbf{X}\left(\mathbf{X}^{\prime} \mathbf{X}\right)^{-1} \mathbf{X}^{\prime}$.
2. Since $\mathbf{C} \boldsymbol{\theta}=\mathbf{C}\left(\mathbf{X}^{\prime} \mathbf{X}\right)^{-1} \mathbf{X}^{\prime} \cdot \mathbf{X} \boldsymbol{\theta}$, we can write $\omega=\mathrm{N}[\mathbf{B}] \cap \Omega$, where $\mathbf{B}=\mathbf{C}\left(\mathbf{X}^{\prime} \mathbf{X}\right)^{-1} \mathbf{X}^{\prime}$. Using (P7),

$$
\omega^{\perp} \cap \Omega=\mathcal{R}\left[\mathbf{P}_{\Omega} \mathbf{B}^{\prime}\right]=\mathcal{R}\left[\mathbf{X}\left(\mathbf{X}^{\prime} \mathbf{X}\right)^{-1} \mathbf{C}^{\prime}\right] .
$$

The final result is obtained by using 1 and (P7).
Consider a special case $\mathbf{C}=\left(\begin{array}{ll}\mathbf{O} & \mathbf{I}_{k-q}\end{array}\right)$. Then $\omega=\mathcal{R}\left[\mathbf{X}_{1}\right]$, where $\mathbf{X}=\left(\begin{array}{ll}\mathbf{X}_{1} & \mathbf{X}_{2}\end{array}\right), \mathbf{X}_{1}: n \times q$. We have the following results:

$$
\begin{aligned}
& \omega^{\perp} \cap \Omega=\mathcal{R}\left[\left(\mathbf{I}_{n}-\mathbf{P}_{\mathrm{X}_{1}}\right) \mathbf{X}_{2}\right], \\
& \mathbf{P}_{\omega^{\perp} \cap \Omega}=\left(\mathbf{I}_{n}-\mathbf{P}_{\mathbf{x}_{1}}\right) \mathbf{X}_{2}\left\{\mathbf{X}_{2}^{\prime}\left(\mathbf{I}_{n}-\mathbf{P}_{\mathbf{x}_{1}}\right) \mathbf{X}_{2}\right\}^{-1} \mathbf{X}_{2}^{\prime}\left(\mathbf{I}_{n}-\mathbf{P}_{\mathbf{x}_{1}}\right) .
\end{aligned}
$$

The expressions (2.11) for $\mathbf{S}_{e}$ and $\mathbf{S}_{h}$ in terms of $\mathbf{S}$ can be obtained from projection matrices based on

$$
\begin{aligned}
\Omega & =\mathcal{R}[\mathbf{X}]=\mathcal{R}\left[\mathbf{1}_{n}\right]+\mathcal{R}\left[\left(\mathbf{I}_{n}-\mathbf{P}_{\mathbf{1}_{n}}\right) \mathrm{X}\right], \\
\omega^{\perp} \cap \Omega & \left.=\mathcal{R}\left[\left(\mathbf{I}_{n}-\mathbf{P}_{\mathbf{1}_{n}}-\mathbf{P}_{\left(\mathbf{I}_{n}-\mathbf{P}_{1_{n}}\right)}\right) \mathbf{X X}_{1}\right) \mathbf{X}_{2}\right] .
\end{aligned}
$$

## 4. General linear hypothesis

In this section, we consider to test a general linear hypothesis

$$
\begin{equation*}
H_{g}: \mathbf{C} \Theta \mathbf{D}=\mathbf{O}, \tag{4.1}
\end{equation*}
$$

against alternatives $K_{g}: \mathbf{C} \Theta \mathbf{D} \neq \mathbf{O}$ under a multivariate linear model given by (2.1), where $\mathbf{C}$ is a $c \times k$ given matrix with rank $c$ and $\mathbf{D}$ is a $p \times d$ given matrix with rank $d$. When $\mathbf{C}=(\mathbf{O}$ $\mathbf{I}_{k-j}$ ) and $\mathbf{D}=\mathbf{I}_{p}$, the hypothesis $H_{g}$ becomes $H: \Theta_{2}=\mathbf{O}$.
For the derivation of LR test of (4.1), we can use the following conventional approach: If $\mathbf{U}=\mathrm{YD}$, then the rows of $\mathbf{U}$ are independent and normally distributed with the identical covariance matrix $\mathbf{D}^{\prime} \Sigma \mathbf{D}$, and

$$
\begin{equation*}
\mathrm{E}(\mathbf{U})=\mathbf{X} \Xi, \tag{4.2}
\end{equation*}
$$

where $\boldsymbol{\Xi}=\boldsymbol{\Theta} \mathbf{D}$. The hypothesis (4.1) is expressed as

$$
\begin{equation*}
H_{g}: \mathbf{C} \Xi=\mathbf{O} . \tag{4.3}
\end{equation*}
$$

Applying a general theory for testing $H_{g}$ in (2.1), we have the LRC $\lambda$ :

$$
\begin{equation*}
\lambda^{2 / n}=\Lambda=\frac{\left|\mathbf{S}_{e}\right|}{\left|\mathbf{S}_{e}+\mathbf{S}_{h}\right|}, \tag{4.4}
\end{equation*}
$$

where

$$
\begin{aligned}
\mathbf{S}_{e} & =\mathbf{U}^{\prime}\left(\mathbf{I}_{n}-\mathbf{P}_{\mathbf{x}}\right) \mathbf{U} \\
& =\mathbf{D}^{\prime} \mathbf{Y}^{\prime}\left(\mathbf{I}_{n}-\mathbf{P}_{\mathbf{A}}\right) \mathbf{Y} \mathbf{D},
\end{aligned}
$$

and

$$
\begin{aligned}
\mathbf{S}_{h} & =\left(\mathbf{C}\left(\mathbf{X}^{\prime} \mathbf{X}\right)^{-1} \mathbf{X}^{\prime} \mathbf{U}\right)^{\prime}\left\{\mathbf{C}\left(\mathbf{X}^{\prime} \mathbf{X}\right)^{-1} \mathbf{C}^{\prime}\right\}^{-1} \mathbf{C}\left(\mathbf{X}^{\prime} \mathbf{X}\right)^{-1} \mathbf{X}^{\prime} \mathbf{U}, \\
& =\left(\mathbf{C}\left(\mathbf{X}^{\prime} \mathbf{X}\right)^{-1} \mathbf{X}^{\prime} \mathbf{Y} D\right)^{\prime}\left\{\mathbf{C}\left(\mathbf{X}^{\prime} \mathbf{X}\right)^{-1} \mathbf{C}^{\prime}\right\}^{-1} \mathbf{C}\left(\mathbf{X}^{\prime} \mathbf{X}\right)^{-1} \mathbf{X}^{\prime} \mathbf{Y} D .
\end{aligned}
$$

Theorem 4.1The statistic $\Lambda$ in (4.4) is an LR statistic for testing (4.1) under (2.1). Further, under $H_{g^{\prime}} \Lambda \sim \Lambda_{d}(c, n-k)$.

Proof. Let $\mathbf{G}=\left(\mathbf{G}_{1} \mathbf{G}_{2}\right)$ be a $p \times p$ matrix such that $\mathbf{G}_{1}=\mathbf{D}, \mathbf{G}_{1}^{\prime} \mathbf{G}_{2}=\mathbf{O}$, and $|\mathbf{G}| \neq 0$. Consider a transformation from $\mathbf{Y}$ to $(\mathbf{U} \mathbf{V})=\mathbf{Y}\left(\mathbf{G}_{1} \mathbf{G}_{2}\right)$.

Then the rows of $(\mathbf{U} \mathbf{V})$ are independently normal with the same covariance matrix

$$
\Psi=\mathbf{G}^{\prime} \Sigma \mathbf{G}=\left(\begin{array}{ll}
\Psi_{11} & \Psi_{12} \\
\Psi_{21} & \Psi_{22}
\end{array}\right), \quad \Psi_{12}: d \times(p-d),
$$

and

$$
\begin{aligned}
E[(\mathbf{U} \mathbf{V})] & =\mathbf{X} \Theta\left(\mathbf{G}_{1} \mathbf{G}_{2}\right) \\
& =\mathbf{X}(\Xi \Delta), \quad \Xi=\Theta \mathbf{G}_{1}, \Delta=\Theta \mathbf{G}_{2} .
\end{aligned}
$$

The conditional of $\mathbf{V}$ given $\mathbf{U}$ is normal. The rows of $\mathbf{V}$ given $\mathbf{U}$ are independently normal with the same covariance matrix $\Psi_{11 \cdot 2}$, and

$$
\begin{aligned}
\mathrm{E}(\mathbf{V} \mid \mathbf{U}) & =\mathbf{X} \Delta+(\mathbf{U}-\mathrm{X} \boldsymbol{\Xi}) \Gamma \\
& =\mathbf{X} \Delta^{*}+\mathbf{U} \Gamma,
\end{aligned}
$$

where $\Delta^{*}=\Delta-\boldsymbol{\Xi} \Gamma$ and $\Gamma=\Psi_{11}^{-1} \Psi_{12}$. We see that the maximum likelihood of $\mathbf{V}$ given $\mathbf{U}$ does not depend on the hypothesis. Therefore, an LR statistic is obtained from the marginal distribution of $\mathbf{U}$, which implies the results required.

## 5. Additional information tests for response variables

We consider a multivariate regression model with an intercept term $x_{0}$ and $k$ explanatory variables $x_{1}, \ldots, x_{k}$ as follows.

$$
\begin{equation*}
\mathbf{Y}=\mathbf{1} \theta^{\prime}+\mathbf{X} \Theta+\mathbf{E}, \tag{5.1}
\end{equation*}
$$

where $\mathbf{Y}$ and $\mathbf{X}$ are the observation matrices on $y=\left(y_{1}, \ldots, y_{p}\right)^{\prime}$ and $x=\left(x_{1}, \ldots, x_{k}\right)^{\prime}$. We assume that the error matrix $\mathbf{E}$ has the same property as in (2.1), and $\operatorname{rank}\left(\mathbf{1}_{n} \mathbf{X}\right)=k+1$. Our interest is
to test a hypothesis $H_{2.1}$ on no additional information of $y_{2}=\left(y_{q+1}, \ldots, y_{p}\right)^{\prime}$ in presence of $y_{1}=\left(y_{1}, \ldots, y_{q}\right)^{\prime}$.

Along the partition of $y$ into $\left(y_{1}{ }^{\prime}, y_{2}{ }^{\prime}\right)$ let $\mathbf{Y}, \theta, \Theta$, and $\Sigma$ partition as

$$
\begin{aligned}
& \mathbf{Y}=\left(\mathbf{Y}_{1} \mathbf{Y}_{2}\right), \quad \Theta=\left(\Theta_{1} \Theta_{2}\right), \\
& \boldsymbol{\theta}=\binom{\theta_{1}}{\theta_{2}}, \quad \Sigma=\left(\begin{array}{ll}
\Sigma_{11} & \Sigma_{12} \\
\Sigma_{21} & \Sigma_{22}
\end{array}\right) .
\end{aligned}
$$

The conditional distribution of $\mathbf{Y}_{2}$ given $\mathbf{Y}_{1}$ is normal with mean

$$
\begin{align*}
\mathrm{E}\left(\mathbf{Y}_{2} \mid \mathbf{Y}_{1}\right) & =\mathbf{1} \theta_{2}^{\prime}+\mathbf{X} \Theta_{2}+\left(\mathbf{Y}_{1}-\mathbf{1}_{n} \theta_{1}^{\prime}-\mathbf{X} \Theta_{1}\right) \Sigma_{11}^{-1} \Sigma_{12} \\
& =\mathbf{1}_{n} \tilde{\theta}_{02}^{\prime}+\mathbf{X} \tilde{\Theta}_{2}+\mathbf{Y}_{1} \Sigma_{11}^{-1} \Sigma_{12}, \tag{5.2}
\end{align*}
$$

and the conditional covariance matrix is expressed as

$$
\begin{equation*}
\operatorname{Var}\left[\operatorname{vec}\left(\mathbf{Y}_{2} \mid \mathbf{Y}_{1}\right)\right]=\boldsymbol{\Sigma}_{22 \cdot 1} \otimes \mathbf{I}_{n}, \tag{5.3}
\end{equation*}
$$

where $\boldsymbol{\Sigma}_{22 \cdot 1}=\boldsymbol{\Sigma}_{22}-\boldsymbol{\Sigma}_{21} \boldsymbol{\Sigma}_{11}^{-1} \boldsymbol{\Sigma}_{12}$, and

$$
\tilde{\theta}_{2}^{\prime}=\theta_{2}^{\prime}-\theta_{1}^{\prime} \Sigma_{11}^{-1} \Sigma_{12}, \quad \tilde{\Theta}_{2}=\Theta_{2}-\Theta_{1} \Sigma_{11}^{-1} \Sigma_{12} .
$$

Here, for an $n \times p$ matrix $\mathbf{Y}=\left(\mathbf{y}_{(1)}, \ldots, \mathbf{y}_{(\mathrm{p},)}\right)$ vec $(\mathbf{Y})$ means an $n p$-vector $\left(\mathbf{y}_{(1)}^{\prime}, \ldots, \mathbf{y}_{(p)}^{\prime}\right)^{\prime}$. Now we define the hypothesis $H_{2 \cdot 1}$ as

$$
\begin{equation*}
H_{2 \cdot 1} \Theta_{2}=\Theta_{1} \Sigma_{11}^{-1} \Sigma_{12} \Leftrightarrow \tilde{\Theta}_{2}=0 . \tag{5.4}
\end{equation*}
$$

The hypothesis $H_{2 \cdot 1}$ means that $y_{2}$ after removing the effects of $y_{1}$ does not depend on $x$. In other words, the relationship between $y_{2}$ and $x$ can be described by the relationship between $y_{1}$ and $\mathbf{x}$. In this sense, $\mathbf{y}_{2}$ is redundant in the relationship between $y$ and $x$.

The LR criterion for testing the hypothesis $H_{2 \cdot 1}$ against alternatives $K_{2 \cdot 1}: \widetilde{\Theta}_{2 \cdot 1} \neq \mathbf{O}$ can be obtained through the following steps.
(D1) The density function of $\mathbf{Y}=\left(\mathbf{Y}_{1} \mathbf{Y}_{2}\right)$ can be expressed as the product of the marginal density function of $\mathbf{Y}_{1}$ and the conditional density function of $\mathbf{Y}_{2}$ given $\mathbf{Y}_{1}$. Note that the density functions of $\mathbf{Y}_{1}$ under $H_{2.1}$ and $K_{2.1}$ are the same.
(D2) The spaces spanned by each column of $E\left(\mathbf{Y}_{2} \mid \mathbf{Y}_{1}\right)$ are the same, and let the spaces under $K_{2 \cdot 1}$ and $H_{2 \cdot 1}$ denote by $\Omega$ and $\omega$, respectively. Then

$$
\Omega=\mathcal{R}\left[\left(\mathbf{1}_{n} \mathbf{Y}_{1} \mathbf{X}\right)\right], \quad \omega=\mathcal{R}\left[\left(\mathbf{1}_{n} \mathbf{Y}_{1}\right)\right],
$$

and $\operatorname{dim}(\Omega)=q+k+1, \operatorname{dim}(\omega)=k+1$.
(D3) The likelihood ratio criterion $\lambda$ is expressed as

$$
\lambda^{2 / n}=\Lambda=\frac{\left|\mathbf{S}_{\Omega}\right|}{\left|\mathbf{S}_{\omega}\right|}=\frac{\left|\mathbf{S}_{\Omega}\right|}{\left|\mathbf{S}_{\Omega}+\left(\mathbf{S}_{\omega}-\mathbf{S}_{\Omega}\right)\right|}
$$

where $\mathbf{S}_{\Omega}=\mathbf{Y}_{2}^{\prime}\left(\mathbf{I}_{n}-\mathbf{P}_{\Omega}\right) \mathbf{Y}_{2}$ and $\mathbf{S}_{\omega}=\mathbf{Y}_{2}^{\prime}\left(\mathbf{I}_{n}-\mathbf{P}_{\omega}\right) \mathbf{Y}_{2}$.
(D4) Note that $E\left(\mathbf{Y}_{2} \mid \mathbf{Y}_{1}\right)^{\prime}\left(\mathbf{P}_{\omega}-\mathbf{P}_{\omega}\right) \mathrm{E}\left(\mathbf{Y}_{2} \mid \mathbf{Y}_{1}\right)=\mathbf{O}$ under $H_{2 \cdot 1}$. The conditional distribution of $\Lambda$ under $H_{2 \cdot 1}$ is $\Lambda_{p-q}(k, n-q-k-1)$, and hence the distribution of $\Lambda$ under $H_{2 \cdot 1}$ is $\Lambda_{p-q}(k, n-q-k-1)$.

Note that the $\Lambda$ statistic is defined through $\mathbf{Y}_{2}^{\prime}\left(\mathbf{I}_{n}-\mathbf{P}_{\Omega}\right) \mathbf{Y}_{2}$ and $\mathbf{Y}_{2}^{\prime}\left(\mathbf{P}_{\Omega}-\mathbf{P}_{\omega}\right) \mathbf{Y}_{2}$, which involve $n \times n$ matrices. We try to write these statistics in terms of the SSP matrix of $\left(y^{\prime}, x^{\prime}\right)^{\prime}$ defined by

$$
\begin{aligned}
\mathbf{S} & =\sum_{i=1}^{n}\binom{\boldsymbol{y}_{i}-\overline{\boldsymbol{y}}}{\boldsymbol{x}_{i}-\overline{\boldsymbol{x}}}\binom{\boldsymbol{y}_{i}-\overline{\boldsymbol{y}}}{\boldsymbol{x}_{i}-\overline{\boldsymbol{x}}} \\
& =\left(\begin{array}{ll}
\mathbf{S}_{y y} & \mathbf{S}_{y x} \\
\mathbf{S}_{x y} & \mathbf{S}_{x x}
\end{array}\right),
\end{aligned}
$$

where $\overline{\mathbf{y}}$ and $\overline{\mathbf{x}}$ are the sample mean vectors. Along the partition of $\mathbf{y}=\left(\mathbf{y}_{1}^{\prime}, \mathbf{y}_{2}^{\prime}\right)^{\prime}$, we partition $\mathbf{S}$ as

$$
\mathbf{S}=\left(\begin{array}{lll}
\mathbf{S}_{11} & \mathbf{S}_{12} & \mathbf{S}_{1 x} \\
\mathbf{S}_{21} & \mathbf{S}_{22} & \mathbf{S}_{2 x} \\
\mathbf{S}_{x 1} & \mathbf{S}_{x 2} & \mathbf{S}_{x x}
\end{array}\right)
$$

We can show that

$$
\begin{aligned}
& \mathbf{S}_{\omega}=\mathbf{S}_{22 \cdot 1}=\mathbf{S}_{22}-\mathbf{S}_{21} \mathbf{S}_{11}^{-1} \mathbf{S}_{12}, \\
& \mathbf{S}_{\Omega}=\mathbf{S}_{22 \cdot 1 x}=\mathbf{S}_{22 \cdot x}-\mathbf{S}_{21 \cdot x} \mathbf{S}_{11 \cdot x}^{-1} \mathbf{S}_{12 \cdot x} .
\end{aligned}
$$

The first result is obtained by using

$$
\omega=\mathcal{R}\left[\mathbf{1}_{n}\right]+\mathcal{R}\left[\left(\mathbf{I}_{n}-\mathbf{P}_{\mathbf{1}_{n}}\right) \mathbf{Y}_{1}\right] .
$$

The second result is obtained by using

$$
\begin{aligned}
\Omega & =\mathcal{R}\left[\mathbf{1}_{n}\right]+\mathcal{R}\left[\left(\tilde{\mathbf{Y}}_{1}, \tilde{\mathbf{X}}\right)\right] \\
& =\mathcal{R}\left[\mathbf{1}_{n}\right]+\mathcal{R}\left[\left(\mathbf{I}_{n}-\mathbf{P}_{0}\right) \mathbf{X}\right]+\mathcal{R}\left[\left(\mathbf{I}_{n}-\mathbf{P}_{\mathbf{X}}\right)\left(\mathbf{I}_{n}-\mathbf{P}_{0}\right) \mathbf{Y}_{1}\right],
\end{aligned}
$$

where $\tilde{\mathbf{Y}}_{1}=\left(\mathbf{I}_{n}-\mathbf{P}_{\mathbf{1}_{n}}\right) \mathbf{Y}_{1}$ and $\tilde{\mathbf{X}}=\left(\mathbf{I}_{n}-\mathbf{P}_{\mathbf{1}_{n}}\right) \mathbf{X}$.
Summarizing the above results, we have the following theorem.
Theorem 5.1In the multivariate regression model (5.1), consider to test the hypothesis $H_{2.1}$ in (5.4) against $K_{2.1}$. Then the LR criterion $\lambda$ is given by

$$
\lambda^{2 / n}=\Lambda=\frac{\left|\mathbf{S}_{22 \cdot 1.1}\right|}{\left|\mathbf{S}_{22 \cdot 1}\right|},
$$

whose null distribution is $\Lambda_{p-q}(k, n-q-k-1)$.
Note that $\mathbf{S}_{22 \cdot 1}$ can be decomposed as

$$
\mathbf{S}_{22 \cdot 1}=\mathbf{S}_{22 \cdot 1 x}+\mathbf{S}_{2 x \cdot 1} \mathbf{S}_{x \times 1}^{-1} \mathbf{S}_{x 2 \cdot 11}
$$

This decomposition is obtained by expressing $\mathbf{S}_{22 \cdot 1 x}$ in terms of $\mathbf{S}_{22 \cdot 1}, \mathbf{S}_{2 x \cdot 1}, \mathbf{S}_{x x \cdot 1}$, and $\mathbf{S}_{x 2 \cdot 1}$ by using an inverse formula

$$
\left(\begin{array}{ll}
\mathbf{H}_{11} & \mathbf{H}_{12}  \tag{5.5}\\
\mathbf{H}_{21} & \mathbf{H}_{22}
\end{array}\right)^{-1}=\left(\begin{array}{cc}
\mathbf{H}_{11}^{-1} & \mathbf{O} \\
\mathbf{O} & \mathbf{O}
\end{array}\right)+\binom{-\mathbf{H}_{11}^{-1} \mathbf{H}_{12}}{\mathbf{I}} \mathbf{H}_{22 \cdot 1}^{-1}\left(\begin{array}{lll}
-\mathbf{H}_{21} \mathbf{H}_{11}^{-1} & \mathbf{I}
\end{array}\right) .
$$

The decomposition is expressed as

$$
\begin{equation*}
\mathbf{S}_{22 \cdot 1}-\mathbf{S}_{22 \cdot 1 x}=\mathbf{S}_{2 x \cdot 1} \mathbf{S}_{x x \cdot 1}^{-1} \mathbf{S}_{x 2 \cdot 1 .} \tag{5.6}
\end{equation*}
$$

The result may be also obtained by the following algebraic method. We have

$$
\begin{aligned}
\mathbf{S}_{22 \cdot 1}-\mathbf{S}_{22 \cdot 1 x} & =\mathbf{Y}_{2}^{\prime}\left(\mathbf{P}_{\Omega}-\mathbf{P}_{\omega}\right) \mathbf{Y}_{2} \\
& =\mathbf{Y}_{2}^{\prime}\left(\mathbf{P}_{\omega \omega^{\perp} \cap \Omega}\right) \mathbf{Y}_{2},
\end{aligned}
$$

and

$$
\Omega=\mathcal{R}\left[\mathbf{1}_{n}\right]+\mathcal{R}\left[\left(\tilde{\mathbf{Y}}_{1} \tilde{\mathbf{X}}\right)\right], \quad \omega=\mathcal{R}\left[\mathbf{1}_{n}\right]+\mathcal{R}\left[\tilde{\mathbf{Y}}_{1}\right] .
$$

Therefore,

$$
\begin{aligned}
\omega^{\perp} \cap \Omega & =\mathcal{R}\left[\left(\mathbf{I}_{n}-\mathbf{P}_{1_{n}}-\mathbf{P}_{\tilde{r}_{1}}\right)\left(\tilde{\mathbf{Y}}_{1} \tilde{\mathbf{X}}\right)\right] \\
& =\mathcal{R}\left[\left(\mathbf{I}_{n}-\mathbf{P}_{1_{n}}-\mathbf{P}_{\tilde{\mathbf{Y}}_{1}}\right) \tilde{\mathbf{X}}\right],
\end{aligned}
$$

which gives an expression for $\mathbf{P}_{\omega^{\perp} \cap \Omega}$ by using Theorem 3.1 (1). This leads to (5.6).

## 6. Tests in discriminant analysis

We consider $q p$-variate normal populations with common covariance matrix $\Sigma$ and the $i$ th population having mean vector $\theta_{i}$. Suppose that a sample of size $n_{i}$ is available from the $i$ th population, and let $y_{i j}$ be the $j$ th observation from the $i$ th population. The observation matrix for all the observations is expressed as

$$
\begin{equation*}
\mathbf{Y}=\left(\boldsymbol{y}_{11}, \ldots, \boldsymbol{y}_{1 n_{1}}, \boldsymbol{y}_{21}, \ldots, \boldsymbol{y}_{q 1}, \ldots, \boldsymbol{y}_{q n_{q}}\right)^{\prime} . \tag{6.1}
\end{equation*}
$$

It is assumed that $y_{i j}$ are independent, and

$$
\begin{equation*}
\boldsymbol{y}_{i j} \sim \mathrm{~N}\left(\theta_{i}, \Sigma\right), j=1, \ldots, n_{i} ; i=1, \ldots, q, \tag{6.2}
\end{equation*}
$$

The model is expressed as

$$
\begin{equation*}
\mathbf{Y}=\mathbf{A} \Theta+\mathbf{E}, \tag{6.3}
\end{equation*}
$$

where

$$
\mathbf{A}=\left(\begin{array}{cccc}
\mathbf{1}_{n_{1}} & \mathbf{0} & \cdots & \mathbf{0} \\
\mathbf{0} & \mathbf{1}_{n_{2}} & \cdots & \mathbf{0} \\
\vdots & \vdots & \ddots & \vdots \\
\mathbf{0} & \mathbf{0} & \cdots & \mathbf{1}_{n_{q}}
\end{array}\right), \quad \Theta=\left(\begin{array}{l}
\theta_{1}^{\prime} \\
\theta_{2}^{\prime} \\
\vdots \\
\theta_{q}^{\prime}
\end{array}\right) .
$$

Here, the error matrix $\mathbf{E}$ has the same property as in (2.1).
First, we consider to test

$$
\begin{equation*}
H: \theta_{1}=\cdots=\theta_{q}(=\theta), \tag{6.4}
\end{equation*}
$$

against alternatives $K$ : $\boldsymbol{\theta}_{i} \neq \boldsymbol{\theta}_{j}$ for some $i, j$. The hypothesis can be expressed as

$$
\begin{equation*}
H: \mathbf{C} \Theta=\mathbf{O}, \quad \mathbf{C}=\left(\mathbf{I}_{q-1},-\mathbf{1}_{q-1}\right) . \tag{6.5}
\end{equation*}
$$

The tests including LRC are based on three basic statistics, the within-group SSP matrix $\mathbf{W}$, the between-group SSP matrix $\mathbf{B}$, and the total SSP matrix $\mathbf{T}$ given by

$$
\begin{align*}
\mathbf{W} & =\sum_{i=1}^{q}\left(n_{i}-1\right) \mathbf{S}_{i}, \quad \mathbf{B}=\sum_{i=1}^{q} n_{i}\left(\overline{\boldsymbol{y}}_{i}-\overline{\boldsymbol{y}}\right)\left(\overline{\boldsymbol{y}}_{i}-\overline{\boldsymbol{y}}\right)^{\prime}, \\
\mathbf{T} & =\mathbf{B}+\mathbf{W}=\sum_{i=1}^{q} \sum_{j=1}^{n_{i}}\left(\boldsymbol{y}_{i j}-\overline{\boldsymbol{y}}\right)\left(\boldsymbol{y}_{i j}-\overline{\boldsymbol{y}}\right)^{\prime}, \tag{6.6}
\end{align*}
$$

where $\overline{\mathbf{y}}_{i}$ and $\mathbf{S}_{i}$ are the mean vector and sample covariance matrix of the $i$ th population, and $\overline{\mathbf{y}}$ is the total mean vector defined by $(1 / n) \sum_{i=1}^{q} n_{i} \overline{\mathbf{y}}_{i}$ and $n=\sum_{i=1}^{q} n_{i}$. In general, $\mathbf{W}$ and $\mathbf{B}$ are independently distributed as a Wishart distribution $W_{p}(n-q, \Sigma)$ and a noncentral Wishart distribution $W_{p}(q-1, \Sigma ; \boldsymbol{\Delta})$, respectively, where

$$
\Delta=\sum_{i=1}^{q} n_{i}\left(\theta_{i}-\bar{\theta}\right)\left(\theta_{i}-\bar{\theta}\right)^{\prime},
$$

where $\bar{\theta}=(1 / n) \sum_{i=1}^{q+1} n_{i} \theta_{i}$. Then, the following theorem is well known.
Theorem 6.1 Let $\lambda=\Lambda^{n / 2}$ be the LRC for testing $H$ in (6.4). Then, $\Lambda$ is expressed as

$$
\begin{equation*}
\Lambda=\frac{|\mathbf{W}|}{|\mathbf{W}+\mathbf{B}|}=\frac{|\mathbf{W}|}{|\mathbf{T}|}, \tag{6.7}
\end{equation*}
$$

where $\mathbf{W}, \mathbf{B}$, and $\mathbf{T}$ are given in (6.6). Further, under $H$, the statistic $\Lambda$ is distributed as a lambda distribution $\Lambda_{p}(q-1, n-q)$.

Now we shall show Theorem 6.1 by an algebraic method. It is easy to see that

$$
\Omega=\mathcal{R}[\mathbf{A}], \quad \omega=\mathrm{N}\left[\mathbf{C}\left(\mathbf{A}^{\prime} \mathbf{A}\right)^{-1} \mathbf{A}^{\prime}\right] \cap \Omega=\mathcal{R}\left[\mathbf{1}_{n}\right] .
$$

The last equality is also checked from that under $H$

$$
\mathrm{E}[\mathbf{Y}]=\mathbf{A 1}_{q} \theta^{\prime}=\mathbf{1}_{n} \theta^{\prime} .
$$

We have

$$
\begin{aligned}
\mathbf{T} & =\mathbf{Y}^{\prime}\left(\mathbf{I}_{n}-\mathbf{P}_{1_{n}}\right) \mathbf{Y} \\
& =\mathbf{Y}^{\prime}\left(\mathbf{I}_{n}-\mathbf{P}_{\mathbf{A}}\right) \mathbf{Y}+\mathbf{Y}^{\prime}\left(\mathbf{P}_{\mathbf{A}}-\mathbf{P}_{1_{n}}\right) \mathbf{Y} \\
& =\mathbf{W}+\mathbf{B} .
\end{aligned}
$$

Further, it is easily checked that

1. $\left(\mathbf{I}_{n}-\mathbf{P}_{\mathrm{A}}\right)^{2}=\mathbf{I}_{n}-\mathbf{P}_{\mathrm{A}^{\prime}} \quad\left(\mathbf{P}_{\mathrm{A}^{\prime}}-\mathbf{P}_{\mathbf{1}_{n}}\right)^{2}=\mathbf{P}_{\mathrm{A}^{\prime}}-\mathbf{P}_{\mathbf{1}_{n}}$.
2. $\left(\mathbf{P}_{\mathrm{A}}-\mathrm{P}_{1_{n}}\right)\left(\mathrm{P}_{\mathrm{A}}-\mathbf{P}_{1_{n}}\right)=\mathbf{O}$.
3. $f_{e}=\operatorname{dim}\left[\mathcal{R}[\mathbf{A}]^{\perp}\right]=\operatorname{tr}\left(\mathbf{I}_{n}-\mathbf{P}_{\mathrm{A}}\right)=n-q$,

$$
f_{h}=\operatorname{dim}\left[\mathcal{R}\left[\mathbf{1}_{n}\right]^{\perp} \cap \mathcal{R}[\mathbf{A}]\right]=\operatorname{tr}\left(\mathbf{P}_{\mathrm{A}}-\mathbf{P}_{\mathbf{1}_{n}}\right)=q-1 .
$$

Related to the test of $H$, we are interested in whether a subset of variables $y_{1}, \ldots, y_{p}$ is sufficient for discriminant analysis, or the set of remainder variables has no additional information or is redundant. Without loss of generality, we consider the sufficiency of a subvector $y_{1}=\left(y_{1}, \ldots, y_{k}\right)^{\prime}$ of $\boldsymbol{y}$, or redundancy of the remainder vector $y_{2}=\left(y_{k+1}, \ldots, y_{p}\right)^{\prime}$. Consider to test

$$
\begin{equation*}
H_{2 \cdot 1}: \theta_{1 ; 2 \cdot 1}=\cdots=\theta_{q ; 2 \cdot 1}\left(=\theta_{2 \cdot 1}\right), \tag{6.8}
\end{equation*}
$$

where

$$
\theta_{i}=\binom{\theta_{i, 1}}{\theta_{i ; 2}}, \quad \theta_{i ; 1} ; k \times 1, \quad i=1, \ldots, q,
$$

and

$$
\theta_{i ; 2,1}=\theta_{i, 2}-\Sigma_{21} \Sigma_{11}^{-1} \theta_{i ; 1,}, \quad i=1, \ldots, q .
$$

The testing problem was considered by [11]. The hypothesis can be formulated in terms of Maharanobis distance and discriminant functions. For its details, see [12, 13]. To obtain a likelihood ratio for $H_{2 \cdot 1}$, we partition the observation matrix as

$$
\mathbf{Y}=\left(\begin{array}{ll}
\mathbf{Y}_{1} & \mathbf{Y}_{2}
\end{array}\right), \quad \mathbf{Y}_{1}: n \times k .
$$

Then the conditional distribution of $\mathbf{Y}_{2}$ given $\mathbf{Y}_{1}$ is normal such that the rows of $\mathbf{Y}_{2}$ are independently distributed with covariance matrix $\boldsymbol{\Sigma}_{22 \cdot 1}=\boldsymbol{\Sigma}_{22}-\boldsymbol{\Sigma}_{21} \boldsymbol{\Sigma}_{11}^{-1} \boldsymbol{\Sigma}_{12}$, and the conditional mean is given by

$$
\begin{equation*}
\mathrm{E}\left(\mathbf{Y}_{2} \mid \mathbf{Y}_{1}\right)=\mathbf{A} \Theta_{2 \cdot 1}+\mathbf{Y}_{1} \boldsymbol{\Sigma}_{11}^{-1} \boldsymbol{\Sigma}_{12}, \tag{6.9}
\end{equation*}
$$

where $\Theta_{2 \mathrm{i}}=\left(\theta_{1 ; 2 \cdot 1}, \ldots, \theta_{q ; 2 \cdot 1}\right)^{\prime}$. The LRC for $H_{2 \cdot 1}$ can be obtained by use of the conditional distribution, and following the steps (D1)-(D4) in Section 5. In fact, the spaces spanned by each column of $\mathrm{E}\left(\mathbf{Y}_{2} \mid \mathbf{Y}_{1}\right)$ are the same, and let the spaces under $K_{2 \cdot 1}$ and $H_{2 \cdot 1}$ denote by $\Omega$ and $\omega$, respectively. Then

$$
\Omega=\mathcal{R}\left[\left(\mathbf{A} \mathbf{Y}_{1}\right)\right], \quad \omega=\mathcal{R}\left[\left(\mathbf{1}_{n} \mathbf{Y}_{1}\right)\right],
$$

$\operatorname{dim}(\Omega)=q+k$, and $\operatorname{dim}(\omega)=q+1$. The likelihood ratio criterion $\lambda$ can be expressed as

$$
\lambda^{2 / n}=\Lambda=\frac{\left|\mathbf{S}_{\Omega}\right|}{\left|\mathbf{S}_{\omega}\right|}=\frac{\left|\mathbf{S}_{\Omega}\right|}{\left|\mathbf{S}_{\Omega}+\left(\mathbf{S}_{\omega}-\mathbf{S}_{\Omega}\right)\right|} .
$$

where $\mathbf{S}_{\Omega}=\mathbf{Y}_{2}^{\prime}\left(\mathbf{I}_{n}-\mathbf{P}_{\Omega}\right) \mathbf{Y}_{2}$ and $\mathbf{S}_{\omega}=\mathbf{Y}_{2}^{\prime}\left(\mathbf{I}_{n}-\mathbf{P}_{\omega}\right) \mathbf{Y}_{2}$. We express the LRC in terms of $\mathbf{W}, \mathbf{B}$, and $\mathbf{T}$. Let us partition $\mathbf{W}, \mathbf{B}$, and $\mathbf{T}$ as

$$
\mathbf{W}=\left(\begin{array}{ll}
\mathbf{W}_{11} & \mathbf{W}_{12}  \tag{6.10}\\
\mathbf{W}_{21} & \mathbf{W}_{22}
\end{array}\right), \quad \mathbf{B}=\left(\begin{array}{ll}
\mathbf{B}_{11} & \mathbf{B}_{12} \\
\mathbf{B}_{21} & \mathbf{B}_{22}
\end{array}\right), \quad \mathbf{T}=\left(\begin{array}{ll}
\mathbf{T}_{11} & \mathbf{T}_{12} \\
\mathbf{T}_{21} & \mathbf{T}_{22}
\end{array}\right),
$$

where $\mathbf{W}_{12}: q \times(p-q), \mathbf{B}_{12}: q \times(p-q)$, and $\mathbf{T}_{12}: q \times(p-q)$. Noting that $\mathbf{P}_{\Omega}=\mathbf{P}_{\mathrm{A}}+\mathbf{P}_{\left(\mathbf{I}_{n}-\mathbf{P}_{\mathrm{A}}\right)} \mathbf{Y}_{1}$, we have

$$
\begin{aligned}
\mathbf{S}_{\Omega} & =\mathbf{Y}_{2}^{\prime}\left\{\mathbf{I}_{n}-\mathbf{P}_{\mathbf{A}}-\left(\mathbf{I}_{n}-\mathbf{P}_{\mathbf{A}}\right) \mathbf{Y}_{1}\left\{\mathbf{Y}_{1}^{\prime}\left(\mathbf{I}_{n}-\mathbf{P}_{\mathbf{A}}\right) \mathbf{Y}_{1}\right\}^{-1} \mathbf{Y}_{1}\left(\mathbf{I}_{n}-\mathbf{P}_{\mathbf{A}}\right)\right\} \mathbf{Y}_{2} \\
& =\mathbf{W}_{22}-\mathbf{W}_{21} \mathbf{W}_{11}^{-1} \mathbf{W}_{12}=\mathbf{W}_{22 \cdot 1} .
\end{aligned}
$$

Similarly, noting that $\mathbf{P}_{\omega}=\mathbf{P}_{\mathbf{1}_{n}}+\mathbf{P}_{\left(\mathbf{I}_{n}-\mathbf{P}_{1}\right)} \mathbf{Y}_{1}$, we have

$$
\begin{aligned}
\mathbf{S}_{\omega} & =\mathbf{Y}_{2}^{\prime}\left\{\mathbf{I}_{n}-\mathbf{P}_{1_{n}}-\left(\mathbf{I}_{n}-\mathbf{P}_{1_{n}}\right) \mathbf{Y}_{1}\left\{\mathbf{Y}_{1}^{\prime}\left(\mathbf{I}_{n}-\mathbf{P}_{1_{n}}\right) \mathbf{Y}_{1}\right\}^{-1} \mathbf{Y}_{1}^{\prime}\left(\mathbf{I}_{n}-\mathbf{P}_{1_{n}}\right)\right\} \mathbf{Y}_{2} \\
& =\mathbf{T}_{22}-\mathbf{T}_{21} \mathbf{T}_{11}^{-1} \mathbf{T}_{12}=\mathbf{T}_{22 \cdot 1 \cdot} .
\end{aligned}
$$

Theorem 6.2Suppose that the observation matrix $\mathbf{Y}$ in (6.1) is a set of samples from $\mathrm{N}_{p}\left(\boldsymbol{\theta}_{i}, \boldsymbol{\Sigma}\right)$, $i=1, \ldots, q$. Then the likelihood ratio criterion $\lambda$ for the hypothesis $H_{2 \cdot 1}$ in (6.8) is given by

$$
\lambda=\left(\frac{\left|\mathbf{W}_{22 \cdot 1}\right|}{\left|\mathbf{T}_{22 \cdot 1}\right|}\right)^{n / 2},
$$

whereW and $\mathbf{T}$ are given by (6.6). Further, under $\mathrm{H}_{2 \cdot 1}$,

$$
\frac{\left|\mathbf{W}_{22 \cdot 1}\right|}{\left|\mathbf{T}_{22 \cdot 1}\right|}: \Lambda_{p-k}(q-1, n-q-k) .
$$

Proof. We consider the conditional distributions of $\mathbf{W}_{22 \cdot 1}$ and $\mathbf{T}_{22 \cdot 1}$ given $\mathbf{Y}_{1}$ by using Theorem 2.3, and see also that they do not depend on $\mathbf{Y}_{1}$. We have seen that

$$
\mathbf{W}_{22 \cdot 1}=\mathbf{Y}_{2}^{\prime} \mathbf{Q}_{1} \mathbf{Y}_{2}, \quad \mathbf{Q}_{1}=\mathbf{I}_{n}-\mathbf{P}_{\mathbf{A}}-\mathbf{P}_{\left(\mathbf{I}_{n}-\mathbf{P}_{\mathbf{A}}\right) \mathbf{Y}_{1}} .
$$

It is easy to see that $\mathbf{Q}_{1}^{2}=\mathbf{Q}_{1}, \operatorname{rank}\left(\mathbf{Q}_{1}\right)=\operatorname{tr} \mathbf{Q}_{1}=n-q-k, \mathbf{Q}_{1} \mathbf{A}=\mathbf{O}, \mathbf{Q}_{1} \mathbf{X}_{1}=\mathbf{O}$, and

$$
\mathrm{E}\left(\mathbf{Y}_{2} \mid \mathbf{Y}_{1}\right)^{\prime} \mathbf{Q}_{1} \mathrm{E}\left(\mathbf{Y}_{2} \mid \mathbf{Y}_{1}\right)=\mathbf{O} .
$$

This implies that $\mathbf{W}_{22 \cdot 1} \mid \mathbf{Y}_{1} \sim \mathbf{W}_{p-k}\left(n-q-k, \boldsymbol{\Sigma}_{22 \cdot 1}\right)$ and hence $\mathbf{W}_{22 \cdot 1} \sim \mathrm{~W}_{p-k}\left(n-q-k, \Sigma_{22 \cdot 1}\right)$. For $\mathrm{T}_{22 \cdot 1}$, we have

$$
\mathbf{T}_{22 \cdot 1}=\mathbf{Y}_{2}^{\prime} \mathbf{Q}_{2} \mathbf{Y}_{2}, \quad \mathbf{Q}_{2}=\mathbf{I}_{n}-\mathbf{P}_{1_{n}}-\mathbf{P}_{\left(\mathbf{I}_{n}-\mathbf{P}_{1_{n}}\right) \mathbf{Y}_{1}},
$$

and hence

$$
\mathbf{T}_{22 \cdot 1}-\mathbf{W}_{22 \cdot 1}=\mathbf{Y}_{2^{\prime}}\left(\mathbf{Q}_{2}-\mathbf{Q}_{1}\right) \mathbf{Y}_{2} .
$$

Similarly, $\mathbf{Q}_{2}$ is idempotent. Using $\mathbf{P}_{\mathbf{1}_{n}} \mathbf{P}_{\mathrm{A}}=\mathbf{P}_{\mathrm{A}} \mathbf{P}_{\mathbf{1}_{n}}=\mathbf{P}_{\mathbf{1}_{n}}$, we have $\mathbf{Q}_{1} \mathbf{Q}_{2}=\mathbf{Q}_{2} \mathbf{Q}_{1}=\mathbf{Q}_{1}$, and hence

$$
\left(\mathbf{Q}_{2}-\mathbf{Q}_{1}\right)^{2}=\mathbf{Q}_{2}-\mathbf{Q}_{1}, \quad \mathbf{Q}_{1} \cdot\left(\mathbf{Q}_{2}-\mathbf{Q}_{1}\right)=\mathbf{0} .
$$

Further, under $H_{2 \cdot 1}$,

$$
\mathrm{E}\left(\mathbf{X}_{2} \mid \mathbf{X}_{1}\right)^{\prime}\left(\mathbf{Q}_{2}-\mathbf{Q}_{1}\right) \mathrm{E}\left(\mathbf{X}_{2} \mid \mathbf{X}_{1}\right)=\mathbf{O}
$$

## 7. General multivariate linear model

In this section, we consider a general multivariate linear model as follows. Let $\mathbf{Y}$ be an $n \times p$ observation matrix whose rows are independently distributed as $p$-variate normal distribution with a common covariance matrix $\Sigma$. Suppose that the mean of $\mathbf{Y}$ is given as

$$
\begin{equation*}
\mathrm{E}(\mathbf{Y})=\mathbf{A} \Theta \mathbf{X}^{\prime} \tag{7.1}
\end{equation*}
$$

where $\mathbf{A}$ is an $n \times k$ given matrix with rank $k, \mathbf{X}$ is a $p \times q$ matrix with rank $q$, and $\Theta$ is a $k \times q$ unknown parameter matrix. For a motivation of (7.1), consider the case when a single variable $y$ is measured at $p$ time points $t_{1}, \ldots, t_{p}$ (or different conditions) on $n$ subjects chosen at random from a group. Suppose that we denote the variable $y$ at time point $t_{j}$ by $y_{j}$. Let the observations $y_{i 1}, \ldots, y_{i p}$ of the $i$ th subject be denoted by

$$
\boldsymbol{y}_{i}=\left(y_{i 1}, \ldots, y_{i p}\right)^{\prime}, \quad i=1, \ldots, n .
$$

If we consider a polynomial regression of degree $q-1$ of $y$ on the time variable $t$, then

$$
\mathrm{E}\left(\boldsymbol{y}_{i}\right)=\mathbf{X} \boldsymbol{\theta},
$$

where

$$
\mathbf{X}=\left(\begin{array}{cccc}
1 & t_{1} & \cdots & t_{1}^{q-1} \\
\vdots & \vdots & \cdots & \vdots \\
1 & t_{p} & \cdots & t_{p}^{q-1}
\end{array}\right), \quad \theta=\left(\begin{array}{c}
\theta_{1} \\
\theta_{2} \\
\vdots \\
\theta_{q}
\end{array}\right) .
$$

If there are $k$ different groups and each group has a polynomial regression of degree $q-1$ of $y$, we have a model given by (7.1). From such motivation, the model (7.1) is also called a growth curve model. For its detail, see [14].

Now, let us consider to derive LRC for a general linear hypothesis

$$
\begin{equation*}
H_{g}: \mathbf{C} \Theta \mathbf{D}=\mathbf{O}, \tag{7.2}
\end{equation*}
$$

against alternatives $K_{g}: \mathbf{C} \Theta \mathbf{D} \neq \mathbf{O}$. Here, $\mathbf{C}$ is a $c \times k$ given matrix with rank $c$, and $\mathbf{D}$ is a $q \times d$ given matrix with rank $d$. This problem was discussed by [15-17]. Here, we obtain LRC by reducing it to the problem of obtaining LRC for a general linear hypothesis in a multivariate linear model. In order to relate the model (7.1) to a multivariate linear model, consider the transformation from $\mathbf{Y}$ to $(\mathbf{U} \mathbf{V})$ :

$$
\begin{equation*}
(\mathbf{U} \mathbf{V})=\mathbf{Y} \mathbf{G}, \quad \mathbf{G}=\left(\mathbf{G}_{1} \mathbf{G}_{2}\right), \tag{7.3}
\end{equation*}
$$

where $\mathbf{G}_{1}=\mathbf{X}\left(\mathbf{X}^{\prime} \mathbf{X}\right)^{-1}, \mathbf{G}_{2}=\widetilde{\mathbf{X}}$, and $\tilde{\mathbf{X}}$ are a $p \times(p-q)$ matrix satisfying $\widetilde{\mathbf{X}}^{\prime} \mathbf{X}=\mathbf{O}$ and $\widetilde{\mathbf{X}}^{\prime} \tilde{\mathbf{X}}=\mathbf{I}_{p-q}$. Then, the rows of $(\mathbf{U} \mathbf{V})$ are independently distributed as $p$-variate normal distributions with means

$$
E[(\mathbf{U} \mathbf{V})]=(\mathbf{A} \Theta \mathbf{O})
$$

and the common covariance matrix

$$
\Psi=\mathbf{G}^{\prime} \Sigma \mathbf{G}=\left(\begin{array}{ll}
\mathbf{G}_{1}^{\prime} \Sigma \mathbf{G}_{1} & \mathbf{G}_{1}^{\prime} \Sigma \mathbf{G}_{2} \\
\mathbf{G}_{2}^{\prime} \Sigma \mathbf{G}_{1} & \mathbf{G}_{2} \Sigma \mathbf{G}_{2}
\end{array}\right)=\left(\begin{array}{ll}
\Psi_{11} & \Psi_{12} \\
\Psi_{21} & \Psi_{22}
\end{array}\right) .
$$

This transformation can be regarded as one from $y=\left(y_{1}, \ldots, y_{p}\right)^{\prime}$ to a $q$-variate main variable $u=\left(u_{1}, \ldots, u_{q}\right)^{\prime}$ and a $(p-q)$-variate auxiliary variable $\boldsymbol{v}=\left(v_{1}, \ldots, v_{p-q}\right)^{\prime}$. The model (7.1) is equivalent to the following joint model of two components:

1. The conditional distribution of $\mathbf{U}$ given $\mathbf{V}$ is

$$
\begin{equation*}
\mathbf{U} \mid \mathbf{V}: \mathrm{N}_{n \times 9}\left(\mathbf{A}^{*} \Xi, \Psi_{11 \cdot 2}\right) \tag{7.4}
\end{equation*}
$$

2. The marginal distribution of $\mathbf{V}$ is

$$
\begin{equation*}
\mathbf{V}: \mathbf{N}_{n \times(p-q)}\left(\mathbf{O}, \Psi_{22}\right), \tag{7.5}
\end{equation*}
$$

where

$$
\begin{aligned}
& \mathbf{A}^{*}=(\mathbf{A} \mathbf{V}), \quad \boldsymbol{\Xi}=\binom{\Theta}{\Gamma}, \\
& \Gamma=\Psi_{22}^{-1} \Psi_{21}, \Psi_{11 \cdot 2}=\Psi_{11}-\Psi_{12} \Psi_{22}^{-1} \Psi_{21} .
\end{aligned}
$$

Before we obtain LRC, first we consider the MLEs in (7.1). Applying a general theory of multivariate linear model to (7.4) and (7.5), the MLEs of $\Xi, \Psi_{11 \cdot 2}$, and $\Psi_{22}$ are given by

$$
\begin{equation*}
\hat{\Xi}=\left(\mathbf{A}^{*} \mathbf{A}^{*}\right)^{-1} \mathbf{A}^{*} \mathbf{U}, \quad n \Psi_{11 \cdot 2}=\mathbf{U}^{\prime}\left(\mathbf{I}_{n}-\mathbf{P}_{\mathbf{A}^{*}}\right) \mathbf{U}, \quad n \Psi_{22}=\mathbf{V}^{\prime} \mathbf{V} \tag{7.6}
\end{equation*}
$$

Let

$$
\mathbf{S}=\mathbf{Y}^{\prime}\left(\mathbf{I}_{n}-\mathbf{P}_{\mathbf{A}}\right) \mathbf{Y}, \quad \mathbf{W}=\mathbf{G}^{\prime} \mathbf{S} \mathbf{G}=(\mathbf{U} \mathbf{V})^{\prime}\left(\mathbf{I}_{n}-\mathbf{P}_{\mathrm{A}}\right)(\mathbf{U} \mathbf{~ V}),
$$

and partition $\mathbf{W}$ as

$$
\mathbf{W}=\left(\begin{array}{ll}
\mathbf{W}_{11} & \mathbf{W}_{12} \\
\mathbf{W}_{21} & \mathbf{W}_{22}
\end{array}\right), \quad \mathbf{W}_{12}: q \times(p-q) .
$$

Theorem 7.1For ann $\times$ pobservation matrix $\mathbf{Y}$, assume a general multivariate linear model given by (7.1). Then:

1. The MLE $\hat{\boldsymbol{\Theta}}$ of $\boldsymbol{\Theta}$ is given by

$$
\hat{\Theta}=\mathbf{A}\left(\mathbf{A}^{\prime} \mathbf{A}\right)^{-1} \mathbf{A}^{\prime} \mathbf{Y} \mathbf{S}^{-1} \mathbf{X}\left(\mathbf{X}^{\prime} \mathbf{S}^{-1} \mathbf{X}\right)^{-1}
$$

2. The MLE $\hat{\Psi}_{11 \cdot 2}$ of $\Psi_{11 \cdot 2}$ is given by

$$
n \Psi_{11 \cdot 2}=\mathbf{W}_{11 \cdot 2}=\left(\mathbf{X}^{\prime} \mathbf{S}^{-1} \mathbf{X}\right)^{-1} .
$$

Proof. The MLE of $\Xi$ is $\hat{\Xi}=\left(\mathbf{A}^{*} \mathbf{A}^{*}\right)^{-1} \mathbf{A}^{*} \mathbf{U}$. The inverse formula (see (5.5)) gives

$$
\begin{aligned}
\mathbf{Q} & =\left(\mathbf{A}^{*} \mathbf{A}^{*}\right)^{-1}=\left(\begin{array}{cc}
\left(\mathbf{A}^{\prime} \mathbf{A}\right)^{-1} & \mathbf{0} \\
\mathbf{O} & \mathbf{0}
\end{array}\right)+\binom{-\left(\mathbf{A}^{\prime} \mathbf{A}\right)^{-1} \mathbf{A}^{\prime} \mathbf{V}}{\mathbf{I}_{p-q}}\left[\mathbf{V}^{\prime}\left(\mathbf{I}_{n}-\mathbf{P}_{\mathbf{A}}\right) \mathbf{V}\right]^{-1}\binom{-\left(\mathbf{A}^{\prime} \mathbf{A}\right)^{-1} \mathbf{A}^{\prime} \mathbf{V}}{\mathbf{I}_{p-q}} \\
& =\left(\begin{array}{ll}
\mathbf{Q}_{11} & \mathbf{Q}_{12} \\
\mathbf{Q}_{21} & \mathbf{Q}_{22}
\end{array}\right) .
\end{aligned}
$$

Therefore, we have

$$
\begin{aligned}
\hat{\boldsymbol{\Theta}} & =\left(\mathbf{Q}_{11} \mathbf{A}^{\prime}+\mathbf{Q}_{12} \mathbf{V}^{\prime}\right) \mathbf{U} \\
& =\left(\mathbf{A}^{\prime} \mathbf{A}\right)^{-1} \mathbf{A}^{\prime} \mathbf{Y G}_{1}-\left(\mathbf{A}^{\prime} \mathbf{A}\right)^{-1} \mathbf{A}^{\prime} \mathbf{Y G}_{2}\left(\mathbf{G}_{2}^{\prime} \mathbf{S G}_{2}\right)^{-1} \mathbf{G}_{2}^{\prime} S G_{1} .
\end{aligned}
$$

Using

$$
\mathbf{G}_{2}\left(\mathbf{G}_{2}^{\prime} \mathbf{S G}_{2}\right)^{-1} \mathbf{G}_{2}^{\prime}=\mathbf{S}^{-1}-\mathbf{G}_{1}\left(\mathbf{G}_{1}^{\prime} \mathbf{S}^{-1} \mathbf{G}_{2}\right)^{-1} \mathbf{G}_{1}^{\prime} \mathbf{S}^{-1},
$$

we obtain 1. For a derivation of 2, let $\mathbf{B}=\left(\mathbf{I}_{n}-\mathbf{P}_{\mathrm{A}}\right) \mathbf{V}$. Then, using $\mathbf{P}_{\mathbf{A}^{*}}=\mathbf{P}_{\mathbf{A}}+\mathbf{P}_{\mathbf{B}}$, the first expression of (1) is obtained. Similarly, the second expression of (2) is obtained.

Theorem 7.2Let $=\Lambda^{n / 2}$ be the LRC for testing the hypothesis (7.2) in the generalized multivariate linear model (7.1). Then,

$$
\Lambda=\left|\mathbf{S}_{e}\right| /\left|\mathbf{S}_{e}+\mathbf{S}_{h}\right|,
$$

where

$$
\mathbf{S}_{e}=\mathbf{D}^{\prime}\left(\mathbf{X}^{\prime} \mathbf{S}^{-1} \mathbf{X}\right)^{-1} \mathbf{D}, \quad \mathbf{S}_{h}=(\mathbf{C} \hat{\Theta} \mathbf{D})\left(\mathbf{C R} \mathbf{C}^{\prime}\right)^{-1} \mathbf{C} \hat{\Theta} \mathbf{D}
$$

and

$$
\begin{aligned}
\mathbf{R}= & \left(\mathbf{A}^{\prime} \mathbf{A}\right)^{-1}+\left(\mathbf{A}^{\prime} \mathbf{A}\right)^{-1} \mathbf{A}^{\prime} \mathbf{Y} \mathbf{S}^{-1}\left\{\mathbf{S}-\mathbf{X}^{\prime}\left(\mathbf{X} \mathbf{S}^{-1} \mathbf{X}^{\prime}\right)^{-1} \mathbf{X}\right\} \\
& \times \mathbf{S}^{-1} \mathbf{Y} \mathbf{} \mathbf{Y}^{\prime} \mathbf{A}\left(\mathbf{A}^{\prime} \mathbf{A}\right)^{-1} .
\end{aligned}
$$

Here $\hat{\boldsymbol{O}}$ is given in Theorem 7.1.1. Further, the null distribution is $\Lambda_{d}(c, n-k-(p-q))$.
Proof. The test of $H_{g}$ in (7.2) against alternatives $K_{g}$ is equivalent to testing

$$
\begin{equation*}
H_{g}: \mathrm{C}^{*} \Xi \mathbf{D}=\mathbf{0} \tag{7.7}
\end{equation*}
$$

under the conditional model (7.4), where $\mathbf{C}^{*}=(\mathbf{C} \mathbf{O})$. Since the distribution of $\mathbf{V}$ does not depend on $H_{g^{\prime}}$ the LR test under the conditional model is the LR test under the unconditional model. Using a general result for a general linear hypothesis given in Theorem 4.1, we obtain

$$
\Lambda=\left|\tilde{\mathbf{S}}_{e}\right| /\left|\tilde{\mathbf{S}}_{e}+\tilde{\mathbf{S}}_{h}\right|,
$$

where

$$
\begin{aligned}
\tilde{\mathbf{S}}_{e} & =\mathbf{D}^{\prime} \mathbf{U}\left(\mathbf{I}_{n}-\mathbf{A}^{*}\left(\mathbf{A}^{* \prime} \mathbf{A}^{*}\right)^{-1} \mathbf{A}^{* \prime}\right) \mathbf{U D}, \\
\tilde{\mathbf{S}}_{h} & =(\mathbf{C} \hat{\boldsymbol{E}} \mathbf{D})\left(\mathbf{C}^{*}\left(\mathbf{A}^{* \prime} \mathbf{A}^{*}\right)^{-1} \mathbf{C}^{*}\right)^{-1} \mathbf{C} \hat{\Xi} D .
\end{aligned}
$$

By reduction similar to those of MLEs, it is seen that $\widetilde{\mathbf{S}}_{e}=\mathbf{S}_{e}$ and $\widetilde{\mathbf{S}}_{h}=\mathbf{S}_{h}$. This completes the proof.

## 8. Concluding remarks

In this chapter, we discuss LRC in multivariate linear model, focusing on the role of projection matrices. Testing problems considered involve the hypotheses on selection of variables or no additional information of a set of variables, in addition to a typical linear hypothesis. It may be noted that various LRCs and their distributions are obtained by algebraic methods.

We have not discussed with LRCs for the hypothesis of selection of variables in canonical correlation analysis, and for dimensionality in multivariate linear model. Some results for these problems can be found in $[3,18]$.

In multivariate analysis, there are some other test criteria such as Lawley-Hotelling trace criterion and Bartlett-Nanda-Pillai trace criterion. For the testing problems treated in this chapter, it is possible to propose such criteria as in [12].

The LRCs for tests of no additional information of a set of variables will be useful in selection of variables. For example, it is possible to propose model selection criteria such as AIC (see [19]).

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