Math tools

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Lecturer: Typed by: Zhuohua Shen

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4 CONTENTS

Preliminary

Statistical inference fundamentals

References: most of the contents are from the undergraduate course STA3020 (by Prof. Jianfeng Mao in 2022-2023 T1, and Prof. Jiasheng Shi in 2023-2024 T2) and postgraduate course STAT5010 (by Kin Wai Keith Chan in 2024-2025 T1), with main textbook Casella and Berger [1]

2.1 Statistical Models

See Chapter 3 of [1]. Suppose $X_i \sim_{\text{iid}} \mathbb{P}_*$, where \mathbb{P}_* refers to the unknown data generating process (DGPg), we find $\widehat{\mathbb{P}} \approx \mathbb{P}_*$. A statistical model is a set of distributions $\mathscr{F} = \{\mathbb{P}_{\theta} : \theta \in \Theta\}$, where Θ is the parameter space. A parametric model is the model with $\dim(\Theta) < \infty$, while a nonparametric model satisfies $\dim(\Theta) = \infty$.

Definition 2.1.1 (Exponential family). A k-dimensional exponential family (EF) $\mathscr{F} = \{f_{\theta} : \theta \in \Theta\}$ is a model consisting of pdfs of the form

$$f_{\theta}(x) = c(\theta)h(x) \exp\left\{ \sum_{j=1}^{k} \eta_{j}(\theta)T_{j}(x) \right\}$$
(2.1)

where $c(\theta), h(x) \geq 0$, $\Theta = \{\theta : c(\theta) \geq 0, \eta_j(\theta) \text{ being well defined for } 1 \leq j \leq k\}$. Let $\eta_j = \eta_j(\theta)$, the canonical form is

$$f_{\eta}(x) = b(\eta)h(x) \exp\left\{\sum_{j=1}^{k} \eta_j T_j(x)\right\}, \qquad (2.2)$$

- k-dim natural exponential family (NEF): $\mathscr{F}' = \{f_{\eta} : \eta \in \Xi\};$
- natural parameter $\eta = (\eta_1, \dots, \eta_k)^{\top}$;
- natural parameter space: $\Xi = \{ \eta \in \mathbb{R}^k : 0 < b(\eta) < \infty \};$
- the NEF \mathscr{F}' is of full rank if Ξ contains an open set in \mathbb{R}^k ;
- the EF is a curved exponential family if $p = \dim(\Theta) < k$.

Properties of EF:

- Let $X \sim f_{\eta}$, where $\eta \in \Xi$ such that (i) f_{η} is of the form (2.2) with $B(\eta) = -\log b(\eta)$, and (ii) Ξ contains an open set in \mathbb{R}^k . Then, for $j, j' = 1, \ldots, k$, $\mathbb{E}\{T_j(X)\} = \partial B(\eta)/\partial \eta_j$ and $\mathbb{C}\text{ov}\{T_j(X), T_{j'}(X)\} = \partial^2 B(\eta)/(\partial \eta_j \partial \eta_{j'})$.
- Stein's identity:

Definition 2.1.2 (Location-scale family). Let f be a density.

- A location-scale family is given by $\mathscr{F} = \{f_{\mu,\sigma} : \mu \in \mathbb{R}, \sigma \in \mathbb{R}^{++}\}$, where $f_{\mu,\sigma}(x) = f((x-\mu)/\sigma)/\sigma$.
- location parameter: μ ; scale parameter: σ ; standard density: f;
- A location family is $\mathscr{F} = \{f_{\mu,1} : \mu \in \mathbb{R}\}.$
- A scale family is $\mathscr{F} = \{f_{0,\sigma} : \sigma \in \mathbb{R}^{++}\}\$

Representation: $X = \mu + \sigma Z$, $Z \sim f_{0,1}(\cdot)$.

- $\bullet\,$ See some examples in Example 3.9, Keith's note 3, and Table 1 in Shi's note L1.
- Transform between location parameter and scale parameter by taking log.

Definition 2.1.3 (Identifiable family). If $\forall \theta_1, \theta_2 \in \Theta$ that

$$\theta_1 \neq \theta_2 \quad \Rightarrow \quad f_{\theta_1}(\cdot) \neq f_{\theta_2}(\cdot),$$

then \mathscr{F} is said to be an identifiable family, or equivalently $\theta \in \Theta$ is identifiable.

- 8 A typical feature of non-identifiable EF is that $GHAPT(G) \stackrel{?}{>} k$ FTATISTICAL INFERENCE FUNDAMENTALS
 - p < k, curved (must).
 - p = k, of full rank.
 - p > k, non-identifiable.

2.2 Principles of Data Reduction

Statistics: $T = T(X_{1:n})$, a function of $X_{1:n}$ and free of any unknown parameter.

2.2.1 Sufficiency Principle

Sufficiency principle: If $T = T(X_{1:n})$ is a "sufficient statistics" for θ , then any inference on θ will depend on $X_{1:n}$ only through T.

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Definition 2.2.1 (Sufficient, minimal sufficient, ancillary, and complete statistics). Suppose X_{1:n} \sim_{iid} \mathbb{P}_{\theta}, where \theta \in \Theta. Let T = T(X_{1:n}) be a statistic. Then T is sufficient (SS) for \theta \Leftrightarrow (def) [X_{1:n} \mid T = t] is free of \theta for each t. \Leftrightarrow (technical lemma) T(x_{1:n}) = T(x'_{1:n}) implies that f_{\theta}(x_{1:n})/f_{\theta}(x'_{1:n}) is free of \theta. \Leftrightarrow (Neyman-Fisher factorization theorem) \forall \theta \in \Theta, x_{1:n} \in \mathcal{X}^n, f_{\theta}(x_{1:n}) = A(t,\theta)B(x_{1:n}). \Leftrightarrow Define \Lambda(\theta',\theta'' \mid x_{1:n}) := f_{\theta'}(x_{1:n})/f_{\theta''}(x_{1:n}). \forall \theta',\theta'' \in \Theta, \exists function C_{\theta',\theta''} such that \Lambda(\theta',\theta'' \mid x_{1:n}) = C_{\theta',\theta''}(t), for all x_{1:n} \in \mathcal{X}^n where t = T(x_{1:n}). T is minimal sufficient (MSS) for \theta \Leftrightarrow (def) (1) T is a SS for \theta; (2) T = g(S) for any other SS S. \Leftrightarrow (1) T is a SS for \theta; (2) S(x_{1:n}) = S(x'_{1:n}) implies T(x_{1:n}) = T(x'_{1:n}) for any SS S. \Leftrightarrow (Lehmann-Scheffé theorem) \forall x_{1:n}, x'_{1:n} \in \mathcal{X}^n, f_{\theta}(x_{1:n})/f_{\theta}(x'_{1:n}) is free of \theta \Leftrightarrow T(x_{1:n}) = T(x'_{1:n}). A = A(X_{1:n}) is ancillary (ANS) if the distribution of A does not depend on \theta. T is complete (CS) if \forall \theta \in \Theta, \mathbb{E}_{\theta}g(T) = 0 implies \forall \theta \in \Theta, \mathbb{P}_{\theta}\{g(T) = 0\} = 1.
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Properties

- (Transformation) If T = r(T'), then (i) T is $SS \Rightarrow T'$ is SS; (ii) T' is $CS \Rightarrow T$ is CS; (iii) r is one-to-one, then if one is SS/MSS/CS, then the another is.
- (Basu's Lemma) $X_i \sim_{iid} \mathbb{P}_{\theta}$, A is ANS and T s CSS, then $A \perp \!\!\! \perp T$.
- (Bahadur's theorem) $X_i \sim_{iid} \mathbb{P}_{\theta}$, if an MSS exists, then any CSS is also an MSS.
 - Then if a CSS exists, then any MSS is also a CSS \Rightarrow CSS=MSS.
 - All or nothing: start with MSS T, check whether T is CS. (i) Yes, it is both CSS and MSS, then the set of MSS=CSS; (ii) No, there is no CSS at all.
- (Exp-family) If $X_i \sim_{\text{iid}} f_{\eta}$ in (2.2), then $T = (\sum_{i=1}^n T_1(X_i), \dots, \sum_{i=1}^n T_k(X_i))$ is a SS, called natural sufficient statistic. If Ξ contains an open set in \mathbb{R}^k (i.e., \mathscr{F}' is of full rank), then T is MSS and CSS.

Proof techniques

- Prove T is not sufficient for θ : show if $\exists x_{1_n}, x'_{1:n} \in \mathcal{X}^n$ and $\theta', \theta'' \in \Theta$, such that $T(x_{1:n}) = T(x'_{1:n})$ and $\Lambda(\theta', \theta'' \mid x_{1:n}) \neq \Lambda(\theta', \theta'' \mid x'_{1:n})$.
- Prove A is an ANS: consider location-scale representation.
- Prove T is a CS: use definition or take $d\mathbb{E}_{\theta}g(T)/d\theta = 0$.
- Disprove T is CS:
 - Construct an ANS S(T) based on T, then $\mathbb{E}S(T)$ is free of θ , then $g(T) = S(T) \mathbb{E}S(T)$ is free of θ but $g(T) \neq 0$ w.p.1.
 - (Cancel the 1st moment) Find two unbiased estiamtors for θ as a function of T. E.g., $X_1, X_2 \sim_{\text{iid}} N(\theta, \theta^2)$, $T = (X_1, X_2), g(T) = X_1 X_2 \sim N(0, 2\theta^2)$.

Remark 2.2.2. • ANS A is useless on its own, but useful together with other information.

• $\mathbb{P}(A(X) \mid \theta)$ is free of θ , but for non-SS T, $\mathbb{P}(A(X) \mid T(X))$ is not necessarily free of θ .

2.2.2 Likelihood principle

Multivariate Inference Fundamentals

Reference:

- Robb J. Muirhead Aspects of multivariate statistical theory [5].
- CUHK STAT4002 Applied Multivariate Analysis (2023 Spring), by Zhixiang Lin.
- CUHK STAT5030 Linear Models (2025 Spring), by Yuanyuan Lin.
- Peng DING Linear Model and Extensions.
- Ronald Christensen Plane answers to complex questions: the theory of linear models [2].

3.1 Random vectors and distributions

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Definition 3.1.1. Let \boldsymbol{x} = (x_1, \dots, x_p)^{\top} \in \mathbb{R}^p be a random vector,

• Mean \mathbb{E}\boldsymbol{x} = \boldsymbol{\mu} = (\mathbb{E}x_1, \dots, \mathbb{E}x_p)^{\top} = (\mu_j).

• Covariance matrix \mathbb{V}\operatorname{ar}(\boldsymbol{x}) = \mathbb{C}\operatorname{ov}(\boldsymbol{x}) = \Sigma = \mathbb{E}[(\boldsymbol{x} - \mathbb{E}\boldsymbol{x})(\boldsymbol{x} - \mathbb{E}\boldsymbol{x})^{\top}] = \mathbb{E}\boldsymbol{x}\boldsymbol{x}^{\top} - \mathbb{E}\boldsymbol{x}\mathbb{E}\boldsymbol{x}^{\top} = (\sigma_{ij}), \ \Sigma \succeq \boldsymbol{0}.

• Correlation matrix R = D^{-1/2}\Sigma D^{-1/2}, where D = \operatorname{diag}(\sigma_{11}, \dots, \sigma_{pp}). We have R_{ij} = \rho_{ij} = \sigma_{ij}/(\sqrt{\sigma_{ii}}\sqrt{\sigma_{jj}}).

• If \boldsymbol{y} \in \mathbb{R}^q random vector, then \mathbb{C}\operatorname{ov}(\boldsymbol{x}, \boldsymbol{y}) = \mathbb{E}[(\boldsymbol{x} - \mathbb{E}\boldsymbol{x})(\boldsymbol{y} - \mathbb{E}\boldsymbol{y})^{\top}] = \mathbb{E}\boldsymbol{x}\boldsymbol{y}^{\top} - \mathbb{E}\boldsymbol{x}\mathbb{E}\boldsymbol{y}^{\top} \in \mathbb{R}^{p \times q}.

If \boldsymbol{Z} = (z_{ij}) \in \mathbb{R}^{p \times q} is a random matrix,

• \mathbb{E}\boldsymbol{Z} = (\mathbb{E}z_{ij}).
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Proposition 3.1.2. Let \mathbf{x} \in \mathbb{R}^p be a random vector, \mathbf{a}, \mathbf{b} \in \mathbb{R}^p be vectors, A \in \mathbb{R}^{r_1 \times p}, B \in \mathbb{R}^{r_2 \times p} be matrices,

• \mathbb{E}\mathbf{a}^{\top}\mathbf{x} = \mathbf{a}^{\top}\mathbb{E}\mathbf{x}, \mathbb{V}\operatorname{ar}(\mathbf{a}^{\top}\mathbf{x}) = \mathbf{a}^{\top}\Sigma\mathbf{a}, and \mathbb{C}\operatorname{ov}(\mathbf{a}^{\top}\mathbf{x}, \mathbf{b}^{\top}\mathbf{x}) = \mathbf{a}^{\top}\Sigma\mathbf{b}.

• \mathbb{E}A\mathbf{x} = A\mathbb{E}\mathbf{x}, \mathbb{V}\operatorname{ar}(A\mathbf{x}) = A\Sigma A^{\top}, and \mathbb{C}\operatorname{ov}(A\mathbf{x}, B\mathbf{x}) = A\Sigma B^{\top}.

• If \mathbf{y} = A\mathbf{x} + \mathbf{b}, where A \in \mathbb{R}^{q \times p}, \mathbf{b} \in \mathbb{R}^q, then \mathbf{\mu_y} = A\mathbf{\mu_x} + \mathbf{b} and \Sigma_{\mathbf{y}} = A\Sigma_{\mathbf{x}}A^{\top}.

• \mathbb{E}(\mathbf{x}^{\top}A\mathbf{x}) = \operatorname{tr}(A\Sigma) + \mathbf{\mu}^{\top}A\mathbf{\mu}.

Let \mathbf{Z} \in \mathbb{R}^{p \times q} be a random matrix, B \in \mathbb{R}^{m \times p}, C \in \mathbb{R}^{q \times n}, and D \in \mathbb{R}^{m \times n} constants, then

• \mathbb{E}(B\mathbf{Z}C + D) = B\mathbb{E}(\mathbf{Z})C + D.
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- The $\Sigma \in \mathbb{R}^{p \times p}$ is a covariance matrix (i.e., $\Sigma = \mathbb{C}\text{ov}(\boldsymbol{x})$ for some random vector $\boldsymbol{x} \in \mathbb{R}^p$) iff $\Sigma \succeq \mathbf{0}$.
 - $-(\Leftarrow)$: suppose $\mathbf{r}(\Sigma) = r \leq p$, write full rank decomposition $\Sigma = CC^{\top}$, $C \in \mathbb{R}^{p \times r}$. Let $\mathbf{y} \sim [\mathbf{0}_r, I_r]$, then $\mathbb{C}\mathrm{ov}(C\mathbf{y}) = \Sigma$.
- If Σ is not PD, then $\exists \boldsymbol{a} \neq \boldsymbol{0}_p$ s.t. $\mathbb{V}\operatorname{ar}(\boldsymbol{a}^{\top}\boldsymbol{x}) = 0$ so w.p.1., $\boldsymbol{a}^{\top}\boldsymbol{x} = k$, i.e., \boldsymbol{x} lies in a hyperplane.

Theorem 3.1.3. If $\mathbf{x} \in \mathbb{R}^p$ random, then its distribution is uniquely determined by the distributions of $\mathbf{a}^{\top}\mathbf{x}$, $\forall \mathbf{a} \in \mathbb{R}^p$.

The proof uses the fact that a distribution in \mathbb{R}^p is uniquely determined by its ch.f., see Theorem 1.2.2. [5].

Definition 3.1.4. Dataset contains p variables and n observations are represented by $\boldsymbol{X} = (\boldsymbol{x}_1, \dots, \boldsymbol{x}_n)^\top$, where the ith row $\boldsymbol{x}_i^\top = (x_{i1}, \dots, x_{ip})$ is the ith observation vector, $i = 1, \dots, n$.

- (Sample mean vector) $\bar{x} = n^{-1} \sum_{i=1}^{n} x_i = (\bar{x}_1, \dots, \bar{x}_p)^{\top}$, where $\bar{x}_j = n^{-1} \sum_{i=1}^{n} x_{ij}$.
- (Sum of squares and cross product (SSCP) matrix) $A = \sum_{i=1}^{n} (x_i \bar{x})(x_i \bar{x})^{\top}$.
- (Sample covariance matrix) $S = (n-1)^{-1}A$.
- (Sample correlation matrix) $\mathbf{R} = D^{-1/2} \mathbf{S} D^{-1/2}$, where $D^{-1/2} = \operatorname{diag}(1/\sqrt{s_{11}}, \dots, 1/\sqrt{s_{pp}})$.

10 • $\bar{x} = n^{-1} X^{\top} \mathbf{1}_n$, and

$$egin{aligned} oldsymbol{A} &= \sum_{i=1}^n (oldsymbol{x}_i - oldsymbol{\mu}) (oldsymbol{x}_i - oldsymbol{\mu})^ op - n(ar{oldsymbol{x}} - oldsymbol{\mu}) (ar{oldsymbol{x}} - oldsymbol{\mu})^ op \ &= (oldsymbol{X} - oldsymbol{1}_n ar{oldsymbol{x}}^ op)^ op (oldsymbol{X} - oldsymbol{1}_n ar{oldsymbol{x}}^ op)^ op (oldsymbol{0}. \end{aligned}$$

• $\mathbb{E}\bar{x} = \mu$, $\mathbb{V}\operatorname{ar}(\bar{x}) = n^{-1}\Sigma$, $\mathbb{E}A = (n-1)\Sigma$, and $\mathbb{E}S = \Sigma$.

3.1.1 Multivariate normal distribution

Definition 3.1.5 (Original definition of multivariate normal). The random vector $x \in \mathbb{R}^p$ is said to have an p-variate normal distribution $(\boldsymbol{x} \sim N_p)$ if $\forall \boldsymbol{a} \in \mathbb{R}^p$, the distribution of $\boldsymbol{a}^\top \boldsymbol{x}$ is univariate normal.

Theorem 3.1.6 (Fundamental properties). Let $x \sim N_p$, we have

- 1. Both $\mu = \mathbb{E}x$ and $\Sigma = \mathbb{C}ov(x)$ exist and the distribution of x is determined by μ and Σ . Write $x \sim N_p(\mu, \Sigma)$.
- 2. (Representation) Let $\Sigma \succeq \mathbf{0}_{p \times p}$, $\mathbf{r}(\Sigma) = r \leq p$, and $u_{1:r} \sim_{\text{iid}} \mathbf{N}(0,1)$, i.e., $\mathbf{u} \sim \mathbf{N}_r(\mathbf{0}_r, I_r)$, then if C is the full rank decomposition of Σ and $\mu \in \mathbb{R}^p$, then $x = Cu + \mu \sim N_p(\mu, \Sigma)$.
 - Let $\Sigma = HDH^{\top}$ be the spectral decomposition, then $\boldsymbol{x} = HD^{1/2}\boldsymbol{z} + \boldsymbol{\mu}$, where $\boldsymbol{z} \sim N_p(\boldsymbol{0}_p, I_p)$.
- 3. If $\mathbf{x} \sim N_p(\boldsymbol{\mu}, \Sigma)$, then its ch.f. $\phi_{\mathbf{x}}(\mathbf{t}) = \exp(i\boldsymbol{\mu}^{\top} \mathbf{t} \mathbf{t}^{\top} \Sigma \mathbf{t}/2)$.
- 4. (Density) $x \sim N_p(\mu, \Sigma)$ with $\Sigma \succ 0$, then x has pdf

$$f(\mathbf{x}) = \frac{1}{(2\pi)^{p/2} |\Sigma|^{1/2}} \exp\left\{-\frac{1}{2} (\mathbf{x} - \boldsymbol{\mu})^T \Sigma^{-1} (\mathbf{x} - \boldsymbol{\mu})\right\}. \tag{3.1}$$

Note that we guarantee the existence of $N_p(\boldsymbol{\mu}, \boldsymbol{\Sigma})$ by means of the representation in point 2. By its density, we have MVN kernel: If

$$f(\boldsymbol{x}) \propto \exp\left\{-\frac{1}{2}(\boldsymbol{x}^{\top}A\boldsymbol{x} - 2\boldsymbol{x}^{\top}B)\right\} = \exp\left\{-\frac{1}{2}(\boldsymbol{x} - A^{-1}B)^{\top}A(\boldsymbol{x} - A^{-1}B) - B^{\top}A^{-1}B\right\},$$

then $x \sim N_p(A^{-1}B, A^{-1})$.

Theorem 3.1.7 (Properties of multivariate normal). If $x \sim N_p(\mu, \Sigma)$, then we have

- 1. (Linearity) Let $B \in \mathbb{R}^{q \times p}$, $\mathbf{b} \in \mathbb{R}^q$ nonrandom, and $B \Sigma B^{\top} \succ \mathbf{0}$, then $B \mathbf{x} + \mathbf{b} \sim \mathrm{N}_q (B \boldsymbol{\mu} + \mathbf{b}, B \Sigma B^{\top})$.
- 2. (Linear combinations) If $\mathbf{x}_k \sim \mathrm{N}_p(\boldsymbol{\mu}_k, \Sigma_k) \perp \text{for } k = 1, \ldots, N$, then for any fixed constants $\alpha_1, \ldots, \alpha_N$, $\sum_{k=1}^N \alpha_k \mathbf{x}_k \sim \mathrm{N}_p(\sum_{k=1}^N \alpha_k \boldsymbol{\mu}_k, \sum_{k=1}^N \alpha_k^2 \Sigma_k)$.

 The sample mean $\bar{\mathbf{x}} \sim \mathrm{N}_p(\boldsymbol{\mu}, \Sigma/N)$.
- 3. (Subset) The marginal distribution of any subset of k < p components of x is k-variate normal.
- 4. (Marginal distribution) Partition

$$m{x} = \left[egin{array}{c} m{x}_1 \ m{x}_2 \end{array}
ight], \quad m{\mu} = \left[egin{array}{c} m{\mu}_1 \ m{\mu}_2 \end{array}
ight], \quad m{\Sigma} = \left[egin{array}{cc} \Sigma_{11} & \Sigma_{12} \ \Sigma_{21} & \Sigma_{22} \end{array}
ight], \quad m{x}_1 \in \mathbb{R}^q, m{x}_2 \in \mathbb{R}^{p-q}, \Sigma_{12} \in \mathbb{R}^{q imes (p-q)}.$$

Then $\boldsymbol{x}_1 \sim N_q(\boldsymbol{\mu}_1, \Sigma_{11}), \ \boldsymbol{x}_1 \perp \boldsymbol{x}_2 \ \text{iff } \Sigma_{12} = \boldsymbol{0}$

- 5. (Conditional distribution) Let Σ_{22}^- be a generalized inverse of Σ_{22} (i.e., $\Sigma_{22}\Sigma_{22}^-\Sigma_{22} = \Sigma_{22}$), then (a) $\mathbf{x}_1 \Sigma_{12}\Sigma_{22}^-\mathbf{x}_2 \sim \mathrm{N}_q(\boldsymbol{\mu}_1 \Sigma_{12}\Sigma_{22}^-\boldsymbol{\mu}_2, \Sigma_{11} \Sigma_{12}\Sigma_{22}^-\Sigma_{21})$, and $\perp \mathbf{x}_2$.
 - (b) $[\mathbf{x}_1 \mid \mathbf{x}_2] \sim N_q(\boldsymbol{\mu}_1 + \Sigma_{12}\Sigma_{22}^-(\mathbf{x}_2 \boldsymbol{\mu}_2), \Sigma_{11} \Sigma_{12}\Sigma_{22}^-\Sigma_{21}).$
- 6. (Cramér) If $p \times 1$ random vectors $\mathbf{x} \perp \mathbf{y}$ and $\mathbf{x} + \mathbf{y} \sim N_p$, then both $\mathbf{x}, \mathbf{y} \sim N_p$.
- 7. (MLE) of (μ, Σ) is $(\bar{x}, A/n)$.
- 8. (Inverse of Σ and conditional independence) Denote $\Sigma^{-1} = (\nu^{jk})_{1 \leq i,k \leq p}$. Then $\forall j \neq k, \ v^{jk} = 0 \Leftrightarrow x_i \perp \!\!\! \perp$ $x_k \mid \boldsymbol{x} \setminus \{x_j, x_k\}.$

For point 3, each component of a random vector is (marginally) normal does not imply that the vector has a multivariate normal distribution. Counterexample: let $U_1, U_2, U_3 \sim_{iid} N(0, 1), Z \perp \!\!\!\perp U_{1:3}$. Define

$$X_1 = \frac{U_1 + ZU_3}{\sqrt{1 + Z^2}}, \quad X_2 = \frac{U_2 + ZU_3}{\sqrt{1 + Z^2}}.$$

Then $[X_1|Z] \sim N(0,1)$, free of Z, so $X_1 \sim N(0,1)$, and $X_2 \sim N(0,1)$. But (X_1,X_2) not normal. The converse is true if the components of x are all independent and normal, or if x consists of independent subvectors, each of which is normally distributed.

For the proof of point 5, we use the lemma: if $\Sigma \succeq \mathbf{0}$, then $\ker(\Sigma_{22}) \subset \ker(\Sigma_{12})$, and $\operatorname{range}(\Sigma_{21}) \subset \operatorname{range}(\Sigma_{22})$. So $\exists B \in \mathbb{R}^{q \times (p-q)} \text{ satisfying } \Sigma_{12} = B\Sigma_{22}.$

Proposition 3.1.8 (MGF). If $x \sim N(\mu, \Sigma)$, A is symmetric and σ IS non-singular, then

$$M_{\boldsymbol{x}^{\top}\boldsymbol{A}\boldsymbol{x}}(t) = |I - 2t\boldsymbol{A}\boldsymbol{\Sigma}|^{-\frac{1}{2}} \exp\left\{-\frac{1}{2}\boldsymbol{\mu}^{\top}[I - (I - 2t\boldsymbol{A}\boldsymbol{\Sigma})^{-1}]\boldsymbol{\Sigma}^{-1}\boldsymbol{\mu}\right\}.$$

Proposition 3.1.9 (Variance). If $x \sim N(\mu, \Sigma)$, then

$$Var(\boldsymbol{x}^{\top} A \boldsymbol{x}) = 2tr(A \Sigma A \Sigma) + 4 \boldsymbol{\mu}^{\top} A \Sigma A \boldsymbol{\mu},$$

$$\mathbb{C}\text{ov}(\boldsymbol{x}^{\top} A_1 \boldsymbol{x}, \boldsymbol{x}^{\top} A_2 \boldsymbol{x}) = 2\text{tr}(A_1 \Sigma A_2 \Sigma) + 4\boldsymbol{\mu}^{\top} A_1 \Sigma A_2 \boldsymbol{\mu}.$$

See Problem B 2.14 in Peng DING.

Theorem 3.1.10 (Independence of quadratic form). Assume $\mathbf{x} \sim N_p(\mu, \Sigma)$.

1. For two symmetric matrix $A, B \in \mathbb{R}^{p \times p}$, $\mathbf{x}^{\top} A \mathbf{x} \perp \mathbf{x}^{\top} B \mathbf{x}$ iff

$$\Sigma A \Sigma B \Sigma = \mathbf{0}, \quad \Sigma A \Sigma B \boldsymbol{\mu} = \Sigma B \Sigma A \boldsymbol{\mu} = \mathbf{0}, \quad \boldsymbol{\mu}^{\top} A \Sigma B \boldsymbol{\mu} = 0.$$

If $\Sigma \succ \mathbf{0}$, then iff $A\Sigma B = \mathbf{0}$.

2. If $\Sigma \succ \mathbf{0}$, $A \in \mathbb{R}^{p \times p}$ symmetric, and $B \in \mathbb{R}^{r \times p}$, then $\mathbf{x}^{\top} A \mathbf{x} \perp \!\!\! \perp B \mathbf{x}$ iff $B \Sigma A = \mathbf{0}$.

See Theorem B.11 in Peng DING, and Problem 1.22–1.23 in [5].

Theorem 3.1.11 (Quadratic form of Σ). If $x, x_{1:N} \sim_{iid} N_p(\mu, \Sigma)$, where Σ is nonsingular, then

- $(x \mu)^{\top} \Sigma^{-1} (x \mu) \sim \chi_p^2$,
- $\boldsymbol{x}^{\top} \Sigma^{-1} \boldsymbol{x} \sim \chi_p^2 (\boldsymbol{\mu}^{\top} \Sigma^{-1} \boldsymbol{\mu}),$
- partition

$$m{x} = egin{bmatrix} m{x}_1 \\ m{x}_2 \end{bmatrix}, \ m{\mu} = egin{bmatrix} m{\mu}_1 \\ m{\mu}_2 \end{bmatrix}, \ \Sigma = egin{bmatrix} \Sigma_{11} & \Sigma_{12} \\ \Sigma_{21} & \Sigma_{22} \end{bmatrix}, \quad m{x}_1, m{\mu}_1 \in \mathbb{R}^k, \ \Sigma_{11} \in \mathbb{R}^{k \times k}, \ then$$

$$Q = (\boldsymbol{x} - \boldsymbol{\mu})^{\top} \Sigma^{-1} (\boldsymbol{x} - \boldsymbol{\mu}) - (\boldsymbol{x}_1 - \boldsymbol{\mu}_1)^{\top} \Sigma_{11}^{-1} (\boldsymbol{x}_1 - \boldsymbol{\mu}_1) \sim \chi_{n-k}^2.$$

- $N(\bar{\boldsymbol{x}}_N \boldsymbol{\mu})^{\top} \Sigma^{-1} (\bar{\boldsymbol{x}}_N \boldsymbol{\mu}) \sim \chi_p^2$,
- the Mahalanobis distance $d_i^2 = (\boldsymbol{x}_i \bar{\boldsymbol{x}}_N)^{\top} \boldsymbol{S}^{-1} (\boldsymbol{x}_i \bar{\boldsymbol{x}}_N) \xrightarrow{\mathrm{d}} \chi_n^2$

If $r(\Sigma) = k \le p$, then

1.

$$(\boldsymbol{x} - \boldsymbol{\mu})^{\top} \Sigma^{-} (\boldsymbol{x} - \boldsymbol{\mu}) \sim \chi_{k}^{2},$$

$$(\boldsymbol{x} - \boldsymbol{\mu})^{\top} \Sigma^{+} (\boldsymbol{x} - \boldsymbol{\mu}) \sim \chi_{k}^{2}.$$

- 2. If Σ is idempotent with $r(\Sigma) = k$, then $(\boldsymbol{x} \boldsymbol{\mu})^{\top}(\boldsymbol{x} \boldsymbol{\mu}) \sim \chi_k^2$.
 - If $\boldsymbol{\mu} \in \operatorname{Col}(\Sigma)$, then $\boldsymbol{x}^{\top} \boldsymbol{x} \sim \chi_k^2(\boldsymbol{\mu}^{\top} \boldsymbol{\mu})$.

Theorem 3.1.12 (Quadratic form of any matrices). If $x \sim N_p(\mu, \Sigma)$,

- 1. if Σ is nonsingular, $B \in \mathbb{R}^{p \times p}$ is symmetric, then $\mathbf{x}^{\top} B \mathbf{x} \sim \chi_k^2(\delta)$ iff $B\Sigma$ is idempotent (equiv., $B\Sigma B = B$), in which case k = r(B) and $\delta = \boldsymbol{\mu}^{\top} B \boldsymbol{\mu}$;
- 2. for $A \in \mathbb{R}^{p \times p}$, $\boldsymbol{x}^{\top} A \boldsymbol{x} \sim \chi^2_{r(A\Sigma)}(\boldsymbol{\mu}^{\top} A \boldsymbol{\mu})$ if

$$(1)\Sigma A\Sigma A\Sigma = \Sigma A\Sigma, \quad (2)\boldsymbol{\mu}^{\top} A\Sigma A\boldsymbol{\mu} = \boldsymbol{\mu}^{\top} A\boldsymbol{\mu}, \quad (3)\Sigma A\Sigma A\boldsymbol{\mu} = \Sigma A\boldsymbol{\mu};$$

- 3. let $A_i \in \mathbb{R}^{p \times p}$ be symmetric with rank k_i for i = 1, ..., m. Denote $A = \sum_{i=1}^m A_i$, which is symmetric with rank k. Then $\mathbf{x}^\top A_i \mathbf{x} \sim \chi_{k_i}^2(\boldsymbol{\mu}^\top A_i \boldsymbol{\mu})$, $\mathbf{x}^\top A_i \mathbf{x}$ are pairwise independent and $\mathbf{x}^\top A_i \mathbf{x} \sim \chi_{k_i}^2(\boldsymbol{\mu}^\top A_i \boldsymbol{\mu})$, iff
 - (I) any two of the following are true: (a) $A_i\Sigma$ idempotent, $\forall i;$ (b) $A_i\Sigma A_j=0, \ \forall i< j;$ and (c) $A\Sigma$ idempotent; OR
 - (II) (c) is true and (d) $k = \sum_{i=1}^{m} k_i$; **OR**
 - (III) (c) is true and (e) $A_1\Sigma, \ldots, A_{m-1}\Sigma$ are idempotent and $A_m\Sigma \succeq \mathbf{0}$.
- 4. (Cochran's theorem) $\mathbf{x} \sim N_p(\mathbf{0}_p, I_p)$ and A_i is symmetric of rank k_i , for i = 1, ..., m with $\sum_{i=1}^m A_i = I_p$, then $\mathbf{x}^\top A_i \mathbf{x} \sim \chi_{k_i}^2$ independently iff $\sum_{i=1}^m k_i = p$.

3.1.2 The noncentral χ^2 and F distribution

3.2 Asymptotic properties CHAPTER 3. MULTIVARIATE INFERENCE FUNDAMENTALS

3.2.1 Asymptotic distributions of sample means and covariance matrices

Refer to section 1.2.2, [5].

Theorem 3.2.1 (CLT for sample means). Let $x_{1:n} \sim_{\text{iid}} [\mu, \Sigma]$, then

$$\sqrt{n}(\bar{\boldsymbol{x}}_n - \boldsymbol{\mu}) = \frac{1}{\sqrt{n}} \sum_{i=1}^n (\boldsymbol{x}_i - \boldsymbol{\mu}) \xrightarrow{\mathrm{d}} \mathrm{N}_p(\boldsymbol{0}_p, \Sigma).$$

Theorem 3.2.2 (CLT for sample covariance matrices). Let $\mathbf{x}_{1:n} \sim_{\text{iid}} [\boldsymbol{\mu}, \Sigma]$ with finite fourth moments, SSCP matrix $\mathbf{A} = \sum_{i=1}^{n} (\mathbf{x}_i - \bar{\mathbf{x}}) (\mathbf{x}_i - \bar{\mathbf{x}})^{\top}$, and $\mathbf{S} = (n-1)^{-1}\mathbf{A}$. Let $V = \mathbb{C}\text{ov}[\text{vec}((\mathbf{x}_1 - \boldsymbol{\mu})(\mathbf{x}_1 - \boldsymbol{\mu})^{\top})]$, then

$$\frac{1}{\sqrt{n}}(\operatorname{vec}(\boldsymbol{A}) - n \cdot \operatorname{vec}(\Sigma)) \xrightarrow{d} \operatorname{N}_{p^2}(\boldsymbol{0}, V),$$

$$\sqrt{n-1}(\operatorname{vec}(\boldsymbol{S}) - \operatorname{vec}(\Sigma)) \xrightarrow{\mathrm{d}} \operatorname{N}_{p^2}(\boldsymbol{0}, V).$$

Note that $V \in \mathbb{R}^{p^2 \times p^2}$ is singular as the LHS vectors above have repeated elements.

Linear Models

4.1 Linear regression for the full-rank model

Consider the linear model

$$Y_n = X_{n \times p} \beta_p + \varepsilon_n, \quad r(X) = p.$$
 (4.1)

- 4.2 Ridge regression
- 4.3 LASSO

Generalized Linear Models

Bayesian Inference

Structural Equation Model (SEM)

Reference:

- CUHK STAT5020 Topics in multivariate analysis (2025 Spring), by Xin Yuan SONG.
- Sik-Yum Lee and Xin-Yuan Song Basic and advanced Bayesian structural equation modeling: With applications in the medical and behavioral sciences [4].

7.1 SEM models

Goal: to examine the relationships among the variables of interest.

Approach: group observed variables to form *latent variables* (a combination of several observed variables).

- Reduce the number of variables compared to direct regression.
- As highly correlated observed variables are grouped into latent variables, the problem induced by multicollinearity is alleviated.
- $\bullet~$ It gives better assessments on the interrelationships of latent constructs.

Definition 7.1.1 (Linear SEMs). Assume that the observed variables $y_p \sim_{\text{iid}} N_p$ with mean μ_p . Let ω_q be latent variables. $\omega_q = (\eta_{q_1}^\top, \xi_{q_2}^\top)^\top$, where η_{q_1} is the key outcome latent variables, and ξ_{q_2} is the explanatory latent variables. Define

(Measurement equation (MeaEq))
$$y_p = \mu_p + \Lambda_{p \times q} \omega_q + \epsilon_p,$$
(Structural equation (StEq))
$$\eta_{q_1} = \Gamma_{q_1 \times q_2} \xi_{q_2} + \delta_{q_1},$$
(7.1)

where Λ is the unknown factor loading matrix, Γ is the unknown matrix of regression coeffcients, and ϵ and δ are measurement (residual) errors.

The basic SEM is formulated by two components:

- 1. (MeaEq) Confirmatory factor analysis (CFA) model groups the highly correlated observed variables into latent variables. We know the information/structure of Λ (e.g., know y_1, y_2 are only affected by η and y_3, y_4 are affected by ξ).
 - Exploratory factor analysis (EFA): we have y and just know they have lower-dimensional structure, but don't know the structure of Λ and q.

It regresses y on a smaller number of latent variables.

2. (StEq) A regression type model, in which η is regressed on ξ .

Assumption 7.1.2 (Standard linear SEMs). For i = 1, ..., n,

- (A1) $\epsilon_i \sim_{\text{iid}} N[\mathbf{0}_p, \mathbf{\Psi}_{\epsilon}]$, where $\mathbf{\Psi}_{\epsilon} \in \mathbb{R}^{p \times p}$ is diagonal.
- (A2) $\xi_i \sim_{\text{iid}} N[\mathbf{0}_{q_2}, \mathbf{\Phi}]$, where $\mathbf{\Phi}$ is a general (we can evaluate the correlations of the latent variables).
- (A3) $\delta_i \sim_{\text{iid}} N[\mathbf{0}_{q_1}, \mathbf{\Psi}_{\delta}]$, where $\mathbf{\Psi}_{\delta}$ is diagonal.
- (A4) $\delta_i \perp \!\!\! \perp \boldsymbol{\xi}_i$, and $\epsilon_i \perp \!\!\! \perp \boldsymbol{\omega}_i, \boldsymbol{\delta}_i$.

These assumptions imply that

$$oldsymbol{\eta}_i \sim_{\mathrm{iid}} \mathrm{N}_{q_1}(\mathbf{0}_{q_1}, oldsymbol{\Gamma} oldsymbol{\Phi} oldsymbol{\Gamma}^ op + oldsymbol{\Psi}_{oldsymbol{\delta}}), \quad oldsymbol{\omega}_i \sim_{\mathrm{iid}} \mathrm{N}_q \left(\mathbf{0}_q, oldsymbol{\Sigma}_{oldsymbol{\omega}} = egin{bmatrix} oldsymbol{\Gamma} oldsymbol{\Phi} oldsymbol{\Gamma}^ op + oldsymbol{\Psi}_{oldsymbol{\delta}} & oldsymbol{\Gamma} oldsymbol{\Phi} \ oldsymbol{\Phi} oldsymbol{\Gamma}^ op & oldsymbol{\Phi} \end{bmatrix} \right), \quad oldsymbol{y}_i \sim_{\mathrm{iid}} \mathrm{N}_p(oldsymbol{\mu}, oldsymbol{\Sigma}(oldsymbol{ heta}) = oldsymbol{\Lambda} oldsymbol{\Sigma}_{oldsymbol{\omega}} oldsymbol{\Lambda}^ op + oldsymbol{\Psi}_{oldsymbol{\epsilon}}).$$

Identifiability issue: The measurement equation is identified if $\forall \theta_1, \theta_2$, MeaEq(θ_1) = MeaEq(θ_2) implies $\theta_1 = \theta_2$. The structural equation is identified if $\forall \theta_1, \theta_2$, StEq(θ_1) = StEq(θ_2) implies $\theta_1 = \theta_2$. The SEM is identified if both of its MeaEq and StEq are identified. For example, consider (7.1) MeaEq:

$$oldsymbol{y} = oldsymbol{\mu} + (oldsymbol{\Lambda} oldsymbol{M})(oldsymbol{M}^{-1}oldsymbol{\omega}) + oldsymbol{\epsilon} = oldsymbol{\mu} + oldsymbol{\Lambda}^*oldsymbol{\omega}^* + oldsymbol{\epsilon}, \quad oldsymbol{\omega}^* \sim \mathrm{N}(oldsymbol{0}_q, oldsymbol{M}^{-1}\mathbb{C}\mathrm{ov}(oldsymbol{\omega})oldsymbol{M}^{- op}).$$

 $\textbf{2We have to impose restrictions on $\boldsymbol{\Lambda}$ and/or $\mathbb{C}\text{ov}(\boldsymbol{\omega})$ scales a transfer on $\boldsymbol{\Sigma}$ The singular MEQUATION MODEL (SEM)}$

• A simple and common method is using a Λ with the *non-overlapping structure*, e.g.,

where 1's are fixed to introduce a scale to latent variables.

• We can also allow λ_{11} , λ_{52} , and/or λ_{83} to be unknowns, and fix diag($\mathbb{C}\text{ov}(\boldsymbol{\omega})$) = I_q , such that $\mathbb{C}\text{ov}(\boldsymbol{\omega})$ is a correlation matrix. But it is not convenient for identifying an SEM with StEq, and induces complication in the Bayesian analysis.

Example 7.1.3. Study the kidney disease of type 2 diabetic patients. We observe: plasma creatine (PCr), urinary albumin creatinine ratio (ACR), systolic blood pressure (SBP), diastolic blood pressure (DBP), body mass index (BMI), waist hip ratio (WHR), glycated hemoglobin (HbAlc), fasting plasma glucose (FPG). Group

- {PCr, ACR}: 'kidney disease (KD)'
- {SBP, DBP}: 'blood pressure (BP)'
- {BMI, WHR}: 'obesity (OB)'
- {HbA1c, FPG}: 'glycemic control (GC)'

 $\boldsymbol{y} = (PCr, ACR, SBP, DBP, BMI, WHR)^{\top}, \ \boldsymbol{\omega} = (KD, BP, OB)^{\top}, \ \boldsymbol{\eta} = KD, \ \boldsymbol{\xi} = (BP, OB)^{\top}, \ p = 6, \ q = 3, q_1 = 1, q_2 = 2.$ Then

$$\text{(MeaEq)} \qquad \begin{bmatrix} \text{PCr} \\ \text{ACR} \\ \text{SBP} \\ \text{DBP} \\ \text{BMI} \\ \text{WHR} \end{bmatrix} = \begin{bmatrix} \mu_1 \\ \mu_2 \\ \mu_3 \\ \mu_4 \\ \mu_5 \\ \mu_6 \end{bmatrix} + \begin{bmatrix} \lambda_{11} & 0 & 0 \\ \lambda_{21} & 0 & 0 \\ 0 & \lambda_{32} & 0 \\ 0 & \lambda_{42} & 0 \\ 0 & 0 & \lambda_{53} \\ 0 & 0 & \lambda_{63} \end{bmatrix} \begin{bmatrix} \text{KD} \\ \text{BP} \\ \text{OB} \end{bmatrix} + \begin{bmatrix} \epsilon_1 \\ \epsilon_2 \\ \epsilon_3 \\ \epsilon_4 \\ \epsilon_5 \\ \epsilon_6 \end{bmatrix}.$$

We know that KD is only linked with PCr and ACR.

(StEq)
$$KD = \gamma_1 BP + \gamma_2 OB + \delta.$$

Suppose we wish to study the effects of $\boldsymbol{\xi}$ on $\boldsymbol{\eta} = (\mathrm{KD}, \eta_A)^{\top}, q_1 = 2,$

$$\begin{pmatrix} \mathrm{KD} \\ \eta_A \end{pmatrix} = \begin{pmatrix} 0 & 0 \\ \pi & 0 \end{pmatrix} \begin{pmatrix} \mathrm{KD} \\ \eta_A \end{pmatrix} + \begin{pmatrix} \gamma_1 & \gamma_2 \\ \gamma_3 & \gamma_4 \end{pmatrix} \begin{pmatrix} \mathrm{BP} \\ \mathrm{OB} \end{pmatrix} + \begin{pmatrix} \delta \\ \delta_A \end{pmatrix}.$$

Definition 7.1.4 (More SEM models). Extensions of StEq:

- Basic StEq: $\eta_{q_1} = \Gamma_{q_1 \times q_2} \boldsymbol{\xi}_{q_2} + \boldsymbol{\delta}_{q_1}$.
- An extension:

$$oldsymbol{\eta}_{q_1} = oldsymbol{\Pi}_{q_1 imes q_1} oldsymbol{\eta}_{q_1} + oldsymbol{\Gamma}_{q_1 imes q_2} oldsymbol{\xi}_{q_2} + oldsymbol{\delta}_{q_1},$$

where Π is a matrix of unknown coefficients such that $I_{q_1} - \Pi$ is nonsingular and the diagonal elements of Π are zero.

• Add fixed known covariates d:

$$oldsymbol{\eta}_{q_1} = oldsymbol{B}_{q_1 imes r_2} oldsymbol{d}_{r_2} + oldsymbol{\Pi}_{q_1 imes q_1} oldsymbol{\eta}_{q_1} + oldsymbol{\Gamma}_{q_1 imes q_2} oldsymbol{\xi}_{q_2} + oldsymbol{\delta}_{q_1}.$$

• Add nonlinear structure:

$$\eta_{q_1} = B_{q_1 \times r_2} d_{r_2} + \Pi_{q_1 \times q_1} \eta_{q_1} + \Gamma_{q_1 \times t} F_t(\xi_{q_2}) + \delta_{q_1}, \tag{7.2}$$

where $F(\boldsymbol{\xi}) = (f_1(\boldsymbol{\xi}), \dots, f_t(\boldsymbol{\xi}))^{\top}$ with nonzero, known, and linearly independent differentiable functions $f_1, \dots, f_t, t \geq q_2$. To identity StEq, structures like $F_1(\boldsymbol{\xi}) = (\xi_1, \xi_2, \xi_1^2, \xi_1^2)^{\top}$ and $F_2(\boldsymbol{\xi}) = (\xi_1, \xi_2, \xi_1 \xi_2, 0)^{\top}$ are not allowed.

• Combine d and F:

$$\eta_{q_1} = \Pi_{q_1 \times q_1} \eta_{q_1} + \Lambda_{\omega, q_1 \times t} G_t(d, \boldsymbol{\xi}) + \delta_{q_1}, \tag{7.3}$$

where $G(d\boldsymbol{\xi}) = (g_1(\boldsymbol{d}, \boldsymbol{\xi}), \dots, g_t(\boldsymbol{d}, \boldsymbol{\xi}))^{\top}$ is a vector-valued function with nonzero, known, and linearly independent differentiable functions. (7.2) can be obtained by letting $\Lambda_{\omega} = (\boldsymbol{B}, \boldsymbol{\Gamma})$ and $G(\boldsymbol{d}, \boldsymbol{\xi}) = (\boldsymbol{d}^{\top}, \boldsymbol{F}(\boldsymbol{\xi})^{\top})^{\top}$. An Extension of MeaEq:

• Add fixed known covariates c (If μ_p is included, then $c = [1, c_2]^{\top}$.)

$$\mathbf{y}_p = \mathbf{A}_{p \times r_1} \mathbf{c}_{r_1} + \mathbf{\Lambda}_{p \times q} \mathbf{\omega}_q + \mathbf{\epsilon}_p. \tag{7.4}$$

Let Λ_k^{\top} be the kth row of Λ , and $\Lambda_k^{\top} = (\Lambda_{k\eta}^{\top}, \Lambda_{k\xi}^{\top})$ be a partition correspondings to the partition of $\boldsymbol{\omega} = (\boldsymbol{\eta}^{\top}, \boldsymbol{\xi}^{\top})^{\top}$. For StEq (7.2) with $r_2 = 0$ and MeaEq (7.1),

$$\mathbb{E}(\boldsymbol{\xi}) = \mathbf{0}_{q_2}, \quad \mathbb{E}(\boldsymbol{\eta}) = [(I_{q_1} - \boldsymbol{\Pi})^{-1} \boldsymbol{\Gamma}] \mathbb{E}(\boldsymbol{F}(\boldsymbol{\xi})),$$

$$\mathbb{E}(y_k) = \mu_k + \boldsymbol{\Lambda}_{k\boldsymbol{\eta}}^{\top} [(I_{q_1} - \boldsymbol{\Pi})^{-1} \boldsymbol{\Gamma}] \mathbb{E}(\boldsymbol{F}(\boldsymbol{\xi})).$$

$$\mathbb{E}(y_k) = \boldsymbol{A}_k^\top \boldsymbol{c} + \boldsymbol{\Lambda}_{k\boldsymbol{\eta}}^\top \mathbb{E}(\boldsymbol{\eta}) = \boldsymbol{A}_k^\top \boldsymbol{c} + \boldsymbol{\Lambda}_{k\boldsymbol{\eta}}^\top [(I_{q_1} - \boldsymbol{\Pi})^{-1} \boldsymbol{\Lambda}_{\boldsymbol{\omega}}] \mathbb{E}(\boldsymbol{G}(\boldsymbol{d}, \boldsymbol{\xi})).$$

Issues of developing a comprehensive SEM:

- 1. Make sure the sample size of y is large enough to achieve accurate statistical results.
- 2. If the size of the proposed SEM and the number of parameters are large, we may encounter diffculties in achieving convergence of the related computing algorithm for obtaining statistical results.
- 3. So far, the most general SEM is MeaEq (7.4) and StEq (7.3):

$$y = Ac + \Lambda\omega + \epsilon, \qquad \eta = \Pi\eta + \Lambda_\omega G(d, \xi) + \delta.$$

It has limitations. As $G(d, \xi)$ does not involve any η_k in η , nonlinear terms related to η_k cannot be used to predict the other $\eta_{k'}$'s, i.e., η_k cannot be accommodated in $G(d, \xi)$ and nonlinear effects of η_k on $\eta_{k'}$ cannot be assessed.

7.2 Bayesian methods for estimating SEM

<u>Motivation</u>: A traditional method is the *covariance structure approach*, which focuses on fitting the covariance structure under the proposed model to the sample covariance matrix S.

• In many complex situations, deriving the explicit covariance structure $\Sigma = \mathbb{C}\text{ov}(\boldsymbol{y})$ or obtaining an appropriate \boldsymbol{S} is difficult (e.g., when \exists missing data).

Advarntages of Bayesian estimates (BE): $\log \mathbb{P}(\theta \mid Y) = \log \mathbb{P}(Y \mid \theta) + \log \mathbb{P}(\theta) + \text{const}$

- 1. Methods are based on the first moment properties of y which are simpler than the second moment properties of S. Hence, it has potential to be applied to more complex situations.
- 2. It estimates ω directly, which cannot be obtained with classical methods.
- 3. As $n \to \infty$, $\log \mathbb{P}(Y \mid \theta)$ could dominate $\log \mathbb{P}(\theta)$, hence the BE have the same optimal properties as the MLE.
- 4. It allows the use of genuine prior information for producing better results. With small or moderate sample sizes, $\mathbb{P}(\theta)$ plays a more substantial role in BE than $\mathbb{P}(Y \mid \theta)$, and is useful for achieving better results.
- 5. It provides more easily assessable statistics for goodness-of-fit and model comparison, and also other useful statistics such as the posterior mean and percentiles. (Don't need to derive the asymptotic distributions.)
- 6. It can give more reliable results for small samples.

7.2.1 Priors for SEM

Informative prior:

Assumption 7.2.1 (Conjugate prior for SEMs). In developing the Bayesian methods for analyzing SEMs, we usually assign fixed known values to the hyperparameters in the conjugate prior distributions. Consider

$$egin{aligned} m{y}_i &= m{\mu} + m{\Lambda} m{\omega}_i + m{\epsilon}_i, \ m{\eta}_i &= m{B} m{d}_i + m{\Pi} m{\eta}_i + m{\Gamma} m{F}(m{\xi}_i) + m{\delta}_i = m{\Lambda}_{m{\omega}} m{G}(m{\omega}_i) + m{\delta}_i, \end{aligned}$$

where $\Lambda_{\omega} = (\boldsymbol{B}, \boldsymbol{\Pi}, \boldsymbol{\Gamma}) \in \mathbb{R}^{q_1 \times (r_2 + q_1 + t)}$, and $\boldsymbol{G}(\omega_i) = (\boldsymbol{d}_i^{\top}, \boldsymbol{\eta}_i^{\top}, \boldsymbol{F}(\boldsymbol{\xi}_i)^{\top})^{\top} \in \mathbb{R}^{r_2 + q_1 + t}$. Assumption 7.1.2 is satisfied. $\boldsymbol{\xi}_i \sim_{\text{iid}} N[\mathbf{0}_{q_2}, \boldsymbol{\Phi}]$, $\boldsymbol{\epsilon}_i \sim_{\text{iid}} N[\mathbf{0}_p, \boldsymbol{\Psi}_{\boldsymbol{\epsilon}} = \text{diag}(\psi_{\boldsymbol{\epsilon}k})]$, and $\boldsymbol{\delta}_i \sim_{\text{iid}} N[\mathbf{0}_{q_1}, \boldsymbol{\Psi}_{\boldsymbol{\delta}} = \text{diag}(\psi_{\boldsymbol{\delta}k})]$.

• Prior (conjugate) for $\theta_{\boldsymbol{y}} = (\boldsymbol{\mu}, \boldsymbol{\Lambda}, \boldsymbol{\Psi}_{\boldsymbol{\epsilon}})$: let $\boldsymbol{\Lambda}_{k}^{\top}$ be the kth row of $\boldsymbol{\Lambda}$

$$\psi_{\epsilon k} \sim \text{IG}(\alpha_{0\epsilon k}, \beta_{0\epsilon k}), \quad [\mathbf{\Lambda}_k \mid \psi_{\epsilon k}] \sim N_q(\mathbf{\Lambda}_{0k}, \psi_{\epsilon k} \mathbf{H}_{0\mathbf{y}k}), \ k = 1, \dots, p, \\
\mathbf{\mu} \sim N_p(\mathbf{\mu}_0, \mathbf{\Sigma}_0).$$

• Prior (conjugate) for $\theta_{\omega} = (\Lambda_{\omega}, \Psi_{\delta}, \Phi)$: let $\Lambda_{\omega k}^{\top}$ be the kth row of $\Lambda_{\omega}, \Psi_{\delta}$,

$$\Phi \sim \mathrm{IW}_{q_2}(\boldsymbol{R}_0^{-1}, \rho_0), \text{ or } \Phi^{-1} \sim \mathrm{W}_{q_2}(\boldsymbol{R}_0, \rho_0)$$

 $\psi_{\delta k} \sim \mathrm{IG}(\alpha_{0\delta k}, \beta_{0\delta k}), \quad [\boldsymbol{\Lambda}_{\omega k} \mid \psi_{\delta k}] \sim \mathrm{N}_{r_2+q_1+t}(\boldsymbol{\Lambda}_{0\omega k}, \psi_{\delta k} \boldsymbol{H}_{0\omega k}), \quad k = 1, \dots, q_1.$

• Assume the prior $\theta_u \perp \!\!\! \perp \theta_\omega$.

Hyperparameter selection: If we have good prior information about a parameter, select the prior distribution with a small variance. E.g.,

- if $\Lambda_k \approx \Lambda_{0k}$, then $H_{0yk} = 0.5I_q$. If not, select the prior with a larger variance;
- since $\epsilon_{ik} \sim N(0, \psi_{\epsilon k})$, if the variation is small, $\psi_{\epsilon k}$ is small, then choose small $\mathbb{E}(\psi_{\epsilon k}) = \beta_{0\epsilon k}/(\alpha_{0\epsilon k} 1)$ and $\mathbb{V}ar(\psi_{\epsilon k}) = \beta_{0\epsilon k}^2/\{(\alpha_{0\epsilon k} 1)^2(\alpha_{0\epsilon k} 2)\};$
- if $\Phi \approx \Phi_0$, since $\mathbb{E}(\Phi) = R_0^{-1}/(\rho_0 q_2 1)$, choose $R_0^{-1} = (\rho_0 q_2 1)\Phi_0$.

Noninformative prior (Jeffrey): If information is not a value and Straughter and Straughter Equation Model (SEM)

$$\mathbb{P}(\boldsymbol{\Lambda}, \boldsymbol{\Psi}_{\boldsymbol{\epsilon}}) \propto \mathbb{P}(\psi_{\boldsymbol{\epsilon}1}, \cdots, \psi_{\boldsymbol{\epsilon}p}) \propto \prod_{k=1}^{p} \psi_{\boldsymbol{\epsilon}k}^{-1}, \quad \mathbb{P}(\boldsymbol{\Lambda}_{\boldsymbol{\omega}}, \boldsymbol{\Psi}_{\boldsymbol{\delta}}) \propto \mathbb{P}(\psi_{\boldsymbol{\delta}1}, \cdots, \psi_{\boldsymbol{\delta}q_1}) \propto \prod_{k=1}^{q_1} \psi_{\boldsymbol{\delta}k}^{-1},$$

$$\mathbb{P}(\boldsymbol{\Phi}) \propto |\boldsymbol{\Phi}|^{-(q_2+1)/2}.$$

If the sample size is large, can use a portion of the data to estimate Λ_{0k} , $\Lambda_{0\omega k}$ and Φ_0 with noninformative priors. If the sample size is moderate, can use the same data twice.

7.2.2 Bayesian estimation using MCMC

Model: Linear SEM with fixed covariates without intercept:

$$egin{aligned} oldsymbol{y}_i &= oldsymbol{\Lambda} oldsymbol{\omega}_i + oldsymbol{\epsilon}_i, \ oldsymbol{\eta}_i &= oldsymbol{B} oldsymbol{d}_i + oldsymbol{\Pi} oldsymbol{\eta}_i + oldsymbol{\Gamma} oldsymbol{\xi}_i + oldsymbol{\delta}_i, \end{aligned}$$

where $\Lambda_{\boldsymbol{\omega}} = (\boldsymbol{B}, \boldsymbol{\Pi}, \boldsymbol{\Gamma}) \in \mathbb{R}^{q_1 \times (r_2 + q_1 + q_2)}$, and $\boldsymbol{v}_i = (\boldsymbol{d}_i^{\top}, \boldsymbol{\eta}_i^{\top}, \boldsymbol{\xi}_i^{\top})^{\top} \in \mathbb{R}^{r_2 + q_1 + q_2}$. That is, assume $\boldsymbol{\mu} = \boldsymbol{0}_p$ and $\boldsymbol{F}(\boldsymbol{\xi}_i) - \boldsymbol{\xi}_i$. Denote data $\boldsymbol{Y} = (\boldsymbol{y}_1, \dots, \boldsymbol{y}_n) = (\boldsymbol{Y}_1, \dots, \boldsymbol{Y}_p)^{\top} \in \mathbb{R}^{p \times n}, \ \boldsymbol{V} = (\boldsymbol{v}_1, \dots, \boldsymbol{v}_n) = (\boldsymbol{V}_1, \dots, \boldsymbol{V}_{r_2 + q_1 + q_2})^{\top}, \ \boldsymbol{\Xi}_k = (\eta_{1k}, \dots, \eta_{nk})^{\top}$ for $k = 1, \dots, q_1$, matrix of latent variables $\boldsymbol{\Omega} = (\boldsymbol{\omega}_1, \dots, \boldsymbol{\omega}_n) \in \mathbb{R}^{q \times n}, \ \boldsymbol{\Omega}_1 = (\boldsymbol{\eta}_1, \dots, \boldsymbol{\eta}_n), \ \boldsymbol{\Omega}_2 = (\boldsymbol{\xi}_1, \dots, \boldsymbol{\xi}_n)$, and

$$\theta = (\boldsymbol{\Lambda}, \boldsymbol{B}, \boldsymbol{\Pi}, \boldsymbol{\Gamma}, \boldsymbol{\Phi}, \boldsymbol{\Psi}_{\boldsymbol{\epsilon}}, \boldsymbol{\Psi}_{\boldsymbol{\delta}}) = (\underbrace{\boldsymbol{\Lambda}, \boldsymbol{\Psi}_{\boldsymbol{\epsilon}}}_{\boldsymbol{\theta}_{\boldsymbol{y}}}, \underbrace{\boldsymbol{\Lambda}_{\boldsymbol{\omega}}, \boldsymbol{\Phi}, \boldsymbol{\Psi}_{\boldsymbol{\delta}}}_{\boldsymbol{\theta}_{\boldsymbol{\omega}}}).$$

Proposition 7.2.2. The above model has the following posterior distributions:

1. Conditional distribution $\mathbb{P}(\Omega \mid Y, \theta) = \prod_{i=1}^{n} \mathbb{P}(\omega_i \mid y_i, \theta) \propto \prod_{i=1}^{n} \mathbb{P}(\omega_i \mid \theta) \mathbb{P}(y_i \mid \omega_i, \theta)$, where

$$egin{aligned} [oldsymbol{\omega}_i \mid oldsymbol{ heta}] &\sim \mathrm{N}_q(oldsymbol{\mu}_{oldsymbol{\omega}_i}, oldsymbol{\Sigma}_{oldsymbol{\omega}}), & [oldsymbol{y}_i \mid oldsymbol{\omega}_i, oldsymbol{ heta}] &\sim \mathrm{N}_p(oldsymbol{\Lambda}oldsymbol{\omega}_i, oldsymbol{\Sigma}_{oldsymbol{\omega}}), oldsymbol{\Sigma}^{*-1}(oldsymbol{\Sigma}_{oldsymbol{\omega}}^{-1} oldsymbol{\mu}_{oldsymbol{\omega}_i} + oldsymbol{\Lambda}^{\top} oldsymbol{\Psi}_{oldsymbol{\epsilon}}^{-1} oldsymbol{y}_i), oldsymbol{\Sigma}^{*-1}) \end{aligned}$$

where

$$egin{aligned} &\mathbf{\Pi}_0 = I_{q_1} - \mathbf{\Pi}, \ \boldsymbol{\mu}_{oldsymbol{\omega}_i} = egin{pmatrix} \mathbf{\Pi}_0^{-1} B oldsymbol{d}_i \ \mathbf{0}_{q_2} \end{pmatrix}, \ \boldsymbol{\Sigma}_{oldsymbol{\omega}} = egin{bmatrix} \mathbf{\Pi}_0^{-1} (oldsymbol{\Gamma} oldsymbol{\Phi} oldsymbol{\Gamma}^{\top} + oldsymbol{\Psi}_{oldsymbol{\delta}}) \mathbf{\Pi}_0^{-\top} & \mathbf{\Pi}_0^{-1} oldsymbol{\Gamma} oldsymbol{\Phi} \ \mathbf{\Phi} \end{bmatrix}, \ & \mathbf{\Sigma}^* = \mathbf{\Sigma}_{oldsymbol{\omega}}^{-1} + \mathbf{\Lambda}^{\top} oldsymbol{\Psi}_{oldsymbol{\epsilon}}^{-1} oldsymbol{\Lambda}. \end{aligned}$$

2. Assume all elements of Λ_k and Λ_{ω} are unknown, conditional distribution $\mathbb{P}(\boldsymbol{\theta} \mid \boldsymbol{Y}, \boldsymbol{\Omega}) = \mathbb{P}(\boldsymbol{\theta_y} \mid \boldsymbol{Y}, \boldsymbol{\Omega})\mathbb{P}(\boldsymbol{\theta_\omega} \mid \boldsymbol{Y}, \boldsymbol{\Omega})$, where we can show $[\Lambda_k, \psi_{\boldsymbol{\epsilon}k} \mid \boldsymbol{Y}, \boldsymbol{\Omega}] \perp \mathcal{I}$ and $[\Lambda_{\omega k}, \psi_{\boldsymbol{\delta}k} \mid \boldsymbol{\Omega}] \perp \mathcal{I}$, and

$$egin{aligned} \mathbb{P}(oldsymbol{ heta_y} \mid oldsymbol{Y}, oldsymbol{\Omega}) & \propto \prod_{k=1}^p \mathbb{P}(oldsymbol{\Lambda}_k, \psi_{oldsymbol{\epsilon}k} \mid oldsymbol{Y}, oldsymbol{\Omega}), \ \mathbb{P}(oldsymbol{ heta_\omega} \mid oldsymbol{Y}, oldsymbol{\Omega}) & \propto \left[\prod_{k=1}^{q_1} \mathbb{P}(oldsymbol{\Lambda}_{oldsymbol{\omega}k}, \psi_{oldsymbol{\delta}k} \mid oldsymbol{\Omega})
ight] \mathbb{P}(oldsymbol{\Phi} \mid oldsymbol{\Omega}_2) \end{aligned}$$

where

$$\mathbb{P}(\boldsymbol{\Lambda}_{k}, \psi_{\boldsymbol{\epsilon}k}^{-1} \mid \boldsymbol{Y}, \boldsymbol{\Omega}) \propto \mathrm{N}_{q}(\boldsymbol{a}_{k}, \psi_{\boldsymbol{\epsilon}k}\boldsymbol{A}_{k}) \cdot \mathrm{Ga}(n/2 + \alpha_{0\boldsymbol{\epsilon}k}, \beta_{\boldsymbol{\epsilon}k}),$$

$$\mathbb{P}(\boldsymbol{\Lambda}_{\boldsymbol{\omega}k}, \psi_{\boldsymbol{\delta}k}^{-1} \mid \boldsymbol{\Omega}) \propto \mathrm{N}_{r_{2}+q_{1}+q_{2}}(\boldsymbol{a}_{\boldsymbol{\omega}k}, \psi_{\boldsymbol{\delta}k}\boldsymbol{A}_{\boldsymbol{\omega}k}) \cdot \mathrm{Ga}(n/2 + \alpha_{0\boldsymbol{\delta}k}, \beta_{\boldsymbol{\delta}k}),$$

$$[\boldsymbol{\Phi} \mid \boldsymbol{\Omega}_{2}] \sim \mathrm{IW}_{q_{2}}[(\boldsymbol{\Omega}_{2}\boldsymbol{\Omega}_{2}^{\top} + \boldsymbol{R}_{0}^{-1}), n + \rho_{0}].$$

where

$$\begin{split} \boldsymbol{A}_k &= (\boldsymbol{H}_{0yk}^{-1} + \boldsymbol{\Omega} \boldsymbol{\Omega}^\top)^{-1}, \ \boldsymbol{a}_k = \boldsymbol{A}_k (\boldsymbol{H}_{0yk}^{-1} \boldsymbol{\Lambda}_{0k} + \boldsymbol{\Omega} \boldsymbol{Y}_k), \\ \boldsymbol{\beta}_{\boldsymbol{\epsilon}k} &= \boldsymbol{\beta}_{0\boldsymbol{\epsilon}k} + \frac{1}{2} (\boldsymbol{Y}_k^\top \boldsymbol{Y}_k - \boldsymbol{a}_k^\top \boldsymbol{A}_k^{-1} \boldsymbol{a}_k + \boldsymbol{\Lambda}_{0k}^\top \boldsymbol{H}_{0yk}^{-1} \boldsymbol{\Lambda}_{0k}), \\ \boldsymbol{A}_{\boldsymbol{\omega}k} &= (\boldsymbol{H}_{0\boldsymbol{\omega}k}^{-1} + \boldsymbol{V}_k \boldsymbol{V}_k^\top)^{-1}, \ \boldsymbol{a}_{\boldsymbol{\omega}k} = \boldsymbol{A}_{\boldsymbol{\omega}k} (\boldsymbol{H}_{0\boldsymbol{\omega}k}^{-1} \boldsymbol{\Lambda}_{0\boldsymbol{\omega}k} + \boldsymbol{V}_k \boldsymbol{\Xi}_k), \\ \boldsymbol{\beta}_{\boldsymbol{\delta}k} &= \boldsymbol{\beta}_{0\boldsymbol{\delta}k} + \frac{1}{2} (\boldsymbol{\Xi}_k^\top \boldsymbol{\Xi}_k - \boldsymbol{a}_{\boldsymbol{\omega}k}^\top \boldsymbol{A}_{\boldsymbol{\omega}k}^{-1} \boldsymbol{a}_{\boldsymbol{\omega}k} + \boldsymbol{\Lambda}_{0\boldsymbol{\omega}k}^\top \boldsymbol{H}_{0\boldsymbol{\omega}k}^{-1} \boldsymbol{\Lambda}_{0\boldsymbol{\omega}k}). \end{split}$$

Remark 7.2.3. For general nonlinear SEMs with MeaEq (7.4) and StEq (7.3), we can defind $\boldsymbol{u} = [\boldsymbol{c}^{\top}, \boldsymbol{\omega}^{\top}]^{\top} \in \mathbb{R}^{r_1+q}$ and use similar procedure to derive the full conditional distributions. But by the nonlinear structure $\boldsymbol{G}(\boldsymbol{\omega})$, $\mathbb{P}(\boldsymbol{\Omega} \mid \boldsymbol{Y}, \boldsymbol{\theta})$ may not have closed form like normal, while $\mathbb{P}(\boldsymbol{\theta} \mid \boldsymbol{Y}, \boldsymbol{\Omega})$ is not affected and keeps normal-Gamma. To handle fixed parameters, see Appendix 3.3 in [4] and Sec 4.3.1 and Appendix 4.3 in [3]. Also see STAT5020 HW1 Q3.

7.3 Bayesian model comparison

Definition 7.4.1 (Two-level nonlinear SEM with mixed type variables). Consider a collection of p-variate random vectors $\boldsymbol{u}_{gi}, i = 1, \dots, N_g$, nested within groups $g = 1, \dots, G$.

(Within-groups)
$$\boldsymbol{u}_{gi} = \boldsymbol{v}_g + \boldsymbol{\Lambda}_{1g}\boldsymbol{\omega}_{1gi} + \boldsymbol{\epsilon}_{1gi}, \quad g = 1, \dots, G, \quad i = 1, \dots, N_g,$$

(Between-groups) $\boldsymbol{v}_g = \boldsymbol{\mu} + \boldsymbol{\Lambda}_2\boldsymbol{\omega}_{2g} + \boldsymbol{\epsilon}_{2g}, \quad g = 1, \dots, G$

where $\mathbf{\Lambda}_{1g} \in \mathbb{R}^{p \times q_1}$, $\boldsymbol{\epsilon}_{1gi} \in \mathbb{R}^{q_1}$, $\boldsymbol{\epsilon}_{1gi} \in \mathbb{R}^p \sim \mathrm{N}(\mathbf{0}, \boldsymbol{\Psi}_{1g})$ independent of $\boldsymbol{\omega}_{1gi}$, where $\boldsymbol{\Psi}_{1g}$ is diagonal. Then

$$\mathbf{u}_{gi} = \mu + \mathbf{\Lambda}_2 \omega_{2g} + \epsilon_{2g} + \mathbf{\Lambda}_{1g} \omega_{1gi} + \epsilon_{1gi}.$$

The SEM is

$$\begin{split} \eta_{1gi} &= \Pi_{1g} \eta_{1gi} + \Gamma_{1g} \mathbf{F}_{1}(\xi_{1gi}) + \delta_{1gi}, \\ \eta_{2g} &= \Pi_{2} \eta_{2g} + \Gamma_{2} \mathbf{F}_{2}(\xi_{2g}) + \delta_{2g}, \end{split}$$

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