

MODULE 7.1

Singular Spectrum Analysis (SSA)

A PRELUDE

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- Time Series(TS) Analysis(TSA) - univariate/scalar
- multivariate/vector
- PCA/EoF (earlier modules)
- SSA is EoF applied to scalar TS
- Multi-channel SSA-(MSSA) : EoF applied to vector TS

- TSA is a chapter in the vast and growing discipline of stochastic process in discrete time
- TSA had its beginnings in 1920's
- Developed by Yule, Walker, Wold, Wiener Kolmogorov, among others

- Given a TS, quantify the linear dependence structure using Auto Correlation Function(ACF) and Partial Auto Correlation Function(PACF)
- Identify a class of (empirical) discrete time stochastic dynamic models that could capture ACF and PACF
- Estimate the parameters of the chosen models

- Compare the adequacy of the models using several measures - residual, Akaike information criterion etc.,
- Choose a very small number of "good" modules
- Develop algorithms for prediction and prediction error

Sources of TSA

- J.D.Hamilton(1995) Time Series Analysis, Princeton university press
- P.S.Brockwell and R.A.Davis(2013) Time Series Theory and Methods Springer
- W.A.Fuller(2009) Introduction to Statistical Time Series, John Wiley and sons

What is SSA?

- SSA is an alternate method to analyze time series data
- It starts by generating a data/trajectory matrix $x \in R^{m \times n}$ from the given scalar TS: $\{y_t / 1 \leq t \leq N\}$
- SSA follows the foot steps of EoF analysis applied to x

What is SSA?

- The idea is to reconstruct the signal, trend, seasonal cycles and predictable parts and noise components using the spectral properties of the covariance matrix $\Sigma = \frac{1}{n}xx^T$
- Since the eigenvalues of the Gramian, $\frac{1}{n}xx^T$ are the singular values of x , the name Singular Spectrum Analysis

Origins of SSA

- SSA was developed rather recently:
- J.M.Colebrook(1978)" Continous plakton records:Zooplakton and environment, North-East Atlantic and North Sea, 1948-1975",Oceanologica ACTA, vol 1, pp 9-23
- D.S.Broomhead and G.P.King(1986) "Extracting Qualitative Dynamics from Experimental Data" Physica, 20D, pp 217-236
- K.Fraedrich(1986) "Estimating the dimension of weather and climate attractor", Journal of Atmospheric Sciences, vol 43, pp 419-43

- J.B.Elsener and A.A.Tsonis(1996) Singular Spectrum Analysis, plenum press, New york
- N.Golyandina, V.Nekrutkin and A.Zhigljavsky(2001) Analysis of Time Series Structure: SSA and related techniques Chapman and Hall
- M.Ghil, et.al.(2002)"Advanced spectral methods for climate time series", Reviews of Geophysics, vol 40, 3-1 to 3-41 - Deals with SSA,MSSA and applications

From TS to Data/trajectory matrix: $x \in R^{m \times n}$

- Let $\{y_t | 1 \leq t \leq N\}$ be the given scalar TS data
- Pick an integer $m : 1 < m \leq \frac{N}{2}$ and let $n = N - m + 1$
- m is called the window length and n is the number of continuous windows that can be formed from the given time series of length N

- The j^{th} column, x_{*j} of x is given by the m entries in the TS starting from location j where $1 \leq j \leq n$:

$$x_{*j} = (y_j, y_{j+1}, \dots, y_{j+m-1})^T \quad (1)$$

- Then

$$x = [x_{*1}, x_{*2}, \dots, x_{*n}] \in R^{m \times n} \quad (2)$$

matrix associated with the TS, y_t

Illustration

- Let $N = 6$ and $m = 3$. Then, $n = N - m + 1 = 4$
- $y_t = \{y_1, y_2, y_3, y_4, y_5, y_6\}$
- $x_{*j} = (y_j, y_{j+1}, \dots, y_{j+m-1})^T$
- $x = \begin{bmatrix} x_{11} & x_{12} & x_{13} & x_{14} \\ x_{21} & x_{22} & x_{23} & x_{24} \\ x_{31} & x_{32} & x_{33} & x_{34} \end{bmatrix} = \begin{bmatrix} y_1 & y_2 & y_3 & y_4 \\ y_2 & y_3 & y_4 & y_5 \\ y_3 & y_4 & y_5 & y_6 \end{bmatrix}$

- From the definition:

$$x = [x_{ij}] \quad \text{and} \quad x_{ij} = y_{i+j-1} \quad (3)$$

- The matrix x is such that elements along the anti diagonal for which $(i + j) = c$ for $2 \leq c \leq m + n$ are the same
- Such matrices are called Hankel matrices

Origin of trajectory matrices: Dynamical system

- Trajectory matrices were introduced to solve a class of inverse problems in dynamical system theory
- Let $f : R^n \rightarrow R^n$ be a smooth vector field and

$$\dot{x} = f(x) \quad (4)$$

be the given dynamical system(DS)

The direct problem

- Given the DS in (4) and an initial condition $x_0 \in R^n$, numerically compute the time series $x(t)$ for $t \geq 0$
- The standard Runge - Kutta methods is often used for this purpose

The inverse problem

- Given only the time series of the i^{th} component, $x_i(t)$ of the solution $\mathbf{x}(t)$ of (4), infer the qualitative properties of the DS that generates $\mathbf{x}(t)$
- These include fixed points and their stability, properties of attractors and their properties etc.,

Origin of the Trajectory matrix

- Trajectory matrix was introduced in the context of solving the above inverse problem: Derive the phase space characterization, from the given TS using the trajectory matrix defined above
- Packard et.al (1980)
- Ruelle (1980)
- Takens (1981)

- There is a rich literature on the study of nonlinear time series in the context of chaotic dynamics
- References are given at the end of this module

- Let $y_t \in R^L$, for some finite integer $L > 1$ be the given vector of time series for $1 \leq t \leq N$
- Define the integers m and n as above
- MSSA starts a trajectory matrix $x \in R^{m \times n}$ by stacking together the trajectory matrices $x(i) \in R^{m \times n}$ for each component y_i, t of the vector $y(t)$

An illustration

- Let $L = 3$, $N = 6$, $m = 3$ and $n = 4$
- The given series: $y_t = (y_{1t}, y_{2t}, y_{3t})^T \in R^3 : 1 \leq t \leq N$
- Then

$$x(1) = \begin{bmatrix} y_{1,1} & y_{1,2} & y_{1,3} & y_{1,4} \\ y_{1,2} & y_{1,3} & y_{1,4} & y_{1,5} \\ y_{1,3} & y_{1,4} & y_{1,5} & y_{1,6} \end{bmatrix} \quad (5)$$

be the 3×4 trajectory matrix built out of the first component $\{y_{1,t} | 1 \leq t \leq 6\}$ of the given vector time series $\{y(t) | 1 \leq t \leq 6\}$

Illustration - continued

- Likewise build trajectory matrices $x(2)$ and $x(3)$ from the second and third components of $\{y_t \in R^3 | 1 \leq t \leq N\}$
- Then

$$x = \begin{matrix} & n \\ m & \begin{bmatrix} x(1) \\ \vdots \\ x(2) \\ \vdots \\ x(3) \end{bmatrix} \end{matrix} \quad (6)$$

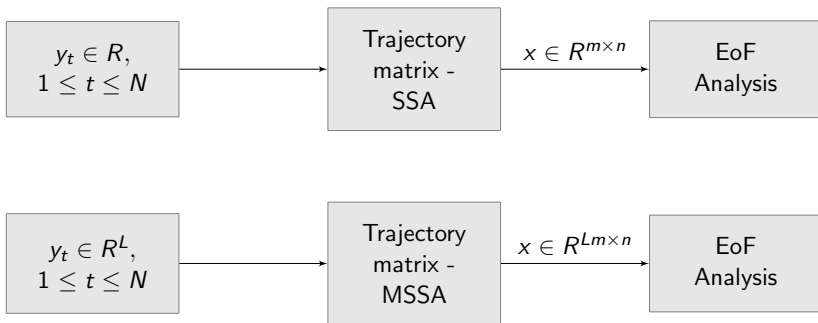
is the trajectory matrix for the MSSA of $\{y_t\}$

- Given $x \in R^{Lm \times n}$, compute the covariance matrix

$$\Sigma = \frac{1}{n}xx^T \in R^{Lm \times Lm} \quad (7)$$

- Σ contains the auto and cross covariances of the components of the vector $\{y_t\}$
- The singular values and vectors of x are closely related to the spectrum of Σ

Summary



- Except for the dimensionality, mathematical analysis of SSA and MSSA are quite similar

References: Use of trajectory matrix

- N.H.Packard, J.P.Crutchfield, J.D.Farmer and R.S.Shaw(1980)" Geometry from a time series", Physical Review Letter, A 45, 712-716
- D.Ruelle (1980)" Strange attractors",Mathematics Intelligence π ,2,37-48
- F.Takens(1981)" Detecting strange attractors in turbulence" in: D.Rand and L.S.Young(Eds) Dynamical Systems and Turbulence, vol 898 of lecture notes in Mathematics, pp 366-381, Springer,Berlin

- H.Kanty and T.Schreiber(1997) Nonlinear Time series Analysis,Cambridge University Press
- H.Tong(1993) Non-linear Time series Analysis: A Dynamical Systems Approach, oxford university press, oxford

MODULE 7.2

SSA: a first look

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Scalar time series

- Let $y = \{y_t | 1 \leq t \leq N\}$ denote the given scalar time series
- Examples: y_t is the global average temperature for year t
- y_t is the total number of deaths due to road accidents in the year t
- y_t is the Facebook stock price at the end of the day t

Trajectory / (lagged) data matrix $x \in R^{m \times n}$

- Let $1 < m < N/2$ and $n = N-m+1$. Then $m < n$
- Define the j^{th} column, x_{*j} of x :

$$x_{*j} = (y_j, y_{j+1}, \dots, y_{j+m-1})^T \quad (1)$$

- The data matrix:

$$x = [x_{ij}] = \frac{1}{\sqrt{n}} [x_{*1}, x_{*2}, \dots, x_{*n}] \in R^{m \times n} \quad (2)$$

where

$$x_{ij} = y_{i+j-1} \quad (3)$$

Example

- $N = 6, m = 3, n = 4, \sqrt{n} = 2$
- $y = \{y_1, y_2, y_3, y_4, y_5, y_6\}$
- $x = \frac{1}{2} \begin{bmatrix} y_1 & y_2 & y_3 & y_4 \\ y_2 & y_3 & y_4 & y_5 \\ y_3 & y_4 & y_5 & y_6 \end{bmatrix} = \begin{bmatrix} x_{11} & x_{12} & x_{13} & x_{14} \\ x_{21} & x_{22} & x_{23} & x_{24} \\ x_{31} & x_{32} & x_{33} & x_{34} \end{bmatrix}$
- Clearly the Hankel structure of x is evident:

$$x_{ij} = y_{i+j-1} \text{ for } i + j = c \quad (4)$$

where $2 \leq c \leq n + m$

A geometric view of x

- It is useful to consider the j^{th} column x_{*j} of x as coordinates of the j^{th} point in R^m
- Thus, the trajectory matrix, x describes the distribution of n points in R^m

Phase space trajectory induced by x

- The n points in R^m can be temporally ordered by connecting the point x_{*j} and $x_{*(j+1)}$ by a line segment for $1 \leq i \leq n$
- The resulting trajectory consisting of $(n-1)$ piece-wise continuous line segments as $n \rightarrow \infty$ provides good amount of qualitative information on the system that generates the original time series

Sample second moment matrix: Version 1

- Let

$$\Sigma(1) = \mathbf{X}\mathbf{X}^T = \frac{1}{n} \begin{bmatrix} x_{*1} & x_{*2} & \dots & x_{*n} \end{bmatrix} \begin{bmatrix} x_{*1}^T \\ x_{*2}^T \\ \vdots \\ x_{*n}^T \end{bmatrix} = \frac{1}{n} \sum_{k=1}^n x_{*k} x_{*k}^T \quad (5)$$

which is symmetric and is the average of the n outer product matrices

Example of Σ

- For the example with $N = 6$, $m = 3$ and $n = 4$:

$$\Sigma(1) = \frac{1}{4} \begin{bmatrix} \sum_{k=1}^4 y_k^2 & \sum_{k=1}^4 y_k y_{k+1} & \sum_{k=1}^4 y_k y_{k+2} \\ \sum_{k=2}^5 y_k y_{k-1} & \sum_{k=2}^5 y_k^2 & \sum_{k=2}^5 y_k y_{k+1} \\ \sum_{k=3}^6 y_k y_{k-2} & \sum_{k=3}^6 y_k y_{k-1} & \sum_{k=3}^6 y_k^2 \end{bmatrix} \quad (6)$$

General structure of Σ : Diagonal elements

- Using $x_{ij} = y_{i+j-1}$:

$$\Sigma_{ii}(1) = \frac{1}{n} \sum_{k=1}^n x_{ik} x_{ik} = \frac{1}{n} \sum_{k=1}^n y_{i+k-1}^2 = \frac{1}{n} \sum_{k=i}^{n+i-1} y_k^2 \quad (7)$$

General structure of Σ : off-diagonal elements



$$\begin{aligned}\Sigma_{ij}(1) &= \frac{1}{n} \sum_{k=1}^n x_{ik} x_{jk} \\ &= \frac{1}{n} \sum_{k=1}^n y_{i+k-1} y_{j+k-1} \\ &= \frac{1}{n} \sum_{k=i}^{n+i-1} y_k y_{k+|j-i|}\end{aligned}\tag{8}$$

- Let N be very large and m be small such that $n \approx N$
- In this case:

$$\Sigma_{ii}(2) = \frac{1}{N} \sum_{k=1}^N y_k^2 = c_0 \quad (9)$$

$$\Sigma_{ij}(2) = \frac{1}{N - |j - i|} \sum_{k=1}^{N - |j - i|} y_k y_{k + |j - i|} = c_{|j - i|} \quad (10)$$

- From (9) and (10)

$$\Sigma(2) \approx \begin{bmatrix} c_0 & c_1 & c_2 & \dots & c_{m-1} \\ c_1 & c_0 & c_1 & \dots & c_{m-2} \\ c_2 & c_1 & c_0 & \dots & c_{m-3} \\ \vdots & \vdots & \vdots & & \vdots \\ c_{m-1} & c_{m-2} & c_{m-3} & \dots & c_0 \end{bmatrix} \quad (11)$$

Symmetric and Toeplitz Structure of $\Sigma(2)$ in (10)

- The second moment matrix $\Sigma(2)$ in (11) in addition to being symmetric, also inherits the Toeplitz structure:

$$\Sigma(2) = [\Sigma_{ij}(2)] \text{ and } \Sigma_{ij}(2) = c_{|j-i|} \quad (12)$$

- That is, elements along the principal diagonal are the same as are those along the diagonals parallel to it

Properties of $\Sigma(1)$ and $\Sigma(2)$

- First choice (Broomhead and King (1986)): $\Sigma(1) \in \mathbb{R}^{m \times m}$ is computed using (7) and (8) which is symmetric
- Second choice (Vautard and Ghil (1989)): $\Sigma(2) \in \mathbb{R}^{m \times m}$ is computed using (9) and (10) which is symmetric and Toeplitz

- The raw/original time series y_t may have non-stationary components such as trend(linear/non-linear)
- Estimate the trend using standard OLS method
- Detrend the given series by subtracting the trend component from the original series

Centering and normalization

- The centered version is obtained by subtracting the over all sample mean from each term of the series
- A centered series is also known as the anomaly series
- Normalized version is obtained by dividing each element by the overall sample standard deviation
- When there is a comparison of methods or different series, it is useful to work with normalized series

- Centering does not remove the trend
- If y_t is a centered second-order stationary series, then for large N

$$\Sigma(2) = \begin{bmatrix} r_0 & r_1 & r_2 & \dots & r_{m-1} \\ r_1 & r_0 & r_2 & \dots & r_{m-2} \\ r_2 & r_1 & r_0 & \dots & r_{m-3} \\ \vdots & \vdots & \vdots & & \vdots \\ r_{m-1} & r_{m-2} & \dots & \dots & r_0 \end{bmatrix} \quad (13)$$

is the symmetric, Toeplitz auto covariance matrix

- If y_t is a normalized second-order stationary series, then for large N

$$\Sigma(2) = \begin{bmatrix} 1 & \rho_1 & \rho_2 & \dots & \rho_{m-1} \\ \rho_1 & 1 & \rho_1 & \dots & \rho_{m-2} \\ \rho_2 & \rho_1 & 1 & \dots & \rho_{m-3} \\ \vdots & \vdots & \vdots & & \vdots \\ \rho_{m-1} & \rho_{m-2} & \rho_{m-3} & \dots & 1 \end{bmatrix} \quad (14)$$

is the symmetric, Toeplitz auto correlation matrix

- Since there is an intrinsic difference between $\Sigma(1)$ and $\Sigma(2)$ especially when N is small, great caution must be exercised in interpreting the results that are dependent on the properties of these matrices

- Another 500 pound Gorilla in the room is the property of the (stochastic) noise, ε_t component that induces randomness to y_t
- In many applications, this noise ε_t is modeled as a white noise: $E(\varepsilon_t) = 0$ $E(\varepsilon_t \varepsilon_s) = 0$ for $t \neq s$
 $= \sigma^2$ for $t = s$

- In many geophysical applications, there is evidence that this corrupting noise ε_t is not a white noise
- Such a red noise is often modeled using a member of the ARMA(p, q) family
- Presence of red noise further complicates the analysis and the conclusions

Challenge of the red noise

- Allen and Smith (1994) provided a comprehensive discussion of the pitfalls in the analysis of time series using the SSA when the noise is not white
- As a further guide to the analysis, they proposed a Monte Carlo SSA that is useful in testing hypothesis relating to the properties of noise

- Let $\Sigma \in R^{m \times m}$ be the matrix computed using either of the two methods using the raw, centered or normalized time series
- Assume that x is of full rank, that is

$$\text{Rank}(x) = \min\{m, n\} = m \quad (15)$$

since $m < n$

- Then, $\Sigma(1)$ is SPD and $\Sigma(2)$ is SPD and Toeplitz

Eigen decomposition of Σ

- Let (λ_i, u_i) , $1 \leq i \leq m$ be the eigenvalue vector pair for Σ , that is, $\Sigma u_i = u_i \lambda_i$ with

$$\lambda_1 \geq \lambda_2 \geq \dots \geq \lambda_m > 0 \quad (16)$$

- Setting $u = [u_1, u_2, \dots, u_m] \in R^{m \times m}$
 $\Lambda = \text{Diag}(\lambda_1, \lambda_2, \dots, \lambda_m)$

-

$$\Sigma u = u \Lambda, uu^T = u^T u = I_m \quad (17)$$

Principal patterns and components

- The eigenvector $u_i, 1 \leq i \leq m$ that constitute an orthogonal basis for R^m are the principal patterns
- The principal components are obtained by projecting the columns of x onto these principal pattern vectors

Principal component matrix: $A \in R^{m \times n}$

- The PC matrix A is given by

$$A = U^T X \in R^{m \times n} \quad (18)$$

- That is,

$$A = [A_{ij}] = \begin{bmatrix} u_1^T \\ u_2^T \\ \vdots \\ u_m^T \end{bmatrix} \begin{bmatrix} x_{*1} & x_{*2} & \dots & x_{*n} \end{bmatrix}$$

-

$$A_{ij} = u_i^T x_{*j} (\text{inner product}) \quad (19)$$

Orthogonality of the rows of A



$$AA^T = u^T x x^T u = u^T \Sigma u = \Lambda \quad (20)$$

- Each of the n columns of A give the coordinates of the n points with respect to new orthonormal coordinate system defined by the columns of u

Example: Trajectory matrix

- Consider a series $\{1, 2, 3, 4, 3, 2, 1\}$ with $N = 7$. Set $m = 3$ and $n = N - m + 1 = 5$
- The trajectory matrix

$$x = \frac{1}{\sqrt{5}} \begin{bmatrix} 1 & 2 & 3 & 4 & 3 \\ 2 & 3 & 4 & 3 & 2 \\ 3 & 4 & 3 & 2 & 1 \end{bmatrix} \in R^{3 \times 5} \quad (21)$$

Example: Σ, u, Λ

- $\Sigma = \Sigma(1) = \begin{bmatrix} 7.8 & 7.6 & 6.2 \\ 7.6 & 8.4 & 7.6 \\ 6.2 & 7.6 & 7.8 \end{bmatrix}$
-

$$\Lambda = \text{Diag}(0.093, 1.600, 22.307) = \text{Diag}(\lambda_3, \lambda_2, \lambda_1) \quad (22)$$

- $u = \begin{bmatrix} -0.4324 & -0.7071 & 0.5595 \\ 0.7912 & -0.0000 & 0.6115 \\ -0.4324 & 0.7071 & 0.5595 \end{bmatrix} = [u_3, u_2, u_1]$

Example: PC matrix A

- $A = \begin{bmatrix} -0.0658 & -0.0987 & 0.2551 & -0.0987 & -0.0658 \\ 0.6325 & 0.6325 & -0.0000 & -0.6325 & -0.6325 \\ 1.5478 & 2.3217 & 2.5952 & 2.3217 & 1.5478 \end{bmatrix}$
- Verify $AA^T = \Lambda$

- In the light of (17), from (18), it is immediate that

$$x = uA = uu^T x \quad (23)$$

- That is, the original matrix, x can be recovered from A using the operation uA

Reconstruction of the signal: $x(s) \in R^{m \times n}$

- In the light of ordering of λ 's in (16), define an integer k such that $1 \leq k < m$ and

$$\sum_{i=1}^k \lambda_i \geq (1 - \beta) \sum_{i=1}^m \lambda_i \quad (24)$$

- Define

$$I_m(m - k) = \text{Diag}(0, 0, \dots, 0, 1, 1, \dots, 1) \quad (25)$$

where 0's are k in number and 1's are $m-k$ in number

- Then

$$\begin{aligned}x(s) &= u[I_m - I_m(m - k)]u^T x \\ &= (\sum_{i=1}^k u_i u_i^T) x\end{aligned}\tag{26}$$

gives the signal component. The noise component is

$$x(n) = x - x(s) = u[I_m(m - k)]u^T x\tag{27}$$

Example - continued

- From (22): $\lambda_1 + \lambda_2 + \lambda_3 = 24$
- $\frac{\lambda_1}{24} = 0.929, \frac{\lambda_2}{24} = 0.0667, \frac{\lambda_3}{24} = 0.00388$
- Verify that $\lambda_1 + \lambda_2 = 23.907 > 0.99(\lambda_1 + \lambda_2 + \lambda_3) = 23.897$ where $\beta = 0.01$

$$\begin{aligned}x(s) &= (u_1 u_1^T + u_2 u_2^T)x \\x(n) &= u_3 u_3^T x\end{aligned}\tag{28}$$

- Verify that (since $n = 5$)

$$\sqrt{5} \times (n) = \begin{bmatrix} 0.9364 & 1.9046 & 2.3467 & 3.9046 & 2.9364 \\ 2.1164 & 3.1747 & 3.5486 & 3.1747 & 2.1164 \\ 2.9364 & 3.9046 & 3.2467 & 1.9046 & 0.9364 \end{bmatrix} \quad (29)$$

is the signal part of the x recovered by the SSA analysis

- The signal recovered in (29) does not inherit the Hankel structure as the original data matrix

- D.S.Broomhead and G.P King (1986) "Extracting qualitative dynamics from experimental data", Physica, D 20, 217-236
- R. Vautard and M.Ghil(1989) "Singular spectrum analysis in nonlinear dynamics with applications to paleoclimatic time series", Physica D, 35, 395-424
- M.R.Allen and L.Smith (1996)"Monte Carlo SSA: Detecting irregular oscillations in the presence of coloured noise", Journal of climate, 93373-3404
- R.Vautard, P.Yiou and M.Ghil (1992), "Singular spectrum analysis: a tool kit for short, noisy and chaotic series", Physica D, 58, 95-126

MODULE 7.3

Singular spectrum Analysis(SSA): A second view

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What is SSA?

- SSA is a powerful tool for the analysis of time series(TS) using ideas from
 - multivariate statistics
 - geometry
 - dynamical system
 - signal processing

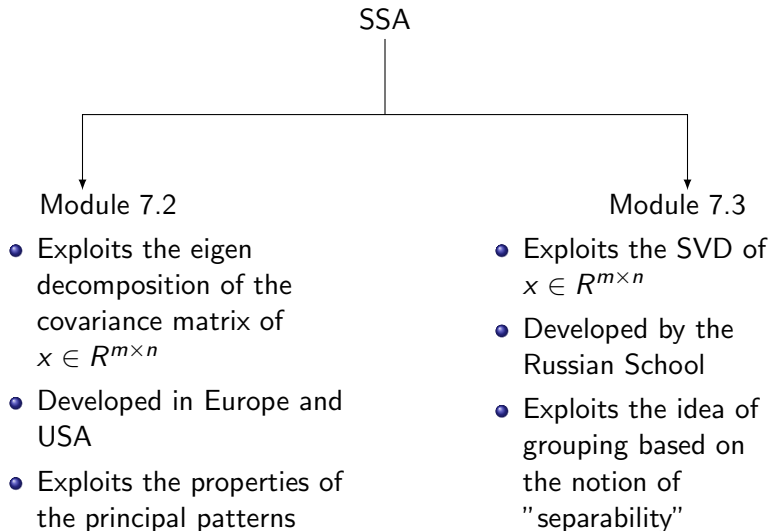
The goal of SSA

- To additively decompose a given TS as a sum of "independent" components that capture the
 - time varying trend
 - oscillatory component
 - noise

Two stages of SSA

- 1 Decomposition stage - two steps
 - 1.1 Embedding
 - 1.2 SVD Analysis
- 2 Reconstruction stage - two steps
 - 2.1 Grouping
 - 2.2 Diagonal averaging

Comments on the two approaches to SSA



- N. Golyandina, V. Nekrutkin and A. Zhigljavsky (2001) Analysis of Time series structure: SSA and Related Techniques, Chapman and Hall/ CRC 305 pages
- J. B. Elsner and A. A. Tsonis (1996) Singular Spectrum Analysis: A new tool in Time Series Analysis, Plenum Press, NY, 164 pages

- H. Hassani (2007) "Singular Spectrum Analysis: Methodology and Comparison", Journal of Data Science, Vol 5, pp239 - 257
- Provides a nice summary of the key steps in the SSA based algorithm and a good comparison with other classical techniques from TSA.

- Let $\{y_t | 1 \leq t \leq N\}$ be the given scalar TS, where N is large.
- Let $1 < m < N/2$ and $n = N - m + 1$
- Define a lagged (Column) vector for $1 \leq j \leq n$:

$$x_{*j} = (y_j, y_{j+1}, \dots, y_{j+m-1})^T \in R^m \quad (1)$$

Step 1.1 Embedding

- The given series is split into n lagged column vectors x_{*j} , $1 \leq j \leq n$
- Define the trajectory matrix: $x \in R^{m \times n}$

$$x = [x_{*1}, x_{*2}, \dots, x_{*n}] \in R^{m \times n} \quad (2)$$

Example

- Let $N = 7$, $m = 3$, $n = 5$
- $y = \{y_1, y_2, y_3, y_4, y_5, y_6, y_7\}$
- $x = \begin{bmatrix} y_1 & y_2 & y_3 & y_4 & y_5 \\ y_2 & y_3 & y_4 & y_5 & y_6 \\ y_3 & y_4 & y_5 & y_6 & y_7 \end{bmatrix} \in \mathbb{R}^{3 \times 5}$
- Verify: $x_{ij} = y_{i+j-1}$ and matrix x inherits the Hankel structure
- elements across the anti-diagonals ($i + j = k$, a constant)
are the same

Step 1.2 : SVD of $x \in \mathbb{R}^{m \times n}$

- Since $1 < m < N/2$ and $n = N - m + 1$, we have $m < n$
- Let x be full rank matrix: $\text{Rank}(x) = \min\{m, n\} = m$
- Then (xx^T) is SPD and

$$(xx^T)u = u\Lambda \quad (3)$$

be the eigen decomposition of the smaller Gramian xx^T

- $u = [u_1, u_2, u_3, \dots, u_m] \in \mathbb{R}^{m \times m}$

$$uu^T = u^T u = I_m \quad (4a)$$

- Let $\Lambda = \text{Diag}(\lambda_1, \lambda_2, \dots, \lambda_m)$:

$$\lambda_1 \geq \lambda_2 \geq \lambda_m \geq \lambda_m \quad (4b)$$

- Define

$$v_i = \frac{1}{\sqrt{\lambda_i}} x^T u_i \quad (5)$$

- $v = \{v_1, v_2, v_3, \dots, v_m\} \in \mathbb{R}^{n \times m}$ and

$$v^T v = I_m \quad (6)$$

$$(x^T x)v = v\Lambda \quad (7)$$

- A dual relation:

$$x^T u_i = v_i \lambda_i^{1/2} \quad \text{and} \quad x v_i = u_i \lambda_i^{1/2} \quad (8)$$

- From (5):

$$x^T u = v \Lambda^{1/2} \quad \text{or} \quad x^T = v \Lambda^{1/2} u^T \quad (9)$$

- From (8): SVD of x :

$$x = u \Lambda^{1/2} v^T \quad (10)$$

- From (10):

$$x = \sum_{i=1}^m \lambda_i u_i v_i^T \quad (11)$$

- (λ_i, u_i, v_i) - i^{th} eigen triple of x , $1 \leq i \leq m$
- These m eigen triples are the basic building blocks for the reconstruction phase used in identifying the trend, oscillatory and noise components.

- Recall:

$$\begin{aligned}\|\lambda_i u_i v_i^T\|_F^2 &= \lambda_i^2 \\ \|\lambda_i u_i v_i^T + \lambda_j u_j v_j^T\|_F^2 &= \lambda_i^2 + \lambda_j^2\end{aligned}\tag{12}$$

- From (11): The total energy in x :

$$\|x\|_F^2 = \sum_{i=1}^m \lambda_i^2\tag{13}$$

Inherent optimality of SVD

- Let, for $1 \leq r \leq m$,

$$x(r) = \sum_{i=1}^r \lambda_i u_i v_i^T \in R^{m \times n} \quad (14)$$

- From (12):

$$\|x(r)\|_F^2 = \sum_{i=1}^r \lambda_i^2 \quad (15)$$

- Hence

$$\|x - x(r)\|^2 = \sum_{i=r+1}^m \lambda_i^2 \quad (16)$$

- In view of the ordering of the λ 's in (4b), it is immediate that $x(r)$ is the best rank r approximation of x

Step 2.1: Grouping of indices

- Let $[m] = \{1, 2, 3, \dots, m\}$
- Let for some $1 \leq p \leq m$, let $s_p \subset [m]$
- Grouping: Let $\{s_1, s_2, \dots, s_p\}$ be partition of $[m]$

$$s_i \cap s_j = \emptyset, \quad \bigcup_{j=1}^p s_j = [m]$$

Resulting decomposition of x

- Consider the i^{th} group s_i of $|s_i|$ indices
- Define a matrix $x(i) \in R^{m \times n}$ of $\text{Rank}(x(i)) = |s_i|$ as

$$x(i) = \sum_{j \in s_i} \lambda_j u_j v_j^T \quad (17)$$

- Clearly,

$$x = \sum_{i=1}^p x(i) \quad (18)$$

where p is the number of sets in the partition of $[m]$

Step 2.2 - Diagonal averaging

- The matrix $x(i)$ in (17) need not be a Hankel matrix
- Hankelization of $x(i)$ relates to creating a Hankel matrix, $\bar{x}(i) = \mathcal{H}x(i)$ for the $x(i)$
- This is the done by replacing each anti-diagonal in $x(i)$ by the average of the elements in that diagonal in $x(i)$

Example of Hankelization

- Let

$$A = \begin{bmatrix} a_{11} & a_{12} & a_{13} & a_{14} \\ a_{21} & a_{22} & a_{23} & a_{24} \\ a_{31} & a_{32} & a_{33} & a_{34} \end{bmatrix} \quad (19)$$

$$Z = \begin{bmatrix} z_1 & z_2 & z_3 & z_4 \\ z_2 & z_3 & z_4 & z_5 \\ z_3 & z_4 & z_5 & z_6 \end{bmatrix} = \mathcal{H}A \quad (20)$$

Example (continued)

- Then

$$\begin{aligned}z_1 &= a_{11} \\z_2 &= \frac{1}{2}(a_{12} + a_{21}) \\z_3 &= \frac{1}{3}(a_{13} + a_{22} + a_{31}) \\z_4 &= \frac{1}{3}(a_{14} + a_{23} + a_{32}) \\z_5 &= \frac{1}{2}(a_{24} + a_{42}) \\z_6 &= a_{34}\end{aligned}\tag{21}$$

Algorithm for Hankelization

- Let $A = [a_{ij}] \in R^{m \times n}$ matrix
- Let $B = \mathcal{H}A = [B_{ij}] \in R^{m \times n}$

$$B_{ij} = \begin{cases} \frac{1}{s-1} \sum_{p=1}^{s-1} a_{p,s-p} & \text{for } 2 \leq s \leq m-1 \\ \frac{1}{m} \sum_{p=1}^m a_{p,s-p} & \text{for } m \leq s \leq n+1 \\ \frac{1}{m+n-s+1} \sum_{p=s-n}^L a_{p,s-p} & \text{for } n+2 \leq s \leq n+m \end{cases} \quad (22)$$

Hankelize the sum in (18)

- Recall that
 - $\mathcal{H}(A+B) = \mathcal{H}(A) + \mathcal{H}(B)$
 - $\mathcal{H}(A) = A$ if A is a Hankel matrix.
- Operating both sides of (18) by \mathcal{H} :

$$x = \sum_{i=1}^p \mathcal{H}(x(i)) = \sum_{i=1}^p \bar{x}(i) \quad (23)$$

Decomposition of y_t

- Recall: Unique relation between time series and Hankel trajectory matrix
- Let $\{y_t | 1 \leq t \leq N\}$ be the TS for x and $\{\bar{y}_t(i) | 1 \leq t \leq N\}$ be the TS for (i)
- Then (23) becomes:

$$y_t = \sum_{i=1}^p \bar{y}_t(i) \text{ for each } 1 \leq t \leq N \quad (24)$$

- The results of the first stage of decomposition critically depends on the window length m
- If we already know that the given TS has periodic components with period T - say using spectral analysis, then L is proportional to this period.
- In any case L must be large but less than $N/2$, half the length of the series

Comments: challenge of grouping

- Of the four steps involved, embedding, SVD and Hankelization are quite algorithmic and can be easily implemented
- Grouping is the most demanding part of this approach to SSA
- Algorithm for optimal grouping is still evasive, extra / supplementary information about the series could be used as a guide to grouping

Guide to grouping - Scree plot

- Identifying breaks in the scree plot - plot of the eigenvalues Vs its rank, could help identify signals from noise
- Presence of white noise corresponds to a constant lower ceiling in the scree plot
- It is known that a harmonic component produces a pair of very close singular values

Guide to grouping - Periodogram

- Compute and plot the periodogram for the TS and identify the frequency with spikes in the spectrum
- We can then search for the eigen triple whose frequencies coincide with those identified by the spikes

Role of separability in grouping

- To partially automate the grouping operation, a notion of "separability" based on a "weighted correlation" is introduced
- Two series are separated from each other if their weighted correlation is low

Choice of Weights

- Let f_t and g_t be the two time series with $1 \leq t \leq N$
- Let $1 \leq m \leq N/2$ and $n = N - m + 1$
- Let

$$w_k = \min\{k, m, N - k + 1\} \text{ for } 1 \leq k \leq N \quad (25)$$

Example

- Let $N = 6$, $m = 3$, $n = 4$
- Then: $w_k = \min\{k, 3, 7-k\}$
- Clearly,

k	1	2	3	4	5	6
w_k	1	2	3	3	2	1

Weighted inner product of two TS

- Let $\{f_t\}$ and $\{g_t\}$ be the two series for $1 \leq t \leq N$
- The weighted inner product between $\{f_t\}$ and $\{g_t\}$ is

$$\langle f, g \rangle_w = \sum_{k=1}^N w_k f_k g_k \quad (26)$$

Norm of a given TS: $||f||_w$

•

$$||f||_w^2 = \langle f, f \rangle_w = \sum_{k=1}^N w_k f_k^2 \quad (27)$$

Weighted correlation between f_t and g_t



$$\rho_w(f, g) = \frac{\langle f, g \rangle_w}{\|f\|_w \|g\|_w} \quad (28)$$

is the weighted correlation

- The idea is: if $\rho_w(f, g)$ is small, then the series $\{f_t\}$ and $\{g_t\}$ are almost w-orthogonal and do not share common information

- Given the m -eigen triples $(\lambda_i^{1/2}, u_i, v_i)$ first compute

$$x(i) = \lambda_i^{1/2} u_i v_i^T \quad 1 \leq i \leq m \quad (29)$$

- Let $\bar{x}(i)$ be the Hankelized $x(i)$ and $\{\bar{y}_t(i) | 1 \leq t \leq N\}$ be the corresponding TS

Compute their W-Correlations

- Given $\bar{x}_t(i)$ for $1 \leq i \leq m$, compute

$$R = [R_{ij}] \in R^{m \times m}$$

where

$$R_{ij} = \rho_w(\bar{x}_t(i), \bar{x}_t(j)) \quad (30)$$

- By examining the off-diagonal entries of this symmetric matrix, we may isolate the correlated pairs from others which in turn provide useful information for grouping

MODULE 17.1

Canonical Correlation Analysis (CCA)

Basic Theory

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- Given a random field $x \in R^m$ (indexed by points in a spatial domain $S \subseteq R^k, k = 1, 2, 3$), the goal of PCA is to decompose x as a linear combination of intrinsic spatial patterns which are related to the eigenvectors of $\text{cov}(x)$ where the variance of the random coefficients in the linear combination are directly related to the corresponding eigenvalues of $\text{cov}(x)$
- If the eigenvalues of $\text{cov}(x)$ are distinct and decreasing, this additive decomposition is often used as basis for "removing the chaff(noise) from the corn(signal)"

- Given two random fields $x_1 \in R^{m_1}$ and $x_2 \in R^{m_2}$ (defined over two spatial domains S_1 and S_2) the goal of CCA is to express each field as a linear combination of its own spatial patterns in such a way that the corresponding spatial patterns exhibit maximum correlation
- An example : x_1 average SST over equatorial pacific during a month and x_2 could be the average rain fall across the USA in that same month

Second-order properties of two random fields

- Let $x_1 \in R^{m_1}$ and $x_2 \in R^{m_2}$ be two random fields of interest



$$E(x) = \mu_1 \quad \text{and} \quad E(y) = \mu_2 \quad (1)$$

$$\text{cov}(x) = E[(x - \mu_1)(x - \mu_1)^T] = \Sigma_{11} \in R^{m_1 \times m_1} \quad (2)$$

$$\text{cov}(y) = E[(y - \mu_2)(y - \mu_2)^T] = \Sigma_{22} \in R^{m_2 \times m_2} \quad (3)$$

- Assume that Σ_{11} and Σ_{22} are SPD

Covariance between x_1 and x_2



$$\begin{aligned} \text{cov}(x_1, x_2) &= E[(x_1 - \mu_1)(x_2 - \mu_2)^T] \\ &= \Sigma_{12} \in R^{m_1 \times m_2} \end{aligned} \quad (4)$$



$$\begin{aligned} \text{cov}(x_2, x_1) &= E[(x_2 - \mu_2)(x_1 - \mu_1)^T] \\ &= \Sigma_{21} = \Sigma_{12}^T \in R^{m_2 \times m_1} \end{aligned} \quad (5)$$

Normalize x_1 and x_2

- Let $D_1 \in R^{m_1 \times m_1} =$ Diagonal matrix with the diagonal elements of Σ_{11}
 $D_2 \in R^{m_2 \times m_2} =$ Diagonal matrix with the diagonal elements of Σ_{22}
- Normalize x_1 and x_2 :
- Then

$$\begin{aligned}\hat{x}_1 &= D_1^{-1/2}(x_1 - \mu_1) \\ \hat{x}_2 &= D_2^{-1/2}(x_2 - \mu_2)\end{aligned}\tag{6}$$

are the centered and normalized versions of x_1 and x_2



$$\begin{aligned} \text{cor}(x_1, x_2) &= \text{cov}(\hat{x}_1, \hat{x}_2) = E[\hat{x}_1(\hat{x}_2)^T] \\ &= D_1^{-1/2} \Sigma_{12} D_2^{-1/2} \in R^{m_1 \times m_2} \end{aligned} \quad (7)$$

Spatial patterns for random fields x_1 and x_2

- Let

$$F = [f_1, f_2, \dots, f_{m_1}] \in R^{m_1 \times m_1} \quad (8)$$

be a matrix whose linearly independent column vectors span R^{m_1} and constitute m_1 distinct spatial patterns for x_1

- Similarly, let

$$G = [g_1, g_2, \dots, g_{m_2}] \in R^{m_2 \times m_2} \quad (9)$$

be that for x_2

Expansion of the random fields using spatial patterns

- By resolving x_1 along each of the spatial patterns f_i , it follows that

$$x_1 = \sum_{i=1}^{m_1} (x_1^T f_i) f_i = \sum_{i=1}^{m_1} (\alpha_i f_i) \quad (10)$$

where the coefficients $\alpha_i = (x_1^T f_i)$ of the linear combination in (10) are random variables since x_1 is random

- Similarly, resolving x_2 along each of the spatial patterns g_j , it follows that

$$x_2 = \sum_{j=1}^{m_2} (x_2^T g_j) g_j = \sum_{j=1}^{m_2} (\beta_j g_j) \quad (11)$$

- Here again, the coefficients $\beta_j = (x_2^T g_j)$ inherit their randomness from that of x_2

Statement of the problem

- The goal is to find the spatial patterns $\{f_i\}$ and $\{g_j\}$ such that

$$\text{cor}(\alpha_1, \beta_1) \geq \text{cor}(\alpha_2, \beta_2) \geq \cdots \geq \text{cor}(\alpha_k, \beta_k) > 0 \quad (12)$$

and

$$\text{cor}(\alpha_i, \beta_i) \text{ is the maximum for } (f_i, g_i) \quad (13)$$

for each $i = 1, 2, \dots, k \leq \min\{m_1, m_2\}$

Expression for typical correlation

- Let $f \in R^{m_1}$ and $g \in R^{m_2}$ be two typical spatial patterns for x_1 and x_2
- Define:

$$\alpha = x_1^T f \quad \text{and} \quad \beta = x_2^T g \quad (14)$$

- Then

$$\rho = \text{cor}(\alpha, \beta) = \frac{\text{cov}(\alpha, \beta)}{[\text{var}(\alpha)\text{var}(\beta)]^{1/2}} \quad (15)$$



$$\begin{aligned}\text{cov}(\alpha, \beta) &= E[(x_1 - \mu_1)^T f (x_2 - \mu_2)^T g] \\ &= f^T \text{cov}(x_1, x_2) g \\ &= f^T \Sigma_{12} g = f^T \Sigma_{21}^T g\end{aligned}\tag{16}$$

Expression for $var(\alpha)$ and $var(\beta)$



$$\begin{aligned} Var(\alpha) &= E[(x_1 - \mu)^T f (x_1 - \mu_1)^T f] \\ &= f^T E[(x_1 - \mu)(x_1 - \mu_1)^T] f \\ &= f^T \Sigma_{11} f > 0 \end{aligned} \tag{17}$$

• Similarly:

$$Var(\beta) = g^T \Sigma_{22} g > 0 \tag{18}$$

- Substituting (16), (17) and (18) in (15):

$$\rho = \text{cor}(\alpha, \beta) = \frac{f^T \Sigma_{12} g}{(f^T \Sigma_{11} f)^{1/2} (g^T \Sigma_{22} g)^{1/2}} \quad (19)$$

- Goal is to find the pattern pair (f,g) that maximizes the right hand side of (19)

Scale invariance of ρ

- Let a, b be two positive real constants
- It can be verified

$$\rho = \text{cov}(\alpha, \beta) = \text{cov}(a\alpha, b\beta) \quad (20)$$

- That is, ρ is invariant under the scaling of the spatial patterns f and g

Normalized spatial patterns

- With out loss of generality, assume that the patterns f and g are normalized:

$$f^T \Sigma_{11} f = 1 \quad \text{and} \quad g^T \Sigma_{22} g = 1 \quad (21)$$

- Define

$$\bar{f} = \Sigma_{11}^{1/2} f \quad \text{and} \quad \bar{g} = \Sigma_{22}^{1/2} g \quad (22)$$

- Then

$$\bar{f}^T \bar{f} = 1 = \bar{g}^T \bar{g} \quad (23)$$

New expression for ρ

- Substituting (21), (22) and (23) in (19), we get a bilinear form ρ given by:

$$\begin{aligned}\rho &= \bar{f}^T (\Sigma_{11}^{-1/2} \Sigma_{12} \Sigma_{22}^{-1/2}) \bar{g} \\ &= \bar{f}^T A \bar{g}\end{aligned}\tag{24}$$

where

$$A = \Sigma_{11}^{-1/2} \Sigma_{12} \Sigma_{22}^{-1/2} \in R^{m_1 \times m_2}\tag{25}$$

- The problem is to find $\bar{f} \in R^{m_1}$ and $\bar{g} \in R^{m_2}$ that maximizes ρ in (24) under the constraints (23)

- Consider the lagragian

$$L(\bar{f}, \bar{g}, a, b) = \bar{f}^T A \bar{g} + a(1 - \bar{f}^T f) + b(1 - \bar{g}^T g) \quad (26)$$

where a and b are the (scalar) unknown Lagrangian multipliers

Conditions for the maximum



$$\nabla_{\bar{f}} L = A\bar{g} - 2a\bar{f} = 0 \quad (27)$$

$$\nabla_{\bar{g}} L = A^T \bar{f} - 2b\bar{g} = 0 \quad (28)$$

$$\nabla_a L = 1 - \bar{f}^T \bar{f} = 0 \quad (29)$$

$$\nabla_b L = 1 - \bar{g}^T \bar{g} = 0 \quad (30)$$

- Conditions (29) and (30) follows from (23)
- Optimal (\bar{f}, \bar{g}) are obtained as the solution of (27) and (28) written as

$$A\bar{g} = 2a\bar{f} \quad (31)$$

$$A^T \bar{f} = 2b\bar{g} \quad (32)$$

A related eigenvalue problem

- Substituting (32) in (31) and vice versa:

$$AA^T \bar{f} = 4ab\bar{f} \quad (33)$$

$$A^T A \bar{g} = 4ab\bar{g} \quad (34)$$

- That is, setting $\lambda = 4ab$, it follows that (λ, \bar{f}) is an eigen pair of AA^T and (λ, \bar{g}) is an eigen pair of $A^T A$

A first look at the solution to the problem

- Indeed, the pair (\bar{f}, \bar{g}) of spatial patterns that maximizes ρ in (24) are given by the eigenvectors of the two Grammian matrices AA^T and $A^T A$ respectively where $A \in R^{m_1 \times m_2}$ is the matrix of the bilinear form in (25)
- Also, notice that the eigenvalues λ are related to the product of the Lagrangian multipliers a and b

Properties of $A^T A$ and AA^T

- For concreteness, let $m_2 < m_1$ and A be of full rank
- $\text{Rank}(A) = \text{Rank}(A^T) = \min\{m_1, m_2\} = m_2$
- $\text{Rank}(AA^T) = \text{Rank}(A^T A) = m_2$
- Thus, smaller Gramian $A^T A$ is of full rank, and SPD but longer Gramian AA^T is rank deficient and symmetric positive semi-definite

Eigen structure of $A^T A$

- Let (λ_i, \bar{g}_i) , $1 \leq i \leq m_2$ be the eigen pairs of $A^T A$. That is, $(A^T A)\bar{g}_i = \bar{g}_i \lambda_i$, $\bar{g}_i \in R^{m_2}$ where

$$\lambda_1 > \lambda_2 > \lambda_3 \cdots > \lambda_{m_2} > 0 \quad (35)$$

- Let

$$\begin{aligned} \bar{G} &= [\bar{g}_1, \bar{g}_2, \dots, \bar{g}_{m_2}] \in R^{m_2 \times m_2} \\ \Lambda &= \text{Diag}(\lambda_1, \lambda_2, \lambda_3 \dots, \lambda_{m_2}) \end{aligned} \quad (36)$$

- Then :

$$\begin{aligned} (A^T A)\bar{G} &= \bar{G}\Lambda \\ \bar{G}^T \bar{G} &= \bar{G}\bar{G}^T = I_{m_2} \end{aligned} \quad (37)$$

Eigen structure of AA^T

- Define

$$\bar{f}_i = \frac{1}{\sqrt{\lambda_i}} A \bar{g}_i \in R^{m_1}, 1 \leq i \leq m_2 \quad (38)$$

- Verify:

$$(AA^T)\bar{f}_i = \lambda_i \bar{f}_i \quad (39)$$

- Set: $\bar{F} = [\bar{f}_1, \bar{f}_2, \dots, \bar{f}_{m_2}] \in R^{m_1 \times m_2}$

- Then

$$\begin{aligned} (AA^T)\bar{F} &= \bar{F}\Lambda \\ \bar{F}^T \bar{F} &= I_{m_2} \end{aligned} \quad (40)$$

- Rewriting (38)

$$A\bar{g}_i = \bar{f}_i\sqrt{\lambda_i}, \quad 1 \leq i \leq m_2$$

- That is,

$$A\bar{G} = \bar{F}\Lambda^{1/2} \quad \text{or} \quad A = \bar{F}\Lambda^{1/2}\bar{G}^T \quad (41)$$

is called the singular value decomposition of A

- λ_i 's are the eigenvalues of $A^T A$ and AA^T and $\sqrt{\lambda_i}$ are the singular values of A

- From (24): For $1 \leq i \leq m_2$,

$$\rho_i = \bar{f}_i^T A \bar{g}_i \quad (42)$$

- Substitute for A using (41):

$$\rho_i = \bar{f}_i^T \bar{F} \Lambda \bar{G}^T g_i \quad (43)$$

Simplification of (43)



$$\begin{aligned}\bar{f}_i^T F &= [\bar{f}_i^T \bar{f}_1, \bar{f}_i^T \bar{f}_2, \dots, \bar{f}_i^T \bar{f}_i, \dots, \bar{f}_i^T \bar{f}_{m_2}] \\ &= (0, 0, \dots, 1, \dots, 0) = e_i^T \in R^{m_2}\end{aligned}\quad (44)$$

the i^{th} unit vector in R^{m_2}



$$\bar{G}^T g_i = \begin{bmatrix} \bar{g}_1^T g_i \\ \bar{g}_2^T g_i \\ \vdots \\ \bar{g}_i^T g_i \\ \vdots \\ \bar{g}_{m_2}^T g_i \end{bmatrix} = \begin{bmatrix} 0 \\ 0 \\ \vdots \\ 1 \\ \vdots \\ 0 \end{bmatrix} = e_i \quad (45)$$

Maximum value of ρ_i along (\bar{f}_i, \bar{g}_i)

- Substituting (44) and (45) in (43):

$$\rho_i = e_i^T \Lambda^{1/2} e_i = \sqrt{\lambda_i} \quad (46)$$

the i^{th} singular value of A

- Since $\lambda_i > 0$ and are ordered as in (35), it follows that ρ_i attains its i^{th} maximum value along the spatial pattern pair (\bar{f}_i, \bar{g}_i)

Singular value and Lagrangian multipliers

- From (33)-(34):

$$\lambda = 4ab \quad \text{or} \quad \sqrt{\lambda} = 2\sqrt{ab}$$

that is singular values are proportional to the square root of the product of the two Lagrangian multipliers

- When $m_1 = m_2$ from (31)-(32), we get

$$0 = \bar{f}^T A \bar{g} - \bar{g}^T \bar{A} \bar{f} = 2\bar{f}^T \bar{g}(a - b) \quad (47)$$

it follows that $a = b$ since $\bar{f}^T \bar{g} \neq 0$

- The books
W.Hardle and L.Simar(2003) Applied Multivariate Statistical Analysis, Springer Verlag
- H.Von Storch and F.Zwiers(1999) Statistical Analysis in climate Research, Cambridge University Press, contain very good introduction to CCA.