

# STATE SPACE MODELS OF STOCHASTIC SYSTEMS

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# STATE SPACE MODELS OF DETERMINISTIC SYSTEMS

$\{\mathbf{y}(t)\}$   $m$ -dimensional output       $\{\mathbf{u}(t)\}$   $r$ -dimensional input

$$\begin{cases} \mathbf{x}(t+1) = A\mathbf{x}(t) + B\mathbf{u}(t) & \mathbf{x}(t_0) = \mathbf{x}_0 \\ \mathbf{y}(t) = C\mathbf{x}(t) + D\mathbf{u}(t) & , \quad t \geq t_0 \end{cases}$$

$\{\mathbf{x}(t)\}$   $n$ -dimensional **state** of the system.  $A, B, C, D$  COSTANT system parameters. May be time-varying (but known).

REACHABILITY       $\text{rank} [BAB \dots, A^{n-1}B] = n$

OBSERVABILITY       $\text{rank} \begin{bmatrix} C \\ CA \\ \vdots \\ CA^{n-1} \end{bmatrix} = n$

**THEOREM 1** *The model is MINIMAL ( $\dim \mathbf{x}(t)$  as small as possible) iff the system is **reachable and observable**.*

# TRANSFER FUNCTION

$$F(z) = C[zI - A]^{-1}B + D$$

$m \times r$  matrix. Elements  $F_{ij}(z)$  are **proper rational functions of  $z$**

$$F(z) = C \frac{\text{Adj}[zI - A]}{\det[zI - A]} B + D$$

Characteristic Polynomial:

$$\Delta(z) = \det[zI - A]$$

$$\Delta(p_k) = 0 \quad \Leftrightarrow \quad \text{POLES OF } W(z)$$

BIBO STABILITY (Bounded inputs  $\Rightarrow$  Bounded outputs)  $\Leftrightarrow |p_k| < 1$

## Z- (Fourier) TRANSFORM

$$\hat{f}(z) := \sum_{t=-\infty}^{+\infty} f(t)z^{-1}, \quad z \in \mathbb{C}$$

Frequency response  $\theta = \omega T$ ,  $F(e^{j\theta})$ .

# SIMILARITY TRANSFORMATION

Warning: use of incongruous units may lead to ill-conditioned models.

State can be transformed  $\hat{x}(t) := Tx(t)$ ;  $T$   $n \times n$  nonsingular matrix

$$\hat{A} = T^{-1}AT, \quad \hat{B} = T^{-1}B \quad \hat{C} = CT$$

$$\begin{cases} \hat{\mathbf{x}}(t+1) = \hat{A}\hat{\mathbf{x}}(t) + \hat{B}\hat{\mathbf{u}}(t) & \hat{\mathbf{x}}(t_0) = \hat{\mathbf{x}}_0 \\ \mathbf{y}(t) = \hat{C}\hat{\mathbf{x}}(t) + D\mathbf{u}(t) & , \quad t \geq t_0 \end{cases}$$

has the same transfer function

$$F(z) = \hat{C} [zI - \hat{A}]^{-1} \hat{B} + D$$

# SINGULAR VALUE DECOMPOSITION (SVD)

**THEOREM 2** *Let  $A \in \mathbb{R}^{m \times p}$  of rank  $n \leq \min(m, p)$ . Can find two orthogonal matrices  $U \in \mathbb{R}^{m \times m}$  and  $V \in \mathbb{R}^{p \times p}$  and positive numbers  $\{\sigma_1 \geq \dots \geq \sigma_n\}$ , the **singular values** of  $A$ , so that*

$$A = U\Delta V^\top \quad \Delta = \begin{bmatrix} \Sigma & 0 \\ 0 & 0 \end{bmatrix}, \quad \Sigma = \text{diag}\{\sigma_1, \dots, \sigma_n\}$$

Full-rank factorization of  $A$

$$A = [u_1, \dots, u_n] \Sigma [v_1, \dots, v_n]^\top := U_n \Sigma V_n^\top$$

where  $U_n, V_n$  submatrices of  $U, V$  keeping only the first  $n$  columns

$$U_n^\top U_n = I_n = V_n^\top V_n$$

$$Ax = \sum_{k=1}^n u_k \sigma_k \langle v_k, x \rangle$$

$U = [u_1, \dots, u_m]$  = normalized eigenvectors of  $AA^\top$ ;

$V := [v_1, \dots, v_p]$  normalized eigenvectors of  $A^\top A$ .

$\{\sigma_1^2 \geq \dots \geq \sigma_n^2\}$  (non zero) eigenvalues of  $AA^\top$  (or of  $A^\top A$ ).



# MATRIX NORMS

**2- norm** of  $A \in \mathbb{R}^{m \times p}$  Let  $\|x\|$  be the Euclidean norm.

$$\|A\|_2 := \sup_{x \neq 0} \frac{\|Ax\|}{\|x\|} = \sigma_1 \quad (\sigma_{MAX}(A))$$

The **Frobenius norm**  $\|A\|_F$  is

$$\|A\|_F^2 = \sum_{i,j} a_{i,j}^2 = \sigma_1^2 + \dots + \sigma_n^2$$

Condition number

$$\kappa(A) = \frac{\|A\|_2}{\|A^{-1}\|_2} = \frac{\sigma_{MAX}(A)}{\sigma_{MIN}(A)}$$

# USEFUL FEATURES OF SVD

Range and Nullspace of  $A$ :

$$\text{Im}(A) = \text{Im}(U_n), \quad Ax = 0 \Leftrightarrow x \in \text{span}([v_{n+1}, \dots, v_p]) = \text{Im} V_n^\perp$$

Approximation properties

$$A_k := \sum_{i=1}^k \sigma_i u_i v_i^\top, \quad k \leq n$$

**is the best approximant of rank  $k$  of  $A$**

$$\min_{\text{rank}(B)=k} \|A - B\|_2 = \|A - A_k\|_2 = \sigma_{k+1}$$

$$\min_{\text{rank}(B)=k} \|A - B\|_F^2 = \|A - A_k\|_F^2 = \sigma_{k+1}^2 + \dots + \sigma_w^2$$

# BALANCING

Assume: Eigenvalues of  $A$  strictly less than 1:  $|\lambda(A)| < 1$ ,  $(A, B)$  reachable +  $(C, A)$  observable.

$$\Pi := \sum_0^{+\infty} A^k B B^\top (A^\top)^k, \quad \Omega := \sum_0^{+\infty} (A^\top)^k C^\top C A^k$$

**Reachability and Observability Gramians**, solutions of the dual LYAPUNOV EQUATIONS

$$\begin{aligned}\Pi &= A \Pi A^\top + B B^\top \\ \Omega &= A^\top \Omega A + C^\top C\end{aligned}$$

**THEOREM 3** *Assume the eigenvalues of  $A$  are strictly less than 1. **System is reachable** if and only if  $\Pi > 0$ . **System is observable** if and only if  $\Omega > 0$ . If both hold the model is MINIMAL ( $\dim \mathbf{x}(t)$  as small as possible).*

# INTERPRETATION OF THE GRAMIANS

Assume we can use only finite energy controls:

$$\|u\|_2^2 := \sum_{k=0}^{+\infty} u(k)^\top u(k) \leq 1$$

Energy of the state  $\mathbf{x}(0) = \sum_0^{+\infty} A^k B \mathbf{u}(-k) := \mathbb{R} \mathbf{u}$

$$\max_{\|\mathbf{u}\| \leq 1} \frac{\|\mathbf{x}(0)\|^2}{\|\mathbf{u}\|^2} = \max_{\|\mathbf{u}\| \leq 1} \frac{\langle \mathbf{u}, \mathbb{R}^* \mathbb{R} \mathbf{u} \rangle}{\|\mathbf{u}\|^2} = \|\mathbb{R}^* \mathbb{R}\|_2 = \|\Pi\|_2 = \lambda_{\max}^2(\Pi)$$

Diagonalize:

$$U_c^\top \Pi U_c \Rightarrow \text{diag}\{\lambda_{c,1}^2, \dots, \lambda_{c,n}^2\} \quad \lambda_{c,1}^2 \geq \dots \geq \lambda_{c,n}^2 > 0$$

Change coordinates  $\mathbf{x}_c(0) := U_c^\top \mathbf{x}(0)$

Along the  $k$ -th eigenvector the maximum energy gain is  $\lambda_{c,k}^2$

Energy ratios of the (orthogonal) state components

$$\frac{\|x_{1c}(0)\|}{\|x_{nc}(0)\|} = \frac{\lambda_{1c}}{\lambda_{nc}} \geq \dots \geq \frac{\|x_{n-1,c}(0)\|}{\|x_{nc}(0)\|} = \frac{\lambda_{n-1,c}}{\lambda_{nc}}$$

$\frac{\lambda_{1c}}{\lambda_{nc}}$  may be large: the effect of the input on certain directions in the state space nearly invisible  $\Rightarrow$  BAD CONDITIONING!

# INTERPRETATION OF THE GRAMIANS (Cont.)

Dual meaning of the observability Gramian: Maximal  $L^2$ -energy ( $\|\mathbf{y}\|_2$ ) of the output  $\mathbf{y}(t) = CA^t\mathbf{x}(0)$  for  $\|\mathbf{x}(0)\| \leq 1$ : maximum singular value of  $\Omega$ .

Diagonalization:

$$U_o^\top \Omega U_o \Rightarrow \text{diag}\{\lambda_{o,1}^2, \dots, \lambda_{o,n}^2\} \quad \lambda_{o,1}^2 \geq \dots \geq \lambda_{o,n}^2 > 0$$

Change coordinates  $\mathbf{x}_o(t) := U_o^\top \mathbf{x}(t)$

Energy of the (orthogonal) state components (for  $t \rightarrow \infty$ )

$$\frac{\|y_{1,o}(0)\|}{\|y_{n,o}(0)\|} = \frac{\lambda_{1,o}}{\lambda_{n,o}} \geq \dots \geq \frac{\|y_{n-1,o}(0)\|}{\|y_{n,o}(0)\|} = \frac{\lambda_{n-1,o}}{\lambda_{no}}$$

$\frac{\lambda_{1,o}}{\lambda_{n,o}}$  may be large: the effect of some states nearly invisible  $\Rightarrow$  BAD CONDITIONING

# (INTERNALLY) BALANCED MODELS

Changing bases can make things better

$$\hat{\Pi} = T^{-1}\Pi T^{-T}, \quad \hat{\Omega} = T^{\top}\Omega T$$

**Definition:** Linear system in **Balanced form** if both  $\hat{\Pi}$  and  $\hat{\Omega}$  are **diagonal and equal**.

**THEOREM 4** *Every linear model with  $|\lambda(A)| < 1$ ,  $(A,B)$  reachable +  $(C,A)$  observable can be transformed to balanced form.*

**ALGORITHM :**

1. Compute  $\Pi$  and  $\Omega$ , solutions of the two dual Lyapunov equations.

2. Compute the SVD

$$\Omega = U\Lambda_o U^\top$$

where  $\Lambda_o$  is the diagonal matrix of eigenvalues of  $\Omega$

3. Change basis  $T_1 := \Lambda_o^{-1/2} U^\top$  so that  $\hat{\Omega} = I$ ; compute

$$\hat{\Pi} = U\Lambda_o^{1/2} \Pi \Lambda_o^{1/2} U^\top$$

4. Compute the SVD

$$\hat{\Pi} = V\Lambda^2 V^\top$$

where  $\Lambda^2$  is diagonal matrix with the (ordered) eigenvalues of  $\hat{\Pi}$



5. Second change of basis defined by  $T_2 := V\Lambda^{1/2}$  so as to make  $\bar{\Pi} := T_2^{-1}\hat{\Pi}T_2^{-T} = \Lambda$ , diagonal.

With this change of basis

$$\bar{\Omega} = T_2^\top \hat{\Omega} T_2 = \Lambda^{1/2} V^\top I V \Lambda^{1/2} = \Lambda$$

The Gramians are diagonal and equal  $\bar{\Pi} = \bar{\Omega} = \Lambda$

# MATLAB

`BALREAL` Balanced state-space realization and model reduction  
`[Ab,Bb,Cb] = BALREAL(A,B,C)` returns a balanced state-space realization of the system  $(A,B,C)$ .

`[Ab,Bb,Cb,G,T] = BALREAL(A,B,C)` also returns a vector  $G$  containing the diagonal of the gramian of the balanced realization, and a matrix  $T$ , the similarity transformation used to convert  $(A,B,C)$  to  $(Ab,Bb,Cb)$ . If the system  $(A,B,C)$  is normalized properly, small elements in gramian  $G$  indicate states that can be removed to reduce the model to lower order.

# SINGULAR VALUES OF A LINEAR SYSTEM

The diagonal matrix  $\Lambda$  is a **system invariant** (does not change if basis is changed). Input-output map (from zero initial conditions)

$$\mathbf{y}(t) = \sum_{k=0}^{t-1} CA^{k-1}B\mathbf{u}(t-k) + D\mathbf{u}(t)$$

In matrix form  $y = \mathbb{H}u$ , where

$$\mathbb{H} := \begin{bmatrix} CB & CAB & CA^2B & \dots \\ CAB & CA^2B & CA^3B & \dots \\ CA^2B & CA^3B & \dots & \dots \\ \dots & \dots & \dots & \dots \end{bmatrix}$$

$\Lambda$  is diag of the **singular values of the Hankel matrix**  $\mathbb{H}$  of the system

$$SVD(\mathbb{H}) = U \begin{bmatrix} \Lambda & 0 \\ 0 & 0 \end{bmatrix} V^T$$

# MODEL REDUCTION BY BALANCED TRUNCATION

How to best approximate a “Large” model (assumed stable + reach + obs)

$$\begin{cases} \mathbf{x}(t+1) = A\mathbf{x}(t) + B\mathbf{u}(t) \\ \mathbf{y}(t) = C\mathbf{x}(t) + D\mathbf{u}(t) \end{cases}$$

Bring it to balanced form. Let  $\Lambda$  be partitioned

$$\Lambda = \begin{bmatrix} \Lambda_1 & 0 \\ 0 & \Lambda_2 \end{bmatrix}$$

$\Lambda_2$   $n_2 \times n_2$  made of small singular values (  $\Lambda_1 \gg \Lambda_2$  )

# BALANCED TRUNCATION

Ideally: Best rank  $n_1$  approximation of  $\mathbb{H}$  ....

$$\begin{cases} \begin{bmatrix} \mathbf{x}_1(t+1) \\ \mathbf{x}_2(t+1) \end{bmatrix} = \begin{bmatrix} A_{11} & A_{12} \\ A_{21} & A_{22} \end{bmatrix} \begin{bmatrix} \mathbf{x}_1(t) \\ \mathbf{x}_2(t) \end{bmatrix} + \begin{bmatrix} B_1 \\ B_2 \end{bmatrix} \mathbf{u}(t) \\ \mathbf{y}(t) = \begin{bmatrix} C_1 & C_2 \end{bmatrix} \begin{bmatrix} \mathbf{x}_1(t) \\ \mathbf{x}_2(t) \end{bmatrix} + D\mathbf{u}(t) \end{cases}$$

$\simeq$

$$\begin{cases} \mathbf{x}_1(t+1) = A_{11}\mathbf{x}_1(t) + B_1\mathbf{u}(t) \\ \mathbf{y}(t) = C_1\mathbf{x}_1(t) + D\mathbf{u}(t) \end{cases}$$

**N.B. STABILITY, REACHABILITY AND OBSERVABILITY ARE PRESERVED.**

# STATE-SPACE MODELS OF RANDOM SIGNALS

$\mathbf{y} = \{\mathbf{y}(t, \omega)\}$  discrete-time  $m$ -dimensional **random signal**  $t \in [t_0, +\infty)$ .

Expected value:  $\mathbb{E} \mathbf{y}(t) = 0 \quad \Leftrightarrow \quad \int_{\Omega} \mathbf{y}(t, \omega) dP = 0$  can be subtracted off. All random quantities **zero mean**.

## STOCHASTIC STATE-SPACE MODEL

$$\begin{cases} \mathbf{x}(t+1) = A\mathbf{x}(t) + B\mathbf{w}(t) & \mathbf{x}(t_0) = \mathbf{x}_0 \\ \mathbf{y}(t) = C\mathbf{x}(t) + D\mathbf{w}(t) & , \quad t \geq t_0 \end{cases}$$

$A, B, C, D$  CONSTANT matrices  $\{\mathbf{w}(t)\}$   $p$ -dimensional **white noise** process of variance

$$\mathbb{E} \mathbf{w}(t) \mathbf{w}(s)^\top = I_p \delta(t-s) \quad \mathbb{E} \mathbf{x}_0 \mathbf{w}(t)^\top = 0 \quad \forall t \geq t_0$$

Initial (random) data  $\mathbb{E} \mathbf{x}_0 = 0$  ,  $\text{Var } \mathbf{x}_0 = \Sigma_0$

# MINIMAL MODELS

REACHABILITY:  $\text{rank} [BAB \dots, A^{n-1}B] = n$

OBSERVABILITY:  $\text{rank} \begin{bmatrix} C \\ CA \\ \vdots \\ CA^{n-1} \end{bmatrix} = n$

are necessary but not enough for minimality.

$$\begin{cases} \mathbf{x}(t+1) &= -a\mathbf{x}(t) + (1-a^2)\mathbf{w}(t) \\ \mathbf{y}(t) &= \mathbf{x}(t) + a\mathbf{w}(t) \end{cases}$$

$\mathbf{y}(t)$  is white noise. Has a minimal representation of order  $n = 0$ .

# STATE SPACE MODELS (Cont.)

Unnormalized white inputs:

$$\mathbf{v}(t) := G\mathbf{w}(t) \quad , \quad \mathbf{w}(t) := D\mathbf{w}(t) \quad ,$$

$$\begin{cases} \mathbf{x}(t+1) = A\mathbf{x}(t) + \mathbf{v}(t) & \mathbf{x}(t_0) = \mathbf{x}_0 \\ \mathbf{y}(t) = C\mathbf{x}(t) + \mathbf{w}(t) & t \geq t_0 \end{cases} \quad ,$$

$$Q := E\{\mathbf{v}(t)\mathbf{v}(t)^\top\} \quad S := E\{\mathbf{v}(t)\mathbf{w}(t)^\top\} \quad R := E\{\mathbf{w}(t)\mathbf{w}(t)^\top\}$$

$$E \left\{ \begin{bmatrix} \mathbf{v}(t) \\ \mathbf{w}(t) \end{bmatrix} \begin{bmatrix} \mathbf{v}(t)^\top & \mathbf{w}(t)^\top \end{bmatrix} \right\} = \begin{bmatrix} Q & S \\ S^\top & R \end{bmatrix} \quad .$$

$\{\mathbf{v}(t)\}$  and  $\{\mathbf{w}(t)\}$  in general *correlated* white noise processes

$$B = [\bar{B} \ 0], \quad D = [0 \ \bar{D}] \quad S = BD^\top = 0$$



# THE STATE PROCESS

$\{\mathbf{x}(t)\}$  is a wide-sense **Markov process**,

$$\hat{\mathbb{E}} [\mathbf{x}(t) \mid \mathbf{x}(\tau) \tau \leq s] = \hat{\mathbb{E}} [\mathbf{x}(t) \mid \mathbf{x}(s)] \quad , \quad \forall t \geq s \quad ,$$

If  $\{\mathbf{w}(t)\}$  and  $\mathbf{x}_0$  jointly Gaussian, then  $\{\mathbf{x}(t)\}$  is Gaussian and Markov in strict sense. **State Variance**

$$\Sigma(t) = \mathbb{E} \mathbf{x}(t) \mathbf{x}(t)^\top := \text{Var}(\mathbf{x}(t))$$

Satisfies a LYAPUNOV DIFFERENCE EQUATION

$$\Sigma(t+1) = A \Sigma(t) A^\top + B B^\top \quad , \quad \Sigma(t_0) = \Sigma_0 \quad .$$

$$\Sigma_{\mathbf{x}}(t, s) = A^{t-s} \Sigma(s) \quad , \quad t \geq s$$

## SECOND-ORDER DESCRIPTION

Joint covariances of  $\{\mathbf{y}(t)\}$  and  $\{\mathbf{x}(t)\}$  are completely determined by the model!

Output Covariance  $\Sigma_{\mathbf{y}}(t, s) = \mathbb{E} \mathbf{y}(t) \mathbf{y}(s)^\top$

$$\Sigma_{\mathbf{x}}(t, s) = \begin{cases} A^{t-s} \Sigma(s) & t \geq s \\ \Sigma(t) (A^\top)^{s-t} & t \leq s \end{cases}$$

$$\Sigma_{\mathbf{y}}(t, s) = \begin{cases} CA^{t-s-1} G(s) & t > s \\ C\Sigma(t)C^\top + DD^\top & t = s \\ G(t)^\top (A^\top)^{s-t-1} C^\top & t < s \end{cases}$$

$$G(s) := A\Sigma(s)C^\top + BD^\top$$

# ASYMPTOTIC STATIONARITY

**Definition:**  $\{\mathbf{y}(t)\}$  is *asymptotically stationary* if for  $t - t_0 \rightarrow \infty$ ,  $\Sigma_{\mathbf{y}}(t, s)$ ,  $t, s \geq t_0$ , tends to depend on the difference  $t - s$ .

If  $A$  **(as.) stable**  $|\lambda(A)| < 1$  then  $\{\mathbf{x}(t)\}$  and  $\{\mathbf{y}(t)\}$  for  $t - t_0 \rightarrow +\infty$ , jointly asympt. stationary

$$\Sigma_{\mathbf{x}}(t - s) = A^{t-s} \bar{\Sigma} \quad , \quad t \geq s \quad ,$$

$$\Sigma_{\mathbf{y}}(t - s) = \begin{cases} CA^{t-s-1} \bar{G} & t > s \\ C\bar{\Sigma}C^{\top} + DD^{\top} & t = s \end{cases}$$

where  $\bar{G} := A\bar{\Sigma}C' + BD'$  and  $\bar{\Sigma} := \lim_{t-t_0 \rightarrow +\infty} \Sigma(t)$  satisfies the LYAPUNOV EQUATION

$$\bar{\Sigma} = A\bar{\Sigma}A^{\top} + BB^{\top} \quad .$$

$\bar{\Sigma}$ , asympt. state variance, does not depend on the initial condition  $\Sigma_0$ .

# THE LYAPUNOV EQUATION

**FACT:** Any two conditions imply the remaining one

i)  $(A, B)$  is reachable

ii)  $A$  is asymptotically stable

iii) The Lyapunov equation

$$X = AXA^T + BB^T$$

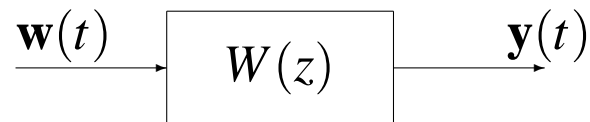
has a unique solution  $P = P^T > 0$

# SHAPING FILTERS AND ARMA MODELS

Assume  $A$  **(as.) stable** i.e.  $|\lambda(A)| < 1$  then  $\{\mathbf{x}(t)\}$  and  $\{\mathbf{y}(t)\}$   $t - t_0 \rightarrow +\infty$ , **jointly asympt. stationary**. Effect of initial conditions disappears

$\{\mathbf{y}(t)\}$  : response to normalized white noise process  $\{\mathbf{w}(t)\}$  of a linear filter (*Shaping Filter*) with transfer function

$$W(z) = C(zI - A)^{-1} B + D$$



$W(z)$  is a rational matrix function. Can be written as a ratio of polynomial matrices

$$W(z) = D(z)^{-1} N(z);$$

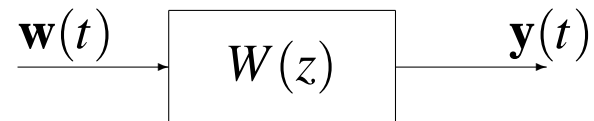
$$D(z) = I z^v + \sum_1^v A_k z^{v-k} \quad N(z) = N_0 z^v + \sum_1^v N_k z^{v-k}$$

$\{\mathbf{y}(t)\}$  may be described by a (multivariable) **ARMA model**

$$\mathbf{y}(t) + \sum_1^v A_k \mathbf{y}(t-k) = N_0 \mathbf{w}(t) + \sum_1^v N_k \mathbf{w}(t-k) \quad .$$

**WARNING:** There are **many** ARMA model representations!

# SHAPING FILTERS AND SPECTRUM



**Wiener-Kintchine formula** gives the **spectral density matrix**  $\Phi(z)$  of  $\{y(t)\}$

$$\Phi(z) = W(z) W(z^{-1})^\top$$

Spectrum is a **rational function** of  $z$ .

Positivity:  $\Phi(e^{j\theta}) = W(e^{j\theta}) W(e^{-j\theta})^\top \geq 0$

# SPECTRAL FACTORIZATION

**FACT:** The shaping filter  $W(z)$  is a **spectral factor** of  $\Phi(z)$

$$\Phi(z) = W(z) W(z^{-1})^{\top}$$

Conversely: modeling by SF is computing (rational) spectral factors from given (rational) spectrum.

Minimal spectral factors (minimal Mc Millan degree)  $\Rightarrow$  Minimal state space models.



# FREQUENCY-DOMAIN COMPUTATIONS

Computing the spectrum of  $\mathbf{y}$ :

$$\Phi(z) = \begin{bmatrix} C(zI - A)^{-1} & I \end{bmatrix} \begin{bmatrix} B \\ D \end{bmatrix} \begin{bmatrix} B^\top & D^\top \end{bmatrix} \begin{bmatrix} (z^{-1}I - A^\top)^{-1}C^\top \\ I \end{bmatrix} .$$

$$\begin{bmatrix} B \\ D \end{bmatrix} \begin{bmatrix} B^\top & D^\top \end{bmatrix} = \begin{bmatrix} BB^\top & BD^\top \\ DB^\top & DD^\top \end{bmatrix} := \begin{bmatrix} Q & S \\ S^\top & R \end{bmatrix} = E \left\{ \begin{bmatrix} \mathbf{v}(t) \\ \mathbf{w}(t) \end{bmatrix} \begin{bmatrix} \mathbf{v}(t)^\top & \mathbf{w}(t)^\top \end{bmatrix} \right\}$$

Scalar process: spectrum from ARMA :

$$\Phi(e^{j\theta}) = \frac{N(e^{j\theta})N(e^{j\theta})^*}{D(e^{j\theta})D(e^{j\theta})^*} = \left| \frac{N(e^{j\theta})}{D(e^{j\theta})} \right|^2$$

## SUMMARY: THREE CLASSES OF MODELS

State Space Models: 
$$\begin{cases} \mathbf{x}(t+1) &= A\mathbf{x}(t) + B\mathbf{w}(t) \\ \mathbf{y}(t) &= C\mathbf{x}(t) + D\mathbf{u}(t) \end{cases}$$

Shaping filters/ARMA : 
$$\mathbf{y}(t) + \sum_1^V A_k \mathbf{y}(t-k) = N_0 \mathbf{w}(t) + \sum_1^V N_k \mathbf{w}(t-k)$$

Spectrum: 
$$\Phi(z) = W(z) W(z^{-1})^\top; \quad W(z) = C(zI - A)^{-1} B + D = D(z)^{-1} N(z)$$

# MODELS FROM COVARIANCE

Assume **given** the covariance function  $\Sigma_{\mathbf{y}}(k) \quad k = 1, 2, \dots$

**PROBLEM (stochastic realization):** From  $\{\Sigma_{\mathbf{y}}(k) \quad k = 0, 1, 2, \dots\}$  compute  $\{A, B, C, D\}$  of a minimal state space model of  $\mathbf{y}$ .

# NECESSARY CONDITIONS

Form the **Hankel Matrix** of  $\Sigma_{\mathbf{y}}$

$$\mathbb{G} := \begin{bmatrix} \Sigma_{\mathbf{y}}(1) & \Sigma_{\mathbf{y}}(2) & \Sigma_{\mathbf{y}}(3) & \dots \\ \Sigma_{\mathbf{y}}(2) & \Sigma_{\mathbf{y}}(3) & \Sigma_{\mathbf{y}}(4) & \dots \\ \Sigma_{\mathbf{y}}(3) & \Sigma_{\mathbf{y}}(4) & \dots & \dots \\ \dots & \dots & \dots & \dots \end{bmatrix}$$

If  $\Sigma_{\mathbf{y}}$  is generated by a linear state space model

$$\Sigma_{\mathbf{y}}(k) = CA^{k-1} \bar{G} \quad \begin{cases} \bar{G} := A\bar{\Sigma}C^{\top} + BD^{\top} \\ \bar{\Sigma} = A\bar{\Sigma}A^{\top} + BB^{\top} \end{cases}$$

$$\Sigma_{\mathbf{y}}(0) = C\bar{\Sigma}C^{\top} + DD^{\top} \quad \text{For } k = 0$$

Must admit a **factorization** of the type

$$\mathbb{G} = \begin{bmatrix} C\bar{G} & CA\bar{G} & CA^2\bar{G} & \dots \\ CA\bar{G} & CA^2\bar{G} & CA^3\bar{G} & \dots \\ CA^2\bar{G} & CA^3\bar{G} & \dots & \dots \\ \dots & \dots & \dots & \dots \end{bmatrix} = \begin{bmatrix} C \\ CA \\ CA^2 \\ \vdots \end{bmatrix} [\bar{G} \quad A\bar{G} \quad A^2\bar{G} \quad \dots]$$

Necessary condition:  $\text{rank } \mathbb{G} = n$

**Positivity** of the function  $k \rightarrow CA^{k-1} \bar{G} \quad (\simeq \Sigma_{\mathbf{y}}(k))$

# TWO-STEPS SOLUTION

**Step 1:** From a *finite submatrix*  $\mathbb{G}_N$  of  $\mathbb{G}$ , of rank =  $n$  compute  $\{A, C, \bar{G}\}$

**ALGORITHM (HO-KALMAN) :**

1. Compute the SVD

$$\mathbb{G}_N = U\Delta V^\top \quad \Delta = \begin{bmatrix} \Sigma & 0 \\ 0 & 0 \end{bmatrix}$$

where  $\Sigma := \text{diag}\{\sigma_1, \dots, \sigma_n\}$  is the diagonal matrix of **nonzero singular values**.

2. Rank  $n$  factorization

$$\mathbb{G}_N = U_n \Sigma V_n^\top = U_n \Sigma^{1/2} \Sigma^{1/2} V_n^\top := \Omega \bar{\Omega}$$

3. Impose

$$\Omega = \begin{bmatrix} C \\ CA \\ CA^2 \\ \vdots \\ CA^{N-1} \end{bmatrix} \quad \bar{\Omega} = [\bar{G} \quad A\bar{G} \quad A^2\bar{G} \quad \dots \quad A^{N-1}\bar{G}]$$

and solve for  $C, A, \bar{G}$ .

## THE HO-KALMAN ALGORITHM (CONT'D)

Computing  $A$ :

$$\mathbf{\Omega} = \begin{bmatrix} C \\ (\downarrow \mathbf{\Omega}) \end{bmatrix} = \begin{bmatrix} (\uparrow \mathbf{\Omega}) \\ CA^{N-1} \end{bmatrix}; \quad (\downarrow \mathbf{\Omega}) = (\uparrow \mathbf{\Omega})A \Rightarrow A = (\uparrow \mathbf{\Omega})^{-L}(\downarrow \mathbf{\Omega})$$

$$\bar{\mathbf{\Omega}} = \begin{bmatrix} \bar{G} & \bar{\mathbf{\Omega}} \end{bmatrix} = \begin{bmatrix} \bar{\mathbf{\Omega}} & A^{N-1}\bar{G} \end{bmatrix}; \quad A\bar{\mathbf{\Omega}} = \bar{\mathbf{\Omega}} \Rightarrow A = (\bar{\mathbf{\Omega}})^{\leftarrow}(\bar{\mathbf{\Omega}})^{\rightarrow R}$$

Found a minimal state space model for the **Causal part** of the spectrum

$$\Phi(z) = \Phi_+(z) + \Phi_+(z^{-1})^\top$$

$$\{A, C, \bar{G}, \Sigma_y(0)\} \Rightarrow \Phi_+(z) = C[zI - A]^{-1}\bar{G} + 1/2\Sigma_y(0)$$



# STOCHASTIC REALIZATION ALGORITHM

Step 2: From the spectrum ( $\Phi_+(z)$ ) to a state-space model

$$\begin{cases} \mathbf{x}(t+1) &= A\mathbf{x}(t) + B\mathbf{w}(t) \\ \mathbf{y}(t) &= C\mathbf{x}(t) + D\mathbf{w}(t) \end{cases},$$

$A, C$ , can be taken the same! Just need to compute  $(B, D)$ .

Recall:  $W(z) := C(zI - A)^{-1}B + D$  is a **spectral factor**  $\Phi(z) = W(z)W(1/z)^\top$

## ALGORITHM:

Given  $(A, C, \bar{G}, \frac{1}{2}\Sigma_{\mathbf{y}}(0))$  a minimal realization of  $\Phi_+(z)$ ,

1. Find  $n \times n$  matrices  $P = P^\top$  solving the *Linear Matrix Inequality*

$$M(P) := \begin{bmatrix} P - APA^\top & \bar{G}^\top - APC^\top \\ \bar{G} - CPA^\top & \Sigma_y(0) - CPC^\top \end{bmatrix} \geq 0$$

2. Compute full column rank matrix factors  $\begin{bmatrix} B \\ D \end{bmatrix}$  of  $M(P)$ ,

$$M(P) = \begin{bmatrix} B \\ D \end{bmatrix} \begin{bmatrix} B^\top & D^\top \end{bmatrix},$$

3.  $W(z) = C(zI - A)^{-1}B + D$ . is a minimal spectral factor (a minimal shaping filter). And conversely....

All symmetric solutions  $P$  of the **LMI** are positive definite : State variance, solution of  $P - APA^\top = BB^\top$

# SOME SPECIAL STOCHASTIC MODELS

Minimal state-space models  $\Leftrightarrow$  set of solutions  $\mathcal{P}$  of the LMI.

If  $\Sigma_{\mathbf{y}}(0) - CPC^{\top} > 0$ , easy to see that  $M(P) \geq 0$  iff  $P$  satisfies the **Algebraic Riccati Inequality**

$$P - APA^{\top} - (\bar{G}^{\top} - APC^{\top})(\Sigma_{\mathbf{y}}(0) - CPC^{\top})^{-1}(\bar{G} - CPA^{\top}) \geq 0.$$

In particular, if  $P$  satisfies the **Algebraic Riccati Equation (ARE)**

$$P = APA^{\top} + (\bar{G}^{\top} - APC^{\top})(\Sigma_{\mathbf{y}}(0) - CPC^{\top})^{-1}(\bar{G} - CPA^{\top}),$$

the corresponding  $W(z)$  is **square**  $m \times m$ .

**FACT:**

Two **special** solutions of the **ARE**:  $P_-, P_+$ , such that  $P_- \leq P \leq P_+$ , for all  $P \in \mathcal{P}$ ,

$$P_- \Rightarrow \begin{bmatrix} B_- \\ D_- \end{bmatrix} \Rightarrow W_-(z) = C(zI - A)^{-1}B_- + D_-$$

The **minimum phase** model: zeros in  $\{|z| \leq 1\}$  i.e. **Causal inverse**

$$P_+ \Rightarrow \begin{bmatrix} B_+ \\ D_+ \end{bmatrix} \Rightarrow W_+(z) = C(zI - A)^{-1}B_+ + D_+$$

The **maximum phase** model: zeros in  $\{|z| \geq 1\}$  i.e. **Anticausal inverse**

NB:  $w(t) = W(z)^{-1}y(t)$  tells how to construct the white noise input!

**WARNING:** with real data the parameters  $\{A, C, \bar{G}\}$  computed by Ho-Kalman may not satisfy the **positivity condition** that  $\Phi_+(z)$  must be the causal part of a power spectrum

$$\Phi_+(e^{j\theta}) + \Phi_+(e^{-j\theta})^\top = W(e^{j\theta}) W(e^{-j\theta})^\top \geq 0$$

This prevents solvability of the Riccati equation.

# THE KALMAN FILTER

PROBLEM: Estimate the state of the linear model

$$\begin{cases} \mathbf{x}(t+1) &= A\mathbf{x}(t) + \mathbf{v}(t) \\ \mathbf{y}(t) &= C\mathbf{x}(t) + \mathbf{w}(t) \end{cases}, \quad t \geq t_0,$$

given past measurements of  $\{\mathbf{y}(t)\}$  ( $m$ -dimensional) up to time  $t$ .

$$\mathbb{E} \left\{ \begin{bmatrix} \mathbf{v}(t) \\ \mathbf{w}(t) \end{bmatrix} \begin{bmatrix} \mathbf{v}(s)^\top, \mathbf{w}(s)^\top \end{bmatrix} \right\} = \begin{bmatrix} Q & S \\ S^\top & R \end{bmatrix} \delta(t-s) \quad R > 0$$

$$\mathbb{E} \mathbf{x}_0 = \mu_0, \quad \text{Var}\{\mathbf{x}_0\} = P_0.$$

KALMAN FILTER (PREDICTOR):

$$\hat{\mathbf{x}}(t+1 | t) = A\hat{\mathbf{x}}(t | t-1) + G(t)\mathbf{e}(t)$$

*The one-step output predictor:*  $\hat{\mathbf{y}}(t | t-1) = C\hat{\mathbf{x}}(t | t-1)$ .

**Innovation process  $\mathbf{e}(t) := \mathbf{y}(t) - C\hat{\mathbf{x}}(t | t-1)$  is white noise !**

# THE KALMAN FILTER (CONTD)

The Kalman gain  $K(t)$

$$K(t) := \left[ AP(t | t-1)C^\top + S \right] \Lambda(t)^{-1}$$

Need **error covariance matrix**,  $P(t | t-1) = \text{Var} \{ \tilde{\mathbf{x}}(t | t-1) := \mathbf{x}(t) - \hat{\mathbf{x}}(t | t-1) \}$

$$P(t | t-1) = \mathbb{E} \tilde{\mathbf{x}}(t | t-1) \tilde{\mathbf{x}}(t | t-1)^\top = P - \hat{P}(t)$$

Innovation covariance  $\Lambda(t) = \mathbb{E} \mathbf{e}(t) \mathbf{e}(t)^\top$ ,

$$\Lambda(t) = CP(t | t-1)C^\top + R$$

**Riccati Equation** for the error covariance  $P(t+1 | t)$ ,

$$P(t+1 | t) = AP(t | t-1)A^\top - K(t)\Lambda(t)K(t)^\top + Q$$

# RICCATI EQUATION

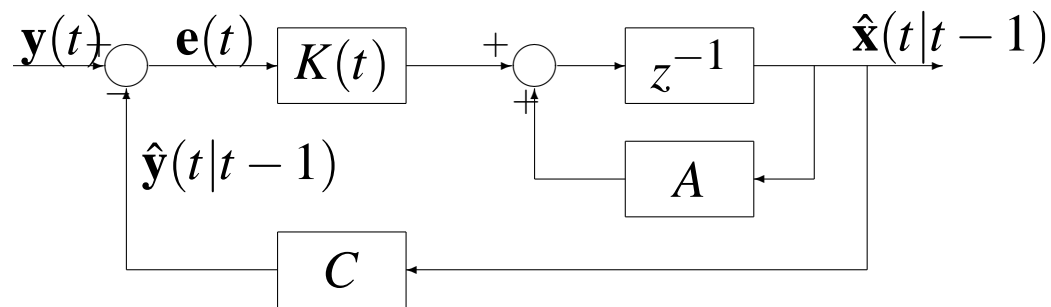
Equivalent form in terms of covariance of  $\hat{\mathbf{x}}(t | t - 1)$

$$\hat{P}(t + 1) = A\hat{P}(t)A^\top + (\bar{G}^\top - A\hat{P}(t)C^\top)(\Sigma_{\mathbf{y}}(0) - C\hat{P}(t)C^\top)^{-1}(\bar{G} - C\hat{P}(t)A^\top),$$



# THE STEADY-STATE KALMAN FILTER

The Kalman filter is an asymptotically stable feedback system!!



Closed loop matrix  $\Gamma(t) = A - K(t)C$ , for  $t - t_0 \rightarrow \infty$  asymptotically stable under very mild conditions

S.S. KALMAN FILTER IS ALSO A STATE MODEL FOR  $y$ !

$$\begin{cases} \hat{\mathbf{x}}(t+1) &= A\hat{\mathbf{x}}(t) + K_{\infty}\mathbf{e}(t) \\ \mathbf{y}(t) &= C\hat{\mathbf{x}}(t) + \mathbf{e}(t) \end{cases}$$

Steady state solution of the RE,  $\lim_{t-t_0 \rightarrow \infty} \hat{P}(t) = P_\infty$ . Solution of the ARE

$$P_\infty = AP_\infty A^\top + (\bar{G}^\top - AP_\infty C^\top)(\Sigma_{\mathbf{y}}(0) - CP_\infty C^\top)^{-1}(\bar{G} - CP_\infty A^\top),$$

SAME RICCATI EQUATION OF STOCHASTIC REALIZATION !!!  $\Rightarrow$  S.S.  
KALMAN FILTER MODEL

$$\begin{cases} \hat{\mathbf{x}}(t+1) &= A\hat{\mathbf{x}}(t) + K_\infty \mathbf{e}(t) \\ \mathbf{y}(t) &= C\hat{\mathbf{x}}(t) + \mathbf{e}(t) \end{cases}$$

# WHAT KIND OF MODEL IS THE STEADY-STATE KALMAN FILTER?

Closed loop matrix  $\Gamma_\infty = A - K_\infty C$ , of the steady state KF is asymptotically stable under very mild conditions.

**Inverse system** (whitening filter)

$$\begin{cases} \hat{\mathbf{x}}(t+1) &= [A - K_\infty C] \hat{\mathbf{x}}(t) + K_\infty \mathbf{y}(t) \\ \mathbf{e}(t) &= -C\hat{\mathbf{x}}(t) + \mathbf{y}(t) \end{cases}$$

Has eigenvalues inside the unit circle. So SSKF is the MINIMUM PHASE MODEL!!

$$P_\infty = P_-$$

# THE BACKWARD KALMAN FILTER

PROBLEM: Estimate the state of the linear model

$$\begin{cases} \mathbf{x}(t+1) &= A\mathbf{x}(t) + B\mathbf{w}(t) \\ \mathbf{y}(t) &= C\mathbf{x}(t) + D\mathbf{w}(t) \end{cases} ,$$

given **future** measurements of  $\{\mathbf{y}(t)\}$  ( $m$ -dimensional) from time  $t$  on.

Backward models.....

# (EARLY) SUBSPACE IDENTIFICATION FOR TIME SERIES [Aoki]

Given observed data (zero mean)

$$\{y_t \mid t = 0, 1, 2, \dots, N\}$$

**Algorithm:**

1. Form covariance estimates

$$\Lambda_k = \frac{1}{N} \sum_{t=0}^{N-k} y_{t+k} y_t^\top \quad (\rightarrow \Sigma_y(k))$$

2. Form the Hankel matrix

$$\mathbb{H}_\Lambda := \begin{bmatrix} \Lambda_1 & \Lambda_2 & \Lambda_3 & \dots & \Lambda_v \\ \Lambda_2 & \Lambda_3 & \Lambda_4 & \dots & \Lambda_{v-1} \\ \Lambda_3 & \Lambda_4 & \dots & \dots & \dots \\ \dots & \dots & \dots & \dots & \dots \\ \Lambda_{v+1} & \dots & \dots & \dots & \Lambda_{2v} \end{bmatrix}$$

Choose  $v$  “large enough” ( $v \geq n$ ).

3. Compute the SVD

$$\mathbb{H}_\Lambda = U\Delta V^\top \quad \Delta = \begin{bmatrix} \Sigma_1 & 0 \\ 0 & \Sigma_2 \end{bmatrix}$$

where  $\Sigma := \text{diag}\{\sigma_1, \dots, \sigma_n\}$  is the diagonal matrix of **dominant singular values**.  $\Sigma_2 \simeq 0$  are neglected.

#### 4. Rank $n$ factorization

$$\mathbb{H}_\Lambda \simeq U_n \Sigma_1 V_n^\top = U_n \Sigma_1^{1/2} \Sigma_1^{1/2} V_n^\top := \Omega \bar{\Omega}$$

#### 5. Impose

$$\Omega = \begin{bmatrix} C \\ CA \\ CA^2 \\ \vdots \\ CA^v \end{bmatrix} \quad \bar{\Omega} = [\bar{G} \quad A\bar{G} \quad A^2\bar{G} \quad \dots \quad A^{v-1}\bar{G}]$$

and get  $C, \bar{G}$  by inspection. Compute  $A$  by solving  $(\downarrow \Omega) = (\uparrow \Omega)A$

$$A = (\uparrow \Omega)^{-L} (\downarrow \Omega) = (\uparrow U_n \Sigma_1^{1/2})^{-L} (\downarrow U_n \Sigma_1^{1/2}) = \Sigma_1^{-1/2} (\uparrow U_n)^\top (\downarrow U_n) \Sigma_1^{1/2}$$

# INNOVATION MODEL IDENTIFICATION

From previous step:  $(A, C, \bar{G})$ . Want  $K, \Lambda = \mathbb{E} \mathbf{e}(t) \mathbf{e}(t)^\top$  in

$$\begin{cases} \hat{\mathbf{x}}(t+1) &= A\hat{\mathbf{x}}(t) + K\mathbf{e}(t) \\ \mathbf{y}(t) &= C\hat{\mathbf{x}}(t) + \mathbf{e}(t) \end{cases}$$

Solve the ARE (minimal solution  $P = P^\top > 0$ )

$$P = APA^\top + (\bar{G}^\top - APC^\top)(\Lambda_0 - CPC^\top)^{-1}(\bar{G} - CPA^\top),$$

$$K = \left[ \bar{G}^\top - APC^\top \right] R(P)^{-1} \quad R(P) = \Lambda_0 - CPC^\top$$



**WARNING:** with real data the parameters  $\{A, C, \bar{G}\}$  computed by Ho-Kalman may not satisfy the **positivity condition** that  $\Phi_+(z)$  must be the causal part of a power spectrum

$$\Phi_+(e^{j\theta}) + \Phi_+(e^{-j\theta})^\top = W(e^{j\theta}) W(e^{-j\theta})^\top \geq 0$$

This prevents solvability of the Riccati equation.

Main drawback of the method: the estimates  $\Lambda(k)$  in general rather poor!

# **SUBSPACE IDENTIFICATION FROM INFINITE/FINITE INPUT-OUTPUT DATA**

OUTLINE OF THE NEXT LECTURES:

1. Some Hilbert space background
2. State construction for stationary processes. Canonical Correlation Analysis (CCA). Stochastic Balancing
3. Realization of stationary processes (no input) with infinite/finite data

4. Subspace algorithms for time series. Relation with to Ho-Kalman algorithm
5. State construction for stationary stoch. systems. Conditional Canonical Correlation Analysis (CCCA)
6. Finite interval realization of stationary stochastic systems with inputs
7. Subspace identification algorithms: CCA, N4SID, MOESP.
8. Numerical aspects

# BASIC IDEA OF SUBSPACE IDENTIFICATION FOR TIME SERIES

Assume we can observe also a **state trajectory**  $\{x_0, x_1, x_2, \dots, x_N\}$  of the model, corresponding to the data

$$\{y_0, y_1, y_2, \dots, y_N\}, \quad y_t \in \mathbb{R}^m$$

Form the “tail” matrices  $\mathbf{Y}_t, \mathbf{X}_t$ ,

$$\begin{aligned} \mathbf{Y}_t &:= [y_t, y_{t+1}, y_{t+2}, \dots] \\ \mathbf{X}_t &:= [x_t, x_{t+1}, x_{t+2}, \dots] \end{aligned}$$

Every sample trajectory  $\{y_t\}, \{x_t\}$  of the system must satisfy the model equations, so there exist  $\{e_t\}$  s.t.

$$\begin{bmatrix} \mathbf{X}_{t+1} \\ \mathbf{Y}_t \end{bmatrix} = \begin{bmatrix} A \\ C \end{bmatrix} \mathbf{X}_t + \begin{bmatrix} K \\ I \end{bmatrix} \mathbf{E}_t$$

## SUBSPACE IDENTIFICATION OF TIME SERIES (cont'd)

$$\begin{bmatrix} \mathbf{X}_{t+1} \\ \mathbf{Y}_t \end{bmatrix} = \begin{bmatrix} A \\ C \end{bmatrix} \mathbf{X}_t + \begin{bmatrix} K \\ I \end{bmatrix} \mathbf{E}_t$$

**Linear Regression ! Solve by Least Squares :**

$$\min_{A,C} \left\| \begin{bmatrix} \mathbf{X}_{t+1} \\ \mathbf{Y}_t \end{bmatrix} - \begin{bmatrix} A \\ C \end{bmatrix} \mathbf{X}_t \right\|$$

getting

$$\begin{bmatrix} \hat{A} \\ \hat{C} \end{bmatrix}_N := \frac{1}{N} \begin{bmatrix} \mathbf{X}_{t+1} \\ \mathbf{Y}_t \end{bmatrix} \mathbf{X}_t^\top \left\{ \frac{1}{N} \mathbf{X}_t \mathbf{X}_t^\top \right\}^{-1}$$

## BASIC IDEA OF SUBSPACE IDENTIFICATION (cont'd)

**Theorem:** If the data are **second order ergodic**, and the inverse exists:

$$\lim_{N \rightarrow \infty} \begin{bmatrix} \hat{A} \\ C \end{bmatrix}_N = \begin{bmatrix} A \\ C \end{bmatrix} \quad (\dagger)$$

**consistent estimate** of  $A, C$ .

**Proof:** HOMEWORK!

# FINITE DATA

**Meaning of Finite Data:** finite string of observed data

$$\{y_0, y_1, y_2, \dots, y_N\}$$

$N$  sufficiently large so that

$$\frac{1}{N+1} \sum_{t=0}^N y_{t+k} y_t^\top \quad k = 1, 2, \dots, T$$

is a “good approximation” of a *finite* set of covariance lags,

$$\{\Lambda(0), \Lambda(1), \dots, \Lambda(T)\},$$

Need to bound  $T$  so that  $T \ll N$ . Rule of thumb is  $T \simeq (1/50)N$

Equivalent to  $\forall a, b \in \mathbb{R}^m$

$$\frac{1}{N+1} \sum_{t=0}^N a^\top y_{t+k} y_{t+j}^\top b \simeq a^\top \mathbb{E} \{ \mathbf{y}(k) \mathbf{y}(j)^\top \} b \quad |k-j| \leq T$$

For  $N \rightarrow \infty$  the **sample covariances**  $\simeq$  **true covariances**.

Assuming  $N$  “very large” numerical TAIL sequences same as **random vectors** !

$$\mathbf{Y}_t \Leftrightarrow \mathbf{y}(t) \quad \frac{1}{N} \mathbf{Y}_t \mathbf{Y}_s^\top \simeq \mathbb{E} \{ \mathbf{y}(t) \mathbf{y}(s)^\top \}$$

EXACTLY THE SAME AS if we had a **finite sequence** of TRULY RANDOM vectors

$$\{ \mathbf{y}(0), \mathbf{y}(1), \mathbf{y}(2), \dots, \mathbf{y}(T) \},$$

extracted from  $\mathbf{y}$ . CAN PRETEND had observations of  $\mathbf{y}$  on the **finite interval**  $[0, T]$ . SAME FORMULAS!



# CONSTRUCTING THE STATE FROM FINITE DATA

Construct  $\hat{\mathbf{x}}(t)$  : **state of transient Kalman filter** on  $[t_0, T]$ :

$$\begin{cases} \hat{\mathbf{x}}(t+1) &= A\hat{\mathbf{x}}(t) + K(t)\hat{\mathbf{e}}(t) \\ \mathbf{y}(t) &= C\hat{\mathbf{x}}(t) + \hat{\mathbf{e}}(t) \\ \hat{\mathbf{x}}(t_0) &= 0 \end{cases}$$

Predictor of finite future based on finite past data :

$$\hat{\mathbf{y}}_t^+ := \mathbb{E} [\mathbf{y}_t^+ | \mathbf{y}_t^-] = \Gamma_k \hat{\mathbf{x}}(t) \quad k = T - t$$

$$\hat{\mathcal{X}}_t = \text{span } \mathbb{E} [\mathbf{y}_t^+ | \mathbf{y}_t^-]$$

# THE STATE BY CANONICAL CORRELATION ANALYSIS

Introduce **Finite past and future** at time  $t$ :

$$\mathbf{y}_t^- := \begin{bmatrix} \mathbf{y}(t_0) \\ \mathbf{y}(t_0 + 1) \\ \vdots \\ \mathbf{y}(t - 1) \end{bmatrix} \simeq \mathbf{Y}_t^- := \begin{bmatrix} \mathbf{Y}_{t_0} \\ \mathbf{Y}_{t_0+1} \\ \vdots \\ \mathbf{Y}_{t-1} \end{bmatrix}$$

$$\mathbf{y}_t^+ := \begin{bmatrix} \mathbf{y}(t) \\ \mathbf{y}(t + 1) \\ \vdots \\ \mathbf{y}(T) \end{bmatrix} \simeq \mathbf{Y}_t^+ := \begin{bmatrix} \mathbf{Y}_t \\ \mathbf{Y}_{t+1} \\ \vdots \\ \mathbf{Y}_T \end{bmatrix}$$

CCA of finite future and past spaces  $\simeq$  CCA of rowspaces of  $\mathbf{Y}_t^-$  and  $\mathbf{Y}_t^+$

# CANONICAL CORRELATION ANALYSIS

CCA is an old concept in statistics. Given two finite-dimensional subspaces  $\mathbf{A}$ ,  $\mathbf{B}$  of zero-mean random variables of dimension  $n$  and  $m$ , one wants to find two special orthonormal bases say  $\{\mathbf{u}_1, \dots, \mathbf{u}_n\}$  for  $\mathbf{A}$ , and  $\{\mathbf{v}_1, \dots, \mathbf{v}_m\}$  for  $\mathbf{B}$  such that

$$\mathbb{E}\{\mathbf{u}_k \mathbf{v}_h\} = \sigma_k \delta_{k,h}, \quad k, h = 1, \dots, \min\{n, m\}$$

This is the same as asking that the correlation matrix of the two random vectors  $\mathbf{u} := [\mathbf{u}_1, \dots, \mathbf{u}_n]'$  and  $\mathbf{v} := [\mathbf{v}_1, \dots, \mathbf{v}_m]'$  made with the elements of the two bases, should be diagonal, i.e. assuming for example that  $n \geq m$ ,

$$\mathbb{E}\{\mathbf{u}\mathbf{v}'\} = \begin{bmatrix} \sigma_1 & 0 & \dots & 0 \\ 0 & \sigma_2 & 0 & \dots \\ \vdots & & \ddots & \\ & & & \sigma_m \\ 0 & 0 & \dots & 0 & 0 \end{bmatrix}$$

To make this choice of basis unique one further requires that all the  $\sigma_k$ 's be nonnegative and ordered in decreasing magnitude.

That two orthonormal bases of this kind always exist follows by considering the singular value decomposition of the projection operator  $\mathbb{E}_{\mathbf{B}}^{\mathbf{A}}$ .

Choosing as orthonormal basis in  $\mathbf{A}$  and in  $\mathbf{B}$  precisely the principal directions  $\{\mathbf{u}_1, \dots, \mathbf{u}_n\}$  and  $\{\mathbf{v}_1, \dots, \mathbf{v}_m\}$  of  $\mathbb{E}_{\mathbf{B}}^{\mathbf{A}}$ , one has

$$\mathbb{E}_{\mathbf{B}}^{\mathbf{A}} \xi = \sum_{k=1}^n \sigma_k \langle \xi, \mathbf{v}_k \rangle \mathbf{u}_k$$

from which it is obvious that the two bases have the required properties. Uniqueness is guaranteed when and only when the singular values  $\{\sigma_k\}$ , which in this context are called *canonical correlation coefficients*, are all distinct.

# CCA ALGORITHM

1. **Normalization:** Form  $T_- := \frac{1}{N} \mathbf{Y}_t^- (\mathbf{Y}_t^-)^\top$      $T_+ := \frac{1}{N} \mathbf{Y}_t^+ (\mathbf{Y}_t^+)^\top$   
 Compute (Cholesky) factors  $T_- = L_- L_-^\top$ ,     $T_+ = L_+ L_+^\top$

$$\hat{\mathbf{Y}}_t^- := L_-^{-1} \mathbf{Y}_t^- \quad \hat{\mathbf{Y}}_t^+ := L_+^{-1} \mathbf{Y}_t^+$$

2. **SVD :**

$$\frac{1}{N} \hat{\mathbf{Y}}_t^+ (\hat{\mathbf{Y}}_t^-)^\top = [\hat{U} \quad \tilde{U}] \begin{bmatrix} \hat{\Sigma} & 0 \\ 0 & \tilde{\Sigma} \end{bmatrix} [\hat{V} \quad \tilde{V}]^\top$$

Can be done (QSVD) without forming the Hankel matrix  $\frac{1}{N} \hat{\mathbf{Y}}_t^+ (\hat{\mathbf{Y}}_t^-)^\top$

**Order estimation:** Choose  $n$  so that  $\hat{\Sigma} \gg \tilde{\Sigma}$

### 3. Canonical Variables

$$\hat{\mathbf{X}}_t := \hat{U}^\top \hat{\mathbf{Y}}_t^- = \hat{U}^\top L_-^{-1} \mathbf{Y}_t^- \quad \hat{\hat{\mathbf{X}}}_t := \hat{V}^\top \hat{\mathbf{Y}}_t^+ = \hat{V}^\top L_+^{-1} \mathbf{Y}_t^+$$

$\hat{\hat{\mathbf{X}}}_t$  basis for the *Backward Kalman filter*.

### 4. Balancing of Canonical Variables

$$\mathbf{Z}_t := \hat{\Sigma}^{1/2} \hat{\mathbf{X}}_t \quad \bar{\mathbf{Z}}_t := \hat{\Sigma}^{1/2} \hat{\hat{\mathbf{X}}}_t \quad \frac{1}{N} \mathbf{Z}_t \mathbf{Z}_t^\top = \hat{\Sigma} = \frac{1}{N} \bar{\mathbf{Z}}_t \bar{\mathbf{Z}}_t^\top$$

5. Repeat for  $t = t + 1$  to get  $\mathbf{Z}_{t+1}$  basis in  $\hat{\mathbf{X}}_{t+1}$  and solve

$$\begin{bmatrix} \mathbf{Z}_{t+1} \\ \mathbf{Y}_t \end{bmatrix} = \begin{bmatrix} A \\ C \end{bmatrix} \mathbf{Z}_t + \begin{bmatrix} K(t) \\ I \end{bmatrix} \hat{\mathbf{E}}_t$$

by Least-Squares.

N.B.  $\mathbf{Z}_{t+1}$  must be a *coherent basis* with  $\mathbf{Z}_t$ .

# THE BACKWARD K.F. AND THE $\bar{G} = \bar{C}$ PARAMETERS

$$\begin{cases} \hat{\mathbf{x}}(t-1) = A^\top \hat{\mathbf{x}}(t) + \bar{K}(t) \bar{\mathbf{e}}(t-1) \\ \mathbf{y}(t-1) = \bar{C} \hat{\mathbf{x}}(t) + \bar{\mathbf{e}}(t-1) \\ \bar{\mathbf{x}}(T) = 0 \end{cases}$$

Predictor of past based on finite future data :

$$\hat{\mathbf{y}}_t^- := \mathbb{E} [\mathbf{y}_t^- | \mathbf{y}_t^+] = \bar{\Gamma}_k \hat{\mathbf{x}}(t) \quad k = t$$

The Backward state space :

$$\hat{\mathcal{X}}_t = \text{span } \mathbb{E} [\mathbf{y}_t^- | \mathbf{y}_t^+]$$

$$\text{Backward covariance } \Sigma_{\mathbf{y}}(-\tau) = \mathbb{E} \mathbf{y}(-\tau) \mathbf{y}(0)^\top = \bar{C} A^{\tau-1} C^\top$$

# ESTIMATING THE $B, D$ PARAMETERS

We have the stationary parameters  $(A, C, \bar{G})$  and  $\Lambda_0 \simeq \Sigma_{\mathbf{y}}(0)$

Solve the Algebraic Riccati Equation

$$P = APA^{\top} + (\bar{G}^{\top} - APC^{\top})(\Lambda_0 - CPC^{\top})^{-1}(\bar{G} - CPA^{\top}) \quad (\text{ARE})$$

To get the minimal (stabilizing) solution  $P_-$

$$K = \left[ \bar{G}^{\top} - AP_-C^{\top} \right] R(P_-)^{-1} \quad R(P_-) = \Lambda_0 - CP_-C^{\top} = D_-D_-^{\top}$$

Equivalently  $B_- = KD_-$

**The ARE has a solution iff  $(A, C, \bar{G}, \Lambda_0)$  is positive real !**



## COMPARISON WITH THE “EARLY” ALGORITHM

Conceptually the algorithm is the same as HO-KALMAN applied to the finite **Normalized** Hankel matrix

$$\hat{\mathbb{H}}_{\Lambda} := L_{+}^{-1} \begin{bmatrix} \Lambda_1 & \Lambda_2 & \Lambda_3 & \dots & \Lambda_v \\ \Lambda_2 & \Lambda_3 & \Lambda_4 & \dots & \Lambda_{v-1} \\ \Lambda_3 & \Lambda_4 & \dots & \dots & \dots \\ \dots & \dots & \dots & \dots & \dots \\ \Lambda_{v+1} & \dots & \dots & \dots & \Lambda_{2v} \end{bmatrix} L_{-}^{-\top}$$

$$\hat{\mathbb{H}}_{\Lambda} \simeq \hat{U} \hat{\Sigma} \hat{V}^{\top} \quad \hat{\Sigma} = \text{diag}\{\sigma_1, \dots, \sigma_n\}$$

Leads to exactly the same formulas as

$$\begin{bmatrix} \hat{A} \\ \hat{C} \end{bmatrix}_N = \begin{bmatrix} \mathbf{Z}_{t+1} \\ \mathbf{Y}_t \end{bmatrix} \mathbf{Z}_t \mathbf{Z}_t^{\top}$$

# NUMERICAL ASPECTS

The **LQ factorization** a key step in subspace identification algorithms.

$$\begin{bmatrix} U \\ Y \end{bmatrix} = \begin{bmatrix} L_{uu} & 0 \\ L_{yu} & L_{yy} \end{bmatrix} \begin{bmatrix} Q_u^\top \\ Q_y^\top \end{bmatrix}$$

where  $Q_u^\top Q_u = I$ ,  $Q_y^\top Q_y = I$ ,  $Q_u^\top Q_y = 0$  and  $L_{uu}$ ,  $L_{yy}$  are lower triangular.

$$\mathbb{E} [Y | \mathcal{U}] = Y Q_u [Q_u^\top Q_u]^{-1} Q_u^\top = L_{yu} Q_u^\top$$

$$\mathbb{E} [Y | \mathcal{U}^\perp] = Y Q_y [Q_y^\top Q_y]^{-1} Q_y^\top = L_{yy} Q_y^\top$$

$Q_y^\top$  an orthonormal basis for the orthogonal complement  $\mathcal{U}^\perp$  in  $\mathcal{U} \vee \mathcal{Y}$ .

## NUMERICAL ASPECTS (Cont'd)

**SVD step:** Compute from LQ factorization

$$\mathbb{E}\{\hat{\mathbf{Y}}_t^+ | \hat{\mathbf{Y}}_t^-\} = [\hat{U} \quad \tilde{U}] \begin{bmatrix} \hat{\Sigma} & 0 \\ 0 & \tilde{\Sigma} \end{bmatrix} [\hat{V} \quad \tilde{V}]^\top$$

**Do order estimation:** pick  $n$  such that  $\tilde{U} \simeq 0$

**Extended Observability matrix** from

$$\mathbb{E}\{\hat{\mathbf{Y}}_t^+ | \hat{\mathbf{Y}}_t^-\} \simeq \hat{U}\hat{\Sigma}^{1/2} \hat{\Sigma}^{1/2}\hat{V}^\top := \mathbf{\Omega}_t \hat{\mathbf{X}}(t)$$

Get  $A, C$  from **Shift-Invariance method** :

$$\hat{U}\hat{\Sigma}^{1/2} = \mathbf{\Omega} = \begin{bmatrix} C \\ (\downarrow \mathbf{\Omega}) \end{bmatrix} = \begin{bmatrix} (\uparrow \mathbf{\Omega}) \\ CA^{N-1} \end{bmatrix}; \quad (\downarrow \mathbf{\Omega}) = (\uparrow \mathbf{\Omega})A \quad \Rightarrow \quad A = (\uparrow \mathbf{\Omega})^{-L}(\downarrow \mathbf{\Omega})$$

N.B. **no need** to compute  $\mathbf{Z}_{t+1}$  *coherent basis in*  $\hat{\mathcal{X}}_{t+1}$  with  $\mathbf{Z}_t$  and solving

$$\begin{bmatrix} \mathbf{Z}_{t+1} \\ \mathbf{Y}_t \end{bmatrix} = \begin{bmatrix} A \\ C \end{bmatrix} \mathbf{Z}_t + \begin{bmatrix} K(t) \\ I \end{bmatrix} \hat{\mathbf{E}}_t$$

by Least-Squares.

Get  $A^\top, \bar{C}$  from the backward procedure

$$\mathbb{E}\{\hat{\mathbf{Y}}_t^- | \hat{\mathbf{Y}}_t^+\} \simeq \hat{V} \hat{\Sigma}^{1/2} \hat{\Sigma}^{1/2} \hat{U}^\top := \bar{\Omega}_t \hat{\mathbf{X}}(t)$$

$$\bar{\Omega} = \begin{bmatrix} \bar{C} \\ (\downarrow \bar{\Omega}) \end{bmatrix} = \begin{bmatrix} (\uparrow \bar{\Omega}) \\ \bar{C}(A^\top)^{N-1} \end{bmatrix}; \quad (\downarrow \bar{\Omega}) = (\uparrow \bar{\Omega})A^\top \Rightarrow A^\top = (\uparrow \bar{\Omega})^{-L}(\downarrow \bar{\Omega})$$

Only need to pick the first block of  $m$  rows to get  $\bar{C}$ .

# ORDER SELECTION

Minimize Akaike-type criterion

$$NIC(n) := \sum_{k=n+1}^{n_{MAX}} \hat{\sigma}_k^2 - d(n) \frac{\log N}{N}$$

where  $d(n)$  = number of additional free parameters in a model of order  $n_{MAX} > n$ .

**Consistency** If data are generated by a true model of order  $n_0$  and  $N \rightarrow \infty$  the minimum  $NIC$  estimate of  $n$  is consistent:

$$\hat{n} \rightarrow n_0 \quad \text{with probability one.}$$

# STATISTICAL PROPERTIES

- **Consistency** If data are generated by a true model
- **Asymptotic Variance of  $A, C$**
- **Efficiency**

# STATE-SPACE MODELS WITH INPUT SIGNALS

$\mathbf{u} = \{\mathbf{u}(t, \omega)\}$  discrete-time  $p$ -dimensional zero-mean **random signal** in  $t \in [t_0, +\infty)$ .

## STOCHASTIC STATE-SPACE MODEL WITH INPUTS

$$\begin{cases} \mathbf{x}(t+1) = A\mathbf{x}(t) + B\mathbf{u}(t) + G\mathbf{w}(t) & \mathbf{x}(t_0) = \mathbf{x}_0 \\ \mathbf{y}(t) = C\mathbf{x}(t) + D\mathbf{u}(t) + J\mathbf{w}(t), & t \geq t_0 \end{cases}$$

$A, B, C, D, G, J$  constant matrices,  $\{\mathbf{x}(t)\}$  is the state process of dimension  $n$ , and  $\{\mathbf{w}(t)\}$  is a normalized white noise process. Assume  $|\lambda(A)| < 1$  (causality).

**N.B: We are not interested in modelling the input  $\{\mathbf{u}(t)\}$ .**

**Assumption: there is no feedback from  $\mathbf{y}$  to  $\mathbf{u}$ . This is the same as: the processes  $\{\mathbf{u}(t)\}$  and  $\{\mathbf{w}(t)\}$  are completely uncorrelated.**

# DETERMINISTIC + STOCHASTIC DECOMPOSITION

State Space Model for  $y$ : *parallel* of models (in general **Non Minimal!**)

$$\text{Stochastic Model} \quad \begin{cases} \mathbf{x}_s(t+1) &= A\mathbf{x}_s(t) + G\mathbf{w}(t) \\ \mathbf{y}_s(t) &= C\mathbf{x}_s(t) + J\mathbf{w}(t) \end{cases}$$

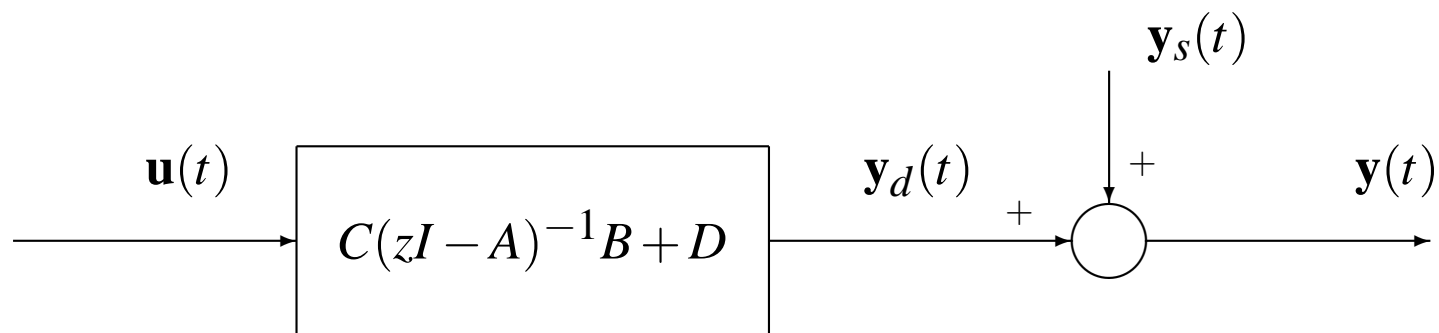
$$\text{Deterministic Model} \quad \begin{cases} \mathbf{x}_d(t+1) &= A\mathbf{x}_d(t) + B\mathbf{u}(t) \\ \mathbf{y}_d(t) &= C\mathbf{x}_d(t) + D\mathbf{u}(t) \end{cases}$$

$$\mathbf{y}(t) = \mathbf{y}_s(t) + \mathbf{y}_d(t) = C [\mathbf{x}_s(t) + \mathbf{x}_d(t)] + D\mathbf{u}(t) + J\mathbf{w}(t)$$

NB.  $\mathbf{x}_s(t)$  **uncorrelated with  $\mathbf{u}$**   $\Rightarrow$   $\mathbf{x}_s(t)$  **uncorrelated with  $\mathbf{x}_d$**  !



# FROM STATE-SPACE TO ARMAX



Deterministic system + “stochastic error” decomposition :

$$\begin{aligned} \mathbf{y}(t) &= [C(zI - A)^{-1}B + D] \mathbf{u}(t) + [C(zI - A)^{-1}G + J] \mathbf{w}(t) \\ &:= F(z)\mathbf{u}(t) + G(z)\mathbf{w}(t) \end{aligned}$$

NB:  $F(z)$  and  $G(z)$  realized with the same  $(A, C)$  pair. In general non-minimal realizations

$F(z)$   $G(z)$  rational. Can be written as a ratio of polynomial matrices with **same denominator**

$$F(z) = A(z)^{-1} B(z); \quad G(z) = A(z)^{-1} C(z)$$

$$A(z) = I z^v + \sum_1^v A_k z^{v-k} \quad B(z) = \sum_1^v B_k z^{v-k} \quad C(z) = C_0 z^v + \sum_1^v C_k z^{v-k}$$

$\{\mathbf{y}(t)\}$  may be described also by the **ARMAX model**

$$\mathbf{y}(t) + \sum_1^v A_k \mathbf{y}(t-k) = \sum_1^v B_k \mathbf{u}(t-k) + C_0 \mathbf{w}(t) + \sum_1^v C_k \mathbf{w}(t-k) \quad .$$

# IDENTIFICATION OF SYSTEMS WITH INPUTS (NO FEEDBACK)

Could be done in two ways. Either identify the **joint model** or first compute  $\mathbf{y}_d(t) = \mathbb{E} \{ \mathbf{y}(t) \mid H(\mathbf{u}) \}$  and identify a **Deterministic Model**

$$\begin{cases} \mathbf{x}_d(t+1) &= A\mathbf{x}_d(t) + B\mathbf{u}(t) \\ \mathbf{y}_d(t) &= C\mathbf{x}_d(t) + D\mathbf{u}(t) \end{cases}$$

Then identify a stochastic model for the **disturbance**

$$\mathbf{y}_s(t) = \mathbf{y}(t) - \mathbf{y}_d(t) = \mathbb{E} \{ \mathbf{y}(t) \mid H(\mathbf{u})^\perp \}$$

Shall do the JOINT model only.

# JOINT INNOVATION MODEL WITH INPUTS

**Steady state Kalman filter:**  $\hat{\mathbf{x}}(t+1) = \mathbb{E} \{ \mathbf{x}(t+1) \mid \mathbf{y}(s), \mathbf{u}(s) s \leq t \}$

**Innovation process (of  $\mathbf{y}$ ):**  $\hat{\mathbf{e}}(t) = \mathbf{y}(t) - C\hat{\mathbf{x}}(t) - D\mathbf{u}(t)$  **white noise !**

$$\begin{bmatrix} \hat{\mathbf{x}}(t+1) \\ \mathbf{y}(t) \end{bmatrix} = \begin{bmatrix} A & B \\ C & D \end{bmatrix} \begin{bmatrix} \hat{\mathbf{x}}(t) \\ \mathbf{u}(t) \end{bmatrix} + \begin{bmatrix} K \\ I \end{bmatrix} \mathbf{e}(t)$$

**No feedback from  $\mathbf{y}$  to  $\mathbf{u}$ :**  $\mathbf{e}(t) \perp \mathbf{u}(\tau) \quad \forall \tau, t$

$$\begin{bmatrix} A & B \\ C & D \end{bmatrix} = \mathbb{E} \left\{ \begin{bmatrix} \hat{\mathbf{x}}(t+1) \\ \mathbf{y}(t) \end{bmatrix} \begin{bmatrix} \hat{\mathbf{x}}(t) \\ \mathbf{u}(t) \end{bmatrix}^\top \right\} \left( \mathbb{E} \left\{ \begin{bmatrix} \hat{\mathbf{x}}(t) \\ \mathbf{u}(t) \end{bmatrix} \begin{bmatrix} \hat{\mathbf{x}}(t) \\ \mathbf{u}(t) \end{bmatrix}^\top \right\} \right)^{-1}$$

**Parameters are uniquely determined by the basis  $\mathbf{x}(t)$  !**

# IDENTIFICATION OF THE DETERMINISTIC SUBSYSTEM

**Problem :** Assume the data are generated by a true stochastic system of order  $n$ . From observed input-output time series

$$\{y_0, y_1, y_2, \dots, y_N\}, \quad y_t \in \mathbb{R}^m \quad \{u_0, u_1, u_2, \dots, u_N\}, \quad u_t \in \mathbb{R}^p$$

find estimates (in a certain basis)  $\begin{bmatrix} A & \hat{B} \\ C & D \end{bmatrix}_N$

such that (**consistency**)

$$\lim_{N \rightarrow \infty} \begin{bmatrix} A & \hat{B} \\ C & D \end{bmatrix}_N = \begin{bmatrix} A & B \\ C & D \end{bmatrix}$$

# BASIC IDEA OF SUBSPACE IDENTIFICATION

Assume we can observe also a **state trajectory**  $\{x_0, x_1, x_2, \dots, x_N\}$ , corresponding to the I/O data

$$\{y_0, y_1, y_2, \dots, y_N\}, \quad y_t \in \mathbb{R}^m \quad \{u_0, u_1, u_2, \dots, u_N\}, \quad u_t \in \mathbb{R}^p$$

Form the “tail” matrices  $\mathbf{Y}_t, \mathbf{X}_t, \mathbf{U}_t$

$$\begin{aligned} \mathbf{Y}_t &:= [y_t, y_{t+1}, y_{t+2}, \dots] \\ \mathbf{X}_t &:= [x_t, x_{t+1}, x_{t+2}, \dots] \\ \mathbf{U}_t &:= [u_t, u_{t+1}, u_{t+2}, \dots] \end{aligned}$$

Every sample trajectory  $\{y_t\}, \{x_t\}, \{u_t\}$  of the system must satisfy the model equations, so

$$\begin{bmatrix} \mathbf{X}_{t+1} \\ \mathbf{Y}_t \end{bmatrix} = \begin{bmatrix} A & B \\ C & D \end{bmatrix} \begin{bmatrix} \mathbf{X}_t \\ \mathbf{U}_t \end{bmatrix} + \begin{bmatrix} K \\ I \end{bmatrix} \mathbf{E}_t$$

## BASIC IDEA OF SUBSPACE IDENTIFICATION (cont'd)

$$\begin{bmatrix} \mathbf{X}_{t+1} \\ \mathbf{Y}_t \end{bmatrix} = \begin{bmatrix} A & B \\ C & D \end{bmatrix} \begin{bmatrix} \mathbf{X}_t \\ \mathbf{U}_t \end{bmatrix} + \begin{bmatrix} K \\ I \end{bmatrix} \mathbf{E}_t$$

**Linear Regression !** Solve by Least Squares :

$$\min_{A,C,B,D} \left\| \begin{bmatrix} \mathbf{X}_{t+1} \\ \mathbf{Y}_t \end{bmatrix} - \begin{bmatrix} A & B \\ C & D \end{bmatrix} \begin{bmatrix} \mathbf{X}_t \\ \mathbf{U}_t \end{bmatrix} \right\|$$

getting

$$\begin{bmatrix} \hat{A} & \hat{B} \\ \hat{C} & \hat{D} \end{bmatrix}_N := \frac{1}{N} \begin{bmatrix} \mathbf{X}_{t+1} \\ \mathbf{Y}_t \end{bmatrix} \begin{bmatrix} \mathbf{X}_t \\ \mathbf{U}_t \end{bmatrix}^\top \left\{ \frac{1}{N} \begin{bmatrix} \mathbf{X}_t \\ \mathbf{U}_t \end{bmatrix} \begin{bmatrix} \mathbf{X}_t \\ \mathbf{U}_t \end{bmatrix}^\top \right\}^{-1}$$

## BASIC IDEA OF SUBSPACE IDENTIFICATION (cont'd)

$$\begin{bmatrix} \hat{A} & \hat{B} \\ \hat{C} & \hat{D} \end{bmatrix}_N := \frac{1}{N} \begin{bmatrix} \mathbf{X}_{t+1} \\ \mathbf{Y}_t \end{bmatrix} \begin{bmatrix} \mathbf{X}_t \\ \mathbf{U}_t \end{bmatrix}^\top \left\{ \frac{1}{N} \begin{bmatrix} \mathbf{X}_t \\ \mathbf{U}_t \end{bmatrix} \begin{bmatrix} \mathbf{X}_t \\ \mathbf{U}_t \end{bmatrix}^\top \right\}^{-1}$$

**Theorem:** If the data are **second order ergodic**, there is no feedback and the inverse exists:

$$\lim_{N \rightarrow \infty} \begin{bmatrix} \hat{A} & \hat{B} \\ \hat{C} & \hat{D} \end{bmatrix}_N = \begin{bmatrix} A & B \\ C & D \end{bmatrix} \quad (\dagger)$$

**consistent estimate** of  $A, B, C, D$ .

**Proof:** HOMEWORK



## SECOND ORDER ERGODICITY

For  $N \rightarrow \infty$  sample covariances converge to true covariances, say

$$\frac{1}{N} \sum_{k=t}^{t+N} \{y_k \mathbf{u}_k^\top\} = \frac{1}{N} \mathbf{Y}_t \mathbf{U}_s^\top \rightarrow \mathbb{E} \{\mathbf{y}(t) \mathbf{u}(s)^\top\} \quad N \rightarrow \infty$$

For  $N \rightarrow \infty$  the **sample covariances can be substituted by the true ones.**

Assuming  $N$  “very large” can use **random variables** instead of numerical sequences!

$$\mathbf{y}(t) \Leftrightarrow \mathbf{Y}_t, \quad \mathbf{u}(t) \Leftrightarrow \mathbf{U}_t, \quad \text{etc.}$$

## COMMENTS

STATE SEQUENCE IS NOT AVAILABLE: NEED TO CONSTRUCT THE STATE FROM INPUT-OUTPUT DATA!

Easy to do if **infinite past data** were available at time  $t$ : want to construct the

**Steady state Kalman filter:**  $\hat{\mathbf{x}}(t) = \mathbb{E} \{ \mathbf{x}(t) \mid \mathbf{y}(s), \mathbf{u}(s) s < t \}$

**Innovation process (of  $\mathbf{y}$ ):**  $\hat{\mathbf{e}}(t) = \mathbf{y}(t) - C\hat{\mathbf{x}}(t) - D\mathbf{u}(t)$  **white noise !**

$$\begin{bmatrix} \hat{\mathbf{x}}(t+1) \\ \mathbf{y}(t) \end{bmatrix} = \begin{bmatrix} A & B \\ C & D \end{bmatrix} \begin{bmatrix} \hat{\mathbf{x}}(t) \\ \mathbf{u}(t) \end{bmatrix} + \begin{bmatrix} K \\ I \end{bmatrix} \mathbf{e}(t)$$

Pick basis vector in the state space of this model : Generalize Akaike procedure by **conditional CCA**

# CONSTRUCTING THE STATE SPACE OF A JOINT STATIONARY MODEL

$\mathbf{X}_t := \text{span} \{ \hat{\mathbf{x}}_1(t), \hat{\mathbf{x}}_2(t) \dots, \hat{\mathbf{x}}_n(t) \}$  State space of a joint innovation model.

Assume we have data starting from  $t = -\infty$

$$\mathcal{P}_t := \mathcal{Y}_t \vee \mathcal{U}_t = \text{span} \{ \mathbf{y}(s), s < t, \mathbf{u}(s), s < t \}$$

$$\mathcal{U}_t^+ := \text{span} \{ \mathbf{u}(s), s \geq t \}$$

**Theorem** If the data are generated by a finite-dimensional stationary model and there is no feedback,

$$\mathbf{X}_t = \text{span} \{ \mathbb{E}_{\|\mathcal{U}_t^+} \{ \mathbf{y}(t+h) \mid \mathcal{P}_t \}; h = 0, 1, \dots, n \}$$

**State-Space = Oblique Predictor Space = Oblique projection of future outputs onto joint past along future inputs**

**Proof:**

$$\begin{aligned}\mathbf{y}(t+h) &= CA^h \hat{\mathbf{x}}(t) + \sum_{k=0}^{h-1} CA^{h-1-k} B \mathbf{u}(t+k) + D \mathbf{u}(t+h) \\ &\quad + \sum_{k=0}^{h-1} CA^{h-1-k} K \mathbf{e}(t+k) + J \mathbf{e}(t+h)\end{aligned}$$

since  $\mathbf{e}(t+k) \perp \mathcal{P}_t$ :

$$\begin{aligned}\mathbb{E}\{\mathbf{y}(t+h) \mid \mathcal{P}_t \vee \mathcal{U}_t^+\} &= \mathbb{E}\{\mathbf{y}(t+h) \mid \mathcal{P}_t \vee \mathcal{U}_{t|t+h}\} \\ &= CA^h \hat{\mathbf{x}}(t) + \sum_{k=0}^{h-1} CA^{h-1-k} B \mathbf{u}(t+k) + D \mathbf{u}(t+h) \\ &= \mathbb{E}_{\|\mathcal{U}_t^+\} \{\mathbf{y}(t+h) \mid \mathcal{P}_t\} + \mathbb{E}_{\|\mathcal{P}_t\} \{\mathbf{y}(t+h) \mid \mathcal{U}_{[tt+h]}\}\end{aligned}$$

HENCE:  $\mathbb{E}_{\|\mathcal{U}_t^+\} \{\mathbf{y}(t+h) \mid \mathcal{P}_t\} = CA^h \hat{\mathbf{x}}(t) \quad h = 0, 1, \dots \quad . \quad \text{QED}$

# OBLIQUE PROJECTIONS

Let  $\mathcal{A} = \text{span}\{\mathbf{a}\}$ ,  $\mathcal{B} = \text{span}\{\mathbf{b}\}$ .

The **oblique projection of  $\mathbf{v}$  onto  $\mathcal{A}$  along  $\mathcal{B}$**  is

$$\mathbb{E}_{\parallel\mathcal{B}}\{\mathbf{v} | \mathcal{A}\} = \begin{bmatrix} \mathbb{E}\{\mathbf{v}\mathbf{a}^\top\} & \mathbb{E}\{\mathbf{v}\mathbf{b}^\top\} \end{bmatrix} \begin{bmatrix} \mathbb{E}\{\mathbf{a}\mathbf{a}^\top\} & \mathbb{E}\{\mathbf{a}\mathbf{b}^\top\} \\ \mathbb{E}\{\mathbf{b}\mathbf{a}^\top\} & \mathbb{E}\{\mathbf{b}\mathbf{b}^\top\} \end{bmatrix}^\dagger \begin{bmatrix} \mathbf{a} \\ 0 \end{bmatrix}$$

If  $\mathcal{A} \perp \mathcal{B}$  i.e.  $\mathcal{A}$  and  $\mathcal{B}$  are **orthogonal**

$$\mathbb{E}_{\parallel\mathcal{B}}\{\mathbf{v} | \mathcal{A}\} = \mathbb{E}\{\mathbf{v} | \mathcal{A}\}$$

If  $\mathcal{A} \cap \mathcal{B} = \{0\}$  i.e.  $\mathcal{A}$  and  $\mathcal{B}$  are in **direct sum** unique decomposition

$$\mathbb{E}\{\mathbf{v} | \mathcal{A} + \mathcal{B}\} = \mathbb{E}_{\parallel\mathcal{B}}\{\mathbf{v} | \mathcal{A}\} + \mathbb{E}_{\parallel\mathcal{A}}\{\mathbf{v} | \mathcal{B}\}$$

# OBLIQUE PROJECTIONS

“Conditional quantities”

$$\mathbb{E} \left\{ \mathbf{v} \mid \mathbf{b}^\perp \right\} := \mathbf{v} - \mathbb{E} \left\{ \mathbf{v} \mid \mathbf{b} \right\}, \quad \mathbb{E} \left\{ \mathbf{a} \mid \mathbf{b}^\perp \right\} := \mathbf{a} - \mathbb{E} \left\{ \mathbf{a} \mid \mathbf{b} \right\}$$

$$\Sigma_{\mathbf{v}\mathbf{a}|\mathbf{b}} := \text{Cov} \left[ \mathbb{E} \left\{ \mathbf{v} \mid \mathbf{b}^\perp \right\}, \mathbb{E} \left\{ \mathbf{a} \mid \mathbf{b}^\perp \right\} \right], \quad \Sigma_{\mathbf{a}\mathbf{a}|\mathbf{b}} := \text{Var} \left[ \mathbb{E} \left\{ \mathbf{a} \mid \mathbf{b}^\perp \right\} \right]$$

**Fact:** Assume  $\mathcal{A} \cap \mathcal{B} = \{0\}$  and  $\mathbf{a}$  and  $\mathbf{b}$  are bases:

$$\mathbb{E}_{\parallel \mathcal{B}} \left\{ \mathbf{v} \mid \mathcal{A} \right\} = \Sigma_{\mathbf{v}\mathbf{a}|\mathbf{b}} \Sigma_{\mathbf{a}\mathbf{a}|\mathbf{b}}^{-1} \mathbf{a} \quad \mathbb{E}_{\parallel \mathcal{A}} \left\{ \mathbf{v} \mid \mathcal{B} \right\} = \Sigma_{\mathbf{v}\mathbf{b}|\mathbf{a}} \Sigma_{\mathbf{b}\mathbf{b}|\mathbf{a}}^{-1} \mathbf{b}$$

# THE STATE SPACE OF A STATIONARY MODEL

Any choice of basis in the **oblique predictor space**

$$\mathbf{X}_t = \mathbb{E}_{\|\mathcal{U}_{[t,t+n]}} \{ \mathcal{Y}_{[t,t+n]} \mid \mathcal{Y}_t^- \vee \mathcal{U}_t^- \} = \text{span} \{ \mathbb{E}_{\|\mathcal{U}_t^+} [\mathbf{y}(t+h) \mid \mathcal{P}_t] ; h=0, 1, 2' \dots \}$$

provides a minimal innovation model ( Steady state Kalman filter)

$$\begin{cases} \hat{\mathbf{x}}(t+1) &= A\hat{\mathbf{x}}(t) + B\mathbf{u}(t) + K\mathbf{e}(t) \\ \mathbf{y}(t) &= C\hat{\mathbf{x}}(t) + D\mathbf{u}(t) + \mathbf{e}(t) \end{cases}$$

Can compute a basis in  $\mathbf{X}_t$  by **conditional CCA**: SVD of the normalized conditional covariance of future outputs  $\mathbf{y}_t^+$  and (joint!) past

$$\mathbf{p}(t) := \begin{bmatrix} \mathbf{u}_t^- \\ \mathbf{y}_t^- \end{bmatrix} \quad [\infty \times 1 \text{ past observations}]$$

given future inputs  $\mathbf{u}_t^+$ .

# STATIONARY CONDITIONAL CCA

If data are described by a true  $n$ -dimensional model, the **Conditional Hankel Matrix**

$$H_{\mathbf{y}^+ \mathbf{p} | \mathbf{u}^+} := \text{Cov} \left[ \mathbb{E} \left\{ \mathbf{y}_t^+ \mid (\mathbf{u}^+)^{\perp} \right\}, \mathbb{E} \left\{ \mathbf{p}(t) \mid (\mathbf{u}_t^+)^{\perp} \right\} \right]$$

has finite rank  $n \Rightarrow \mathbf{y}_t^+$  and  $\mathbf{u}_t^+$  can be taken to be finite dimensional vectors.

Cholesky factors

$$H_{\mathbf{y}^+ \mathbf{y}^+ | \mathbf{u}^+} := \text{Var} \left[ \mathbb{E} \left\{ \mathbf{y}_t^+ \mid (\mathbf{u}^+)^{\perp} \right\} \right] = L_{\mathbf{y}^+ | \mathbf{u}^+} L_{\mathbf{y}^+ | \mathbf{u}^+}^{\top},$$

$$H_{\mathbf{p} \mathbf{p} | \mathbf{u}^+} := \text{Var} \left[ \mathbb{E} \left\{ \mathbf{p}(t) \mid (\mathbf{u}^+)^{\perp} \right\} \right] = L_{\mathbf{p} | \mathbf{u}^+} L_{\mathbf{p} | \mathbf{u}^+}^{\top}$$

Do SVD of the *normalized conditional Hankel matrix*

$$\hat{H}_{\mathbf{y}^+ \mathbf{p} | \mathbf{u}^+} := L_{\mathbf{y}^+ | \mathbf{u}^+}^{-1} H_{\mathbf{y}^+ \mathbf{p} | \mathbf{u}^+} L_{\mathbf{p} | \mathbf{u}^+}^{-\top}$$



Order estimation

$$\hat{H}_{\mathbf{y}^+|\mathbf{p}^+|\mathbf{u}^+} := [\hat{U} \quad \tilde{U}] \begin{bmatrix} \hat{\Sigma} & 0 \\ 0 & \tilde{\Sigma} \end{bmatrix} [\hat{V} \quad \tilde{V}]^\top \simeq \hat{U} \hat{\Sigma} \hat{V}^\top$$

Canonical state

$$\mathbf{z}(t) = \hat{\Sigma}^{1/2} \hat{V}^\top L_{\mathbf{p}^+|\mathbf{u}^+}^{-1} \mathbf{p}(t)$$

MAIN DIFFICULTY: The **The infinite past**  $\mathbf{p}(t)$  spanning  $\mathcal{Y}_{-\infty|t} \vee \mathcal{U}_{-\infty|t}$  is not available !! Approximation with available **finite past** yields biased estimates. Bias may be large if the zeros of the true system are far from the unit circle.

## ONLY FINITE DATA ARE AVAILABLE!

Infinite past approximation leads to **errors (bias)** in the estimate which do not  $\rightarrow 0$  as  $N \rightarrow \infty$ .

Bias can be made arbitrarily large taking the zeros of the stochastic subsystem arbitrarily close to the unit circle.

For consistency with finite regression data: NEED FINITE-INTERVAL (NON-STATIONARY) STOCHASTIC REALIZATION

# FINITE-INTERVAL INNOVATION MODEL

Modeling using “data” on a finite-interval  $[t_0, T]$ . The estimate

$$\hat{\mathbf{x}}(t) := \mathbb{E} \left[ \mathbf{x}(t) \mid \mathcal{P}_{[t_0, t)} \vee \mathcal{U}_{[t, T]} \right]$$

satisfies the *transient conditional Kalman filter equation*

$$\begin{cases} \hat{\mathbf{x}}(t+1) &= A\hat{\mathbf{x}}(t) + B\mathbf{u}(t) + K(t)\hat{\mathbf{e}}(t) \\ \mathbf{y}(t) &= C\hat{\mathbf{x}}(t) + D\mathbf{u}(t) + \hat{\mathbf{e}}(t) \\ \hat{\mathbf{x}}(t_0) &= \mathbb{E} \left[ \mathbf{x}(t_0) \mid \mathcal{U}_{[t_0, T]} \right] \end{cases}$$

**How to construct  $\hat{\mathbf{x}}(t)$  ?**

Is  $\hat{\mathbf{x}}(t)$  a basis in some predictor space? e.g.  $\mathbb{E}_{\|\mathcal{U}_{[t, T]}} \left[ \mathcal{Y}_{[t, T]} \mid \mathcal{P}_{[t_0, t)} \right]$  ?

Cannot be  $\mathbb{E} \left[ \mathcal{Y}_{[t, T]} \mid \mathcal{P}_{[t_0, t)} \right]$ ; would introduce innovation of  $\mathbf{u}$  !!

# FINITE-INTERVAL REALIZATION THEORY

$$H_d := \begin{bmatrix} D & 0 & \dots & 0 & 0 \\ CB & D & \dots & 0 & 0 \\ \vdots & & & \ddots & \vdots \\ CA^{v-1}B & CA^{v-2}B & \dots & CB & D \end{bmatrix},$$

$$\begin{aligned} & \mathbb{E} \left[ \mathbf{y}(t+h) \mid \mathcal{P}_{[t_0,t)} \vee \mathcal{U}_{[t,T]} \right] = \\ & = \mathbb{E}_{\|\mathcal{U}_{[t,T]}} \left[ \mathbf{y}(t+h) \mid \mathcal{P}_{[t_0,t)} \right] + \mathbb{E}_{\|\mathcal{P}_{[t_0,t)}} \left[ \mathbf{y}(t+h) \mid \mathcal{U}_{[t,T]} \right] \\ & = CA^h \mathbb{E}_{\|\mathcal{U}_{[t,T]}} \left[ \hat{\mathbf{x}}(t) \mid \mathcal{P}_{[t_0,t)} \right] + CA^h \mathbb{E}_{\|\mathcal{P}_{[t_0,t)}} \left[ \hat{\mathbf{x}}(t) \mid \mathcal{U}_{[t,T]} \right] + H_{d,h} \mathbf{u}_t^+ \end{aligned}$$

Projecting along  $\mathcal{U}_{[t,T]}$  kills one piece of

$$\hat{\mathbf{x}}(t) = \mathbb{E}_{\|\mathcal{U}_{[t,T]}} \left[ \hat{\mathbf{x}}(t) \mid \mathcal{P}_{[t_0,t)} \right] + \mathbb{E}_{\|\mathcal{P}_{[t_0,t)}} \left[ \hat{\mathbf{x}}(t) \mid \mathcal{U}_{[t,T]} \right]$$

# PATCHING UP

**Causal component of the state:**  $\xi(t) := \mathbb{E}_{\|\mathcal{U}_{[t,T]}} \left[ \hat{\mathbf{x}}(t) \mid \mathcal{P}_{[t_0,t)} \right]$

$$\mathbb{E}_{\|\mathcal{U}_{[t,T]}} \left[ \mathbf{y}(t+h) \mid \mathcal{P}_{[t_0,t)} \vee \mathcal{U}_{[t,T]} \right] := CA^h \xi(t)$$

From an oblique projection can recover the **Observability Matrix**  $\Gamma_k$ :

$$\hat{\mathbf{y}}_t^+ = \mathbb{E} \left\{ \left[ \begin{array}{c} \mathbf{y}(t) \\ \mathbf{y}(t+1) \\ \vdots \\ \mathbf{y}(t+k) \end{array} \right] \mid \mathcal{P}_{[t_0,t)} \vee \mathcal{U}_{[t,T]} \right\} = \Gamma_k \xi(t) + \text{part in } \mathcal{U}_{[t,T]}$$

Need

$$\mathbb{E} \{ \xi(t) \xi(t)^\top \} > 0 \quad (\text{consistency condition})$$

# THE VAN OVERSCHEE-DE MOOR MODEL

**Pseudostate:**  $\bar{\mathbf{x}}(t) := \Gamma_k^{-L} \hat{\mathbf{y}}_t^+ = \hat{\mathbf{x}}(t) + \Gamma_k^{-L} H_k \mathbf{u}_t^+$

$$\begin{bmatrix} \bar{\mathbf{x}}(t+1) \\ \mathbf{y}(t) \end{bmatrix} = \begin{bmatrix} A \\ C \end{bmatrix} \bar{\mathbf{x}}(t) + \begin{bmatrix} \mathcal{K}_1 \\ \mathcal{K}_2 \end{bmatrix} \mathbf{u}_t^+ + \mathbf{w}_t^\perp \quad (*)$$

$\mathcal{K}_1$   $\mathcal{K}_2$  known linear functions of  $(B, D)$ .

$$\begin{bmatrix} A \\ C \end{bmatrix} \Sigma_{\bar{\mathbf{x}}\bar{\mathbf{x}}|\mathbf{u}^+} = \begin{bmatrix} \Sigma_{\bar{\mathbf{x}}_1\bar{\mathbf{x}}|\mathbf{u}^+} \\ \Sigma_{\bar{\mathbf{y}}\bar{\mathbf{x}}|\mathbf{u}^+} \end{bmatrix}$$

$$\begin{bmatrix} \mathcal{K}_1 \\ \mathcal{K}_2 \end{bmatrix} \Sigma_{\mathbf{u}^+\mathbf{u}^+|\bar{\mathbf{x}}} = \begin{bmatrix} \Sigma_{\bar{\mathbf{x}}_1\mathbf{u}^+|\bar{\mathbf{x}}} \\ \Sigma_{\mathbf{y}\mathbf{u}^+|\bar{\mathbf{x}}} \end{bmatrix}$$

Solve in terms of **Conditional Covariances:**

$$\Sigma_{\bar{\mathbf{x}}\bar{\mathbf{x}}|\mathbf{u}^+} = E\{[\bar{\mathbf{x}}(t) - E(\bar{\mathbf{x}}(t) | \mathbf{u}_t^+)] [\bar{\mathbf{x}}(t) - E(\bar{\mathbf{x}}(t) | \mathbf{u}_t^+)]^\top\} = \Sigma_{\hat{\mathbf{x}}\hat{\mathbf{x}}|\mathbf{u}^+}$$

$$\Sigma_{\mathbf{u}^+\mathbf{u}^+|\bar{\mathbf{x}}} = E\{[\mathbf{u}_t^+ - E(\mathbf{u}_t^+ | \bar{\mathbf{x}}(t))] [\mathbf{u}_t^+ - E(\mathbf{u}_t^+ | \bar{\mathbf{x}}(t))]^\top\}$$

# THE N4SID ALGORITHM

## [vanOverschee-DeMoor94]

1. Predictor matrix based on joint input-output data

$$\hat{Y}_{[t, T-1]} := \mathbb{E} \left[ Y_{[t, T-1]} \mid Y_{[t_0, t)} \vee U_{[t_0, T]} \right]$$

(projection onto the joint rowspace).

2. Compute the oblique projection along  $U_{[t, T]}$

$$Z(t) := \mathbb{E}_{\parallel \mathcal{U}_{[t, T]}} \left[ \hat{Y}_{[t, T-1]} \mid Y_{[t_0, t)} \vee U_{[t_0, t)} \right]$$

to get an estimate of  $\Gamma_k \Xi_t$

- 3 Estimate the order and the observability matrix  $\Gamma_k$  by SVD factorization.

4. The “Pseudostate”  $\bar{X}_t := \Gamma_k^{-L} \hat{Y}_{[t, T-1]}$  obeys the recursion

$$\begin{bmatrix} \bar{X}_{t+1} \\ Y_t \end{bmatrix} = \begin{bmatrix} A \\ C \end{bmatrix} \bar{X}_t + \begin{bmatrix} \mathcal{K}_1 \\ \mathcal{K}_2 \end{bmatrix} U_{[t, T]} + W^\perp$$

5. Compute a coherent pseudostate at time  $t + 1$ :  $\bar{X}_{t+1} := \Gamma_k^{-L} \hat{Y}_{[t+1, T]}$

6. Solve by LS for the unknown parameters  $(A, C)$  and  $(\mathcal{K}_1, \mathcal{K}_2)$

7. Estimate  $(B, D)$  from  $(\mathcal{K}_1, \mathcal{K}_2)$ .



## “ MOESP ”

Start from the stationary innovation model (state vector  $\mathbf{x}(t)$  ) future horizon:  $k = T - t$ ,

$$\begin{bmatrix} \mathbf{y}(t) \\ \mathbf{y}(t+1) \\ \vdots \\ \mathbf{y}(T) \end{bmatrix} = \begin{bmatrix} C \\ CA \\ \vdots \\ CA^k \end{bmatrix} \mathbf{x}(t) + \begin{bmatrix} D & & & 0 & 0 \\ CB & D & & & 0 \\ \vdots & \ddots & \ddots & & \\ CA^{k-1}B & \dots & CB & D & \end{bmatrix} \begin{bmatrix} \mathbf{u}(t) \\ \mathbf{u}(t+1) \\ \vdots \\ \mathbf{u}(T) \end{bmatrix} \\ + \begin{bmatrix} I & & & 0 & 0 \\ CK & I & & & 0 \\ \vdots & \ddots & \ddots & & \\ CA^{k-1}K & \dots & CK & I & \end{bmatrix} \begin{bmatrix} \mathbf{e}(t) \\ \mathbf{e}(t+1) \\ \vdots \\ \mathbf{e}(T) \end{bmatrix}$$

$$\mathbf{y}_t^+ = \Gamma_k \mathbf{x}(t) + H_d \mathbf{u}_t^+ + H_e \mathbf{e}_t^+$$

Want to kill the last two pieces.

## “ MOESP” (Cont’d)

1. Project orthogonally onto  $\mathcal{Y}_{[t_0, t)} \vee \mathcal{U}_{[t_0, T]}$

$$\hat{\mathbf{y}}_t^+ := \mathbb{E} \left[ \mathbf{y}_t^+ \mid \mathcal{P}_{[t_0, t)} \vee \mathcal{U}_{[t, T]} \right] = \Gamma_k \hat{\mathbf{x}}(t) + H_d \mathbf{u}_t^+$$

2. Project onto the orthogonal complement  $\mathcal{U}_{[t, T]}^\perp$

$$\hat{\mathbf{z}}_t^+ := \hat{\mathbf{y}}_t^+ - \mathbb{E} \left[ \hat{\mathbf{y}}_t^+ \mid \mathcal{U}_{[t, T]} \right] = \Gamma_k \hat{\mathbf{x}}^c(t)$$

$$\hat{\mathbf{x}}^c(t) = \hat{\mathbf{x}}(t) - \mathbb{E} \left[ \hat{\mathbf{x}}(t) \mid \mathcal{U}_{[t, T]} \right]$$

3. Factorize  $\hat{\mathbf{z}}_t^+$ , i.e. the matrix

$$Z^c(t) := \mathbb{E} \left[ \hat{\mathbf{Y}}_{[t, T]} \mid \mathcal{U}_{[t, T]}^\perp \right]$$

by SVD to get an estimate of the order  $n$  and of  $\Gamma_k$  e.g.  $\hat{\Gamma}_k$ .

4. Estimate  $(A, C)$  from the estimated observability matrix by the “shift-Invariance method”.

5. Construct a (projection) matrix  $\hat{\Gamma}_k^\perp$  such that  $\hat{\Gamma}_k^\perp \hat{\Gamma}_k = 0$

6. Compute

$$\hat{\Gamma}_k^\perp \hat{\mathbf{y}}_t^+ = \Gamma_k^\perp H_d \mathbf{u}_t^+ + \text{noise}$$

7. A linear function of  $(B, D)$  : estimate  $(B, D)$  by linear regression

$$\hat{\Gamma}_k^\perp \hat{\mathbf{y}}_t^+ = L(A, C) \text{vec}(B, D) + \text{noise}$$

# MOESP $\equiv$ ORTHOGONALIZING REGRESSORS

$$\begin{bmatrix} \bar{X}_{t+1} \\ Y_t \end{bmatrix} = \begin{bmatrix} A \\ C \end{bmatrix} \bar{X}_t + \begin{bmatrix} \mathcal{K}_1 \\ \mathcal{K}_2 \end{bmatrix} U_{[t,T]} + W^\perp$$

Introduce orthogonal regressors  $\bar{X}_t^c := \bar{X}_t - \mathbb{E}\{\bar{X}_t \mid U_{[t,T]}\}$

$$\begin{bmatrix} \bar{X}_{t+1}^c \\ Y_t \end{bmatrix} = \begin{bmatrix} A \\ C \end{bmatrix} \bar{X}_t^c + \begin{bmatrix} \mathcal{K}_1^c \\ \mathcal{K}_2^c \end{bmatrix} U_{[t,T]} + W^\perp$$

Least squares: right-multiply by  $(\bar{X}_t^c)^\top$ ,  $\mathbb{E} U_{[t,T]} (\bar{X}_t^c)^\top \rightarrow 0$

$$\mathbb{E} \bar{X}_t^c (\bar{X}_t^c)^\top \rightarrow \Sigma_{\mathbf{x}^c, \mathbf{x}^c} = \Sigma_{\bar{\mathbf{x}}\bar{\mathbf{x}}|\mathbf{u}^+} = \Sigma_{\hat{\mathbf{x}}\hat{\mathbf{x}}|\mathbf{u}^+}$$

# WEIGHTING

Both N4SID (“Robust”) and MOESP use the preliminary orthogonalization

basic step is the SVD of

$$Z^c(t) := \mathbb{E} \left[ \hat{Y}_{[t,T]} \mid \mathcal{U}_{[t,T]}^\perp \right] \simeq \mathbb{E} \left[ \hat{\mathbf{y}}_t^+ \mid (\mathbf{u}_t^+)^\perp \right] = \mathbb{E} \left[ \mathbf{y}_t^+ \mid (\mathbf{p}(t) - \mathbb{E} \{ \mathbf{p}(t) \mid \mathbf{u}_t^+ \}) \right]$$

$$\mathbf{z}^c(t) = \mathbb{E} \left[ \mathbf{y}_t^+ (\mathbb{E} \{ \mathbf{p}(t) \mid (\mathbf{u}_t^+)^\perp \})^\top \right] \text{Var} \{ \mathbb{E} \{ \mathbf{p}(t) \mid (\mathbf{u}_t^+)^\perp \}^{-1} \mathbb{E} \{ \mathbf{p}(t) \mid (\mathbf{u}_t^+)^\perp \}$$

## Conditional Hankel Matrix

$$H_{\mathbf{y}^+ \mathbf{p} | \mathbf{u}^+} := \text{Cov} \left[ \mathbb{E} \left\{ \mathbf{y}_t^+ \mid (\mathbf{u}^+)^\perp \right\}, \mathbb{E} \left\{ \mathbf{p}(t) \mid (\mathbf{u}_t^+)^\perp \right\} \right]$$

has finite rank  $n \Rightarrow \mathbf{y}_t^+$  and  $\mathbf{u}_t^+$  are finite dimensional vectors.

CCA : introduce Cholesky factors

$$H_{\mathbf{y}^+ \mathbf{y}^+ | \mathbf{u}^+} := \text{Var} \left[ \mathbb{E} \left\{ \mathbf{y}_t^+ \mid (\mathbf{u}^+)^{\perp} \right\} \right] = L_{\mathbf{y}^+ | \mathbf{u}^+} L_{\mathbf{y}^+ | \mathbf{u}^+}^{\top},$$

$$H_{\mathbf{p} \mathbf{p} | \mathbf{u}^+} := \text{Var} \left[ \mathbb{E} \left\{ \mathbf{p}(t) \mid (\mathbf{u}^+)^{\perp} \right\} \right] = L_{\mathbf{p} | \mathbf{u}^+} L_{\mathbf{p} | \mathbf{u}^+}^{\top}$$

CCA: Doing SVD of the *normalized conditional Hankel matrix*

$$\hat{H}_{\mathbf{y}^+ \mathbf{p} | \mathbf{u}^+} := L_{\mathbf{y}^+ | \mathbf{u}^+}^{-1} H_{\mathbf{y}^+ \mathbf{p} | \mathbf{u}^+} L_{\mathbf{p} | \mathbf{u}^+}^{-\top}$$

Same as introducing a **weighting matrix on the left**

$$\mathbf{z}^c(t) \rightarrow L_{\mathbf{y}^+ | \mathbf{u}^+}^{-1} \mathbf{z}^c(t)$$

Weighting on the right side does not make sense....

# ASYMPTOTIC SOLUTION (CONSISTENCY)

For  $N \rightarrow \infty$  the estimates tend to satisfy

$$\begin{bmatrix} A \\ C \end{bmatrix} \Sigma_{\hat{\mathbf{x}}\hat{\mathbf{x}}|\mathbf{u}^+} = \begin{bmatrix} \Sigma_{\bar{\mathbf{x}}_1\bar{\mathbf{x}}|\mathbf{u}^+} \\ \Sigma_{\bar{\mathbf{y}}\bar{\mathbf{x}}|\mathbf{u}^+} \end{bmatrix}$$

$$\begin{bmatrix} \mathcal{K}_1 \\ \mathcal{K}_2 \end{bmatrix} \Sigma_{\mathbf{u}^+\mathbf{u}^+|\bar{\mathbf{x}}} = \begin{bmatrix} \Sigma_{\bar{\mathbf{x}}_1\mathbf{u}^+|\bar{\mathbf{x}}} \\ \Sigma_{\mathbf{y}\mathbf{u}^+|\bar{\mathbf{x}}} \end{bmatrix}$$

Solve in terms of **Conditional Covariances**:

$$\Sigma_{\hat{\mathbf{x}}\hat{\mathbf{x}}|\mathbf{u}^+} = \Sigma_{\hat{\mathbf{x}}^c\hat{\mathbf{x}}^c}$$

$$\Sigma_{\mathbf{u}^+\mathbf{u}^+|\bar{\mathbf{x}}} = E\{[\mathbf{u}_t^+ - E(\mathbf{u}_t^+ | \bar{\mathbf{x}}(t))] [\mathbf{u}_t^+ - E(\mathbf{u}_t^+ | \bar{\mathbf{x}}(t))]^\top\}$$

# CONSISTENCY CONDITION AND ILL-CONDITIONING

**Jansson-Wahlberg consistency condition:**

$$\Sigma_{\hat{\mathbf{x}}\hat{\mathbf{x}}|\mathbf{u}^+} = \Sigma_{\hat{\mathbf{x}}^c\hat{\mathbf{x}}^c} \quad \text{MUST BE NON SINGULAR!}$$

$\Sigma_{\hat{\mathbf{x}}^c\hat{\mathbf{x}}^c} (= \Sigma_{\hat{\mathbf{x}}\hat{\mathbf{x}}|\mathbf{u}^+})$  may be ILL-CONDITIONED!  $\Rightarrow$

The computation of the parameters  $(A, C)$  of the regression will be **ill-conditioned: random fluctuation errors in the data will be amplified.**

$\Sigma_{\hat{\mathbf{x}}\hat{\mathbf{x}}|\mathbf{u}^+}$  ILL-CONDITIONED  $\Leftrightarrow$  Rowspace of  $\hat{\mathbf{X}}_t$  and  $\mathbf{U}_{[t,T]}$  are “NEARLY PARALLEL”

Similar analysis holds for  $(\mathcal{K}_1, \mathcal{K}_2)$  and  $\Sigma_{\mathbf{u}^+\mathbf{u}^+|\bar{\mathbf{x}}}$ .



## Crucial question: how “parallel” are the rowspaces of

$$U_{[t,T]} \quad \text{and} \quad \bar{X}_t = \hat{X}_t + \Gamma_k^{-L} H_k U_{[t,T]}$$

If some (canonical) angles are nearly zero  $\Rightarrow$  the computation of the parameters  $(A, C)$  and  $(\mathcal{K}_1, \mathcal{K}_2)$  of the regression will be **ill-conditioned** (large errors).

# CONDITIONING OF SUBSPACE IDENTIFICATION

The conditioning of the problem (\*) is determined by the singular values of the conditional covariances

$$\Sigma_{\hat{\mathbf{x}}\hat{\mathbf{x}}|\mathbf{u}^+} \quad \text{and} \quad \Sigma_{\mathbf{u}^+|\bar{\mathbf{x}}}$$

$$\Pi := E \left[ \mathbf{u}_t^+ \hat{\mathbf{x}}(t)^\top \right] \quad \bar{\Pi} := E \left[ \mathbf{u}_t^+ \hat{\mathbf{x}}(\mathbf{t})^\top \right] \quad \Lambda_u = \text{Cov} \left[ \mathbf{u}_t^+ \right]$$

$$\hat{\Pi} := L_{\mathbf{u}^+}^{-1} \Pi L_{\hat{\mathbf{x}}}^{-\top} \quad \hat{\bar{\Pi}} := L_{\mathbf{u}^+}^{-1} \bar{\Pi} L_{\bar{\mathbf{x}}}^{-\top}$$

Singular values of  $\hat{\Pi}$  are cosines of the **canonical angles** between  $\hat{\mathcal{X}}_t$  and  $\mathcal{U}_{[t,T]}$ . **Condition numbers:**

$$\kappa \left( \Sigma_{\hat{\mathbf{x}}\hat{\mathbf{x}}|\mathbf{u}^+} \right) \leq \kappa \left( \Sigma_{\hat{\mathbf{x}}} \right) \frac{1 - \sigma_{\min}^2(\hat{\Pi})}{1 - \sigma_{\max}^2(\hat{\Pi})}$$

$$\kappa\left(\Sigma_{\mathbf{u}+\mathbf{u}^+|\bar{\mathbf{x}}}\right) \leq \kappa(\Lambda_u) \frac{1}{1 - \sigma_{max}^2(\hat{\Pi})}$$

# CONDITIONING OF SUBSPACE IDENTIFICATION (back to)

- singular values of  $\hat{\Pi} =$  singular values of  $\hat{\Pi}_d \equiv E [\mathbf{u}_t^+ \hat{\mathbf{x}}_d(t)^\top]$  cosines of the canonical angles of the spaces spanned by  $\mathbf{u}_t^+$  and the **deterministic state**  $\hat{\mathbf{x}}_d(t)$
- singular values of  $\bar{\Pi} \equiv E [\mathbf{u}_t^+ \hat{\mathbf{x}}(t)^\top]$  cosines of the canonical angles of the spaces spanned by  $\mathbf{u}_t^+$  and  $\bar{\mathbf{x}}(t)$
- conditioning of the input  $\kappa(\Lambda_{\mathbf{u}^+})$  large when the amplitude of the spectrum of  $\mathbf{u}$  varies widely

## PRINCIPAL ANGLES (Canonical correlations)

Introduce Cholesky factors:

$$\Sigma_{\hat{\mathbf{x}}\hat{\mathbf{x}}} = L_{\hat{\mathbf{x}}}L_{\hat{\mathbf{x}}}^{\top} \quad \Sigma_{\mathbf{u}^+, \mathbf{u}^+} = L_{\mathbf{u}^+}L_{\mathbf{u}^+}^{\top}$$

Normalized Cross-Covariance (Correlation Matrix)

$$\Pi := L_{\mathbf{u}^+}^{-1} \text{Cov} \{ \mathbf{u}_t^+ \hat{\mathbf{x}}(t) \} L_{\hat{\mathbf{x}}}^{-\top}$$

Singular values of  $\Pi$  are cosines of the **canonical angles** between  $\hat{\mathbf{X}}_t$  and  $\mathbf{U}_{[t, T]}$ .

$$\Sigma_{\hat{\mathbf{x}}\hat{\mathbf{x}}|\mathbf{u}^+} = L_{\hat{\mathbf{x}}} \left[ I - \Pi^{\top} \Pi \right] L_{\hat{\mathbf{x}}}^{\top}$$

$$\sigma_{\text{MAX}}\{ \hat{\mathbf{X}}_t, \mathbf{U}_{[t, T]} \} \simeq 1 \Leftrightarrow \Sigma_{\hat{\mathbf{x}}\hat{\mathbf{x}}|\mathbf{u}^+} \text{ Nearly Singular !}$$

# CCA OF FUTURE INPUTS AND STATE SPACE

Given an input with assigned spectrum  $\Phi_{\mathbf{u}}$ . Which systems  $F(z)$  have the SMALLEST canonical angles of (the spaces spanned by)  $\mathbf{u}_t^+$  and  $\mathbf{x}_d(t)$  (worst conditioning of the identification problem) ??

$\sigma_k(\mathcal{X}_d, \mathcal{U}^+)$  cosines of **Canonical Angles** btw. the subspaces

$$\mathcal{U}^+ \quad \text{and} \quad \mathcal{X}_d := \text{span} \{ \mathbf{x}_d(t) \} \subset \mathcal{U}^-$$

**Fact:**

$$\sigma_k(\mathcal{X}_d, \mathcal{U}^+) \leq \sigma_k(\mathcal{U}^-, \mathcal{U}^+) \quad , \quad k = 1, 2, \dots$$

Maximal when

$$\sigma_k(\mathcal{X}_d, \mathcal{U}^+) = \sigma_k(\mathcal{U}^-, \mathcal{U}^+) \quad k = 1, 2, \dots, n_d$$

**if and only if the first  $n_d$  principal directions of  $\mathcal{U}^-$  for the pair of subspaces  $(\mathcal{U}^-, \mathcal{U}^+)$  span  $\mathcal{X}_d$ .**

## PROBING INPUTS (ASYMPTOTICS FOR $N, T \rightarrow \infty$ )

**Theorem** Assume  $\mathbf{u}$  has given rational spectral density matrix  $\Phi_u$ . The maximal canonical correlation coefficients  $\sigma_k(\mathbf{X}, \mathbf{U}^+)$  are obtained when, and only when there are  $n_d$  principal zeros of the spectral density matrix  $\Phi_u$  of  $\mathbf{u}$  cancelling all the poles of the deterministic transfer function  $F(z) = C(zI - A)^{-1}B + D$ .

How to deal with ill-conditioning? Sometimes Decoupling + Orthogonalization helps.

## ASYMPTOTIC VARIANCE OF $A, C$

**THEOREM 5** *Under standard assumptions on the true innovation noise, the estimation errors  $\tilde{A}_N := \hat{A}_N - A$ ,  $\tilde{C}_N := \hat{C}_N - C$  are asymptotically Normal,*

$$\begin{aligned} \lim_{N \rightarrow \infty} N \mathbb{E} \left\{ \text{vec}(\tilde{A}_N) \text{vec}(\tilde{A}_N)^\top \right\} &= \left\{ \Sigma_{\hat{\mathbf{x}}^c \hat{\mathbf{x}}^c}^{-1} \otimes [M H_s] \right\} \cdot \\ &\cdot \sum_{|\tau| \leq k} \Sigma_{\hat{\mathbf{x}}^c \hat{\mathbf{x}}^c}(\tau) \otimes \Sigma_{\bar{\mathbf{e}} + \bar{\mathbf{e}} +}(\tau) \cdot \left\{ \Sigma_{\hat{\mathbf{x}}^c \hat{\mathbf{x}}^c}^{-1} \otimes [M H_s] \right\}^\top \\ \lim_{N \rightarrow \infty} N \mathbb{E} \left\{ \text{vec}(\tilde{C}_N) \text{vec}(\tilde{C}_N)^\top \right\} &= \left\{ \Sigma_{\hat{\mathbf{x}}^c \hat{\mathbf{x}}^c}^{-1} \otimes [R H_s] \right\} \cdot \\ &\cdot \sum_{|\tau| < k} \Sigma_{\hat{\mathbf{x}}^c \hat{\mathbf{x}}^c}(\tau) \otimes \Sigma_{\mathbf{e} + \mathbf{e} +}(\tau) \cdot \left\{ \Sigma_{\hat{\mathbf{x}}^c \hat{\mathbf{x}}^c}^{-1} \otimes [R H_s] \right\}^\top \end{aligned}$$



## NOTATIONS

$$M := \begin{bmatrix} K & \Gamma^\dagger \end{bmatrix} - A \begin{bmatrix} \Gamma^\dagger & 0_{n \times m} \end{bmatrix} \quad R := \begin{bmatrix} I_m & 0_{m \times m(k-1)} \end{bmatrix} - C\Gamma^\dagger$$

$\Gamma$  the observability matrix in a certain basis.

$$H_s \quad : \quad = \begin{bmatrix} I & 0 & \dots & 0 & 0 \\ CK & I & \dots & 0 & 0 \\ \vdots & & & \ddots & \vdots \\ CA^{k-1}K & CA^{k-2}K & \dots & CK & I \end{bmatrix}$$

$$\mathbf{e}_t^+ \quad := \quad \begin{bmatrix} \mathbf{e}(t) \\ \mathbf{e}(t+1) \\ \vdots \\ \mathbf{e}(T-1) \end{bmatrix} \quad \bar{\mathbf{e}}_t^+ := \begin{bmatrix} \mathbf{e}_t^+ \\ \mathbf{e}(T) \end{bmatrix}$$

$$\Sigma_{\mathbf{e}^+ \mathbf{e}^+}(\tau) \quad := \quad \mathbb{E} \{ \mathbf{e}_{t+\tau}^+ (\mathbf{e}_t^+)^{\top} \} \quad \Sigma_{\bar{\mathbf{e}}^+ \bar{\mathbf{e}}^+}(\tau) = \mathbb{E} \{ \bar{\mathbf{e}}_{t+\tau}^+ (\bar{\mathbf{e}}_t^+)^{\top} \}$$

Valid for N4SID, MOESP, and also CCA.

- $\Sigma_{\hat{\mathbf{x}}^c \hat{\mathbf{x}}^c}^{-1} = \Sigma_{\hat{\mathbf{x}} \hat{\mathbf{x}} | \mathbf{u}^+}^{-1}$  Very “large” for ill-conditioned problems, the variance of the estimation errors will also be large.
- No (or white) input:  $\Sigma_{\hat{\mathbf{x}} \hat{\mathbf{x}} | \mathbf{u}^+} \equiv \Sigma_{\hat{\mathbf{x}} \hat{\mathbf{x}}}$

## PREVIOUS AVAILABLE RESULTS

[Bauer, Bauer-Ljung, Bauer-Jansson]: asymptotic formulas valid for  $N \rightarrow \infty$   
AND  $p := t - t_0$  (past data horizon), tending to infinity with  $N$  at a certain rate

Estimates neglect transient due to FINITE-INTERVAL DATA. Consistency only for  $p \rightarrow \infty$

Different asymptotic formulas for different methods, CCA, MOESP, N4SID etc. Complicated and difficult to use.

Asymptotic formulas should be valid for FINITE  $p$  and “transient” estimates ( in practice can only regress on **finite past**). Stationary approximations are biased for finite  $p$ .

## APPLICATIONS

Assume for simplicity that  $A$  has simple eigenvalues.

here is an eigenvalue  $\lambda^i$  of  $A$  such that the difference between the  $i$ -th eigenvalue of  $\hat{A}_N$ ,  $\hat{\lambda}_N^i$ , and  $\lambda^i$ , satisfies

$$\hat{\lambda}_N^i - \lambda^i \simeq \frac{v_i^\top \tilde{A}_N u_i}{v_i^\top u_i} + O(\|\tilde{A}_N\|^2)$$

where  $v_i$  and  $u_i$  are the normalized left and right eigenvectors of  $A$  corresponding to  $\lambda^i$ .

$$N\mathbb{E}(\hat{\lambda}_N^i - \lambda^i)^2 = \frac{1}{(v_i^\top u_i)^2} (u_i^\top \otimes v_i^\top) N\mathbb{E} \left\{ \text{vec}(\tilde{A}_N) \text{vec}(\tilde{A}_N)^\top \right\} (u_i \otimes v_i)$$

Note that  $(v_i^\top u_i)^2$  is the square of the cosine of the angle between the two eigenvectors and is equal to one if the matrix  $A$  is symmetric (in which case  $v_i = u_i$ ).

## ASYMPTOTIC VARIANCE OF $(B, D)$

The vectorized parameter estimates  $\text{vec}(\hat{\mathcal{K}}_{1,N})$   $\text{vec}(\mathcal{K}_{2,N})$  form an asymptotically Gaussian sequence

$$\text{AsVar} \left( \sqrt{N} \text{vec}(\hat{\mathcal{K}}_{1,N}) \right) = \bar{G} \left\{ \sum_{|\tau| \leq k} \Sigma_{\mathbf{u}^+ \mathbf{u}^+ | \bar{\mathbf{x}}}(\tau) \otimes \Sigma_{\mathbf{e}^+ \mathbf{e}^+}(\tau) \right\} \bar{G}^\top$$

$$\text{AsVar} \left( \sqrt{N} \text{vec}(\hat{\mathcal{K}}_{2,N}) \right) = G \left\{ \sum_{|\tau| < k} \Sigma_{\mathbf{u}^+ \mathbf{u}^+ | \bar{\mathbf{x}}}(\tau) \otimes \Sigma_{\mathbf{e}^+ \mathbf{e}^+}(\tau) \right\} G^\top$$

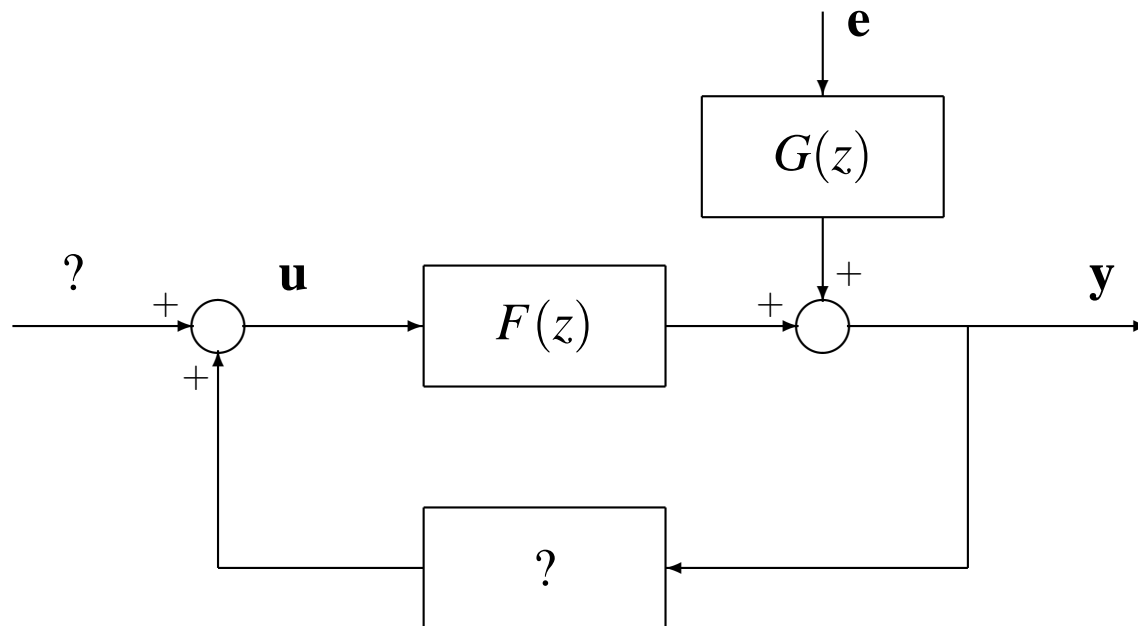
$$G := \Sigma_{\mathbf{u}^+ \mathbf{u}^+ | \bar{\mathbf{x}}}^{-1} \otimes [RH_s], \quad \bar{G} := \Sigma_{\mathbf{u}^+ \mathbf{u}^+ | \bar{\mathbf{x}}}^{-1} \otimes [M\bar{H}_s]$$

$R$  and  $M$  being as before, and,

$$\Sigma_{\bar{\mathbf{u}}^+ \bar{\mathbf{u}}^+ | \bar{\mathbf{x}}}(\tau) := \mathbb{E} \left\{ \tilde{\mathbf{u}}_{t+\tau}^+ (\tilde{\mathbf{u}}_t^+)^{\top} \right\}, \quad \Sigma_{\mathbf{e}^+ \mathbf{e}^+}(\tau) = E \left\{ \mathbf{e}_{t+\tau}^+ (\mathbf{e}_t^+)^{\top} \right\}$$

$\tilde{\mathbf{u}}_{t+\tau}^+$  the  $\tau$ -steps ahead stationary shift of the random vector  $\tilde{\mathbf{u}}_t^+ := \mathbf{u}_t^+ - \mathbb{E} [\mathbf{u}_t^+ | \bar{\mathbf{x}}(t)]$ .

# SUBSPACE IDENTIFICATION WITH FEEDBACK



$+ F(\infty) = 0.$

## PROBLEMS WITH STATE CONSTRUCTION

$$\mathbf{y}(t+h) = CA^h \mathbf{x}(t) + \text{“terms in } \mathbf{U}_t^+ \text{”} + \text{“terms in } \mathbf{E}_t^+ \text{”} \quad h = 0, 1, \dots, k$$

Classical (N4SID, CVA, MOESP) construct the state space via the oblique projection

$$E_{\parallel \mathbf{U}_t^+} [\mathbf{Y}_t^+ \mid \mathbf{Y}_t^- \vee \mathbf{U}_t^-]$$

Needs  $\mathbf{E}_t^+ \perp \mathbf{U}_t^+$  which is equivalent to *Absence of Feedback* from  $\mathbf{y}$  to  $\mathbf{u}$ .  
(Granger)

Need an **alternative way to construct the state space**, see the discussion in *Ljung-McKelvey 1996*

## REMEDY (Jansson 2003/Chiuso-Picci 2004)

**FACT:**  $\mathbf{x}(t)$  is also the state space of the **predictor model**

$$\begin{cases} \mathbf{x}(t+1) &= (A - KC)\mathbf{x}(t) + Bu(t) + Ky(t) \\ \hat{\mathbf{y}}(t | t-1) &= C\mathbf{x}(t) \end{cases}$$

$$\hat{\mathbf{y}}(t+h | t) = C(A - KC)^h \mathbf{x}(t) + \text{“terms in } \mathbf{U}_t^+ \vee \mathbf{Y}_t^+ \text{”}$$

$$\mathbf{X}_t^{+/-} = E_{\|\mathbf{U}_t^+ \vee \mathbf{Y}_t^+} [\hat{\mathbf{Y}}_t^+ | \mathbf{U}_t^- \vee \mathbf{Y}_t^-]$$

**Jansson 2003** Compute predictor space removing the effect of undesired terms pre-estimating Markov parameters of predictor using an ARX model.



## “PREDICTOR IDENTIFICATION ALGORITHM:

1. Compute the oblique predictors

$$\hat{\mathbf{y}}(t+h | t) := E_{\| \mathbf{U}_{[t,t+h]} \vee \mathbf{Y}_{[t,t+h]}} \left[ \mathbf{y}(t+h) | \mathbf{Y}_{[t_0,t]} \vee \mathbf{U}_{[t_0,t]} \right]$$

2. Compute  $\hat{\mathbf{X}}_t^{+/-}$  as “best”  $n$ -dimensional approximation of the space spanned by  $\hat{\mathbf{y}}(t+h | t)$ ,  $h = 0, \dots, k$ , repeat for  $\hat{\mathbf{X}}_{t+1}^{+/-}$
3. Solve regression in the least squares sense to get  $\hat{A}$ ,  $\hat{B}$ ,  $\hat{C}$ ,  $\hat{K}$ .

## COMMENTS:

- The classical subspace procedure to construct the state space turns out to be WRONG if data are collected in closed-loop.
- Subspace methods based on the **predictor model** work also with feedback !
- Predictor is always stable (joint spectrum bounded away from zero  $\Rightarrow |\lambda(A - KC)| < 1.$ )
- Ideally predictor space can be constructed without any assumption on feedback channel.

## REMARKS

1. Predictor identification “ideally” yields consistent estimators
2. Practically need to work with **finite past** starting from a certain time  $t_0$ .
3. If number of data points ( $[y_t, y_{t+1}, \dots, y_{t+N}]$ )  $N \rightarrow \infty$ , but  $t - t_0$  **fixed and finite** Consistency not guaranteed.
4. “Transient” predictors (transient Kalman filter) involve also the dynamics of  $\mathbf{u}$  !

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