AA216/CME345: MODEL REDUCTION

Balanced Truncation

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These slides are based on the recommended textbook: A.C. Antoulas, "Approximation of Large-Scale Dynamical Systems," Advances in Design and Control, SIAM, ISBN-0-89871-529-6

Outline

- Consider the following **stable**, high-dimensional, LTI system

$$\boxed{\begin{aligned} \frac{d\mathbf{w}}{dt}(t) &= \mathbf{A}\mathbf{w}(t) + \mathbf{B}\mathbf{u}(t) \\ \mathbf{y}(t) &= \mathbf{C}\mathbf{w}(t) \\ \mathbf{w}(0) &= \mathbf{w}_0 \end{aligned}}$$

  - $\mathbf{w} \in \mathbb{R}^N$: State variables
  - $\mathbf{u} \in \mathbb{R}^p$: Input variables, typically $p \ll N$
  - $\mathbf{y} \in \mathbb{R}^q$: Output variables, typically $q \ll N$

- Recall that the solution $\mathbf{w}(t)$ of the above linear ODE can be written as

$$\boxed{\mathbf{w}(t) = \phi(t, \mathbf{u}; t_0, \mathbf{w}_0) = e^{\mathbf{A}(t-t_0)}\mathbf{w}(t_0) + \int_{t_0}^{t} e^{\mathbf{A}(t-\tau)}\mathbf{B}\mathbf{u}(\tau)d\tau, \quad \forall t \geq t_0}$$

Definition (1)

A state $\mathbf{w} \in \mathbb{R}^N$ is **reachable** if there exists an input function $\mathbf{u}(.)$ of finite energy and a time $T < \infty$ such that under this input and zero initial condition, the state of the system is $\mathbf{w}$

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Definition (2)

A state $\mathbf{w} \in \mathbb{R}^N$ is **controllable** to the zero state if there exist an input function $\mathbf{u}(.)$ and a time $T < \infty$ such that

$$\phi(T, \mathbf{u}; 0, \mathbf{w}) = \mathbf{0}_N$$

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Definition (3)

A state $\mathbf{w} \in \mathbb{R}^N$ is **unobservable** if for all $t \geq 0$,

$$\mathbf{y}(t) = \mathbf{C}\phi(t, \mathbf{0}; 0, \mathbf{w}) = \mathbf{0}_q$$

Definition (4, - R.E. Kalman, 1963)

A linear dynamical system $(\mathbf{A}, \mathbf{B}, \mathbf{C})$ is **completely controllable** at time t_0 if it is not equivalent, for all $t \geq t_0$, to a system of the type

$$\begin{aligned} \frac{d\mathbf{w}^{(1)}}{dt} &= \mathbf{A}^{(1,1)}\mathbf{w}^{(1)} + \mathbf{A}^{(1,2)}\mathbf{w}^{(2)} + \mathbf{B}^{(1)}\mathbf{u} \\ \frac{d\mathbf{w}^{(2)}}{dt} &= \mathbf{A}^{(2,2)}\mathbf{w}^{(2)} \\ \mathbf{y}(t) &= \mathbf{C}^{(1)}\mathbf{w}^{(1)} + \mathbf{C}^{(2)}\mathbf{w}^{(2)} \end{aligned}$$

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- This definition can be extended to linear time-variant systems

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A linear dynamical system $(\mathbf{A}, \mathbf{B}, \mathbf{C})$ is **completely observable** at time t_0 if it is not equivalent, for all $t \leq t_0$, to any system of the type

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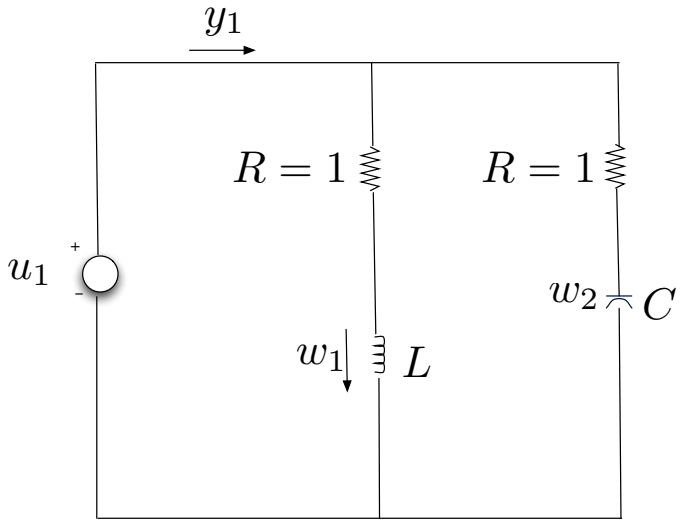

- The equation of the network in terms of the current w_1 flowing through the inductor and the potential w_2 across the capacitor ($LC = R = 1$) is given by

$$\begin{aligned} \frac{dw_1}{dt} &= -\frac{1}{L}w_1 + u_1 \\ \frac{dw_2}{dt} &= -\frac{1}{C}w_2 + u_1 \\ y_1 &= \frac{1}{L}w_1 + \frac{1}{C}w_2 + u_1 \end{aligned}$$

- Under the change of variable $\tilde{w}_1 = (w_1 + w_2)/2$ and $\tilde{w}_2 = (w_1 - w_2)/2$, the previous dynamical system becomes

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- $\tilde{w}_1$ is controllable but not observable
- $\tilde{w}_2$ is observable but not controllable

Theorem (Kalman 1961)

Consider a dynamical system $(\mathbf{A}, \mathbf{B}, \mathbf{C})$. *Then:*
(i) There is a state space coordinate system in which the components of the state vector can be decomposed into four parts

$$\mathbf{w} = [\mathbf{w}^{(a)} \ \mathbf{w}^{(b)} \ \mathbf{w}^{(c)} \ \mathbf{w}^{(d)}]^T$$

(ii) The sizes N_a, N_b, N_c *and* N_d *of these vectors do not depend on the choice of basis*
(iii) The system matrices take the form

$$\mathbf{A} = \begin{bmatrix} \mathbf{A}^{(a,a)} & \mathbf{A}^{(a,b)} & \mathbf{A}^{(a,c)} & \mathbf{A}^{(a,d)} \\ 0 & \mathbf{A}^{(b,b)} & 0 & \mathbf{A}^{(b,d)} \\ 0 & 0 & \mathbf{A}^{(c,c)} & \mathbf{A}^{(c,d)} \\ 0 & 0 & 0 & \mathbf{A}^{(d,d)} \end{bmatrix}, \mathbf{B} = \begin{bmatrix} \mathbf{B}^{(a)} \\ \mathbf{B}^{(b)} \\ 0 \\ 0 \end{bmatrix},$$

$$\mathbf{C} = \begin{bmatrix} 0 & \mathbf{C}^{(b)} & 0 & \mathbf{C}^{(d)} \end{bmatrix}$$

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  - part (d) is uncontrollable and completely observable
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Definition (6)

The **reachable subspace** $\mathbb{W}_{reach} \subset \mathbb{R}^N$ of a system $(\mathbf{A}, \mathbf{B}, \mathbf{C})$ is the set containing all reachable states of the system.

$$\mathcal{R}(\mathbf{A}\ \mathbf{B}) = [\mathbf{B}\ \mathbf{AB}\ \cdots\ \mathbf{A}^{N-1}\mathbf{B}\ \cdots]$$

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Definition (7)

The **controllable subspace** $\mathbb{W}_{contr} \subset \mathbb{R}^N$ of a system $(\mathbf{A}, \mathbf{B}, \mathbf{C})$ is the set containing all controllable states of the system

Theorem

Consider the system $(\mathbf{A}, \mathbf{B}, \mathbf{C})$. *Then*

$$\mathbb{W}_{reach} = Im\ \mathcal{R}(\mathbf{A}, \mathbf{B})$$

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Corollary

(i) $\mathbf{A}\mathbb{W}_{reach} \subset \mathbb{W}_{reach}$
(ii) The system is completely reachable if and only if rank $\mathcal{R}(\mathbf{A}, \mathbf{B}) = N$
(iii) Reachability is basis independent

Definition

The **reachability (controllability) Gramian** at time $t < \infty$ is defined as the $N \times N$ symmetric positive semi-definite matrix

$$\mathcal{P}(t) = \int_0^t e^{\mathbf{A}\tau} \mathbf{B}\mathbf{B}^* e^{\mathbf{A}^*\tau} d\tau$$

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Definition

The **observability Gramian** at time $t < \infty$ is defined as the $N \times N$ symmetric positive semi-definite matrix

$$\mathcal{Q}(t) = \int_0^t e^{\mathbf{A}^*\tau} \mathbf{C}^* \mathbf{C} e^{\mathbf{A}\tau} d\tau$$

Proposition

The columns of $\mathcal{P}(t)$ *span the reachability subspace* $\mathbb{W}_{reach} = Im\ \mathcal{R}(\mathbf{A}, \mathbf{B})$

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Corollary

A system $(\mathbf{A}, \mathbf{B}, \mathbf{C})$ is reachable if and only if $\mathcal{P}(t)$ is symmetric positive definite at some time $t > 0$

Theorem

For continuous linear dynamical systems, the notions of controllability and reachability are equivalent – that is,

$$\mathbb{W}_{reach} = \mathbb{W}_{contr}$$

Definition

The **unobservability subspace** $\mathbb{W}_{unobs} \subset \mathbb{R}^N$ is the set of all unobservable states of the system. The matrix

$$\mathcal{O}(\mathbf{C}, \mathbf{A}) = [\mathbf{C}^* \ \mathbf{A}^*\mathbf{C}^* \ \cdots \ (\mathbf{A}^*)^i\mathbf{C}^* \ \cdots]^*$$

is the **observability matrix** of the system.

Theorem

Consider the system $(\mathbf{A}, \mathbf{B}, \mathbf{C})$. *Then*

$$\mathbb{W}_{unobs} = Ker\ \mathcal{O}(\mathbf{C}, \mathbf{A})$$

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Corollary

(i) $\mathbf{A}\mathbb{W}_{unobs} \subset \mathbb{W}_{unobs}$
(ii) The system is completely observable if and only if rank $\mathcal{O}(\mathbf{C}, \mathbf{A}) = N$
(iii) Observability is basis independent

Definition

The **infinite reachability (controllability) Gramian** is defined for a **stable** LTI system as the $N \times N$ symmetric positive semi-definite matrix

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The **infinite observability Gramian** is defined for a **stable** LTI system as the $N \times N$ symmetric positive semi-definite matrix

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$$\boxed{\mathcal{P} = \frac{1}{2\pi}\int_{-\infty}^{\infty}(j\omega\mathbf{I}_N - \mathbf{A})^{-1}\mathbf{B}\mathbf{B}^*(-j\omega\mathbf{I}_N - \mathbf{A}^*)^{-1}d\omega}$$

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$$\|\mathbf{w}\|^2_{\mathcal{P}^{-1}} = \mathbf{w}^T \mathcal{P}^{-1} \mathbf{w} \qquad \left(\leq \int_0^t (\mathbf{B}\mathbf{u}(\tau))^* \, \mathbf{B}\mathbf{u}(\tau) d\tau \right)$$

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- The inner product based on $\mathcal{Q}$ indicates the maximal energy produced by observing the output of the system corresponding to an initial state $\mathbf{w}_0$ when no input is applied

$$\|\mathbf{w}\|^2_{\mathcal{Q}} = \mathbf{w}^T \mathcal{Q} \mathbf{w}$$

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- The above notions are **basis-dependent**
- One would like to consider a basis where both concepts are **equivalent**, i.e., where the system is **balanced**

- Balancing requires changing the basis for the state using a transformation $\mathbf{T} \in \mathrm{GL}(N)$

$$\tilde{\mathbf{w}} = \mathbf{T}\mathbf{w}$$

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  - the reachability Gramian becomes

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  - the reachability Gramian becomes

$$\boxed{\tilde{\mathcal{P}} = \mathbf{T}\mathcal{P}\mathbf{T}^*}$$

  - $\mathbf{C}e^{\mathbf{A}t} \Rightarrow \mathbf{C}\mathbf{T}^{-1}\mathbf{T}e^{\mathbf{A}t}\mathbf{T}^{-1} = \mathbf{C}e^{\mathbf{A}t}\mathbf{T}^{-1}$

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- The balancing transformations $\mathbf{T}_{\text{bal}}$ and $\mathbf{T}_{\text{bal}}^{-1}$ can be computed as follows
  1. compute the Cholesky factorization $\mathcal{P} = \mathbf{U}\mathbf{U}^*$
  2. compute the eigenvalue decomposition of $\mathbf{U}^*\mathcal{Q}\mathbf{U}$

     $$\mathbf{U}^*\mathcal{Q}\mathbf{U} = \mathbf{K}\boldsymbol{\Sigma}^2\mathbf{K}^*$$

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- Then, one can check that balancing is achieved

$$\boxed{\mathbf{T}_{\text{bal}}\mathcal{P}\mathbf{T}_{\text{bal}}^* = \mathbf{T}_{\text{bal}}^{-*}\mathcal{Q}\mathbf{T}_{\text{bal}}^{-1} = \mathbf{\Sigma}}$$

Definition (Hankel Singular Values)

$\mathbf{\Sigma} = \text{diag}(\sigma_1, \cdots, \sigma_N)$ contains the N **Hankel singular values of the system**

- Computing the balancing transformation $\mathbf{T}_{\text{bal}}$ is equivalent to minimizing the following function

$$\boxed{\min_{\mathbf{T}\in \mathrm{GL}(N)} f(\mathbf{T}) = \min_{\mathbf{T}\in \mathrm{GL}(N)} \mathrm{trace}(\mathbf{T}\mathcal{P}\mathbf{T}^* + \mathbf{T}^{-*}\mathcal{Q}\mathbf{T}^{-1})}$$

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- For the optimal transformation $\mathbf{T}_{\text{bal}}$, the function takes the value

$$f(\mathbf{T}_{\text{bal}}) = 2\text{tr}(\mathbf{\Sigma}) = 2\sum_{i=1}^{N} \sigma_i$$

where $\{\sigma_i\}_{i=1}^{N}$ are the Hankel singular values

- Applying the change of variable $\tilde{\mathbf{w}} = \mathbf{T}_{\text{bal}}\mathbf{w}$ transforms the given dynamical system into $(\tilde{\mathbf{A}}, \tilde{\mathbf{B}}, \tilde{\mathbf{C}})$ where

$$\mathbf{T}_{\text{bal}}\mathbf{A}\mathbf{T}_{\text{bal}}^{-1} = \tilde{\mathbf{A}}, \quad \mathbf{T}_{\text{bal}}\mathbf{B} = \tilde{\mathbf{B}}, \quad \mathbf{C}\mathbf{T}_{\text{bal}}^{-1} = \tilde{\mathbf{C}}$$

- Let $1 \leq k \leq N$; the system can be partitioned in blocks as

$$\tilde{\mathbf{A}} = \begin{bmatrix} \tilde{\mathbf{A}}_{11} & \tilde{\mathbf{A}}_{12} \\ \tilde{\mathbf{A}}_{21} & \tilde{\mathbf{A}}_{22} \end{bmatrix}, \qquad \tilde{\mathbf{B}} = \begin{bmatrix} \tilde{\mathbf{B}}_1 \\ \tilde{\mathbf{B}}_2 \end{bmatrix}, \qquad \tilde{\mathbf{C}} = \begin{bmatrix} \tilde{\mathbf{C}}_1 & \tilde{\mathbf{C}}_2 \end{bmatrix}$$

- The subscripts 1 and 2 denote the dimensions k and $N - k$, respectively

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- Then, the following lower-dimensional model of size k

$$(\mathbf{A}_r, \mathbf{B}_r, \mathbf{C}_r) = (\tilde{\mathbf{A}}_{11}, \tilde{\mathbf{B}}_1, \tilde{\mathbf{C}}_1) \in \mathbb{R}^{k\times k} \times \mathbb{R}^{k\times p} \times \mathbb{R}^{q\times k}$$

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- The **left and right** ROBs are

$$\mathbf{W} = \mathbf{T}^*_{\text{bal}}(:, 1:k) \quad \text{and} \quad \mathbf{V} = \mathbf{S}_{\text{bal}}(:, 1:k), \text{ respectively,}$$

  where $\mathbf{S}_{\text{bal}} = \mathbf{T}^{-1}_{\text{bal}}$

Definition (The Hardy space $\mathcal{H}_\infty$)

The $\mathcal{H}_\infty$-norm associated with a system characterized by a transfer function $\mathbf{G}(\cdot)$ is defined as

$$\|\mathbf{G}\|_{\mathcal{H}_\infty} = \sup_{\mathbf{z}\in\mathbb{C}_+} \|\mathbf{G}(z)\|_\infty = \sup_{\mathbf{z}\in\mathbb{C}_+} \sigma_{\max}(\mathbf{G}(z))$$

where $z \in \mathbb{C}_+$ if $z \in \mathbb{C}$ and $\mathfrak{Im}(z) > 0$.

Proposition

$$(i) \quad \|\mathbf{G}\|_{\mathcal{H}_\infty} = \sup_{\omega\in\mathbb{R}} \|\mathbf{G}(i\omega)\|_\infty = \sup_{\omega\in\mathbb{R}} \sigma_{\max}(\mathbf{G}(i\omega))$$

$$(ii) \quad \|\mathbf{G}\|_{\mathcal{H}_\infty} = \sup_{\mathbf{u}\neq\mathbf{0}} \frac{\|\mathbf{y}(\cdot)\|_2}{\|\mathbf{u}(\cdot)\|_2} = \sup_{\mathbf{u}\neq\mathbf{0}} \sqrt{\frac{\int_0^\infty \|\mathbf{y}(t)\|_2^2 dt}{\int_0^\infty \|\mathbf{u}(t)\|_2^2 dt}}$$

- The $\mathcal{H}_\infty$ norm of the error between the HDM- and PROM-based solutions will be used as an error criterion

Theorem (Error Bounds)

The Balanced Truncation procedure yields the following error bound for the output of interest.
Let $\{\tilde{\sigma}_i\}_{i=1}^{N_{SV}} \subseteq \{\sigma_i\}_{i=1}^{N}$ *denote the distinct Hankel singular values of the system and* $\{\tilde{\sigma}_i\}_{i=N_k+1}^{N_{SV}}$ *the ones that have been truncated. Then*

$$\|\mathbf{y}(\cdot) - \mathbf{y}_r(\cdot)\|_2 \leq 2 \sum_{i=N_k+1}^{N_{SV}} \tilde{\sigma}_i \|\mathbf{u}(\cdot)\|_2$$

Equivalently, the above result can be written in terms of the $\mathcal{H}_\infty$*-norm of the system error as follows*

$$\|\mathbf{G}(\cdot) - \mathbf{G}_r(\cdot)\|_{\mathcal{H}_\infty} \leq 2 \sum_{i=N_k+1}^{N_{SV}} \tilde{\sigma}_i$$

where $\mathbf{G}$ *and* $\mathbf{G}_r$ *are the full- and reduced-order transfer functions. Equality holds when* $\tilde{\sigma}_{N_k+1} = \tilde{\sigma}_{N_{SV}}$ *(all truncated singular values are equal).*

Proof. The proof proceeds in 3 steps:

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3. Use this result to show that in the general case

$$\|\mathbf{E}(\cdot)\|_{\mathcal{H}_\infty} \leq 2 \sum_{i=N_k+1}^{n_{SV}} \tilde{\sigma}_i$$

Theorem (Stability Preservation)

Consider $(\mathbf{A}_r, \mathbf{B}_r, \mathbf{C}_r) = (\tilde{\mathbf{A}}_{11}, \tilde{\mathbf{B}}_1, \tilde{\mathbf{C}}_1)$*, a PROM obtained by Balanced Truncation. Then*
(i) $\mathbf{A}_r = \tilde{\mathbf{A}}_1$ *has no eigenvalues in the open right half plane*

(ii) Furthermore, if the systems $(\tilde{\mathbf{A}}_{11}, \tilde{\mathbf{B}}_1, \tilde{\mathbf{C}}_1)$ *and* $(\tilde{\mathbf{A}}_{22}, \tilde{\mathbf{B}}_2, \tilde{\mathbf{C}}_2)$ *have no Hankel singular values in common,* $\mathbf{A}_r$ *has no eigenvalues on the imaginary axis*

- Because of numerical stability issues, computing the transformations

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  3. construct the matrices

$$\boxed{\mathbf{T}_{\text{bal}} = \mathbf{\Sigma}^{-\frac{1}{2}}\mathbf{V}^*\mathbf{Z}^* \quad \text{and} \quad \mathbf{T}_{\text{bal}}^{-1} = \mathbf{U}\mathbf{W}\mathbf{\Sigma}^{-\frac{1}{2}}}$$

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  - the computation of a Gramian is intensive as it requires the solution of a Lyapunov equation ($\mathcal{O}(N^3)$ operations)
  - for this reason, Balanced Truncation is in general impractical for large systems – say $N \gtrsim 10^5$ (but monitor progress in the literature if interested)

Recall the theorem underlying the construction of a POD basis

Theorem

Let $\widehat{\mathbf{K}} \in \mathbb{R}^{N\times N}$ *be the real symmetric semi-definite positive matrix defined as*

$$\hat{\mathbf{K}} = \int_0^{\mathcal{T}} \mathbf{w}(t)\mathbf{w}(t)^T dt$$

Let $\hat{\lambda}_1 \geq \hat{\lambda}_2 \geq \cdots \geq \hat{\lambda}_N \geq 0$ *denote the ordered eigenvalues of* $\widehat{\mathbf{K}}$ *and let* $\widehat{\phi}_i \in \mathbb{R}^N,\ i = 1, \cdots, N$, *denote the associated eigenvectors*

$$\widehat{\mathbf{K}}\widehat{\phi}_i = \hat{\lambda}_i \widehat{\phi}_i,\ i = 1, \cdots, N.$$

The subspace $\widehat{\mathcal{V}} = range(\widehat{\mathbf{V}})$ *of dimension* k *minimizing* $J(\mathbf{\Pi}_{\mathbf{V},\mathbf{V}})$ *is the invariant subspace of* $\widehat{\mathbf{K}}$ *associated with the eigenvalues* $\hat{\lambda}_1, \cdots, \hat{\lambda}_k$

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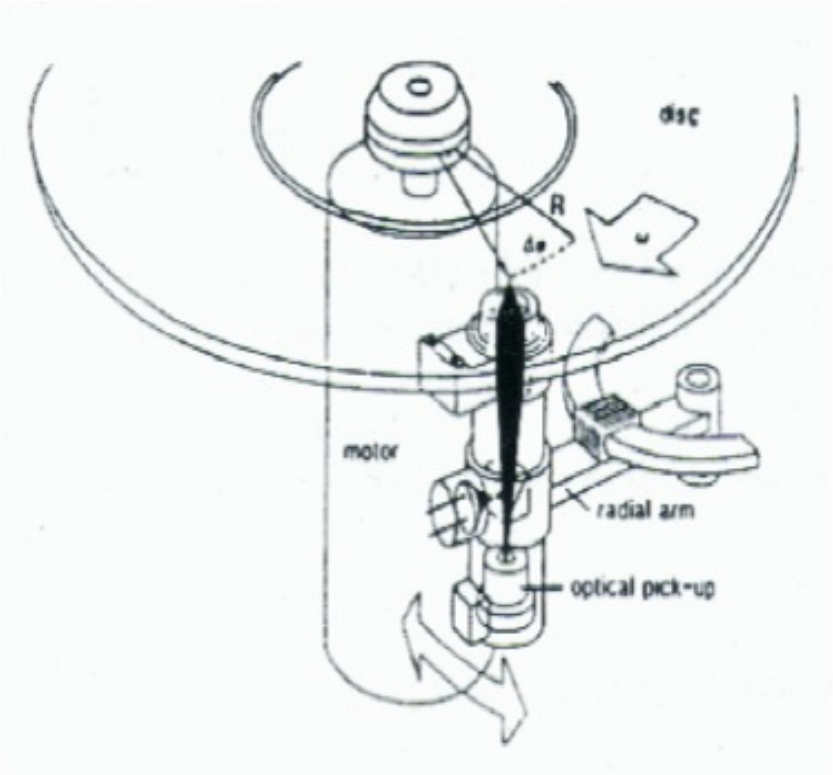
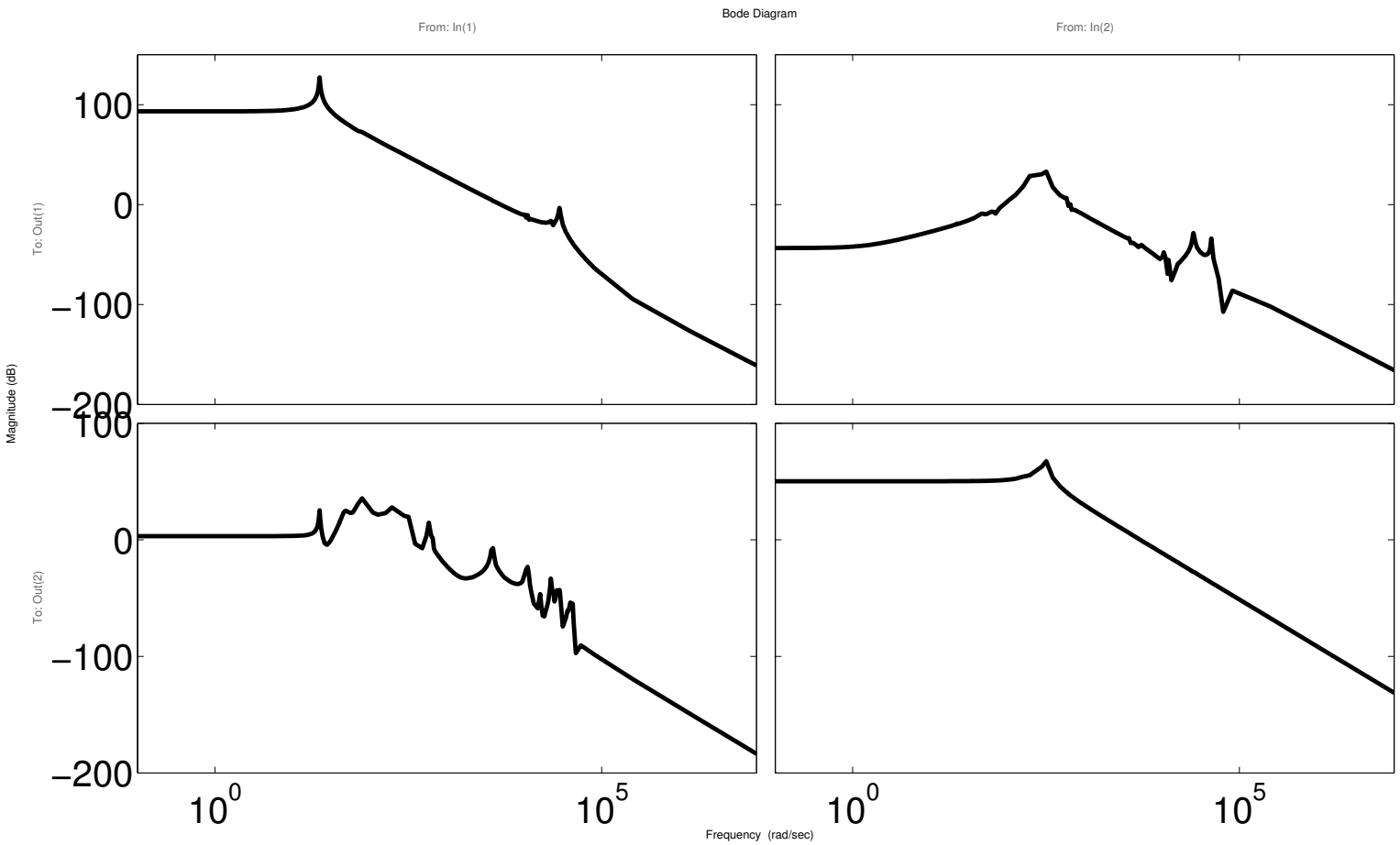
- Unlike the Balanced Truncation method, the POD method does not take into account the observability Gramian to determine the PROM: therefore, every observable state may be truncated

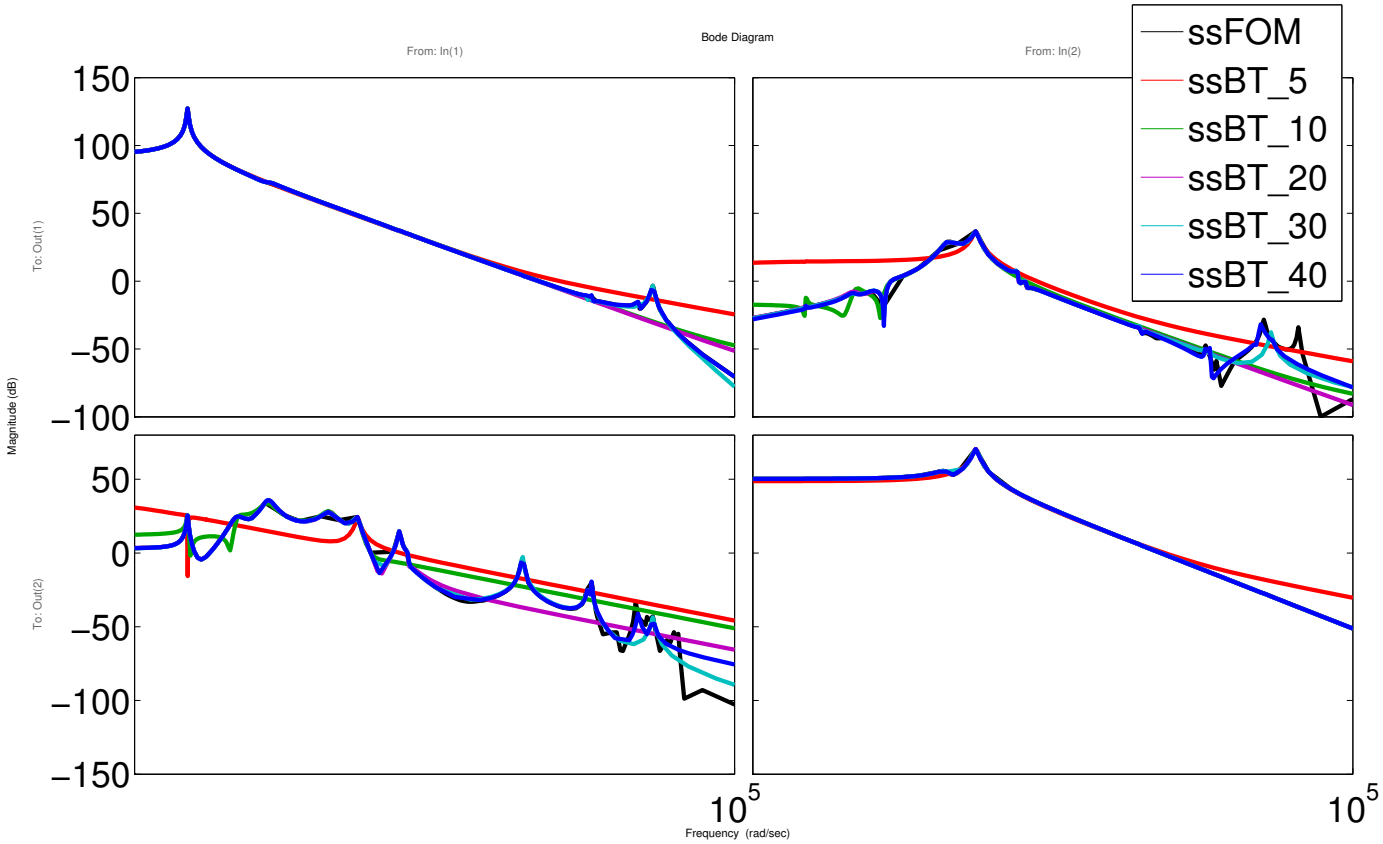


- Objective: model the position of the lens controlled by a swing arm
- System with $p = 2$ inputs and $q = 2$ outputs

Bode plots associated with the HDM-based solution ($N = 120$): Each column represents one input and each row represents a different output



Bode plots associated with the PROM-based (Balanced Truncation) solution: Each column represents one input and each row represents a different output



- The Balanced POD (BPOD) method generates two sets of snapshots: the standard POD solution snapshots; and the dual POD snapshots introduced below

$$\begin{aligned} \mathbf{S} &= \left[(j\omega_1\mathbf{I}-\mathbf{A})^{-1}\mathbf{B} \ \cdots \ (j\omega_{N_{\text{snap}}}\mathbf{I}-\mathbf{A})^{-1}\mathbf{B}\right] \\ \mathbf{S}_{\text{dual}} &= \left[(-j\omega_1\mathbf{I}-\mathbf{A}^*)^{-1}\mathbf{C}^* \ \cdots \ (-j\omega_{N_{\text{snap}}}\mathbf{I}-\mathbf{A}^*)^{-1}\mathbf{C}^*\right] \end{aligned}$$

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- Next, BPOD computes right and left ROBs as follows

$$\begin{aligned} \mathbf{S}_{\text{dual}}^T\mathbf{S} &= \mathbf{U}\boldsymbol{\Sigma}\mathbf{Z}^T \ \ \text{(SVD)} \\ \mathbf{V} &= \mathbf{S}\mathbf{Z}_k\boldsymbol{\Sigma}_k^{-1/2} \\ \mathbf{W} &= \mathbf{S}_{\text{dual}}\mathbf{U}_k\boldsymbol{\Sigma}_k^{-1/2} \end{aligned}$$

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  where the subscript k designates the first k terms of the singular value decomposition
- If no truncation is performed, the result is equivalent to two-sided moment matching at $s_i \in \{\omega_1, \cdots, \omega_{N_{\text{snap}}}\}$ (see later)

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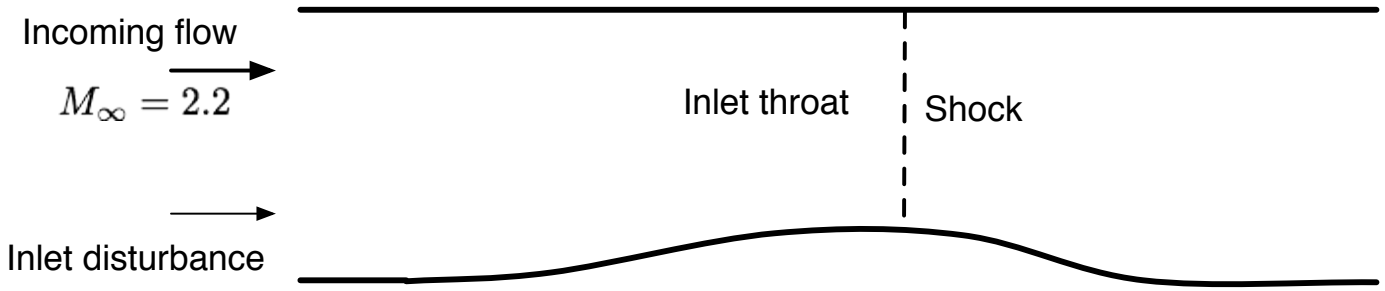
- The POD method in the time domain is based solely on the reachability concept
- However, the BPOD method
  - introduces also the notion of observability in the construction of a PROM
  - is tractable for very large-scale systems
  - provides an approximation to the Balanced Truncation method
  - does not however guarantee in general the stability of the resulting PROM

- Supersonic Inlet Problem (part of the Oberwolfach Model Reduction Benchmark Collection repository)



$$\begin{aligned} \mathbf{E}\frac{d\mathbf{w}}{dt}(t) &= \mathbf{A}\mathbf{w}(t) + \mathbf{B}\mathbf{u}(t) \\ \mathbf{y}(t) &= \mathbf{C}\mathbf{w}(t) \end{aligned}$$

- $N = 11,370$ (2D Euler equations)
- $p = 1$ input (density disturbance of the inlet flow)
- $q = 1$ output (average Mach number at the diffuser throat)

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- In both cases, same frequency sampling for the computation of solution snapshots
- Plot of the magnitude of the relative error in the transfer function (within the sampled frequency interval) as a function of the dimension k of the constructed PROM