# Linear Dynamic Systems and Control 

## 최신 제어 기법

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$$
\begin{aligned}
\dot{x}(t) & =A(t) x(t)+B(t) u(t) \\
y(t) & =C(t) x(t)+D(t) u(t)
\end{aligned}
$$

선형 제어 시스템을 학습하기 위한 기초 및 기본 개념을 제공하는 강좌로, system model에 대한 기본 개념을 상태변수 공간에서 학습하고, 선형대수학의 기초를 다시 살펴 본 후, 상태공간 방정식의 해, 시스템의 안정도 (Lyapunov stability), 제어 가능한 시스템 (controllability), 관측 가능한 시스템 (observability), controllable and observable canonidal form, duality, 상태변수 궤환 제어기 설계, 상태변수 관측기 설계, 및 출력 궤환 제어 방법을 학습한다.

1. Systems and State
2. Vector Spaces
2.1 Vector Spaces
2.2 Linear Transformations and Matrices
2.3 Eigenvalues and Eigenvectors, Diagonalization
2.4 Cayley-Hamilton THM, Functions of a Square Matrix
2.5 Inner Product Spaces, Normed Spaces
3. State Space
3.1 State Transition matrix
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5. Controllability and Observability
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5.2 Controllability of Linear Systems
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5.5 Observability of Linear Time-Invariant Systems
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6. State Feedback and Observers
6.1 State Feedback
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6.3 Feedback Control Systems Using Observers
7. Realization
7.1 Minimal Realizations
7.2 Controllable and Observable Canonical Forms
7.3 Realizability of Transfer Matrices

Textbook:
Lecture note written by Prof. Jin Heon Seo
References:
W. Brogan, Modern Control Theory, 3rd Edition, Prentice-Hall, 1990.
C.T. Chen, Linear System Theory and Design, 4th Edition, Oxford publishing, 2012.

## Definition (Causality)

A system is called causal (or nonanticipative) if, for any $t, y(t)$ does not depend on any $u\left(t_{1}\right)$ for $t_{1}>t$. A system is called noncausal (or anticipative) if it is not causal.

## Definition (Dynamic Systems)

A system is called dynamic (or is said to have memory) if, for some $t_{0}, y\left(t_{0}\right)$ depends on $u(t)$ for $t \neq t_{0}$. A system is called instantaneous (static, memoryless or is said to have zero memory) if, for any $t_{0}, y\left(t_{0}\right)$ does not depend on $u(t)$ for $t \neq t_{0}$.

A system is said to be relaxed (or at rest) at time $t_{0}$ if the output $y_{\left[t_{0}, \infty\right)}$ is solely excited and uniquely determined by the input $u_{\left[t_{0}, \infty\right)}$. We assume that every system is relaxed at time $t=-\infty$.

Definition (State)
The state $x(t)$ of a system at time $t$ is the information at time $t$ that is sufficient to uniquely specify the output $y_{[t, \infty)}$ given the input $u_{[t, \infty)}$.

Definition (Linear State Space Systems)
A system is said to be linear if for every $t_{0}$ and any admissible two input-state-output pairs

$$
\begin{array}{r}
\left\{u_{\left[t_{0}, \infty\right)}^{1}, x^{1}\left(t_{0}\right)\right\} \longrightarrow y_{\left[t_{0}, \infty\right)}^{1} \\
\left\{u_{\left[t_{0}, \infty\right)}^{2}, x^{2}\left(t_{0}\right)\right\} \longrightarrow y_{\left[t_{0}, \infty\right)}^{2}
\end{array}
$$

and real (complex) scalar $\alpha$, the property of additivity
$\left\{u_{\left[t_{0}, \infty\right)}^{1}+u_{\left[t_{0}, \infty\right)}^{2}, x^{1}\left(t_{0}\right)+x^{2}\left(t_{0}\right)\right\} \longrightarrow y_{\left[t_{0}, \infty\right)}^{1}+y_{\left[t_{0}, \infty\right)}^{2}$
and the property of homogeneity

$$
\left\{\alpha u_{\left[t_{0}, \infty\right)}^{1}, \alpha x^{1}\left(t_{0}\right)\right\} \longrightarrow \alpha y_{\left[t_{0}, \infty\right)}^{1}
$$

hold. Otherwise, the system is said to be nonlinear.

Definition (Time Invariance)
A system is said to be time-invariant if for every $t_{0}$ and every admissible input-state-output pair on the interval $\left[t_{0}, \infty\right)$

$$
\left\{u_{\left[t_{0}, \infty\right)}^{1}, x^{1}\left(t_{0}\right)=x^{0}\right\} \longrightarrow y_{\left[t_{0}, \infty\right)}^{1}
$$

and any $T$, the input-state-output pair on the interval $\left[t_{0}+T, \infty\right)$

$$
\left\{u_{\left[t_{0}+T, \infty\right)}^{2}, x^{2}\left(t_{0}+T\right)=x^{0}\right\} \longrightarrow y_{\left[t_{0}+T, \infty\right)}^{2}
$$

where
$u^{2}(t)=u^{1}(t-T), \quad y^{2}(t)=y^{1}(t-T), \quad t \geq t_{0}+T$
is also admissible.

## Definition (Fields)

A field consists of a set, denoted by $\mathcal{F}$, of elements called scalars and two operations called addition " + " and multiplication "." which are defined over $\mathcal{F}$ such that they satisfy the following conditions:

1. To every pair of elements $\alpha, \beta \in \mathcal{F}$, there is an associated unique element $\alpha+\beta$ in $\mathcal{F}$, called the sum of $\alpha$ and $\beta$.
2. There exists an element, denoted by 0 , such that $\alpha+0=\alpha$ for all $\alpha \in \mathcal{F}$.
3. Addition is commutative; $\alpha+\beta=\beta+\alpha$.
4. Addition is associative;
$(\alpha+\beta)+\gamma=\alpha+(\beta+\gamma)$.
5. For each $\alpha \in \mathcal{F}$, there exists an element $\beta$ such that $\alpha+\beta=0$, which is called the additive inverse and denoted by $-\alpha$.
6 . To every pair of elements $\alpha, \beta \in \mathcal{F}$, there is associated a unique element $\alpha \cdot \beta$ in $\mathcal{F}$, called the product of $\alpha$ and $\beta$.
6. There exists an element, denoted by 1 , such that $\alpha \cdot 1=\alpha$ for all $\alpha \in \mathcal{F}$.
7. Multiplication is commutative; $\alpha \cdot \beta=\beta \cdot \alpha$.
8. Multiplication is associative;
$(\alpha \cdot \beta) \cdot \gamma=\alpha \cdot(\beta \cdot \gamma)$.
10 . For each $\alpha \in \mathcal{F}, \alpha \neq 0$, there exists an element $\gamma$ such that $\alpha \cdot \gamma=1$, which is called the multiplicative inverse and denoted by $\alpha^{-1}$.
9. Multiplication is distributive with respect to addition; $(\alpha+\beta) \cdot \gamma=\alpha \cdot \beta+\alpha \cdot \gamma$.

## Definition (Vector Spaces)

A vector space over a field $\mathcal{F}$, denoted by $(\mathcal{X}, \mathcal{F})$, consists of a set, denoted by $\mathcal{X}$, of elements called vectors, a field $\mathcal{F}$, and two operations called vector addition " + " and scalar multiplication ".", which are defined over $\mathcal{X}$ and $\mathcal{F}$ such that they satisfy the following conditions:

1. To every pair of elements $x_{1}, x_{2} \in \mathcal{X}$, there is an associated unique element $x_{1}+x_{2}$ in $\mathcal{X}$, called the sum of $x_{1}$ and $x_{2}$.
2. Addition is commutative; $x_{1}+x_{2}=x_{2}+x_{1}$.
3. Addition is associative;
$\left(x_{1}+x_{2}\right)+x_{3}=x_{1}+\left(x_{2}+x_{3}\right)$.
4. There exists a vector, denoted by 0 , such that $x+0=x$ for all $x \in \mathcal{X}$.
5. For each $x \in \mathcal{X}$, there exists an element $y \in \mathcal{X}$ such that $x+y=0$, which is denoted by $-x$.
6. To every scalar $\alpha \in \mathcal{F}$ and vector $x \in \mathcal{X}$, there is associated a unique element $\alpha \cdot x \in \mathcal{X}$, called the scalar product of $\alpha$ and $x$.
7. Scalar multiplication is associative;
$(\alpha \beta) \cdot x=\alpha \cdot(\beta \cdot x)$.
8. Scalar multiplication is distributive with respect to vector addition; $\alpha \cdot\left(x_{1}+x_{2}\right)=\alpha \cdot x_{1}+\alpha \cdot x_{2}$.
9. Scalar multiplication is distributive with respect to scalar addition; $(\alpha+\beta) \cdot x=\alpha \cdot x+\beta \cdot x$.
10. For $1 \in \mathcal{F}$ and $x \in \mathcal{X}, 1 \cdot x=x$.

## Definition (Subspaces)

Let $(\mathcal{X}, \mathcal{F})$ be a vector space and $\mathcal{Y} \subset \mathcal{X}$. Then, $(\mathcal{Y}, \mathcal{F})$ is said to be a subspace of $(\mathcal{X}, \mathcal{F})$ if under the operations of $(\mathcal{X}, \mathcal{F}), \mathcal{Y}$ forms a vector space over $\mathcal{F}$, that is, for each pair of elements $x_{1}, x_{2} \in \mathcal{Y}$, $x_{1}+x_{2} \in \mathcal{Y}$, and for each scalar $\alpha \in \mathcal{F}$ and vector $x \in \mathcal{Y}, \alpha x \in \mathcal{Y}$.

## Definition (Sum of Subspaces)

Let $(\mathcal{X}, \mathcal{F})$ be a vector space and $\mathcal{Y}, \mathcal{Z} \subset \mathcal{X}$ be subspaces of $\mathcal{X}$. Then, the sum of two subspaces $\mathcal{Y}, \mathcal{Z} \subset \mathcal{X}$ is the set

$$
\mathcal{Y}+\mathcal{Z}=\{x: x=y+z, y \in \mathcal{Y}, z \in \mathcal{Z}\}
$$

## Definition (Intersection of Subspaces)

Let $(\mathcal{X}, \mathcal{F})$ be a vector space and $\mathcal{Y}, \mathcal{Z} \subset \mathcal{X}$ be subspaces of $\mathcal{X}$. Then, the intersection of two subspaces $\mathcal{Y}, \mathcal{Z} \subset \mathcal{X}$ is the set

$$
\mathcal{Y} \cap \mathcal{Z}=\{x: x \in \mathcal{Y}, x \in \mathcal{Z}\} .
$$

## Definition (Direct Sum of Subspaces)

Let $(\mathcal{X}, \mathcal{F})$ be a vector space and $\mathcal{Y}, \mathcal{Z} \subset \mathcal{X}$ be subspaces of $\mathcal{X}$ such that

$$
\begin{gathered}
\mathcal{Y}+\mathcal{Z}=\mathcal{X} \\
\mathcal{Y} \cap \mathcal{Z}=\{0\}
\end{gathered}
$$

In this case, $\mathcal{X}$ is said to be a direct sum of $\mathcal{Y}$ and $\mathcal{Z}$, and we write $\mathcal{X}=\mathcal{Y} \oplus \mathcal{Z}$.

Definition (Linear Combinations)
Let $A$ be a set in a vector space $(\mathcal{X}, \mathcal{F})$. A vector $x \in \mathcal{X}$ is said to be a linear combination of elements in $A$ if there exist a finite set of vectors
$\left\{x_{1}, x_{2}, \ldots, x_{n}\right\}$ in $A$ and a finite set of scalars
$\left\{\alpha_{1}, \alpha_{2}, \ldots, \alpha_{n}\right\}$ such that

$$
x=\alpha_{1} x_{1}+\alpha_{2} x_{2}+\cdots+\alpha_{n} x_{n} .
$$

## Definition (Span)

Let $A$ be a nonempty subset in a vector space $(\mathcal{X}, \mathcal{F})$. The span of $A$, denoted by $\operatorname{span}(A)$, is the set consisting of all linear combinations of elements in $A$. For convenience, $\operatorname{span}(\emptyset)=\{0\}$.

Definition (Linear Independence)
A nonempty set $A$ in a vector space $(\mathcal{X}, \mathcal{F})$ is said to be linearly dependent if there exists a finite set of distinct elements $\left\{x_{1}, x_{2}, \ldots, x_{n}\right\}$ in $A$ and a finite set of scalars $\left\{\alpha_{1}, \alpha_{2}, \ldots, \alpha_{n}\right\}$, not all zero, such that

$$
\alpha_{1} x_{1}+\alpha_{2} x_{2}+\cdots+\alpha_{n} x_{n}=0 .
$$

A set $A$ in a vector space $(\mathcal{X}, \mathcal{F})$ is said to be linearly independent if it is not linearly dependent. In other words, a set $A$ in a vector space $(\mathcal{X}, \mathcal{F})$ is linearly independent if for each nonempty finite subset of distinct elements in $A$, say, $\left\{x_{1}, x_{2}, \ldots, x_{n}\right\}$ in $A$, the only $n$-tuple of scalars satisfying the equation

$$
\alpha_{1} x_{1}+\alpha_{2} x_{2}+\cdots+\alpha_{n} x_{n}=0
$$

is the trivial solution $\alpha_{1}=\alpha_{2}=\cdots=\alpha_{n}=0$.

Definition (Basis)
A set $B$ of linearly independent vectors in a vector space $(\mathcal{X}, \mathcal{F})$ is said to be a basis of $\mathcal{X}$ if every
vector in $\mathcal{X}$ can be expressed as a linear combination of vectors in $B$, that is, $\operatorname{span}(B)=\mathcal{X}$.

## Theorem

Let $(\mathcal{X}, \mathcal{F})$ be a vector space and
$V=\left\{v_{1}, v_{2}, \ldots, v_{n}\right\}$ be a subset of $\mathcal{X}$. Then, $V$ is a basis for $\mathcal{X}$ if and only if each vector $x$ in $\mathcal{X}$ can be uniquely expressed as a linear combination of vectors in $V$, that is, can be expressed in the form

$$
x=\alpha_{1} v_{1}+\alpha_{2} v_{2}+\cdots+\alpha_{n} v_{n}
$$

for unique scalars $\alpha_{1}, \alpha_{2}, \ldots, \alpha_{n}$.

## Theorem (Replacement Theorem)

 Let $(\mathcal{X}, \mathcal{F})$ be a vector space and $V=\left\{v_{1}, v_{2}, \ldots, v_{n}\right\} \subset \mathcal{X}$ be a basis for $\mathcal{X}$. Let $U=\left\{u_{1}, u_{2}, \ldots, u_{m}\right\} \subset \mathcal{X}$ be linearly independent. Then, $m \leq n$.Theorem
Let $(\mathcal{X}, \mathcal{F})$ be a vector space and
$V=\left\{v_{1}, v_{2}, \ldots, v_{n}\right\} \subset \mathcal{X}$ be a basis for $\mathcal{X}$. Let $U=\left\{u_{1}, u_{2}, \ldots, u_{m}\right\} \subset \mathcal{X}$ be another basis for $\mathcal{X}$.
Then, $m=n$.

Definition (Dimension)
A vector space $(\mathcal{X}, \mathcal{F})$ is called finite-dimensional if it has a basis $B$ consisting of a finite number of elements. The unique number of elements in each basis for $\mathcal{X}$ is called the dimension of $\mathcal{X}$ and is denoted by $\operatorname{dim}(\mathcal{X})$. A vector space that is not finite-dimensional is called infinite-dimensional.

Definition (Linear Transformation)
Let $(\mathcal{X}, \mathcal{F})$ and $(\mathcal{Y}, \mathcal{F})$ be vector spaces. A function $L$ which maps $(\mathcal{X}, \mathcal{F})$ into $(\mathcal{Y}, \mathcal{F})$ is called a linear transformation if

$$
\begin{aligned}
& L(\alpha x)=\alpha L x \quad \text { for all } x \in \mathcal{X} \text { and all } \alpha \in \mathcal{F} \\
& L\left(x_{1}+x_{2}\right)=L x_{1}+L x_{2} \quad \text { for all } x_{1}, x_{2} \in \mathcal{X}
\end{aligned}
$$

Definition (Range space)
The range space $\mathcal{R}(L)$ of a linear transformation $L$ from $(\mathcal{X}, \mathcal{F})$ into $(\mathcal{Y}, \mathcal{F})$ is a subset of $(\mathcal{Y}, \mathcal{F})$ defined by

$$
\mathcal{R}(L)=\{y \in \mathcal{Y}: y=L x \text { for some } x \in \mathcal{X}\}
$$

Definition (Null space)
The null space $\mathcal{N}(L)$ of a linear transformation $L$ from $(\mathcal{X}, \mathcal{F})$ into $(\mathcal{Y}, \mathcal{F})$ is a subset of $(\mathcal{X}, \mathcal{F})$ defined by

$$
\mathcal{N}(L)=\{x \in \mathcal{X}: L x=0\}
$$

Theorem
Consider a linear transformation $L$ from $(\mathcal{X}, \mathcal{F})$ into $(\mathcal{Y}, \mathcal{F})$. The following statements are equivalent.

1. The mapping $L$ is one-to-one.
2. Null space of $L$ is trivial, that is, $\mathcal{N}(L)=\{0\}$.
3. $L$ maps linearly independent vectors in $\mathcal{X}$ into a linearly independent vectors in $\mathcal{Y}$.

Definition (Nullity and rank)
Let $L$ be a linear transformation from $(\mathcal{X}, \mathcal{F})$ into $(\mathcal{Y}, \mathcal{F})$. If $\mathcal{N}(L)$ and $\mathcal{R}(L)$ are finite dimensional, we define the nullity of $L$, denoted by nullity $(L)$, and rank of $L$, denoted by $\operatorname{rank}(L)$, to be the dimensions of $\mathcal{N}(L)$ and $\mathcal{R}(L)$, respectively.

Theorem (Dimension Theorem)
Let $L$ be a linear transformation from $(\mathcal{X}, \mathcal{F})$ into $(\mathcal{Y}, \mathcal{F})$. If $\mathcal{X}$ is finite dimensional, then $\operatorname{nullity}(L)+\operatorname{rank}(L)=\operatorname{dim}(\mathcal{X})$.

Definition (Similarity transformation)
Two square matrices $A$ and $B$ are said to be similar if a nonsingular matrix $T$ exists such that $A=T^{-1} B T$. The matrix $T$ is called a similarity transformation.

Definition (Invariant subspace)
Let $L$ be a linear transformation of $(\mathcal{X}, \mathcal{F})$ into itself.
A subspace $\mathcal{Y}$ of $\mathcal{X}$ is said to be an invariant subspace of $\mathcal{X}$ under $L$, or an $L$-invariant subspace of $\mathcal{X}$, if $L(\mathcal{Y}) \subset \mathcal{Y}$, which implies that $L y \in \mathcal{Y}$ for all $y \in \mathcal{Y}$.

Definition (Eigenvalue and eigenvector)
Let $A$ be a linear transformation from $\left(\mathrm{C}^{n}, \mathrm{C}\right)$ into ( $\mathrm{C}^{n}, \mathrm{C}$ ). A scalar $\lambda$ in C is called an eigenvalue of $A$ if there exists a nonzero vector $x$ in $\mathrm{C}^{n}$ such that $A x=\lambda x$. Any nonzero vector $x$ satisfying $A x=\lambda x$ is called an eigenvector of $A$ associated with the eigenvalue $\lambda$.

## Theorem

Let $\lambda_{1}, \lambda_{2}, \ldots, \lambda_{d}$ be distinct eigenvalues of $A$ and let $v_{i}$ be an eigenvector of $A$ associated with $\lambda_{i}$. Then, the set $\left\{v_{1}, v_{2}, \ldots, v_{d}\right\}$ is linearly independent.

## Theorem

Let $\lambda_{1}, \lambda_{2}, \ldots, \lambda_{d}$ be distinct eigenvalues of $A$ and let $\left\{v_{i 1}, v_{i 2}, \ldots, v_{i g_{i}}\right\}$ be a set of linearly independent eigenvectors of $A$ associated with $\lambda_{i}$. Then, the set

$$
\left\{v_{11}, \ldots, v_{1 g_{1}}, v_{21}, \ldots, v_{2 g_{2}}, \ldots, v_{d 1}, \ldots, v_{d g_{d}}\right\}
$$

is linearly independent.

Theorem
$A$ is diagonalizable if and only if it has $n$ linearly independent eigenvectors, that is, $g_{i}=m_{i}$, $i=1,2, \ldots, d$.

## Definition (Generalized eigenspace)

The generalized eigenspace of $A$ corresponding to the eigenvalue $\lambda$ is the subset defined by

$$
\left\{x:(\lambda I-A)^{k} x=0 \text { for some } k \geq 1\right\} .
$$

Theorem
The dimension of the generalized eigenspace of $A$ corresponding to the eigenvalue $\lambda_{i}$ is equal to $m_{i}$.

Definition (Generalized eigenvector)
A vector $v$ is called a generalized eigenvector of grade $k(k \geq 1)$ of $A$ associated with eigenvalue $\lambda$ if

$$
\begin{aligned}
(A-\lambda I)^{k} v & =0 \\
(A-\lambda I)^{k-1} v & \neq 0
\end{aligned}
$$

Theorem
The union over all the eigenvalues of the set of generalized eigenvectors associated with each different eigenvalue is linearly independent.

$$
A=\left[\begin{array}{ccccc}
1 & 0 & -1 & 1 & 0 \\
-4 & 1 & -3 & 2 & 1 \\
-2 & -1 & 0 & 1 & 1 \\
-3 & -1 & -3 & 4 & 1 \\
-8 & -2 & -7 & 5 & 4
\end{array}\right], \quad \operatorname{det}(\lambda I-A)=(\lambda-2)^{5}
$$

$\operatorname{dim}(\mathcal{N}(A-2 I))=2, \operatorname{dim}\left(\mathcal{N}(A-2 I)^{2}\right)=4$, $\operatorname{dim}\left(\mathcal{N}(A-2 I)^{3}\right)=5$

$$
u_{11}=\left[\begin{array}{c}
0 \\
-1 \\
1 \\
1 \\
0
\end{array}\right], u_{12}=\left[\begin{array}{l}
0 \\
1 \\
0 \\
0 \\
1
\end{array}\right], u_{21}=\left[\begin{array}{l}
0 \\
1 \\
0 \\
0 \\
0
\end{array}\right], u_{22}=\left[\begin{array}{c}
-1 \\
0 \\
1 \\
0 \\
0
\end{array}\right]
$$

$$
u_{31}=\left[\begin{array}{l}
0 \\
0 \\
0 \\
1 \\
0
\end{array}\right]=v_{31}, v_{21}=(A-2 I) v_{31}=\left[\begin{array}{l}
1 \\
2 \\
1 \\
2 \\
5
\end{array}\right]
$$

Choose $v_{22}$ s.t. $\left\{u_{11}, u_{12}, v_{21}, v_{22}\right\}$ is a basis of $\mathcal{N}(A-2 I)^{2}$ :

$$
v_{22}=\left[\begin{array}{c}
-1 \\
0 \\
1 \\
0 \\
0
\end{array}\right]
$$

$$
v_{11}=(A-2 I) v_{21}, \quad v_{12}=(A-2 I) v_{22}
$$

$$
A=\left[\begin{array}{ccc}
5 & 19 & 9 \\
0 & 0 & 1 \\
-1 & -4 & -2
\end{array}\right], \operatorname{det}(\lambda I-A)=(\lambda-1)^{3}
$$

$\operatorname{dim}(\mathcal{N}(A-I))=1$

$$
v_{1}=[7,-1,-1]^{T}
$$

$$
(A-I) v_{2}=v_{1}
$$

$$
v_{2}=[-3,1,0]^{T}
$$

$$
(A-I) v_{3}=v_{2}
$$

$$
v_{3}=[-3,0,1]^{T}
$$

$$
A=\left[\begin{array}{ccc}
-2 & -9 & -12 \\
1 & 4 & 4 \\
0 & 0 & 1
\end{array}\right], \operatorname{det}(\lambda I-A)=(\lambda-1)^{3}
$$

$\operatorname{dim}(\mathcal{N}(A-I))=2$

$$
\begin{array}{r}
v_{11}=[3,-1,0]^{T}, v_{12}=[4,0,-1]^{T} \\
(A-I) v_{2}=?
\end{array}
$$

## Theorem (Cayley-Hamilton THM)

Let
$\gamma(\lambda)=\operatorname{det}(\lambda I-A)=\lambda^{n}+c_{n-1} \lambda^{n-1}+\cdots+c_{1} \lambda+c_{0}$ be the characteristic polynomial of $A$. Then,

$$
\gamma(A)=A^{n}+c_{n-1} A^{n-1}+\cdots+c_{1} A+c_{0} I=0
$$

Definition (Minimal Polynomial)
The minimal polynomial of a matrix $A$ is the monic polynomial $\mu(\lambda)$ of the least degree such that $\mu(A)=0$.

Theorem
The minimal polynomial $\mu(\lambda)$ of $J$ is given as

$$
\mu(\lambda)=\prod_{i=1}^{d}\left(\lambda-\lambda_{i}\right)^{\eta_{i}}
$$

where $\eta_{i}$ is the index of $\lambda_{i}$ in $J$.

Definition (Function of a matrix)
Let $f(\lambda)$ be a function that is defined on the spectrum of $A$. If $p(\lambda)$ is a polynomial that has the same values as $f(\lambda)$ on the spectrum of $A$, then $f(A)$ is defined as $f(A)=p(A)$.

## Lemma

Given distinct numbers $\lambda_{1}, \lambda_{2}, \ldots, \lambda_{d}$, positive integers $\eta_{1}, \eta_{2}, \ldots, \eta_{d}$ with $\eta=\sum_{i=1}^{d} \eta_{i}$, and a set of numbers

$$
f_{i, 0}, f_{i, 1}, \ldots, f_{i, \eta_{i}-1}, \quad i=1,2, \ldots, d
$$

there exists a polynomial $p(\lambda)$ of degree less than $\eta$ such that
$p\left(\lambda_{i}\right)=f_{i, 0}, p^{(1)}\left(\lambda_{i}\right)=f_{i, 1}, \ldots, p^{\left(\eta_{i}-1\right)}\left(\lambda_{i}\right)=f_{i, \eta_{i}-1}$,
Furthermore, such a polynomial $p(\lambda)$ is unique.

$$
a_{1}=\frac{1}{4}\left(e^{-t}-e^{-5 t}\right) \text { and } a_{0}=\frac{5}{4} e^{-t}-\frac{1}{4} e^{-5 t}
$$

$$
\begin{aligned}
e^{A t} & =\left(\frac{5}{4} e^{-t}-\frac{1}{4} e^{-5 t}\right) I+\frac{1}{4}\left(e^{-t}-e^{-5 t}\right) A \\
& =\left[\begin{array}{cc}
2 e^{-t}-e^{-5 t} & e^{-t}-e^{-5 t} \\
\frac{1}{2}\left(e^{-t}-e^{-5 t}\right) & \frac{3}{2} e^{-t}-\frac{1}{2} e^{-5 t}
\end{array}\right]
\end{aligned}
$$

$$
\begin{aligned}
f(1)=p(1) & e^{t}=a_{0}+a_{1}+a_{2} \\
f^{(1)}(1)=p^{(1)}(1) & t e^{t}=a_{1}+2 a_{2} \\
f(2)=p(2) & e^{2 t}=a_{0}+2 a_{1}+4 a_{2}
\end{aligned}
$$

$$
e^{A t}=\left(-2 t e^{t}+e^{2 t}\right) I+\left(3 t e^{t}+2 e^{t}-2 e^{2 t}\right) A+\left(e^{2 t}-e^{t}-t e^{t}\right) A^{2}
$$

$$
\begin{gathered}
e^{t}=a_{0}+a_{1} \\
e^{2 t}=a_{0}+2 a_{1} \\
e^{A t}=\left(2 e^{t}-e^{2 t}\right) I+\left(e^{2 t}-e^{t}\right) A
\end{gathered}
$$

cf.

$$
e^{A t}=\left(-2 t e^{t}+e^{2 t}\right) I+\left(3 t e^{t}+2 e^{t}-2 e^{2 t}\right) A+\left(e^{2 t}-e^{t}-t e^{t}\right) A^{2}
$$

$$
\begin{gathered}
A=\left[\begin{array}{lll}
1 & 1 & 2 \\
0 & 1 & 3 \\
0 & 0 & 2
\end{array}\right]=P J P^{-1}=P\left[\begin{array}{lll}
1 & 1 & 0 \\
0 & 1 & 0 \\
0 & 0 & 2
\end{array}\right] P^{-1} \\
P=\left[\begin{array}{lll}
1 & 0 & 5 \\
0 & 1 & 3 \\
0 & 0 & 1
\end{array}\right] \\
f(\lambda)=e^{\lambda t}, f^{(1)}(\lambda)=t e^{\lambda t} \\
f(J)=e^{J t}=\left[\begin{array}{ccc}
e^{t} & t e^{t} & 0 \\
0 & e^{t} & 0 \\
0 & 0 & e^{2 t}
\end{array}\right] \\
f(A)=e^{A t}=P e^{J t} P^{-1}=P\left[\begin{array}{ccc}
e^{t} & t e^{t} & 0 \\
0 & e^{t} & 0 \\
0 & 0 & e^{2 t}
\end{array}\right] P^{-1}
\end{gathered}
$$

## Definition (Inner Product)

Let $(\mathcal{X}, \mathcal{F})$ be a vector space. The inner product of two vectors $x$ and $y$ in $\mathcal{X}$ denoted by $<x, y\rangle$ takes the value in $\mathcal{F}$ and satisfies the following properties for all $a, b \in \mathcal{F}$ :

$$
\begin{aligned}
& \text { 1. }<x, y>=<y, x> \\
& \text { 2. }<x, a y+b z>=a<x, y>+b<x, z> \\
& \text { 3. }<x, x>\geq 0 \text { for all } x \text {, and }<x, x>=0 \text { if and } \\
& \text { only if } x=0 \text {. }
\end{aligned}
$$

Definition (Orthogonality)
Two vectors $x$ and $y$ in an inner product space $\mathcal{X}$ are said to be orthogonal if their inner product is zero, that is, $\langle x, y\rangle=0$. If $x$ and $y$ are orthogonal, this is denoted by $x \perp y$.
Two subsets $A$ and $B$ in an inner product space $\mathcal{X}$ are said to be orthogonal if $x \perp y$ for all $x$ in $A$ and $y$ in $B$. If $A$ and $B$ are orthogonal, this is denoted by $A \perp B$.

## Definition (Orthogonal Complement)

Let $\mathcal{X}$ be an inner product space and $\mathcal{Y}$ be any subspace of $\mathcal{X}$. The set

$$
\mathcal{Y}^{\perp}=\{x:<x, y>=0 \text { for all } y \in \mathcal{Y}\}
$$

is called the orthogonal complement of $\mathcal{Y}$.

Definition (Norm)
A real-valued function $\|x\|$ defined on a vector space $(\mathcal{X}, \mathcal{F})$, where $x \in \mathcal{X}$, is called a norm if for all $x, y \in \mathcal{X}$ and $\alpha \in \mathcal{F}$
(i) $\|x\| \geq 0$ for all $x \in \mathcal{X}$, and $\|x\|=0$ if and only if $x=0$
(ii) $\quad\|\alpha x\|=|\alpha|\|x\|$ for all $x \in \mathcal{X}$ and $\alpha \in \mathcal{F}$
(iii) $\|x+y\| \leq\|x\|+\|y\|$ for all $x, y \in \mathcal{X}$.

## Theorem (Schwarz Inequality)

If, on an inner product space, $\|x\|$ is defined by $\|x\|=\left\langle x, x>^{1 / 2}\right.$, then

$$
|<x, y>| \leq\|x\|\|y\| .
$$

Definition (Orthonormal Basis)
Let $\mathcal{X}$ be an inner product space. A subset $B$ is an orthonormal basis if it is a basis that is orthonormal.

## Theorem (Gram-Schmidt Orthonormalization Process)

Let $\mathcal{X}$ be an inner product space, and $W=\left\{w_{1}, w_{2}, \ldots, w_{n}\right\}$ be a linearly independent subset of $\mathcal{X}$. Then, there exists an orthonormal set $V=\left\{v_{1}, v_{2}, \ldots, v_{n}\right\}$ such that $v_{k}=\sum_{i=1}^{k} a_{i k} w_{i}$, $k=1,2, \ldots, n$.

Let $w_{1}=[1,1,0]^{T}, w_{2}=[2,0,1]^{T}, w_{3}=[2,2,1]^{T}$.
Then,

$$
\begin{gathered}
v_{1}=\frac{1}{\left\|w_{1}\right\|} w_{1}=\left[\frac{1}{\sqrt{2}}, \frac{1}{\sqrt{2}}, 0\right]^{T} \\
u_{2}=w_{2}-<v_{1}, w_{2}>v_{1}=[2,0,1]^{T}-[1,1,0]^{T} \\
=[1,-1,1]^{T} \\
v_{2}=\left[\frac{1}{\sqrt{3}},-\frac{1}{\sqrt{3}}, \frac{1}{\sqrt{3}}\right]^{T} \\
u_{3}=w_{3}-<v_{1}, w_{3}>v_{1}-<v_{2}, w_{3}>v_{2} \\
v_{3}=\frac{1}{\left\|u_{3}\right\|} u_{3}=\left[-\frac{1}{\sqrt{6}}, \frac{1}{\sqrt{6}}, \frac{2}{\sqrt{6}}\right]^{T}
\end{gathered}
$$

## Corollary

Let $\mathcal{X}$ be an inner product space and $\mathcal{Y}$ be a finite dimensional subspace of $\mathcal{X}$. For each $x \in \mathcal{X}$, there exist unique vectors $y \in \mathcal{Y}$ and $z \in \mathcal{Y}^{\perp}$ such that $x=y+z$. Furthermore, if $\left\{v_{1}, v_{2}, \ldots, v_{k}\right\}$ is an orthonormal basis of $\mathcal{Y}$, then

$$
y=\sum_{i=1}^{k}<v_{i}, x>v_{i}
$$

The vector $y$ is the unique vector such that $\|x-y\|<\|x-u\|$ for any $u \in \mathcal{Y}$ such that $u \neq y$.

## Theorem (Properties of induced norm)

For $m \times n$ matrices $A$ and $B$ mapping normed space $\mathcal{R}^{n}$ into normed space $\mathcal{R}^{m}$, and $n \times l$ matrix $C$ mapping normed space $\mathcal{R}^{l}$ into normed space $\mathcal{R}^{n}$, the following hold:

$$
\begin{aligned}
& \text { 1. }\|A\| \geq 0 \text {, and }\|A\|=0 \text { if and only if } A=0 \text {, } \\
& \text { 2. }\|\alpha A\|=|\alpha|\|A\| \text {, } \\
& \text { 3. }\|A+B\| \leq\|A\|+\|B\| \text {, } \\
& \text { 4. }\|A C\| \leq\|A\|\|C\| \text {. }
\end{aligned}
$$

## Theorem

Let $A$ be an $m \times n$ matrix mapping $\mathcal{R}^{n}$ into $\mathcal{R}^{m}$ defined by $y=A x$ with the adjoint transformation given by $x=A^{T} y$.

1. $\mathcal{N}\left(A^{T}\right)$ is an orthogonal complement of $\mathcal{R}(A)$, that is, $\mathcal{N}\left(A^{T}\right)=\mathcal{R}(A)^{\perp}$.
2. $\mathcal{R}(A)$ is an orthogonal complement of $\mathcal{N}\left(A^{T}\right)$, that is, $\mathcal{R}(A)=\mathcal{N}\left(A^{T}\right)^{\perp}$.

## Theorem

For a Hermitian matrix $H$,

1. all the eigenvalues are real,
2. there are $n$ linearly independent eigenvectors,
3. the eigenvectors corresponding to different eigenvalues are orthogonal,
4. $\exists$ a unitary matrix $Q$ such that

$$
H=Q \Lambda Q^{-1}=Q \Lambda \bar{Q}^{T}
$$

where $\Lambda$ is a diagonal matrix.

## Definition (Positive definiteness)

A quadratic form $x^{T} M x$ is said to be

1. positive definite if $x^{T} M x>0$ for all $x \neq 0$,
2. positive semidefinite if $x^{T} M x \geq 0$ for all $x$,
3. negative definite if $x^{T} M x<0$ for all $x \neq 0$,
4. negative semidefinite if $x^{T} M x \leq 0$ for all $x$.

A symmetric matrix $M$ is said to be positive definite (positive semidefinite, negative definite, negative semidefinite, respectively) if the quadratic form $x^{T} M x$ is so.

## Theorem

A symmetric matrix $M$ is positive definite (positive semidefinite) if and only if all the eigenvalues of $M$ are positive (nonnegative).
Theorem
If a symmetric matrix $M$ is positive definite (positive semidefinite), then $\operatorname{det}(M)>0(\operatorname{det}(M) \geq 0)$.

## Theorem

A symmetric matrix $M$ is positive definite (positive semidefinite) if and only if all the leading principal minors (all the principal minors) of $M$ are positive (nonnegative).


$$
\begin{gathered}
(M+m) \ddot{x}+m l\left(\ddot{\theta} \cos \theta-\dot{\theta}^{2} \sin \theta\right)=u \\
m(\ddot{x} \cos \theta+l \ddot{\theta}-g \sin \theta)=0 \\
\Downarrow \\
(M+m) \ddot{x}+m l \ddot{\theta}=u \\
\ddot{x}+l \ddot{\theta}-g \theta=0
\end{gathered}
$$

$(M+m) \ddot{x}+m l\left(\ddot{\theta} \cos \theta-\dot{\theta}^{2} \sin \theta\right)=u$ $m(\ddot{x} \cos \theta+l \ddot{\theta}-g \sin \theta)=0$
$\Downarrow$

$$
\begin{gathered}
{\left[\begin{array}{cc}
M+m & m l \\
1 & l
\end{array}\right]\left[\begin{array}{l}
\ddot{x} \\
\ddot{\theta}
\end{array}\right]=\left[\begin{array}{c}
u \\
g \theta
\end{array}\right]} \\
\Downarrow \\
{\left[\begin{array}{c}
\ddot{x} \\
\ddot{\theta}
\end{array}\right]=\frac{1}{M l}\left[\begin{array}{cc}
l & -m l \\
-1 & M+m
\end{array}\right]\left[\begin{array}{c}
u \\
g \theta
\end{array}\right]}
\end{gathered}
$$

## Theorem

If $A$ is an $n \times n$ matrix function whose entries are continuous functions of time on the interval $I=\left[t_{l}, t_{u}\right]$, then there exists the unique solution to the initial value problem

$$
\dot{x}(t)=A(t) x(t), \quad x\left(t_{0}\right)=x^{0}, \quad t_{0} \in I=\left[t_{l}, t_{u}\right] .
$$

$$
\frac{d}{d t} x(t)=\left[\begin{array}{ll}
1 & 0 \\
t & 1
\end{array}\right] x(t)=A(t) x(t)
$$

Show

$$
\Phi\left(t, t_{0}\right)=\left[\begin{array}{cc}
e^{\left(t-t_{0}\right)} & 0 \\
\frac{1}{2}\left(t^{2}-t_{0}^{2}\right) e^{\left(t-t_{0}\right)} & e^{\left(t-t_{0}\right)}
\end{array}\right]
$$

is the state transition matrix. Answer: $\Phi\left(t_{0}, t_{0}\right)=I$ and

$$
\begin{aligned}
\frac{d}{d t} \Phi\left(t, t_{0}\right) & =\left[\begin{array}{cc}
e^{\left(t-t_{0}\right)} & 0 \\
\frac{1}{2}\left(t^{2}-t_{0}^{2}\right) e^{\left(t-t_{0}\right)}+t e^{\left(t-t_{0}\right)} & e^{\left(t-t_{0}\right)}
\end{array}\right] \\
& =\left[\begin{array}{ll}
1 & 0 \\
t & 1
\end{array}\right]\left[\begin{array}{cc}
e^{\left(t-t_{0}\right)} & 0 \\
\frac{1}{2}\left(t^{2}-t_{0}^{2}\right) e^{\left(t-t_{0}\right)} & e^{\left(t-t_{0}\right)}
\end{array}\right] \\
& =\left[\begin{array}{ll}
1 & 0 \\
t & 1
\end{array}\right] \Phi\left(t, t_{0}\right)=A(t) \Phi\left(t, t_{0}\right)
\end{aligned}
$$

## Definition (Stability)

For $\dot{x}=f(t, x)$, the equilibrium point $x_{e}$ (i.e., $\left.f\left(t, x_{e}\right)=0,{ }^{\forall} t\right)$ is

1. stable i.s.L. (in the sense of Lyapunov) if for each $t_{0}$ and each $\epsilon>0$, there exists $\delta\left(\epsilon, t_{0}\right)>0$ s.t. if $\left\|x\left(t_{0}\right)-x_{e}\right\|<\delta$ then

$$
\left\|x(t)-x_{e}\right\|<\epsilon \quad \text { for all } t \geq t_{0}
$$

2. uniformly stable i.s.L. if, ${ }^{\forall} \epsilon>0, \exists \delta=\delta(\epsilon)>0$ s.t. if $\left\|x\left(t_{0}\right)-x_{e}\right\|<\delta$ then

$$
\left\|x(t)-x_{e}\right\|<\epsilon \quad \text { for all } t \geq t_{0}
$$

3. unstable if it is not stable
4. asymptotically stable if it is stable i.s.L. and for each $t_{0}$, there is a positive constant $c$ such that if $\left\|x\left(t_{0}\right)-x_{e}\right\|<c$, then $x(t) \rightarrow x_{e}$ as $t \rightarrow \infty$.
5. globally asymptotically stable if it is stable i.s.L. and for each $t_{0}$ and each $x\left(t_{0}\right), x(t) \rightarrow x_{e}$ as $t \rightarrow \infty$.

$$
\begin{gathered}
\dot{x}=\left[\begin{array}{cc}
0 & \omega \\
-\omega & 0
\end{array}\right] x(t) \\
\Phi\left(t, t_{0}\right)=\left[\begin{array}{cc}
\cos \omega\left(t-t_{0}\right) & \sin \omega\left(t-t_{0}\right) \\
-\sin \omega\left(t-t_{0}\right) & \cos \omega\left(t-t_{0}\right)
\end{array}\right]
\end{gathered}
$$

$$
\begin{aligned}
& x_{1}(t)=\cos \omega\left(t-t_{0}\right) x_{1}\left(t_{0}\right)+\sin \omega\left(t-t_{0}\right) x_{2}\left(t_{0}\right) \\
& x_{2}(t)=-\sin \omega\left(t-t_{0}\right) x_{1}\left(t_{0}\right)+\cos \omega\left(t-t_{0}\right) x_{2}\left(t_{0}\right)
\end{aligned}
$$

## Theorem

For linear continuous time system, the equilibrium point at the origin is

1. stable if and only if for each $t_{0}$, there exists a constant $\kappa\left(t_{0}\right)$ such that

$$
\left\|\Phi\left(t, t_{0}\right)\right\| \leq \kappa\left(t_{0}\right) \quad \text { for all } t \geq t_{0}
$$

2. asymptotically stable if and only if for each $t_{0}$, $\left\|\Phi\left(t, t_{0}\right)\right\| \rightarrow 0$ as $t \rightarrow \infty$.

$$
\begin{gathered}
\dot{x}=\left[\begin{array}{cc}
-1 & e^{2 t} \\
0 & -1
\end{array}\right] x \\
\Phi\left(t, t_{0}\right)=\left[\begin{array}{cc}
e^{-\left(t-t_{0}\right)} & \frac{1}{2} e^{-\left(t-t_{0}\right)}\left(e^{2 t}-e^{2 t_{0}}\right) \\
0 & e^{-\left(t-t_{0}\right)}
\end{array}\right]
\end{gathered}
$$

## Theorem

For the linear time invariant system with system matrix $A$, the equilibrium point at the origin is

1. asymptotically stable if and only if all the eigenvalues of $A$ have negative real parts,
2. stable if and only if all the eigenvalues of $A$ have nonpositive real parts, and those eigenvalues with zero real parts are distinct roots of the minimal polynomial of $A$ (or, equivalently, have indices equal to 1 ),
3. unstable if there exist eigenvalues with positive real parts or eigenvalues with zero real parts which are not distinct roots of the minimal polynomial of $A$ (or equivalently, have indices greater than 1).

## Definition (BIBO stability)

The input-output system is said to be bounded-input-bounded-output (BIBO) stable if for any bounded input $u(t),\|u(t)\| \leq M$ for all $t$, there exists a finite constant $N(M)$ such that $\|y(t)\| \leq N$ for all $t$.

Theorem
For a linear time-invariant system, the zero-state response is BIBO stable if and only if all the poles of the transfer function are located in the open left-half complex plane.
Theorem
For a linear time-invariant system, the zero-state response is BIBO stable if for all the eigenvalues $\lambda_{i}$ of the system matrix $A, \operatorname{Re} \lambda_{i}<0$.

## Theorem

The matrix $A$ is Hurwitz, or equivalently, the zero state of $\dot{x}=A x$ is asymptotically stable if and only if for any given symmetric positive definite matrix $Q$, the matrix equation

$$
A^{T} P+P A=-Q
$$

has a unique symmetric positive definite solution $P$.

$$
\begin{gathered}
A=\left[\begin{array}{cc}
-1 & -2 \\
1 & -4
\end{array}\right], \quad Q=I \\
P=\left[\begin{array}{cc}
23 / 60 & -7 / 60 \\
-7 / 60 & 11 / 60
\end{array}\right]
\end{gathered}
$$

Definition (Controllability)
A linear system (or the pair $(A(t), B(t))$ ) is said to be controllable on the interval $\left[t_{0}, t_{1}\right]$ if for any $x^{0}$ in the state space $\mathcal{S}$ and any $x^{1}$ in $\mathcal{S}$, there exists an input $u_{\left[t_{0}, t_{1}\right]}$ which transfers the state $x\left(t_{0}\right)=x^{0}$ to the state $x\left(t_{1}\right)=x^{1}$ at time $t_{1}$.

## Theorem

A linear system (or the pair $(A(t), B(t))$ ) is controllable on the interval $\left[t_{0}, t_{1}\right]$ if and only if the controllability Gramian

$$
G_{c}\left(t_{0}, t_{1}\right)=\int_{t_{0}}^{t_{1}} \Phi\left(t_{0}, t\right) B(t) B^{T}(t) \Phi^{T}\left(t_{0}, t\right) d t
$$

is nonsingular.

## Theorem

A linear system is controllable on the interval $\left[t_{0}, t_{1}\right]$ if there exists $\tau \in\left[t_{0}, t_{1}\right]$ such that

$$
\operatorname{rank}\left(\left[\begin{array}{llll}
M_{0}(\tau) & M_{1}(\tau) & \cdots & M_{n-1}(\tau)
\end{array}\right]\right)=n
$$

where $M_{0}(\tau)=B(\tau)$ and

$$
M_{j}(\tau)=-A(\tau) M_{j-1}(\tau)+\frac{d}{d \tau} M_{j-1}(\tau)
$$

where $j=1,2, \ldots, n-1$.

## Theorem

A linear time-invariant system (or the pair $(A, B)$ ) is controllable on the interval $\left[t_{0}, t_{1}\right]$ if and only if the controllability matrix

$$
\mathcal{P}=\left[\begin{array}{llll}
B & A B & \cdots & A^{n-1} B
\end{array}\right]
$$

has rank $n$.

## Lemma

The range space and null space of

$$
G_{c}\left(t_{0}, t_{1}\right)=\int_{t_{0}}^{t_{1}} e^{A\left(t_{0}-t\right)} B B^{T} e^{A^{T}\left(t_{0}-t\right)} d t
$$

coincide with the range space and null space of $\mathcal{P} \mathcal{P}^{T}$.

## Definition (Observability)

A linear system (or the pair $(C(t), A(t))$ ) is said to be observable on the interval $\left[t_{0}, t_{1}\right]$ if for any initial state $x\left(t_{0}\right)$ in the state space $\mathcal{S}$, the knowledge of the input $u_{\left[t_{0}, t_{1}\right]}$ and the output $y_{\left[t_{0}, t_{1}\right]}$ is sufficient to uniquely solve for $x\left(t_{0}\right)$.

## Theorem

A linear system is observable on the interval $\left[t_{0}, t_{1}\right]$ if there exists $\tau \in\left[t_{0}, t_{1}\right]$ such that

$$
\operatorname{rank}\left(\left[\begin{array}{c}
N_{0}(\tau) \\
N_{1}(\tau) \\
\vdots \\
N_{n-1}(\tau)
\end{array}\right]\right)=n
$$

where $N_{0}(\tau)=C(\tau)$ and

$$
N_{j}(\tau)=N_{j-1}(\tau) A(\tau)+\frac{d}{d \tau} N_{j-1}(\tau)
$$

## Theorem

For a single-input LTI system

$$
\dot{x}=A x+b u
$$

with $\operatorname{det}(s I-A)=s^{n}+a_{n-1} s^{n-1}+\cdots+a_{1} s+a_{0}$, there exists a change of coordinates $x_{c}=T x$ s.t.

$$
\dot{x}_{c}=\left[\begin{array}{cccc}
0 & 1 & \cdots & 0 \\
0 & 0 & \cdots & 0 \\
\vdots & \vdots & \ddots & \vdots \\
0 & 0 & \cdots & 1 \\
-a_{0} & -a_{1} & \cdots & -a_{n-1}
\end{array}\right] x_{c}+\left[\begin{array}{c}
0 \\
0 \\
\vdots \\
0 \\
1
\end{array}\right] u
$$

if and only if the pair $(A, b)$ is controllable.

## Theorem

For a single-output LTI system

$$
\dot{x}=A x, \quad y=c x
$$

with $\operatorname{det}(s I-A)=s^{n}+a_{n-1} s^{n-1}+\cdots+a_{1} s+a_{0}$, there exists a change of coordinates $x_{o}=T x$ s.t.

$$
\begin{aligned}
\dot{x}_{o} & =\left[\begin{array}{ccccc}
0 & 0 & \cdots & 0 & -a_{0} \\
1 & 0 & \cdots & 0 & -a_{1} \\
\vdots & \vdots & \ddots & \vdots & \vdots \\
0 & 0 & \cdots & 0 & -a_{n-2} \\
0 & 0 & \cdots & 1 & -a_{n-1}
\end{array}\right] x_{o} \\
y & =\left[\begin{array}{lllll}
0 & 0 & \cdots & 0 & 1
\end{array}\right] x_{o}
\end{aligned}
$$

if and only if the pair $(c, A)$ is observable.

Definition (Controllable subspace)
The set of the initial states that can be transferred to the zero state in finite time is called the controllable subspace and is denoted by $\mathcal{C}$.

Definition (Stabilizability)
The linear time invariant system is said to be stabilizable if its unstable subspace is contained in its controllable subspace, that is, any vector $x$ in the unstable subspace is also in the controllable subspace.

Definition (Unobservable subspace)
The set of the initial states that produce zero-input responses $\bar{y}(t)$ which are identically zero on any finite time interval is called the unobservable subspace and is denoted by $\mathcal{O}$.

## Definition (Detectability)

The linear time invariant system is said to be detectable if its unobservable subspace is contained in its stable subspace, that is, any vector $x$ in the unobservable subspace is also in the stable subspace.

Theorem
The pair $(A, B)$ is controllable if and only if by the state feedback $u(t)=K x(t)+r(t)$, the eigenvalues of $(A+B K)$ can be arbitrarily assigned provided that complex conjugate eigenvalues appear in pair.

Theorem
The pair $(C, A)$ is observable if and only if the eigenvalues of $A-L C$ can be arbitrarily assigned by a proper choice of the matrix $L$ provided that complex conjugate eigenvalues appear in pair.

## Definition (Realization)

A realization of a transfer matrix $G(s)$ is any state space model $(A, B, C, D)$ such that
$G(s)=C(s I-A)^{-1} B+D$. If such a state space model exists, then $G(s)$ is said to be realizable.

Definition (Minimal realization)
A realization $(A, B, C, D)$ of a transfer matrix $G(s)$ is called a minimal realization if there exists no other realization with state space of smaller dimension.

