## Chapter 2

## Linear Dynamical Systems

## Topics :

1. Solution of Uncontrolled System
2. Solution of Controlled System
3. Time-Varying Systems
4. Relationship between State Space and Classical Forms

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A control system (with outputs) $\Sigma=\left(\mathbb{R}^{m}, \mathbb{R}^{\ell}, \mathcal{U}, F, \mathbb{R}^{n}, h\right)$ is linear if the dynamics $F$ is linear in $(x, u)$, and the measurement function $h$ is linear, for each fixed $t \in \mathbb{R}$. Such a control system is described by (state equation and observation equation)

$$
\dot{x}=A(t) x+B(t) u(t) \quad \text { and } \quad y=C(t) x
$$

where $A(t) \in \mathbb{R}^{m \times m}, B(t) \in \mathbb{R}^{m \times \ell}$, and $C(t) \in \mathbb{R}^{n \times m}$, each of whose entries is a (continuous) function of time. The system is called time-invariant if the structure is independent of time. A system that is not necessarily time-invariant is sometimes called, to emphasize the fact, a time-varying system. Sets of (scalar) state equations describing a linear (time-invariant) control system are the easiest to manage analytically and numerically, and the first model of a situation is often constructed to be linear for this reason.

### 2.1 Solution of Uncontrolled System

To begin with we shall consider dynamical systems (i.e. systems without the presence of control variables). We may also refer to such systems as uncontrolled (or unforced) systems.
We discuss methods of finding the solution (state vector) $x(t)=\left[\begin{array}{c}x_{1}(t) \\ \vdots \\ x_{m}(t)\end{array}\right] \in$ $\mathbb{R}^{m \times 1}$ of the (initialized) linear dynamical system described by

$$
\begin{equation*}
\dot{x}=A x, \quad x(0)=x_{0} . \tag{2.1}
\end{equation*}
$$

Here $A=\left[a_{i j}\right] \in \mathbb{R}^{m \times m}(x \mapsto A x$ represents the linear dynamics $F)$ and $x_{0} \in \mathbb{R}^{m}$ is the initial state.
Note : We identify the column matrix (or vector) $\left[\begin{array}{c}a_{1} \\ \vdots \\ a_{m}\end{array}\right] \in \mathbb{R}^{m \times 1}$ with the $m$-tuple (or point) $\left(a_{1}, a_{2}, \ldots, a_{m}\right) \in \mathbb{R}^{m}$, whenever appropriate. However, we do not identify the row matrix (or covector) $\left[\begin{array}{llll}a_{1} & a_{2} & \cdots & a_{n}\end{array}\right] \in \mathbb{R}^{1 \times n}$ with $\left(a_{1}, a_{2}, \ldots, a_{n}\right)$ (but rather with the linear functional $\left.\left(x_{1}, x_{2}, \cdots, x_{n}\right) \mapsto a_{1} x_{1}+a_{2} x_{2}+\cdots+a_{n} x_{n}\right)$.

$$
\text { We shall assume that all the eigenvalues } \lambda_{1}, \lambda_{2}, \ldots, \lambda_{m} \text { of } A \text { are distinct. }
$$

Note : In fact, in real-life situations, this is not too severe a restriction, since if $A$ does have repeated eigenvalues, very small perturbations in a few of its elements (which will only be known to a certain degree of accuracy) will suffice to separate these equal eigenvalues.

## Spectral form

If $w_{i}$ is an eigenvector corresponding to $\lambda_{i}$, then $w_{1}, w_{2}, \ldots, w_{m}$ are linearly independent (see Exercise 4), so we can express the solution of (2.1)
as

$$
\begin{equation*}
x(t)=c_{1}(t) w_{1}+c_{2}(t) w_{2}+\cdots+c_{m}(t) w_{m} \tag{2.2}
\end{equation*}
$$

where $c_{i}=c_{i}(t), \quad i=1,2, \ldots, m$ are scalar functions of time.
Differentiation of (2.2) and substitution into (2.1) gives :

$$
\sum_{i=1}^{m} \dot{c}_{i}(t) w_{i}=A \sum_{i=1}^{m} c_{i}(t) w_{i}=\sum_{i=1}^{m} c_{i}(t) \lambda_{i} w_{i}
$$

Hence, by the independence of the $w_{i} \quad(i=1,2, \ldots, m)$

$$
\dot{c}_{i}=\lambda_{i} c_{i}, \quad i=1,2, \ldots, m
$$

and these equations have the solution

$$
c_{i}(t)=c_{i}(0) e^{\lambda_{i} t}, \quad i=1,2, \ldots, m
$$

giving

$$
\begin{equation*}
x(t)=\sum_{i=1}^{m} c_{i}(0) e^{\lambda_{i} t} w_{i} \tag{2.3}
\end{equation*}
$$

Let $W$ denote the matrix whose columns are $w_{1}, w_{2}, \ldots, w_{m}$; that is,

$$
W=\left[\begin{array}{llll}
w_{1} & w_{2} & \ldots & w_{m}
\end{array}\right]
$$

We shall denote by $v_{1}, v_{2}, \ldots, v_{m}$ the rows of the matrix $W^{-1}$; that is,

$$
\left[\begin{array}{llll}
w_{1} & w_{2} & \ldots & w_{m}
\end{array}\right]^{-1}=\left[\begin{array}{c}
v_{1} \\
v_{2} \\
\vdots \\
v_{m}
\end{array}\right]
$$

Since we have

$$
v_{i} w_{j}=\delta_{i j}= \begin{cases}1 & \text { if } i=j \\ 0 & \text { if } i \neq j\end{cases}
$$

multiplying (2.3) on the left by $v_{i}$ and setting $t=0$ in the resulting expression gives

$$
v_{i} x(0)=c_{i}(0), \quad i=1,2, \ldots, m
$$

Thus the solution of (2.1) is

$$
\begin{equation*}
x(t)=\sum_{i=1}^{m}\left(v_{i} x(0)\right) e^{\lambda_{i} t} w_{i} \tag{2.4}
\end{equation*}
$$

Expression (2.4) depends only upon the initial condition and the eigenvalues and eigenvectors of $A$, and for this reason is referred to as the spectral form solution.
2.1.1 Example. Find the general solution of the uncontrolled system :

$$
\left[\begin{array}{c}
\dot{x}_{1} \\
\dot{x}_{2}
\end{array}\right]=\left[\begin{array}{rr}
0 & 1 \\
-2 & -3
\end{array}\right]\left[\begin{array}{l}
x_{1} \\
x_{2}
\end{array}\right] .
$$

Solution : The characteristic equation of $A$ is

$$
\left|\lambda I_{2}-A\right|=\left|\begin{array}{rr}
\lambda & -1 \\
2 & \lambda+3
\end{array}\right|=0 \Longleftrightarrow \lambda^{2}+3 \lambda+2=0
$$

which gives

$$
\lambda_{1}=-2 \quad \text { and } \quad \lambda_{2}=-1
$$

(i) $\lambda=-2$.

$$
\left[\begin{array}{rr}
2 & 1 \\
-2 & -1
\end{array}\right]\left[\begin{array}{l}
w_{11} \\
w_{21}
\end{array}\right]=\left[\begin{array}{l}
0 \\
0
\end{array}\right] \Longleftrightarrow\left\{\begin{array}{r}
2 w_{11}+w_{21}=0 \\
-2 w_{11}-w_{21}=0
\end{array}\right.
$$

which implies

$$
w_{21}=-2 w_{11}
$$

and thus (we can choose)

$$
w_{1}=\left[\begin{array}{r}
1 \\
-2
\end{array}\right] .
$$

(ii) $\lambda=-1$.

$$
\left[\begin{array}{rr}
1 & 1 \\
-2 & -2
\end{array}\right]\left[\begin{array}{l}
w_{12} \\
w_{22}
\end{array}\right]=\left[\begin{array}{l}
0 \\
0
\end{array}\right] \Longleftrightarrow\left\{\begin{array}{r}
w_{12}+w_{22}=0 \\
-2 w_{12}-2 w_{22}=0
\end{array}\right.
$$

which implies

$$
w_{22}=-w_{12}
$$

and thus (we can choose)

$$
w_{2}=\left[\begin{array}{r}
1 \\
-1
\end{array}\right] .
$$

We have

$$
W=\left[\begin{array}{ll}
w_{1} & w_{2}
\end{array}\right]=\left[\begin{array}{rr}
1 & 1 \\
-2 & -1
\end{array}\right] \Rightarrow W^{-1}=\left[\begin{array}{rr}
-1 & -1 \\
2 & 1
\end{array}\right]
$$

so

$$
v_{1}=\left[\begin{array}{ll}
-1 & -1
\end{array}\right] \quad \text { and } \quad v_{2}=\left[\begin{array}{cc}
2 & 1
\end{array}\right] .
$$

Finally, we get

$$
\begin{aligned}
x(t) & =\left(v_{1} x(0)\right) e^{-2 t} w_{1}+\left(v_{2} x(0)\right) e^{-t} w_{2} \\
& =\left(-x_{1}(0)-x_{2}(0)\right) e^{-2 t}\left[\begin{array}{r}
1 \\
-2
\end{array}\right]+\left(2 x_{1}(0)+x_{2}(0)\right) e^{-t}\left[\begin{array}{r}
1 \\
-1
\end{array}\right]
\end{aligned}
$$

or

$$
\left\{\begin{array}{l}
x_{1}(t)=-\left(x_{1}(0)+x_{2}(0)\right) e^{-2 t}+\left(2 x_{1}(0)+x_{2}(0)\right) e^{-t} \\
x_{2}(t)=2\left(x_{1}(0)+x_{2}(0)\right) e^{-2 t}-\left(2 x_{1}(0)+x_{2}(0)\right) e^{-t}
\end{array}\right.
$$

## Exponential form

We now present a different approach to solving equation (2.1) which avoids the need to calculate the eigenvectors of $A$.

Recall the definition of the matrix exponential

$$
\begin{equation*}
\exp (t A):=I_{m}+t A+\frac{t^{2}}{2!} A^{2}+\frac{t^{3}}{3!} A^{3}+\cdots \tag{2.5}
\end{equation*}
$$

2.1.2 Lemma. Let $A \in \mathbb{R}^{m \times m}$. Then

$$
\frac{d}{d t}(\exp (t A))=A \exp (t A)=\exp (t A) A
$$

Proof: We have

$$
\begin{aligned}
\frac{d}{d t} \exp (t A) & =\lim _{h \rightarrow 0} \frac{1}{h}(\exp ((t+h) A)-\exp (t A)) \\
& =\lim _{h \rightarrow 0} \frac{1}{h}(\exp (t A) \cdot \exp (h A)-\exp (t A)) \\
& =\exp (t A) \lim _{h \rightarrow 0} \frac{1}{h}\left(\exp (h A)-I_{m}\right) \\
& =\exp (t A) \lim _{h \rightarrow 0} \lim _{k \rightarrow \infty}\left(A+\frac{h}{2!} A+\cdots+\frac{h^{k-1}}{k!} A^{k}\right) \\
& =\exp (t A) A .
\end{aligned}
$$

(Two convergent limit processes can be interchanged if one of them converges uniformly.) Observe that $A$ commutes with each term of the (absolutely convergent) series for $\exp (t A)$, hence with $\exp (t A)$. This proves the lemma.

By the preceding lemma, if $x(t)=\exp (t A) x_{0}$, then

$$
\dot{x}(t)=\frac{d}{d t} \exp (t A) x_{0}=A \exp (t A) x_{0}=A x(t)
$$

for all $t \in \mathbb{R}$. Also,

$$
x(0)=I_{m} x_{0}=x_{0} .
$$

Thus $x(t)=\exp (t A) x_{0}$ is $a$ solution of (2.1).
To see that this is the only solution, let $x(\cdot)$ be any solution of (2.1) and set

$$
y(t)=\exp (-t A) x(t) .
$$

Then (from the above lemma and the fact that $x(\cdot)$ is a solution of (2.1))

$$
\begin{aligned}
\dot{y}(t) & =-A \exp (-t A) x(t)+\exp (-t A) \dot{x}(t) \\
& =-A \exp (-t A) x(t)+\exp (-t A) A x(t) \\
& =0
\end{aligned}
$$

for all $t \in \mathbb{R}$ since $\exp (-t A)$ and $A$ commute. Thus, $y(t)$ is a constant. Setting $t=0$ shows that $y(t)=x_{0}$ and therefore any solution of (2.1) is given by

$$
x(t)=\exp (t A) y(t)=\exp (t A) x_{0}
$$

Hence

$$
\begin{equation*}
x(t)=\exp (t A) x_{0} \tag{2.6}
\end{equation*}
$$

does represent the solution of (2.1).
Note : In case the initial condition $x(0)=x_{0}$ is replaced by the slighty more general one $x\left(t_{0}\right)=x_{0}$, the solution is often written as

$$
\begin{equation*}
x(t)=\Phi\left(t, t_{0}\right) x_{0} . \tag{2.7}
\end{equation*}
$$

One refers to (the matrix)

$$
\begin{equation*}
\Phi\left(t, t_{0}\right):=\exp \left(\left(t-t_{0}\right) A\right) \tag{2.8}
\end{equation*}
$$

as the state transition matrix (since it relates the state at any time $t$ to the state at any other time $t_{0}$ ).
2.1.3 Proposition. The state transition matrix $\Phi\left(t, t_{0}\right)$ has the following properties :
(a) $\frac{d}{d t} \Phi\left(t, t_{0}\right)=A \Phi\left(t, t_{0}\right)$.
(b) $\Phi\left(t_{0}, t_{0}\right)=I_{m}$.
(c) $\Phi\left(t_{0}, t\right)=\Phi^{-1}\left(t, t_{0}\right)$.
(d) $\Phi\left(t, t_{0}\right)=\Phi\left(t, t_{1}\right) \Phi\left(t_{1}, t_{0}\right)$.

Proof : We have
(a) $\frac{d}{d t} \Phi\left(t, t_{0}\right)=\frac{d}{d t} \exp \left(\left(t-t_{0}\right) A\right)=A \exp \left(\left(t-t_{0}\right) A\right)=A \Phi\left(t, t_{0}\right)$.
(b) $\Phi(t, t)=\exp ((t-t) A)=\exp (0)=I_{m}$.
(c) $\Phi^{-1}\left(t, t_{0}\right)=\left(\exp \left(\left(t-t_{0}\right) A\right)\right)^{-1}=\exp \left(-\left(t-t_{0}\right) A\right)=\exp \left(\left(t_{0}-t\right) A\right)=$ $\Phi\left(t_{0}, t\right)$.
(d) $\Phi\left(t, t_{1}\right) \Phi\left(t_{1}, t_{0}\right)=\exp \left(\left(t-t_{1}\right) A\right) \cdot \exp \left(\left(t_{1}-t_{0}\right) A\right)=\exp \left(\left(t-t_{1}+t_{1}-t_{0}\right) A\right)=$ $\exp \left(\left(t-t_{0}\right) A\right)=\Phi\left(t, t_{0}\right)$.

Note : The matrix-valued mapping $X(t)=\Phi\left(t, t_{0}\right)$ (or curve in $\mathbb{R}^{m \times m}$ ) is the unique solution of the matrix differential equation

$$
\dot{X}=A X, \quad X \in \mathbb{R}^{m \times m}
$$

subject to the initial condition $X\left(t_{0}\right)=I_{m}$.
2.1.4 Example. (Simple harmonic motion) Consider a unit mass connected to a support through a spring whose spring constant is unity. If $z$ measures the displacement of the mass from equilibrium, then

$$
\ddot{z}+z=0 .
$$

Letting $x_{1}=z$ and $x_{2}=\dot{z}$ gives

$$
\left[\begin{array}{c}
\dot{x}_{1}(t) \\
\dot{x}_{2}(t)
\end{array}\right]=\left[\begin{array}{rr}
0 & 1 \\
-1 & 0
\end{array}\right]\left[\begin{array}{l}
x_{1}(t) \\
x_{2}(t)
\end{array}\right] .
$$

Note : The associated transition matrix $\Phi\left(t, t_{0}\right)$ has the form

$$
\Phi\left(t, t_{0}\right)=\left[\begin{array}{ll}
\phi_{11}\left(t, t_{0}\right) & \phi_{12}\left(t, t_{0}\right) \\
\phi_{21}\left(t, t_{0}\right) & \phi_{22}\left(t, t_{0}\right)
\end{array}\right]
$$

and therefore satisfies

$$
\left[\begin{array}{ll}
\dot{\phi}_{11}\left(t, t_{0}\right) & \dot{\phi}_{12}\left(t, t_{0}\right) \\
\dot{\phi}_{21}\left(t, t_{0}\right) & \dot{\phi}_{22}\left(t, t_{0}\right)
\end{array}\right]=\left[\begin{array}{cc}
0 & 1 \\
-1 & 0
\end{array}\right]\left[\begin{array}{ll}
\phi_{11}\left(t, t_{0}\right) & \phi_{12}\left(t, t_{0}\right) \\
\phi_{21}\left(t, t_{0}\right) & \phi_{22}\left(t, t_{0}\right)
\end{array}\right]
$$

with the initial condition

$$
\Phi\left(t_{0}, t_{0}\right)=I_{m} .
$$

What is the physical interpretation of $\Phi\left(t, t_{0}\right)$ in this case? The first column of $\Phi\left(t, t_{0}\right)$ has as its first entry the position as a function of time which results when
the mass is displaced by one unit and released at $t_{0}$ with zero velocity. The second entry in the first column is the corresponding velocity. The second column of $\Phi\left(t, t_{0}\right)$ has as its first entry the position as a function of time which results when the mass is started from zero displacement but with unit velocity at $t=t_{0}$. The second entry in the second column is the corresponding velocity.

The series for computing $\Phi\left(t, t_{0}\right)$ in this case is easily summed because

$$
A^{k}=\left\{\begin{aligned}
A & \text { if } k=4 p+1 \\
-I_{2} & \text { if } k=4 p+2 \\
-A & \text { if } k=4 p+3 \\
I_{2} & \text { if } k=4 p
\end{aligned}\right.
$$

A short calculation gives

$$
\Phi\left(t, t_{0}\right)=\left[\begin{array}{rr}
\cos \left(t-t_{0}\right) & \sin \left(t-t_{0}\right) \\
-\sin \left(t-t_{0}\right) & \cos \left(t-t_{0}\right)
\end{array}\right]
$$

Exercise 21 Work out the preceding computation.
2.1.5 Example. (Satellite problem) In section 1.3 we introduced the equations of a unit mass in an inverse square law force field. These were then linearized about a circular orbit to get

$$
\left[\begin{array}{c}
\dot{x}_{1}(t) \\
\dot{x}_{2}(t) \\
\dot{x}_{3}(t) \\
\dot{x}_{4}(t)
\end{array}\right]=\left[\begin{array}{cccc}
0 & 1 & 0 & 0 \\
3 \omega^{2} & 0 & 0 & 2 \omega \\
0 & 0 & 0 & 1 \\
0 & -2 \omega & 0 & 0
\end{array}\right]\left[\begin{array}{c}
x_{1}(t) \\
x_{2}(t) \\
x_{3}(t) \\
x_{4}(t)
\end{array}\right]+\left[\begin{array}{c}
0 \\
u_{1}(t) \\
0 \\
u_{2}(t)
\end{array}\right]
$$

The series for computing $\Phi(t, 0)$ can be summed to get

$$
\Phi(t, 0)=\left[\begin{array}{cccc}
4-3 \cos \omega t & \frac{\sin \omega t}{\omega} & 0 & \frac{2(1-\cos \omega t)}{\omega} \\
3 \omega \sin \omega t & \cos \omega t & 0 & 2 \sin \omega t \\
6(-\omega t+\sin \omega t) & -\frac{2(1-\cos \omega t)}{\omega} & 1 & \frac{-3 \omega t+4 \sin \omega t}{\omega} \\
6 \omega(-1+\cos \omega t) & -2 \sin \omega t & 0 & -3+4 \cos \omega t
\end{array}\right]
$$

## Evaluation of the matrix exponential

Evaluation of $\exp (t A)$, when all the eigenvalues $\lambda_{1}, \lambda_{2}, \ldots, \lambda_{m}$ are distinct, can be achieved by Sylvester's formula which gives

$$
\begin{equation*}
\exp (t A)=\sum_{k=1}^{m} e^{\lambda_{k} t} Z_{k} \tag{2.9}
\end{equation*}
$$

where

$$
Z_{k}=\prod_{\substack{i=1 \\ i \neq k}}^{m} \frac{A-\lambda_{i} I_{m}}{\lambda_{k}-\lambda_{i}}, \quad k=1,2, \ldots, m .
$$

Note : Since the $Z_{k}(k=1,2, \ldots, m)$ in (2.9) are constant matrices depending only on $A$ and its eigenvalues, the solution in the form given in (2.9) requires calculation of only the eigenvalues of $A$.
2.1.6 Example. Consider again the uncontrolled system

$$
\left[\begin{array}{c}
\dot{x}_{1} \\
\dot{x}_{2}
\end{array}\right]=\left[\begin{array}{rr}
0 & 1 \\
-2 & -3
\end{array}\right]\left[\begin{array}{l}
x_{1} \\
x_{2}
\end{array}\right] .
$$

The solution is

$$
x(t)=\exp (t A) x_{0}=\left(\sum_{k=1}^{2} e^{\lambda_{k} t} Z_{k}\right) x_{0}
$$

where

$$
Z_{k}=\prod_{\substack{i=1 \\ i \neq k}}^{2} \frac{A-\lambda_{i} I_{2}}{\lambda_{k}-\lambda_{i}}, \quad k=1,2 .
$$

We have

$$
Z_{1}=\frac{A-(-1) I_{2}}{-2-(-1)}=\left[\begin{array}{rr}
-1 & -1 \\
2 & 2
\end{array}\right] ; \quad Z_{2}=\frac{A-(-2) I_{2}}{-1-(-2)}=\left[\begin{array}{rr}
2 & 1 \\
-2 & -1
\end{array}\right] .
$$

Hence,

$$
\begin{aligned}
x(t) & =\left(e^{-2 t} Z_{1}+e^{-t} Z_{2}\right) x_{0} \\
& =\left(e^{-2 t}\left[\begin{array}{rr}
-1 & -1 \\
2 & 2
\end{array}\right]+e^{-t}\left[\begin{array}{rr}
2 & 1 \\
-2 & -1
\end{array}\right]\right)\left[\begin{array}{l}
x_{1}(0) \\
x_{2}(0)
\end{array}\right]
\end{aligned}
$$

or

$$
\left\{\begin{array}{l}
x_{1}(t)=-\left(x_{1}(0)+x_{2}(0)\right) e^{-2 t}+\left(2 x_{1}(0)+x_{2}(0)\right) e^{-t} \\
x_{2}(t)=2\left(x_{1}(0)+x_{2}(0)\right) e^{-2 t}-\left(2 x_{1}(0)+x_{2}(0)\right) e^{-t}
\end{array}\right.
$$

An alternative way of evaluating $\exp (t A)$ (again, when all the eigenvalues $\lambda_{1}, \lambda_{2}, \ldots, \lambda_{m}$ are distinct), is as follows.

We can write

$$
\begin{equation*}
e^{t \lambda}=q(\lambda) \cdot \operatorname{char}_{A}(\lambda)+r(\lambda) \tag{2.10}
\end{equation*}
$$

where $\operatorname{deg}(r)<m$. Since (2.10) is an identity, we have

$$
\exp (t A) \equiv q(A) \cdot \operatorname{char}_{A}(A)+r(A)
$$

which by the Cayley-Hamilton Theorem reduces to

$$
\exp (t A) \equiv r(A)
$$

showing that $\exp (t A)$ can be represented by a finite sum of powers of $A$ of degree not exceeding $m-1$. Then $m$ coefficients of $r(\lambda)$ are functions of $t$ obtained from the solution of the system of $m$ linear equations

$$
e^{\lambda_{i} t}=r\left(\lambda_{i}\right), \quad i=1,2, \ldots, m .
$$

2.1.7 Example. Consider once again the uncontrolled system

$$
\left[\begin{array}{c}
\dot{x}_{1} \\
\dot{x}_{2}
\end{array}\right]=\left[\begin{array}{rr}
0 & 1 \\
-2 & -3
\end{array}\right]\left[\begin{array}{l}
x_{1} \\
x_{2}
\end{array}\right] .
$$

Since $m=2$, the polynomial $r(\lambda)$ can be written

$$
r(\lambda)=r_{0} \lambda+r_{1}
$$

and so we have

$$
\left\{\begin{array}{l}
e^{-t}=r_{1}-r_{0} \\
e^{-2 t}=r_{1}-2 r_{0}
\end{array}\right.
$$

which gives $r_{0}=e^{-t}-e^{-2 t}$ and $r_{1}=2 e^{-t}-e^{-2 t}$. Hence, the solution is

$$
\begin{aligned}
x(t) & =\exp (t A) x_{0}=\left(r_{0} A+r_{1} I_{2}\right) x_{0} \\
& =\left(\left(e^{-t}-e^{-2 t}\right)\left[\begin{array}{rr}
0 & 1 \\
-2 & -3
\end{array}\right]+\left(2 e^{-t}-e^{-2 t}\right)\left[\begin{array}{ll}
1 & 0 \\
0 & 1
\end{array}\right]\right)\left[\begin{array}{l}
x_{1}(0) \\
x_{2}(0)
\end{array}\right]
\end{aligned}
$$

or

$$
\left\{\begin{array}{l}
x_{1}(t)=-\left(x_{1}(0)+x_{2}(0)\right) e^{-2 t}+\left(2 x_{1}(0)+x_{2}(0)\right) e^{-t} \\
x_{2}(t)=2\left(x_{1}(0)+x_{2}(0)\right) e^{-2 t}-\left(2 x_{1}(0)+x_{2}(0)\right) e^{-t}
\end{array}\right.
$$

### 2.2 Solution of Controlled System

Consider the (initialized) linear control system, written in state space form,

$$
\begin{equation*}
\dot{x}=A x+B u(t), \quad x(0)=x_{0} \tag{2.11}
\end{equation*}
$$

where $A \in \mathbb{R}^{m \times m}, B \in \mathbb{R}^{m \times \ell}$, and $\ell \leq m$.
After multiplication of both sides of (2.11), on the left, by $\exp (-t A)$, the equation can be written

$$
\begin{equation*}
\frac{d}{d t}(\exp (-t A) x)=\exp (-t A) B u \tag{2.12}
\end{equation*}
$$

which produces

$$
\begin{equation*}
x(t)=\exp (t A)\left[x_{0}+\int_{0}^{t} \exp (-\tau A) B u(\tau) d \tau\right] . \tag{2.13}
\end{equation*}
$$

If the initial condition is $x\left(t_{0}\right)=x_{0}$, then integration of (2.12) from $t_{0}$ to $t$ and use of the definition of $\Phi$ gives

$$
\begin{equation*}
x(t)=\Phi\left(t, t_{0}\right)\left[x_{0}+\int_{t_{0}}^{t} \Phi\left(t_{0}, \tau\right) B u(\tau) d \tau\right] \tag{2.14}
\end{equation*}
$$

Note : If $u(t)$ is known for $t \geq t_{0}$, then $x(t)$ can be determined by finding the state transition matrix and carrying out the integration in (2.14).
2.2.1 Example. Consider the equation of motion

$$
\ddot{z}=u(t)
$$

of a unit mass moving in a straight line, subject to an external force $u(t), z(t)$ being the displacement from some fixed point. In state space form, taking

$$
x_{1}=z \quad \text { and } \quad x_{2}=\dot{z}
$$

as state variables, this becomes

$$
\dot{x}=\left[\begin{array}{c}
\dot{x}_{1} \\
\dot{x}_{2}
\end{array}\right]=\left[\begin{array}{ll}
0 & 1 \\
0 & 0
\end{array}\right]\left[\begin{array}{l}
x_{1} \\
x_{2}
\end{array}\right]+\left[\begin{array}{l}
0 \\
1
\end{array}\right] u(t)=A x+B u(t) .
$$

Since here we have $A^{2}=0, \exp (t A)=I_{2}+t A$, and so

$$
\left[\begin{array}{l}
x_{1}(t) \\
x_{2}(t)
\end{array}\right]=\left[\begin{array}{ll}
1 & t \\
0 & 1
\end{array}\right]\left[\begin{array}{l}
x_{1}(0) \\
x_{2}(0)
\end{array}\right]+\left[\begin{array}{ll}
1 & t \\
0 & 1
\end{array}\right] \int_{0}^{t}\left[\begin{array}{rr}
1 & -\tau \\
0 & 1
\end{array}\right]\left[\begin{array}{l}
0 \\
1
\end{array}\right] u(\tau) d \tau
$$

Solving for $x_{1}(t)$ leads to

$$
z(t)=z(0)+t \dot{z}(0)+\int_{o}^{t}(t-\tau) u(\tau) d \tau
$$

where $\dot{z}(0)$ denotes the initial velocity of the mass.
2.2.2 Example. We are now in a position to express the solution of the linearized equations describing the motion of a satellite in a near circular orbit. We have

$$
\left.\begin{array}{l}
{\left[\begin{array}{l}
x_{1}(t) \\
x_{2}(t) \\
x_{3}(t) \\
x_{4}(t)
\end{array}\right]=\left[\begin{array}{cccc}
4-3 \cos \omega t & \frac{\sin \omega t}{\omega} & 0 & \frac{2(1-\cos \omega t)}{\omega} \\
3 \omega \sin \omega t & \cos \omega t & 0 & 2 \sin \omega t \\
6(-\omega t+\sin \omega t) & -\frac{2(1-\cos \omega t)}{\omega} & 1 & \frac{-3 \omega t+4 \sin \omega t}{\omega} \\
6 \omega(-1+\cos \omega t) & -2 \sin \omega t & 0 & -3+4 \cos \omega t
\end{array}\right]\left[\begin{array}{l}
x_{1}(0) \\
x_{2}(0) \\
x_{0} \\
+x_{3}(0) \\
x_{4}(0)
\end{array}\right]+} \\
-\frac{2(1-\cos \omega(t-\tau))}{\omega} \\
-2\left[\begin{array}{c}
\frac{2(1-\cos \omega(t-\tau))}{\omega} \\
-2 \sin \omega(t-\tau)
\end{array}\right] u_{1}(\tau)+\left[\begin{array}{c}
2 \sin \omega(t-\tau) \\
\frac{-3 \omega(t-\tau)+4 \sin \omega(t-\tau)}{\omega} \\
-3+4 \cos \omega(t-\tau)
\end{array}\right] u_{2}(\tau)
\end{array}\right] d \tau .
$$

### 2.3 Time-varying Systems

Of considerable importance in many applications are linear systems in which the elements of $A$ and $B$ are (continuous) functions of time for $t \geq 0$.

Note : In general, it will not be possible to give explicit expressions for solutions and we shall content ourselves with obtaining some general properties.

We first consider the uncontrolled case

$$
\begin{equation*}
\dot{x}=A(t) x, \quad x(0)=x_{0} . \tag{2.15}
\end{equation*}
$$

2.3.1 Theorem. (Existence and Uniqueness Theorem) If the matrixvalued mapping

$$
A:[0, \infty) \rightarrow \mathbb{R}^{m \times m}, \quad t \mapsto A(t)
$$

is continuous, then (2.15) has a unique solution

$$
x(t)=X(t) x_{0}, \quad t \geq 0
$$

where $X(\cdot)$ is the unique matrix-valued mapping (or curve in $\mathbb{R}^{m \times m}$ ) satisfying

$$
\begin{equation*}
\dot{X}=A(t) X, \quad X(0)=I_{m} \tag{2.16}
\end{equation*}
$$

Proof : We shall use the method of successive approximations to establish the existence of a solution of (2.16). In place of (2.16), we consider the integral equation

$$
\begin{equation*}
X=I_{m}+\int_{0}^{t} A(\tau) X d \tau \tag{2.17}
\end{equation*}
$$

Define the sequence $\left(X_{k}\right)_{k \geq 0}$ of matrices (in fact, of matrix-valued mappings) as follows :

$$
\begin{aligned}
X_{0} & =I_{m} \\
X_{k+1} & =I_{m}+\int_{0}^{t} A(\tau) X_{k} d \tau, \quad k=0,1,2, \ldots
\end{aligned}
$$

Then we have

$$
X_{k+1}-X_{k}=\int_{0}^{t} A(\tau)\left(X_{k}-X_{k-1}\right) d \tau, \quad k=1,2, \ldots
$$

Let

$$
\nu=\max _{0 \leq t \leq t_{1}}\|A(t)\|
$$

where

$$
\|A(t)\|:=\sum_{i, j=1}^{m}\left|a_{i j}(t)\right|
$$

Note : Any matrix norm (on $\mathbb{R}^{m \times m}$ ) will do.
We have

$$
\begin{aligned}
\left\|X_{k+1}-X_{k}\right\| & =\left\|\int_{0}^{t} A(\tau)\left(X_{k}-X_{k-1}\right) d \tau\right\| \\
& \leq \int_{0}^{t}\|A(\tau)\|\left\|X_{k}-X_{k-1}\right\| d \tau \\
& \leq \nu \int_{0}^{t}\left\|X_{k}-X_{k-1}\right\| d \tau
\end{aligned}
$$

for $0 \leq t \leq t_{1}$. Since, in this same interval,

$$
\left\|X_{1}-X_{0}\right\| \leq \int_{0}^{t}\|A(\tau)\| d \tau \leq \nu t
$$

we have inductively

$$
\left\|X_{k+1}-X_{k}\right\| \leq M_{k+1}:=\frac{\nu^{k+1} t_{1}^{k+1}}{(k+1)!} \quad \text { for } \quad 0 \leq t \leq t_{1}
$$

Note : The Weierstrass M-Test states that if $\xi_{k}:\left[t_{0}, t_{1}\right] \rightarrow \mathbb{R}^{m \times m}$ are continuous and

- $\left\|\xi_{k}(t)\right\| \leq M_{k}$ for every $k$
- $\sum_{k=0}^{\infty} M_{k}<\infty$
then the series $\sum_{k \geq 0} \xi_{k}(t)$ converges uniformly and absolutely on the interval $\left[t_{0}, t_{1}\right]$.
Hence, the (matrix-valued mapping) series

$$
X_{0}+\sum_{k \geq 0}\left(X_{k+1}-X_{k}\right)
$$

converges uniformly for $0 \leq t \leq t_{1}$. Consequently, $\left(X_{k}\right)$ converges uniformly and absolutely to a matrix-valued mapping $X(\cdot)$, which satisfies (2.17), and thus (2.16).

Since, by assumption, $A(\cdot)$ is continuous for $t \geq 0$, we must take $t_{1}$ arbitrarily large. We thus obtain a solution valid for $t \geq 0$.

It is easy to verify that $x(t)=X(t) x_{0}$ is a solution of (2.15), satisfying the required initial condition.

Let us now establish uniqueness of this solution. Let $Y$ be another solution of (2.16). Then $Y$ satisfies (2.17), and thus we have the relation

$$
X-Y=\int_{0}^{t} A(\tau)(X(\tau)-Y(\tau)) d \tau
$$

Hence

$$
\|X-Y\| \leq \int_{0}^{t}\|A(\tau)\|\|X(\tau)-Y(\tau)\| d \tau
$$

Since $Y$ is differentiable, hence continuous, define

$$
\nu_{1}:=\max _{0 \leq t \leq t_{1}}\|X(t)-Y(t)\| .
$$

We obtain

$$
\|X-Y\| \leq \nu_{1} \int_{0}^{t}\|A(\tau)\| d \tau, \quad 0 \leq t \leq t_{1}
$$

Using this bound, we obtain

$$
\begin{aligned}
\|X-Y\| & \leq \nu_{1} \int_{0}^{t}\|A(\tau)\|\left(\int_{0}^{\tau}\|A(\sigma)\| d \sigma\right) d \tau \\
& \leq \frac{\nu_{1}\left(\int_{0}^{t}\|A(\tau)\| d \tau\right)^{2}}{2} .
\end{aligned}
$$

Iterating, we get

$$
\|X-Y\| \leq \frac{\nu_{1}^{k}\left(\int_{0}^{t}\|A(\tau)\| d \tau\right)^{k+1}}{(k+1)!}
$$

Letting $k \rightarrow \infty$, we see that $\|X-Y\| \leq 0$. Hence $X \equiv Y$.
Exercise 22 Show that

$$
\lim _{k \rightarrow \infty} \frac{\alpha^{k}}{(k+1)!}=0 \quad(\alpha>0) .
$$

Having obtained the matrix $X$, it is easy to see that $x(t)=X(t) x_{0}$ is a solution of (2.15). Since the uniqueness of solutions of (2.15) is readily established by means of the same argument as above, it is easy to see that $x(t)=X(t) x_{0}$ is the solution.

Note : We can no longer define a matrix exponential, but there is a result corresponding to the fact that $\exp (t A)$ is nonsingular when $A$ is constant. We can write $x(t)=\Phi(t, 0) x_{0}$, where $\Phi(t, 0)$ has the form
$I_{m}+\int_{0}^{t} A(\tau) d \tau+\int_{0}^{t} A\left(\tau_{1}\right) \int_{0}^{\tau_{1}} A\left(\tau_{2}\right) d \tau_{2} d \tau_{1}+\int_{0}^{t} A\left(\tau_{1}\right) \int_{0}^{\tau_{1}} A\left(\tau_{2}\right) \int_{0}^{\tau_{2}} A\left(\tau_{3}\right) d \tau_{3} d \tau_{2} d \tau_{1} \cdots$ (the Peano-Baker series).

## Some remarks and corollaries

2.3.2 Proposition. In the Existence and Uniqueness Theorem the matrix $X(t)$ is nonsingular (for every $t \geq 0$ ).

Proof : Define a matrix-valued mapping $Y(\cdot)$ as the solution of

$$
\begin{equation*}
\dot{Y}=-Y A(t), \quad Y(0)=I_{m} \tag{2.18}
\end{equation*}
$$

(Such a mapping exists and is unique by an argument virtually identical to that which is used in the proof of the Existence and Uniqueness Theorem.)

Now

$$
\frac{d}{d t}(Y X)=\dot{Y} X+Y \dot{X}=-Y A X+Y A X=0
$$

so $Y(t) X(t)$ is equal to a constant matrix, which must be the unit matrix because of condition at $t=0$.

Exercise 23 Show that (for every $t \geq 0$ )

$$
\operatorname{det}(X(t)) \neq 0
$$

Hence $X(t)$ is nonsingular (and its inverse is in fact $Y(t)$ ).

We can also generalize the idea of state transition matrix by writing

$$
\begin{equation*}
\Phi\left(t, t_{0}\right):=X(t) X^{-1}\left(t_{0}\right) \tag{2.19}
\end{equation*}
$$

which exists for all $t, t_{0} \geq 0$. It is easy to verify that

$$
\begin{equation*}
x(t)=\Phi\left(t, t_{0}\right) x_{0} \tag{2.20}
\end{equation*}
$$

is the solution of $(2.15)$. Also, $\Phi\left(t, t_{0}\right)^{-1}=\Phi\left(t_{0}, t\right)$.
Note : The expression (2.20) has the same form as that for the time invariant case. However, it is most interesting that although, in general, it is not possible to obtain an analytic expression for the solution of (2.16), and therefore for $\Phi\left(t, t_{0}\right)$ in (2.20), this latter matrix possesses precisely the same properties as those for the constant case.

When $m=1$, we can see that

$$
\Phi\left(t, t_{0}\right)=\exp \left(\int_{t_{0}}^{t} A(\tau) d \tau\right) .
$$

Note that the formula above is generally not true. However, one can show that it does hold if

$$
A(t) \int_{t_{0}}^{t} A(\tau) d \tau=\left(\int_{t_{0}}^{t} A(\tau) d \tau\right) A(t) \quad \text { for all } t .
$$

Otherwise this is not necessarily true and the state transition matrix is not necessarily the exponential of the integral of $A$.

The following result is interesting. We shall omit the proof.
2.3.3 Proposition. If $\Phi\left(t, t_{0}\right)$ is the state transition matrix for

$$
\dot{x}=A(t) x, \quad x\left(t_{0}\right)=x_{0}
$$

then

$$
\operatorname{det}\left(\Phi\left(t, t_{0}\right)\right)=e^{\int_{t_{0}}^{t} \operatorname{tr}(A(\tau)) d \tau}
$$

A further correspondence with the time-invariant case is the following result.
2.3.4 Proposition. The solution of

$$
\begin{equation*}
\dot{x}=A(t) x+B(t) u(t), \quad x\left(t_{0}\right)=x_{0} \tag{2.21}
\end{equation*}
$$

is given by

$$
\begin{equation*}
x(t)=\Phi\left(t, t_{0}\right)\left[x_{0}+\int_{t_{0}}^{t} \Phi\left(t_{0}, \tau\right) B(\tau) u(\tau) d \tau\right] \tag{2.22}
\end{equation*}
$$

where $\Phi\left(t, t_{0}\right)$ is defined in (2.19).
Proof: Put

$$
x=X(t) w \Longleftrightarrow X^{-1}(t) x=w
$$

where $X(t)$ is defined in (2.16). Substitution into (2.21) produces

$$
\dot{x}=A X w+X \dot{w}=A X w+B u(t)
$$

Hence $X(t) \dot{w}=B u(t)$ and so $\dot{w}=X^{-1}(t) B u(t)$, which gives

$$
w(t)=w\left(t_{0}\right)+\int_{t_{0}}^{t} X^{-1}(\tau) B(\tau) u(\tau) d \tau
$$

The desired expression then follows using $x_{0}=X\left(t_{0}\right) w\left(t_{0}\right)$ and (2.19). Indeed,

$$
x_{0}=x\left(t_{0}\right)=X\left(t_{0}\right) w\left(t_{0}\right) \Rightarrow w\left(t_{0}\right)=X^{-1}\left(t_{0}\right) x_{0}
$$

and we have

$$
\begin{aligned}
x(t)=X(t) w & =X(t)\left[X^{-1}\left(t_{0}\right) x_{0}+\int_{t_{0}}^{t} X^{-1}(\tau) B(\tau) u(\tau) d \tau\right] \\
& =X(t) X^{-1}\left(t_{0}\right)\left[x_{0}+\int_{t_{0}}^{t} X\left(t_{0}\right) X^{-1}(\tau) B(\tau) u(\tau) d \tau\right] \\
& =\Phi\left(t, t_{0}\right)\left[x_{0}+\int_{t_{0}}^{t} \Phi\left(t_{0}, \tau\right) B(\tau) u(\tau) d \tau\right]
\end{aligned}
$$

### 2.4 Relationship between State Space and Classical Forms

Classical linear control theory deals with scalar ODEs of the form

$$
z^{(m)}+k_{1} z^{(m-1)}+\cdots+k_{m-1} z^{(1)}+k_{m} z=\beta_{0} u^{(\ell)}+\beta_{1} u^{(\ell-1)}+\cdots+\beta_{\ell} u
$$

where $k_{1}, k_{2}, \ldots, k_{m}$ and $\beta_{0}, \beta_{1}, \ldots, \beta_{\ell}$ are constants; it is assumed that $\ell<$ $m$.

We shall consider a simplified form

$$
\begin{equation*}
z^{(m)}+k_{1} z^{(m-1)}+\cdots+k_{m} z=u(t) \tag{2.23}
\end{equation*}
$$

where $u(\cdot)$ is the single control variable.

It is easy to write (2.23) in matrix form by taking as state variables

$$
\begin{equation*}
w_{1}=z, \quad w_{2}=z^{(1)}, \quad \ldots, \quad w_{m}=z^{(m-1)} \tag{2.24}
\end{equation*}
$$

Since

$$
\dot{w}_{i}=w_{i+1}, \quad i=1,2, \ldots, m-1
$$

(2.23) and (2.24) lead to the state space form

$$
\begin{equation*}
\dot{w}=C w+d u(t) \tag{2.25}
\end{equation*}
$$

where

$$
C=\left[\begin{array}{ccccc}
0 & 1 & 0 & \ldots & 0  \tag{2.26}\\
0 & 0 & 1 & \ldots & 0 \\
0 & 0 & 0 & \ldots & 0 \\
\vdots & \vdots & \vdots & & \vdots \\
0 & 0 & 0 & \ldots & 1 \\
-k_{m} & -k_{m-1} & -k_{m-2} & \ldots & -k_{1}
\end{array}\right], w=\left[\begin{array}{c}
w_{1} \\
w_{2} \\
\vdots \\
w_{m}
\end{array}\right], d=\left[\begin{array}{c}
0 \\
0 \\
\vdots \\
1
\end{array}\right]
$$

The matrix $C$ is called the companion form matrix.

Exercise 24 Show that the characteristic polynomial of $C$ is

$$
\begin{equation*}
\operatorname{char}_{C}(\lambda)=\lambda^{m}+k_{1} \lambda^{m-1}+k_{2} \lambda^{m-2}+\cdots+k_{m} \tag{2.27}
\end{equation*}
$$

which has the same coefficients as those in (2.23).

The state space form (2.25) is very special; we call it the canonical form.
Note : The classical form (2.23) and the canonical form are equivalent.
Having seen that (2.23) can be put into matrix form, a natural question is to ask whether the converse hold : can any linear system in state space form with a single control variable

$$
\dot{x}=A x+b u(t)
$$

be put into the classical form (2.23)?
2.4.1 Example. Consider the linear system in state space form with a single control variable

$$
\left[\begin{array}{l}
\dot{x}_{1}  \tag{2.28}\\
\dot{x}_{2}
\end{array}\right]=\left[\begin{array}{rr}
-2 & 2 \\
1 & -1
\end{array}\right]\left[\begin{array}{l}
x_{1} \\
x_{2}
\end{array}\right]+\left[\begin{array}{l}
1 \\
0
\end{array}\right] u(t) .
$$

The (state) equations describing the system are

$$
\left\{\begin{array}{l}
\dot{x}_{1}=-2 x_{1}+2 x_{2}+u(t) \\
\dot{x}_{2}=x_{1}-x_{2}
\end{array}\right.
$$

Differentiating the second equation we get

$$
\ddot{x}_{2}=\dot{x}_{1}-\dot{x}_{2}=\left(-2 x_{1}+2 x_{2}+u(t)\right)-\left(x_{1}-x_{2}\right)=-3 x_{1}+3 x_{2}+u(t) .
$$

Hence

$$
\ddot{x}_{2}+3 \dot{x}_{2}=u(t) .
$$

This second-order ODE for $x_{2}$ has the form (2.23) for $m=2$. Its associated canonical form is

$$
\dot{w}=\left[\begin{array}{rr}
0 & 1  \tag{2.29}\\
0 & -3
\end{array}\right] w+\left[\begin{array}{l}
0 \\
1
\end{array}\right] u(t) .
$$

We expect that there is a transformation on $\mathbb{R}^{2}$ that transforms our original control system (2.28) to the canonical form (2.29). Since the differential equations are linear, we expect that the transformation is linear, say $w=T x$. Differentiation than gives

$$
\dot{w}=T A T^{-1} w+T b u(t) .
$$

If we set

$$
w_{1}=x_{2} \quad \text { and } \quad w_{2}=\dot{x}_{2}
$$

then

$$
\left\{\begin{array}{l}
w_{1}=x_{2} \\
w_{2}=x_{1}-x_{2}
\end{array}\right.
$$

so a transformation transforming (2.28) into (2.29) is given by the matrix

$$
T=\left[\begin{array}{rr}
0 & 1 \\
1 & -1
\end{array}\right]
$$

2.4.2 EXAMPLE. The control system in state space form

$$
\dot{x}=\left[\begin{array}{ll}
1 & 0 \\
0 & 1
\end{array}\right] x+b u(t)
$$

(where $b \in \mathbb{R}^{2 \times 1}$ ) cannot be transformed to the canonical form (or, equivalently, to the classical form).

Exercise 25 Prove the preceding statement.

We would like to determine when and how such a procedure can be carried out in general. Thus the natural questions concerning existence, uniqueness, and computation of $T$ arise. The answer to these questions is provided by the following result.
2.4.3 Theorem. A linear control system in state space form

$$
\dot{x}=A x+b u(t)
$$

(where $A \in \mathbb{R}^{m \times m}$ and $0 \neq b \in \mathbb{R}^{m \times 1}$ ) can be transformed by a linear transformation (i.e. invertible linear mapping)

$$
w=T x
$$

into the canonical form

$$
\dot{w}=C w+d u(t)
$$

where $C$ and $d$ are given by (2.26), provided

$$
\operatorname{rank}\left[\begin{array}{lllll}
b & A b & A^{2} b & \ldots & A^{m-1} b \tag{2.30}
\end{array}\right]=m
$$

Conversely, if such a transformation $T$ exists, then (2.30) holds.

Proof : $\quad(\Leftarrow)$ Sufficiency. Substitution of $w=T x$ into

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

produces

$$
\dot{w}=T A T^{-1} w+T b u
$$

We take

$$
T=\left[\begin{array}{c}
\tau \\
\tau A \\
\tau A^{2} \\
\vdots \\
\tau A^{m-1}
\end{array}\right]
$$

where $\tau$ is any row $m$-vector such that $T$ is nonsingular, assuming for the present that at least one suitable $\tau$ exists.

Denote the columns of $T^{-1}$ by $s_{1}, s_{2}, \ldots, s_{m}$ and consider

$$
T A T^{-1}=\left[\begin{array}{cccc}
\tau A s_{1} & \tau A s_{2} & \ldots & \tau A s_{m} \\
\tau A^{2} s_{1} & \tau A^{2} s_{2} & \ldots & \tau A^{2} s_{m} \\
\vdots & \vdots & & \vdots \\
\tau A^{m} s_{1} & \tau A^{m} s_{2} & \ldots & \tau A^{m} s_{m}
\end{array}\right]
$$

Comparison with the identity $T T^{-1}=I_{m}$ (that is,

$$
\left.\left[\begin{array}{cccc}
\tau s_{1} & \tau s_{2} & \ldots & \tau s_{m} \\
\tau A s_{1} & \tau A s_{2} & \ldots & \tau A s_{m} \\
\vdots & \vdots & & \vdots \\
\tau A^{m-1} s_{1} & \tau A^{m-1} s_{2} & \ldots & \tau A^{m-1} s_{m}
\end{array}\right]=\left[\begin{array}{cccc}
1 & 0 & \ldots & 0 \\
0 & 1 & \ldots & 0 \\
\vdots & \vdots & & \vdots \\
0 & 0 & \ldots & 1
\end{array}\right]\right)
$$

establishes that the $i^{\text {th }}$ row of $T A T^{-1}$ is the $(i+1)^{t h}$ row of $I_{m} \quad(1=$ $1,2, \ldots, m-1)$, so $T A T^{-1}$ has the same form as $C$ in $(2.26)$, with last row given by

$$
k_{i}=-\tau A^{m} s_{m-i+1}, \quad i=1,2, \ldots, m .
$$

For

$$
\dot{w}=T A T^{-1} w+T b u
$$

to be identical to

$$
\dot{w}=C w+d u
$$

we must also have

$$
T b=d
$$

and substitution of $T$ into this relation gives

$$
\tau b=0, \quad \tau A b=0, \quad \ldots, \quad \tau A^{m-2} b=0, \quad \tau A^{m-1} b=1
$$

or

$$
\tau\left[\begin{array}{lllll}
b & A b & A^{2} b & \ldots & A^{m-1} b
\end{array}\right]=d^{T}
$$

which has a unique solution for $\tau$ (in view of condition (2.30)).
It remains to prove that (the matrix) $T$ is nonsingular; we shall show that its rows are linearly independent. Suppose that

$$
\alpha_{1} \tau+\alpha_{2} \tau A+\cdots+\alpha_{m} \tau A^{m-1}=0
$$

for some scalars $\alpha_{i}, i=1,2, \ldots, m$. Multiplying this relation on the right by $b$ gives $\alpha_{m}=0$. Similarly, multiplying on the right successively by $A b, A^{2} b, \ldots, A^{m-1} b$ gives $\alpha_{m-1}=0, \ldots, \alpha_{1}=0$. Thus, the rows of $T$ are linearly independent.
$(\Rightarrow)$ Necessity. Conversely, if such a transformation $T$ exists, then

$$
\begin{aligned}
\operatorname{rank}\left[\begin{array}{llll}
b & A b & \ldots & A^{m-1} b
\end{array}\right] & =\operatorname{rank}\left[\begin{array}{lllll}
T b & T A b & T A^{2} b & \ldots & T A^{m-1} b
\end{array}\right] \\
& =\operatorname{rank}\left[\begin{array}{lllll}
T b & \left(T A T^{-1}\right) T b & \ldots & \left(T A T^{-1}\right)^{m-1} T b
\end{array}\right] \\
& =\operatorname{rank}\left[\begin{array}{lllll}
d & C d & C^{2} d & \ldots & C^{m-1} d
\end{array}\right] .
\end{aligned}
$$

It is easy to verify that this last matrix has the triangular form

$$
\left[\begin{array}{cccccc}
0 & 0 & 0 & \ldots & 0 & 1 \\
0 & 0 & 0 & \ldots & 1 & \theta_{1} \\
\vdots & \vdots & \vdots & & \vdots & \vdots \\
0 & 0 & 1 & \ldots & \theta_{m-4} & \theta_{m-3} \\
0 & 1 & \theta_{1} & \ldots & \theta_{m-3} & \theta_{m-2} \\
1 & \theta_{1} & \theta_{2} & \ldots & \theta_{m-2} & \theta_{m-1}
\end{array}\right]
$$

and therefore has full rank. This completes the proof.

Note: $T$ can be constructed using

$$
T=\left[\begin{array}{c}
\tau \\
\tau A \\
\vdots \\
\tau A^{m-1}
\end{array}\right] \quad \text { and } \quad \tau\left[\begin{array}{lllll}
b & A b & A^{2} b & \ldots & A^{m-1} b
\end{array}\right]=d^{T} .
$$

However, we can also give an explicit expression for the matrix in the transformation $x=T^{-1} w$. We have seen that

$$
T\left[\begin{array}{lllll}
b & A b & A^{2} b & \ldots & A^{m-1} b
\end{array}\right]=\left[\begin{array}{cccccc}
0 & 0 & 0 & \ldots & 0 & 1 \\
0 & 0 & 0 & \ldots & 1 & \theta_{1} \\
\vdots & \vdots & \vdots & & \vdots & \vdots \\
0 & 0 & 1 & \ldots & \theta_{m-4} & \theta_{m-3} \\
0 & 1 & \theta_{1} & \ldots & \theta_{m-3} & \theta_{m-2} \\
1 & \theta_{1} & \theta_{2} & \ldots & \theta_{m-2} & \theta_{m-1}
\end{array}\right] .
$$

This latter matrix has elements given by

$$
\theta_{i}=-k_{1} \theta_{i-1}-k_{2} \theta_{i-2}-\cdots-k_{i}, \quad i=1,2, \ldots, m-1 .
$$

It is straightforward to verify that

$$
\left[\begin{array}{cccccc}
0 & 0 & 0 & \ldots & 0 & 1 \\
0 & 0 & 0 & \ldots & 1 & \theta_{1} \\
\vdots & \vdots & \vdots & & \vdots & \vdots \\
0 & 0 & 1 & \ldots & \theta_{m-4} & \theta_{m-3} \\
0 & 1 & \theta_{1} & \ldots & \theta_{m-3} & \theta_{m-2} \\
1 & \theta_{1} & \theta_{2} & \ldots & \theta_{m-2} & \theta_{m-1}
\end{array}\right]^{-1}=\left[\begin{array}{ccccc}
k_{m-1} & k_{m-2} & \ldots & k_{1} & 1 \\
k_{m-2} & k_{m-3} & \ldots & 1 & 0 \\
\vdots & \vdots & & \vdots & \vdots \\
k_{1} & 1 & \ldots & 0 & 0 \\
1 & 0 & \ldots & 0 & 0
\end{array}\right]
$$

Finally,

$$
T^{-1}=\left[\begin{array}{lllll}
b & A b & A^{2} b & \ldots & A^{m-1} b
\end{array}\right]\left[\begin{array}{ccccc}
k_{m-1} & k_{m-2} & \ldots & k_{1} & 1  \tag{2.31}\\
k_{m-2} & k_{m-3} & \ldots & 1 & 0 \\
\vdots & \vdots & & \vdots & \vdots \\
k_{1} & 1 & \ldots & 0 & 0 \\
1 & 0 & \ldots & 0 & 0
\end{array}\right]
$$

$\left(k_{1}, k_{2}, \ldots, k_{m-1}\right.$ are the coefficients in the characteristic equation of $C$.)
2.4.4 Example. Consider a system in the form

$$
\dot{x}=A x+b u(t)
$$

where

$$
A=\left[\begin{array}{rr}
1 & -3 \\
4 & 2
\end{array}\right] \quad \text { and } \quad b=\left[\begin{array}{l}
1 \\
1
\end{array}\right] .
$$

Find the matrix $T$ and the transformed system.
Solution : From

$$
\tau b=0, \quad \tau A b=1
$$

with $\tau=\left[\begin{array}{ll}\tau_{1} & \tau_{2}\end{array}\right]$ we have

$$
\left\{\begin{aligned}
\tau_{1}+\tau_{2} & =0 \\
-2 \tau_{1}+6 \tau_{2} & =1
\end{aligned}\right.
$$

whence

$$
\tau_{1}=-\frac{1}{8} \quad \text { and } \quad \tau_{2}=\frac{1}{8}
$$

Now

$$
T=\left[\begin{array}{c}
\tau \\
\tau A
\end{array}\right]=\frac{1}{8}\left[\begin{array}{rr}
-1 & 1 \\
3 & 5
\end{array}\right]
$$

and

$$
T^{-1}=\left[\begin{array}{rr}
-5 & 1 \\
3 & 1
\end{array}\right]
$$

Then

$$
T A T^{-1}=\frac{1}{8}\left[\begin{array}{rr}
-1 & 1 \\
3 & 5
\end{array}\right]\left[\begin{array}{rr}
1 & -3 \\
4 & 2
\end{array}\right]\left[\begin{array}{rr}
-5 & 1 \\
3 & 1
\end{array}\right]=\left[\begin{array}{rr}
0 & 1 \\
-14 & 3
\end{array}\right]
$$

Thus the transformed system is

$$
\dot{w}=T A T^{-1} w+T b u=\left[\begin{array}{rr}
0 & 1 \\
-14 & 3
\end{array}\right]\left[\begin{array}{l}
w_{1} \\
w_{2}
\end{array}\right]+\frac{1}{8}\left[\begin{array}{rr}
-1 & 1 \\
3 & 5
\end{array}\right]\left[\begin{array}{l}
1 \\
1
\end{array}\right] u
$$

that is,

$$
\left\{\begin{array}{l}
\dot{w}_{1}=w_{2} \\
\dot{w}_{2}=-14 w_{1}+3 w_{2}+u
\end{array}\right.
$$

or $\left(\right.$ for $\left.w_{1}=z, w_{2}=\dot{z}\right)$

$$
\ddot{z}-3 \dot{z}+14 z=u .
$$

A result similar to Theorem 2.4 .3 can be obtained for systems having zero input and scalar output, so that the system equations are

$$
\left\{\begin{array}{l}
\dot{x}=A x  \tag{2.32}\\
y=c x
\end{array}\right.
$$

where $A \in \mathbb{R}^{m \times m}, c \in \mathbb{R}^{1 \times m}$, and $y(\cdot)$ is the output variable.
2.4.5 Theorem. Any system described by (2.32) can be transformed by $x=S v$, with $S$ nonsingular, into the canonical form

$$
\begin{equation*}
\dot{v}=E v, \quad y=f v \tag{2.33}
\end{equation*}
$$

where

$$
E=\left[\begin{array}{cccccc}
0 & 0 & 0 & \ldots & 0 & -e_{m}  \tag{2.34}\\
1 & 0 & 0 & \ldots & 0 & -e_{m-1} \\
0 & 1 & 0 & \ldots & 0 & -e_{m-2} \\
\vdots & \vdots & \vdots & & \vdots & \vdots \\
0 & 0 & 0 & \ldots & 1 & -e_{1}
\end{array}\right] \text { and } f=\left[\begin{array}{cccc}
0 & 0 & \ldots & 1
\end{array}\right]
$$

provided that

$$
\operatorname{rank}\left[\begin{array}{c}
c  \tag{2.35}\\
c A \\
c A^{2} \\
\vdots \\
c A^{m-1}
\end{array}\right]=m
$$

Conversely, if such a transformation $S$ exists, then condition (2.35) holds.
The proof is very similar to that of Theorem 2.4.3 and will be omitted.
NOTE : $E$ is also a companion matrix because its characteristic polynomial is

$$
\operatorname{char}_{E}(\lambda)=\operatorname{det}\left(\lambda I_{m}-E\right)=\lambda^{m}+e_{1} \lambda^{m-1}+\cdots+e_{m}
$$

which again is identical to the characteristic polynomial of $A$.

### 2.5 Exercises

Exercise 26 Find the general solution, in spectral form, of the (initialized) uncontrolled system

$$
\dot{x}=A x, \quad x(0)=x_{0}
$$

in each of the following cases:
(a) $A=\left[\begin{array}{rr}-1 & -1 \\ 2 & -4\end{array}\right]$.
(b) $A=\left[\begin{array}{rr}2 & -1 \\ -1 & 2\end{array}\right]$.
(c) $A=\left[\begin{array}{ll}0 & 1 \\ 1 & 0\end{array}\right]$.
(d) $A=\left[\begin{array}{rrr}1 & 0 & -1 \\ 1 & 2 & 1 \\ 2 & 2 & 3\end{array}\right]$.
(e) $A=\left[\begin{array}{lll}2 & 1 & 1 \\ 0 & 0 & 1 \\ 0 & 1 & 0\end{array}\right]$.

Exercise 27 Find the general solution, in exponential form, of the (initialized) uncontrolled system

$$
\dot{x}=A x, \quad x(0)=x_{0}
$$

in each of the cases given in Exercise 26.

Exercise 28 Consider the equation of simple harmonic motion

$$
\ddot{z}+\omega^{2} z=0 .
$$

Take as state variables $x_{1}=z$ and $x_{2}=\frac{1}{\omega} \dot{z}$, and find the state transition matrix $\Phi(t, 0)$.

Exercise 29 Use the exponential matrix to solve the rabbit-fox environment problem

$$
\left[\begin{array}{c}
\dot{x}_{1} \\
\dot{x}_{2}
\end{array}\right]=\left[\begin{array}{cc}
a_{1} & -a_{2} \\
a_{3} & -a_{4}
\end{array}\right]\left[\begin{array}{c}
x_{1} \\
x_{2}
\end{array}\right] \quad\left(a_{1}, a_{2}, a_{3}, a_{4}>0\right)
$$

subject to the condition

$$
\frac{a_{1}}{a_{3}}=\frac{a_{2}}{a_{4}} .
$$

Show that for arbitrary initial conditions, the populations will attain a steady state as $t \rightarrow \infty$ only if $a_{1}-a_{4}<0$, and give an expression for the ultimate size of the rabbit population in this case. Finally, deduce that if the environment is to reach a steady state in which both rabbits and foxes are present, then $x_{1}(0)>\frac{a_{1}}{a_{3}} x_{2}(0)$.

Exercise 30 A linear control system is described by the equations

$$
\left\{\begin{array}{l}
\dot{x}_{1}=x_{1}+4 x_{2}+u(t) \\
\dot{x}_{2}=3 x_{1}+2 x_{2} .
\end{array}\right.
$$

Determine the state transition matrix and write down the general solution.

Exercise 31 A linear control system is described by

$$
\ddot{z}+3 \dot{z}+2 z=u, \quad z(0)=\dot{z}(0)=0
$$

where

$$
u(t)= \begin{cases}1 & \text { if } 0 \leq t<1 \\ 0 & \text { if } t \geq 1\end{cases}
$$

Calculate the state transition matrix and determine $z(2)$.

Exercise 32 Verify that the solution of the matrix differential equation

$$
\dot{W}=A W+W B, \quad W(0)=C
$$

(where $A, B \in \mathbb{R}^{m \times m}$ ) is

$$
W(t)=\exp (t A) C \exp (t B)
$$

Exercise 33 Consider the linear control system

$$
\dot{x}=A x+b u(t)
$$

where

$$
A=\left[\begin{array}{ll}
1 & 2 \\
0 & 1
\end{array}\right], \quad b=\left[\begin{array}{l}
2 \\
1
\end{array}\right]
$$

and take

$$
u(t)= \begin{cases}1 & \text { if } t \geq 0 \\ 0 & \text { if } t<0\end{cases}
$$

Evaluate $\exp (t A)$ and show that the solution of this problem, subject to

$$
x(0)=\left[\begin{array}{l}
1 \\
0
\end{array}\right]
$$

is

$$
x(t)=\left[\begin{array}{c}
1+3 t+\frac{5}{2} t^{2}+\frac{7}{6} t^{3}+\cdots \\
t+\frac{1}{2} t^{2}+\frac{1}{6} t^{3}+\cdots
\end{array}\right] .
$$

Exercise 34 Consider the matrix differential equation

$$
\dot{X}=A(t) X, \quad X(0)=I_{m}
$$

Show that, when $m=2$,

$$
\frac{d}{d t} \operatorname{det}(X(t))=\operatorname{tr}(A(t)) \cdot \operatorname{det}(X(t))
$$

and hence deduce that $X(t)$ is nonsingular, $t \geq 0$.

Exercise 35 Verify that the properties of the state transition matrix
(a) $\frac{d}{d t} \Phi\left(t, t_{0}\right)=A \Phi\left(t, t_{0}\right)$.
(b) $\Phi\left(t_{0}, t_{0}\right)=I_{m}$.
(c) $\Phi\left(t_{0}, t\right)=\Phi^{-1}\left(t, t_{0}\right)$.
(d) $\Phi\left(t, t_{0}\right)=\Phi\left(t, t_{1}\right) \Phi\left(t_{1}, t_{0}\right)$
do carry over to the time varying case.
Exercise 36 Consider the (initialized) uncontrolled system

$$
\dot{x}=A(t) x, \quad x(0)=x_{0} .
$$

If $B(t)=\int_{0}^{t} A(\tau) d \tau$, show that the solution in this case is

$$
x(t)=\exp (B(t)) x_{0}
$$

provided $B(t)$ and $A(t)$ commute with each other (for all $t \geq 0$ ).

Exercise 37 Verify that the solution of the matrix differential equation

$$
\dot{W}=A(t) W+W A^{T}(t), \quad W\left(t_{0}\right)=C
$$

is

$$
W(t)=\Phi\left(t, t_{0}\right) C \Phi^{T}\left(t, t_{0}\right)
$$

Exercise 38 For the linear control system

$$
\dot{x}=\left[\begin{array}{ll}
-1 & -4 \\
-1 & -1
\end{array}\right] x+\left[\begin{array}{l}
1 \\
1
\end{array}\right] u(t)
$$

determine $\Phi(t, 0)$. If

$$
x(0)=\left[\begin{array}{l}
1 \\
2
\end{array}\right] \quad \text { and } \quad u(t)=e^{2 t}, t \geq 0
$$

use formula

$$
x(t)=\Phi(t, 0)\left[x_{0}+\int_{0}^{t} \Phi(0, \tau) B u(\tau) d \tau\right]
$$

to obtain the expression for $x(t)$.

Exercise 39 Consider a single-input control system, written in state space form,

$$
\dot{x}=A x+b u(t) .
$$

Find the matrix $T$ of the linear tansformation $w=T x$ and the transformed system (the system put into canonical form)

$$
\dot{w}=C w+d u(t)
$$

for each of the following cases :
(a) $A=\left[\begin{array}{rr}-1 & -1 \\ 2 & -4\end{array}\right], \quad b=\left[\begin{array}{l}1 \\ 3\end{array}\right]$.
(b) $A=\left[\begin{array}{rrr}1 & 0 & -1 \\ 1 & 2 & 1 \\ 2 & 2 & 3\end{array}\right], \quad b=\left[\begin{array}{l}1 \\ 0 \\ 1\end{array}\right]$.

Exercise 40 Consider a single-output (uncontrolled) system, written in state space form,

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

Find the matrix $P$ of the linear transformation $x=P v$ and the transformed system (the system put into the canonical form)

$$
\dot{v}=E v, \quad y=f v
$$

when

$$
A=\left[\begin{array}{rrr}
1 & 2 & 0 \\
3 & -1 & 1 \\
0 & 2 & 0
\end{array}\right], \quad c=\left[\begin{array}{lll}
0 & 0 & 2
\end{array}\right]
$$

