C24 Dynamical Systems Notes based on Ron Daniel's lectures and lecture notes in MT 2014

1 Linear algebra

Definition 1.1. Let **A** be an $n \times n$ matrix. Then

$$e^{\mathbf{A}} = \sum_{k=0}^{\infty} \frac{\mathbf{A}^k}{k!} \tag{1}$$

Proposition 1.1. If **S** and **T** are linear transformations on \mathbb{R}^n which commute, i.e. $\mathbf{ST} = \mathbf{TS}$, then $e^{\mathbf{S}+\mathbf{T}} = e^{\mathbf{ST}}$.

Proof. Not interesting. \Box

Proposition 1.2. If \mathbf{P}, \mathbf{T} are linear transformations on \mathbb{R}^n (i.e. $\mathbf{P}, \mathbf{T} \in \mathbb{R}^{n \times n}$) and $\mathbf{S} = \mathbf{P}\mathbf{T}\mathbf{P}^{-1}$ then $e^{\mathbf{S}} = \mathbf{P}e^{\mathbf{T}}\mathbf{P}^{-1}$.

Proof.

$$e^{\mathbf{S}} = \lim_{n \to \infty} \sum_{k=0}^{n} \frac{(\mathbf{PTP}^{-1})^{k}}{k!}$$

$$= \lim_{n \to \infty} \sum_{k=0}^{n} \frac{\mathbf{PT}^{k} \mathbf{P}^{-1}}{k!}$$

$$= \mathbf{P} \left(\lim_{n \to \infty} \sum_{k=0}^{n} \frac{T^{k}}{k!} \right) \mathbf{P}^{-1}$$

$$= \mathbf{P} e^{\mathbf{T}} \mathbf{P}^{-1}$$

Proposition 1.3. If

$$\mathbf{A} = \begin{bmatrix} a & 0 \\ 0 & b \end{bmatrix}$$

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then

$$e^{\mathbf{A}} = \begin{bmatrix} e^a & 0\\ 0 & e^b \end{bmatrix}$$

Proof. By induction

$$\mathbf{A}^k = \begin{bmatrix} a^k & 0 \\ 0 & b^k \end{bmatrix}$$

for $k = 0, 1, \ldots$ Then we can write

$$e^{\mathbf{A}} = \sum_{k=0}^{\infty} \frac{\mathbf{A}^k}{k!}$$

$$= \sum_{k=0}^{\infty} \frac{1}{k!} \begin{bmatrix} a^k & 0\\ 0 & b^k \end{bmatrix}$$

$$= \sum_{k=0}^{\infty} \begin{bmatrix} a^k/k! & 0\\ 0 & b^k/k! \end{bmatrix}$$

$$= \begin{bmatrix} e^a & 0\\ 0 & e^b \end{bmatrix}$$

Proposition 1.4. If

$$\mathbf{A} = \begin{bmatrix} a & b \\ 0 & a \end{bmatrix}$$

then

$$e^{\mathbf{A}} = e^a \begin{bmatrix} 1 & b \\ 0 & 1 \end{bmatrix}$$

Proof. Write $\mathbf{A} = a\mathbf{I} + \mathbf{B}$ where

$$\mathbf{B} = \begin{bmatrix} 0 & b \\ 0 & 0 \end{bmatrix}$$

Then $a\mathbf{I}$ commutes with \mathbf{B} and by Proposition 1.1,

$$e^{\mathbf{A}} = e^{a\mathbf{I}}e^{\mathbf{B}} = e^a e^{\mathbf{B}}$$

And from the definition

$$e^{\mathbf{B}} = \mathbf{I} + \mathbf{B} + \mathbf{B}^2/2! + \dots = \mathbf{I} + \mathbf{B}$$

since by direct computation $\mathbf{B}^2 = \mathbf{B}^3 = \cdots = 0$.

Proposition 1.5. If

$$\mathbf{A} = \begin{bmatrix} a & -b \\ b & a \end{bmatrix}$$

then

$$e^{\mathbf{A}} = e^a \begin{bmatrix} \cos b & -\sin b \\ \sin b & \cos b \end{bmatrix}$$

Proof. If $\lambda = a + ib$, it follows by induction that

$$\begin{bmatrix} a & -b \\ b & a \end{bmatrix}^k = \begin{bmatrix} \operatorname{Re}(\lambda^k) & -\operatorname{Im}(\lambda^k) \\ \operatorname{Im}(\lambda^k) & \operatorname{Re}(\lambda^k) \end{bmatrix}$$

explicitly: it is true for k = 0 and assuming true for k - 1, we can write

$$\lambda^{k-1} = a_{k-1} + ib_{k-1}$$

$$\lambda^k = \lambda^{k-1}\lambda$$

$$= (a_{k-1} + ib_{k-1})(a+ib)$$

$$= (a_{k-1}a - b_{k-1}b) + i(b_{k-1}a + a_{k-1}b)$$

and so

$$\begin{bmatrix} a & -b \\ b & a \end{bmatrix}^k = \begin{bmatrix} \operatorname{Re}(\lambda^{k-1}) & -\operatorname{Im}(\lambda^{k-1}) \\ \operatorname{Im}(\lambda^{k-1}) & \operatorname{Re}(\lambda^{k-1}) \end{bmatrix} \begin{bmatrix} a & -b \\ b & a \end{bmatrix}$$
$$= \begin{bmatrix} a_{k-1} & -b_{k-1} \\ b_{k-1} & a_{k-1} \end{bmatrix} \begin{bmatrix} a & -b \\ b & a \end{bmatrix}$$
$$= \begin{bmatrix} (a_{k-1}a - b_{k-1}b) & (-a_{k-1}b - b_{k-1}a) \\ (b_{k-1}a + a_{k-1}b) & (-b_{k-1}b + a_{k-1}a) \end{bmatrix}$$
$$= \begin{bmatrix} \operatorname{Re}(\lambda^k) & -\operatorname{Im}(\lambda^k) \\ \operatorname{Im}(\lambda^k) & \operatorname{Re}(\lambda^k) \end{bmatrix}$$

Using this

$$e^{\mathbf{A}} = \sum_{k=0}^{\infty} \begin{bmatrix} \operatorname{Re}(\frac{\lambda_k}{k!}) & -\operatorname{Im}(\frac{\lambda_k}{k!}) \\ \operatorname{Im}(\frac{\lambda_k}{k!}) & \operatorname{Re}(\frac{\lambda_k}{k!}) \end{bmatrix}$$
$$= \begin{bmatrix} \operatorname{Re}(e^{\lambda}) & -\operatorname{Im}(e^{\lambda}) \\ \operatorname{Im}(e^{\lambda}) & \operatorname{Re}(e^{\lambda}) \end{bmatrix}$$
$$= e^{a} \begin{bmatrix} \cos b & -\sin b \\ \sin b & \cos b \end{bmatrix}$$

Theorem 1.1 (The Jordan Canonical Form). Let $\mathbf{A} \in \mathbb{R}^{(2n-k)\times(2n-k)}$ be a real matrix with

- real eigenvalues $\lambda_j, j = 1, \dots, k$ and
- complex eigenvalues $\lambda_j = a_j + ib_j$ and $\bar{\lambda}_j = a_j ib_j, j = k + 1, \dots, n$.

Then there exists a basis $\{\mathbf{v}_1,\ldots,\mathbf{v}_k,\mathbf{v}_{k+1},\mathbf{u}_{k+1},\ldots,\mathbf{v}_n,\mathbf{u}_n\}$ for \mathbb{R}^{2n-k} , where

• $\mathbf{v}_{i}, j = 1, \dots, k \text{ and }$

•
$$\mathbf{u}_i + i\mathbf{v}_i, j = k+1, \ldots, n$$

are generalised eigenvectors of **A** such that the matrix

$$\mathbf{W} = [\mathbf{v}_1, \dots, \mathbf{v}_k, \mathbf{v}_{k+1}, \mathbf{u}_{k+1}, \dots, \mathbf{v}_n, \mathbf{u}_n]$$

is invertible and

$$\mathbf{\Lambda} = \mathbf{W}^{-1} \mathbf{A} \mathbf{W} = \begin{bmatrix} \mathbf{J}_1 & & \\ & \ddots & \\ & & \mathbf{J}_r \end{bmatrix}$$
 (2)

where the elementary Jordan blocks $\mathbf{J} = \mathbf{J}_i, j = 1, \dots, r$ are either of the form

$$\mathbf{J} = \begin{bmatrix} \lambda & 1 & 0 & \cdots & 0 \\ 0 & \lambda & 1 & \cdots & 0 \\ \cdots & & & & \\ 0 & \cdots & \lambda & 1 \\ 0 & \cdots & 0 & \lambda \end{bmatrix}$$
(3)

for one of the real eigenvalues of \mathbf{A} , λ , or of the form

$$\mathbf{J} = \begin{bmatrix} \mathbf{D} & \mathbf{I} & \mathbf{0} & \cdots & \mathbf{0} \\ \mathbf{0} & \mathbf{D} & \mathbf{I} & \cdots & \mathbf{0} \\ \cdots & & & & \\ \mathbf{0} & \cdots & \mathbf{D} & \mathbf{I} \\ \mathbf{0} & \cdots & \mathbf{0} & \mathbf{D} \end{bmatrix}$$
(4)

with

$$\mathbf{D} = \begin{bmatrix} a & -b \\ b & a \end{bmatrix}, \mathbf{I} = \begin{bmatrix} 1 & 0 \\ 0 & 1 \end{bmatrix}, \mathbf{0} = \begin{bmatrix} 0 & 0 \\ 0 & 0 \end{bmatrix}$$

for one of the complex eigenvalues of \mathbf{A} , $\lambda = a + ib$.

Corollary 1.1. For any $\mathbf{A} \in \mathbb{R}^{2\times 2}$, there exists an invertible matrix $\mathbf{W} \in \mathbb{R}^{2\times 2}$ (described in the proof) such that the matrix

$$\Lambda = \mathbf{W}^{-1} \mathbf{A} \mathbf{W}$$

has one of the following forms

$$\mathbf{\Lambda} = \begin{bmatrix} \lambda & 0 \\ 0 & \mu \end{bmatrix}, \mathbf{\Lambda} = \begin{bmatrix} \lambda & 1 \\ 0 & \lambda \end{bmatrix}, or \mathbf{\Lambda} = \begin{bmatrix} a & -b \\ b & a \end{bmatrix}$$

It then follows from Propositions 1.3, 1.4 and 1.5 that

$$e^{\mathbf{\Lambda}t} = \begin{bmatrix} e^{\lambda t} & 0 \\ 0 & e^{\mu t} \end{bmatrix}, e^{\mathbf{\Lambda}t} = e^{\lambda t} \begin{bmatrix} 1 & t \\ 0 & 1 \end{bmatrix}, or \ e^{\mathbf{\Lambda}t} = e^{at} \begin{bmatrix} \cos bt & -\sin bt \\ \sin bt & \cos bt \end{bmatrix}$$

respectively. And by Proposition 1.2, the matrix $e^{\mathbf{A}t}$ is then given by

$$e^{\mathbf{A}t} = \mathbf{W}e^{\mathbf{\Lambda}t}\mathbf{W}^{-1}$$

Proof. We analyse the Jordan Canonical Form of three possible cases of eigendecomposition of $\mathbf{A} \in \mathbb{R}^{2 \times 2}$ (characteristic equation $|\mathbf{A} - \lambda \mathbf{I}|$ a quadratic equation in λ):

1. Eigenvalues are real and distinct: λ, μ . Eigenvectors are $\mathbf{w}_1, \mathbf{w}_2$. In this case n = k = 2. For the matrix $\mathbf{W} = [\mathbf{w}_1, \mathbf{w}_2]$:

$$\mathbf{W}^{-1}\mathbf{A}\mathbf{W} = \begin{bmatrix} J_1 & 0\\ 0 & J_2 \end{bmatrix}$$

where $J_1 = \lambda, J_2 = \mu$.

2. Eigenvalues are real and equal: λ . There is only one eigenvector, \mathbf{w} . In this case n = k = 2. For the matrix $\mathbf{W} = [\mathbf{w}, \mathbf{w}]$:

$$\mathbf{W}^{-1}\mathbf{A}\mathbf{W} = [\mathbf{J}_1]$$

where

$$\mathbf{J}_1 = \begin{bmatrix} \lambda & 1 \\ 0 & \lambda \end{bmatrix}$$

3. Eigenvalues are complex conjugates: $\lambda_1 = a + ib$, $\lambda_2 = a - ib$. Eigenvectors are $\mathbf{w}_1 = \mathbf{u} + i\mathbf{v}$, $\mathbf{w}_2 = \mathbf{u} - i\mathbf{v}$.

In this case k = 0, n = 1. For the matrix $\mathbf{W} = [\mathbf{v}, \mathbf{u}]$:

$$\mathbf{W}^{-1}\mathbf{A}\mathbf{W} = [\mathbf{D}]$$

where

$$\mathbf{D} = \begin{bmatrix} a & -b \\ b & a \end{bmatrix}$$

2 Introduction

We have: state $\mathbf{x}(t) \in \mathbb{R}^n$, function $\mathbf{f}: D \to \mathbb{R}^n, D \subseteq \mathbb{R}^n$. Then a dynamical system can be written as

$$\dot{\mathbf{x}} = \mathbf{f}(\mathbf{x}) \tag{5}$$

In this lecture course, we assume that the solution exists and is unique locally, i.e. we will assume that the vector fields are sufficiently smooth to allow this.

We call $\dot{\mathbf{x}} = \mathbf{A}\mathbf{x}$ a linear autonomous system.

Theorem 2.1 (The Fundamental Theorem for Linear Systems). Let **A** be an $n \times n$ matrix. Then for a given $\mathbf{x}_0 \in \mathbb{R}^n$, the initial value problem

$$\dot{\mathbf{x}} = \mathbf{A}\mathbf{x} \tag{6}$$

$$\mathbf{x}(0) = \mathbf{x}_0 \tag{7}$$

has a unique solution given by

$$\mathbf{x}(t) = e^{\mathbf{A}t}\mathbf{x}_0 \tag{8}$$

3 Equilibria and stability

Equilibrium solution of $\dot{\mathbf{x}} = \mathbf{f}(\mathbf{x})$ is a solution $\mathbf{x}^* \in \mathbb{R}^n$ which is constant, i.e. $\mathbf{f}(\mathbf{x}^*) = \mathbf{0}$. \mathbf{x}^* is called a fixed point, stationary point, rest point, critical point or steady state.

For maps, $\mathbf{x}_{k+1} = \mathbf{F}(\mathbf{x}_k)$, the equilibrium solution is a solution $\mathbf{x}^* \in \mathbb{R}^n$ such that $\mathbf{x}^* = \mathbf{F}(\mathbf{x}^*)$.

Definition 3.1. An equilibrium point \mathbf{x}^* is said to be stable if, given $\epsilon > 0$, there exists a real number $\delta > 0$ such that for any other solution $\mathbf{y}(t)$ satisfying $|\mathbf{y}(0) - \mathbf{x}^*| < \delta$, $|\mathbf{y}(t) - \mathbf{x}^*| < \epsilon$ for all $t \geq 0$. Otherwise it is called unstable.

Definition 3.2. An equilibrium point \mathbf{x}^* is said to be asymptotically stable if it is stable and there is a b > 0 such that if $|\mathbf{y}(0) - \mathbf{x}^*| < b$ then $\lim_{t \to \infty} |\mathbf{y}(t) - \mathbf{x}^*| = 0$.

Definition 3.3. An equilibrium point \mathbf{x}^* is said to be exponentially stable if it is asymptotically stable and there exist finite $\alpha, \beta, \delta > 0$ such that if $|\mathbf{y}(0) - \mathbf{x}^*| < \delta$ then $|\mathbf{y}(t) - \mathbf{x}^*| \le \alpha e^{-\beta t} |\mathbf{y}(0) - \mathbf{x}^*|$ for $t \ge 0$.

3.1 Analysing stability in 2x2 systems

In a linear system

$$\dot{\mathbf{x}} = \mathbf{A}\mathbf{x}$$
$$\mathbf{x}(0) = \mathbf{x}_0$$

where $\mathbf{x} \in \mathbb{R}^2$ where we can decompose $\mathbf{A} = \mathbf{W} \mathbf{\Lambda} \mathbf{W}^{-1}$ according to Corollary 1.1, we can use the coordinate transformation $\mathbf{z} = \mathbf{W}^{-1} \mathbf{x}$ and $\mathbf{z}_0 = \mathbf{W}^{-1} \mathbf{x}_0$ to analyse the system in the transformed coordinates \mathbf{z} :

$$\dot{\mathbf{z}} = \mathbf{\Lambda}\mathbf{z}$$
$$\mathbf{z}(0) = \mathbf{z}_0$$

hence according to Theorem 2.1:

$$\mathbf{z}(t) = e^{\mathbf{\Lambda}t} \mathbf{z}_0$$

We now analyse the three possible cases of Λ :

Real, distinct eigenvalues. We have for $\lambda, \mu, \lambda \neq \mu$:

$$\mathbf{\Lambda} = \begin{bmatrix} \lambda & 0 \\ 0 & \mu \end{bmatrix}$$

and

$$e^{\mathbf{\Lambda}t} = \begin{bmatrix} e^{\lambda t} & 0\\ 0 & e^{\mu t} \end{bmatrix}$$

hence

$$\mathbf{z}(t) = \begin{bmatrix} e^{\lambda t} & 0\\ 0 & e^{\mu t} \end{bmatrix} \mathbf{z}_0$$
$$= z_{1,0} e^{\lambda t} + z_{2,0} e^{\mu t}$$

There can be three subcases:

- $\lambda < \mu < 0$: stable.
- $\lambda < 0 < \mu$: saddle shape (separatrices).
- $0 < \lambda < \mu$: unstable (opposite arrows to stable case).

Duplicate real eigenvalue. We have for $\lambda \neq 0$:

$$\mathbf{\Lambda} = \begin{bmatrix} \lambda & 1 \\ 0 & \lambda \end{bmatrix}$$

and

$$e^{\mathbf{\Lambda}t} = e^{\lambda t} \begin{bmatrix} 1 & t \\ 0 & 1 \end{bmatrix}$$

hence

$$\mathbf{z}(t) = e^{\lambda t} \begin{bmatrix} 1 & t \\ 0 & 1 \end{bmatrix} \mathbf{z}_0$$

There can be two subcases:

- $\lambda < 0$: stable.
- $\lambda > 0$: unstable.

Complex eigenvalues. We have for $\lambda_1 = a + ib$, $\lambda_2 = a - ib$:

$$\mathbf{\Lambda} = \begin{bmatrix} a & -b \\ b & a \end{bmatrix}$$

and

$$e^{\mathbf{\Lambda}t} = e^{at} \begin{bmatrix} \cos bt & -\sin bt \\ \sin bt & \cos bt \end{bmatrix}$$

hence

$$\mathbf{z}(t) = e^{at} \begin{bmatrix} \cos bt & -\sin bt \\ \sin bt & \cos bt \end{bmatrix} \mathbf{z}_0$$

We can look at the norm of $\mathbf{z}(t)$ to analyse the shape of the solution:

$$\begin{aligned} \|\mathbf{z}(t)\|^2 &= |e^{at}|^2 \left\| \begin{bmatrix} z_{1,0}\cos bt - z_{2,0}\sin bt \\ z_{1,0}\sin bt + z_{2,0}\cos bt \end{bmatrix} \right\|^2 \\ &= |e^{at}|^2 (z_{1,0}^2\cos^2 bt - 2z_{1,0}z_{2,0}\sin bt\cos bt + z_{2,0}^2\sin^2 bt + z_{1,0}^2\sin^2 bt + 2z_{1,0}z_{2,0}\sin bt\cos bt + z_{2,0}^2\cos^2 bt) \\ &= |e^{at}|^2 (z_{1,0}^2(\sin^2 bt + \cos^2 bt) + z_{2,0}^2(\sin^2 bt + \cos^2 bt)) \\ &= |e^{at}|^2 (z_{1,0}^2 + z_{2,0}^2) \\ &= |e^{at}|^2 \|\mathbf{z}_0\|^2 \end{aligned}$$

i.e. $\|\mathbf{z}(t)\| = |e^{at}| \|\mathbf{z}_0\|$. So

- a=0: circle.
- a < 0: stable spiral.
- a > 0: unstable spiral.
- b > 0: anticlockwise.
- b < 0: clockwise.

3.2 Nonlinear systems

Consider $\dot{\mathbf{x}} = \mathbf{f}(\mathbf{x})$ with an equilibrium \mathbf{x}^* , such that $\mathbf{f}(\mathbf{x}^*) = \mathbf{0}$. Let $\mathbf{x} = \mathbf{x}^* + \mathbf{w}$ and assume \mathbf{f} is differentiable and not that the equilibrium point doesn't change with time. Then for $j = 1, \ldots, n$:

$$f_j(\mathbf{x} + \mathbf{w}) = f_j(\mathbf{x}^*) + (\operatorname{grad}_{\mathbf{x}} f_j(\mathbf{x}^*))^T \mathbf{w} + \cdots$$

Hence we can write

$$\dot{\mathbf{x}} = \dot{\mathbf{w}} = D\{\mathbf{f}(\mathbf{x}^*)\}\mathbf{w} + \cdots$$

where

$$D\{\mathbf{f}(\mathbf{x}^*)\} = \begin{bmatrix} (\operatorname{grad}_{\mathbf{x}} f_1(\mathbf{x}^*))^T \\ \vdots \\ (\operatorname{grad}_{\mathbf{x}} f_n(\mathbf{x}^*))^T \end{bmatrix} \in \mathbb{R}^{n \times n}$$

Hence we can analyse the following small perturbation model

$$\dot{\mathbf{w}} \approx D\{\mathbf{f}(\mathbf{x}^*)\}\mathbf{w} = \mathbf{A}\mathbf{w}$$

Definition 3.4. Let \mathbf{x}^* be an equilibrium of $\dot{\mathbf{x}} = \mathbf{f}(\mathbf{x})$. Then \mathbf{x}^* is called a hyperbolic fixed point if none of the eigenvalues of $D\{\mathbf{f}(\mathbf{x}^*)\}$ have zero real part.

Theorem 3.1. Suppose \mathbf{x}^* is a hyperbolic fixed point and all the real parts of the eigenvalues are negative. Then the equilibrium solution $\mathbf{x} = \mathbf{x}^*$ is asymptotically stable.

4 Invariant Manifolds

4.1 Linear systems: Stable, Unstable and Centre subspaces

In an $n \times n$ linear system

$$\dot{\mathbf{x}} = \mathbf{A}\mathbf{x}$$

if $\mathbf{A} \in \mathbb{R}^n$, n = 2m - k is eigendecomposed to

- real eigenvalues $\lambda_j = a_j$ with corresponding eigenvectors $\mathbf{w}_j = \mathbf{u}_j, \mathbf{v}_j = \mathbf{0}$ for $j = 1, \dots, k$, and
- complex eigenvalues $\lambda_j = a_j + ib_j$ and $\bar{\lambda}_j = a_j ib_j$ with corresponding eigenvectors $\mathbf{w}_j = \mathbf{u}_j + i\mathbf{v}_j$ and $\mathbf{w}_j = \mathbf{u}_j i\mathbf{v}_j$ for $j = k+1, \ldots, m$

then

$$\mathbf{W} = [\mathbf{u}_1, \dots, \mathbf{u}_k, \mathbf{v}_{k+1}, \mathbf{u}_{k+1}, \dots, \mathbf{v}_n, \mathbf{u}_n]$$
 $\mathbf{\Lambda} = \mathbf{W}^{-1} \mathbf{A} \mathbf{W}$

$$= \begin{bmatrix} \mathbf{J}_1 & & & \\ & \ddots & & \\ & & \mathbf{J}_r \end{bmatrix}$$

by Theorem 1.1.

Definition 4.1. E^s , E^c and E^u are the subspaces of \mathbb{R}^n spanned by the real and imaginary parts of the generalised eigenvectors corresponding to eigenvalues with negative, zero and positive real parts respectively:

$$E^{s} = \operatorname{span}\{\mathbf{u}_{j}, \mathbf{v}_{j} : a_{j} < 0\}$$

$$E^{u} = \operatorname{span}\{\mathbf{u}_{j}, \mathbf{v}_{j} : a_{j} > 0\}$$

$$E^{c} = \operatorname{span}\{\mathbf{u}_{j}, \mathbf{v}_{j} : a_{j} = 0\}$$

Definition 4.2. If all eigenvalues of the $n \times n$ matrix \mathbf{A} have nonzero real part, then the flow $e^{\mathbf{A}t} : \mathbb{R}^n \to \mathbb{R}^n$ (mapping \mathbf{x}_0 to $\mathbf{x}(t)$) is called a hyperbolic flow and $\dot{\mathbf{x}} = \mathbf{A}\mathbf{x}$ is called a hyperbolic linear system.

Definition 4.3. A subspace $E \subseteq \mathbb{R}^n$ is said to be invariant with respect to the flow $e^{\mathbf{A}t}$ if for any $\mathbf{x} \in E$: $e^{\mathbf{A}t}\mathbf{x} \in E$, $\forall t$ (in other words, E is closed under the linear transformation $e^{\mathbf{A}t}$ for all t; or: what starts in E stays in E).

Theorem 4.1. Let **A** be a real $n \times n$ matrix. Then

$$\mathbb{R}^n = E^s \oplus E^u \oplus E^c$$

where E^s , E^u and E^c are the stable, unstable and center subspaces of $\dot{\mathbf{x}} = \mathbf{A}\mathbf{x}$ respectively; furthermore, E^s , E^u and E^c are invariant with respect to the flow $e^{\mathbf{A}t}$ of $\dot{\mathbf{x}} = \mathbf{A}\mathbf{x}$ respectively.

Proof. (fake (Perko p55 for the real deal)) Let

$$\{\mathbf{w}_{i}^{s}, i = 1, \dots, n_{s}\} = \{\mathbf{u}_{j}, \mathbf{v}_{j} : a_{j} < 0, \mathbf{u}_{j}, \mathbf{v}_{j} \neq \mathbf{0}\}$$

$$\{\mathbf{w}_{i}^{u}, i = 1, \dots, n_{u}\} = \{\mathbf{u}_{j}, \mathbf{v}_{j} : a_{j} > 0, \mathbf{u}_{j}, \mathbf{v}_{j} \neq \mathbf{0}\}$$

$$\{\mathbf{w}_{i}^{c}, i = 1, \dots, n_{c}\} = \{\mathbf{u}_{j}, \mathbf{v}_{j} : a_{j} = 0, \mathbf{u}_{j}, \mathbf{v}_{j} \neq \mathbf{0}\}$$

then we can rearrange (questionable) W and Λ in the following way

$$\begin{split} \mathbf{W} &= [\mathbf{w}_1^s, \dots, \mathbf{w}_{n_s}^s, \mathbf{w}_1^u, \dots, \mathbf{w}_{n_u}^u, \mathbf{w}_1^c, \dots, \mathbf{w}_{n_c}^c] \\ \mathbf{\Lambda} &= \begin{bmatrix} \mathbf{J}_1^s & & & & & \\ & \mathbf{J}_r^s & & & & \\ & & \mathbf{J}_1^u & & & \\ & & & \mathbf{J}_1^c & & & \\ & & & & \mathbf{J}_1^c & & \\ & & & & \mathbf{J}_r^c \end{bmatrix} \\ &= \begin{bmatrix} \mathbf{\Lambda}^s & & & & & \\ & \mathbf{\Lambda}^u & & & \\ & & \mathbf{\Lambda}^c \end{bmatrix} \end{split}$$

where $\mathbf{\Lambda}^s \in \mathbb{R}^{n_s \times n_s}$, $\mathbf{\Lambda}^u \in \mathbb{R}^{n_u \times n_u}$ and $\mathbf{\Lambda}^c \in \mathbb{R}^{n_c \times n_c}$ such that $\mathbf{\Lambda} = \mathbf{W}^{-1} \mathbf{A} \mathbf{W}$ still holds.

Using the coordinate transformation $\mathbf{z} = \mathbf{W}^{-1}\mathbf{x}$, we get $\dot{\mathbf{z}} = \mathbf{\Lambda}\mathbf{z}$ hence

$$egin{aligned} \dot{\mathbf{z}}^s &= \mathbf{\Lambda}^s \mathbf{z}^s \ \dot{\mathbf{z}}^u &= \mathbf{\Lambda}^u \mathbf{z}^u \ \dot{\mathbf{z}}^c &= \mathbf{\Lambda}^c \mathbf{z}^c \end{aligned}$$

where

$$\begin{aligned} \mathbf{z}^s &= [z_1, \dots, z_{n_s}]^T \\ &= [z_1^s, \dots, z_{n_s}^s]^T \\ \mathbf{z}^u &= [z_{n_s+1}, \dots, z_{n_s+n_u}]^T \\ &= [z_1^u, \dots, z_{n_u}^u]^T \\ \mathbf{z}^c &= [z_{n_s+n_u+1}, \dots, z_{n_s+n_u+n_c}]^T \\ &= [z_1^c, \dots, z_{n_c}^c]^T \end{aligned}$$

Fix $\mathbf{x}_0 \in E^s$. We can then write (questionable)

$$\begin{split} \mathbf{x}_0 &= \mathbf{W} \mathbf{z}(0) \\ &= \sum_{i=1}^{n_s} \mathbf{w}_i^s z_i^s(0) + \sum_{j=1}^{n_u} \mathbf{w}_j^u z_j^u(0) + \sum_{k=1}^{n_c} \mathbf{w}_k^c z_k^c(0) \\ &= \sum_{i=1}^{n_s} \mathbf{w}_i^s z_i^s(0) + \mathbf{0} + \mathbf{0} \end{split}$$

hence (questionable) $\mathbf{z}^{u}(0) = \mathbf{z}^{c}(0) = \mathbf{0}$. The solution in $\mathbf{z}(t)$ then becomes

$$\mathbf{z}^{s}(t) = e^{\mathbf{\Lambda}^{s} t} \mathbf{z}^{s}(0)$$
$$\mathbf{z}^{u}(t) = \mathbf{0}$$
$$\mathbf{z}^{c}(t) = \mathbf{0}$$

which corresponds to

$$\begin{split} \mathbf{x}(t) &= \mathbf{W} \mathbf{z}(t) \\ &= \sum_{i=1}^{n_s} \mathbf{w}_i^s z_i^s(t) + \sum_{j=1}^{n_u} \mathbf{w}_j^u z_j^u(t) + \sum_{k=1}^{n_c} \mathbf{w}_k^c z_k^c(t) \\ &= \sum_{i=1}^{n_s} \mathbf{w}_i^s z_i^s(t) + \mathbf{0} + \mathbf{0} \end{split}$$

which is $\in E^s$. Hence E^s is with respect to the flow $e^{\mathbf{A}t}$. Similarly for E^u and E^c .

Theorem 4.2. If $\mathbf{x}_0 \in E^s$, then $e^{\mathbf{A}t}\mathbf{x}_0 \in E^s$ for all $t \in \mathbb{R}$ and

$$\lim_{t \to \infty} e^{\mathbf{A}t} \mathbf{x}_0 = \mathbf{0}.$$

And if $\mathbf{x}_0 \in E^u$, then $e^{\mathbf{A}t}\mathbf{x}_0 \in E^u$ for all $t \in \mathbb{R}$ and

$$\lim_{t \to -\infty} e^{\mathbf{A}t} \mathbf{x}_0 = \mathbf{0}.$$

4.2 Nonlinear system local theory

Consider the linearised model

$$\dot{\mathbf{w}} \approx D\{\mathbf{f}(\mathbf{x}^*)\}\mathbf{w} = \mathbf{A}\mathbf{w}$$

Theorem 4.3 (Hartman-Grobman). \forall hyperbolic equilibrium points $\mathbf{x}^* \in A$ (A is an open set)

 \exists bi-continuous (a mapping that is continous and whose inverse is also continuous) function $H: A \to B$ where

B is the open set containing the origin of the linearised model

so that trajectories are mapped exactly and the parameterisation of time is preserved.

4.3 Conservative systems

If there exists a non-constant function $V(\mathbf{x})$ such that dV/dt = 0 along a solution of the nonlinear differential equation $\dot{\mathbf{x}} = \mathbf{f}(\mathbf{x})$, the equations are called conservative.

If $\mathbf{x} = \mathbf{x}^*$ is an isolated equlibrium and there is a $V(\mathbf{x})$ that is a local minimum/maximum at \mathbf{x}^* then there is a region around that point that contains a closed orbit.

5 Lyapunov Functions

Theorem 5.1. Consider a nonlinear system of the form

$$\dot{\mathbf{x}} = \mathbf{f}(\mathbf{x})$$

and let \mathbf{x}^* be an equilibrium point, i.e. $\mathbf{f}(\mathbf{x}^*) = \mathbf{0}$. Let $V : D \to \mathbb{R}$ (unrelated to V in conservative systems) be a continuously differentiable function, defined on a neighbourhood D of (open set containing) \mathbf{x}^* such that

- $V(\mathbf{x}^*) = 0$ and $V(\mathbf{x}) > 0$ if $x \neq x^*$.
- $\dot{V}(\mathbf{x}) = (\operatorname{grad}_{\mathbf{x}} V(\mathbf{x}))^T \mathbf{f}(\mathbf{x}) \le 0 \text{ in } D \setminus {\mathbf{x}^*}.$

Then \mathbf{x}^* is stable. If moreover

• $\dot{V}(\mathbf{x}) < 0$ in $D \setminus {\mathbf{x}^*}$.

then \mathbf{x}^* is asymptotically stable.

5.1 Vector fields possesing an integral

5.2 Hamiltonian Systems

Let $\mathbf{x}, y \in \mathbb{R}^n$, and $H = H(\mathbf{x}, \mathbf{y})$. A system of the form

$$\dot{\mathbf{x}} = \frac{\partial H}{\partial \mathbf{y}}$$

$$\dot{\mathbf{y}} = -\frac{\partial H}{\partial \mathbf{x}},$$

is called a Hamiltonian system with n degrees of freedom.

Theorem 5.2 (Conservation of Energy). The total energy $H(\mathbf{x}, \mathbf{y})$ of a Hamiltonian system remains constant along its trajectories.

Proof. The total derivative of the Hamiltonian function $H(\mathbf{x}, \mathbf{y})$ along a solution trajectory $\mathbf{x}(t), \mathbf{y}(t)$ is

$$\begin{split} \frac{\mathrm{d}H}{\mathrm{d}t} &= \frac{\partial H}{\partial \mathbf{x}} \cdot \dot{\mathbf{x}} + \frac{\partial H}{\partial \mathbf{y}} \cdot \dot{\mathbf{y}} \\ &= \frac{\partial H}{\partial \mathbf{x}} \cdot \frac{\partial H}{\partial \mathbf{y}} - \frac{\partial H}{\partial \mathbf{y}} \cdot \frac{\partial H}{\partial \mathbf{x}} \\ &= 0 \end{split}$$

5.3 Gradient Systems

A system of the form

$$\dot{\mathbf{x}} = -\operatorname{grad}V(\mathbf{x})\tag{9}$$

is called a gradient system.

Theorem 5.3. At regular points of the function $V(\mathbf{x})$, trajectories of a gradient system cross the level surfaces $V(\mathbf{x}) = \text{constant orthogonally}$. And strict local minima of the function $V(\mathbf{x})$ are asymptotically stable equilibrium points of this gradient system.

Also, $V(\mathbf{x}) - V(\mathbf{x}_0)$ is a strict Lyapunov function for this system, where \mathbf{x}_0 is a strict local minimum of $V(\mathbf{x})$.

5.4 A relationship between Gradient and Hamiltonian Systems

They are orthogonal.

6 Asymptotic behaviour

Let the non-linear system

$$\dot{\mathbf{x}} = \mathbf{f}(\mathbf{x}) \tag{10}$$

have a flow

$$\phi(t, \mathbf{x}) : \mathbf{f}(\mathbf{x}) = \frac{\mathrm{d}}{\mathrm{d}t} \phi(0, \mathbf{x})$$

within some open set D around the point \mathbf{x}_0 . We require the solution to pass through the point \mathbf{x}_0 at t=0. This solution defines a path $\Gamma_{\mathbf{x}_0}$ which is the curve

$$\Gamma_{\mathbf{x}_0} = \{\mathbf{x} \in D : \mathbf{x} = \phi(t, \mathbf{x}_0), t \in R\}.$$

- 7 Limit Cycles and Index Theory
- 8 Local Bifurcations
- 9 Chaos