

Project Number B-06-613
Center Number 10/24-6-R7876-0A0
Project Director HALE, JACK
Project Unit CTR DYN SY
Sponsor NATL SCIENCE FOUNDATION/GENERAL
Division Id 3393
Contract Number DMS-9306265 Contract Entity GTRC
Prime Contract Number
Title TOPICS IN DYNAMICAL SYSTEMS
Effective Completion Date $30-J U N-1998$ (Performance) 30-SEP-1998 (Reports)

| Closeout Action: | $Y / N$ | Date |
| :--- | :--- | :--- |
|  |  | Submitted |

Final Invoice or Copy of Final Invoice ..... N
Final Report of Inventions and/or Subcontracts ..... N
Government Property Inventory and Related Certificate ..... N
Classified Material Certificate ..... N
Release and Assignment ..... N
Other ..... N
Comments
LETTER OF CREDIT APPLIES. 98A SATISFIES PATENT REPORT.
Distribution Required:
Project Director/Principal Investigator ..... $Y$
Research Administrative Network ..... $Y$
Accounting ..... $Y$
Research Security Department ..... N
Reports Coordinator ..... $\mathbf{Y}$
Research Property Team ..... $\mathbf{Y}$
Supply Services Department/Procurement ..... $Y$
Georgia Tech Research Corporation ..... Y
Project File ..... Y

# TOPICS IN DYNAMICAL SYSTEMS 

## ANNUAL PROGRESS REPORT

Period Covered: 7/15/93-6/30/94

Shui-Nee Chow and Jack K. Hale, PI's<br>CDSNS<br>Georgia Tech<br>Atlanta, GA 30332-0190

August, 1994

Dr. Shui-Nee Chow, Co-Principal Investigator on this project wrote a paper with Hugo Leiva, a graduate student partially supported on this NSF Grant. Attached, please find a copy of these papers. They have been submitted to two journals for consideration of publication.

# Existence and roughness of the exponential dichotomy for skew-product semiflow in Banach spaces 

by S-N. Chow and H. Leiva

# Existence and Roughness of the Exponential Dichotomy for Skew-Product Semiflow in Banach Spaces * 

Shui-Nee Chow ${ }^{\dagger}$ and Hugo Leiva ${ }^{\ddagger}$


#### Abstract

In this paper we introduce a concept of exponential dichotomy for skew-product semiflow in infinite dimensional Banach spaces which is an extension of the classic concept for evolution operators. This concept is used to study the roughness property of the skew-product semiflow. Also, we introduce the concept of discrete skew-product and give a necessary and sufficient condition for this discrete skewproduct to have a Discrete Dichotomy. After that, we give necessary and sufficient conditions for the existence of exponential dichotomy for skew-product semiflow. Moreover we prove that the exponential dichotomy for skew-product semiflow is not destroyed by small perturbation. Finally, we apply these results to parabolic partial differential equations and functional differential equations.


Key words. skew-product semiflow, exponential dichotomy, discrete skew-product, discrete dichotomy, roughness.

AMS(MOS) subject classifications. primary: 34G10; secondary: 35B40.

[^0]
## Contents

1 Introduction ..... 2
2 Preliminaries ..... 5
2.1 Linear Skew-Product Semiflow ..... 5
2.2 Projectors and Subbundles ..... 7
2.3 The Stable, the Unstable and Initial Bounded Sets. ..... 8
2.4 Exponential Dichotomy for Linear Skew-Product Semiflow ..... 9
3 Discrete Skew-Product ..... 11
3.1 Discrete Dichotomy ..... 12
3.2 Necessary and Sufficient Conditions for Discrete Dichotomy ..... 16
3.3 Equivalence Between Pointwise and Uniform Discrete Dichotomy. ..... 24
3.4 Roughness for Discrete Dichotomy. ..... 27
4 Roughness ..... 30
5 Applications ..... 39

## 1 Introduction

The concept of exponential dichotomy of linear differential equations was introduced by Perron in 1930 [24], which is concerned with the problem of conditional stability of a system

$$
\begin{equation*}
\dot{x}=A(t) x \tag{1.1}
\end{equation*}
$$

and the conection with the existence of bounded solutions of the equation

$$
\begin{equation*}
\dot{x}=A(t) x+f(x, t) \tag{1.2}
\end{equation*}
$$

where the state space is a Banch space $X$ and $t \rightarrow A(t): \mathbb{R} \rightarrow L(X)$ is bounded, continuous in the strong operator topology. For related work, see Massera and Shaffer [21], Hale [8], Levinson [15], Coppel [6], Sacker and Sell [28], [29], [30] and Palmer [22].

One of the important problems of exponential dichotomies of the equation (1.1) is its roughness. That is, they are not destroyed by small pertubations of the bounded operator $A(t)$. This was first proved by Massera and Shaffer [21] under the assuption that the original operator $A(t)$ is a bounded matrix; for the case that $A(t)$ is not a matrix the results still true if $A(t) \in L(X)$ where $X$ is and infinite dimensional Banach space and can be found in Daleckii and Krein [7]. Also, Palmer [22] proved the following Lemma:

Lemma 1.1 Let $A(t)$ and $B(t)$ be $n \times n$ matrix functions, bounded and continuous on $\left[t_{0}, \infty\right)$. Suppose the system (1.1) has an exponential dichotomy on $\left[t_{0}, \infty\right)$ with projection matrix function $P(t)$ and $B(t) \rightarrow 0$ as $t \rightarrow \infty$. Then the perturbed system

$$
\begin{equation*}
\dot{x}=[A(t)+B(t)] x \tag{1.3}
\end{equation*}
$$

also has an exponential dichotomy on $\left[t_{0}, \infty\right)$ and if $Q(t)$ is a corresponding projection matrix function

$$
\begin{equation*}
\|Q(t)-P(t)\| \rightarrow 0 \quad \text { as } \quad t \rightarrow \infty \tag{1.4}
\end{equation*}
$$

The need for a new approach arose from the fact that for the linear differential equation (1.1) with unbounded operator $A(t)$, the solutions generally speaking, either can not be extended in the direction of the negative times, or can be extended, but not uniquely. For example, for parabolic partial differential equations many authors have studied these problems, (see for example, Henry [11], Kolesov [12] and Xiao-Biao Lin [17]). In the case of Henry, he studied the existence and roughness of the exponential dichotomy of the following linear differential equation

$$
\begin{equation*}
\dot{x}=\left(A_{0}+A(t)\right) x \tag{1.5}
\end{equation*}
$$

where $t \rightarrow A(t): \mathbb{R} \rightarrow L(X)$ is bounded, continuous in the strong operator topology and $A_{0}$ is the infinitesimal generator of an analitic semigroup of bounded linear operator in the Banach space $X$. Henry's results has been generalized by A. Carvalho [2], to the case when $A_{0}$ is the infinitesimal generator of a $C_{0}$ semigroup instead of an analitic semigroup. Both of them used the relation between discrete dichotomies and exponential dichotomies.

For the case of functional differential equations we can see the work done by HaleX.B.Lin [9], H.Rodrigues [27] and M.Lizana [18].

In the case of M.Lizana, he proved that, if the linear nonautonomous functional differential equation

$$
\begin{equation*}
\dot{x}(t)=L(t) x_{t}, \quad t \geq 0 \tag{1.6}
\end{equation*}
$$

admits and exponential dichotomy on $\mathbb{R}_{+}$where $L(t): C=C\left[[-r, 0], \mathbb{R}^{n}\right] \rightarrow \mathbb{R}^{n}$ is a linear operator which is uniformly continuous and bounded with respect $t$ on $\mathbb{R}_{+}$in the operator norm of the space $\mathrm{L}\left(C, \mathbb{R}^{n}\right)$, then it is not destroyed by small perturbation. Basically he proved the following:

Lemma 1.2 If (1.6) has an exponential dichotomy on $\mathbb{R}_{+}$and $B(t)$ is a linear bounded operator for $t \geq 0$ and continuous with respect to $t$ in the operator norm of $L\left(C, \mathbb{R}^{n}\right)$, $\|B(t)\| \leq \epsilon$ for all $t \geq 0$, then

$$
\begin{equation*}
\dot{x}(t)=(L(t)+B(t)) x_{t} \tag{1.7}
\end{equation*}
$$

has an exponential dichotomy on $\mathbb{R}_{+}$if $\epsilon$ is sufficiently small.
The prove of the above Lemma is similar to that for ordinary differential equations (see [8], Lemma 5.2, p.p 125-127).

Also, Lizana proved the following Theorem
Theorem 1.1 Suppose the above mentioned hypothesis on $L$ are satisfied. If $L(t)$ is globally Lipschitz in $t$ with a constant $\rho>0$ sufficiently small and all the roots of the characteristic equation $\operatorname{det} \Delta(t, \lambda)=0$, where

$$
\begin{equation*}
\Delta(t, \lambda)=\lambda I-L(t)\left(e^{\lambda} \cdot I\right) \tag{1.8}
\end{equation*}
$$

verify the condition

$$
\begin{equation*}
|\operatorname{Re} e \lambda(t)| \geq \alpha>0, \quad \forall t \geq 0 \tag{1.9}
\end{equation*}
$$

where $\alpha$ is independent of $t$, then $\dot{x}(t)=L(t) x_{t}, \quad t \geq 0$ has an exponential dichotomy on $\mathbb{R}^{+}$, with some projection operators $P(\sigma)$ and $Q(\sigma)$, the subspace $Q(\sigma) X$ is finite dimensional and the dimension is independent of $\sigma$ on $\mathbb{R}^{+}$.

All the problems above can be treated in the unifield setting of a linear skew-product semiflow. In [31] Sacker-Sell use a concept of exponential dichotomy for skew-product semiflow with the restriction that the unstable manifold has finite dimension and they give a sufficient condition for the existence of exponential dichotomies for skew-product semiflow, which is given by the following Theorem.

Theorem 1.2 Let $\pi=(\Phi, \sigma)$ be a weakly hyperbolic skew-product semiflow on $\mathcal{E}=$ $X \times \Theta$. If $\operatorname{codimS}(\theta)=k, \quad \theta \in \Theta \quad(\mathcal{S}(\theta)$-stable manifold), then $\pi$ has an exponential dichotomy over $\Theta$.

This concept is also used by Magalhães in [19] and [20]. It is not hard to prove that the concept of exponential dichotomy used by Sacker-Sell and Magalhães is stronger than the copcept we use here.

A different characterization of the exponential dichotomy for skew-product flow in infinite Banach spaces appears in Rau [25] and [26]. He associates a strongly continuous group to the skew-product flow $\pi=(\Phi, \sigma)$ in the following way:

Given a skew-product flow $\pi=(\Phi, \sigma)$ on $\mathcal{E}=X \times \Theta$ we can associate a family $\{T(t)\}_{t \in \mathbb{R}}$ of linear operators on the Banach space $C(\Theta, X)$ defined by

$$
\begin{equation*}
T(t) f(\theta)=\Phi(\theta \cdot(-t), t) f(\theta \cdot(-t)), \quad \forall t \in \mathbb{R} \tag{1.10}
\end{equation*}
$$

for all $\theta \in \Theta$ and $f \in C(\Theta, X)$.
Proposition 1.1 The operator family $\{T(t)\}_{t \in \mathbb{R}}$ given by (1.10) is a strongly continuous group on $C(\Theta, X)$.

Since $\pi=(\Phi, \sigma)$ is a flow (two side flow) the definition of exponential dichotomy is the same as in the finite dimensional case; this allows Rau in [25] to give the following sufficient and necessary condition for the existence of exponential dichotomy for skew-product flow $\pi=(\Phi, \sigma)$.

Theorem 1.3 (Theorema 12 in [25]) Let $\{T(t)\}_{t \in \mathbb{R}}$ be the strongly continuous group given by (1.10). Then the following statements are equivalent:
(A) $\pi=(\Phi, \sigma)$ has an exponential dichotomy over $\Theta$.
(B) $\sigma(T(t)) \cap \Gamma=\emptyset, \quad \forall t \neq 0$
where $\Gamma$ denotes the unit circle in $C$.

A similar result as Rau's can be found in Latushkin-Stepin [13], [14] and Antonevich [1].
In this paper, first of all, we introduce a concept of exponential dichotomy for skewproduct semiflow weaker than the concepts used by Sacker-Sell and Magalhäes. Here we allow the unstable manifold to have infinite dimension. On the other hand, ours concept is an extension of the classic concept of exponential dichotomy for evolution operators used by Henry in [11].

Second, we introduce the concept of Skew-Product Sequence and give a necessary and sufficient condition for this sequence to have discrete dichotomy.

Third, we apply the above results to prove roughness. We also give a necessary and sufficient condition for a linear skew-product semiflow to has an exponential dichotomy.

Finally, we consider some examples as applications of these results.

## 2 Preliminaries

In this section we shall present some definitions, notations and results about Skew Product Semiflow in infinite Banach spaces, as well as the definition of exponential dichotomy for the skew-product semiflow which is one of the most important concept of this paper.

### 2.1 Linear Skew-Product Semiflow

We begin with the notion of skew-product semiflow on the trivial Banach bundle $\mathcal{E}=$ $X \times \Theta$ where $X$ is a fixed a Banach space (the state space) and $\Theta$ is a compact Hausdorff space.

Definition 2.1 Suppose that $\sigma(\theta, t)=\theta \cdot t$ is a flow on $\Theta$, i.e., the mapping $(\theta, t) \rightarrow \theta \cdot t$ is continuous, $\theta \cdot 0=\theta$ and $\theta \cdot(s+t)=(\theta \cdot s) \cdot t$, for all $s, t \in \mathbb{R}$.

A linear skew-product semiflow $\pi=(\Phi, \sigma)$ on $\mathcal{E}=X \times \Theta$ is a mapping $\pi(x, \theta, t)=$ ( $\Phi(\theta, t) x, \theta . t)$ for $t \geq 0$, with the following properties:
(1) $\Phi(\theta, 0)=I$, the identity operator on $X$, for all $\theta \in \Theta$
(2) $\lim _{t \rightarrow 0^{+}} \Phi(\theta, t) x=x$, uniformly in $\theta$. This means that for every $x \in X$ and every $\epsilon>0$ there is a $\delta=\delta(x, \epsilon)>0$ such that $\|\Phi(\theta, t) x-x\| \leq \epsilon$, for all $\theta \in \Theta$ and $0 \leq t \leq \delta$.
(3) $\Phi(\theta, t)$ is a bounded linear operator from $X$ into $X$ that satisfies the cocycle identity:

$$
\begin{equation*}
\Phi(\theta, t+s)=\Phi(\theta . t, s) \Phi(\theta, t) \quad \theta \in \Theta, \quad 0 \leq s, t \tag{2.1}
\end{equation*}
$$

(4) for all $t \geq 0$ the mapping from $\mathcal{E}$ into $X$ given by

$$
(x, \theta) \rightarrow \Phi(\theta, t) x
$$

is continuous.
The properties (2) and (3) imply that for each $(x, \theta) \in \mathcal{E}$ the solution operator $t \rightarrow$ $\Phi(\theta, t) x$ is right continuous for $t \geq 0$. In fact :

$$
\|\Phi(\theta, t+h) x-\Phi(\theta, t) x\|=\|[\Phi(\theta \cdot t, h)-I] \Phi(\theta, t) x\|
$$

which goes to 0 as $h$ goes to $0^{+}$.
For any subset $\mathcal{F} \subset \mathcal{E}$ we define the fiber

$$
\begin{equation*}
\mathcal{F}(\theta):=\{x \in X:(x, \theta) \in \mathcal{F}\}, \quad \theta \in \Theta . \tag{2.2}
\end{equation*}
$$

So $\mathcal{E}(\theta)=X \times\{\theta\}, \quad \theta \in \Theta$. If $U \subset \Theta$, then we define

$$
\mathcal{F}(U):=\bigcup_{\theta \in U} \mathcal{F}(\theta)
$$

Also we define

$$
\mathcal{E}_{0}=\{(x, \theta) \in \mathcal{E}: x=0\}
$$

the zers fiber.
Proposition 2.1 Let $\pi=(\Phi, \sigma)$ be a linear skew-product semiflow on $\mathcal{E}$. Then there exists constants $M \geq 1, \quad a>0$ such that

$$
\|\Phi(\theta, t)\| \leq M e^{a t} \quad \theta \in \Theta, \quad t \in \mathbb{R}_{+}
$$

Proof First we claim that there is $\delta>0$ such that

$$
M=\sup \{\|\Phi(\theta, t)\|: \quad \theta \in \Theta, \quad 0 \leq t \leq \delta\}<\infty
$$

For the purpose of contradiction. Let us assume that there are sequences $\theta_{n} \in \Theta, t_{n} \in \mathbb{R}_{+}$ such that $t_{n} \rightarrow 0^{+}$and $\left\|\Phi\left(\theta_{n}, t_{n}\right)\right\|>n$. The Banach-Steinhaus Theorem (Uniform Boundedness Principle) implies that there is $x \in X$ such that

$$
\left\{\left\|\Phi\left(\theta_{n}, t_{n}\right) x\right\|: n \in \mathbb{N}\right\}
$$

is unbounded. This contradicts the fact that

$$
\lim _{t \rightarrow 0^{+}} \Phi(\theta, t) x=x
$$

uniformly in $\theta \in \Theta$. Therefore $M<\infty$.
On the other hand we have that $\Phi(\theta, 0)=I$, so $M \geq 1$.
Now fix $t \in \mathbb{R}_{+}$. Let $m$ be an integer satisfying $m \leq t / \delta \leq m+1$, i.e., $\delta m \leq t \leq \delta m+\delta$. Then for any $\theta \in \Theta$ we have

$$
\begin{aligned}
\|\Phi(\theta, t)\| & =\|\Phi(\theta, t-\delta m+\delta m)\| \\
& =\|\Phi(\theta \cdot \delta,(t-\delta m)+\delta(m-1)) \Phi(\theta, \delta)\|
\end{aligned}
$$

Now putting $\theta_{0}=\theta, \quad \theta_{1}=\theta_{0} \cdot \delta, \quad \theta_{2}=\theta_{1} \cdot \delta, \cdots \theta_{m}=\theta_{m-1} \cdot \delta$; we get the following:

$$
\begin{aligned}
\|\Phi(\theta, t)\| & =\left\|\Phi\left(\theta_{m}, t-\theta m\right) \Phi\left(\theta_{m-1}, \theta\right) \cdots \Phi\left(\theta_{1}, \delta\right) \Phi\left(\theta_{0}, \delta\right)\right\| \\
& \leq M^{m+1}=M \cdot M^{m} \leq M \cdot M^{t / \delta}
\end{aligned}
$$

If we put $a=\frac{1}{\delta} \ln (M)$, then

$$
\|\Phi(\theta, t)\| \leq M e^{a t}
$$

Remark 2.1 The theory we present here extends easily to Banach bundles which are locally product spaces. A Banach bundle $\mathcal{E}$ with fiber $X$ over a base space $\Theta$ with projection $\mathcal{P}$ is denoted by $(\mathcal{E}, X, \Theta, \mathcal{P})$, or $\mathcal{E}$ for short, and is defined as follows:
(1) $X$ is a fixed Banach space and $\Theta$ is a compact Hausdorff space.
(2) The mapping $\mathcal{P}: \mathcal{E} \rightarrow \Theta$ is continuous.
(3) For each $\theta \in \Theta, \mathcal{P}^{-1}(\theta)=\mathcal{E}(\theta)$ is a Banach space, which is referred to as the fiber over $\theta$.
(4) For each $\theta \in \Theta$, there is an open neighborhood $U$ of $\theta$ in $\Theta$ and a homeomorphism $\tau: \mathcal{P}^{-1}(U) \rightarrow X \times U$ such that, for each $\eta \in U, \quad \mathcal{P}^{-1}$ is mapped onto $X \times\{\eta\}$ and $\tau: \mathcal{P}^{-1}(\eta) \rightarrow X \times\{\eta\}$ is a linear isomorphism.
(5) The norms $\|\cdot\|=\|\cdot\|_{\theta}$ on the fiber $p^{-1}(\theta)$ vary continuously in $\theta$. One can use the local coordinate notation $(x, \theta)$ to denote a typical point in a Banach bundle $\mathcal{E}$.

### 2.2 Projectors and Subbundles

A mapping $\mathbf{P}: \mathcal{E} \rightarrow \mathcal{E}$ is said to be a projector if $\mathbf{P}$ is continuous and has the form $\mathbf{P}(x, \theta)=(P(\theta) x, \theta)$, where $P(\theta)$ is a bounded linear projection on the fiber $\mathcal{E}(\theta)$.
For any projector $\mathbf{P}$ we define the range and null space by

$$
\begin{aligned}
& \mathcal{R}=\mathcal{R}(\mathbf{P})=\{(x, \theta) \in \mathcal{E}: P(\theta) x=x\} \\
& \mathcal{N}=\mathcal{N}(\mathbf{P})=\{(x, \theta) \in E: P(\theta) x=0\}
\end{aligned}
$$

The continuity of $\mathbf{P}$ implies that the fibers $\mathcal{R}(\theta)$ and $\mathcal{N}(\theta)$ vary continuously in $\theta$. This also means that $P(\theta)$ is strongly continuous in $\theta$. The following result can be found in Sacker-Sell [31].

Lemma 2.1 Let $\mathbf{P}$ be a projector on $\mathcal{E}$, then $\mathcal{R}$ and $\mathcal{N}$ are closed subsets in $\mathcal{E}$ and we have

$$
\mathcal{R}(\theta) \cap \mathcal{N}=\{0\}, \quad \mathcal{R}(\theta)+\mathcal{N}(\theta)=\mathcal{E}(\theta) \quad \theta \in \Theta .
$$

Definition 2.2 A subset $\mathcal{V}$ is said to be a subbundle of $\mathcal{E}$, if there is a projector $\mathbf{P}$ on $\mathcal{E}$ with the property that $\mathcal{R}(\mathbf{P})=\mathcal{V}$; in this case $\mathcal{W}=\mathcal{N}(\mathbf{P})$ is a complementary subbundle. i.e., $\mathcal{E}=\mathcal{V}+\mathcal{W}$ as a Whitney sum of subbundles.

For the proof of the following lemma see [31].

Lemma 2.2 Let $\mathcal{V} \subset \mathcal{E}$ with the properties:
(1) $\mathcal{V}$ is closed.
(2) $\mathcal{V}(\theta)$ is a linear subspace of $\mathcal{E}(\theta)$ for all $\theta \in \Theta$.
(3) $\operatorname{codim} \mathcal{V}(\theta)$ is finite for all $\theta \in \Theta$.
(4) $\operatorname{codim\mathcal {V}}(\theta)$ is locally constant on $\Theta$.

Then $\mathcal{V}$ is a subbundle of $\mathcal{E}$.

### 2.3 The Stable, the Unstable and Initial Bounded Sets.

Definition 2.3 A point $(x, \theta) \in \mathcal{E}$ is said to have a negative continuation with respect to $\pi$ if there exists a continuous functions $\phi:(-\infty, 0] \rightarrow \mathcal{E}$ satisfying the following properties:
(1) $\phi(t)=\left(\phi^{x}(t), \theta \cdot t\right)$ where $\phi^{x}:(-\infty, 0] \rightarrow X$
(2) $\phi(0)=(x, \theta)$
(3) $\phi(s) \in \mathcal{E}(\theta \cdot s)$ for each $s \leq 0$
(4) $\pi(\phi(s), t)=\phi(s+t)$ for each $s \leq 0$, and $0 \leq t \leq-s$

In this case the function $\phi$ is said to be a negative continuation of the point $(x, \theta)$. For any negative continuation $\phi$ and any $\tau \leq 0$, we define $\phi_{\tau}(t)=\phi(\tau+t)$ for $-\infty<$ $t<-\tau$.

Now we shall define the following sets:
$\mathcal{M}:=\{(x, \theta) \in \mathcal{E}:(x, \theta)$ has a negative continuation $\phi\}$
$\mathcal{U}:=\left\{(x, \theta) \in \mathcal{M}:\right.$ there is a negative continuation $\phi$ of $(x, \theta)$ satisfying $\left\|\phi^{x}(t)\right\| \rightarrow$ 0 as $t \rightarrow-\infty\}$
$\mathcal{B}^{+}:=\left\{(x, \theta) \in \mathcal{E}: \sup _{t \geq 0}\|\Phi(\theta, t) x\|<\infty\right\}$
$\mathcal{B}_{u}^{-}:=\{(x, \theta) \in \mathcal{M}:(x, \theta)$ has a unique bounded negative continuation $\phi\}$
$\mathcal{B}^{-}:=\left\{(x, \theta) \in \mathcal{M}:\right.$ there is a negative continuation $\phi$ of $(x, \theta)$ satisfying $\sup _{t \leq 0}\left\|\phi^{x}(t)\right\|<$ $\infty\}$
$\mathcal{S}:=\{(x, \theta) \in \mathcal{E}:\|\Phi(\theta, t) x\| \rightarrow 0$ as $t \rightarrow \infty\}$
$\mathcal{B}:=\mathcal{B}^{+} \cap \mathcal{B}^{-}$
The set $\mathcal{U}$ is called unstable set, $\mathcal{S}$ is the stable set and $\mathcal{B}$ is the initial bounded set.

Remark 2.2 The theory described here does allow for the possibility that the linear operator $\Phi(\theta, t)$ need not be one-to-one for some $t \leq 0$, i.e., $\Phi(\theta, t)$ may has a nontrivial null space. Because of this, it maybe possible for a point $(x, \theta) \in \mathcal{E}$ to have more than one negative continuation. It is easy to see that if $\Phi(\theta, t)$ is one-to-one for all $t>0$, then every negative continuation is unique. Uniqueness of negative continuations is a common feature in the study of partial differential equations, see for example, Temam [32] and Hale [10].

For all $(x, \theta) \in \mathcal{B}_{u}^{-}$we shall denote the unique bounded negative continuation, by $\Phi(\theta, t) x, t \leq 0$; this defines an extension of the mapping $\Phi$. It is clear that for each $\theta \in \Theta$ the fiber $\mathcal{B}_{u}^{-}(\theta)$ is a linear subspaces of $\mathcal{E}(\theta)$ and $\Phi(\theta, t) x$ is linear in $x$ for each $t \geq 0$, i.e., $\Phi(\theta, t)$ is a linear mapping from $\mathcal{B}_{u}^{-}(\theta)$ to $\mathcal{B}_{u}^{-}(\theta \cdot t)$ for $t \leq 0$.
Moreover, the cocycle identity

$$
\Phi(\theta, t+s) x=\Phi(\theta \cdot t, s) \Phi(\theta, t) x \quad s, t \in \mathbb{R}
$$

is valid for all $(x, \theta) \in \mathcal{B}_{u}^{-}$.

### 2.4 Exponential Dichotomy for Linear Skew-Product Semiflow

Now we shall introduce a new concept of exponential dichotomy for skew-product semiflow in infinite dimensional Banach spaces which is an extension of the concept given by Henry in [11].

Definition 2.4 A projector $\mathbf{P}$ on $\mathcal{E}$ is said to be invariant if it satisfies the following property

$$
\begin{equation*}
P(\theta \cdot t) \Phi(\theta, t)=\Phi(\theta, t) P(\theta) \quad t \geq 0, \quad \theta \in \Theta \tag{2.3}
\end{equation*}
$$

Definition 2.5 We shall said that a linear skew-product semiflow $\pi$ on $\mathcal{E}$ has an exponential dichotomy over an invariant set $\hat{\Theta}$, where $\hat{\Theta} \subset \Theta$, if there are constants $k \geq 1, \beta>0$ and invariant projector $\mathbf{P}$ such that for all $\theta \in \hat{\Theta}$ we have the following:
(1) $\Phi(\theta, t): \mathcal{N}(P(\theta)) \rightarrow \mathcal{N}(P(\theta \cdot t)), \quad t \geq 0$ is an isomorphism with inverse:

$$
\Phi(\theta \cdot t,-t): \mathcal{N}(P(\theta \cdot t)) \rightarrow \mathcal{N}(P(\theta)), \quad t \geq 0
$$

(2) $\|\Phi(\theta, t) P(\theta)\| \leq k e^{-\beta t}, \quad t \geq 0$
(3) $\| \Phi(\theta, t)\left(I-P(\theta) \| \leq k e^{\beta t}, \quad t \leq 0\right.$
where $\Phi(\theta, t)(I-P(\theta))$ is well defined for $t \leq 0$ since $\mathcal{N}(P(\theta))=\mathcal{R}(I-P(\theta))$.
Proposition 2.2 If $\pi=(\Phi, \sigma)$ is a linear skew-product semiflow on $\mathcal{E}=X \times \Theta$ which admits an exponential dichotomy over $\Theta$, then one has that $\mathcal{B}=\mathcal{E}_{0}$ and the corresponding projector $\mathbf{P}$ is such that

$$
\mathcal{R}=\mathcal{S}(\Theta), \quad \mathcal{N}=\mathcal{U}(\Theta)
$$

and

$$
\mathcal{E}=\mathcal{R}(\mathbf{P})+\mathcal{N}(\mathbf{P})=\mathcal{S}(\Theta)+\mathcal{U}(\Theta)
$$

(the Whitney sum of two subbundles)
Proof Let $\mathbf{P}$ be the corresponding projector and consider $(x, \theta) \in \mathcal{B}$. Set $y=P(\theta) x$ and $z=(I-P(\theta)) x$. Then each one of the trajectories $\Phi(\theta, t) x, \quad \Phi(\theta, t) y$ and $\Phi(\theta, t) z$ has a negative contiuation $\phi^{x}, \phi^{y}$ and $\phi^{z}$ respectively. In fact, if $\phi^{x}$ is the negative continuation of $\Phi(\theta, t)$, then:
(1) $\phi^{x}(0)=x$
(2) $\phi^{x}(t+s)=\Phi(\theta \cdot s, t) \phi^{x}(s), \quad 0 \leq t \leq-s$
(3) $\Phi(\theta \cdot s, t) \phi^{x}(s)=\Phi(\theta, t+s) x, \quad 0 \leq-s \leq t$.

From here it is easy to prove that

$$
\phi^{y}(s)=P(\theta \cdot s) \phi^{x}(s) \text { and } \phi^{z}(s)=(I-P(\theta \cdot s)) \phi^{x}(s)
$$

for all $s \leq 0$. Therefore

$$
\phi^{x}(s)=\phi^{y}(s)+\phi^{2}(s), \quad(-\infty, 0] .
$$

So

$$
\|y\| \leq k e^{\beta t}\left\|\phi^{x}\right\|, \quad \forall t \leq 0
$$

Since $\phi^{x}(t)$ is bounded and $\beta>0$, then $y=0$.
From the definition of exponential dichotomy we have that

$$
\Phi(\theta, t): \mathcal{N}(P(\theta)) \rightarrow \mathcal{N}(P(\theta \cdot t)), \quad t \geq 0
$$

is an isomorphism with inverse:

$$
\Phi(\theta \cdot t,-t): \mathcal{N}(P(\theta \cdot t)) \rightarrow \mathcal{N}(P(\theta)), \quad t \geq 0
$$

Since $\mathcal{R}(I-P(\theta))=\mathcal{N}(P(\theta))$, then we get the following: $z=(I-P(\theta)) x \in \mathcal{N}(P(\theta))$ and

$$
\begin{aligned}
z & =\Phi(\theta \cdot t,-t) \Phi(\theta, t) z \\
& =\Phi(\theta \cdot t,-t) \Phi(\theta, t)(I-P(\theta)) x \\
& =\Phi(\theta \cdot t,-t)(I-P(\theta \cdot t)) \Phi(\theta, t) x, \quad t \geq 0 .
\end{aligned}
$$

Hence

$$
\begin{aligned}
\|z\| & \leq\|\Phi(\theta \cdot t,-t)(I-P(\theta \cdot t))\|\|\Phi(\theta, t) x\| \\
& \leq k e^{-\beta t}\|\Phi(\theta, t) x\|, \quad t \geq 0 .
\end{aligned}
$$

Since $\Phi(\theta, t) x$ is bounded for $t \geq 0$ and $\beta>0$, then $z=0$. Therefore $x=y+z$. which means that $\mathcal{B}=\mathcal{E}_{0}$. This implies that $\mathcal{S} \cap \mathcal{U}=\mathcal{E}_{0}$.

Clearly we have that

$$
\mathcal{R}(\mathbf{P}) \subset \mathcal{S}, \quad \mathcal{N}(\mathbf{P}) \subset \mathcal{U}
$$

and

$$
\mathcal{E}=X \times \Theta=\mathcal{R}(\mathbf{P})+\mathcal{N}(\mathbf{P})
$$

Then

$$
\mathcal{R}(\mathbf{P})=\mathcal{S} \text { and } \mathcal{N}(\mathbf{P})=\mathcal{U}
$$

Corolary 2.1 The projector $\mathbf{P}$ associated with the exponential dichotomy of $\pi=(\Phi, \sigma)$ is uniquely determined.

Proof Assume that we have another projector $\mathbf{Q}$ associated with the exponential dichotomy of $\pi=(\Phi, \sigma)$. Then $\mathbf{Q}(x, \theta)=(Q(\theta) x, \theta)$ where $Q(\theta)$ is a linear bounded Projection on $X$.
From the previous lemma we get that

$$
\mathcal{R}(P(\theta))=\mathcal{R}(Q(\theta)), \quad \mathcal{N}(P(\theta))=\mathcal{N}(Q(\theta)), \quad \theta \in \Theta
$$

This implies that $P(\theta)=Q(\theta), \quad \theta \in \Theta$. Therefore $\mathbf{P}(x, \theta)=\mathbf{Q}(x, \theta),(x, \theta) \in \mathcal{E}$.

Lemma 2.3 Assume $\pi=(\Phi, \sigma)$ has an exponential dichotomy over $\Theta$ with constants $\beta, k$ and corresponding projector $\mathbf{P}$. Then the following holds:
(1) $\| P(\theta)) \| \leq k, \quad \theta \in \Theta$
(2) If we define

$$
G(\theta, t, s)= \begin{cases}\Phi(\theta \cdot s, t-s) P(\theta \cdot s), & \text { if } t \geq s  \tag{2.4}\\ -\Phi(\theta \cdot s, t-s)(I-P(\theta \cdot s)), & \text { if } t \leq s\end{cases}
$$

then

$$
\begin{equation*}
\|G(\theta, t, s)\| \leq k e^{-\beta|t-s|}, \quad t, s \in \mathbb{R} \tag{2.5}
\end{equation*}
$$

Proof Follows from the definition of exponential dichotomy.

## 3 Discrete Skew-Product

Here we shall introduce the concept of discrete skew-product and present some results relative to discrete dichotomy for such skew-product, so as existence and preserving of the discrete dichotomy under small perturbation.

### 3.1 Discrete Dichotomy

We begin this subsection with the concept of discrete skew-product which appears in natural way when we discretize the skew-product semiflow on time.

Definition 3.1 Let $\sigma(\theta, t)=\theta \cdot t, t \in \mathbb{R}$ be a flow on $\Theta$. A mapping $\hat{\pi}: \mathcal{E} \times Z \rightarrow \mathcal{E}$ is called discrete skew-product if it can be written as follows:

$$
\hat{\pi}(x, \theta, n)=\left(\Phi_{n}(\theta) x, \theta \cdot n\right), \quad(x, \theta) \in \mathcal{E}, \quad n \in Z
$$

where $\Phi_{n}(\theta) \in L(X)$ and has the following properties
(1) there exists $\rho>0$ such that

$$
\left\|\Phi_{n}(\theta)\right\|<\rho, \text { for all } n \in Z, \theta \in \Theta
$$

(2) for each $n \in Z$ the mapping from $\mathcal{E}$ to $X$ given by ( $x, \theta) \rightarrow \Phi_{n}(\theta) x$ is continuous.

Remark 3.1 In principle a discrete skew-product $\hat{\pi}$ is not a semiflow since there is not relation between $\Phi_{n}(\theta)$ 's. However, in section 4 we shall define a discrete skew-product for which we do have relation between the $\Phi_{n}(\theta)$ 's. Also, we do not need to assume that $\Phi_{n}(\theta)$ is invertible in whole the space $X$.

Now we shall introduce two concepts of discrete dichotomy over $\Theta$.
Definition 3.2 (pointwise discrete dichotomy) We shall say that a discrete skewproduct $\hat{\pi}$ has a pointwise discrete dichotomy over $\Theta$ if for each $\theta \in \Theta$, there exist $M_{\theta}, \alpha_{\theta}<1$ and a family of projections $P_{n}(\theta)$ on $X$ such that
(1) $\Phi_{n}(\theta) P_{n}(\theta)=P_{n+1}(\theta) \Phi_{n}(\theta), \quad n \in Z$.
(2) $\Phi_{n}(\theta): \mathcal{N}\left(P_{n}(\theta)\right) \rightarrow \mathcal{N}\left(P_{n+1}(\theta)\right), \quad n \in Z$
is an isomorphism with inverse:

$$
\Phi_{n}(\theta)^{-1}: \mathcal{N}\left(P_{n+1}(\theta)\right) \rightarrow \mathcal{N}\left(P_{n}(\theta)\right)
$$

(3) Define $\Phi_{n, m}(\theta):=\Phi_{n-1}(\theta) \Phi_{n-2}(\theta) \ldots \Phi_{m}(\theta)$, and $\Phi_{m, m}(\theta):=I, n>m$. Then

$$
\left\|\Phi_{n, m}(\theta) P_{m}(\theta) x\right\| \leq M_{\theta} \alpha_{\theta}^{n-m}\|x\|, \quad n \leq m \in Z
$$

(4) $\left\|\Phi_{n, m}(\theta)\left(I-P_{m}(\theta)\right) x\right\| \leq M_{\theta} \alpha_{\theta}^{m-n}\|x\|, \quad n<m ; \quad n, m \in Z$
where $\Phi_{n, m}(\theta)\left(I-P_{m}(\theta)\right)$ is well defined by (2) since $\mathcal{R}\left(I-P_{m}(\theta)\right)=\mathcal{N}\left(P_{m}(\theta)\right)$. In particular, if $m=0$ we get

$$
\left\|\Phi_{n, 0}(\theta) P_{0}(\theta) x\right\| \leq M_{\theta} \alpha_{\theta}^{n}\|x\|, \quad n \geq 0
$$

and

$$
\left\|\Phi_{n, 0}(\theta)\left(I-P_{0}(\theta)\right) x\right\| \leq M_{\theta} \alpha_{\theta}^{-n}\|x\|, \quad n \leq 0
$$

Definition 3.3 (uniform discrete dichotomy) We shall say that a discrete skewproduct $\hat{\pi}$ has a uniform discrete dichotomy over $\Theta$, if in definition $3.2 M_{\theta}$ and $\alpha_{\theta}$ are independent of $\theta \in \Theta\left(M_{\theta}=M, \quad \alpha_{\theta}=\alpha\right)$.

Remark 3.2 It is easy to check the following properties:

$$
\begin{equation*}
\Phi_{n, m}(\theta) \Phi_{m, k}(\theta)=\Phi_{n, k}(\theta) ; \quad \Phi_{n}(\theta) \Phi_{n, m}(\theta)=\Phi_{n+1, m}(\theta) \tag{3.1}
\end{equation*}
$$

for $n>m \geq k$.

$$
\begin{gather*}
\Phi_{n, m}(\theta) P_{m}(\theta)=P_{n}(\theta) \Phi_{n, m}(\theta)  \tag{3.2}\\
\Phi_{n, n-1}(\theta)=\Phi_{n-1}(\theta) ; \quad \Phi_{n, r}(\theta) \Phi_{r-1}(\theta)=\Phi_{n, r-1}(\theta) . \tag{3.3}
\end{gather*}
$$

for all $r \leq n, m, k \in Z$ and $\theta \in \Theta$.
Lemma 3.1 Let $\hat{\pi}: \mathcal{E} \times Z \rightarrow \mathcal{E}$ be a discrete skew-product. If for a sequence $y=$ $\left\{y_{n}\right\}$ in $X$ and $\theta \in \Theta$ we have

$$
\begin{equation*}
x_{n+1}=\Phi_{n}(\theta) x_{n}+y_{n}, \quad n \in Z \tag{3.4}
\end{equation*}
$$

Then

$$
\begin{equation*}
x_{n}=\Phi_{n, m}(\theta) x_{m}+\sum_{k=m}^{n-1} \Phi_{n, k+1}(\theta) y_{k}, \quad n>m \in Z \tag{3.5}
\end{equation*}
$$

## Proof

The case $m=n-1$ is easy. In fact, from part (3) of definiton 3.2 we have that

$$
x_{n}=\Phi_{n-1}(\theta) x_{n-1}+y_{n-1}=\Phi_{n-1}(\theta) x_{n-1}+\sum_{k=n-1}^{n-1} \Phi_{n, k+1}(\theta) y_{k}
$$

Assume that (3.5) is true for $m=r<n$. Then we shall prove it is true also for $m=r-1$. In fact,

$$
x_{n}=\Phi_{n, r}(\theta) x_{r}+\sum_{k=r}^{n-1} \Phi_{n, k+1}(\theta) y_{k}
$$

and $x_{r}=\Phi_{r-1} x_{r-1}+y_{r-1}$. Therefore, using Remark 3.1 we get the following

$$
\begin{aligned}
x_{n} & =\Phi_{n, r}(\theta)\left[\Phi_{r-1}(\theta) x_{r-1}+y_{r-1}\right]+\sum_{k=r}^{n-1} \Phi_{n, k+1}(\theta) y_{k} \\
& =\Phi_{n, r}(\theta) \Phi_{r-1}(\theta) x_{r-1}+\Phi_{n, r}(\theta) y_{r-1}+\sum_{k=r}^{n-1} \Phi_{n, k+1}(\theta) y_{k} \\
& =\Phi_{n, r-1}(\theta) x_{r-1}+\sum_{k=r-1}^{n-1} \Phi_{n, k+1}(\theta) y_{k}
\end{aligned}
$$

Lemma 3.2 Assume the discrete skew-product $\hat{\pi}: \mathcal{E} \times Z \rightarrow \mathcal{E}$ has a pointwise discrete dichotomy over $\Theta$ with the corresponding $M_{\theta}, \alpha_{\theta}$ and projections $P_{n}(\theta)$. Then the following holds:
(a) $\left\|P_{n}(\theta)\right\| \leq M_{\theta}, \quad n \in Z, \theta \in \Theta$
(b) We define the discrete Green's function as follows

$$
G_{n, m}(\theta)= \begin{cases}\Phi_{n, m}(\theta) P_{m}(\theta), & n \geq m  \tag{3.6}\\ -\Phi_{n, m}(\theta)\left(I-P_{m}(\theta)\right), & n<m\end{cases}
$$

Then

$$
\begin{equation*}
\left\|G_{n, m}(\theta)\right\| \leq M_{\theta} \alpha_{\theta}^{|n-m|}, \quad \theta \in \Theta, n, m \in Z \tag{3.7}
\end{equation*}
$$

Proof It follows from the definition 3.2 of pointwise discrete dichotomy.

Corolary 3.1 Assume the discrete skew-product $\hat{\pi}: \mathcal{E} \times Z \rightarrow \mathcal{E}$ has a uniform discrete dichotomy over $\Theta$ with constants $M$ and $\alpha<1$. Let $\left\{P_{n}(\theta)\right\}, n \in Z, \theta \in \Theta$ be the corresponding family of projections. Then we have:
(a) $\left\|P_{n}(\theta)\right\| \leq M, \quad n \in Z, \quad \theta \in \Theta$
(b) $\left\|G_{n, m}(\theta)\right\| \leq M \alpha^{|n-m|}, \quad n \in Z, \quad \theta \in \Theta$.

Lemma 3.3 Let $\hat{\pi}: \mathcal{E} \times Z \rightarrow \mathcal{E}$ be a discrete skew-product which has a pointwise discrete dichotomy over $\Theta$. Consider $\theta \in \Theta$ and $y=\left\{y_{n}\right\}$ in $l_{\infty}(Z, X)$. Then $\left\{x_{n}\right\} \in l_{\infty}(Z, X)$ is solution of

$$
\begin{equation*}
x_{n+1}=\Phi_{n}(\theta) x_{n}+y_{n}, \quad n \in Z \tag{3.8}
\end{equation*}
$$

if and only if

$$
\begin{equation*}
x_{n}=\sum_{-\infty}^{\infty} G_{n, k+1}(\theta) y_{k} \tag{3.9}
\end{equation*}
$$

Proof We have the following

$$
x_{r+1}=\Phi_{r}(\theta) x_{r}+y_{r}
$$

and using (3.5) we get

$$
x_{r}=\Phi_{n, r}(\theta) x_{n}+\sum_{k=n}^{r-1} \Phi_{r, k+1}(\theta) y_{k}, \quad r>n
$$

Then

$$
\begin{array}{r}
\Phi_{n, r}(\theta)\left[I-P_{r}(\theta)\right] x_{r}=\Phi_{n, r}(\theta) x_{r}-\Phi_{n, r}(\theta) P_{r}(\theta) x_{r} \\
\quad=\Phi_{n, r}(\theta)\left[\Phi_{r, n}(\theta) x_{n}+\sum_{k=n}^{r-1} \Phi_{r, k+1}(\theta) y_{k}\right]- \\
\\
\Phi_{n, r}(\theta) P_{r}(\theta)\left[\Phi_{r, n}(\theta) x_{n}+\sum_{k=n}^{r-1} \Phi_{r, k+1}(\theta) y_{k}\right]
\end{array}
$$

From Remark 3.1 we get

$$
\begin{aligned}
\Phi_{n, r}(\theta)\left[I-P_{r}(\theta)\right] x_{r} & =x_{n}+\sum_{k=n}^{r-1} \Phi_{n, k+1}(\theta) y_{k}-P_{n}(\theta) x_{n} \\
& -\sum_{k=n}^{r-1} \Phi_{n, k+1}(\theta) P_{k+1}(\theta) y_{k} .
\end{aligned}
$$

Hence

$$
\Phi_{n, r}(\theta)\left[I-P_{r}(\theta)\right] x_{r}=\left[I-P_{n}(\theta)\right] x_{n}+\sum_{k=n}^{r-1} \Phi_{n, k+1}(\theta)\left[I-P_{k+1}(\theta)\right] y_{k}
$$

Since

$$
\left\|\Phi_{n, r}(\theta)\left[I-P_{r}(\theta)\right] x_{r}\right\| \leq M_{\theta} \alpha_{\theta}^{r-n}\left\|x_{r}\right\|
$$

and $\left\|x_{r}\right\| \leq C$ for all $r>n$, then

$$
\left\|\Phi_{n, r}(\theta)\left[I-P_{r}(\theta)\right] x_{r}\right\| \rightarrow 0, \text { as } r \rightarrow+\infty
$$

and the series

$$
\sum_{k=n}^{\infty} \Phi_{n, k+1}(\theta)\left[I-P_{k+1}(\theta)\right] y_{k} \text { converges. }
$$

Therefore

$$
\begin{equation*}
\left[I-P_{n}(\theta)\right] x_{n}=-\sum_{k=n}^{\infty} \Phi_{n, k+1}(\theta)\left[I-P_{k+1}(\theta)\right] y_{k} \tag{3.10}
\end{equation*}
$$

For $n>m$ we have the following:

$$
x_{n}=\Phi_{n, m}(\theta) x_{m}+\sum_{k=m}^{n-1} \Phi_{n, k+1}(\theta) y_{k} .
$$

Then

$$
\begin{aligned}
P_{n}(\theta) x_{n} & =P_{n}(\theta) \Phi_{n, m}(\theta) x_{m}+\sum_{k=m}^{n-1} P_{n}(\theta) \Phi_{n, k+1}(\theta) y_{k} \\
& =\Phi_{n, m}(\theta) P_{m}(\theta) x_{m}+\sum_{k=m}^{n-1} \Phi_{n, k+1}(\theta) P_{k+1}(\theta) y_{k} .
\end{aligned}
$$

On the other hand, from definition 3.2 we have

$$
\left\|\Phi_{n, m}(\theta) P_{m}(\theta) x_{m}\right\| \leq M_{\theta} \alpha_{\theta}^{n-m}\left\|x_{m}\right\|, \quad n \geq m
$$

Since $\left\{x_{n}\right\}$ is bounded and $\alpha_{\theta}<1$, then

$$
\left\|\Phi_{n, m}(\theta) P_{m}(\theta) x_{m}\right\| \rightarrow 0, \text { as } m \rightarrow-\infty .
$$

Therefore the series

$$
\sum_{-\infty}^{n-1} \Phi_{n, k+1}(\theta) P_{k+1}(\theta) y_{k} \quad \text { converges. }
$$

Hence

$$
\begin{equation*}
P_{n}(\theta) x_{n}=\sum_{-\infty}^{n-1} \Phi_{n, k+1}(\theta) P_{k+1}(\theta) y_{k} \tag{3.11}
\end{equation*}
$$

From (3.10) and (3.11) we get:

$$
x_{n}=P_{n}(\theta) x_{n}+\left(I-P_{n}(\theta)\right) x_{n}=\sum_{-\infty}^{\infty} G_{n, k+1}(\theta) y_{k} .
$$

Now we shall prove that the sequence $\left\{x_{n}\right\}$ given by (3.9) belongs to $l_{\infty}(Z, X)$ and is a solution of (3.8). In fact,

$$
\begin{aligned}
x_{n+1} & =\sum_{-\infty}^{\infty} G_{n+1, k+1}(\theta) y_{k} \\
& =\sum_{-\infty}^{n} \Phi_{n+1, k+1}(\theta) P_{k+1}(\theta) y_{k}-\sum_{n+1}^{\infty} \Phi_{n+1, k+1}(\theta)\left[I-P_{k+1}(\theta)\right] y_{k} .
\end{aligned}
$$

Then

$$
\begin{aligned}
x_{n+1} & =\sum_{-\infty}^{n-1} \Phi_{n+1, k+1}(\theta) P_{k+1}(\theta) y_{k}-\sum_{n}^{\infty} \Phi_{n+1, k+1}(\theta)\left[I-P_{k+1}(\theta)\right] y_{k}+y_{n} \\
& =\Phi_{n}(\theta)\left[\sum_{-\infty}^{n-1} \Phi_{n, k+1}(\theta) P_{k+1}(\theta) y_{k}-\sum_{n}^{\infty} \Phi_{n, k+1}(\theta)\left[I-P_{k+1}(\theta)\right] y_{k}\right]+y_{n} \\
& =\Phi_{n}(\theta) \sum_{-\infty}^{\infty} G_{n, k+1}(\theta) y_{k}+y_{n} .
\end{aligned}
$$

Hence

$$
x_{n+1}=\Phi_{n}(\theta) x_{n}+y_{n} .
$$

Now we see that $\left\{x_{n}\right\}$ is bounded respect $\theta \in \Theta$. In fact,

$$
\left\|x_{n}\right\| \leq M_{\theta} \sum_{-\infty}^{\infty} \alpha_{\theta}^{|n-(k+1)|}\left\|y_{k}\right\| \leq M_{\theta}\|y\| \frac{1+\alpha_{\theta}}{1-\alpha_{\theta}}
$$

### 3.2 Necessary and Sufficient Conditions for Discrete Dichotomy

We begin this subsection with a theorem which is an extension of Theorem 7.6.5 [11], to the case of discrete skew-product.

Theorem 3.1 Let $\hat{\pi}: \mathcal{E} \times Z \rightarrow \mathcal{E}$ be a discrete skew-product in $\mathcal{E}=X \times \Theta$. Then the following statements are equivalent:
(A) $\hat{\pi}$ has a pointwise discrete dichotomy over $\Theta$.
(B) for $f=\left\{f_{n}\right\}$ in $l_{\infty}(Z, X)$ and $\theta \in \Theta$, there exists a unique bounded solution $\hat{x}=$ $\hat{x}(\theta)=\left\{x_{n}\right\}$ of the equation

$$
\begin{equation*}
x_{n+1}=\Phi_{n}(\theta) x_{n}+f_{n} \tag{3.12}
\end{equation*}
$$

Proof If (A) holds, then by using Lemma 3.3 we get that (B) holds with

$$
x_{n}=x_{n}(\theta)=\sum_{-\infty}^{\infty} G_{n, k+1}(\theta) y_{k}
$$

and

$$
\left\|x_{n}\right\| \leq M_{\theta} \frac{1+\alpha_{\theta}}{1-\alpha_{\theta}}\|f\|_{\infty}, \quad n \in Z
$$

Suppose (B) holds and put $\mathbb{B}=l_{\infty}(Z, X)$. Consider the linear operator $S_{\theta}: \mathbb{B} \rightarrow \mathbb{B}$ given by

$$
\hat{x}=\left\{x_{n}\right\} \in \mathbb{B} \rightarrow\left\{x_{n+1}-\Phi_{n}(\theta) x_{n}\right\} .
$$

Since $\left\|\Phi_{n}(\theta)\right\|<\rho$ for all $\theta \in \Theta$ and $n \in Z$, then $S_{\theta} \hat{x} \in \mathbb{B}$ for all $\hat{x} \in \mathbb{B}$. We shall prove the Theorem in several claims.
Claim 1. $S_{\theta}$ is a bounded linear isomorphism and

$$
\left\|S_{\theta}\right\| \leq 1+\rho, \quad \forall \theta \in \Theta
$$

It follows from (B) and Part (2) of the definition 3.1.
Claim 2. $S_{\theta}^{-1}$ is a bounded linear operator with

$$
\left\|S_{\theta}^{-1}\right\| \geq \frac{1}{1+\rho}
$$

In fact. From the open mapping Theorem we get that $S_{\theta}^{-1}$ is a bounded linear operator. On the other hand, we have the following:

$$
\left\|S_{\theta}^{-1}\right\|\left\|S_{\theta}\right\| \geq\left\|S_{\theta}^{-1} S_{\theta}\right\|=\|I\|=1 \Rightarrow(1+\rho)\left\|S_{\theta}^{-} 1\right\| \geq 1
$$

Claim 3. Define $G_{\theta}=S_{\theta}^{-1}$. Then $G_{\theta}$ can be written as follows

$$
\left(G_{\theta} f\right)_{n}=\sum_{-\infty}^{\infty} G_{n, k+1}(\theta) f_{k}, \quad n \in Z
$$

at least for sequence $\left\{f_{k}\right\}$ with $f_{k}=0$ for all large $|k|$. Where $G_{n, m}(\theta) \in L(X)$ with

$$
\left\|G_{n, m}(\theta)\right\|_{L(X)} \leq\left\|G_{\theta}\right\|_{L(\mathbb{B})}
$$

and
(a) $G_{n+1, m}(\theta)-\Phi_{n}(\theta) G_{n, m}(\theta)=0$, if $n \neq m-1$
(b) $G_{m, m}(\theta)-\Phi_{m-1}(\theta) G_{m-1, m}(\theta)=I$

Proof of the claim 3. Consider the sequence $f=\{\cdots 0, x, 0, \cdots\}$ with $x$ the $(m-1)$ th component of $f$. Define the following operator

$$
G_{n, m}(\theta) x=\left(G_{\theta} f\right)_{n}, \text { where } G_{\theta} f=\left\{\left(G_{\theta} f\right)_{n}: n \in Z\right\}
$$

Now consider $f=\left\{\cdots 0, x_{1}, x_{2}, 0, \cdots\right\}$ with $x_{1}$ the $(m-1)$ th component of $f$. Then

$$
\left(G_{\theta} f\right)_{n}=\left(G_{\theta} f_{1}\right)_{n}+\left(G_{\theta} f_{2}\right)_{n}=G_{n, m}(\theta) x_{1}+G_{n, m+1}(\theta) x_{2}
$$

where $f_{1}=\left\{\cdots 0, x_{1}, 0 \cdots\right\}$ and $f_{2}=f-f_{1}$.
In this way for sequence $f \in \mathbb{B}$ such that $f_{k}=0$ for all large $|k|$ we get that

$$
\left(G_{\theta} f\right)_{n}=\sum_{-\infty}^{\infty} G_{n, k+1}(\theta) f_{k}, \quad n \in Z
$$

It is easy to see that:

$$
G_{n, m}(\theta) \in L(X) \text { and }\left\|G_{n, m}(\theta)\right\|_{L(X)} \leq\left\|G_{\theta}\right\|_{L(\mathbb{B})}
$$

Let $\hat{x}=\left\{x_{n}\right\}$ and $f=\left\{f_{n}\right\}$ in $\mathbb{B}$ such that $G_{\theta} f=\hat{x}$.
Then we have

$$
x_{n+1}=\Phi_{n}(\theta) x_{n}+f_{n} \Longleftrightarrow\left(G_{\theta} f\right)_{n+1}=\Phi_{n}(\theta)\left(G_{\theta} f\right)_{n}+f_{n}
$$

In particular if $f=\{\cdots 0, x, 0, \cdots\}$ with $x$ the $(m-1)$ th component of $f$, then

$$
G_{n+1, m}(\theta) x=\Phi_{n}(\theta) G_{n, m}(\theta) x+f_{n}
$$

Therefore

$$
\begin{aligned}
G_{n+1, m}(\theta)-\Phi_{n}(\theta) G_{n, m}(\theta) & =0 \text { if } n \neq m-1 \\
G_{m, m}(\theta)-\Phi_{m-1}(\theta) G_{m-1, m}(\theta) & =I .
\end{aligned}
$$

Now consider

$$
\begin{align*}
P_{m}(\theta)=G_{m, m}(\theta) & =\Phi_{m-1} G_{m-1, m}(\theta)+I ., \text { i.e., }  \tag{3.13}\\
I-P_{m}(\theta) & =-\Phi_{m-1} G_{m-1, m}(\theta)
\end{align*}
$$

## Claim 4.

(c) $G_{n, m}(\theta)=\Phi_{n, m}(\theta) P_{m}(\theta)$, if $n \geq m$
(d) $\Phi_{m, n}(\theta) G_{n, m}(\theta)=-\left(I-P_{m}(\theta)\right)$, if $n<m$

Proof of the claim 4. By induction and using parts (a) and (b) of claim 3. In fact. For $n=m$ (c) is true. Suppose that (c) is true for $k>m$, i.e., $G_{k, m}(\theta)=\Phi_{k, m}(\theta) P_{m}(\theta)$.

Since $k>m$, then $k \neq m-1$. Therefore

$$
G_{k+1, m}(\theta)=\Phi_{k}(\theta) G_{k, m}(\theta)=\Phi_{k}(\theta) \Phi_{k, m}(\theta) P_{m}(\theta)
$$

So

$$
G_{k+1, m}(\theta)=\Phi_{k+1, m}(\theta) P_{m}(\theta) .
$$

(d) if $n=m-1$, then

$$
G_{m, m}(\theta)=\Phi_{m-1}(\theta) G_{m-1, m}(\theta)+I
$$

so

$$
\Phi_{m-1}(\theta) G_{m-1, m}(\theta)=-\left(I-P_{m}(\theta)\right) .
$$

Since $\Phi_{m-1}(\theta)=\Phi_{m, m-1}(\theta)$, then we get

$$
\Phi_{m, m-1}(\theta) G_{m-1, m}(\theta)=-\left(I-P_{m}(\theta)\right)
$$

Suppose the relation (d) is true for $n=k<m$ : $\Phi_{m, k}(\theta) G_{k, m}(\theta)=-\left(I-P_{m}(\theta)\right)$. Since $k-1 \neq m-1$, from Claim 3 we get

$$
\begin{aligned}
\Phi_{m, k-1}(\theta) G_{k-1, m}(\theta) & =\Phi_{m, k}(\theta) \Phi_{k-1}(\theta) G_{k-1, m}(\theta) \\
=\Phi_{m, k}(\theta) G_{k, m}(\theta) & =-\left(I-P_{m}(\theta)\right) .
\end{aligned}
$$

Claim 5. If $x_{n+1}=\Phi_{n}(\theta) x_{n}, n \geq m$ defines a bounded sequence, then $\left(I-P_{m}(\theta)\right) x_{m}=0$. Putting $x_{n}=0$ for $n<m$ we get the following

$$
x_{n+1}=\Phi_{n}(\theta) x_{n}, \text { if } n \neq m-1 \text { and } x_{m}-\Phi_{m-1}(\theta) x_{m-1}=x_{m} .
$$

Therefore $\hat{x}=\left\{x_{n}\right\}$ is a bounded solution of the equation $x_{n+1}=\Phi_{n}(\theta) x_{n}+f_{n}$ where $f=\left\{\cdots, 0, x_{m}, 0, \cdots\right\}$ with $x_{m}$ the $(m-1) t h$ componet of f .
On the other hand we know that

$$
y_{n}=\left(G_{\theta} f\right)_{n}=G_{n, m}(\theta) x_{m}
$$

is also solution of $x_{n+1}=\Phi_{n}(\theta) x_{n}+f_{n}$. Hence , $x_{n}=G_{n, m}(\theta) x_{m}$.
In particular:

$$
x_{m}=G_{m, m}(\theta) x_{m}=P_{m}(\theta) x_{m} \Rightarrow\left(I-P_{m}(\theta)\right) x_{m}=0,
$$

and the claim 5 is proved.
Claim 6. $P_{m}(\theta)$ is a projection.
In fact. Consider $x \in X$ and $x_{n}=G_{n, m}(\theta) x$. Then $x_{n+1}=\Phi_{n}(\theta) x_{n}$ if $n \geq m$ and $\left\|x_{n}\right\| \leq\left\|G_{\theta}\right\|\|x\|$. It follows from the foregoing Claim that:

$$
\left(I-P_{m}(\theta)\right) G_{m, m}(\theta) x=\left(I-P_{m}(\theta)\right) P_{m}(\theta) x=0 . \text { So } P_{m}(\theta) x=P_{m}^{2}(\theta) x, \quad \forall x \in X
$$

Claim 7. If $\left(I-P_{m}(\theta)\right) x=0$, then $\Phi_{m}(\theta) P_{m}(\theta) x=P_{m+1}(\theta) \Phi_{m}(\theta) x$.

If $\left(I-P_{m}(\theta)\right) x=0$, then

$$
\left(I-G_{m, m}(\theta)\right) x=0 \Longleftrightarrow x=G_{m, m}(\theta) x=x_{m}
$$

Consider $x_{n}=G_{n, m}(\theta) x$. Then

$$
x_{n+1}=\Phi(\theta) x_{n} \text { if } n \geq m
$$

and $\left\{x_{n}\right\}$ is bounded. Therefore

$$
\begin{aligned}
& 0=\left(I-P_{m+1}(\theta)\right) x_{m+1}=\left(I-P_{m+1}(\theta)\right) \Phi_{m}(\theta) x_{m} \\
& =\left(I-P_{m+1}(\theta)\right) \Phi_{m}(\theta) x=\Phi_{m}(\theta)\left(I-P_{m}(\theta)\right) x_{m} \\
& =\Phi_{m}(\theta)\left(I-P_{m}(\theta)\right) x
\end{aligned}
$$

So

$$
P_{m+1}(\theta) \Phi_{m}(\theta) x=\Phi_{m}(\theta) P_{m}(\theta) x
$$

and the claim 7 is established.
Claim 8. $\Phi_{m}(\theta): \mathcal{N}\left(P_{m}(\theta)\right) \rightarrow \mathcal{N}\left(P_{m+1}(\theta)\right)$ is an isomorphism for all $\theta \in \Theta$.
In fact. Assume $\left\{x_{n}: n \leq m\right\}$ is bounded and $x_{n+1}=\Phi_{n}(\theta) x_{n}, \quad n<m$.
Then by putting $x_{n}=0$ for $n>m$ we get that $\left\{x_{n}\right\}$ is solution of

$$
x_{n+1}=\Phi_{n}(\theta) x_{n}+f_{n} \text { with } f=\left\{\cdots 0,0,-\Phi_{m}(\theta) x_{m}, 0, \cdots\right\}
$$

where $-\Phi_{m}(\theta) x_{m}$ is the $m$ th component of $f$. Therefore

$$
x_{n}=G_{n, m+1}\left(-\Phi_{m}(\theta) x_{m}\right)=-G_{n, m+1}(\theta) \Phi_{m}(\theta) x_{m}
$$

In particular

$$
0=x_{m+1}=-G_{m+1, m+1}(\theta) \Phi_{m}(\theta) x_{m}=-P_{m+1}(\theta) \Phi_{m}(\theta) x_{m}
$$

so $\Phi_{m}(\theta) x_{m} \in \mathcal{N}\left(P_{m+1}(\theta)\right)$.
Now suppose that $x \in \mathcal{N}\left(P_{m}(\theta)\right)$ and let $y_{n}=G_{n, m}(\theta) x$. Then $\left\{y_{n}\right\}$ is bounded and

$$
\left\{\begin{array}{l}
y_{n+1}=\Phi_{n}(\theta) y_{n} \text { if } n \neq m-1 \\
y_{m-1}=\Phi_{m-2}(\theta) y_{m-2} \\
y_{m}==\Phi_{m-1}(\theta) y_{m-1}+x
\end{array}\right.
$$

Then the set $\left\{y_{n}: n \leq m-2\right\}$ is bounded and $y_{n+1}=\Phi_{n}(\theta) x_{n}, \quad n<m-2$. Therefore, if we put $y_{n}=0$ for $n>m-2$ we get that $\left\{y_{n}\right\}$ is solution of

$$
y_{n+1}=\Phi_{n}(\theta) y_{n}+f_{n} \text { with } f=\left\{\cdots 0,0,-\Phi_{m-2}(\theta) y_{m-2}, 0,0, \cdots\right\}
$$

where $-\Phi_{m-2}(\theta) y_{m-2}$ is the $(m-2) t h$ component of $f$.
So in the same way as before we get that $0=y_{m-1}=-G_{m-1, m-1}(\theta) \Phi_{m-2}(\theta) y_{m-2}$, which implies

$$
P_{m-1}(\theta) \Phi_{m-2}(\theta) y_{m-2}=0 \Longleftrightarrow y_{m-1}=\Phi_{m-2}(\theta) y_{m-2} \in \mathcal{N}\left(P_{m-1}(\theta)\right)
$$

Since $x \in \mathcal{N}\left(P_{m}(\theta)\right)$, then we get: $y_{m}=G_{m, m}(\theta) x=P_{m}(\theta) x=0$ and

$$
\begin{aligned}
& -x=-\left(I-P_{m}(\theta)\right) x=y_{m}-x=\left(G_{m, m}(\theta)-I\right) x \\
& =\Phi_{m-1}(\theta) G_{m-1, m}(\theta) x=\Phi_{m-1} y_{m-1} \in \Phi_{m-1}(\theta) \mathcal{N}\left(P_{m-1}(\theta)\right.
\end{aligned}
$$

Then

$$
\begin{equation*}
x \in \Phi_{m-1}(\theta) \mathcal{N}\left(P _ { m - 1 } ( \theta ) , \text { so } \mathcal { N } \left(P _ { m } ( \theta ) \subseteq \Phi _ { m - 1 } ( \theta ) \mathcal { N } \left(P_{m-1}(\theta)\right.\right.\right. \tag{3.14}
\end{equation*}
$$

Now consider $\bar{y}_{n}=y_{n}=G_{n, m}(\theta) x$, for $n<m-1$ and

$$
\bar{y}_{n+1}=y_{n}=\Phi_{n}(\theta) \bar{y}_{n}, \quad n+1 \geq m-1 \Longleftrightarrow n \geq m-2 .
$$

Then

$$
\bar{y}_{m-1}=\Phi_{m-2}(\theta) \bar{y}_{m-2}=\Phi_{m-2}(\theta) y_{m-2}=y_{m-1}
$$

and

$$
\bar{y}_{m}=\Phi_{m-1}(\theta) \bar{y}_{m-1}=\Phi_{m-1}(\theta) y_{m-1}=-x .
$$

So we have that $\left\{\bar{y}_{n}: n \leq m\right\}$ is bounded and

$$
\bar{y}_{n+1}=\Phi_{n}(\theta) \bar{y}_{n}, \quad n<m
$$

In the same way as before we get the following:

$$
0=-P_{m+1}(\theta) \Phi_{m}(\theta) \bar{y}_{m}=-P_{m+1}(\theta) \Phi_{m}(\theta)(-x)
$$

then

$$
\begin{equation*}
\Phi_{m}(\theta) x \in \mathcal{N}\left(P_{m+1}(\theta)\right) \tag{3.15}
\end{equation*}
$$

If $\Phi_{m}(\theta) x=0$, then $\bar{y}_{m+1}=\Phi_{m}(\theta) \bar{y}_{m}=-\Phi_{m}(\theta) x=0$. Therefore $\bar{y}_{n}=0$ for all $n \leq m$.
Thus $\left\{\bar{y}_{n}\right\}$ is bounded for $n \geq m$ and $\bar{y}_{n+1}=\Phi_{n}(\theta) \bar{y}_{n}, \quad n \leq m$.
So from Claim 5 we get the following

$$
\begin{equation*}
\left(I-P_{m}(\theta)\right) \bar{y}_{m}=\left(I-P_{m}(\theta)\right)(-x)=0 \Rightarrow x-P_{m}(\theta) x=x=0 \tag{3.16}
\end{equation*}
$$

In conclusion we have proved the following:
From (3.14) $\mathcal{N}\left(P_{m+1}(\theta)\right) \subseteq \Phi_{m}(\theta)\left(\mathcal{N}\left(P_{m}(\theta)\right)\right.$
From (3.15) $\Phi_{m}(\theta)\left(\mathcal{N}\left(P_{m}(\theta)\right) \subseteq \mathcal{N}\left(P_{m+1}(\theta)\right)\right.$
From (3.16) $\Phi_{m}(\theta)$ is one to one. Therefore

$$
\Phi_{m}(\theta): \mathcal{N}\left(P_{m}(\theta)\right) \rightarrow \mathcal{N}\left(P_{m+1}(\theta)\right)
$$

is an isomorphism.
Now we can write the Claim 4 in the following way

$$
\begin{cases}(c) G_{n, m}(\theta)=\Phi_{n, m}(\theta) P_{m}(\theta), & \text { if } n \geq m \\ (d) G_{n, m}(\theta)=-\Phi_{n, m}(\theta)\left(I-P_{m}(\theta)\right), & \text { if } n<m\end{cases}
$$

Claim 9. If $P_{m}(\theta) x=0$, then

$$
\Phi_{m}(\theta) P_{m}(\theta) x=P_{m+1}(\theta) \Phi_{m}(\theta) x
$$

If $x \in \mathcal{N}\left(P_{m}(\theta)\right) \Rightarrow \Phi_{m}(\theta) P_{m}(\theta) x=0$ and $\Phi_{m}(\theta) x \in \mathcal{N}\left(P_{m+1}(\theta)\right)$, so $P_{m+1}(\theta) \Phi_{m}(\theta) x=$ 0 . Therefore

$$
\Phi_{m}(\theta) P_{m}(\theta) x=P_{m+1}(\theta) \Phi_{m}(\theta) x
$$

Claim 10.

$$
\Phi_{m}(\theta) P_{m}(\theta) x=P_{m+1}(\theta) \Phi_{m}(\theta) x, \quad \theta \in \Theta, \quad m \in Z .
$$

It follows from Claims 7,9 and $I=P_{m}(\theta)+\left(I-P_{m}(\theta)\right)$.
Now we are ready to prove part (3) and (4) of Definition 3.2 of pointwise discrete dichotomy.
Claim 11. For all $\theta \in \Theta$ we have:

$$
\begin{equation*}
\left\|\Phi_{n, m}(\theta) P_{m}(\theta) x\right\| \leq\left\|G_{\theta}\right\|^{2}\left(1-\left\|G_{\theta}\right\|_{-1}\right)^{n-m}, \quad n \geq m \tag{3.17}
\end{equation*}
$$

If $\Phi_{n, m}(\theta) P_{m}(\theta) x=0$ for some $n \geq m$, then $\Phi_{p, m}(\theta) P_{m}(\theta) x=0$, for $p \geq n$; it follows from: $\Phi_{n+1, m}(\theta) x=\Phi_{n}(\theta) \Phi_{n, m}(\theta) x$.
Assume that $\Phi_{n, m}(\theta) P_{m}(\theta) x \neq 0$. Then

$$
\Phi_{k}^{-1}(\theta):=\left\|\Phi_{k, m}(\theta) P_{m}(\theta) x\right\|>0, \quad m \leq k \leq n .
$$

From Claim 4 we have:

$$
G_{n, k}(\theta)=\Phi_{n, k}(\theta) P_{k}(\theta) \text { and } V_{k}(\theta):=\Phi_{k, m}(\theta) P_{m}(\theta) x \cdot \phi_{k}(\theta), \quad m \leq k \leq n
$$

Then

$$
\begin{aligned}
\sum_{k=m}^{n} G_{n, k}(\theta) V_{k}(\theta) & =\sum_{k=m}^{n} \Phi_{n, k}(\theta) P_{k}(\theta) \Phi_{k, m}(\theta) P_{m}(\theta) x \cdot \phi_{k}(\theta) \\
& =\sum_{k=m}^{n} \Phi_{n, k}(\theta) \Phi_{k, m}(\theta) P_{m}(\theta) P_{m}(\theta) x \cdot \phi_{k}(\theta) \\
& =\sum_{k=m}^{n} \Phi_{n, m}(\theta) P_{m}(\theta) x \cdot \phi_{k}(\theta) \\
& =\Phi_{n, m}(\theta) P_{m}(\theta) x \cdot \sum_{k=m}^{n} \phi_{k}(\theta) .
\end{aligned}
$$

Let $f=\left\{\cdots 0, V_{m}(\theta), \cdots, V_{n}(\theta), 0, \cdots\right\}$ with $V_{m}$ the $(m-1)$ th component of $f$. Then

$$
\|f\|=\sup _{k \in Z}\left\|f_{k}\right\|=\sup _{m \leq k \leq n}\left\|V_{k}(\theta)\right\|=1
$$

and

$$
\begin{array}{r}
\left\|\sum_{k=m}^{n} G_{n, k}(\theta) V_{k}(\theta)\right\|=\left\|\sum_{k=m}^{n} G_{n, k}(\theta) f_{k}\right\| \\
=\left\|\left(G_{\theta} f\right)_{n}\right\| \leq\left\|G_{\theta}\right\|\|f\|=\left\|G_{\theta}\right\| .
\end{array}
$$

Therefore

$$
\Phi_{n}^{-1}(\theta) \cdot \sum_{k=m}^{n} \Phi_{k}(\theta)=\left\|\Phi_{n, m}(\theta) P_{m}(\theta) x \cdot \sum_{k=m}^{n} \Phi_{k}(\theta)\right\| \leq\left\|G_{\theta}\right\| .
$$

Then we have:

$$
\frac{-1}{\Phi_{n}^{-1}(\theta) \cdot \sum_{k=m}^{n} \Phi_{k}(\theta)} \leq-\left\|G_{\theta}\right\|^{-1}, \text { and }\left\|G_{\theta}\right\|>1
$$

Define $\Psi_{n}(\theta):=\sum_{k=m}^{n} \Phi_{k}(\theta)$. Then we get the following

$$
\begin{gathered}
\Psi_{n-1}(\theta)=\sum_{k=m}^{n-1} \Phi_{k}(\theta)=\sum_{k=m}^{n} \Phi_{k}(\theta)-\phi_{n}(\theta) \\
=\Psi_{n}(\theta)\left(1-\frac{1}{\phi_{n}^{-1}(\theta) \cdot \Psi_{n}(\theta)}\right) \leq \Psi_{n}(\theta)\left(1-\left\|G_{\theta}\right\|^{-1}\right) .
\end{gathered}
$$

Then $\Psi_{n}(\theta)-\Psi_{n-1}(\theta) \geq \Psi_{n}(\theta)\left\|G_{\theta}\right\|^{-1}$. So we get

$$
\phi_{n}(\theta) \geq\left\|G_{\theta}\right\|^{-1}\left(1-\left\|G_{\theta}\right\|^{-1}\right)^{-1} \Psi_{n-1}(\theta)
$$

Thus

$$
\left.\phi_{n}(\theta) \geq\left\|G_{\theta}\right\|^{-1}\left(1-\left\|G_{\theta}\right\|^{-1}\right)^{-1}\right)^{-(n-m)} \Psi_{m}(\theta)
$$

i.e.,

$$
\left.\phi_{n}(\theta) \geq\left\|G_{\theta}\right\|^{-1}\left(1-\left\|G_{\theta}\right\|^{-1}\right)^{-1}\right)^{-(n-m)} \phi_{m}(\theta)
$$

i.e.,

$$
\left.\phi_{n}^{-1}(\theta) \leq\left\|G_{\theta}\right\|^{-1}\left(1-\left\|G_{\theta}\right\|^{-1}\right)^{-1}\right)^{n-m} \phi_{m}^{-1}(\theta)
$$

Hence

$$
\begin{aligned}
\left\|\Phi_{n, m}(\theta) P_{m}(\theta) x\right\| & \leq\left\|G_{\theta}\right\|\left(1-\left\|G_{\theta}\right\|^{-1}\right)^{n-m}\left\|P_{m}(\theta) x\right\| \\
& =\left\|G_{\theta}\right\|\left(1-\left\|G_{\theta}\right\|^{-1}\right)^{n-m}\left\|G_{n, m}(\theta) x\right\| \\
& \leq\left\|G_{\theta}\right\|^{2}\left(1-\left\|G_{\theta}\right\|^{-1}\right)^{n-m}\|x\|, \quad n \geq m
\end{aligned}
$$

We have proved this by assuming that the left hand of (3.17) was not trivial, in other case it is evidently true. This completes the proof of Claim 11.

Claim 12. For all $\theta \in \Theta$ we have:

$$
\begin{equation*}
\left\|\Phi_{n, m}(\theta)\left(I-P_{m}(\theta)\right) x\right\| \leq\left(1+\left\|G_{\theta}\right\|\right)^{2}\left(\frac{\left\|G_{\theta}\right\|}{1+\left\|G_{\theta}\right\|}\right)^{m-n}\|x\|, \text { if } n<m . \tag{3.18}
\end{equation*}
$$

In fact, it can be proved in the same way as (3.17).
Finally if we define for all $\theta \in \Theta$ the following numbers:

$$
\begin{equation*}
M_{\theta}:=\left(1+\left\|G_{\theta}\right\|\right)^{2}>\left\|G_{\theta}\right\|^{2} \tag{3.19}
\end{equation*}
$$

$$
\begin{equation*}
\alpha_{\theta}:=\frac{\left\|G_{\theta}\right\|}{1+\left\|G_{\theta}\right\|} \geq 1-\left\|G_{\theta}\right\|^{-1} . \tag{3.20}
\end{equation*}
$$

Then (3.17) and (3.18) imply that $\hat{\pi}: \mathcal{E} \times \boldsymbol{Z} \rightarrow \mathcal{E}$ has pointwise discrete dichotomy over $\Theta$ with constants $M_{\theta}$ and $\alpha_{\theta}$ given by (3.19) and (3.20).

### 3.3 Equivalence Between Pointwise and Uniform Discrete Dichotomy.

Now we want to know when pointwise discrete dichotomy and uniform discrete dichotomy are equivalent. The answer fo this question is given by the following lemmas.

Lemma 3.4 Let $\hat{\pi}: \mathcal{E} \times Z \rightarrow \mathcal{E}$ be a discrete skew-product and suppose there exists $1<\delta<2$ such that $\left\|\Phi_{n}(\theta) x\right\| \geq \delta\|x\|$ for all $n \in Z, \quad \theta \in \Theta, x \in X$. Then pointwise discrete dichotomy and uniform discrete dichotomy of $\hat{\pi}$ over $\Theta$ are equivalent.
Proof Clearly uniform discrete dichotomy implies pointwise discrete dichotomy .
Assume that $\hat{\pi}$ has pointwise discrete dichotomy with constants $M_{\theta}, \alpha_{\theta}<1$ for all $\theta \in \Theta$. Then from Theorem 3.1 we have that $M_{\theta}, \alpha_{\theta}$ are given by the formulas (3.19) and (3.20) respectively. Our goal is to find constants $M, \alpha<1$ indenpendent of $\theta \in \Theta$.
Let $x \in \mathbb{B}=l_{\infty}(Z, X)$ and consider

$$
\begin{aligned}
\left\|S_{\theta} \hat{x}\right\|_{\infty} & =\sup _{n \in Z}\left\{\left\|x_{n+1}-\Phi_{n}(\theta) x_{n}\right\|\right\} \\
& \geq \sup _{n \in Z}\left\{\left\|\Phi_{n}(\theta) x_{n}\right\|-\left\|x_{n+1}\right\|\right\} \\
& \geq \sup _{n \in Z}\left\{\delta\left\|x_{n}\right\|-\left\|x_{n+1}\right\|\right\} \\
& \geq(\delta-1)\|\hat{x}\|_{\infty}=\gamma\|x\|_{\infty}, \quad \gamma=\delta-1<1
\end{aligned}
$$

Therefore, $\quad\left\|S_{\theta} \hat{x}\right\|_{\infty} \geq \gamma\|\hat{x}\|, \hat{x} \in \mathbb{B}$. Then

$$
\|\hat{x}\|=\left\|S_{\theta} S_{\theta}^{-1} \hat{x}\right\| \geq \gamma\left\|S_{\theta}^{-1} \hat{x}\right\|, \quad \hat{x} \in \mathbb{B}
$$

So $\quad\left\|S_{\theta}^{-1} \hat{x}\right\| \leq \frac{1}{\gamma}\|\hat{x}\| \Rightarrow\left\|S_{\theta}^{-1}\right\| \leq \frac{1}{\gamma}, \quad \theta \in \Theta$.
Hence the operator $G_{\theta}=S_{\theta}^{-1}$ given by Theorem 3.1 has the following property

$$
1<\left\|G_{\theta}\right\| \leq \frac{1}{\gamma}
$$

Then, from (3.19) and (3.20) it is enough to take $M$ and $\alpha$ as follows:

$$
\begin{aligned}
M & :=\left(1+\frac{1}{\gamma}\right)^{2} \geq\left(1+\left\|G_{\theta}\right\|\right)^{2} \geq\left\|G_{\theta}\right\|^{2} \\
1>\alpha & :=\frac{\frac{1}{\gamma}}{1+\frac{1}{\gamma}} \geq \frac{\left\|G_{\theta}\right\|}{1+\left\|G_{\theta}\right\|} \geq 1-\left\|G_{\theta}\right\|^{-1}
\end{aligned}
$$

or

$$
\begin{aligned}
M & =\sup _{\theta \in \Theta}\left(1+\left\|G_{\theta}\right\|\right)^{2}<\infty \\
\alpha & =\sup _{\theta \in \Theta} \frac{\left\|G_{\theta}\right\|}{1+\left\|G_{\theta}\right\|}<1
\end{aligned}
$$

Lemma 3.5 Let $\hat{\pi}: \mathcal{E} \times Z \rightarrow \mathcal{E}$ be a discrete skew-product and suppose there exist $0<\delta_{1}<\delta_{2}<1$ such that

$$
\begin{equation*}
\delta_{1}\|x\| \leq\left\|\Phi_{n}(\theta) x\right\| \leq \delta_{2}\|x\|, \quad n \in Z, \quad \theta \in \Theta, \quad x \in X . \tag{3.21}
\end{equation*}
$$

Then pointwise discrete dichotomy and uniform discrete dichotomy of $\hat{\pi}$ over $\Theta$ are equivalent.

Proof Clearly uniform discrete dichotomy implies pointwise discrete dichotomy .
Assume that $\hat{\pi}$ has pointwise discrete dichotomy with constants $M_{\theta}, \alpha_{\theta}<1$ for all $\theta \in \Theta$. Then from Theorem 3.1 we have that $M_{\theta}, \alpha_{\theta}$ are given by the formulas (3.19) and (3.20) respectively. Again, our goal is to find constants $M, \alpha<1$ indenpendent of $\theta \in \Theta$. In order to do that, we shall consider the following linear bounded operators:

$$
\begin{equation*}
T(\theta): \mathbb{B} \rightarrow \mathbb{B}, \quad(T(\theta) \hat{x})_{n}=\Phi_{n}(\theta) x_{n}, \quad \theta \in \Theta \tag{3.22}
\end{equation*}
$$

$$
L: \mathbb{B} \rightarrow \mathbb{B}, \quad(L \hat{x})_{n}=x_{n+1}, \quad \text { - the shift operator. }
$$

Then the operator $S_{\theta}: \mathbb{B} \rightarrow \mathbb{B}$ given by Theorem 3.1 can be written as follows

$$
\begin{equation*}
S_{\theta}=L-\Phi(\theta) \tag{3.23}
\end{equation*}
$$

It is easy to see that L is an isomorphism with $\|L \hat{x}\|=\|\hat{x}\|, \hat{x} \in X$. So $\|L\|=\left\|L^{-1}\right\|=1$. On the other hand, we have that: $\|T(\theta) \hat{x}\|_{\infty}=\sup _{n \in Z}\left\|\Phi_{n}(\theta) x_{n}\right\|$. Hence, from (3.21) we get

$$
\delta_{1} \leq\|T(\theta)\| \leq \delta_{2}<1, \quad \theta \in \Theta
$$

Therefore $\left\|L^{-1} T(\theta)\right\|<1$. It is well known that $\left(I-L^{-1} T(\theta)\right)$ is invertible and

$$
\left(I-L^{-1} T(\theta)\right)^{-1}=\sum_{k=0}^{\infty}\left(L^{-1} T(\theta)\right)^{k}
$$

From (3.23) we have the following

$$
S_{\theta}^{-1}=\left[L\left(I-L^{-1} T(\theta)\right)\right]^{-1}=\left(I-L^{-1} T(\theta)\right)^{-1} L^{-1}
$$

Hence

$$
\begin{aligned}
\left\|S_{\theta}^{-1}\right\| & \leq\left\|\left(I-L^{-1} T(\theta)\right)^{-1}\right\| \leq \sum_{k=0}^{\infty}\|T(\theta)\|^{k} \\
& =\frac{1}{1-\|T(\theta)\|} \leq \frac{1}{1-\delta_{1}}=\frac{1}{\gamma}, \quad \gamma=1-\delta_{1}, \quad \theta \in \Theta .
\end{aligned}
$$

From here the proof follows as in Lemma 3.4.

Corolary 3.2 If $\hat{\pi}: \mathcal{E} \times Z \rightarrow \mathcal{E}$ has pointwise discrete dichotomy over $\Theta$, then the family of projections $\left\{P_{n}(\theta): \theta \in \Theta, n \in Z\right\}$ associated with the pointwise discrete dichotomy of $\hat{\pi}$ is unique.

Proof Assume that $\left\{Q_{n}(\theta)\right\}$ is another family of projection associated with the pointwise discrete dichotomy of $\hat{\pi}$. Let $x \in X$ and consider the sequence $f=\left\{f_{n}\right\}$ such that $f_{m-1}=x$, and $f_{n}=0$ for $n \neq m-1$. For $\theta \in \Theta$ consider $\left\{x_{n}\right\}$ the unique bounded solution of

$$
x_{n+1}=\Phi_{n}(\theta) x_{n}+f_{n} .
$$

Then

$$
x_{n}=\sum_{-\infty}^{\infty} G_{n, k+1}(\theta) f_{k}=\sum_{-\infty}^{\infty} \bar{G}_{n, k+1}(\theta) f_{K}
$$

then $G_{n, m}(\theta) x=\bar{G}_{n, m}(\theta) x$. Where

$$
\begin{aligned}
& G_{n, m}(\theta)=\left\{\begin{array}{cc}
\Phi_{n, m}(\theta) P_{m}(\theta), & \text { if } n \geq m \\
-\Phi_{n, m}(\theta)\left[I-P_{m}(\theta)\right], & \text { if } n<m
\end{array}\right. \\
& \bar{G}_{n, m}(\theta)=\left\{\begin{array}{cc}
\begin{array}{ll}
\Phi_{n, m}(\theta) Q_{m}(\theta), & \text { if } n \geq m \\
-\Phi_{n, m}(\theta)\left[I-Q_{m}(\theta)\right], & \text { if } n<m
\end{array}
\end{array}\right.
\end{aligned}
$$

So

$$
G_{m, m}(\theta) x=\bar{G}_{m, m}(\theta) x \Longleftrightarrow P_{m}(\theta) x=Q_{m, m}(\theta) x
$$

Therefore $P_{m}(\theta)=Q_{m}(\theta), \quad \theta \in \Theta, \quad m \in Z$.

Corolary 3.3 For $\gamma>0$ consider

$$
S_{\gamma}=\left\{\hat{x}=\left\{x_{n}\right\} \in \mathbb{B}: \sup _{n \in Z} \gamma^{n}\left\|x_{n}\right\|<\infty\right\}
$$

where $\mathbb{B}=l_{\infty}(Z, X)$. Suppose that for each $\theta \in \Theta$ we have $0<\alpha_{\theta}<1, \alpha_{\theta} \leq \gamma \leq \frac{1}{\alpha_{\theta}}$ and for each $f \in S_{\gamma}$ there exists a unique solution of

$$
x_{n+1}=\Phi_{n}(\theta) x_{n}+f_{n} \text { in } S_{\gamma} .
$$

Then the skew-product sequence $\hat{\pi}$ has pointwise discrete dichotomy over $\Theta$.
Moreover, if $\left\{G_{n, m}(\theta)\right\}$ is the corresponding Green's family, then $\left\|G_{n, m}(\theta)\right\| \leq M \alpha_{\theta}^{|n-m|}$.

### 3.4 Roughness for Discrete Dichotomy.

In this section we shall prove that discrete dichotomy for discrete skew-product $s$ is not destroyed by small perturbations. For the proof of the following Theorem we shall use a Lemma from Henry [11], which says

Lemma 3.6 If $a \geq 0, \quad b \geq 0, \quad 0<r<r_{1}, \quad r_{2}<1, \quad b<\frac{r_{j}-r}{1+r r_{j}}$ for $j=1,2$ and $\left\{g_{n}\right\}$ is a nonnegative sequence in $\mathbb{R}$ with $g_{n}=O\left(r_{2}^{-|n|}\right)$ as $|n| \rightarrow \infty$, and

$$
g_{n} \leq a r_{1}^{|n|}+b \sum_{-\infty}^{\infty} r^{|n-k-1|} g_{k}, \quad n \in Z
$$

then

$$
g_{n} \leq \frac{a r_{1}^{|n|}}{\frac{1-b\left(1+r r_{1}\right)}{r_{1}-r}}, n \in Z
$$

Theorem 3.2 Assume the discrete skew-product $\hat{\pi}: \mathcal{E} \times Z \rightarrow \mathcal{E}$ has a pointwise discrete dichotomy over $\Theta$ with constants $M_{\theta}, \alpha_{\theta}<1$. If $\tilde{M}_{\theta}>M_{\theta}$ and $\alpha_{\theta}<\tilde{\alpha}_{\theta}<1$, then there exists $\epsilon_{\theta}>0$, such that for any discrete skew-product $\hat{\pi}_{1}: \mathcal{E} \times Z \rightarrow \mathcal{E}, \quad \hat{\pi}_{1}=\left(\Psi_{n}(\cdot), \sigma\right)$ with

$$
\sup _{n \in \boldsymbol{Z}}\left\|\Phi_{n}(\theta)-\Psi_{n}(\theta)\right\|<\epsilon_{\theta}, \quad \theta \in \Theta
$$

has a pointwise discrete dichotomy over $\theta \in \Theta$.
Proof From Lemma 3.6 the requirement on $\epsilon_{\theta}$ is

$$
\epsilon_{\theta} M_{\theta}<\frac{\tilde{\alpha}_{\theta}-\alpha_{\theta}}{1+\alpha_{\theta} \tilde{\alpha}_{\theta}}
$$

and

$$
\tilde{M}_{\theta}\left(1-\frac{\epsilon_{\theta} M_{\theta}\left(1+\alpha_{\theta} \tilde{\alpha}_{\theta}\right)}{\tilde{\alpha}_{\theta} \alpha_{\theta}}\right) \geq M_{\theta}, \quad \theta \in \Theta
$$

For each $\theta \in \Theta$ and $f \in \mathbb{B}, x_{n+1}=\Psi_{n}(\theta) x_{n}+f_{n}$ has a unique solution $\hat{x} \in \mathbb{B}$ if and only if

$$
x_{n}=\sum_{\infty}^{\infty} G_{n, k+1}(\theta)\left[\left(\Psi_{k}(\theta)-\Phi_{n}(\theta)\right) x_{k}+f_{k}\right]
$$

is also solvable for each $f \in \mathbb{B}$; and it is true provided that

$$
\sup _{n} \sum_{\infty}^{\infty}\left\|G_{n, K+1}(\theta)\left(\Psi_{k}(\theta)-\Phi_{n}(\theta)\right)\right\| \leq \epsilon_{\theta} M_{\theta} \frac{1+\alpha_{\theta}}{1-\alpha_{\theta}}<1 .
$$

In this case Theorem 3.1 shows that $\hat{\pi}_{1}$ has pointwise discrete dichotomy over $\Theta$.

Let $\left\{\hat{G}_{n, m}(\theta)\right\}$ be the corresponding discrete Green'sfunction. Then for all $n, m$ we have:

$$
\hat{G}_{n, m}(\theta)=G_{n, m}(\theta)+\sum_{\infty}^{\infty} G_{n, K+1}(\theta)\left(\Psi_{k}(\theta)-\Phi_{n}(\theta)\right) \hat{G}_{n, m}(\theta)
$$

So

$$
\left\|\hat{G}_{n, m}(\theta)\right\| \leq M_{\theta} \alpha_{\theta}^{|n-m|}+\epsilon_{\theta} M_{\theta} \sum_{-\infty}^{\infty} \alpha_{\theta}^{|n-k-1|}\left\|\hat{G}_{n, m}(\theta)\right\|
$$

and $\left\|\hat{G}_{n, m}(\theta)\right\|$ is bounded.
From Lemma 3.6 above we get the following: If $\alpha_{\theta}<\tilde{\alpha}_{\theta}<1$ and

$$
\epsilon_{\theta} M_{\theta}<\frac{\tilde{\alpha}_{\theta}-\alpha_{\theta}}{1+\alpha_{\theta} \tilde{\alpha}_{\theta}}<\frac{1-\alpha_{\theta}}{1+\alpha_{\theta}}<1
$$

then

$$
\left\|\hat{G}_{n, m}(\theta)\right\| \leq \frac{M_{\theta} \tilde{\alpha}_{\theta}^{|n-m|}}{\frac{1-\epsilon_{\theta} M_{\theta}\left(1+\alpha_{\theta} \tilde{\alpha}_{\theta}\right)}{\tilde{\alpha}_{\theta}-\alpha_{\theta}}} \leq \tilde{M}_{\theta} \tilde{\alpha}_{\theta}^{|n-m|}, \text { for small } \epsilon_{\theta}
$$

Corolary 3.4 Assume the discrete skew-product $\hat{\pi}: \mathcal{E} \times Z \rightarrow \mathcal{E}$ has uniform discrete dichotomy over $\Theta$ with constants $M, \alpha<1$, then there exist $\epsilon>0$, such that any discrete skew-product $\hat{\pi}_{1}: \mathcal{E} \times Z \rightarrow \mathcal{E}$ with

$$
\sup _{n \in Z}\left\|\Phi_{n}(\theta)-\Psi_{n}(\theta)\right\|<\epsilon, \quad \theta \in \Theta
$$

has uniform discrete dichotomy over $\Theta$.
Theorem 3.3 Assume that the discrete skew-product $s \hat{\pi}_{1}, \hat{\pi}_{2}$ have pointwise discrete dichotomy with discrete Green'sfunction satisfying

$$
\left\|G_{n, m}^{i}(\theta)\right\| \leq M_{\theta} \alpha_{\theta}^{|n-m|}, \quad i=1,2 . \quad \alpha_{\theta}<1
$$

Then
(a) if $\left\|\Phi_{n}^{1}(\theta)-\Phi_{n}^{2}(\theta)\right\| \leq \epsilon$ for $|n-m| \leq N$ and $\left\|\Phi_{n}^{1}(\theta)-\Phi_{n}^{2}(\theta)\right\| \leq B \quad n \in Z$; then

$$
\left\|P_{m}^{1}(\theta)-P_{m}^{2}(\theta)\right\| \leq \frac{2 M_{\theta}^{2}}{1-\alpha_{\theta}^{2}}\left(\epsilon+B \alpha_{\theta}^{2 N+1}\right)
$$

(b) if $b<\alpha_{\theta}^{-2}$ and $\left\|\Phi_{n}^{1}(\theta)-\Phi_{n}^{2}(\theta)\right\| \leq \epsilon b^{|n-m|}, \quad n \in Z$; then

$$
\left\|P_{m}^{1}(\theta)-P_{m}^{2}(\theta)\right\| \leq \alpha_{\theta}(1+b) \frac{M_{\theta}^{2} \epsilon}{1-b \alpha_{\theta}^{2}}
$$

Lemma 3.7 Assume that $\hat{\pi}: \mathcal{E} \times Z_{+} \rightarrow \mathcal{E}$ has a pointwise discrete dichotomy over $\Theta$ with constants $M_{\theta}, \alpha_{\theta}<1$ and $G_{n, m}(\theta), n, m \geq 0$ the corresponding discrete Green'sfunction. Define the extension of $\hat{\pi}$ to $\mathcal{E} \times \boldsymbol{Z}$ as follows: For all $\theta \in \Theta$ and $n<0$

$$
P_{n}(\theta)=P_{0}(\theta) \text { and } \Phi_{n}(\theta):=\alpha_{\theta}\left(I-P_{0}(\theta)\right)+\alpha_{\theta}^{-1} P_{0}(\theta) .
$$

Then $\hat{\pi}: \mathcal{E} \times \boldsymbol{Z} \rightarrow \mathcal{E}$ has pointwise discrete dichotomy over $\Theta$ with constants $M_{\theta}, \alpha_{\theta}<1$.
Proof We shall check the conditions of Definition 3.2. In fact,
(a) $\Phi_{n}(\theta) P_{n}(\theta)=P_{n+1}(\theta) \Phi_{n}(\theta) \quad n \geq 0, \quad \theta \in \Theta$.
which is true by hypothesis. On the other hand for $n<0$ we have:

$$
\begin{aligned}
& \Phi_{n}(\theta) P_{n}(\theta)=\left(\alpha_{\theta}\left(I-P_{0}(\theta)\right)+\alpha_{\theta}^{-1} P_{0}(\theta)\right) P_{0}(\theta)= \\
& P_{0}(\theta)\left(\alpha_{\theta}\left(I-P_{0}(\theta)\right)+\alpha_{\theta}^{-1} P_{0}(\theta)\right)=P_{n+1}(\theta) \Phi_{n}(\theta)
\end{aligned}
$$

(b) $\Phi_{n}(\theta): \mathcal{N}\left(P_{n}(\theta)\right) \rightarrow \mathcal{N}\left(P_{n+1}(\theta)\right)$ is an isomorphism for $n \geq 0$.

If $n<0$, then $\Phi_{n}(\theta): \mathcal{N}\left(P_{n}(\theta)\right) \rightarrow \mathcal{N}\left(P_{n+1}(\theta)\right)$ is also an isomorphism. In fact, if $x \in \mathcal{N}\left(P_{n}(\theta)\right)=\mathcal{N}\left(P_{0}(\theta)\right)$, then

$$
\Phi_{n}(\theta) x=\alpha_{\theta}\left(I-P_{0}(\theta)\right) x=\alpha_{\theta} x
$$

which is an isomorphism.
(c) $\left\|\Phi_{n, m}(\theta) P_{m}(\theta)\right\| \leq M \alpha_{\theta}^{n-m}\|x\|, \quad n \geq m \geq 0$.

If $n \geq 0>m$, then

$$
\begin{aligned}
\Phi_{n, m}(\theta) & =\Phi_{n, 0}(\theta) \Phi_{0, m}(\theta) \text { and } P_{m}(\theta)=P_{0}(\theta) \\
\Phi_{0, m}(\theta) P_{m}(\theta) & =\Phi_{-1}(\theta) \Phi_{-2}(\theta) \cdots \Phi_{m}(\theta) P_{0}(\theta) \\
& =\Phi_{-1}(\theta) \Phi_{-2}(\theta) \cdots \Phi_{m+1}(\theta) \alpha_{\theta}^{-1} P_{0}(\theta) \\
& =\alpha_{\theta}^{-m} P_{0}(\theta)
\end{aligned}
$$

Then

$$
\begin{aligned}
\left\|\Phi_{n, m}(\theta) P_{m}(\theta)\right\| & =\left\|\Phi_{n, 0}(\theta) \alpha_{\theta}^{-m} P_{0}(\theta) x\right\| \\
& \leq M_{\theta} \alpha_{\theta}^{n-m}\|x\| .
\end{aligned}
$$

If $0>n \geq m$, then

$$
\Phi_{n, m}(\theta) P_{m}(\theta) x=\alpha_{\theta}^{-m} P_{0}(\theta) x
$$

So

$$
\left\|\Phi_{n, m}(\theta) P_{m}(\theta) x\right\| \leq \alpha_{\theta}^{|n-m|}\|x\|
$$

(d) If $m>n \geq 0$, then

$$
\left\|\Phi_{n, m}(\theta) P_{m}(\theta) x\right\| \leq M_{\theta} \alpha_{\theta}^{m-n}\|x\|
$$

If $m \geq 0>n$, then we have the following

$$
\begin{aligned}
\Phi_{n, m}(\theta)\left(I-P_{m}(\theta)\right) x & =\Phi_{n, 0}(\theta) \Phi_{0, m}(\theta)\left(I-P_{m}(\theta)\right) x \\
& =\Phi_{0, n}^{-1}(\theta) \Phi_{0, m}(\theta)\left(I-P_{m}(\theta)\right) x \\
& =\alpha_{\theta}^{-n} \Phi_{0, m}(\theta)\left(I-P_{m}(\theta)\right) x .
\end{aligned}
$$

Therefore

$$
\begin{aligned}
\left\|\Phi_{n, m}(\theta)\left(I-P_{m}(\theta)\right) x\right\| & \leq \alpha_{\theta}^{-1}\left\|\Phi_{0, m}(\theta)\left(I-P_{m}(\theta)\right) x\right\| \\
& \leq M_{\theta} \alpha_{\theta}^{m-n}\|x\|
\end{aligned}
$$

If $0>m>n$, then

$$
\Phi_{n, m}(\theta)\left(I-P_{m}(\theta)\right) x=\alpha_{\theta}^{m-n}\left(I-P_{0}(\theta)\right) x
$$

So

$$
\left\|\Phi_{n, m}(\theta)\left(I-P_{m}(\theta)\right) x\right\| \leq M_{\theta} \alpha_{\theta}^{m-n}\|x\| .
$$

Theorem 3.4 Assume $\hat{\pi}: \mathcal{E} \times Z_{+} \rightarrow \mathcal{E}$ has pointwise discrete dichotomy with constants $M_{\theta}, \alpha_{\theta}<1$. If $\tilde{M}_{\theta}>M_{\theta}$ and $\alpha_{\theta}<\tilde{\alpha}_{\theta}<1$, then there exists $\mathcal{E}_{\theta}>0$ such that for any discrete skew-product $\hat{\pi}: \mathcal{E} \times Z_{+} \rightarrow \mathcal{E}$ with

$$
\sup _{n \geq 0}\left\|\Phi_{n}(\theta)-\Phi_{n}^{1}(\theta)\right\|<\mathcal{E}_{\theta}, \quad \text { for } \quad \theta \in \Theta
$$

has pointwise discrete dichotomy over $\Theta$, here $\hat{\pi}=\left(\Phi_{n}(\cdot), \sigma\right), \quad \hat{\pi}_{1}=\left(\Phi_{n}^{1}(\cdot), \sigma\right)$.

## 4 Roughness

In this section we shall give a necessary and sufficient condition for skew-product semiflow to have an exponential dichotomy over $\Theta$. Also, we will prove that the exponential dichotomy is not destroyed by small perturbation. In order to do so we will use the concept of discrete dichotomy introduced in the foregoing section.

We begin this section with a proposition on the relation between skew-product semiflow and skew-product sequence.

Proposition 4.1 If $\pi=(\Phi, \sigma)$ is a skew-product semi-flow on $\mathcal{E}$, then the mapping $\hat{\pi}: \mathcal{E} \times Z \rightarrow \mathcal{E}$ given by

$$
\begin{equation*}
\hat{\pi}(x, \theta, n):=\left(\Phi_{n}(\theta), \theta \cdot n\right) \tag{4.1}
\end{equation*}
$$

where

$$
\Phi_{n}(\theta)=\Phi(\theta \cdot n, 1)
$$

is a skew-product sequence.

Proof Follows directly from Definition 3.1
Remark 4.1 Even though $\Phi(\theta, t)$ is defined in principle only for $t \geq 0, \quad \Phi_{n}(\theta)=$ $\Phi(\theta \cdot n, 1)$ is well defined for all $n \in Z$ since $\sigma(\theta, t)=\theta \cdot t$ is defined for all $t \in \mathbb{R}$. This could be one of the advantage of using discrete skew-product.

Lemma 4.1 Let $\pi=(\Phi, \sigma)$ be a skew-product semi-flow on $\mathcal{E}$. Suppose that $\pi$ has an exponential dichotomy over $\Theta$ (with exponent $\beta$ and $k$ ). Then the skew-product sequence $\hat{\pi}$ given by (4.1) has uniformdiscrete dichotomy over $\Theta$, with constants $M=K$ and $\alpha=e^{-\beta}$.
Proof Let $\mathbf{P}$ be the projector associated to the exponential dichotomy of $\pi$ over $\Theta$. Define the family of projections $\left\{P_{n}(\theta)\right\}$ as follows:

$$
P_{n}(\theta):=P(\theta \cdot n), \text { for all } \theta \in \Theta \text { and } n \in Z
$$

we shall prove the condition (a) - (d) of Definitions 3.2 and 3.3.
In order to see (a), consider

$$
\begin{gathered}
\Phi_{n}(\theta) P_{n}(\theta)=\Phi(\theta \cdot n, 1) P(\theta \cdot n)= \\
P(\theta \cdot n \cdot 1) \Phi(\theta \cdot n, 1)=P(\theta \cdot(n+1)) \Phi_{n}(\theta) \\
=P_{n+1}(\theta) P_{n}(\theta), \quad n \in Z, \quad \theta \in \Theta
\end{gathered}
$$

(b) $\Phi_{n}(\theta): \mathcal{N}\left(P_{n}(\theta)\right) \rightarrow \mathcal{N}\left(P_{n+1}(\theta)\right)$ is an isomorphism.

In fact, from the Definition 2.5 we know that:

$$
\Phi(\theta \cdot n, 1): \mathcal{N}(P(\theta \cdot n)) \rightarrow \mathcal{N}(P(\theta \cdot(n+1))
$$

is an isomorphism.
Therefore

$$
\Phi_{n}(\theta): \mathcal{N}\left(P_{n}(\theta)\right) \rightarrow \mathcal{N}\left(P_{n+1}(\theta)\right)
$$

is an isomorphism.
(c) If we put $\Phi_{n, m}(\theta):=\Phi_{n-1}(\theta) \cdots \Phi_{m}(\theta), n>m$ and $\Phi_{m, m}(\theta):=I, \quad \theta \in \Theta$ Then

$$
\begin{gathered}
\left\|\Phi_{n, m}(\theta) P_{m}(\theta)\right\|=\|\Phi(\theta \cdot(n-1), 1) \cdots \Phi(\theta \cdot m, 1) \cdot P(\theta \cdot m)\| \\
=\|\Phi(\theta \cdot m, n-m) P(\theta \cdot m)\| \leq M e^{\beta(n-m)} \text { for all } \theta \in \Theta \text { and } n \geq m .
\end{gathered}
$$

Hence

$$
\left\|\Phi_{n, m}(\theta) P_{m}(\theta)\right\| \leq M \alpha^{n-m}, \text { for all } \theta \in \Theta \text { and } n \geq m
$$

In the same way we can prove (d), this means

$$
\left\|\Phi_{n, m}(\theta)\left(I-P_{m}(\theta)\right)\right\| \leq M \alpha^{m-n}, \quad m>n \text { and } \theta \in \Theta
$$

This completes the proof of the Theorem.
The following theorem give us a sufficient condition for the skew-product semiflow on $\mathcal{E}$ to have an exponential dichotomy over $\Theta$.

Theorem 4.1 Assume that $\pi=(\Phi, \sigma)$ is a skew-product semiflow on $\mathcal{E}$ and $M ; \alpha$ are positive constant such that $\alpha=\exp (-\beta)<1$. Consider

$$
\begin{equation*}
L:=\sup \{\|\Phi(\theta, t)\|: 0 \leq t \leq 1, \quad \theta \in \Theta\}<\infty \tag{4.2}
\end{equation*}
$$

If the skew-product sequence $\hat{\pi}$ given by (4.1) has uniform discrete dichotomy over $\Theta$ with constant $M$ and $\alpha$, then $\pi=(\Phi, \sigma)$ has exponential dichotomy with exponent $\beta$ and constant $K M$ where

$$
\begin{equation*}
K=\sup \left\{\|\Phi(\theta, t)\| e^{\beta t}: \quad 0 \leq t \leq 1, \quad \theta \in \Theta\right\} \tag{4.3}
\end{equation*}
$$

Proof Suppose that $\left\{P_{n}(\theta)\right\}$ is the family of projections associated with the uniform discrete dichotomy of $\hat{\pi}=\left(\Phi_{n}(\cdot), \sigma\right)$. Define $P: \mathcal{E} \rightarrow \mathcal{E}$ as follows:

$$
P(x, \theta):=(P(\theta) x, \theta) ; \quad P(\theta):=P_{0}(\theta) .
$$

Claim 1. $P_{k}(\theta)=P(\theta \cdot k)=P_{0}(\theta \cdot k)$ for all $\theta \in Z$ and $k \in Z$
In fact. Consider the skew-product sequence $\hat{\pi}_{k}: \mathcal{E} \times Z \rightarrow \mathcal{E}$ given by

$$
\begin{equation*}
\hat{\pi}_{k}(x, \theta, n):=\left(\Phi_{n}(\theta \cdot k), \theta \cdot n\right) \tag{4.4}
\end{equation*}
$$

The $\hat{\pi}$ has uniform discrete dichotomy with family of projections $\left\{P_{n}(\theta \cdot k)\right\}$ and $\Phi_{n}(\theta$. $k)=\Phi_{n+k}(\theta)$. Therefore $\hat{\pi}_{k}$ has uniform discrete dichotomy with family of projections $\left\{P_{n+k}(\theta)\right\}$. Then from the uniqueness we get that $P_{n+k}(\theta)=P_{n}(\theta \cdot k)$. In particular $P_{k}(\theta)=P_{n}(\theta . K)=P(\theta \cdot k)$
Claim 2. For all $n \geq K$ and $\theta \in \Theta$ we have the following:

$$
P(\theta \cdot n) \Phi(\theta \cdot k, n-k)=\Phi(\theta \cdot k, n-k) P(\theta \cdot k)
$$

In fact. It follows from the relation.

$$
P_{n}(\theta) \Phi_{n, k}(\theta)=\Phi_{k}(\theta) P_{k}(\theta)
$$

Claim 3. For any $t \geq 0$ and $\theta \in \Theta$ we have

$$
\|\Phi(\theta, t) P(\theta)\| \leq K M e^{-\beta t}
$$

In fact. Let $n \in Z_{+}$be such that $n \leq t \leq n+1$. Then

$$
\begin{aligned}
& \|\Phi(\theta, t) P(\theta)\|=\|\Phi(\theta, t+n-n) P(\theta)\| \\
& =\|\Phi(\theta \cdot n, t-n) \Phi(\theta, n) P(\theta)\| \leq\|\Phi(\theta \cdot n, t-n)\| M \alpha^{n} \\
& =\|\Phi(\theta \cdot n, t-n)\| M e^{-\beta n} \\
& \leq\|\Phi(\theta \cdot n, t-n)\| e^{\beta(t-h)} M e^{-\beta t} \leq K M e^{-\beta t} .
\end{aligned}
$$

Here we have used the fact that:

$$
\Phi(\theta, n)=\Phi_{n, 0}(\theta)=\Phi_{n-1}(\theta) \cdots \Phi_{0}(\theta)
$$

For $x \in \mathcal{N}(P(\theta)), \quad t \leq 0$ and $n \in Z$ with $n \leq t<n+1$. We shall define:

$$
\begin{equation*}
\Phi(\theta, t) x:=\Phi(\theta \cdot n, t-n) \Phi(\theta, n) x \tag{4.5}
\end{equation*}
$$

where

$$
\Phi(\theta, n)=\Phi_{n, 0}(\theta)=\left[\Phi_{0, n}(\theta)\right]^{-1}
$$

and

$$
\Phi_{0, n}(\theta): \mathcal{N}(P(\theta)) \rightarrow \mathcal{N}(P(\theta \cdot n))
$$

is an isomorphism.
Claim 4. For $t \leq 0, \quad \Phi(\theta, t): \mathcal{N}(P(\theta)) \rightarrow X$ defined by (4.5) satisfies the following:

$$
\| \Phi(\theta, t)\left(I-P(\theta) x\left\|\leq K M e^{\beta t}\right\| x \|\right.
$$

In fact.

$$
\begin{aligned}
\|\Phi(\theta, t)[I-P(\theta)] x\| & \leq\|\Phi(\theta \cdot n, t-n)\|\|\Phi(\theta, n)[I-P(\theta)] x\| \\
& \leq\|\Phi(\theta \cdot n, t-n)\| M \alpha^{-n}\|x\|=\|\Phi(\theta \cdot n, t-n)\| M e^{\beta n}\|x\| \\
& =\|\Phi(\theta \cdot n, t-n)\| e^{\beta n-\beta t} M e^{\beta t}\|x\| \\
& \leq\|\Phi(\theta \cdot n, t-n)\| e^{\beta(t-n)} M e^{\beta t}\|x\| K M e^{\beta t}
\end{aligned}
$$

Claim 5. For all $\theta \in \Theta \mathcal{B}^{+}(\theta)=\mathcal{R}(P(\theta))$.
In fact. From 2.3 we have that:

$$
\mathcal{B}^{+}(\theta)=\left\{x \in X: \sup _{t \geq 0}\|\Phi(\theta, t) x\|<\infty\right\}
$$

If $x \in \mathcal{R}(P(\theta))$, then $P(\theta) x=x$. So from Claim 3 we get

$$
\|\Phi(\theta, t) x\|=\|\Phi(\theta, t) P(\theta) x\| \leq K M e^{-\beta t}
$$

Hence $\Phi(\theta, \cdot) x$ is bounded for $t \geq 0$. This means that $\mathcal{R}(P(\theta)) \subset \mathcal{B}^{+}(\theta), \quad \forall \theta \in \Theta$. Now we shall prove the following implication:

$$
\text { If } x \in \mathcal{B}^{+}(\theta) \text {, then } x \in \mathcal{R}(P(\theta))
$$

In fact. If $x \notin \mathcal{R}(P(\theta))$, then

$$
\Phi(\theta, t) x=\Phi(\theta, t) P(\theta) x+\Phi(\theta, t)(I-P(\theta)) x
$$

We already know that $\Phi(\theta, \cdot) P(\theta) x$ is bounded for $t \geq 0$. Hence, we only need to prove that $\Phi(\theta, t)(I-P(\theta)) x$ is unbounded for $t \geq 0$. In fact. We have the following:

$$
\|\Phi(\theta \cdot n,-n)(I-P(\theta \cdot n)) v\| \leq M e^{-\beta n}\|v\|, \quad n \geq 0
$$

where:

$$
\Phi(\theta, n)=\Phi_{n, 0}(\theta): \mathcal{N}(P(\theta)) \rightarrow \mathcal{N}(P(\theta \cdot n))
$$

is an isomorphism with inverse:

$$
\begin{equation*}
\Phi(\theta \cdot n,-n):=\Phi_{0, n}(\theta): \mathcal{N}(P(\theta \cdot n) \rightarrow \mathcal{N}(P(\theta)) \tag{4.6}
\end{equation*}
$$

Putting $\quad v=\Phi(\theta, n)(I-P(\theta) x$ we get:

$$
\begin{aligned}
\|(I-P(\theta)) x\| & =\| \Phi(\theta \cdot n,-n)(I-P(\theta \cdot n)) \Phi(\theta, n)(I-P(\theta) x \| \\
& \leq M e^{-\beta n}\|\Phi(\theta, n)(I-P(\theta)) x\|
\end{aligned}
$$

then

$$
\|\Phi(\theta, n)(I-P(\theta)) x\| \geq M^{-1} e^{\beta n}\|(I-P(\theta)) x\|
$$

Since $(I-P(\theta)) x \neq 0$, then $\|\Phi(\theta, n)(I-P(\theta)) x\| \rightarrow \infty$ as $n \rightarrow \infty$. So $x \notin \mathcal{B}^{+}(\theta)$.
Claim 6.

$$
\Phi(\theta, t) \mathcal{R}(P(\theta)) \subseteq \mathcal{R}(P(\theta \cdot t)), \quad t \geq 0
$$

i.e.,

$$
\Phi(\theta, t) \mathcal{R}(I-P(\theta)) \subseteq \mathcal{R}(I-P(\theta \cdot t))
$$

i.e.,

$$
\Phi(\theta, t) \mathcal{N}(P(\theta)) \subseteq \mathcal{N}(P(\theta \cdot t)
$$

In fact. $x \in \mathcal{R}(P(\theta)) \Longleftrightarrow P(\theta) x=x$. Then $\Phi(\theta, t) x$ is bounded for $t \geq 0$ and

$$
\Phi(\theta, s) x=\Phi(\theta \cdot t, s-t) \Phi(\theta, t) x, \quad s \geq t \geq 0
$$

Putting $r=s-t \geq$ and $z=\Phi(\theta, t) x$. We get that $\Phi(\theta \cdot t, r) z$ is bounded for $r \geq 0$. Then from Claim 5 we get that

$$
\Phi(\theta, t) x=z \in \mathcal{R}(P(\theta \cdot t)) .
$$

## Claim 7.

$$
\Phi(\theta, t): \mathcal{N}(P(\theta)) \rightarrow \mathcal{N}(P)(\theta \cdot t)), \quad t \geq 0 \text { is one to one }
$$

In fact. For the purpose of contradiction let us suppose that there is $x \neq 0$ such that

$$
x \in \mathcal{N}(P(\theta)) \text { and } \Phi(\theta, t) x=0
$$

Consider $n \geq t$ and $\Phi(\theta, n) x=\Phi(\theta \cdot t, n-t) \Phi(\theta, t) x=0$. Hence

$$
\Phi(\theta, n)=\Phi_{n, 0}(\theta): \mathcal{N}(P(\theta)) \rightarrow \mathcal{N}(P(\theta \cdot n))
$$

is not one to one; which is a contradiction with the uniform discrete dichotomy of $\hat{\pi}$.

## Claim 8.

$$
\Phi(\theta \cdot t, s-t): \Phi(\theta, t) \mathcal{N}(P(\theta)) \rightarrow \mathcal{N}(P(\theta \cdot s)) s \geq t \geq 0 \text { is one to one }
$$

In fact. Suppose thre is $y \in \Phi(\theta, t) \mathcal{N}(P(\theta))$ such that $y \neq 0$ and $\Phi(\theta \cdot t, s-t) y=0$. Let $n \geq s$ and $x \in \mathcal{N}(p(\theta))$ such that $\Phi(\theta, t) x=y$ Then

$$
\begin{aligned}
\Phi(\theta, n) x & =\Phi(\theta \cdot s, n-s) \Phi(\theta, s) x \\
& =\Phi(\theta \cdot s, n-s) \Phi(\theta \cdot t, s-t) \Phi(\theta, t) x \\
& =\Phi(\theta \cdot s, n-s) \Phi(\theta \cdot t, s-t) y=0
\end{aligned}
$$

Hence, $\Phi(\theta, n) x=0$ and $x \in \mathcal{N}(P(\theta))$ with $x \neq 0$. Which is a contradiction with the fact that $\Phi(\theta, n): \mathcal{N}(P(\theta)) \rightarrow \mathcal{N}(P(\theta \cdot n))$ is an isomorphism.

Let $x \in X$ and $t_{1}<0$. If there exists a unique $y \in X$ such that $\Phi\left(\theta \cdot t_{1},-t_{1}\right) y=x$ we shall define: $\Phi\left(\theta, t_{1}\right) x:=y$. In particular, if $x \in \mathcal{N}(P(\theta))$ and $t \leq 0$ we have the following

$$
\Phi(\theta, t) x:=\Phi(\theta \cdot n, t-n) \Phi(\theta, n) x=y
$$

where $0 \leq t-n<1$ and $\Phi(\theta, n): \mathcal{N}(P(\theta)) \rightarrow \mathcal{N}(P(\theta \cdot n))$ is an isomorphims. Indeed

$$
\begin{aligned}
\Phi(\theta \cdot t,-t) y & =\Phi(\theta \cdot t,-t) \Phi(\theta \cdot n, t-n) \Phi(\theta, n) x \\
& =\Phi(\theta \cdot n \cdot(t-n),-t) \Phi(\theta \cdot n, t-n) \Phi(\theta, n) x \\
& =\Phi(\theta \cdot n, t-n-t) \Phi(\theta, n) x \\
& =\Phi(\theta \cdot n,-n) \Phi(\theta, n) x=x
\end{aligned}
$$

Claim 9. $\quad \mathcal{N}(P(\theta))=\mathcal{B}_{u}^{-}(\theta)$.
In fact. From 2.3 we have that

$$
\mathcal{B}_{u}^{-}(\theta)=\{x \in X: \Phi(\theta, t) x \text { is well defined and bounded for } t \leq 0\}
$$

If $x \in \mathcal{N}(P(\theta))$, then

$$
\Phi(\theta, t) x:=\Phi(\theta \cdot n, t-n) \Phi(\theta, n)[I-P(\theta)] x
$$

and from Claim 4 we get that

$$
\|\Phi(\theta, t) x\| \leq K M e^{\beta t}\|x\|, \quad t \leq 0
$$

Which implies that $\Phi(\theta, t) x$ is bounded for $t \leq 0$. Suppose that $x \notin \mathcal{N}(P(\theta))$, then putting

$$
y=P(\theta) x \text { and } z=(I-P(\theta)) x \neq 0
$$

we have two cases:
a) If $\Phi(\theta, t) x$ is not well defined, then $x \notin \mathcal{B}_{u}^{-}(\theta)$.
b) If $\Phi(\theta, t) x$ is well defined for $t \leq 0$, then $\Phi(\theta, t) z$ and $\Phi(\theta, t) x$ are well defined by using (4.5). Hence $\Phi(\theta, t) y$ is also well defined for $t \leq 0$. Therefore:

$$
\Phi(\theta, t) x=\Phi(\theta, t) z+\Phi(\theta, t) y
$$

We shall prove that $\Phi(\theta, t) y=\Phi(\theta, t) P(\theta) x$ is unbounded for $t \leq 0$. In fact

$$
\begin{aligned}
P(\theta) x & =\Phi(\theta \cdot n,-n) \Phi(\theta, n) P(\theta) x \\
& =\Phi(\theta \cdot n,-n) P(\theta \cdot n) \Phi(\theta, n) P(\theta) x, \quad n \leq 0
\end{aligned}
$$

From Claim 3 we get that:

$$
\|P(\theta) x\| \leq K M e^{\beta n}\|\Phi(\theta, n) P(\theta) x\|
$$

So

$$
\|\Phi(\theta, n) P(\theta) x\| \geq(K M)^{-1}\|P(\theta) x\| e^{-\beta n} \quad \rightarrow \infty \text { as } n \rightarrow-\infty
$$

Since $\Phi(\theta, t)(I-P(\theta)) x$ is bounded, we get that $\Phi(\theta, t) x$ is unbounded.
Thus

$$
\mathcal{N}(P(\theta))=\mathcal{B}_{u}^{-}(\theta), \text { for all } \theta \in \Theta
$$

Claim 10.

$$
\Phi(\theta, t) \mathcal{N}(P(\theta)) \subset \mathcal{N}(P(\theta \cdot t)), \text { for } t \leq 0
$$

In fact. $x \in \mathcal{N}(P(\theta)) \Longleftrightarrow x \in \mathcal{B}_{u}^{-}(\theta)$. Therefore, for all $x \in \mathcal{N}(P(\theta))$ and $s \leq t \leq 0$ we have

$$
\Phi(\theta, s) x=\Phi(\theta \cdot t, s-t) \Phi(\theta, t) x
$$

So $\Phi(\theta \cdot t, s-t) \Phi(\theta, t) x$ is bounded for $s \leq t$. Then

$$
\Phi(\theta, t) x \in \mathcal{B}_{u}^{-}(\theta \cdot t)=\mathcal{N}(P(\theta \cdot t))
$$

## Claim 11.

$$
\Phi(\theta, t): \mathcal{N}(P(\theta)) \rightarrow \mathcal{N}(P(\theta \cdot t)), \quad t \geq 0
$$

is an isomorphism. In fact, for all $x \in \mathcal{N}(P(\theta))$ we have

$$
\Phi(\theta, t) x=\Phi(\theta \cdot n, t-n) \Phi(\theta, n) x
$$

Where $0 \leq n \leq t<n+1$. Since $\Phi(\theta, n) \mathcal{N}(P(\theta))=\mathcal{N}(P(\theta \cdot n))$ and using Claim 6 we get

$$
\Phi(\theta, t) \mathcal{N}(P(\theta)=\Phi(\theta \cdot n, t-n) \mathcal{N}(P(\theta \cdot n)) \subseteq \mathcal{N}(P(\theta \cdot n \cdot(t-n)))=\mathcal{N}(P(\theta \cdot t))
$$

Now, if $\quad x \in \mathcal{N}(P(\theta \cdot t))=\mathcal{B}_{u}^{-}(\theta \cdot t)$, then

$$
\Phi(\theta \cdot t, r+s) x=\Phi(\theta \cdot t \cdot r, s) \Phi(\theta \cdot t, r) x \text { for all } r, s \in \mathbb{R}
$$

So

$$
\Phi(\theta, t) \Phi(\theta \cdot t,-t) x=x
$$

If $y=\Phi(\theta \cdot t,-t) x \Rightarrow y \in \mathcal{N}(P(\theta)) \quad$ and $\quad \Phi(\theta, t) y=x$.
Hence, from Claim 7 we get that

$$
\Phi(\theta, t): \mathcal{N}(P(\theta)) \rightarrow \mathcal{N}(P(\theta \cdot t)) \text { is an isomorphism. }
$$

Claim 12.

$$
\Phi(\theta, t) P(\theta)=P(\theta \cdot t) \Phi(\theta, t), \quad \theta \in \Theta, \quad t \geq 0
$$

In fact. Clearly that for $t \geq 0$ we get

$$
\begin{gathered}
\Phi(\theta, t) P(\theta) x=P(\theta \cdot t) \Phi(\theta, t) x, \quad x \in \mathcal{N}(P(\theta) \\
\Phi(\theta, t)(I-P(\theta)) x=(I-P(\theta \cdot t)) \Phi(\theta, t) x, \quad x \in \mathcal{R}(P(\theta))
\end{gathered}
$$

i.e.,

$$
\Phi(\theta, t) P(\theta) x=P(\theta \cdot t) \Phi(\theta, t) x, \quad x \in \mathcal{R}(P(\theta))
$$

It follows that

$$
\Phi(\theta, t) P(\theta)=P(\theta \cdot t) \Phi(\theta, t), \text { for all } \theta \in \Theta \text { and } t \geq 0
$$

This completes the Proof of the Theorem
The following Theorem also gives a sufficient condition for the existence of exponential dichotomy for skew-product semiflow.

Theorem 4.2 Let $\pi=(\Phi, \sigma)$ be a skew-product semiflow on $\mathcal{E}$. Assume there is an invariant projector $\mathbf{P}$ and for each $\theta \in \Theta$ there are constants $k_{\theta} \geq 1, \beta_{\theta}>0$ such that (1) $\Phi(\theta, t): \mathcal{N}(P(\theta)) \rightarrow \mathcal{N}(P(\theta \cdot t)), \quad t \geq 0$ is an isomorphism with inverse:

$$
\Phi(\theta \cdot t,-t): \mathcal{N}(P(\theta \cdot t)) \rightarrow \mathcal{N}(P(\theta)), \quad t \geq 0
$$

(2) $\|\Phi(\theta, t) P(\theta)\| \leq k_{\theta} e^{-\beta_{\theta} t}, \quad t \geq 0$
(3) $\| \Phi(\theta, t)\left(I-P(\theta) \| \leq k_{\theta} e^{\beta_{\theta} t}, \quad t \leq 0\right.$.

If one of the following conditions holds:
(4) there exists $1<\delta<2$ such that

$$
\|\Phi(\theta, 1) x\| \geq \delta\|x\|, \quad \theta \in \Theta, x \in X
$$

(5) there exist $0<\delta_{1}<\delta_{2}<1$ such that

$$
\delta_{1}\|x\| \leq\|\Phi(\theta, 1) x\| \leq \delta_{2}\|x\|, \quad \theta \in \Theta, \quad x \in X
$$

then $\pi$ has an exponential dichotomy over $\Theta$.

Proof It follows from Lemmas 3.4 and 3.5.
The following Theorem says that the exponential dichotomy of the skew-product semiflow is not destroyed by small perturbation.

Theorem 4.3 Suppose $\pi=(\Phi, \sigma)$ is a skew-product semiflow on $\mathcal{E}$ which has a exponential dichotomy (with exponent $\beta$ and constant $M$ ). If

$$
L=\sup \{\|\Phi(\theta, t)\|: \quad 0 \leq t \leq 1, \quad \theta \in \Theta\}
$$

and $M e^{-\beta}<e^{-\beta_{1}}, M_{1}>M$, then there exists $\epsilon=\epsilon\left(\beta, \beta_{1}, M, M_{1}, L\right)>0$ such that any skew-product semiflow $\tilde{\pi}=(\Psi, \sigma)$ on $\mathcal{E}$ satisfying

$$
\sup \{\|\Phi(\theta, t)-\Psi(\theta, t)\|: \quad 0 \leq t \leq 1, \quad \theta \in \Theta\} \leq \epsilon
$$

has exponential dichotomy with exponent $\beta_{1}$ and constant $M_{1}$.
Proof Let $\hat{\pi}, \check{\pi}: \mathcal{E} \times Z \rightarrow \mathcal{E}$ be the skew-product sequences given by:

$$
\hat{\pi}(x, \theta, n):=\left(\Phi_{n}(\theta) x, \theta \cdot n\right) ; \quad \check{\pi}(x, \theta, n):=\left(\Psi_{n}(\theta) x, \theta \cdot n\right)
$$

where

$$
\Phi_{n}(\theta):=\left(\Phi(\theta \cdot n, 1) ; \quad \Psi_{n}(\theta):=\Psi(\theta \cdot n, 1)\right.
$$

Clearly $\hat{\pi}$ has uniform discrete dichotomy with constants $M$ and $\alpha=e^{-\beta}$. Consider for $t \geq 0$ the following:

$$
\begin{aligned}
\|\Phi(\theta, t)-\Psi(\theta, t)\| & =\|[\Phi(\theta \cdot k, t-k)-\Psi(\theta \cdot k, t-k)] \Psi(\theta, k) \\
& -\Phi(\theta \cdot k, t-k)[\Psi(\theta, k)-\Psi(\theta, k)] \|
\end{aligned}
$$

for $k \leq t \leq k+1, \quad k \geq 0$. Since $0 \leq t-k \leq 1$, we get:

$$
\|\Phi(\theta \cdot k, t-k)-\Psi(\theta \cdot k, t-k)\| \leq \epsilon
$$

and

$$
\begin{aligned}
\|\Psi(\theta, k)\| & =\|\Psi(\theta \cdot(k-1), 1) \Psi(\theta \cdot(k-2), 1) \cdots \Psi(\theta, 1)\| \\
& =\left\|\Psi_{k-1}(\theta) \Psi_{k-2}(\theta) \cdots \Psi_{0}(\theta)\right\| \leq(L+\epsilon)^{k} .
\end{aligned}
$$

Hence

$$
\|\Phi(\theta, t)-\Psi(\theta, t)\| \leq(L+\epsilon)^{k} \cdot \epsilon+L\|\Psi(\theta, k)-\Phi(\theta, k)\|
$$

In the same way we get that
$\|\Phi(\theta, k)-\Psi(\theta, k)\|=$
$\|[\Phi(\theta .(k-1), 1)-\Psi(\theta \cdot(k-1), 1)] \Psi(\theta, k-1)-\Phi(\theta .(k-1), 1)[\Psi(\theta, k-1)-\Phi(\theta, k-1] \|$

$$
\leq(L+\epsilon)^{k-1} \epsilon+L\|\Psi(\theta, k-1)-\Phi(\theta, k-1)\|
$$

Therefore

$$
\|\Phi(\theta, t)-\Psi(\theta, t)\| \leq\left[(L+\epsilon)^{K}+L(L+\epsilon)^{K-1}+\cdots+L^{K}\right] \epsilon
$$

Then

$$
\|\Psi(\theta, t)-\Psi(\theta, t)\| \leq C \epsilon, \quad 0 \leq t \leq 1, \quad \theta \in \Theta
$$

Then

$$
\|\Phi(\theta \cdot n, 1)-\Psi(\theta \cdot n, 1)\|<C \epsilon, \quad \theta \in \Theta
$$

i.e.,

$$
\left\|\Phi_{n}(\theta)-\Psi_{n}(\theta)\right\|<C \epsilon, \quad \theta \in \Theta
$$

Since $M_{1}>M$ and $1>\alpha_{1}=e^{-\beta_{1}}>e^{\beta}=\alpha$, then for $\epsilon$ small enough Corolary 3.6 implies that $\check{\pi}$ has uniform discrete dichotomy over $\Theta$. Hence, it follows from Theorem 4.1 that $\tilde{\pi}=(\Psi, \sigma)$ has exponential dichotomy with constant $K M_{1}$ and exponent $\beta_{1}$.

## 5 Applications

In this section we shall consider a linear time dependent differential equation which generates a linear skew-product semiflow on the trivial Banach bundle $\mathcal{E}=X \times \Theta$, where $X$ is a Banach space and $\Theta$ is a compact topological Hausdorff space.

Consider the following linear time dependent system

$$
\begin{equation*}
\dot{x}(t)=A(\theta \cdot t) x(t), \quad t>0 \tag{5.1}
\end{equation*}
$$

where $A(\theta \cdot t)=A+B(\theta \cdot t), \quad A$ is the infinitesimal generator of the strongly continuous semigroup $\{T(t)\}_{t \geq 0} ; \quad \sigma(\theta, t)=\theta \cdot t$ is a flow on $\Theta$ and $B(\theta) \in L(X), \quad t \in \mathbb{R}$.

Lemma 5.1 If $B(\cdot): \Theta \rightarrow L(X)$ is strongly continuous, then the set $\{\|B(\theta)\|: \theta \in \Theta\}$ is bounded.

Proof Consider the following sets

$$
H=\{\|B(\theta)\|: \theta \in \Theta\}, \quad H(x)=\{\|B(\theta) x\|: \theta \in \Theta\}
$$

Since $\theta \rightarrow B(\theta) x$ is continuous and $\Theta$ is compact, then for each $x \in X$ we get that $H(x)$ is bounded. Hence, by the Uniform Boundedness Theorem we obtain that $H$ is bounded.

Lemma 5.2 If $B(\cdot): \Theta \rightarrow L(X)$ is strongly continuous and $x(\cdot): \mathbb{R} \rightarrow X$ is a continuous function, then for each $\theta \in \Theta$ the mapping $t \rightarrow B(\theta \cdot t) x(t)$ is continuous.

Proof Fix $t \in \mathbb{R}$. Then

$$
\begin{gathered}
\|B(\theta(t+h)) x(t+h)-B(\theta \cdot t) x(t)\| \\
=\|B(\theta(t+h))[x(t+h)-x(t)]-[B(\theta \cdot(t+h))-B(\theta \cdot t)] x(t)\| \\
\leq L\|x(t+h)-x(t)\|+\|[B(\theta \cdot(t+h))-B(\theta \cdot t)] x(t)\|
\end{gathered}
$$

where $L=\sup \{\|B(\theta)\|: \theta \in \Theta\} \quad$ and $\quad \theta \cdot t \in \Theta$ for $t \in \mathbb{R}$.
Therefore

$$
\|B(\theta(t+h)) x(t+h)-B(\theta \cdot t) x(t)\| \rightarrow 0, \text { as } h \rightarrow 0
$$

To be precise in which sense the equation (5.1) generates a linear skew-product semiflow, we shall consider the following family of integral differential equations:

$$
\begin{equation*}
x(t)=T(t) x_{0}+\int_{0}^{t} T(t-s) B(\theta \cdot s) x(s) d s . \quad t \geq 0 \quad \theta \in \Theta . \tag{5.2}
\end{equation*}
$$

Definition 5.1 (Mild Solution). A solution $x(t)=x(t, \theta)$ of the equation (5.2) is called Mild Solution of (5.1).

The proof of the following theorem can be found in Chow and Leiva [3].
Theorem 5.1 Let $A$ be the infinitesimal generator of an strongly continuous semigroup $\{T(t)\}_{t \geq 0}$ on $X$ and $B(\cdot): \Theta \rightarrow L(X)$ is also strongly continuous. Then for each $\theta \in \Theta$ and $x_{0} \in X$ the problem

$$
\begin{equation*}
\dot{x}(t)=A(\theta \cdot t) x=(A+B(\theta \cdot t)) x(t) ; \quad x(0)=x_{0} \tag{5.3}
\end{equation*}
$$

has a unique mild solution $\Phi(\theta, t) x_{0}$ given by

$$
\begin{equation*}
\Phi(\theta, t) x_{0}=T(t) x_{0}+\int_{0}^{t} T(t-s) B(\theta \cdot s) \Phi(\theta, s) x_{0} d s \tag{5.4}
\end{equation*}
$$

If

$$
\|T(t)\| \leq M e^{W t}, \quad t \geq 0
$$

then

$$
\begin{equation*}
\|\Phi(\theta, t)\| \leq M e^{(W+L M) t}, \quad t \geq 0 \tag{5.5}
\end{equation*}
$$

where

$$
L=\sup \{\|B(\theta)\|: \theta \in \Theta\}
$$

Moreover, the mapping $\pi: \mathcal{E} \times \mathbb{R}_{+} \rightarrow \mathcal{E}$ given by

$$
\begin{equation*}
\pi(x, \theta, t)=(\Phi(\theta, t) x, \theta \cdot t) \tag{5.6}
\end{equation*}
$$

is a linear skew-product semiflow on $\mathcal{E}=X \times \Theta$.

Theorem 5.2 Assume $A_{0}$ is the infinitesimal generator of a strongly continuous semigroup $\left\{T_{0}(t)\right\}_{t \geq 0}$, the mapping $\theta \rightarrow A(\theta)-A_{0}: \Theta \rightarrow L(X)$ is strongly continuous and the equation $\quad \dot{x}(t)=A(\theta \cdot t) x(t) \quad$ has exponential dichotomy over $\Theta$ with exponent $\beta>0$ and constant $M$. If $0<\beta_{1}<\beta$ and $M_{1}>M$, then there exist $\epsilon>0$ such that for any mapping $\theta \rightarrow B(\theta): \Theta \rightarrow L(X)$ strongly continuous and $\|B(\theta)\|<\epsilon, \theta \in \Theta$ the equation

$$
\dot{x}(t)=(A(\theta \cdot t)+B(\theta \cdot t)) x(t)
$$

has exponential dichotomy over $\Theta$ with exponent $\beta_{1}$ and constant $M_{1}$.
Proof From Theorem 5.1 we get that the equations

$$
\dot{x}(t)=A(\theta \cdot t) x(t)=\left(A_{0}+\left(A(\theta \cdot t)-A_{0}\right)\right) x(t)
$$

and

$$
\dot{x}(t)=(A(\theta \cdot t)+B(\theta \cdot t)) x(t)=\left(A_{0}+\left(A(\theta \cdot t)-A_{0}+B(\theta \cdot t)\right)\right) x(t)
$$

generate respectively the skew-product semiflows:

$$
\pi=(\Phi, \sigma), \tilde{\pi}=(\Psi, \sigma): \mathcal{E} \times \mathbb{R}_{+} \rightarrow \mathcal{E}, \quad \mathcal{E}=X \times \Theta
$$

given by:

$$
\begin{gathered}
\Phi(\theta, t) x=T_{0}(t) x+\int_{0}^{t} T_{0}(t-s)\left(A(\theta \cdot s)-A_{0}\right) \Phi(\theta, s) x d s \\
\Psi(\theta, t) x=T_{0}(t) x+\int_{0}^{t} T_{0}(t-s)\left(A(\theta \cdot s) B(\theta \cdot s)-A_{0}\right) \Psi(\theta, s) x d s
\end{gathered}
$$

Then

$$
\|\Phi(\theta, t) x\| \leq M_{0}\|x\|+\int_{0}^{t} M_{2}\|\Phi(\theta, s) x\| d s, \quad 0 \leq t \leq 1
$$

where

$$
M_{0}=\sup \left\{\left\|T_{0}(t)\right\|: \quad 0 \leq s \leq 1\right\}, \quad M_{2}=M_{0} \sup \left\{\left\|A(\theta)-A_{0}\right\|: \quad \theta \in \Theta\right\}
$$

From Gronwall's Lemma we get $\|\Phi(\theta, t)\| \leq M_{0} e^{M_{2}}=L_{1}, \quad 0 \leq t \leq 1, \quad \theta \in \Theta$. In the same way we get that $\|\Psi(\theta, t)\| \leq L_{2}, \quad 0 \leq t \leq 1, \quad \theta \in \Theta$. On the other hand, we have the followig

$$
\begin{aligned}
\|\Phi(\theta, t) x-\Psi(\theta, t) x\| & \leq \int_{0}^{t} M_{2}\|\Phi(\theta, s) x-\Psi(\theta, s) x\| d s+\int_{0}^{t} M_{0} L_{2}\|B(\theta \cdot s) x\| d s \\
& \leq M_{0} L_{2} \epsilon\|x\|+M_{2} \int_{0}^{t}\|\Phi(\theta, s) x-\Psi(\theta, s) x\| d s, \quad 0 \leq t \leq 1
\end{aligned}
$$

Hence, from Gronwall's Lemma we get

$$
\|\Phi(\theta, t)-\Psi(\theta, t)\| \leq M_{0} L_{2} \epsilon e^{M_{2} t} \leq M_{0} L_{2} e^{M_{2}} \epsilon, \quad 0 \leq t \leq 1
$$

From Theorem 4.1 we get that the equation $\dot{x}(t)=(A(\theta \cdot t)+B(\theta \cdot t)) x(t)$ has exponential dichotomy over $\Theta$ with exponent $\beta_{1}$ and constant $M_{1}$.

Theorem 5.3 Let $-A$ be a sectorial operator with a sector

$$
S_{-a, \phi}=\{\lambda \in C: \phi \leq|\arg (\lambda+a)| \leq \pi, \quad \lambda \neq-a\} \subset \rho(-A)
$$

where $\rho(-A)$ is the resolvent set of $-A$ and $a \in \mathbb{R}, \quad 0<\phi<\pi / 2$. Suppose that the spectrum $\sigma(A)$ does not intersect the strip region $\{\lambda \in C: \beta \leq R e \lambda \leq \alpha\}$, where $\beta<\lambda<a$. If for all $x \in X$ the mapping $\theta \rightarrow B(\theta) x: \Theta \rightarrow L(X)$ is continuous and $\sup \{\|B(\theta)\|: \theta \in \Theta\}$ is small enough, then the equation

$$
\begin{equation*}
\dot{x}(t)=(A+B(\theta \cdot t)) x(t) \tag{5.7}
\end{equation*}
$$

has an exponential dichotomy over $\theta \in \Theta$.
Proof It follows from Theorem 2.3 in [33] and Theorem 1.5.2 in [11] that the equation $\dot{x}=A x$ has an exponential dichotomy over $\Theta$. Finally, we can apply Theorem 5.2 to the perturbed equation (5.7).

Now, consider the following family of functional differential equations which is more general than equation (1.6) studied in [18].

$$
\begin{equation*}
\dot{x}(t)=L(\theta \cdot t) x_{t}, \quad \theta \in \Theta \tag{5.8}
\end{equation*}
$$

where $\theta \cdot t$ is a flow on the compact Hausdorff set $\Theta$ which depend continuously on $\theta$ uniformly on compact interval of the time $t ; x_{t}$ denotes the function $s \rightarrow x(t+s), \quad-r \leq$ $s \leq 0$. We assume that $\left.L(\theta): C=C\left[[-r, 0], \mathbb{R}^{n}\right]\right] \rightarrow \mathbb{R}^{n}$ is linear and bounded operator, and for all $\phi \in C$ the mapping $\theta \rightarrow L(\theta) \phi$ is continuous.

Lemma 5.3 If the equation (5.8) has an exponential dichotomy on $\Theta$ and $\theta \rightarrow M(\theta) \in$ $L\left(C, \mathbb{R}^{n}\right)$ is strongly continuous. If $\|M(\theta)\|$ is small enough, then the equation

$$
\begin{equation*}
\dot{x}(t)=(L(\theta \cdot t)+M(\theta \cdot t)) x_{t}, \quad \theta \in \Theta \tag{5.9}
\end{equation*}
$$

has an exponential dichotomy over $\Theta$.
Proof Using the same idea of [18] we can write the equation (5.9) in the following abstract way

$$
\begin{equation*}
\dot{z}(t)=(A(\theta \cdot t)+B(\theta \cdot t)) z(t), \quad \theta \in \Theta, \quad z \in Z \tag{5.10}
\end{equation*}
$$

where $Z=\mathbb{R}^{n} \times L_{2}^{n}[-r, 0]$ is a Hilbert space with the inner product

$$
<\cdot, \cdot>_{z}=<\cdot, \cdot>_{\mathbb{R}^{n}}+<\cdot, \cdot>_{L_{2}}
$$

The operators $A(\theta)$ are defined on a commun domain $D$ given by

$$
D=\operatorname{Dom}(A(\theta))=\left\{(v, \phi) \in Z: \phi \in W^{1,2}[-r, 0], \quad \phi(0)=v\right\}, \quad \theta \in \Theta
$$

by

$$
A(\theta)(v, \phi):=(L(\theta) \phi, \dot{\phi})
$$

and $B(\theta)$ is a linear bounded operator from $Z$ to $Z$ given by

$$
B(\theta) z=(M(\theta) \phi, 0), \quad z=(v, \phi) \in Z
$$

So, for all $z \in Z$ mapping $\theta \rightarrow B(\theta) z: \Theta \rightarrow L(Z)$ is continuous and $\|B(\theta)\|$ is small enough. On the other hand, if we consider the operator

$$
A_{0}(v, \phi)=(0, \dot{\phi}), \quad(v, \phi) \in D
$$

The operator $A_{0}$ is the infinitesimal generator of a strongly continuous semigroup $\left\{T_{0}(t)\right\}_{t \geq 0}$ and the mapping $\theta \rightarrow A(\theta)-A_{0}: \Theta \rightarrow L(Z)$ is strongly continuous. Now, we finish the proof by applying the Theorem 5.2 to the equation (5.10).

## References

[1] A.B. Antonevich (1985), "Two methods for investigating the invertibility of operators from C* -algebras generated by dynamical systems", Math. USSR-Sbornik vol.52, pp. 1-20.
[2] A.Carvalho (1987), "Perturbations of exponential dichotomies for evolution operators", USP-Brazil, resume of Master Tesis.
[3] S.N.Chow and H.Leiva (1993), "Dynamical Spectrum for Time Dependent Linear Systems in Banach Spaces", (to appear in JJIAM).
[4] S.N.Chow and H.Leiva (1994), "Dynamical Spectrum for Skew-Product Flow in Banach Spaces", (to appear in J. Integral Differential Equations).
[5] S.N. Chow and K.Lu (1988), "Invariant manifolds for flows in Banach spaces", J. Differential. Equations, 74:285-317.
[6] W.A. Coppel (1978), "Dichotommies in stability theory", Lect. Notes in Math, vol.629, Springer Verlag, New York.
[7] J.L.Daleckii and M.G.Krein (1974), "Stability of solutions of differential equations in Banach space", Trans. of Math. Monographs, vol.43, AMS.
[8] J.K.Hale (1984), "Introduction to dynamical bifurcation", Lect. Notes in Math, vol.1057, pp. 107-151.
[9] J.Hale and X.B.Lin (1986), "Heteroclinic orbits for retarded functional differential equations", J.Differential Equations, vol.65, pp.175-202.
[10] J.K.Hale (1988), "Asymtotic Behavior of dissipative systems", Math.Surveys and Monographs, Vol.25, Amer.Soc., Providence, R.I.
[11] D.Henry (1981), "Geometric theory of semilinear parabolic equations", Springer New York.
[12] J.S.Kolesov (1979), "On the stability of solutions of linear differential equations of parabolic type with almost periodic coefficients", Trans.Moscow Math. Soc, Issue 2.
[13] Y.D.Latushkin and A.M.Stepin, "Linear skew-product flows and semigroups of weighted composition operators", (to appear).
[14] Y.D.Latushkin and A.M.Stepin (1991), "Weighted translations operators and linear extensions of dynamical systems", Russian Math. Surveys, vol.46, pp. 95-165.
[15] N.Levinson (1946), "The asymptotic behavior of system of linear differential equations", Amer.J.Math. vol.68, pp.1-6.
[16] X.B.Lin (1984), "Exponential dichotomies and homoclinic orbits in functional Differential Equations", LCDS Report No. 84-17, Brown University, J. Differential Equations, in press.
[17] X.B.Lin (1991), "Exponential dichotomies in intermediate spaces with applications to a diffusively perturbed predator-prey", Depatment of Mathematics, North Carolina State University.
[18] M.P.Lizana (1992), "Exponential dichotomy singularly perturbed linear functional differential equations with small delays", Applicable Analysis, vol.47, pp. 213-225.
[19] L.T.Magalhães (1987), "The sprectrum of invariant sets for dissipative semiflows, in dynamics of infinite dimensional systems", NATO ASI series, No.F-37, Springer Verlag, New York.
[20] L.T.Magalhães (1987), "Persistence and smoothness of hyperbolic invariant manifold for functional differential equations", SIAM J. Math. Anal., 18, no. 3 pp. 670-693.
[21] J.L.Massera and J.J.Shaffer (1966) , "Linear differential equations and function spaces", Academic Press, New York.
[22] K.J.Palmer (1984), "Exponential dichotomies and transversal homiclinic points", J. Differential Equations, vol.55, pp.225-256.
[23] A.Pazy (1983), "Semigroups of linear operators and applications to partial differential equations", Applied Mathematical Sciences, vol.44, Springer Verlag, New York.
[24] O.Perron (1930), "Die stabilit'atsfrage bei differentialgleichungen", Math.Z vol.32, pp. 703-728.
[25] R.T.Rau (1993), "Hyperbolic linear skew-product semiflows", Arbeitsbereich Funktionalanalysis, Mathematishes Institut der Universit'a, Germany.
[26] R.T.Rau (1994), "Hyperbolic evolutionary semigroups on vector valued function spaces", Semigroup Forum, vol.48, pp. 107-118.
[27] H.Rodrigues (1980), "On gorwth and decay of solutions of perturbed retarded linear equations", Tôhoku Math.Journ, 32, pp. 593-605.
[28] R.J.Sacker and G.R.Sell (1974), "Existence of dichotomies and invariant splitting for linear differential systems $I^{\prime \prime}$ J.Differential Equations, 15, pp. 429-458.
[29] R.J.Sacker and G.R.Sell (1976 A), "Existence of dichotomies and invariant splitting for linear differential systems II", J.Differential Equations, 22, pp.478-496.
[30] R.J.Sacker and G.R.Sell (1976 B), "Existence of dichotomies and invariant splittings for linear differential systems III", J.Differential Equations, 22, pp. 497-52.
[31] R.J.Sacker and G.R.Sell, "Dichotomies for linear evolutionary equations in Banach spaces", J. Dynamics Differential Equations (to appear).
[32] R.Temam (1988), "Infinite dimensional dynamical systems in Mechanics and Physics", Springer Verlag, New York.
[33] W.Zhang (1993), "Generalized exponential dichotomies and invariant manifolds for differential equations", Advances in Mathematics, vol. 22, No. 1.

# CENTER FOR DYNAMICAL SYSTEMS AND NONLINEAR STUDIES Report Series 

September 1993 -

## CDSNS92-

102 A normally elliptic Hamiltonian bifurcation, H.W. Broer, S.-N. Chow, Y. Kim and G. Vegter.

103 Large time behavior of an explicit finite difference scheme for an equation arising from compressible flow through porous media, H. Fan.

104 On the perturbation of the kernel for delay systems with continuous kernels, G. Hines.

105 Manuscript withdrawn.
106 Shape index and other indices of Conley type for local maps on locally compact Hausdorff spaces, M. Mrozek.

107 Differentiability with respect to boundary conditions and deviating argument for functional differential systems, J. Ehme, P.W. Eloe and J. Henderson.

108 A discontinuous semilinear elliptic problem without a growth condition, M. Bouguima and A. Boucherif.

109 Equivalent dynamics for a structured population model and a related functional differential equation, H.L. Smith.

110 A degenerate singularity generating geometric Lorenz attractors, F. Dumortier, H. Kokubu and H. Oka.

111 Statistical properties of the periodic Lorentz gas. Multidimensional case, N.I. Chernov.

## CDSNS93-

112 The symmetric Hartman-Grobman Theorem, H.M. Rodrigues.
113 Variation of constants for hybrid systems of FDE, J.K. Hale and W. Huang.
114 Neumann eigenvalue problems on exterior perturbations of the domain, J.M. Arrieta.

115 Bifurcation of equilibria for one-dimensional semilinear equation of the thermoelasticity, L.A.F. de Oliveira and A. Perissinotto, Jr.

116 Isolating neighborhoods and chaos, K. Mischaikow and M. Mrozek.

117 Limits of semigroups depending on parameters, J.K. Hale and G. Raugel.
118 On the sign-variations of solutions of nonlinear two-point boundary value problems, A. Boucherif and B.A. Slimani.

119 Two-point boundary value problems for fourth order ordinary differential equations, A. Boucherif and J. Henderson.

120 The structure of isolated invariant sets and the Conley Index, K. Mischaikow.
121 Competition for a single limiting resource in continuous culture: The variableyield model, H.L. Smith and P. Waltman.

122 Inertial manifolds and the cone condition, J.C. Robinson.
123 The melnikov method and elliptic equations with critical exponent, R.A. Johnson, X . Pan and Y. Yi.

124 Chaos in the Lorenz Equations: a Computer Assisted Proof, K. Mischaikow and M. Mrozek

125 Invariant manifolds and foliations for quasiperiodic systems, S.-N. Chow and K. Lu.

126 Gradient-like structure and Morse decompositions for time-periodic onedimensional parabolic equations, X.-Y. Chen and P. Poláčik.

127 Homoclinics and subharmonics of nonlinear two dimensional systems. Boundedness of generalized inverses, H.M. Rodrigues and J.G. Ruas-Filho.

128 Periodic boundary value problems and a priori bounds on solutions, A. Boucherif.
129 Attractors in inhomogeneous conservation laws and parabolic regularizations, H. Fan and J.K. Hale.

130 Periodic boundary value problems with Caratheodory nonlinearities, A. Boucherif.
131 Attractors and convergence of PDE on thin L-shaped domains, J.K. Hale and G. Raugel.

132 Attractor of a semigroup of multi-valued mappings corresponding to an elliptic equation, A.V. Babin.

133 Rotators, periodicity and absence of diffusion in cyclic cellular automata, L.A. Bunimovich and S.E.Troubetzkoy.

134 About completeness for a class of unbounded operators appearing in delay equations, S.M. Verduyn Lunel.

135 A new invariant manifold with an application to a smooth conjugacy at a node, W.M. Rivera.

136 A new invariant manifold withan application to bifurcation of smoothness, W.M. Rivera.

137 Attractors and inertial manifolds for the dynamics of a closed thermosyphon, A. Rodriguez-Bernal.
$138 C^{k}$-smoothness of invariant curves in a global saddle-node bifurcation, T. Young.
139 Partial neutral functional differential equations, J.K. Hale.
140 Numerical dynamics, J.K. Hale.
141 A reaction diffusion equation on a thin L-shaped domain, J.K. Hale and G. Raugel.
142 Criteria of spatial chaos in lattice dynamical systems, V.S. Afraimovich and S.-N. Chow.

143 Statistical properties of 2-D generalized hyperbolic attractors, V.S. Afraimovich, N.L. Chernov and E.A. Sataev.

144 Long-time behavior as $\lambda \rightarrow 0$ of solutions of parabolic equations depending on $\lambda t$, A. Babin and S.-N. Chow.

145 Some remarks on PDEs with nonlinear Neumann boundary conditions, A. Rodriguez-Bernal.

146 Structure of the global attractor of cyclic feedback systems, T. Gedeon and K. Mischaikow.

147 Dynamics of almost periodic scalar parabolic equations, W. Shen and Y. Yi.
148 Existence of a global attractor for the sunflower equation with small delay, M. Lizana.

149 Reaction-diffusion systems on domains with thin channels, S.M. Oliva.
150 Existence and partial characterization of the global attractor for the sunflower equation, M. Lizana.

151 Asymptotic almost periodicity of scalar parabolic equations with almost periodic time dependence, W. Shen and Y. Yi.

152 Zeta functions, periodic trajectories, and the Conley index, C. McCord, K. Mischaikow and M. Mrozek.

153 Center manifold and stability for skew-product flows, S.-N. Chow and Y. Yi.
154 Partial and complete linearizations at stationary points of infinite dimensional dynamical systems with foliations and applications, W.M. Rivera.

155 Dynamical spectrum for skew product flow in Banach spaces, S.N. Chow and H. Leiva.

156 Existence and roughness of the exponential dichotomy for skew-product semiflow in Banach spaces, S-N. Chow and H. Leiva.

## CDSNS94-

157 Rates of eigenvalues on a dumbbell domain. Simple eigenvalue case, J.M. Arrieta.
158 Existence of standing waves for competition-diffusion equations, Y. Kan-On.
159 Fisher-type property of travelling waves for competition-diffusion equations, Y. Kan-On.

160 Uniform ultimate boundedness and synchronization, H.M. Rodrigues.
161 Complete families of pseudotrajectories and shape of attractors, S.Yu. Pilyugin.
162 Aubry-Mather Theorem and quasiperiodic orbits for time dependent reversible systems, S.-N. Chow and M.-L. Pei.

163 Slowly-migrating transition layers for the discrete Allen-Cahn and Cahn-Hilliard equations, C.P. Grant and E.S. Van Vleck.

164 Nontrivial partially hyperbolic sets from a co-dimension one bifurcation, T. Young.
165 Asymptotically autonomous semiflows: Chain recurrence and Lyapunov functions, K. Mischaikow, H. Smith and H.R. Thieme.

166 On the problem of stability in lattice dynamical systems, L.A. Bunimovich and E.A. Carlen.

167 Remark on continuous dependence of attractors on the shape of domain, A.V. Babin and S.Yu. Pilyugin.

168 Special pseudotrajectories for lattice dynamical systems, V.S. Afraimovich and S.Yu. Pilyugin.

169 Proof and generalization of Kaplan-Yorke conjecture on periodic solution of differential delay equations, $\mathrm{J} . \mathrm{Li}$ and $\mathrm{X}-\mathrm{Z}$. He.

170 On the construction of periodic solutions of Kaplan-Yorke type for some differential delay equations, K. Gopalsamy, J. Li and X-Z. He.

171 Uniformly accurate schemes for hyperbolic systems with relaxation, R.E. Caflisch, S. Jin and G. Russo.

172 Numerical integrations of systems of conservation laws of mixed type, S. Jin.

173 A complex algorithm for computing Lyapunov values, R. Mao and D. Wang.
174 On minimal sets of scalar parabolic equations with skew-product structures, W. Shen and Y. Yi.

175 Lorenz type attractors from codimensional-one bifurcation, V. Afraimovich, S.-N. Chow and W. Liu.

176 Dynamics in a discrete Nagumo equation - spatial chaos, S.-N. Chow and W. Shen.

177 Stability and bifurcation of traveling wave solutions in coupled map lattices, S.-N. Chow and W. Shen.

178 Density of defects and spatial entropy in extended systems, V.S. Afraimovich and L.A. Bunimovich.

179 On the second eigenvalue of the Laplace operator penalized by curvature, E.M. Harrell II.

180 Singular limits for travelling waves for a pair of equations, V. Hutson and K. Mischaikow.

181 Conley Index Theory: Some recent developments, K. Mischaikow.
182 Variational principle for periodic trajectories of hyperbolic billiards, L.A. Bunimovich.

183 Epidemic waves: A diffusion model for fox rabies, W.M. Rivera.
184 High complexity of spatial patterns in gradient reaction-diffusion systems, V. Afraimovich, A. Babin and S.-N. Chow.

185 Spatial chaotic structure of attractors of reaction-diffusion systems, V. Afraimovich, A. Babin and S.-N. Chow.

Dynamical spectrum for skew product flow in Banach spaces
by S.N. Chow and H. Leiva

CDSNS93-155

# DYNAMICAL SPECTRUM FOR SKEW PRODUCT FLOW IN BANACH SPACES 

S.N. CHOW AND H.LEIVA<br>School of Mathematics-CDSNS, Georgia Tech<br>Atlanta, Ga 30332

Abstract. In an earlier paper [1] we characterized the dynamical spectrum for Linear Skew-Product Semiflow in infinite dimensional Banach spaces. It was proved the spectrum is always a closed set, but it could be empty. Also we investigate the relation between the dynamical spectrum and the Lyapunov exponents. In this paper we shall characterize the dynamical spectrum for Linear Skew-Product Flow $\pi=(\Phi, \sigma)$ in a Banach space $X$. The fact that $\pi$ is a flow allows us to prove that the spectrum is a nonempty compact set and get more information about it, also we can tell more about the Lyapunov exponents. Finally ours results can be applied to hyperbolic partial differential equations and neutral functional differential equations.

1

[^1]
## 1 Introduction

In an earlier paper [1] we began an investigation of the dynamical spectrum for time dependent systems in infinite dimensional Banach spaces, using the concept of skew-product semiflow. That is the case for parabolic partial differential equations and functional differential equations; for that reason in [1] we use the concept, of negative continuation and exponential dichotomy used by Sacker-Sell in [18]. In this definition of exponential dichotomy we assume that the unstable mainfold has finite dimension and is contained in the set of points that have a unique negative continuation. We characterized the dynamical spectrum and proved that the spectrum is always a closed set, but it could be empty; also, we investigate the relation between the dynamical spectrum and the Lyapunow exponents.

This paper is concerned with the dynamical spectrum for time dependent linear systems whose solutions are globally defined in $\mathbb{R}$. This is the case for hyperbolic partial differential equations, neutral functional differential equations and abstract ordinary differential equation $\dot{x}=A(t) x$ with bounded operator $A(t)$.

To study this problems we shall use the unified setting of a linear skewproduct flow $\Phi(\theta, t): \mathcal{E} \times \mathbb{R} \rightarrow \mathcal{E}$ (see Definition 2.1), where $\mathcal{E}=X \times \Theta, X$ is a Banach space, $\Theta$ is a compact topological Hausdorff space and $\pi$ is given by

$$
\pi(x, \theta, t)=(\Phi(\theta, t) x, \theta . t), \quad t \in \mathbb{R}, \quad x \in X, \quad \theta \in \Theta
$$

Many people have worked with skew-product flow, and semiflows in infinite dimensional Banach spaces. For example Sacker-Sell in [18] studied the existence of exponential dichotomy for the skew-product semiflow $\pi=(\Phi, \sigma)$ on $\mathcal{E}=X \times \Theta$.

Also, Magalhães in [11] give some characterization of the dynamical spectrum. A different characterization of the dynamical spectrum for skew-product flow in infinite Banach spaces appears in R. T. Rau [13]. He associates a strongly continuous group to the skew-product flow $\pi=(\Phi, \sigma)$ in the following way: Given a skew-product flow $\pi=(\Phi, \sigma)$ on $\mathcal{E}=X \times \Theta$ we can associate a family $\{T(t)\}_{t \in \mathbb{R}}$ of linear operators on the Banach space $C(\Theta, X)$ defined by

$$
\begin{equation*}
T(t) f(\theta)=\Phi(\theta \cdot(-t), t) f(\theta \cdot(-t)), \quad \forall t \in \mathbb{R} \tag{1.1}
\end{equation*}
$$

for all $\theta \in \Theta$ and $f \in C(\Theta, X)$.
It is shown that the operator family $\{T(t)\}_{t \in \mathbb{R}}$ given by (1.1) is a strongly continuous group on $C(\Theta, X)$.

Since $\pi=(\Phi, \sigma)$ is a flow (two side flow) the definition of exponential dichotomy is the same as in the finite dimensional case; this allows Rau in [13] to give the following characterization of the dynamical spectrum

Theorem 1.1 (Rau [13]) Let $\pi=(\Phi, \sigma)$ be a skew-product flow on $\mathcal{E}=X \times \Theta$ and denote $G$ the infinitesimal generator of the evolution group $\{T(t)\}_{t \in \mathbb{R}}$ given by (1.1). Then

$$
\Sigma(\Theta)=\ln |\sigma(T(1)) \backslash\{0\}| \supseteq G
$$

A similar result as Rau's can be found in Latushki-Stepin [9]
In this paper we shall give a different characterization of the dynamical spectrum. Our characterization is an extension of the Sacker and Sell Theorem given in [18]. Here we proved that the spectrum can be written as a countable union of nonempty close disjoints intervals. we show the relation between the spectrum and the spectral subbundles associated with the corresponding spectral interval.

Also, this spectral decomposition can be used to study invariant manifold around an invariant set.

In this paper we follow closely the work done by Sacker and Sell in [18] for the finite dimensional case and the notation used for them in [19].

Since $\pi$ is a flow (two side flow) the definitions of exponential dichotomy is very simple and we don't have to worry about negative continuation, like in [19], [1] and [11]. This allows us to prove that the spectrum is a nonempty compact set and we can give a simple characterization of the Lyapunov exponents in terms of the dynamical spectrum.

Finally we present some examples of Skew-Product Flow arising from hyperbolic partial differential equations and neutral functional differential equations.

## 2 Preliminaries

### 2.1 Linear Skew-Product Flow

In this section we shall present some definitions, notations, and results about Skew Product flows on Banach Bundles that we will use in the next sections.

Definition 2.1 Let $\mathcal{E}=X \times \Theta$ be given where $X$ is a fixed Banach space (state space) and $\Theta$ is a compact Hausdorff Space. Assume that $\sigma(\theta, t)=\theta . t$ is a flow on $\Theta$, i.e., the mapping $(\theta, t) \rightarrow \theta \cdot t$ is continuous, $\theta .0=\theta$, and we have $\theta \cdot(s+t)=(\theta \cdot s) \cdot t$ for all $s, t \in \mathbb{R}$. Then we shall call a Linear Skew-Product Flow $\pi=(\Phi, \sigma)$ on $\mathcal{E}$ as a mapping

$$
\pi(x, \theta, t)=(\Phi(\theta, t) x, \theta \cdot t), \quad \forall t \in \mathbb{R}
$$

with the following properties:
(1) $\Phi(\theta, 0)=I$, the identity operator, for all $\theta \in \theta$.
(2) $\lim _{t \rightarrow 0} \Phi(\theta, t) x=x$, and this limit is uniform in $\theta$. This mean that for every $x \in X$ and every $\epsilon>0$ there is a $\delta=\delta(x, \epsilon)>0$ such that $\|\Phi(\theta, t) x-x\| \leq \epsilon$, for all $\theta \in \Theta$ whenever $0 \leq t \leq \delta$.
(3) $\Phi(\theta, t)$ is a bounded linear mapping from $X$ into $X$ that satisfies the cocycle identity:

$$
\begin{equation*}
\Phi(\theta, s+t)=\Phi(\theta \cdot t, s) \Phi(\theta, t) \quad \theta \in \Theta ; \quad s, t \in \mathbb{R} \tag{2.1}
\end{equation*}
$$

(4) For each $t \in \mathbb{R}$ the mapping of $\mathcal{E}$ into $X$ given by $(x, \theta) \rightarrow \Phi(\theta, t) x$ is continuous.

Properties (2) and (3) imply that for each $(x, \theta) \in \mathcal{E}$ the solution operator $t \rightarrow$ $\Phi(\theta, t) x$ is continuous for $t \in \mathbb{R}$. Indeed one has

$$
\|\Phi(\theta, t+h) x-\Phi(\theta, t) x\|=\|[\Phi(\theta \cdot t, h)-I] \Phi(\theta, t) x\|
$$

which goes to zero as $h \rightarrow 0$. The cocycle identity (2.1) implies that $\Phi(\theta, t)$ is an isomorphism with inverse

$$
\Phi^{-1}(\theta, t)=\Phi(\theta \cdot t,-t) \quad \forall t \in \mathbb{R}
$$

Proposition 2.1 Let $\pi=(\Phi, \sigma)$ be a linear skew-product flow on $\mathcal{E}$. Then there exist constants $M \geq 1, \quad a>0$ such that

$$
\begin{equation*}
\|\Phi(\theta, t)\| \leq M e^{a|t|}, \quad \theta \in \Theta, \quad t \in \mathbb{R} \tag{2.2}
\end{equation*}
$$

Proof First we claim that there is a $\delta>0$ such that

$$
M=\sup \{\|\Phi(\theta, t)\|: \theta \in \Theta,-\delta \leq t \leq \delta\}<\infty
$$

For the purpose of contradiction. Let us assume that there are sequences $\theta_{n} \in$ $\Theta, t_{n} \in \mathbb{R}$ such that $t_{n} \rightarrow 0$ and $\left\|\Phi\left(\theta_{n}, t_{n}\right)\right\|>n$.

The Banach-Steinhaus Theorem (Uniform Boundedness Principle) implies that there is $x \in X$ such that

$$
\left\{\left\|\Phi\left(\theta_{n}, t_{n}\right) x\right\|: n \in \mathbb{N}\right\}
$$

is unbounded. This contradicts the fact that

$$
\lim _{t \rightarrow 0} \Phi(\theta, t) x=x
$$

unifomly for $\theta \in \Theta$. Therefore $M<\infty$. On the other hand we have that $\Phi(\theta, 0)=$ $I$, so $M \geq 1$.

Now fix $t \in \mathbb{R}$. Assume that $t \geq 0$ (a similar argument will take care of the case $t \leq 0$ ). Let $m$ be an integer satisfying $m \leq t / \delta \leq m+1 \Longleftrightarrow \delta_{m} \leq t \leq m \delta+\delta$.

Then for any $\theta \in \Theta$ one has

$$
\|\Phi(\theta, t)\|=\|\Phi(\theta, t-\delta m+\delta m)\|=\|\Phi(\theta \cdot \delta,(t-\delta m)+\delta(m-1)) \Phi(\theta, \delta)\|
$$

Now putting:

$$
\theta_{0}=\theta, \quad \theta_{1}=\theta_{0} \cdot \delta, \quad \theta_{2}=\theta_{1} \cdot \delta, \cdots, \theta_{m}=\theta_{m-1} \cdot \delta
$$

we get the following:

$$
\begin{aligned}
\|\Phi(\theta, t)\| & =\left\|\Phi\left(\theta_{m}, t-\delta m\right) \Phi\left(\theta_{m-1}, \delta\right) \ldots \Phi\left(\theta_{0}, \delta\right)\right\| \\
& \leq M^{m+1}=M \cdot M^{m} \leq M \cdot M^{t / \delta}
\end{aligned}
$$

If we put

$$
a=1 / \delta \ln (M)
$$

then

$$
\|\Phi(\theta, t)\| \leq M e^{a t}
$$

### 2.2 Projectors and Subbundles

A Banach bundle $\mathcal{E}$ with fiber $X$ over a base space $\Theta$ with projection $\mathcal{P}$ is denoted by $(\mathcal{E}, X, \Theta, \mathcal{P})$, or $\mathcal{E}$ for short, and is defined as follows:
(1) $X$ is a fixed Banach space and $\theta$ is a compact Hausdorff space.
(2) The mapping $\mathcal{P}: \mathcal{E} \rightarrow \Theta$ is a continuous mapping.
(3) For each $\theta \in \Theta, \mathcal{P}^{-1}(\theta)=\mathcal{E}(\theta)$ is a Banach space, which is referred to as the fiber over $\theta$.
(4) For each $\theta \in \Theta$, there is an open neighborhood $U$ of $\theta$ in $\Theta$ and a homeomorphism
$\tau: \mathcal{P}^{-1}(U) \rightarrow X \times U$ such that for each $\mathcal{N} \in U, \mathcal{P}^{-1}(\mathcal{N})$ is a mapped onto $X \times\{\mathcal{N}\}$ and $\tau: \mathcal{P}^{-1}(\mathcal{N}) \rightarrow X \times\{\mathcal{N}\}$ is a linear isomorphism.
(5) The norms $\|\cdot\|=\|\cdot\|_{\theta}$ on the fiber $\mathcal{P}^{-1}(\theta)$ vary continuously in $\theta$.

O e can use the local coordinate notation $(x, \theta)$ to denote a typical point in a Banach bundle $\mathcal{E}$. By this we mean that $(x, \theta) \in \mathcal{E}$. This is a shortened way to refer to property (4) above.

For any subset $\mathcal{F} \subset \mathcal{E}$ we define the fiber:

$$
\begin{gathered}
\mathcal{F}(\theta):=\{x \in X:(x, \theta) \in \mathcal{F}\} \\
\mathcal{E}(\theta):=X \times\{\theta\}
\end{gathered}
$$

If $U \subset \Theta$, then we shall define the following set

$$
\mathcal{F}(U):=\cup_{\theta \in U} \mathcal{F}(\theta)
$$

$$
\mathcal{E}_{0}=\{(x, \theta) \in \mathcal{E}: x=0\}
$$

$\mathcal{E}_{0}$ is called the zero fiber.
A mapping $\mathbf{P}: \mathcal{E} \rightarrow \mathcal{E}$ is said to be a projection if $\mathbf{P}$ is continuous and has the form $\mathbf{P}(x, \theta)=(P(\theta), \theta)$ where $P(\theta)$ is a bounded linear projection on the fiber $\mathcal{E}(\theta)$. For any projector P we define the range and the null space by

$$
\mathcal{R}=\mathcal{R}(\mathbf{P})=\{(x, \theta) \in \mathcal{E}: P(\theta) x=x\}
$$

and

$$
\mathcal{N}=\mathcal{N}(\mathbf{P})=\{(x, \theta) \in \mathcal{E}: P(\theta) x=0\}
$$

Since $\mathbf{P}$ is continuous, this means that the fibers $\mathcal{R}(\theta)$ and $\mathcal{N}(\theta)$ vary continuously in $\theta$. This also means that $P(\theta)$ varies continuous in the operator norm. The following result can be found in Sacker- Sell [18].

Lemma 2.1 Let $\mathbf{P}$ be a projector on $\mathcal{E}$. Then $\mathcal{R}$ and $\mathcal{N}$ are closed in $\mathcal{E}$, and one has

$$
\mathcal{R}(\theta) \cap \mathcal{N}(\theta)=\{0\}, \quad \mathcal{R}(\theta)+\mathcal{N}(\theta)=\mathcal{E}
$$

for all $\theta \in \Theta$.

Definition 2.2 A subset $\nu$ is said to be a subbundle of $\mathcal{E}$ if there is a projector $P$ on $\mathcal{E}$ with the property that $\mathcal{R}(\mathbf{P})=\mathcal{V}$

In this case $\mathcal{W}=\mathcal{N}(P)$ is a Complementary subbundle, i.e., $\mathcal{E}=\mathcal{V}+\mathcal{W}$ as a Whitney Sum .

Lemma 2.2 Let $\mathcal{V} \subset \mathcal{E}$ with the properties:
(A) $\mathcal{V}$ is closed.
(B) $\mathcal{V}(\theta)$ is a linear subspace of $\mathcal{E}(\theta)$ for all $\theta \in \Theta$.
(C) $\operatorname{codim} \mathcal{V}(\theta)$ is finite for all $\theta \in \Theta$.
(D) $\operatorname{codim} \mathcal{V}(\theta)$ is locally constamt on $\Theta$.

Then $\mathcal{V}$ is a subbundle of $\mathcal{E}$.

Proof See [18].

### 2.3 Stable, Unstable and the Initial Bounded Sets

Let $\pi=(\Phi, \sigma)$ be a given linear skew-product flow defined on $\mathcal{E}=X \times \Theta$. For $\lambda \in \mathbb{R}$ we define the shifted flow as follows:

$$
\begin{gathered}
\pi_{\lambda}=\left(\Phi_{\lambda}, \sigma\right), \Phi_{\lambda}(\theta, t)=e^{\lambda t} \Phi(\theta, t) \text { for } t \in \mathbb{R}, \quad \theta \in \Theta . \\
\mathcal{B}_{\lambda}=\left\{(x, \theta) \in \mathcal{E}: \sup _{t \in \mathbb{R}}\left\|e^{-\lambda t} \Phi(\theta, t) x\right\|<+\infty\right\} \\
\mathcal{S}_{\lambda}=\left\{(x, \theta) \in \mathcal{E}:\left\|e^{-\lambda t} \Phi(\theta, t) x\right\| \rightarrow 0, t \rightarrow+\infty\right\} \\
\mathcal{U}_{\lambda}=\left\{(x, \theta) \in \mathcal{E}:\left\|e^{-\lambda t} \Phi(\theta, t) x\right\| \rightarrow 0, t \rightarrow-\infty\right\}
\end{gathered}
$$

The set $\mathcal{U}_{\lambda}$ is the unstable set, $\mathcal{S}_{\lambda}$ is the stable set, and $\mathcal{B}_{\lambda}$ is the initial bounded set corresponding to $\pi_{\lambda}$. If $\lambda=0$ we shall denote $\mathcal{B}=\mathcal{B}_{0}, \mathcal{U}=$ $\mathcal{U}_{0}$ and $\mathcal{S}=\mathcal{S}_{0}$.

We are interested in knowing when $\mathcal{S}_{\boldsymbol{\lambda}}$ and $\mathcal{U}_{\boldsymbol{\lambda}}$ are complementary invariant subbundles of $\mathcal{E}$. The answer of this quation can be formulated in terms of dichotomies.

Definition 2.3 A project $\mathbf{P}$ on $\mathcal{E}$ is said to be invariant if we have

$$
\begin{equation*}
P(\theta . t) \Phi(\theta, t)=\Phi(\theta, t) P(\theta), \quad t \in \mathbb{R}, \quad \theta \in \Theta \tag{2.3}
\end{equation*}
$$

which is equivalent to:

$$
\begin{equation*}
P(\theta . t)=\Phi(\theta, t) P(\theta) \Phi^{-1}(\theta, t), \quad t \in \mathbb{R}, \quad \theta \in \Theta \tag{2.4}
\end{equation*}
$$

Definition 2.4 We shall say that a linear skew-product flow $\pi=(\Phi, \sigma)$ on $\mathcal{E}$ has an exponential dichotomy (ED) over an invariant set $\hat{\Theta}$, were $\hat{\Theta} \subset \Theta$, if there is an invariant Projector $\mathbf{P}$ on $\mathcal{E}$ and constants $k \geq 1, \beta>0$ such that

$$
\begin{gather*}
\left\|\Phi(\theta, t) P(\theta) \Phi^{-1}(\theta, s)\right\| \leq k e^{-\beta(t-s)}, \quad s \leq t  \tag{2.5}\\
\left\|\Phi(\theta, t)[I-P(\theta)] \Phi^{-1}(\theta, s)\right\| \leq k e^{\beta(t-s)}, \quad s \geq t \tag{2.6}
\end{gather*}
$$

for all $\theta \in \hat{\Theta}$

## Remark 2.1

(1) If $\hat{\Theta}=\{\theta\}$, then E.D corresponds to the usual concept of dichotomy
(2) If $\hat{\Theta}=\Theta$, then E.D over $\hat{\Theta}$ is equivalent to the splitting of $\mathcal{E}$.
(3) $P(\theta)$ varies continuously over $\hat{\Theta}$.
(4) $k, \beta$ depend of $\hat{\Theta}$.

Proposition 2.2 If $\pi$ is a linear skew-product flow on $\mathcal{E}=X \times \Theta$ admits an exponential dichotony over $\Theta$, then one has that the initial bounded set $\mathcal{B}=\mathcal{E}_{0}$ and the correspondent Projector $\mathbf{P}$ is such that:

$$
\begin{gathered}
\mathcal{R}(\mathbf{P})=\mathcal{S}(\hat{\Theta}), \quad \mathcal{N}(\mathbf{P})=\mathcal{U}(\hat{\Theta}) \\
\mathcal{E}=\mathcal{R}(\mathbf{P}) \oplus \mathcal{N}(\mathbf{P})=\mathcal{S}(\hat{\Theta}) \oplus \mathcal{U}(\hat{\Theta})
\end{gathered}
$$

(The whitney sum of two bundles)

Proof Let $\mathbf{P}$ be the correspondent Projector. Consider $(x, \theta) \in \mathcal{B}$ and set

$$
y=P(\theta) x, \quad z=(I-P(\theta)) x
$$

Then we get the following:

$$
\|\Phi(\theta, t)\|<N, \quad \forall t \in \mathbb{R}
$$

and

$$
y=P(\theta) \Phi^{-1}(\theta, t) \Phi(\theta, t) x=\Phi(\theta, 0) P(\theta) \Phi^{-1}(\theta, t) x
$$

Therefore, from (2.5) we get:

$$
\begin{gathered}
\|y\| \leq k e^{\beta t} N, \quad t<0 \Rightarrow\|y\|=0 \\
z=\Phi(\theta, 0)[I-P(\theta)] \Phi^{-1}(\theta, t) \Phi(\theta, t) x
\end{gathered}
$$

so

$$
\|z\| \leq k e^{-\beta t} N, \quad t \in \mathbb{R} \Rightarrow\|z\|=0
$$

Hence

$$
x=P(\theta) x+[I-P(\theta)] x=y+z=0
$$

It is easy to prove that :

$$
\mathcal{R}(\mathrm{P}) \subset \mathcal{S} \quad \mathcal{N}(\mathrm{P}) \subset \mathcal{U}
$$

Now, since $\quad \mathcal{E}=\mathcal{R}(\mathbf{P})+\mathcal{N}(\mathbf{P})$, then $\mathcal{S}=\mathcal{R}(\mathbf{P}), \quad \mathcal{U}=\mathcal{N}(\mathbf{P})$.

## 3 The Dynamical Spectrum

Let $\hat{\Theta}$ be an invariant subset of $\Theta$ under the flow $\sigma$. Then the resolvent $\rho(\hat{\Theta})$ of $\hat{\Theta}$ under $\pi$ is defined as follows:
$\rho(\hat{\Theta}):=\left\{\lambda \in \mathbb{R}: \pi_{\lambda}\right.$ admits and exponential dichotomy over $\left.\hat{\Theta}\right\}$.

The spectrum $\Sigma(\hat{\Theta})$ of $\hat{\Theta}$ under $\pi$ is defined as follows

$$
\Sigma(\hat{\Theta})=\mathbb{R} \backslash \rho(\hat{\Theta})
$$

Our main results are the following Theorems:
Theorem 3.1 Let $\pi=(\Phi, \sigma)$ be a skew-product flow on $\mathcal{E}=X \times \Theta$ and $\hat{\Theta}$ a compact connected invariant subset of $\Theta$. Then the following statements are valid:
(A) There is $a>0$ such that

$$
\|\Phi(\theta, t)\| \leq M e^{a|t|}, \quad \forall \theta \in \Theta \quad \text { and } t \in \mathbb{R},
$$

and $\quad \Sigma(\hat{\Theta}) \neq \emptyset, \quad \Sigma(\hat{\Theta}) \subset[-a, a]$
(B) For each set $\left\{\lambda_{0}<\lambda_{1}<\cdots<\lambda_{m}\right\} \subset \rho(\hat{\Theta})$ with $\lambda_{0}<-a$ and $a<\lambda_{m}$ such that $\Sigma(\hat{\Theta}) \cap\left(\lambda_{i_{1}}, \lambda_{i}\right) \neq \phi$ we get that

$$
\mathcal{V}_{i}:=\mathcal{V}_{i}(\hat{\Theta})=\mathcal{S}_{\lambda_{i}}(\hat{\Theta}), \quad i=1,2, \cdots, m
$$

are invariant subbundles of $\mathcal{E}(\hat{\Theta})$.
(C) Let $\pi^{i}$ be the restriction of $\pi$ to $\mathcal{V}_{i}$ and $\Sigma_{i}(\hat{\Theta})$ the spectrum of $\left(\mathcal{V}_{i}, \pi^{i}\right)$ over $\hat{\Theta}$. Then one has that

$$
\Sigma_{i}(\hat{\Theta})=\Sigma(\hat{\Theta}) \cap\left(\lambda_{i-1}, \lambda_{i}\right), \quad i=1,2, \cdots, m .
$$

(D) $\Sigma_{i}(\hat{\Theta})=\Sigma(\hat{\Theta}) \cap\left(\lambda_{i-1}, \lambda_{i}\right), \quad i=1,2, \cdots, m$
(E) $\Sigma(\hat{\Theta})=\bigcup_{i=1}^{m} \Sigma_{i}(\hat{\Theta})$
(F) $\mathcal{V}_{i}(\hat{\Theta}) \cap \mathcal{V}_{j}(\hat{\Theta})=\mathcal{E}_{0}(\hat{\Theta}), \quad i \neq j$
(G) $\mathcal{E}(\hat{\Theta})=\mathcal{V}_{1}(\hat{\Theta})+\mathcal{V}_{2}(\hat{\Theta})+\cdots+\mathcal{V}_{m}(\hat{\Theta})$ (Whitney sum)

In order to get more information about the spectrum we shall put some restriction on the unstable manifold $\mathcal{U}_{\lambda}$ for some $\lambda \in \rho(\hat{\Theta})$. Also we will need the following notation:

For $\lambda \in \rho(\hat{\Theta})$ we shall define:

$$
\begin{equation*}
\Sigma_{\lambda}(\hat{\Theta}):=\Sigma(\hat{\Theta}) \cap(-\infty, \lambda) \tag{3.1}
\end{equation*}
$$

Theorem 3.2 Assume that $\operatorname{dim} \mathcal{E}(\hat{\Theta})=\infty, \quad \lambda \in \rho(\hat{\Theta}), \quad a>0$ is as in part (A) of Theorem 3.1. Then the following statements are valid:
(A) If $\operatorname{dim} \mathcal{U}_{\lambda}=n(\lambda)<\infty$, then $\lambda \geq-a$
(B) If $\operatorname{dim} \mathcal{S}_{\lambda}=m(\lambda)<\infty$, then $\lambda \leq a$
(C) If $\lambda \geq-a$ and $\mathcal{U}_{\lambda} \neq \mathcal{E}_{0}(\hat{\Theta})$, then $\lambda \leq a$
(D) If $\lambda \leq a$ and $\operatorname{dim} \mathcal{S}_{\lambda} \neq \mathcal{E}_{0}(\hat{\Theta})$, then $\lambda \geq-a$
(E) If $\operatorname{dim} \mathcal{U}_{\lambda} \rightarrow n(-a)<\infty$, as $\lambda \rightarrow-a^{+}$
then

$$
-a \in \Sigma(\hat{\Theta}) \text { and } \operatorname{dim} \mathcal{U}_{\lambda}=n(\lambda)<\infty \quad \forall \lambda \in[-a+\infty) \cap \rho(\hat{\Theta})
$$

(F) If $\operatorname{dim} S_{\lambda} \rightarrow m(a)<\infty$, as $\lambda \rightarrow a^{+}$
then

$$
a \in \Sigma(\hat{\Theta}) \text { and } \operatorname{dim} \mathcal{S}_{\lambda}=m(\lambda) \quad \forall \lambda \in(-\infty, a] \cap \rho(\hat{\Theta})
$$

(G) If $1 \leq \operatorname{dim} \mathcal{U}_{\lambda_{0}}=n\left(\lambda_{0}\right)<\infty$, then $\lambda_{0} \in[-a, a]$ and

$$
\begin{equation*}
\Sigma(\hat{\Theta})=\Sigma_{\lambda_{0}}(\hat{\Theta}) \cup\left(\bigcap_{i=1}^{m}\left[a_{i}, b_{i}\right]\right) \tag{3.2}
\end{equation*}
$$

$$
m \leq \operatorname{dim} \mathcal{U}_{\lambda_{0}}=n\left(\lambda_{0}\right)
$$

## Moreover:

(H) $\Sigma_{\lambda_{0}}(\hat{\Theta})=\left[a, \lambda_{0}\right) \cap \Sigma(\hat{\Theta})$
(I) If $\Sigma(\hat{\Theta}) \subset\left[\lambda_{0}, a\right] \Rightarrow \Sigma(\hat{\Theta})=\bigcup_{i=1}^{m}\left[a_{i}, b_{i}\right]$.
(J) $\mathcal{E}(\hat{\Theta})=\mathcal{S}_{\lambda_{0}}(\hat{\Theta})+\mathcal{V}_{1}(\hat{\Theta})+\ldots+\mathcal{V}_{m}(\hat{\Theta})$,
where

$$
\mathcal{V}_{i}(\hat{\Theta})=\mathcal{U}_{\lambda_{i}-1} \cap \mathcal{S}_{\lambda_{i}}, \quad i=1,2, \ldots, m
$$

and

$$
\left\{\lambda_{0}<\lambda_{1}<\ldots<\lambda_{m}\right\} \subset \rho(\hat{\Theta})
$$

with

$$
a<\lambda_{m}, \quad\left(\lambda_{i-1}, \lambda_{i}\right) \cap \Sigma(\hat{\Theta}) \neq \phi
$$

### 3.1 Lemmas

Here we shall derive a number of properties of the spectrum and the resolvent set which will be used in the proof of the main theorems.

Lemma 3.1 Let $\hat{\Theta}$ be a compact invariant set in $\Theta$ and $\lambda \in \mathbb{R}$. Then the following statements are valid:
(A) If $\left\|\Phi_{\lambda}(\theta, t)\right\| \rightarrow 0$ as $t \rightarrow+\infty$ for each $\theta \in \hat{\Theta}$, then $\lambda \in \rho(\hat{\Theta}), \quad \Sigma(\hat{\Theta}) \subseteq$ $(-\infty, \lambda)$, and $\mathcal{S}_{\mu}(\hat{\Theta})=\mathcal{E}(\hat{\Theta})$ for all $\mu \geq \lambda$.
(B) If $\left\|\Phi_{\lambda}(\theta, t)\right\| \rightarrow 0$ as $t \rightarrow-\infty$ for each $\theta \in \hat{\Theta}$, then $\lambda \in \rho(\hat{\Theta}), \quad \Sigma(\hat{\Theta}) \subseteq$ $(\lambda,+\infty)$,
and $\mathcal{U}_{\mu}(\hat{\Theta})=\mathcal{E}(\hat{\Theta})$ for all $\mu \leq \lambda$.

Proof We shall prove (A). The proof of $(B)$ is similar. For each $\theta \in \hat{\Theta}$, the is $T(\theta)>0$ such that $\left\|\Phi_{\lambda}(\theta, t)\right\|<\frac{1}{2}, \quad t \geq T(\theta)$.

Consider $x \in X$ fixed with $\|x\|=1$. By the continuity of $\Phi_{\lambda}(\theta, t) x$ with respect to $\theta$ there exist a neighborhood $N_{x}(\theta)$ of $\theta$ such that $\left\|\Phi_{\lambda}(\bar{\theta}, T(\theta)) x\right\|<\frac{1}{2}$, for all $\tilde{\theta} \in N_{x}(\theta)$. Then by the compactness of $\tilde{\theta}$ we have the following:

$$
\hat{\Theta} \subset \cup_{i=1}^{m} N_{x}\left(\theta_{i}\right), \quad T\left(\theta_{1}\right) \leq T\left(\theta_{2}\right) \leq \cdots \leq T\left(\theta_{m}\right)
$$

We shall put $T_{j}=T\left(\theta_{j}\right), \quad J=1,2, \cdots m$.
We claim the following:

$$
k=\sup \left\{\left\|\Phi_{\lambda}(\theta, t)\right\|: \theta \in \hat{\Theta}, \quad 0 \leq t \leq T_{m}\right\}<\infty .
$$

In fact. Assume that $K=\infty$. Then there are sequences $\left\{\theta_{n}\right\} \subset \hat{\Theta},\left\{t_{n}\right\} \subset$ $\left[0, T_{m}\right]$ such that

$$
\left\|\Phi_{\lambda}\left(\theta_{n}\right), t_{n}\right\|>n, \quad n=1,2,3 \ldots
$$

Since $\hat{\Theta}$ and $\left[0, T_{m}\right]$ are compact sets we can assume that $\left\{\theta_{n}\right\}$ converges to $\theta_{0} \in \hat{\Theta}$ and $\left\{t_{n}\right\}$ converge to $t^{*} \in\left[0, T_{m}\right]$. Then by Banach-Steinhaus Theorem there must be an element $x_{0} \in X$ so that the set:

$$
\left.\left\{\| \Phi_{\lambda}\left(\theta_{n}, t_{n}\right) x_{0}\right) \|: n=1,2,3, \ldots\right\}
$$

is unbounded. On the other hand the definition of skew-product flow implies that

$$
\lim _{n \rightarrow \infty}\left\|\Phi_{\lambda}\left(\theta_{n}, t_{n}\right) x_{0}\right\|=\left\|\Phi\left(\theta_{0}, t^{*}\right) x_{0}\right\|
$$

Which is a contradiction.
Now fix $t \geq 0$ and let $\theta \in \hat{\Theta}$. Then $\theta \in N_{x}\left(\theta_{j_{1}}\right)$ for some $J_{1}$ and $\left\|\Phi_{\lambda}\left(\theta, T_{j_{1}}\right) x\right\|<$ $\left(\frac{1}{2}\right)^{2}$. In the same way $\theta \cdot t_{j_{1}} \in N_{x}\left(\theta_{j_{2}}\right)$ for some $J_{2}$ and

$$
\left\|\Phi_{\lambda}\left(\theta, T_{j_{1}}+T_{j_{2}}\right) x\right\|=\left\|\Phi_{\lambda}\left(\theta \cdot t_{j_{1}}, T_{j_{2}}\right) \Phi\left(\theta, T_{j_{2}}\right) x\right\|<\left(\frac{1}{2}\right)^{2}
$$

Now continue this process until one has

$$
\tau=T_{j_{1}}+\cdots+T_{j_{l}} \leq t<\tau+T_{j_{(l+1)}} \text { and }\left\|\Phi_{\lambda}(\theta, \mathrm{t}) \mathrm{x}\right\| \leq\left(\frac{1}{2}\right)^{1}
$$

Since $l T_{1} \leq \tau \leq t$ and $0 \leq t-\tau \leq \mathrm{T}_{\mathrm{m}}$, we get the following

$$
\left\|\Phi_{\lambda}(\theta, t) x\right\|=\left\|\Phi_{\lambda}(\theta \cdot \tau, t-\tau) \Phi_{\lambda}(\theta, \tau) x\right\| \leq k\left(\frac{1}{2}\right)^{l} \leq k\left(\frac{1}{2}\right)^{t / T_{1}}=k e^{-\alpha t}
$$

where $\alpha=-\frac{1}{T_{1}} \ln \left(\frac{1}{2}\right)>0$.
Therefore we have gotten the following

$$
\left\|\Phi_{\lambda}(\theta, t) x\right\| \leq k e^{-\alpha t}, \quad \theta \in \hat{\Theta}, \quad t \geq 0
$$

Since $k$ and $\alpha$ do not depends of $x$ we get that

$$
\left\|\Phi_{\lambda}(\theta, t)\right\| \leq k e^{-\alpha t}, \quad \theta \in \hat{\Theta}, \quad t \geq 0
$$

From here we get that the skew-product flow $\pi_{\lambda}=\left(\Phi_{\lambda}, \sigma\right)$ has ED over $\hat{\Theta}$, with projections $P(\theta)=I$, e.i., $\lambda \in \rho(\hat{\Theta})$. On the other hand, if $\mu \geq \lambda$, then

$$
\left\|\Phi_{\mu}(\theta, t)\right\|=e^{(\lambda-\mu) t}\left\|\Phi_{\lambda}(\theta, t)\right\| \leq k e^{-\alpha t}, \quad \theta \in \hat{\Theta}, \quad t \geq 0
$$

Therefore,

$$
\mu \in \rho(\hat{\Theta}) \text { and } \mathcal{S}_{\lambda}(\hat{\Theta})=\mathcal{S}_{\mu}(\hat{\Theta})=\mathcal{E}(\hat{\Theta})
$$

Lemma 3.2 Let $\hat{\Theta}$ be a compact invariant set in $\Theta$. Then the resolvent $\rho(\hat{\Theta})$ is open. Moreover

$$
\text { if } \quad \lambda \in \rho(\hat{\Theta}), \text { then } \mathcal{S}_{\lambda}=\mathcal{S}_{\mu} \text { and } \mathcal{U}_{\lambda}=\mathcal{U}_{\mu}
$$

for all $\mu$ in a neighborhood of $\lambda$.

Proof Fix $\lambda \in \rho(\hat{\Theta})$. Then by definition $\pi_{\alpha}$ admits an exponential dichotomy over $\hat{\Theta}$.

Hence there exists an invariant Projector $\mathrm{P}: \mathcal{E} \rightarrow \mathcal{E}$ and positive constants $k$ and $\beta$ such that

$$
\begin{gathered}
\left\|\Phi_{\lambda}(\theta, t) P(\theta) \Phi_{\lambda}^{-1}(\theta, s)\right\| \leq k e^{-\beta(t-s)}, \quad s \leq t \\
\left\|\Phi_{\lambda}(\theta, t)[I-P(\theta)] \Phi_{\lambda}^{-1}(\theta, s)\right\| \leq k e^{\beta(t-s)}, \quad s \geq t
\end{gathered}
$$

Claim. If $\|\lambda-\mu\|<\alpha$ where $\alpha=\beta / 2$, then
(a) $\left\|\Phi_{\mu}(\theta, t) P(\theta) \Phi^{-1}(\theta, s)\right\| \leq k e^{-\alpha(t-s)}, \quad t \geq s$
(b) $\| \Phi_{\mu}(\theta, t)[I-P(\theta)] \Phi_{\mu}^{-1}(\theta, s) \leq k e^{\alpha(t-s)}, \quad s \geq t$.

In fact:

$$
\Phi_{\mu}(\theta, t) P(\theta) \Phi_{\mu}^{-1}(\theta, s)=e^{(\lambda-\mu)(t-s)} \Phi_{\lambda}(\theta, t) P(\theta) \Phi_{\lambda}^{-1}(\theta, s)
$$

Therefore

$$
\left\|\Phi_{\mu}(\theta, t) P(\theta) \Phi_{\mu}^{-1}(\theta, s)\right\| \leq k e^{(\lambda-\mu)(t-s)} \cdot e^{\beta(t-s)}, \quad s \leq t
$$

From the fact that $-\beta / 2<\lambda-\mu<\beta / 2$ we get

$$
\left\|\Phi_{\mu}(\theta, t) P(\theta) \Phi^{-1}(\theta, s)\right\| \leq k e^{-\alpha(t-s)}, \quad s \geq t
$$

In the same way, we get that

$$
\left\|\Phi_{\mu}(\theta, t)[I-P(\theta)] \Phi^{-1}(\theta, s)\right\| \leq k e^{\alpha(t-s)}, \quad s \geq t
$$

Since the exponential dichotomies of $\pi_{\lambda}, \pi_{\mu}$ involve the same projector $\mathbf{P}$ we have that

$$
\mathcal{S}_{\lambda}(\hat{\Theta}) \text { and } \mathcal{U}_{\lambda}(\hat{\Theta})=\mathcal{U}_{\mu}(\hat{\Theta}) \text { for }\|\lambda-\mu\|<\alpha
$$

Lemma 3.3 Let $\hat{\Theta}$ be a compact invariant set in $\Theta$. Then the spectrum $\Sigma(\hat{\Theta})$ is compact. More specifically, there exists an $a>0$ such that, if $\lambda>a$, then $\lambda \in \rho(\hat{\Theta})$, and $\mathcal{S}_{\lambda}=\mathcal{E}$ and if $\lambda<-a$ then $\lambda \in \rho(\hat{\Theta})$ and $\mathcal{U}_{\lambda}=\mathcal{E}$.

Proof Because of Lemma 3.2 we need only to prove that $\Sigma(\hat{\Theta})$ is bounded. Thanks to the Proposition 2.1 we get $k \geq 1$ and $a>0$ such that

$$
\|\Phi(\theta, t)\| \leq e^{a|t|}, \quad \theta \in \hat{\Theta}, \quad t \in \mathbb{R}
$$

Then if $\lambda>a$ we get

$$
\left\|\Phi_{\lambda}(\theta, t)\right\| \leq k e^{(a-\lambda) t} \rightarrow 0, \text { as } t \rightarrow+\infty
$$

For all $\theta \in \hat{\Theta}$. Therefore, by Lemma 3.1 one has that

$$
(a, \infty) \subset \rho(\hat{\Theta}) \Leftrightarrow \Sigma(\hat{\Theta}) \subset(-\infty, a]
$$

Similarly, if $\lambda<-a$ one has $\left\|\Phi_{\lambda}(\theta, t)\right\| \rightarrow 0$ as $t \rightarrow-\infty$ for all $\theta \in \hat{\Theta}$. Consequently

$$
(-\infty,-a) \subset \rho(\hat{\Theta}) \Leftrightarrow \Sigma(\hat{\Theta}) \subset[-a, \infty)
$$

Hence

$$
\Sigma(\hat{\Theta}) \subset[-a, a]
$$

Lemma 3.4 Let $\hat{\Theta}$ be a nonempty set in $\Theta$ and assume that $\operatorname{dim} \mathcal{E} \geq 1$. Then the spectrum $\Sigma(\hat{\Theta})$ is nonempty.

Proof Pick $\theta_{0} \in \hat{\Theta}$ and set $M_{0}=H\left(\theta_{0}\right)$ where

$$
H\left(\theta_{0}\right)=\operatorname{cl}\left\{\theta_{0} . t: t \in \mathbb{R}\right\}
$$

Then $M_{0}$ is a compact invariant set and clearly $\Sigma\left(M_{0}\right) \subseteq \Sigma(\hat{\Theta})$. It will be sufficient to show that $\Sigma\left(M_{0}\right)$ is noempty. From Proposition 2.1, we have $k \geq 1$ and $a>0$ such that

$$
\|\Phi(\theta, t)\| \leq k e^{a|t|}, \quad \theta \in M_{0}, \quad t \in \mathbb{R}
$$

By Lemma 3.1 we get the following :
(a) If $\lambda>a$, then $\lambda \in \rho\left(M_{0}\right), \quad S_{\lambda}\left(M_{0}\right)=\mathcal{E}\left(M_{0}\right)$ and $\mathcal{U}_{\lambda}\left(M_{0}\right)=\mathcal{E}_{0}\left(M_{0}\right)$.
(b) If $\lambda<-a$, then $\lambda \in \rho\left(M_{0}\right), \mathcal{U}_{\lambda}\left(M_{0}\right)=\mathcal{E}\left(M_{0}\right)$ and $S_{\lambda}\left(M_{0}\right)=\mathcal{E}_{0}\left(M_{0}\right)$.

Therefore $\Sigma\left(M_{0}\right) \subset[-a, a]$.
Next define

$$
\lambda_{0}=\inf \left\{\lambda \in \rho\left(M_{0}\right): \mathcal{S}_{\lambda}\left(M_{0}\right)=\mathcal{E}\left(M_{0}\right)\right\}
$$

Then $-a \leq \lambda \leq a$.
For the purpose of contradiction, let us assume that $\lambda_{0} \in \rho\left(M_{0}\right)$. Then there are two cases to consider:
(i) $\mathcal{S}_{\lambda_{0}}\left(M_{0}\right)=\mathcal{E}\left(M_{0}\right)$
(ii) $\mathcal{S}_{\lambda_{0}}\left(M_{0}\right) \neq \mathcal{E}\left(M_{0}\right)$.

For the case (i) the Lemma 3.2 implies that $S_{\lambda}\left(M_{0}\right)=\mathcal{E}\left(M_{0}\right)$ for $\lambda$ in a neighborhood of $\lambda_{0}$, which contradicts the definition of $\lambda_{0}$.

For the case (ii) one must have $\mathcal{U}_{\lambda_{0}}\left(M_{0}\right) \neq \mathcal{E}_{0}\left(M_{0}\right)$. From Proposition 2.2 we get that

$$
\mathcal{E}(\Theta)=\mathcal{U}_{\lambda_{0}}(\theta)+\mathcal{S}_{\lambda_{0}}(\theta), \quad \forall \theta \in M_{0}
$$

Then by using Lemma 3.2 one again, we have $\mathcal{U}_{\lambda}\left(M_{0}\right) \neq \mathcal{E}_{0}\left(M_{0}\right)$ in a neighborhood of $\lambda_{0}$. This contradicts the fact that $\mathcal{U}_{\lambda}\left(M_{0}\right)=\mathcal{E}_{0}\left(M_{0}\right)$ for $\lambda>\lambda_{0}$ close enough to $\lambda_{0}$. Therefore $\lambda_{0} \in \Sigma\left(M_{0}\right)$.

Lemma 3.5 Let $\hat{\Theta}$ be a compact invariant set in $\Theta$. Consider $\lambda_{1}, \lambda_{2} \in \rho(\hat{\Theta})$ with $\lambda_{1}<\lambda_{2}$.

$$
\text { If } \quad \mathcal{S}_{\lambda_{1}}(\hat{\Theta})=\mathcal{S}_{\lambda_{2}}(\hat{\Theta}) \text { and } \mathcal{U}_{\lambda_{1}}(\hat{\Theta})=\mathcal{U}_{\lambda_{2}}(\hat{\Theta})
$$

then

$$
\left[\lambda_{1}, \lambda_{2}\right] \subset \rho(\hat{\Theta}) \text { and } S_{\lambda}(\hat{\Theta})=\mathcal{S}_{\lambda_{1}}(\hat{\Theta}), \quad \mathcal{U}_{\lambda}(\hat{\Theta})=\mathcal{U}_{\lambda_{1}}(\hat{\Theta})
$$

for all $\lambda \in\left[\lambda_{1}, \lambda_{2}\right]$.
Proof In the same way as the proof of Lemma 8 in [17].

The following Propositions are easy to prove.

Proposition 3.1 Let $A, B$ and $C$ be subspaces of $X$. If $C \subseteq A$ then

$$
A \cap(B+C)=(A \cap B)+(A \cap C)
$$

Proposition 3.2 If $\lambda_{1}, \lambda_{2} \in \rho(\hat{\Theta})$ and $\lambda_{1}<\lambda_{2}$, then

$$
\begin{equation*}
\mathcal{E}(\hat{\Theta})=\mathcal{U}_{\lambda_{1}}(\hat{\Theta}) \cap \mathcal{S}_{\lambda_{2}}(\hat{\Theta})+\mathcal{U}_{\lambda_{2}}(\hat{\Theta})+\mathcal{S}_{\lambda_{1}}(\hat{\Theta}) \tag{3.3}
\end{equation*}
$$

Proof By simplicity we shall use the following notation: $\mathcal{S}_{\boldsymbol{\lambda}}:=\mathcal{S}_{\boldsymbol{\lambda}}(\hat{\Theta})$ and $\mathcal{U}:=$ $\mathcal{U}_{\lambda}(\hat{\Theta})$.

It is easy to prove the following property:

$$
\begin{equation*}
\text { If } \lambda_{1} \leq \lambda_{2} \Rightarrow \mathcal{S}_{\lambda_{1}} \subseteq \mathcal{S}_{\lambda_{2}} \text { and } \mathcal{U}_{\lambda_{2}} \subseteq \mathcal{U}_{\lambda_{1}} \tag{3.4}
\end{equation*}
$$

From Proposition 2.2 we already know that

$$
\mathcal{E}(\hat{\Theta})=\mathcal{U}_{\lambda_{i}}+\mathcal{S}_{\lambda_{i}}, \quad i=1,2
$$

Therefore, from Proposition 3.1 we get

$$
\mathcal{U}_{\lambda_{1}}=\mathcal{U}_{\lambda_{1}} \cap\left(\mathcal{U}_{\lambda_{2}}+\mathcal{S}_{\lambda_{2}}\right)=\mathcal{U}_{\lambda_{2}}+\mathcal{U}_{\lambda_{1}} \cap \mathcal{S}_{\lambda_{2}}
$$

and

$$
\mathcal{S}_{\lambda_{2}}=\mathcal{S}_{\lambda_{2}} \cap\left(\mathcal{S}_{\lambda_{1}}+\mathcal{U}_{\lambda_{1}}\right)=\mathcal{S}_{\lambda_{1}}+\mathcal{U}_{\lambda_{1}} \cap \mathcal{S}_{\lambda_{2}}
$$

Then by using (3.2) we get that

$$
\mathcal{E}(\hat{\Theta})=\mathcal{U}_{\lambda_{1}}+\mathcal{S}_{\lambda_{1}}=\mathcal{U}_{\lambda_{2}}+\mathcal{S}_{\lambda_{1}}+\mathcal{U}_{\lambda_{1}} \cap \mathcal{S}_{\lambda_{2}}
$$

Lemma 3.6 Let $\hat{\Theta}$ be a compact invariant set in $\Theta$ and $\lambda_{1}, \lambda_{2} \in \rho(\hat{\Theta})$ with $\lambda_{1}<\lambda_{2}$.

Then the following statements are equivalent:
(A) There is a $\mu \in\left(\lambda_{1}, \lambda_{2}\right) \cap \Sigma(\hat{\Theta})$
(B) $\mathcal{U}_{\lambda_{1}}(\hat{\Theta}) \cap \mathcal{S}_{\lambda_{2}}(\hat{\Theta}) \neq \mathcal{E}_{0}(\hat{\Theta})$.

Moreover, $\mathcal{F}=\mathcal{U}_{\lambda_{1}}(\hat{\Theta}) \cap \mathcal{U}_{\lambda_{2}}(\hat{\Theta})$ is an invariant subbundle of $\mathcal{E}$

Proof $(A) \Rightarrow(B)$
From Proposition 3.2 we have the following

$$
\mathcal{E}(\hat{\Theta})=\mathcal{U}_{\lambda_{1}}(\hat{\Theta}) \cap \mathcal{S}_{\lambda_{2}}(\hat{\Theta})+\mathcal{U}_{\lambda_{2}}(\hat{\Theta})+\mathcal{S}_{\lambda_{1}}(\hat{\Theta})
$$

If $\mathcal{U}_{\lambda_{1}} \cap \mathcal{S}_{\lambda_{2}}=\mathcal{E}_{0}, \quad$ then we get

$$
\left.\mathcal{E}=\mathcal{U}_{\lambda_{2}}(\hat{\Theta})+\mathcal{S}_{\lambda_{1}}(\hat{\Theta})=\mathcal{U}_{\lambda_{1}}(\hat{\Theta})+\mathcal{S}_{\lambda_{1}}(\hat{\Theta})=\mathcal{U}_{\lambda_{2}}(\hat{\Theta})+\mathcal{S}_{\lambda_{2}} \hat{\Theta}\right)
$$

Since $\mathcal{U}_{\lambda_{2}}(\hat{\Theta}) \subseteq \mathcal{U}_{\lambda_{1}}(\hat{\Theta})$ and $\mathcal{S}_{\lambda_{1}}(\hat{\Theta}) \subseteq \mathcal{S}_{\lambda_{2}}(\hat{\Theta})$
then

$$
\mathcal{S}_{\lambda_{1}}(\hat{\Theta})=\mathcal{S}_{\lambda_{2}}(\hat{\Theta}) \text { and } \mathcal{U}_{\lambda_{1}}(\hat{\Theta})=\mathcal{U}_{\lambda_{2}}(\hat{\Theta})
$$

Now we can apply Lemma 3.5, it means that

$$
\mathcal{U}_{\lambda_{1}}(\hat{\Theta})=\mathcal{U}_{\lambda}(\hat{\Theta}) \text { and } S_{\lambda_{2}}(\hat{\Theta})=\mathcal{S}_{\lambda}(\hat{\Theta})
$$

for $\lambda \in\left[\lambda_{1}, \lambda_{2}\right] \subset \rho(\hat{\Theta})$. Therefore $\left(\lambda_{1}, \lambda_{2}\right) \cap \Sigma(\hat{\Theta})=\emptyset$ which contradicts $(A)$
$(B) \Rightarrow(A)$. Define

$$
\mu:=\inf \left\{\lambda \in \rho(\hat{\Theta}): \operatorname{cal}_{\lambda}(\hat{\Theta})=\mathcal{S}_{\lambda_{2}}(\hat{\Theta})\right\}
$$

Then $\lambda_{1}<\mu<\lambda_{2}$ and $\mu \in \Sigma(\hat{\Theta})$.
In fact. Lemma 3.2 implies that $\mu<\lambda_{2}$.
For the purpose of contradiction, let us assume that $\mu \in \rho(\hat{\Theta})$. Then there exists a neighborhood of $\mu$ such that for all $\lambda$ in that neighborhood $\mathcal{S}_{\lambda}(\hat{\Theta})=$ $\mathcal{S}_{\mu}(\hat{\Theta})$. Hence, $\mu \in \Sigma(\hat{\Theta})$.

Assume that $\mu<\lambda_{1}$. Then we get

$$
\mathcal{U}_{\mu}(\hat{\Theta}) \subseteq \mathcal{S}_{\lambda_{1}}(\hat{\Theta}) \subseteq \mathcal{S}_{\lambda_{2}}(\hat{\Theta})
$$

Then applying Lemma 3.2 and the definition of $\mu$ we get that $S_{\lambda_{1}}(\hat{\Theta})=S_{\lambda_{2}}(\hat{\Theta})$. From $(B)$ we have that $\mathcal{S}_{\lambda_{2}}(\hat{\Theta}) \cap \mathcal{U}_{\lambda_{1}}(\hat{\Theta}) \neq \mathcal{E}_{0}(\hat{\Theta})$. So $\mathcal{U}_{\lambda_{1}}(\hat{\Theta}) \cap \mathcal{S}_{\lambda_{1}}(\hat{\Theta}) \neq \mathcal{E}_{0}(\hat{\Theta})$. Which is a contradiction with the fact that $\lambda_{1} \in \rho(\hat{\Theta})$. Thus $\mu \in\left(\lambda_{1}, \lambda_{2}\right)$

Finally, since both $\mathcal{U}_{\lambda_{1}}(\hat{\Theta})$ and $\mathcal{S}_{\lambda_{2}}(\hat{\Theta})$ are invariant subbundles of $\mathcal{E}(\hat{\Theta})$, it follows that

$$
\mathcal{F}=\mathcal{U}_{\lambda_{1}}(\hat{\Theta}) \cap \mathcal{S}_{\lambda_{2}}(\hat{\Theta})
$$

is also invariant subbundle.

Lemma 3.7 Let $\hat{\Theta}$ be a compact invariant set in $\Theta$ and let $\lambda_{1}, \lambda_{2}$ be chosen so $\left(\lambda_{1}, \lambda_{2}\right) \cup \Sigma(\hat{\Theta}) \neq \phi$. Let

$$
\mathcal{F}:=\mathcal{U}_{\lambda_{1}}(\hat{\Theta}) \cap \mathcal{S}_{\lambda_{2}}(\hat{\Theta})
$$

and $\hat{\pi}$ the restriction of $\pi$ to $\mathcal{F}$. Let $\hat{\Sigma}(\hat{\Theta}):=\Sigma(\mathcal{F})$ denote the spectrum of $(\mathcal{F}, \hat{\pi})$ over $\hat{\Theta}$. Then

$$
\hat{\Sigma}(\hat{\Theta})=\Sigma(\hat{\Theta}) \cap\left(\lambda_{1}, \lambda_{2}\right)
$$

Proof We shall give the proof in three steps.
Step 1. Consider $\lambda \in \rho(\hat{\Theta})$ and define

$$
\mathcal{F}_{s}(\lambda):=\mathcal{F} \cap \mathcal{S}_{\lambda}(\hat{\Theta}) \text { and } \mathcal{F}_{u}(\lambda):=\mathcal{F} \cap \mathcal{U}_{\lambda}(\hat{\Theta})
$$

Then $\mathcal{F}_{s}(\lambda)$ and $\mathcal{F}_{u}(\lambda)$ are invariant subbundles and

$$
\mathcal{F}=\mathcal{F}_{\mathbf{a}}(\lambda)+\mathcal{F}_{u}(\lambda)
$$

In fact, suppose that $\lambda<\lambda_{1}$. Then

$$
\mathcal{S}_{\lambda}(\hat{\Theta}) \subset \mathcal{S}_{\lambda_{1}}(\hat{\Theta}) \text { and } \mathcal{U}_{\lambda_{1}}(\hat{\Theta}) \subset \mathcal{U}_{\lambda}(\hat{\Theta})
$$

From now on we shall omit the argument $\hat{\Theta}$ if it is neccesary, in order to simplify the computation. So

$$
\mathcal{F}_{\Delta}(\lambda)=\mathcal{F} \cap \mathcal{S}_{\lambda}=\mathcal{U}_{\lambda_{1}} \cap \mathcal{S}_{\lambda_{2}} \cap \mathcal{S}_{\lambda}=\mathcal{E}_{0} \cap \mathcal{S}_{\lambda_{2}}=\mathcal{F}_{0}
$$

and

$$
\mathcal{F}_{u}(\lambda)=\mathcal{F} \cap \mathcal{U}_{\lambda}=\mathcal{U}_{\lambda_{1}} \cap \mathcal{S}_{\lambda_{2}} \cap \mathcal{U}_{\lambda}=\mathcal{U}_{\lambda_{1}} \cap \mathcal{S}_{\lambda_{2}}=\mathcal{F}
$$

Hence $\mathcal{F}=\mathcal{F}_{s}(\lambda)+\mathcal{F}_{u}(\lambda)$.

Similarly, if $\lambda>\lambda_{2}$, then

$$
\mathcal{F}_{s}(\lambda)=\mathcal{F} \quad \text { and } \quad \mathcal{F}_{u}(\lambda)=\mathcal{F}_{0}
$$

For all $\lambda \in\left[\lambda_{1}, \lambda_{2}\right]$ we shall use Proposition 3.2 and the fact that $\mathcal{E}=\mathcal{U}_{\lambda_{2}}+\mathcal{S}_{\lambda_{2}}=$ $\mathcal{U}_{\lambda_{1}}+\mathcal{S}_{\lambda_{1}}$. So we have:

$$
\begin{aligned}
\mathcal{U}_{\lambda_{1}}=\mathcal{U}_{\lambda_{1}} \cap\left(\mathcal{U}_{\lambda_{2}}+\mathcal{S}_{\lambda_{2}}\right) & =\mathcal{U}_{\lambda_{2}}+\mathcal{F} \\
\mathcal{U}_{\lambda_{1}}=\mathcal{U}_{\lambda_{1}} \cap\left(\mathcal{U}_{\lambda}+\mathcal{S}_{\lambda}\right) & =\mathcal{U}_{\lambda}+\mathcal{F}_{s}(\lambda) \\
\mathcal{U}_{\lambda}=\mathcal{U}_{\lambda} \cap\left(\mathcal{U}_{\lambda_{2}}+\mathcal{S}_{\lambda_{2}}\right) & =\mathcal{U}_{\lambda_{2}}+\mathcal{F}_{u}(\lambda)
\end{aligned}
$$

Therefore

$$
\mathcal{U}_{\lambda_{1}}=\mathcal{U}_{\lambda_{2}}+\mathcal{F}=\mathcal{U}_{\lambda_{2}}+\mathcal{F}_{u}(\lambda)+\mathcal{F}_{s}(\lambda)
$$

Since $\quad \mathcal{F}_{u}(\lambda)+\mathcal{F}_{s}(\lambda) \subseteq \mathcal{F}$, then $\mathcal{F}=\mathcal{F}_{u}(\lambda)+\mathcal{F}_{s}(\lambda)$
Step 2. $\rho(\hat{\Theta}) \subseteq \hat{\rho}(\hat{\Theta})=\mathbb{R} \backslash \hat{\Sigma}(\hat{\Theta})$. In fact, if $\lambda \in \rho(\hat{\Theta})$, then there is a projector $\mathbf{P}: \mathcal{E}(\hat{\boldsymbol{\Theta}}) \rightarrow \mathcal{E}(\hat{\boldsymbol{\Theta}})$ and positive constants $k$ and $\beta$ such that

$$
\mathcal{R}(\mathbf{P})=\mathcal{S}_{\lambda}(\hat{\Theta}), \quad \mathcal{N}(\mathbf{P})=\mathcal{U}_{\lambda}(\hat{\Theta})
$$

and

$$
\begin{gathered}
\left\|\Phi_{\lambda}(\theta, t) P(\theta) \Phi_{\lambda}^{-1}(\theta, s) x\right\| \leq k\|x\| e^{-\beta(t-s)}, \quad t \geq s, \quad \theta \in \hat{\Theta} \\
\left\|\Phi_{\lambda}(\theta, t)(I-P(\theta)) \Phi_{\lambda}^{-1}(\theta, s) x\right\| \leq k\|x\| e^{\beta(t-s)}, \quad t \leq s, \quad \theta \in \hat{\Theta}
\end{gathered}
$$

If $\hat{\mathbf{P}}$ is the restiction of $\mathbf{P}$ to $\mathcal{F}$, then

$$
\mathcal{R}(\hat{\mathbf{P}})=\mathcal{F}_{s}(\lambda), \quad \mathcal{N}(\hat{\mathbf{P}})=\mathcal{F}_{u}(\lambda), \quad \mathcal{F}=\mathcal{R}(\hat{\mathbf{P}})+\mathcal{N}(\hat{\mathbf{P}})
$$

Therefore, by restricting the above inequalities to all $(x, \theta) \in \mathcal{F}$ we obtain

$$
\begin{gathered}
\left\|\hat{\Phi}_{\lambda}(\theta, t) \hat{P}(\theta) \hat{\Phi}_{\lambda}^{-1}(\theta, s) x\right\| \leq k\|x\| e^{-\beta(t-s)}, \quad t \geq s \\
\left\|\hat{\Phi}_{\lambda}(\theta, t)(I-\hat{P}(\theta)) \hat{\Phi}_{\lambda}^{-1}(\theta, s) x\right\| \leq k\|x\| e^{\beta(t-s)}, \quad t \leq s
\end{gathered}
$$

Thus $\hat{\pi}_{\lambda}$ admits an exponential dichotomy over $\hat{\Theta}$, i.e., $\lambda \in \hat{\rho}(\hat{\Theta})$. So $\rho(\hat{\Theta}) \subset \hat{\rho}(\hat{\Theta})$ and therefore $\hat{\Sigma}(\hat{\Theta}) \subset \Sigma(\hat{\Theta})$. This completes the step 1 .

Since $\mathcal{F} \subseteq \mathcal{S}_{\lambda_{2}}(\hat{\Theta})$, then

$$
\left\|\hat{\Phi}_{\lambda_{2}}(\theta, t) x\right\| \rightarrow 0 \quad \text { as } \quad t \rightarrow \infty
$$

for all $(x, \theta) \in \mathcal{F}$ uniformly in $x$. Hence, applying Lemma 3.1 we obtain that $\hat{\Sigma}(\hat{\Theta}) \subseteq\left(-\infty, \lambda_{2}\right)$ and $\mathcal{F}_{s}(\lambda)=\mathcal{F}, \quad \lambda \geq \lambda_{2}$. In the same way we obtain $\mathcal{F} \subseteq$ $\mathcal{U}_{\lambda_{1}}(\hat{\Theta})$. This implies that

$$
\left\|\hat{\Phi}_{\lambda_{1}}(\theta, t)\right\| \rightarrow 0, \quad \text { as } t \rightarrow-\infty
$$

## Using Lemma 3.1 again, we get that

$\hat{\Sigma}(\hat{\Theta}) \subseteq\left(\lambda_{1}, \infty\right)$ and $\mathcal{F}_{u}=\mathcal{F}$. Then

$$
\hat{\Sigma}(\hat{\Theta}) \subseteq\left(\lambda_{1}, \lambda_{2}\right) \cap \Sigma(\hat{\Theta})
$$

In order to prove the opposite inclusion, it is sufficient to show that: if $\lambda \in$ $\hat{\rho}(\hat{\Theta}) \cap\left(\lambda_{1}, \lambda_{2}\right), \quad$ then $\lambda \in \rho(\hat{\Theta})$. In fact, suppose that

$$
\hat{\rho}(\hat{\Theta}) \cap\left(\lambda_{1}, \lambda_{2}\right) \subset \rho(\hat{\Theta})
$$

which is equivalent to :

$$
\rho^{c}(\hat{\Theta}) \subset(\hat{\rho}(\hat{\Theta}))^{c} \cup\left(-\infty, \lambda_{1}\right] \cup\left[\lambda_{2}, \infty\right) \Leftrightarrow \Sigma(\hat{\Theta}) \subset \hat{\Sigma}(\hat{\Theta}) \cup\left(-\infty, \lambda_{1}\right] \cup\left[\lambda_{2}, \infty\right)
$$

On the other hand, we already know that $\hat{\Sigma}(\hat{\Theta}) \subset\left(\lambda_{1}, \lambda_{2}\right)$. Therefore,

$$
\Sigma(\hat{\Theta}) \cap\left(\lambda_{1}, \lambda_{2}\right) \subset \hat{\Sigma}(\hat{\Theta})
$$

Step 3. $\hat{\rho}(\hat{\Theta}) \cap\left(\lambda_{1}, \lambda_{2}\right) \subset \rho(\hat{\Theta})$. In fact, if $\lambda \in \hat{\rho}(\hat{\Theta}) \cap\left(\lambda_{1}, \lambda_{2}\right)$, then there is a projector $\mathbf{Q}: \mathcal{F} \rightarrow \mathcal{F}$ and positive contants $k$ and $\beta$ such that

$$
\begin{gather*}
\left\|\hat{\Phi}_{\lambda}(\theta, t) Q(\theta) \hat{\Phi}_{\lambda}^{-1}(\theta, s) x\right\| \leq k\|x\| e^{-\beta(t-s)}, \quad t \geq s  \tag{3.5}\\
\left\|\hat{\Phi}_{\lambda}(\theta, t)(I-Q(\theta)) \hat{\Phi}_{\lambda}^{-1}(\theta, s) x\right\| \leq k\|x\| e^{\beta(t-s)}, \quad t \leq s \tag{3.6}
\end{gather*}
$$

and for all $\theta \in \hat{\Theta}$
Consider the projectors:

$$
\mathbf{P}_{1}, \mathbf{P}: \mathcal{E}(\hat{\Theta}) \rightarrow \mathcal{E}(\hat{\Theta}) \quad \text { such that } \quad \mathcal{R}\left(\mathbf{P}_{1}\right)=\mathcal{S}_{\lambda_{1}}, \quad \mathcal{N}\left(\mathbf{P}_{1}\right)=\mathcal{U}_{\lambda_{1}}
$$

and

$$
\mathcal{R}(\mathbf{P})=\mathcal{F}=\mathcal{U}_{\lambda_{1}} \cap \mathcal{S}_{\lambda_{2}}, \quad \mathcal{N}(\mathbf{P})=\mathcal{U}_{\lambda_{2}}+\mathcal{S}_{\lambda_{1}} .
$$

From the Proposition 3.2 we have

$$
\mathcal{E}=\mathcal{U}_{\lambda_{2}}+\mathcal{S}_{\lambda_{1}}+\mathcal{U}_{\lambda_{1}} \cap \mathcal{S}_{\lambda_{2}}=\mathcal{R}(\mathbf{P})+\mathcal{N}(\mathbf{P})
$$

$\mathbf{P}_{\lambda}=\mathbf{P}_{1}+\mathbf{Q P}$ is a projector on $\mathcal{E}$ such that

$$
\mathcal{R}\left(\mathbf{P}_{\lambda}\right)=\mathcal{S}_{\lambda_{1}}+\hat{\mathcal{S}}_{\lambda}, \quad \mathcal{N}\left(\mathbf{P}_{\lambda}\right)=\hat{\mathcal{U}}_{\lambda}+\mathcal{U}_{\lambda_{2}}
$$

where $\quad \hat{\mathcal{S}}_{\lambda}=\mathcal{R}(\mathbf{Q}) \quad$ and $\quad \hat{\mathcal{U}}_{\lambda}=\mathcal{N}(\mathbf{Q})$.
Using (3.5) and (3.6) and the fact that $\lambda \in\left(\lambda_{1}, \lambda_{2}\right), \lambda_{1}, \lambda_{2} \in \rho(\hat{\Theta})$ one can show that there are positive constants $l$ and $\alpha$ such that

$$
\begin{gathered}
\left\|\Phi_{\lambda}(\theta, t) P_{\lambda}(\theta) \Phi_{\lambda}^{-1}(\theta, s)\right\| \leq l e^{-\alpha(t-s)}, \quad t \geq s, \quad \theta \in \hat{\Theta} \\
\left\|\Phi_{\lambda}(\theta, t)\left(I-P_{\lambda}(\theta)\right) \Phi_{\lambda}^{-1}(\theta, s)\right\| \leq l e^{\alpha(t-s)}, \quad t \leq s, \quad \theta \in \hat{\Theta}
\end{gathered}
$$

noindent which implies that $\lambda \in \rho(\hat{\Theta})$.

### 3.2 Proof of Main Theorems

## Proof of Theorem 3.1

The statement (A), (B) and (C) follow from Lemma 3.3, 3.4 and 3.6 respectively. The statement (D) follows from Lemma 3.7.
(E) From Lemma 3.3 we have that

$$
\Sigma(\hat{\Theta}) \subset[-a, a]
$$

Therefore

$$
\Sigma(\hat{\Theta})=\Sigma(\hat{\Theta}) \cap\left(\bigcup_{i=1}^{m}\left(\lambda_{i-1}, \lambda_{i}\right)\right)=\bigcup_{i=1}^{m} \Sigma(\hat{\Theta}) \cap\left(\lambda_{i-1}, \lambda_{i}\right)=\bigcup_{i=1} \Sigma_{i}(\hat{\Theta})
$$

(F) Consider $i+1 \leq j$

$$
\mathcal{V}_{i}=\mathcal{U}_{\lambda_{i}-1} \cap \mathcal{S}_{\lambda_{i}} \Rightarrow \mathcal{V}_{i} \subseteq \mathcal{S}_{\lambda_{i}},
$$

by the monotocity of $\mathcal{U}_{\lambda}$ we get that $\mathcal{V}_{j} \subset \mathcal{U}_{\lambda_{i}}$. On the other hand we know that $\mathcal{U}_{\lambda_{i}} \cap \mathcal{S}_{\lambda_{i}}=\mathcal{E}_{0}$.
(G) From Lemma 3.3 we have that:
if $\lambda>a \Rightarrow \lambda \in \rho(\hat{\Theta}), \quad \mathcal{S}_{\lambda}(\hat{\Theta})=\mathcal{E}(\hat{\Theta})$ and $\mathcal{U}_{\lambda}(\hat{\Theta})=\mathcal{E}_{0}(\hat{\Theta})$.
if $\lambda<-a \Rightarrow \lambda \in \rho(\hat{\Theta}), \quad \mathcal{U}_{\lambda}(\hat{\Theta})=\mathcal{E}(\hat{\Theta})$ and $\mathcal{S}_{\lambda}(\hat{\Theta})=\mathcal{E}_{0}(\hat{\Theta})$.
Also we know that

$$
\mathcal{E}(\hat{\Theta})=\mathcal{S}_{\lambda_{i}}(\hat{\Theta})+\mathcal{U}_{\lambda_{i}}(\hat{\Theta})
$$

Therefore

$$
\mathcal{E}(\hat{\Theta})=\mathcal{U}_{\lambda_{0}}=\mathcal{U}_{\lambda_{0}} \cap\left(\mathcal{S}_{\lambda_{1}}+\mathcal{U}_{\lambda_{1}}\right)
$$

$$
\begin{aligned}
& =\mathcal{U}_{\lambda_{0}} \cap \mathcal{S}_{\lambda_{1}}+\mathcal{U}_{\lambda_{1}} \\
& =\mathcal{V}_{1}+\mathcal{U}_{\lambda_{1}}=\mathcal{V}_{1}+\mathcal{U}_{\lambda_{1}} \cap\left(\mathcal{S}_{\lambda_{2}}+\mathcal{U}_{\lambda_{2}}\right) \\
& =\mathcal{V}_{1}+\mathcal{U}_{\lambda_{1}} \cap \mathcal{S}_{\lambda_{2}}+\mathcal{U}_{\lambda_{2}} \\
& =\mathcal{V}_{1}+\mathcal{V}_{2}+\mathcal{U}_{\lambda_{2}} \cap\left(\mathcal{S}_{\lambda_{3}}+\mathcal{U}_{\lambda_{3}}\right) \\
& =\ldots \ldots \ldots \ldots \ldots \\
& =\ldots \ldots \ldots \ldots \ldots \\
& =\ldots \ldots \ldots \ldots \ldots \\
& =\mathcal{V}_{1}+\mathcal{V}_{2}+\mathcal{V}_{3}+\ldots+\mathcal{V}_{m}+\mathcal{U}_{\lambda_{m}} \\
& =\mathcal{V}_{1}+\mathcal{V}_{2}+\ldots+\mathcal{V}_{m} . \quad\left(\mathcal{U}_{\lambda_{m}}=\mathcal{E}_{0}(\hat{\Theta})\right)
\end{aligned}
$$

## Proof of Theorem 3.2

(A) - (F) follow easily from Lemmas 3.2 and 3.3. Proof of (G), from Lemma 3.3 we have the following:

If $\lambda_{0}<-a$, then $\mathcal{U}_{\lambda_{0}}=\mathcal{E}$ and $\operatorname{dim} \mathcal{U}_{\lambda_{0}}=\infty$
If $\lambda_{0}>a$, then $\mathcal{U}_{\lambda_{0}}=\mathcal{E}_{0}$ and $\operatorname{dim} \mathcal{U}_{\lambda_{0}}=0$
Therefore $\lambda_{0} \in[-a, a]$.
Now consider the set

$$
\left\{\lambda_{0}<\lambda_{1}<\lambda_{2}<\ldots<\lambda_{m}\right\} \subset \rho(\hat{\Theta})
$$

such that

$$
a<\lambda_{m}, \quad\left(\lambda_{i-1}, \lambda_{i}\right) \cap \Sigma(\hat{\Theta}) \neq \phi
$$

From Proposition 2.2. we get that: $\mathcal{E}(\hat{\Theta})=\mathcal{S}_{\lambda_{0}}(\hat{\Theta})+\mathcal{U}_{\lambda_{0}}(\hat{\Theta})$. Then let us denote by $\pi^{\lambda_{0}}$ the restriction of $\Pi$ to $\mathcal{U}_{\lambda_{0}}$ and $\Sigma^{\lambda_{0}}(\hat{\Theta})$ the spectrum of $\left(\mathcal{U}_{\lambda_{0}}, \pi^{\lambda_{0}}\right)$ over $\hat{\Theta}$.

So by using Lemma 3.7 we obtain that

$$
\Sigma^{\lambda_{0}}(\hat{\Theta})=\left(\lambda_{0},+\infty\right) \cap \Sigma(\hat{\Theta})=\left(\lambda_{0}, a\right] \cap \Sigma(\hat{\Theta})
$$

Therefore $\operatorname{dim} \mathcal{V}_{i} \geq 1$, which implies that $m \leq \operatorname{dim} \mathcal{U}_{\lambda_{0}}$. Now we shall show that the resolvent $\rho(\hat{\Theta})$ of $\left(\mathcal{U}_{\lambda_{0}}, \pi^{\lambda_{0}}\right)$ consist of $(k+1)$ intervals where $k \leq n\left(\lambda_{0}\right)$. If $\rho_{\lambda_{0}}(\hat{\Theta})$ consists of $\left(n\left(\lambda_{0}\right)+2\right)$ intervals :

$$
\rho_{\lambda_{0}}(\hat{\Theta})=\bigcup_{i=1}^{n\left(\lambda_{0}\right)+2}\left(c_{i}, d_{i}\right)
$$

So we get that

$$
\Sigma_{\lambda_{0}}(\hat{\Theta}) \cap\left(\lambda_{i-1}, \lambda_{i}\right) \neq \phi
$$

Then

$$
\mathcal{V}_{i}=\mathcal{U}_{\lambda_{i-1}} \cap \mathcal{S}_{\lambda_{i}} \neq \mathcal{E}_{0}(\hat{\Theta}), \quad \text { in }\left(\mathcal{U}_{\lambda_{0}}, \pi^{\lambda_{0}}\right)
$$

Therefore

$$
\operatorname{dim} \mathcal{U}_{\lambda_{0}} \geq n_{1}+n_{2}+\ldots+N_{n+1} \geq n+1
$$

Which is a contradiction. Hence $\rho_{\lambda_{0}}(\hat{\Theta})$ consists of $k+1$ intervals whith $k \leq n$. So $\Sigma_{\lambda_{0}}(\hat{\Theta})$ is the union of $k$ compact intervals. The remainder of the proof is easy .

## 4 Lyapunov Exponents

In this section we shall investigate the relation between the Dynamic Spectrum and the Lyapunov characteristic exponents.

For this purpose we shall assume that there exists $\lambda_{0} \in \rho(\hat{\Theta})$ such that $1 \leq$ $\operatorname{dim} \mathcal{U}_{\lambda_{0}}<\infty$.

Consider $\left\{\lambda_{0}<\ldots<\lambda_{m}\right\} \subset \rho(\hat{\Theta})$, such that $0<\lambda_{m}$ and $\Sigma(\hat{\Theta}) \cap\left(\lambda_{i-1}, \lambda_{i}\right) \neq \phi$. Then from Theorem 3.2 we get:

$$
\begin{gather*}
\Sigma_{\lambda_{0}}(\hat{\Theta})=\left(\lambda_{0}, a\right] \cap \Sigma(\hat{\Theta})  \tag{4.1}\\
\Sigma(\hat{\Theta})=\Sigma_{\lambda_{0}}(\hat{\Theta}) \cup\left(U_{i=1}^{m}\left[a_{i}, b_{i}\right]\right)  \tag{4.2}\\
m \leq \operatorname{dim} \mathcal{U}_{\lambda_{0}}, \quad \lambda_{0} \in[-a, a]  \tag{4.3}\\
\mathcal{E}(\hat{\Theta})=S_{\lambda_{0}}(\hat{\Theta})+\mathcal{V}_{1}(\hat{\Theta})+\ldots+\mathcal{V}_{m}(\hat{\Theta})  \tag{4.4}\\
\mathcal{V}_{i}=S_{\lambda_{i}} \cap \mathcal{U}_{\lambda_{i}-1}, \quad i=1,2, \ldots, m  \tag{4.5}\\
\text { If } \Sigma(\hat{\Theta}) \subset\left[\lambda_{0}, a\right] \Rightarrow \Sigma(\hat{\Theta})=U_{i=1}^{m}\left[a_{i}, b_{i}\right] \tag{4.6}
\end{gather*}
$$

The spectral intervals $\left[a_{i}, b_{i}\right]$ have been ordered so that $b_{i} \leq a_{i+1}, i=$ $1,2, \ldots, m-1$.

Let $\quad P_{i}: \mathcal{E}(\hat{\Theta}) \rightarrow \mathcal{E}(\hat{\Theta})$ denote a projector corresponding with the descomposition (4.4) such that Range $\left(\mathbf{P}_{\mathbf{i}}\right)=\mathcal{R}\left(\mathbf{P}_{\mathbf{i}}\right)=\mathcal{V}_{\mathbf{i}}$ and the null space being the sum of the remain $\mathcal{V}_{j}$ and $\mathcal{S}_{\lambda_{0}}$ for $j \neq i$

Then if $\mathrm{P}_{\lambda_{\mathbf{0}}}: \mathcal{E}(\hat{\Theta}) \rightarrow \mathcal{E}(\hat{\Theta})$ is the projector on $\mathcal{E}$ such that

$$
\mathcal{R}\left(\mathbf{P}_{\lambda_{0}}\right)=\mathcal{S}_{\lambda_{0}} \text { and } \mathcal{N}\left(\mathbf{P}_{\lambda_{0}}\right)=\mathcal{V}_{1}+\mathcal{V}_{2}+\ldots+\mathcal{V}_{m}
$$

Hence

$$
I=P_{\lambda_{0}}(\theta)+P_{1}(\theta)+\ldots+P_{m}(\theta), \quad \forall \theta \in \hat{\Theta}
$$

Given a point $(x, \theta) \in \mathcal{E}, x \neq 0$, we shall define the four Lyapunov exponents of $(x, \theta)$ as follows:

$$
\begin{align*}
& \lambda_{s}^{+}(x, \theta)=\varlimsup_{t \rightarrow+\infty} \frac{1}{t} \ln \|\Phi(\theta, t) x\|  \tag{4.7}\\
& \lambda_{i}^{+}(x, \theta)=\varliminf_{t \rightarrow+\infty} \frac{1}{t} \ln \|\Phi(\theta, t) x\| \tag{4.8}
\end{align*}
$$

$$
\begin{equation*}
\lambda_{s}^{-}(x, \theta)=\varlimsup_{t \rightarrow-\infty} \frac{1}{t} \ln \|\Phi(\theta, t) x\| \tag{4.9}
\end{equation*}
$$

$$
\begin{equation*}
\lambda_{i}^{-}(x, \theta)=\varliminf_{t \rightarrow-\infty} \frac{1}{t} \ln \|\Phi(\theta, t) x\| \tag{4.10}
\end{equation*}
$$

Theorem 4.1 (A) If $(x, \theta) \in \mathcal{V}_{i}$ where $\mathcal{V}_{i}$, is the spectral subbundle associated with $\lambda_{0}$ and $\left[a_{i}, b_{i}\right]$, and $x \neq 0$, then the four Lyapunov exponents agree and the limits

$$
\lim _{t \rightarrow+\infty} \frac{1}{t} \ln \|\Phi(\theta, t) x\|=\lim _{t \rightarrow-\infty} \frac{1}{t} \ln \|\Phi(\theta, t) x\|
$$

exist and are equal to $a_{i}$.
(B) If $(x, \theta) \in \mathcal{S}_{\lambda_{0}}, \quad x \neq 0$ then the two Lyapunov exponents (4.7) - (4.8) agree and the limits

$$
\varlimsup_{t \rightarrow+\infty} \frac{1}{t} \ln \|\Phi(\theta, t) x\|=\underline{\lim }_{t \rightarrow+\infty} \frac{1}{t} \ln \|\Phi(\theta, t) x\|=\lambda_{0}
$$

Proof The proof of (A) is similar to the prove of Theorem 3 in []. In order to prove (B) let us consider $(x, \theta) \in S_{\lambda_{0}}$ with $x \neq 0$. Then

$$
\lim _{t \rightarrow \infty}\left\|\Phi_{\lambda_{0}}(\theta, t) x\right\|=\lim _{t \rightarrow+\infty}\left\|e^{-\lambda_{0} t} \Phi(\theta, t) x\right\|=0
$$

Therefore, there is a constant $M>0$ such that

$$
\ln \left\|e^{-\lambda_{0} t} \Phi(\theta, t) x\right\| \leq M, \quad t \geq 0
$$

Then

$$
0=\lim \frac{1}{t} \ln \left\|e^{-\lambda_{0} t} \Phi(\theta, t) x\right\|=\lim _{t \rightarrow \infty}\left[-\lambda_{0}+\frac{1}{t} \ln \|\Phi(\theta, t) x\|\right]
$$

So

$$
\lim _{t \rightarrow \infty} \frac{1}{t} \ln \|\Phi(\theta, t) x\|=\lambda_{0}
$$

Definition 4.1 For all $\theta \in \Theta$ we define the upper Lyapunov exponent $\lambda_{s}^{+}(\theta)$ and the lower Lyapunov exponent $\lambda_{i}^{+}(\theta)$ as follows:

$$
\begin{aligned}
& \lambda_{s}^{+}(\theta):=\sup \left\{\lambda_{s}^{+}(x, \theta): x \in X, \quad x \neq 0\right\} \\
& \lambda_{s}^{+}(\theta):=\inf \left\{\lambda_{i}^{+}(x, \theta): x \in X, \quad x \neq 0\right\}
\end{aligned}
$$

Theorem 4.2 The upper Lyapunov exponent $\lambda_{s}^{+}(\theta)$ and the lower Lyapunov exponent $\lambda_{i}^{+}(\theta)$ associated to $\theta \in \Theta$ are given by :

$$
\begin{gather*}
\lambda_{s}^{+}(\theta)=\varlimsup_{t \rightarrow \infty} \frac{1}{t} \ln \|\Phi(\theta, t)\|  \tag{4.11}\\
\lambda_{i}^{+}(\theta)=-\varlimsup_{t \rightarrow \infty} \frac{1}{t} \ln \|\Phi(\theta \cdot(-t), t)\| \tag{4.12}
\end{gather*}
$$

Proof We shall denote by $\gamma(\theta)$ the right side of (4.11). We have the following

$$
\frac{1}{t} \ln \|\Phi(\theta, t) x\| \leq \frac{1}{t} \ln \|\Phi(\theta, t)\|+\frac{1}{t} \ln \|x\|
$$

which implies that

$$
\varlimsup_{t \rightarrow \infty} \frac{1}{t} \ln \|\Phi(\theta, t) x\| \leq \varlimsup_{t \rightarrow \infty} \frac{1}{t} \ln \|\Phi(\theta, t)\|
$$

Hence $\lambda_{;}^{+}(x, \theta) \leq \gamma(\theta)$, for all $x \in X, x \neq 0$; therefore $\lambda_{s}^{+}(\theta) \leq \gamma(\theta)$.
In order to prove the opposite inequality we shall use the following fact: For each $\epsilon>0$ and $x \in X$ there is $N_{\epsilon, x}>0$ such that

$$
\|\Phi(\theta, t) x\| \leq N_{\epsilon, x} e^{\left(\lambda_{b}^{+}(\theta)+\epsilon\right) t}, \quad t \geq 0
$$

The above inequality can be writing as follows

$$
\left\|\Phi(\theta, t) x e^{-\left(\lambda_{0}^{+}(\theta)+\epsilon\right) t}\right\| \leq N_{\epsilon, x}
$$

Then the operators family

$$
\left\{\Phi(\theta, t) e^{-\left(\lambda_{5}^{+}(\theta)+\epsilon\right) t}: \quad t \geq 0\right\}
$$

is bounded for each $x \in X$. It follows from Uniform Boundness Theorem that there exists $N_{\epsilon}>0$ such that

$$
\left\|\Phi(\theta, t) e^{-\left(\lambda_{\dot{\theta}}^{+}(\theta)+\epsilon\right) t}\right\| \leq N_{\epsilon}, \quad t \geq 0
$$

Hence

$$
\|\Phi(\theta, t)\| \leq N_{\epsilon} e^{\left(\lambda_{0}^{+}(\theta)+\epsilon\right) t}, \quad t \geq 0
$$

Therefore

$$
\frac{1}{t} \ln \|\Phi(\theta, t) x\| \leq \frac{1}{t} \ln \left(N_{\epsilon}\right)+\left(\lambda_{s}^{+}(\theta)+\epsilon\right)
$$

So

$$
\varlimsup_{t \rightarrow \infty} \frac{1}{t} \ln \|\Phi(\theta, t) x\| \leq \lambda_{s}^{+}(\theta)+\epsilon
$$

Hence

$$
\gamma(\theta) \leq \lambda_{s}^{+}(\theta)
$$

The proof of (4.12) is not hard to do.

## 5 Examples

In this section we shall present some examples in which we are able to locate the Dynamical Spectrum.

The following example is sufficiently general to be considered here.

Example 5.1 In this example we shall consider a linear time dependent differential equation that generates a linear skew-product flow on the trivial Banach bundle $\mathcal{E}=X \times \Theta$ where $X$ is a Banach space and $\Theta$ is a compact topological Hausdorff space.

More precisely, we shall study the following linear time dependent system

$$
\begin{equation*}
\dot{x}(t)=A(\theta \cdot t) x(t), \quad t \in \mathbb{R} \tag{5.1}
\end{equation*}
$$

where $A(\theta \cdot t)=A+B(\theta \cdot t), \quad A$ is the infinitesimal generator of the strongly continuous group $\{T(t)\}_{t \in \mathbb{R}} ; \quad \sigma(\theta, t)=\theta \cdot t \quad$ is a flow on $\quad \Theta$ and $B(\theta) \in$ $L(X), \quad t \in \mathbb{R}$.

Lemma 5.1 If $B(\cdot): \Theta \rightarrow L(X)$ is strongly continuous, then the set $\{\|B(\theta)\|:$ $\theta \in \Theta\}$ is bounded.

Proof Consider the following sets

$$
H=\{\|B(\theta)\|: \theta \in \Theta\}, \quad H(x)=\{\|B(\theta) x\|: \theta \in \Theta\}
$$

Since $\theta \rightarrow B(\theta) x$ is continuous and $\Theta$ is compact, then for each $x \in X$ we get that $H(x)$ is bounded. Hence, by the Uniform Boundedness Theorem we obtain that $H$ is bounded.

Lemma 5.2 If $B(\cdot): \Theta \rightarrow L(X)$ is strongly continuous and $x(\cdot): \mathbb{R} \rightarrow X$ is a continuous function, then for each $\theta \in \Theta$ the mapping $t \rightarrow B(\theta \cdot t) x(t)$ is continuous.

Proof Fix $t \in \mathbb{R}$. Then

$$
\begin{gathered}
\|B(\theta(t+h)) x(t+h)-B(\theta \cdot t) x(t)\| \\
=\|B(\theta(t+h))[x(t+h)-x(t)]-[B(\theta \cdot(t+h))-B(\theta \cdot t)] x(t)\| \\
\leq L\|x(t+h)-x(t)\|+\|[B(\theta \cdot(t+h))-B(\theta \cdot t)] x(t)\|
\end{gathered}
$$

where $L=\sup \{\|B(\theta)\|: \theta \in \Theta\} \quad$ and $\quad \theta . t \in \Theta$ for $t \in \mathbb{R}$.
Therefore

$$
\|B(\theta(t+h)) x(t+h)-B(\theta \cdot t) x(t)\| \rightarrow 0, \quad \text { as } \quad h \rightarrow 0
$$

To be precise in which sense the equation (5.1) generates a linear skew-product flow, we shall consider the following family of integral differential equations:

$$
\begin{equation*}
x(t)=T(t) x_{0}+\int_{0}^{t} T(t-s) B(\theta \cdot s) x(s) d s . \quad t \in \mathbb{R} \quad \theta \in \Theta \tag{5.2}
\end{equation*}
$$

Definition 5.1 A solution $x(t)=x(t, \theta)$ of the equation (5.2) is called a Mild Solution of (5.1).

Theorem 5.1 Let $A$ be the infinitesimal generator of an strongly continuous group $\{T(t)\}_{t \in \mathbb{R}}$ on $X$ and $B(\cdot): \Theta \rightarrow L(X)$ is also strongly continuous. Then for each $\theta \in \Theta$ and $x_{0} \in X$ the problem

$$
\begin{equation*}
\dot{x}(t)=A(\theta \cdot t) x=(A+B(\theta \cdot t)) x(t) ; \quad x(0)=x_{0} \tag{5.3}
\end{equation*}
$$

has a unique mild solution $\Phi(\theta, t) x_{0}$ given by

$$
\begin{equation*}
\Phi(\theta, t) x_{0}=T(t) x_{0}+\int_{0}^{t} T(t-s) B(\theta \cdot s) \Phi(\theta, s) x_{0} d s \tag{5.4}
\end{equation*}
$$

If

$$
\|T(t)\| \leq M e^{W|t|}, \quad t \in \mathbb{R}
$$

then

$$
\begin{equation*}
\|\Phi(\theta, t)\| \leq M e^{(W+L M)|t|}, \quad t \in \mathbb{R} \tag{5.5}
\end{equation*}
$$

where

$$
L=\sup \{\|B(\theta)\|: \theta \in \Theta\}
$$

Moreover, the mapping $\pi: \mathcal{E} \times \mathbb{R} \rightarrow \mathcal{E}$ given by

$$
\begin{equation*}
\pi(x, \theta, t)=(\Phi(\theta, t) x, \theta \cdot t) \tag{5.6}
\end{equation*}
$$

is a linear skew-product flow on $\mathcal{E}=X \times \Theta$, maybe without the condition (4) of Definition 2.1.

The following Proposition corresponds to the part (4) of Definition 2.1

Proposition 5.1 If the flow $\sigma(\theta, t)=\theta \cdot t$ depends continuously on $\theta$ in compacts intervals, then for all fixed $t \in \mathbb{R}$, the mapping from $\mathcal{E}$ to $X$ given by $(x, \theta) \rightarrow$ $\Phi(\theta, t) x$ is continuous.

Proposition 5.2 Let $\pi=(\Phi, \sigma)$ be the linear skew-product flow defined by (5.6). Then

$$
\Sigma(\hat{\Theta}) \subset[-(W+L M), W+L M]
$$

where

$$
L=\sup \{\|B(\theta)\|: \theta \in \Theta\}, \quad\|T(t)\| \leq M e^{W|t|}, \quad t \in \mathbb{R}
$$

Example 5.2 (The perturbed wave equation). The simplest one dimensional hyperbolic equation with a time dependent perturbation of the form

$$
\begin{align*}
u_{t t} & =u_{x x}+b(t) u_{x}, \quad t \in \mathbb{R},-\infty<x<\infty  \tag{5.7}\\
u(0, x) & =u_{1}(x) \\
u_{t}(0, x) & =u_{2}(x)
\end{align*}
$$

where $b(\cdot): \mathbb{R} \rightarrow \mathbb{R}$ is uniformly continuous and bounded function.
Let $u_{1}, u_{2} \in L_{2}(\mathbb{R})$. Define:

$$
z_{1}(t, x)=u_{t}(t, x) \text { and } z_{2}(t, x)=u_{x}(t, x) .
$$

Then the equation can be rewritten as follow

$$
\begin{align*}
& \frac{\partial z_{1}(t, x)}{\partial t}=\frac{\partial z_{2}(t, x)}{\partial x}+b(t) z_{2}(t, x)  \tag{5.8}\\
& \frac{\partial z_{2}(t, x)}{\partial t}=\frac{\partial z_{1}(t, x)}{\partial x} \tag{5.9}
\end{align*}
$$

Hence, letting $z(t, x)$ denote the column vector with components $z_{1}(t, x), z_{2}(t, x)$ , we can write the equations (5.8) and (5.9) as follow

$$
\begin{equation*}
\frac{\partial z(t, x)}{\partial t}=A z(t, x)+B(t) z(t, x) \tag{5.10}
\end{equation*}
$$

where

$$
A=\left[\begin{array}{cc}
0 & \frac{\partial}{\partial x} \\
\frac{\partial}{\partial x} & 0
\end{array}\right] \quad B(t)=\left[\begin{array}{ll}
0 & b(t)
\end{array}\right]
$$

Therefore the operator $A$ is now defined on the class of functions $\operatorname{col}\left[z_{1}, z_{2}\right]$ in the product Hilbert spaces $Z=L_{2}(\mathbb{R}) \times L_{2}(\mathbb{R})$ such that their derivatives are also in $Z$ and $t \rightarrow B(t) \in L(Z)$ is a uniformly continuous and bounded function.

It is well known that $A$ generates a strongly continuous contraction semigroup $\{T(t)\}_{t \geq} \subset L(Z)$. Actually for $z \in D(A)$

$$
\frac{d}{d t}[T(t) z, T(t) z]=0
$$

Since the domain of $A$ is dense this means that $\|T(t) z\|=\|z\|$ for every $z \in Z$ or $T(t)$ is an "isometry". We have actually a group; $T(t)$ has bounded inverse for every $t$. Either by working with Fourier transforms (or formally taking $e^{A t}$ ) it can be verified that

$$
T(t) z(x)=\left[\begin{array}{l}
\frac{z_{1}(x+t)+z_{1}(x-t)}{2}+\frac{z_{2}(x+t)-z_{2}(x-t)}{2} \\
\frac{z_{2}(x+t)+z_{2}(x-t)}{2}+\frac{z_{1}(x+t)-z_{1}(x-t)}{2}
\end{array}\right]
$$

The function $B(\cdot)$ belong to the space $W$ of the function $C: \mathbb{R} \rightarrow \mathbb{R}^{2}$ which are continuous and bounded, endow with the topology of the uniform convergence on compact sets of $\mathbb{R}$. For each $C \in W$ and $\tau \in \mathbb{R}$ we define the $\tau$ - translation of $C$ as the function $C_{\tau} \in W$ given by $C_{\tau}(t)=C(\tau+t)$. Then we consider the Hull of $B(\cdot)$ as the following set

$$
\Theta=\operatorname{Hull}(B)=c l\left\{B_{\tau}: \tau \in \mathbb{R}\right\}
$$

where $c l$ denotes the closure in the topology of $W$. The set $\Theta$ is compact in $W$ due to Proposition 1.1 in [1].

We define on $\Theta$ the flow $\sigma(C, t)=C_{t}=\theta \cdot t, \theta=C$. Then $\Theta$ is invariant under $\sigma$.

Instead of concentrating on the single aquation (5.10) we consider the family of equations

$$
\begin{equation*}
\frac{d}{d t} z(t)=A z(t)+B(\theta \cdot t) z(t), \quad \theta \in \Theta \tag{5.11}
\end{equation*}
$$

Here we have abused notation and written $B(\theta \cdot t):=\sigma(\theta \cdot t)$.
Then using Theorema 5.1 we get that the mapping $\pi: Z \times \Theta \times \mathbb{R} \rightarrow \Theta \times Z$ given by

$$
\pi(x, \theta, t)=(\Phi(\theta, t) x, \theta . t)
$$

where $\Phi(\theta, t)$ is given by

$$
\Phi(\theta, t) z_{0}=T(t) z_{0}+\int_{0}^{t} T(t-s) B(\theta \cdot s) \Phi(\theta, s) z_{0} d s
$$

is an skew-product flow on $Z \times \Theta$.
Moreover, if $L=\sup \{\|B(t)\|: t \in \mathbb{R}\}$, then we get that

$$
\Sigma(\hat{\Theta}) \subset[-L, L]
$$

since $\|T(t)\|=1$.

Example 5.3 Using the same idea as in [2] we can consider the equation (5.9) with more general perturbation function $b(t, x)$ which is only bounded and measurable depending on $t, x$ in. In fact. Consider the equation

$$
\begin{equation*}
u_{t t}=u_{x x}+b(t, x) u_{x}, \quad t \in \mathbb{R}, \quad-\infty<x<\infty \tag{5.12}
\end{equation*}
$$

$$
\begin{aligned}
u(0, x) & =u_{1}(x) \\
u_{t}(0, x) & =u_{2}(x)
\end{aligned}
$$

In the same way as equation (5.9) we can write equation (5.12) as follow

$$
\begin{align*}
& \frac{\partial z_{1}(t, x)}{\partial t}=\frac{\partial z_{2}(t, x)}{\partial x}+b(t, x) z_{2}(t, x)  \tag{5.13}\\
& \frac{\partial z_{2}(t, x)}{\partial t}=\frac{\partial z_{1}(t, x)}{\partial x} \tag{5.14}
\end{align*}
$$

Hence, letting $z(t, x)$ denote the column vector with components $z_{1}(t, x), z_{2}(t, x)$ , we can write the equations (5.13) and (5.14) as follow

$$
\begin{equation*}
\frac{\partial z(t, x)}{\partial t}=A z(t, x)+B(t) z(t, x) \tag{5.15}
\end{equation*}
$$

where

$$
A=\left[\begin{array}{cc}
0 & \frac{\partial}{\partial x} \\
\frac{\partial}{\partial x} & 0
\end{array}\right] \quad B(t) z(x)=\left[\begin{array}{cc}
0 & b(t, x) z_{2}(x)
\end{array}\right]
$$

Therefore the operator $A$ is now defined on the class of functions $\operatorname{col}\left[z_{1}, z_{2}\right]$ in the product Hilbert spaces $Z=L_{2}(\mathbb{R}) \times L_{2}(\mathbb{R})$ such that their derivatives are also in $Z$ and the function $B(t, x)=\left[\begin{array}{ll}0 & b(t, x)\end{array}\right]=b(t, x)$ belong to the space $L^{\infty}=l^{\infty}\left(\mathbb{R}^{2}\right)=\left\{b: \mathbb{R}^{2} \rightarrow \mathbb{R}^{2}: b(\cdot, \cdot)\right.$ is measurable and bounded almost everywhere $\}$
with the essential supremum norm, i.e.,

$$
\|b\|_{\infty}=\text { ess.sup }\{\|b(t, x)\|: \quad t, x \in \mathbb{R}\}
$$

For $R>0$ we consider

$$
B_{R}=\left\{b(\cdot, \cdot) \in L^{\infty}:\|b\|_{\infty} \leq R\right\}
$$

So, $B_{R}$ denote the close ball in $L^{\infty}$ of radiu $R$. It is basic fact that the set $B_{R}$ endowed with the weak* topology of $L^{\infty}$ is a compact metrizable topological space (Interesting, the whole space $L^{\infty}$ with the weak* topology is not metrizable). Since $L_{1}^{*}=L_{\infty}$, then a sequence $\left\{f_{n}\right\}$ converges to $f \in L_{\infty}$ if and only if, for every $g \in L_{1}$ one has:

$$
<f_{n}, f>\rightarrow<f, g>\Longleftrightarrow \int f_{n} \cdot g d \mu \rightarrow \int f . g d \mu, \text { as } n \rightarrow \infty
$$

where $<\cdot, \cdot>$ is the duality between $L_{\infty}$ and $L_{1}$.
Moreover, the translation mapping $\sigma(\cdot, \cdot): R_{R} \times \mathbb{R} \rightarrow B_{R}$ given by

$$
\sigma(b, t)=b_{t}=b(t+\cdot, \cdot)
$$

is continuous in this topoly and the mapping $\sigma$ is a flow on $B_{R}$. Then we shall consider the Hull of $b$ as the following set:

$$
\Theta=\operatorname{Hull}(b)=c l\left\{b_{t}: \quad t \in \mathbb{R}\right\} \subseteq B_{R}
$$

where $c l$ denotes the closure in the topology of $B_{R}$. Clearly the set $\Theta$ is compact in $B_{R}$.

We define on $\Theta$ the flow $\sigma(b, t)=b_{t}=\theta \cdot t, \theta=b$. Then $\Theta$ is invariant under $\sigma$. Instead of concentrating on the single aquation (5.15) we consider the family of equations

$$
\begin{equation*}
\frac{d}{d t} z(t)=A z(t)+B(\theta \cdot t) z(t), \quad \theta \in \Theta \tag{5.16}
\end{equation*}
$$

Here we have abused notation and written $B(\theta \cdot t):=\sigma(\theta \cdot t)$.
Then using Theorema 5.1 we get that the mapping $\pi: Z \times \Theta \times \mathbb{R} \rightarrow \Theta \times Z$ given by

$$
\pi(x, \theta, t)=(\Phi(\theta, t) x, \theta . t)
$$

where $\Phi(\theta, t)$ is given by

$$
\Phi(\theta, t) z_{0}=T(t) z_{0}+\int_{0}^{t} T(t-s) B(\theta \cdot s) \Phi(\theta, s) z_{0} d s
$$

is an skew-product flow on $Z \times \Theta$.
Moreover, if $L=\|b(t, x)\|_{\infty}$, then we get that

$$
\Sigma(\hat{\Theta}) \subset[-L, L]
$$

since $\|T(t)\|=1$.

Example 5.4 Linear skew-product flow defined by a Linear Neutral Functional Differential Equation.

For a fixed $r>0$, the space of continuous functions $g:[-r, 0] \rightarrow \mathbb{R}$ with the usual uniform norm is denoted by $C=C([-r, 0], \mathbb{R})$.

Then we shall consider the folling Linear Neutral Functional differential Equation (LNFDE)

$$
\begin{equation*}
\frac{d}{d t}[x(t)-c x(t-r)]=b x(t-r)+a(t) x(t) \tag{5.17}
\end{equation*}
$$

where $c, b \in \mathbb{R}$ and $a: \mathbb{R} \rightarrow \mathbb{R}$ is a uniformly continuous and bounded function.
The equation (5.17) can be written as follow

$$
\begin{equation*}
\frac{d}{d t}[D x(t)]=L x_{t}+f\left(t, x_{t}\right) \tag{5.18}
\end{equation*}
$$

where $x_{t}(\theta)=x(t+\theta)$, for $\quad-r \leq \theta \geq 0, \quad D \phi=\phi(0)-c \phi(-r), \quad L \phi=$ $b \phi(-r)$ and $f\left(t, x_{t}\right)=a(t) x(t)$.

From Theorem 7.2 in [3] we have that for each $\phi \in C$ there is a unique solution $x(t, \phi, a)$ of (5.18) defined on $\mathbb{R}$ which satisfies the initial condition $x_{0}=\phi$.

Moreover, if we consider the equation

$$
\begin{equation*}
\frac{d}{d t}[D x(t)]=L x_{t} \quad(a=0) \tag{5.19}
\end{equation*}
$$

and the corresponding solution map

$$
T(t): C \rightarrow C, \quad T(t) \phi=x_{t}(\cdot, \phi), \quad t \in \mathbb{R}
$$

then from Theorem 7.4 of [3] we get that the solutions of (5.18) are given by the following variational constant formula

$$
\begin{equation*}
x(t, \phi, a)=T(t) \phi(0)+\int_{0}^{t} X(t-s) a(s) x(s, \phi, a) d s \tag{5.20}
\end{equation*}
$$

where $X(t)$ is the solution of the equation

$$
\dot{X}(t)-c \dot{X}(t-r)=b X(t-r)
$$

with the initial condition $X(t)=0, \quad t<0, \quad X(0)=1$. From Theroem 7.6 and corollary 7.2 in [3] we get the following: If $\alpha_{0}=\sup \left\{\operatorname{Re} \lambda: \lambda\left(1-c e^{-\lambda r}\right)=b e^{-\lambda r}\right\}$, then for any $\alpha>\alpha_{0}$ there is $k=k(\alpha)$ such that

$$
\begin{equation*}
\|T(t)\| \leq k e^{\alpha|t|}\|\phi\| \text { and }|X(t)| \leq k e^{\alpha|t|}, \quad \forall t \in \mathbb{R} \tag{5.21}
\end{equation*}
$$

where $\|\phi\|=s u p_{-r \leq s \leq 0}|\phi(s)|$.
The function $a(\cdot)$ belong to the space $W$ of the functions $g: \mathbb{R} \rightarrow \mathbb{R}$ which are continuous and bounded, endow with the topology of the uniform convergence on compact sets of $\mathbb{R}$. For each $g \in W$ and $\tau \in \mathbb{R}$ we define the $\tau$ - translation of $g$ as the function $g_{\tau} \in W$ given by $g_{\tau}(t)=g(\tau+t)$. Then we consider the Hull of $a(\cdot)$ as the following set

$$
\Theta=\operatorname{Hull}(a)=\operatorname{cl}\left\{a_{\tau}: \tau \in \mathbb{R}\right\}
$$

where $c l$ denotes the closure in the topology of $W$. The set $\Theta$ is compact in $W$ and invariant under the flow $\sigma(g, t)=g_{t}=\theta \cdot t, \quad \theta=g$, due to Proposition 1.1 in [1].

Instead of concentrating on the single aquation (5.18) we consider the family of equations

$$
\begin{equation*}
\frac{d}{d t}[D x(t)]=L x_{t}+a(\theta \cdot t) x(t), \quad \theta \in \Theta \tag{5.22}
\end{equation*}
$$

here we have abused notation and written $a(\theta \cdot t):=\sigma(\theta \cdot t)$.
Then we get that the mapping $\pi: C \times \Theta \times \mathbb{R} \rightarrow C \times \Theta$ given by

$$
\pi(\phi, \theta, t)=(\Phi(\theta, t) \phi, \theta \cdot t)
$$

where $\Phi(\theta, t) \phi$ is given by

$$
\Phi(\theta, t) \phi=T(t) \phi+\int_{0}^{t} X_{t-s} a(s) \Phi(\theta, s) \phi d s
$$

is an skew-product flow on $C \times \Theta$.
Moreover, if $L=\sup \{\|a(t)\|: t \in \mathbb{R}\}$, then from (5.21) we get that

$$
\|\Phi(\theta, t)\| \leq k e^{(\alpha+k L)|t|}, \quad \forall t \in \mathbb{R} \text { and } \Sigma(\hat{\Theta}) \subset[-(\alpha+k L), \alpha+k L]
$$

## References

[1] S.N.Chow, and H.Leiva, Dynamical spectrum for time dependent linear systems in banach spaces, October, 1993.
[2] S.N.Chow, K.Lu, and J.Mallet-Paret, Floquet bundles for scalar parabolic equations, to appear, October, 1993.
[3] J.K.Hale, Theory of functional differential equations, Appl.Math.Sci., Springer Verlag, New York, 1977.
[4] J.K.Hale, Asymtotic Behavior of dissipative systems, Math.Surveys and Monographs, Vol.25, Amer.Soc., Providence, R.I, 1988.
[5] J.K.Hale, and G.Raugel, Upper semicontinuity of the attractor for a singularly pertubed hyperbolic equation, J.Differential Equations. 73 (1988), 197214.
[6] A.Haraux, Two remarks on dissipative hyperbolic systems, in College de France Seminaire, Pitman, Boston, 1984.
[7] A.Haraux, Attractors of asymptotically compact processes and applications to nonlinear partial differential equations, Comm. in PDE. 13 (1988), 13831414.
[8] D.Henry, Geometric theory of semilinear parabolic equations, Springer- verlag, New York, 1981.
[9] Y.D.Latushkin, and A.M.Stepin, Linear skew-product flows and semigroups of weighted composition operators, (to appear).
[10] M.L.Peña, Exponential dichotomy singularly perturbed linear functional differential equations with small delays, Applicable Analysis.Vol. 47 (1992), 213225.
[11] L.T.Magalhães, The sprectrum of invariant sets for dissipative semiflows, in dynamics of infinite dimensional systems, NATO ASI series, No.F-37, Spriager Verlag, New York, 1987.
[12] A.Pazy, Semigroups of linear operators and applications to partial differential equations, Applied Mathematical Sciences, Vol.44, Springer Verlag, New York, 1983.
[13] R.T.Rau, Hyperbolic linear skew-product semiflows, Arbeitsbereich Funktional analysis, Mathematishes Institut der Universit'a, Germany, 1993.
[14] D.Ruelle, Characteristic exponents and invariant manifolds in Hilbert space, Annals of Math. 115 (1982), 243-290.
[15] R.J.Sacker and G.R.Sell, Existence of dichotomies and invariant splitting for linear differential systems I, J.Differential Equations. 15 (1974), 429-458.
[16] R.J.Sacker and G.R.Sell, Existence of dichotomies and invariant splitting for linear differential systems II, J.Differential Equations. 22 (1976 A), 478-496.
[17] R.J.Sacker and G.R.Sell, Existence of dichotomies and invariant splittings for linear differential systems III, J.Differential Equations. 22 (1976 B), 497-552.
[18] R.J.Sacker and G.R.Sell, A spectral theory for linear differential systems, J. Differential Equations. 27 (1978), 320-358.
[19] R.J.Sacker and G.R.Sell, Dichotomies for linear evolutionary equations in Banach spaces, J. Dynamics Differential Equations (to appear).
[20] R.J.Sacker and G.R.Sell, The spectrum of an invariant submanifold", J.Differential Equations. 38 (1980), 135-160.
[21] G.R.Sell, The structure of a flow in the vicinity of an almost periodic motion, J.Differential Equations. 27 (1980), 359-393.
[22] G.R.Sell, Smooth linearization near a fixed point, Amer.J.Math. 107 (1985), 1035-1091.
[23] R.Temam, Infinite dimensional dynamical systems in Mechanics and Physics", Springer Verlag, New York, 1988.
[24] Y.Yi, Bifurcation of higher dimensional tory under generic conditions, Ph.D Thesis, 1990.
[25] Y.Yi, Generalized integral manifold theorem, J. Diff. Eqns. (to appear), 1992.
[26] Y.Yi, Stability of integral manifold and orbital attraction of quasi-periodic motion, J. Diff. Eqns.(to appear), 1992.

## CENTER FOR DYNAMICAL SYSTEMS AND NONLINEAR STUDIES

## Report Series

## September 1993 -

## CDSNS92-

102 A normally elliptic Hamiltonian bifurcation, H.W. Broer, S.-N. Chow, Y. Kim and G. Vegter.

103 Large time behavior of an explicit finite difference scheme for an equation arising from compressible flow through porous media, H. Fan.

104 On the perturbation of the kernel for delay systems with continuous kernels, G. Hines.

105 Manuscript withdrawn.
106 Shape index and other indices of Conley type for local maps on locally compact Hausdorff spaces, M. Mrozek.

107 Differentiability with respect to boundary conditions and deviating argument for functional differential systems, J. Ehme, P.W. Eloe and J. Henderson.

108 A discontinuous semilinear elliptic problem without a growth condition, M. Bouguima and A. Boucherif.

109 Equivalent dynamics for a structured population model and a related functional differential equation, H.L. Smith.

110 A degenerate singularity generating geometric Lorenz attractors, F. Dumortier, H. Kokubu and H. Oka.

111 Statistical properties of the periodic Lorentz gas. Multidimensional case, N.I. Chernov.

## CDSNS93-

112 The symmetric Hartman-Grobman Theorem, H.M. Rodrigues.
113 Variation of constants for hybrid systems of FDE, J.K. Hale and W. Huang.
114 Neumann eigenvalue problems on exterior perturbations of the domain, J.M. Arrieta.

115 Bifurcation of equilibria for one-dimensional semilinear equation of the thermoelasticity, L.A.F. de Oliveira and A. Perissinotto, Jr.

116 Isolating neighborhoods and chaos, K. Mischaikow and M. Mrozek.

117 Limits of semigroups depending on parameters, J.K. Hale and G. Raugel.
118 On the sign-variations of solutions of nonlinear two-point boundary value problems, A. Boucherif and B.A. Slimani.

119 Two-point boundary value problems for fourth order ordinary differential equations, A. Boucherif and J. Henderson.

120 The structure of isolated invariant sets and the Conley Index, K. Mischaikow.
121 Competition for a single limiting resource in continuous culture: The variableyield model, H.L. Smith and P. Waltman.

122 Inertial manifolds and the cone condition, J.C. Robinson.
123 The melnikov method and elliptic equations with critical exponent, R.A. Johnson, X. Pan and Y. Yi.

124 Chaos in the Lorenz Equations: a Computer Assisted Proof, K. Mischaikow and M. Mrozek

125 Invariant manifolds and foliations for quasiperiodic systems, S.-N. Chow and K. Lu.

126 Gradient-like structure and Morse decompositions for time-periodic onedimensional parabolic equations, X.-Y. Chen and P. Polácik.

127 Homoclinics and subharmonics of nonlinear two dimensional systems. Boundedness of generalized inverses, H.M. Rodrigues and J.G. Ruas-Filho.

128 Periodic boundary value problems and a priori bounds on solutions, A. Boucherif.
129 Attractors in inhomogeneous conservation laws and parabolic regularizations, H . Fan and J.K. Hale.

130 Periodic boundary value problems with Caratheodory nonlinearities, A. Boucherif.
131 Attractors and convergence of PDE on thin L-shaped domains, J.K. Hale and G. Raugel.

132 Attractor of a semigroup of multi-valued mappings corresponding to an elliptic equation, A.V. Babin.

133 Rotators, periodicity and absence of diffusion in cyclic cellular automata, L.A. Bunimovich and S.E.Troubetzkoy.

134 About completeness for a class of unbounded operators appearing in delay equations, S.M. Verduyn Lunel.

135 A new invariant manifold with an application to a smooth conjugacy at a node, W.M. Rivera.

136 A new invariant manifold withan application to bifurcation of smoothness, W.M. Rivera.

137 Attractors and inertial manifolds for the dynamics of a closed thermosyphon, A. Rodriguez-Bernal.
$138 C^{k}$-smoothness of invariant curves in a global saddle-node bifurcation, T. Young.
139 Partial neutral functional differential equations, J.K. Hale.
140 Numerical dynamics, J.K. Hale.
141 A reaction diffusion equation on a thin L-shaped domain, J.K. Hale and G. Raugel.
142 Criteria of spatial chaos in lattice dynamical systems, V.S. Afraimovich and S.-N. Chow.

143 Statistical properties of 2-D generalized hyperbolic attractors, V.S. Afraimovich, N.L. Chernov and E.A. Sataev.

144 Long-time behavior as $\lambda \rightarrow 0$ of solutions of parabolic equations depending on $\lambda t$, A. Babin and S.-N. Chow.

145 Some remarks on PDEs with nonlinear Neumann boundary conditions, A. Rodríguez-Bernal.

146 Structure of the global attractor of cyclic feedback systems, T. Gedeon and K. Mischaikow.

147 Dynamics of almost periodic scalar parabolic equations, W. Shen and Y. Yi.
148 Existence of a global attractor for the sunflower equation with small delay, M. Lizana.

149 Reaction-diffusion systems on domains with thin channels, S.M. Oliva.
150 Existence and partial characterization of the global attractor for the sunflower equation, M. Lizana.

151 Asymptotic almost periodicity of scalar parabolic equations with almost periodic time dependence, W. Shen and Y. Yi.

152 Zeta functions, periodic trajectories, and the Conley index, C. McCord, K. Mischaikow and M. Mrozek.

153 Center manifold and stability for skew-product flows, S.-N. Chow and Y. Yi.
154 Partial and complete linearizations at stationary points of infinite dimensional dynamical systems with foliations and applications, W.M. Rivera.

155 Dynamical spectrum for skew product flow in Banach spaces, S.N. Chow and H. Leiva.

156 Existence and roughness of the exponential dichotomy for skew-product semiflow in Banach spaces, S-N. Chow and H. Leiva.

## CDSNS94-

157 Rates of eigenvalues on a dumbbell domain. Simple eigenvalue case, J.M. Arrieta.
158 Existence of standing waves for competition-diffusion equations, Y. Kan-On.
159 Fisher-type property of travelling waves for competition-diffusion equations, Y. Kan-On.

160 Uniform ultimate boundedness and synchronization, H.M. Rodrigues.
161 Complete families of pseudotrajectories and shape of attractors, S.Yu. Pilyugin.
162 Aubry-Mather Theorem and quasiperiodic orbits for time dependent reversible systems, S.-N. Chow and M.-L. Pei.

163 Slowly-migrating transition layers for the discrete Allen-Cahn and Cahn-Hilliard equations, C.P. Grant and E.S. Van Vleck.

164 Nontrivial partially hyperbolic sets from a co-dimension one bifurcation, T. Young.
165 Asymptotically autonomous semiflows: Chain recurrence and Lyapunov functions, K. Mischaikow, H. Smith and H.R. Thieme.

166 On the problem of stability in lattice dynamical systems, L.A. Bunimovich and E.A. Carlen.

167 Remark on continuous dependence of attractors on the shape of domain, A.V. Babin and S.Yu. Pilyugin.

168 Special pseudotrajectories for lattice dynamical'systems, V.S. Afraimovich and S.Yu. Pilyugin.

169 Proof and generalization of Kaplan-Yorke conjecture on periodic solution of differential delay equations, J. Li and X-Z. He.

170 On the construction of periodic solutions of Kaplan-Yorke type for some differential delay equations, K. Gopalsamy, J. Li and X-Z. He.

171 Uniformly accurate schemes for hyperbolic systems with relaxation, R.E. Caflisch, S. Jin and G. Russo.

172 Numerical integrations of systems of conservation laws of mixed type, S. Jin.

173 A complex algorithm for computing Lyapunov values, R. Mao and D. Wang.
174 On minimal sets of scalar parabolic equations with skew-product structures, W. Shen and Y. Yi.

175 Lorenz type attractors from codimensional-one bifurcation, V. Afraimovich, S.-N. Chow and W. Liu.

176 Dynamics in a discrete Nagumo equation - spatial chaos, S.-N. Chow and W. Shen.

177 Stability and bifurcation of traveling wave solutions in coupled map lattices, S.-N. Chow and W. Shen.

178 Density of defects and spatial entropy in extended systems, V.S. Afraimovich and L.A. Bunimovich.

179 On the second eigenvalue of the Laplace operator penalized by curvature, E.M. Harrell II.

180 Singular limits for travelling waves for a pair of equations, V. Hutson and K. Mischaikow.

181 Conley Index Theory: Some recent developments, K. Mischaikow.
182 Variational principle for periodic trajectories of hyperbolic billiards, L.A. Bunimovich.

183 Epidemic waves: A diffusion model for fox rabies, W.M. Rivera.
184 High complexity of spatial patterns in gradient reaction-diffusion systems, V. Afraimovich, A. Babin and S.-N. Chow.

185 Spatial chaotic structure of attractors of reaction-diffusion systems, V. Afraimovich, A. Babin and S.-N. Chow.

# TOPICS IN DYNAMICAL SYSTEMS 

## ANNUAL PROGRESS REPORT

Period Covered: 7/1/94-6/30/95

Project Number DMS-9306265
Shui-Nee Chow and Jack K. Hale, PI's
CDSNS
Georgia Tech
Atlanta, GA 30332-0190

April, 1995

## NSF Annual Progress Report

## TOPICS IN DYNAMICAL SYSTEMS

## Summary of Progress:

Two papers were written that received full or partial support from this project. Please find attached copies of these papers and written below are titles and abstracts. Also, J. Pinto finished his Ph.D. thesis on slow motion manifolds. A manuscript is in preparation.
V.S. Afraimovich and S-N. Chow, Synchronizations in lattices of nonlinear Duffing's oscillators

Abstract: In this work we prove the possibility of stochastic synchronization in twodimensional lattice of coupled Duffing's oscillators with external periodic forces. The synchronization occurs provided coupling is dissipative and the coefficients of coupling is greater than some critical values. These values depend on parameters of individual subsystems and on the size of the lattice.

## J.K. Hale and W. Huang, Periodic solutions of singularly perturbed delay equations

Abstract: We consider the singularly perturbed delay differential system

$$
\begin{align*}
\epsilon \dot{x}(t) & =f_{\lambda}(x(t), y(t))  \tag{1.1}\\
y(t) & =g_{\lambda}(x(t), y(t), x(t-1), y(t-1)),
\end{align*}
$$

where $\epsilon>0, \lambda$ are small real parameters, $x \in \mathbb{R}^{m}, y \in \mathbb{R}^{n}$ are vectors and the functions $f_{\lambda}(x, y)=f(x, y, \lambda)$ and $g_{\lambda}(x, y, z, w)=g(x, y, z, w, \lambda)$ are smooth vector valued functions which vanish for all variables equal to zero.

In many situations, for fixed $\lambda>0$, there is an $\epsilon(\lambda)>0$ such that (1.1) undergoes a generic Hopf bifurcation at $(\lambda, \epsilon, x, y)=(\lambda, \epsilon(\lambda), 0,0)$ to a periodic solution $\left(\tilde{x}_{\epsilon \lambda}^{*}, \tilde{y}_{\epsilon \lambda}^{*}\right)$. Let us assume that there is a sectorial region $S$ in a neighborhood $U$ of $(\lambda, \epsilon)=(0,0)$ such that $\left(\tilde{x}_{\epsilon \lambda}^{*}, \tilde{y}_{\epsilon \lambda}^{*}\right)$ exists for all $(\lambda, \epsilon>0)$ in $S$. Our objective is to understand the behavior of the profile of $\left(\tilde{x}_{\epsilon \lambda}^{*}, \tilde{y}_{\epsilon \lambda}^{*}\right)$ as $\epsilon \rightarrow 0$. We are not able to do this in the general context described, but we can say something if we impose more conditions on the functions $f_{\lambda}, g_{\lambda}$.

It is to be expected that the limiting profile is in some way related to the equation obtained by putting $\epsilon=0$ in (1.1). For $\epsilon=0$, we suppose that the resulting equation (1.1) defines a map on $\mathbb{R}^{m} \times \mathbb{R}^{n}$ :

$$
\begin{equation*}
(x, y) \in \mathbb{R}^{m} \times \mathbb{R}^{n} \mapsto T_{\lambda}(x, y) \in \mathbb{R}^{m} \times \mathbb{R}^{n} \tag{1.2}
\end{equation*}
$$

for which the origin is stable for $\lambda<0$ and unstable for $\lambda>0$. Let us also suppose that $T_{\lambda}(x, y)$ undergoes a generic period doubling bifurcation at $(x, y, \lambda)=(0,0,0)$ with the period two points being $d_{j \lambda} \in \mathbb{R}^{m} \times \mathbb{R}^{n}, j=1,2$. If the bifurcation is supercritical, the period two orbit is stable and, if the bifurcation is subcritical, the period two orbit is unstable. In the supercritical case, there is a natural stable periodic function of period two given by the square wave $\left(x_{\lambda}^{s}(t), y_{\lambda}^{s}(t)\right)$ which alternately takes the values $d_{1 \lambda}$ and $d_{2 \lambda}$ on intervals of length one. In the subcritical case, the period two points are unstable and the natural periodic function of period two is a pulse wave $\left(x_{\lambda}^{p}(t), y_{\lambda}^{p}(t)\right)$ which is zero except on the integers and it alternately takes the values $d_{1 \lambda}$ and $d_{2 \lambda}$ on the integers.

Under some conditions on the functions $f_{\lambda}, g_{\lambda}$, the generic period doubling bifurcation of $T_{\lambda}(x, y)$ leads to a generic Hopf bifurcation in (1.1) which is supercritical (resp. subcritical) if the period doubling bifurcation is supercritical (resp. subcritical). A natural question is the following: Is it possible that the limiting profile of the periodic solution $\left(\tilde{x}_{\epsilon \lambda}^{*}, \tilde{y}_{\epsilon \lambda}^{*}\right)$ obtained through the Hopf bifurcation is either the square wave or pulse wave? We present situations in which this is true.

# Periodic solutions of singularly perturbed delay equations 

## by

J.K. Hale and W. Huang

CDSNS94-201

# Periodic solutions of singularly perturbed delay equations 

by

Jack K. Hale * and Wenzhang Huang

1. Introduction. In this paper, we consider the singularly perturbed delay differential system

$$
\begin{align*}
\epsilon \dot{x}(t) & =f_{\lambda}(x(t), y(t))  \tag{1.1}\\
y(t) & =g_{\lambda}(x(t), y(t), x(t-1), y(t-1))
\end{align*}
$$

where $\epsilon>0, \lambda$ are small real parameters, $x \in \mathbb{R}^{m}, y \in \mathbb{R}^{n}$ are vectors and the functions $f_{\lambda}(x, y)=f(x, y, \lambda)$ and $g_{\lambda}(x, y, z, w)=g(x, y, z, w, \lambda)$ are smooth vector valued functions which vanish for all variables equal to zero.

In many situations, for fixed $\lambda>0$, there is an $\epsilon(\lambda)>0$ such that (1.1) undergoes a generic Hopf bifurcation at $(\lambda, \epsilon, x, y)=(\lambda, \epsilon(\lambda), 0,0)$ to a periodic solution $\left(\tilde{x}_{\epsilon \lambda}^{*}, \tilde{y}_{\epsilon \lambda}^{*}\right)$. Let us assume that there is a sectorial region $S$ in a neighborhood $U$ of $(\lambda, \epsilon)=(0,0)$ such that $\left(\tilde{x}_{\epsilon \lambda}^{*}, \tilde{y}_{\epsilon \lambda}^{*}\right)$ exists for all $(\lambda, \epsilon>0)$ in $S$. Our objective is to understand the behavior of the profile of $\left(\tilde{x}_{\epsilon \lambda}^{*}, \tilde{y}_{\epsilon \lambda}^{*}\right)$ as $\epsilon \rightarrow 0$. We are not able to do this in the general context described, but we can say something if we impose more conditions on the functions $f_{\lambda}, g_{\lambda}$.

It is to be expected that the limiting profile is in some way related to the equation obtained by putting $\epsilon=0$ in (1.1). For $\epsilon=0$, we suppose that the resulting equation (1.1) defines a map on $\mathbb{R}^{m} \times \mathbb{R}^{n}$ :

$$
\begin{equation*}
(x, y) \in \mathbb{R}^{m} \times \mathbb{R}^{n} \mapsto T_{\lambda}(x, y) \in \mathbb{R}^{m} \times \mathbb{R}^{n} \tag{1.2}
\end{equation*}
$$

for which the origin is stable for $\lambda<0$ and unstable for $\lambda>0$. Let us also suppose that $T_{\lambda}(x, y)$ undergoes a generic period doubling bifurcation at $(x, y, \lambda)=(0,0,0)$ with the period two points being $d_{j \lambda} \in \mathbb{R}^{m} \times \mathbb{R}^{n}, j=1,2$. If the bifurcation is supercritical, the period two orbit is stable and, if the bifurcation is subcritical, the period two orbit is unstable. In the supercritical case, there is a natural stable periodic function of period two given by the square wave $\left(x_{\lambda}^{s}(t), y_{\lambda}^{s}(t)\right)$ which alternately takes the values $d_{1 \lambda}$ and $d_{2 \lambda}$ on intervals of length one. In the subcritical case, the period two points are unstable and the natural periodic function of period two is a pulse wave $\left(x_{\lambda}^{p}(t), y_{\lambda}^{p}(t)\right)$ which is zero except on the integers and it alternately takes the values $d_{1 \lambda}$ and $d_{2 \lambda}$ on the integers.

Under some conditions on the functions $f_{\lambda}, g_{\lambda}$, the generic period doubling bifurcation of $T_{\lambda}(x, y)$ leads to a generic Hopf bifurcation in (1.1) which is supercritical (resp. subcritical) if the period doubling bifurcation is supercritical (resp. subcritical). A natural question is the following: Is it possible that the limiting profile of the periodic solution $\left(\tilde{x}_{\epsilon \lambda}^{*}, \tilde{y}_{\epsilon \lambda}^{*}\right)$ obtained through the Hopf bifurcation is either the square wave or pulse wave? We will give some situations in which this is true.

Let us mention that equations of the form (1.1) occur frequently in the applications. The delay differential equation

$$
\begin{equation*}
\epsilon \dot{x}(t)=-x(t)+f_{\lambda}(x(t-1)) \tag{1.3}
\end{equation*}
$$

[^2]has often served as a model for physiological control systems [6], [14], [20] and for the transmission of light through a ring cavity [1], [5], [12], [13]. This can be considered as a special case of (1.1) by letting $y(t)=x(t-1)$. Under the conditions that the map $x \mapsto f_{\lambda}(x)$ undergoes a generic supercritical periodic doubling bifurcation at $(x, \lambda)=(0,0)$, it was shown in [2] that there is a periodic solution of (1.3) which approaches the square wave mentioned above. If the bifurcation is subcritical, it was shown in [7] that there is a periodic solution of (1.3) which approaches a pulse wave but the pulses are larger that the values $d_{1 \lambda}, d_{2 \lambda}$ mentioned above. The explanation for this will be given later. The limiting behavior of large amplitude periodic solutions of (1.3) also has been investigated in [15], [16] when the limiting profile is a square wave.

The more general equation

$$
\begin{equation*}
\left(\epsilon_{m} \frac{d}{d t}+1\right) \cdots\left(\epsilon_{1} \frac{d}{d t}+1\right) y(t)=h_{\lambda}(y(t-1)) \tag{1.4}
\end{equation*}
$$

where each $\epsilon_{j}>0$ is a small parameter, also has been proposed as a model for transmission of light through a ring cavity (see [18], [19]). Equation (1.4) is equivalent to the system

$$
\begin{align*}
& \epsilon_{1} \dot{x}_{1}(t)+x_{1}(t)=x_{2}(t) \\
& \cdots \cdots  \tag{1.5}\\
& \epsilon_{m-1} \dot{x}_{m-1}(t)+x_{m-1}(t)=x_{m}(t) \\
& \epsilon_{m} \dot{x}_{m}(t)+x_{m}(t)=h_{\lambda}\left(x_{1}(t-1)\right)
\end{align*}
$$

If we scale the $\epsilon_{j}$ as $\epsilon_{j}=\epsilon \alpha_{j}^{-1}, j=1,2, \ldots, m$, then we obtain an equivalent matrix equation

$$
\begin{equation*}
\epsilon \dot{x}(t)+A x(t)=A f_{\lambda}(x(t-1)), \tag{1.6}
\end{equation*}
$$

where $A$ and $f$ are given by

$$
A=\left[\begin{array}{cccccc}
\alpha_{1} & -\alpha_{1} & 0 & \ldots & 0 & 0 \\
0 & \alpha_{2} & -\alpha_{2} & \ldots & 0 & 0 \\
\cdot & \cdot & \cdot & \ldots & \cdot & \cdot \\
\cdot & \cdot & \cdot & \ldots & \cdot & \cdot \\
\cdot & \cdot & \cdot & \ldots & \cdot & \cdot \\
0 & 0 & 0 & \ldots & \alpha_{m-1} & -\alpha_{m-1} \\
0 & 0 & 0 & \ldots & 0 & \alpha_{m}
\end{array}\right], \quad f_{\lambda}(x)=\left[\begin{array}{c}
h_{\lambda}\left(x_{1}\right) \\
h_{\lambda}\left(x_{1}\right) \\
\cdot \\
\cdot \\
\cdot \\
h_{\lambda}\left(x_{1}\right) \\
h_{\lambda}\left(x_{1}\right)
\end{array}\right]
$$

This also is a special case of (1.1) with $y(t)=x_{1}(t-1)$.
Under the assumption that the map $x \mapsto h_{\lambda}(x)$ undergoes a generic period doubling, it was shown in [8] that there are periodic solutions of (1.5) resembling either a square wave or a pulse wave depending upon whether the bifurcation was supercritical or subcritical. For the general matrix equation (1.6) without assuming that it arises from (1.5) and under generic conditions on the function $f_{\lambda}$ in (1.6), similar results were given in [8].

The proofs of the above results exploited the properties of flows on center manifolds for some scaled equations which will be mentioned below.

The equation

$$
\begin{align*}
\epsilon \dot{x}(t) & +A x(t)=A f_{\lambda}(y(t))  \tag{1.7}\\
y(t) & =g_{\lambda}(x(t-1), y(t-1))
\end{align*}
$$

where $x \in \mathbb{R}^{m}, y \in \mathbb{R}^{n}$ are vectors, the $m \times m$ matrix $A$ has an inverse and the functions $f_{\lambda}(y)$ and $g_{\lambda}(x, y)$ are smooth vector valued functions, has been used in [12], [13] as a model of a ring cavity containing a nonlinear dielectric medium for which part of the transmitted light is fed back into the medium. For $\epsilon=0$, the map $T_{\lambda}(x, y)=\left(f_{\lambda}(y), g_{\lambda}\left(f_{\lambda}(y), y\right)\right.$. Under the assumption that there is a supercritical period doubling bifurcation of the map $y \rightarrow g_{\lambda}\left(f_{\lambda}(y), y\right)$ at $(y, \lambda)=(0,0)$ and some other generic conditions, it was shown in [3] that there is a periodic solution of (1.7) which is similar to the square wave for $\epsilon$ small. The proof of this result used functional analytic methods based on exponential dichotomies.

Our objective in this paper is to show that center manifold techniques can be applied to discuss (1.7) when the period doubling bifurcation is either supercritical or subcritical. The limiting profile of the corresponding periodic solutions is either a square wave or a pulse wave. The precise statements of the results on convergence to square waves or pulse waves for the case (1.7) are given in Section 2. Sections 3 and 4 are devoted to the proof. The hypotheses imposed imply the existence of periodic solutions which are related to the period two points of the map. However, that analysis does not address the important question of the relataionship between the existence of generic period doubling of the map and the existence and stability of a generic Hopf bifurcation. These problems are addressed in Sections 5 and 6.

We remark that the equation (1.1) arises also in the theory of transmission lines. If the lines are lossless and described by the telegraph equations with the boundary conditions reflecting Kirchoff's laws, it has been known for a long time that the flow can be described by an equivalent set of neutral delay differential equations (see, for example, [11] for a discussion and references). Many of these same problems also can be written in the form (1.1). For example, in [17], the equations for a transmission line with a tunnel diode and a lumped parallel capacitance can be written as

$$
\begin{align*}
\epsilon \dot{x}(t) & =y(t)-g(x(t)) \\
y(t) & =\alpha+K y(t-1)-x(t)-L x(t-1) \tag{1.8}
\end{align*}
$$

where $(x(t), y(t))$ represent the voltage and current at one end of the line, all constants are positive and represent physical parameters. Under reasonable assumptions in the model, the parameter $\epsilon$ can be considered to be very small. In [17], several wave forms were observed numerically which compared reasonably well with experimental results. Some of these wave forms were very similar to square waves.

It is natural to enquire if it is theoretically possible to prove that, for some values of the parameters, there are periodic solutions of (1.8) which are similar to a square wave for $\epsilon>0$ small. If we try to follow the same procedure as before, we first investigate the possibility of a supercritical period doubling bifurcation of some map. For $\epsilon=0$, we must have $y(t)=g(x(t))$ and thus

$$
\begin{equation*}
x(t)+g(x(t))=\alpha+K g(x(t-1))-L x(t-1) \tag{1.9}
\end{equation*}
$$

For a given function $g$ similar to the cubic polynomial, it is possible to show that there are values of the parameters in the equation and a constant $x_{0}$ such that (1.9) defines a map which undergoes a supercritical period doubling bifurcation near $x_{0}$. With the methods that we develop below, it should be possible to prove that there is a solution of (1.8) which approaches as $\epsilon \rightarrow 0$ either a square wave or a pulse wave. However, we do not discuss this problem here.
2. A ring cavity model. In this section, we consider the ring cavity model (1.7). We will impose some generic conditions on the functions $f, g$ which will ensure that the associated map for $\epsilon=0$ undergoes a generic period doubling at $(y, \lambda)=(0,0)$. We search for periodic solutions of (1.7) of period approximately 2 and impose additional generic conditions on $f, g$ which will permit the reduction of the problem to a discussion of certain periodic solutions of a two-dimensional vector field on a center manifold of specific scaled equations. We discuss also how the conditions on $f, g$ are related to the generic Hopf bifurcation of periodic solutions of (1.7) for fixed $\epsilon>0$ and $\lambda$ being the bifurcation parameter. Since this section is devoted entirely to (1.7), we repeat some of the previous formulas and hypotheses since some are stated in other terms.

We consider the equation

$$
\begin{align*}
& \epsilon \dot{x}(t)+A x(t)=A f_{\lambda}(y(t)) \\
& y(t)=g_{\lambda}(x(t-1), y(t-1)) \tag{2.1}
\end{align*}
$$

where $\epsilon>0, \lambda$ are small parameters, $x \in \mathbb{R}^{m}, y \in \mathbb{R}^{n}$ are vectors, $A$ is an $m \times m$ constant matrix and the functions $f_{\lambda}(y)$ and $g_{\lambda}(x, y)$ are smooth vector valued functions. We suppose that

$$
\begin{equation*}
A^{-1} \text { exists } \quad f_{\lambda}(0)=0, \quad g_{\lambda}(0,0)=0 \tag{H1}
\end{equation*}
$$

For $\epsilon=0$, there is the map

$$
\begin{equation*}
y \in \mathbb{R}^{n} \mapsto \mathcal{F}_{\lambda}(y) \equiv g_{\lambda}\left(f_{\lambda}(y), y\right) \in \mathbb{R}^{n} \tag{2.2}
\end{equation*}
$$

We first want to impose some conditions on the functions $f_{\lambda}, g_{\lambda}$ to ensure that $\mathcal{F}_{\lambda}$ undergoes a generic period doubling at $(y, \lambda)=(0,0)$ and that the period two point is either stable or of index 1 ; that is, has an unstable manifold of dimension 1. For notation, we let

$$
\begin{align*}
A_{2}(\lambda) & =D_{y} f_{\lambda}(0), \quad B_{1}(\lambda)=D_{x} g_{\lambda}(0,0) ; \quad B_{2}(\lambda)=D_{y} g_{\lambda}(0,0)  \tag{2.3}\\
C(\lambda) & \equiv B_{1}(\lambda) A_{2}(\lambda)+B_{2}(\lambda)=D_{y} \mathcal{F}_{\lambda}(0)
\end{align*}
$$

We also let $\sigma(C)$ denote the spectrum of a square matrix $C$ and let $B_{\rho}=\{z \in \mathbb{C}:|z| \leq \rho\}$. We assume that

$$
\begin{gather*}
-(1+\lambda) \in \sigma(C(\lambda)) \text { is a simple eigenvalue }  \tag{H2}\\
\sigma(C(\lambda)) \backslash\{-(1+\lambda)\} \subset B_{\rho} \text { with } \rho<1
\end{gather*}
$$

We remark that the theory below could be developed for the case in which $\sigma(C(\lambda) \backslash\{-(1+$ $\lambda)\}$ has no eigenvalues of modulus 1 , but in this case, if there is an eigenvalue with modulus $>1$, then the period two bifurcation always will be unstable, a situation that is not the desirable one from the point of view of applications.

If (H2) is satisfied, then we can introduce a change of variables in $y$ to obtain

$$
C(\lambda)=\left[\begin{array}{cc}
-(1+\lambda) & 0  \tag{2.4}\\
0 & H_{0}(\lambda)
\end{array}\right], \quad \sigma\left(H_{0}(\lambda) \subset B_{\rho}, \quad \rho<1\right.
$$

If we let $y=\operatorname{col}\left(y_{1}, y_{2}\right) \in \mathbb{R}^{1} \times \mathbb{R}^{n-1}, \mathcal{F}_{\lambda}=\operatorname{col}\left(\mathcal{F}_{1 \lambda}, \mathcal{F}_{2 \lambda}\right) \in \mathbb{R}^{1} \times \mathbb{R}^{n-1}$, then we can use (2.4) to write

$$
\begin{align*}
& \mathcal{F}_{1 \lambda}(y)=-(1+\lambda) y_{1}+k_{1}(\lambda) y_{1}^{2}+y_{1} k_{2}(\lambda) y_{2}+k_{3}(\lambda) y_{1}^{3}+O\left(\left\|y_{2}\right\|^{2}+y_{1}^{2}\left\|y_{2}\right\|+y_{1}^{4}\right)  \tag{2.5}\\
& \mathcal{F}_{2 \lambda}(y)=H_{0}(\lambda) y_{2}+y_{1}^{2} H_{1}(\lambda)+y_{1} H_{2}(\lambda)+y_{2} O\left(\left\|y_{2}\right\|^{2}+y_{1}^{2}\left\|y_{2}\right\|+\|y\|^{3}\right)
\end{align*}
$$

as $(y, \lambda) \rightarrow(0,0)$.
We assume that

$$
\begin{equation*}
R_{1} \equiv k_{2}(0)\left[I_{n-1}-H_{0}(0)\right]^{-1} H_{1}(0)+k_{1}^{2}(0)+k_{3}(0) \neq 0 \tag{H3}
\end{equation*}
$$

where, for any integer $k, I_{k}$ denotes the identity matrix in $\mathbb{R}^{k}$. In the following, we often omit this index if there is no reason for confusion.

Lemma 2.1. If ( H 1$)-(\mathrm{H} 3)$ are satisfied, then there is a generic period two doubling bifurcation of $\mathcal{F}_{\lambda}$ at $(y, \lambda)=(0,0)$. More precisely, if $R_{1} \lambda>0$, then there are period two points $d_{1 \lambda}, d_{2 \lambda}$ of $\mathcal{F}_{\lambda}$ such that $\mathcal{F}_{\lambda}\left(d_{1 \lambda}\right)=d_{2 \lambda}, \mathcal{F}_{\lambda}\left(d_{2 \lambda}\right)=d_{1 \lambda}$, which is stable for $R_{1}>0$ (supercritical bifurcation) and unstable for $R_{1}<0$ (subcritical bifurcation).

The proof is a standard application of the method of Lyapunov-Schmidt and is omitted.

To relate the period doubling of the map $\mathcal{F}$ to (2.1), we seek periodic solutons of (2.1) with a period $2+2\left(r_{0}+h\right) \epsilon$, where $r_{0}$ is a fixed parameter (to be determined later) which depends only upon the matrices $A, A_{2}(0), B_{1}(0), B_{2}(0)$ and $h$ will be determined as a function of $\epsilon$. If $(x(t), y(t))$ is such a solution of (2.1), we introduce the transformation (originally used in [4] for a scalar equation)

$$
\begin{array}{cc}
u_{1}(t)=x\left(-\epsilon\left(r_{0}+h\right) t\right), & u_{2}(t)=x\left(-\epsilon\left(r_{0}+h\right) t+1+\epsilon\left(r_{0}+h\right)\right)  \tag{2.6}\\
v_{1}(t)=y\left(-\epsilon\left(r_{0}+h\right) t\right), & v_{2}(t)=y\left(-\epsilon\left(r_{0}+h\right) t+1+\epsilon\left(r_{0}+h\right)\right)
\end{array} .
$$

Since $x(t)$ and $y(t)$ have period $2+2\left(r_{0}+h\right) \epsilon$, we see that

$$
\begin{align*}
& u_{2}(t-1)=x\left(-\epsilon\left(r_{0}+h\right) t-1\right)  \tag{2.7}\\
& v_{2}(t-1)=y\left(-\epsilon\left(r_{0}+h\right) t-1\right)
\end{align*}
$$

If we use (2.6) and (2.7) in (2.1), we deduce that

$$
\begin{align*}
& \dot{u}_{1}(t)=\left(r_{0}+h\right) A u_{1}(t)-\left(r_{0}+h\right) A f_{\lambda}\left(v_{1}(t)\right) \\
& \dot{u}_{2}(t)=\left(r_{0}+h\right) A u_{2}(t)-\left(r_{0}+h\right) A f_{\lambda}\left(v_{2}(t)\right) \\
& v_{1}(t)=g_{\lambda}\left(u_{2}(t-1), v_{2}(t-1)\right)  \tag{2.8}\\
& v_{2}(t)=g_{\lambda}\left(u_{1}(t-1), v_{1}(t-1)\right) .
\end{align*}
$$

This equation is now independent of $\epsilon$ and we can consider $h, \lambda$ as the bifurcation parameters, assuming, of course, that we know $r_{0}$.

It is convenient to write these equations in a more compact form by letting $u=\mathrm{col}$ $\left(u_{1}, u_{2}\right), v=\operatorname{col}\left(v_{1}, v_{2}\right)$,

$$
\begin{align*}
\hat{A} & =\left[\begin{array}{cc}
A & 0 \\
0 & A
\end{array}\right], \quad \hat{A}_{2}=\left[\begin{array}{cc}
A_{2}(0) & 0 \\
0 & A_{2}(0)
\end{array}\right],  \tag{2.9}\\
\hat{B}_{j} & =\left[\begin{array}{cc}
0 & B_{j}(0) \\
B_{j}(0) & 0
\end{array}\right], j=1,2 . \\
F_{h \lambda}(u, v) & =-\left[\begin{array}{l}
\left(r_{0}+h\right) f_{\lambda}\left(v_{1}\right)-r_{0} D_{y} f_{\lambda}(0) v_{1}-h A u_{1} \\
\left(r_{0}+h\right) f_{\lambda}\left(v_{2}\right)-r_{0} D_{y} f_{\lambda}(0) v_{2}-h A u_{2}
\end{array}\right] \\
G_{h \lambda}(u, v) & =\left[\begin{array}{l}
g_{\lambda}\left(u_{2}, v_{2}\right)-D_{x} g_{\lambda}(0,0) u_{2}-D_{y} g_{\lambda}(0,0) v_{2} \\
g_{\lambda}\left(u_{1}, v_{1}\right)-D_{x} g_{\lambda}(0,0) u_{1}-D_{y} g_{\lambda}(0,0) v_{1}
\end{array}\right] . \tag{2.10}
\end{align*}
$$

If we use (2.9) and (2.10), then (2.8) can be written as

$$
\begin{align*}
\dot{u}(t) & =r_{0} \hat{A} u(t)-r_{0} \hat{A} \hat{A}_{2} v(t)+F_{h \lambda}(u(t), v(t)) \\
v(t) & =\hat{B}_{1} u(t-1)+\hat{B}_{2} v(t-1)+G_{h \lambda}(u(t-1), v(t-1)) \tag{2.11}
\end{align*}
$$

The linear variational equation of (2.11) for $(h, \lambda)=(0,0)$ is given by

$$
\begin{align*}
& \dot{u}(t)=r_{0} \hat{A} u(t)-r_{0} \hat{A} \hat{A}_{2} v(t)  \tag{2.12}\\
& v(t)=\hat{B}_{1} u(t-1)+\hat{B}_{2} v(t-1)
\end{align*}
$$

The eigenvalues of (2.12) are the solutions of the characteristic equation

$$
\operatorname{det} \Delta\left(\mu, r_{0}\right)=0, \quad \Delta\left(\mu, r_{0}\right)=\left[\begin{array}{cc}
\mu-r_{0} \hat{A} & r_{0} \hat{A} \hat{A}_{2}  \tag{2.13}\\
-\hat{B}_{1} e^{-\mu} & I-\hat{B}_{2} e^{-\mu}
\end{array}\right]
$$

Because of (2.3), (2.4) and the symmetry in (2.12), zero always is an eigenvalue. We impose conditions on the coefficients to ensure that, for a suitable choice of $r_{0}, 0$ is an eigenvalue of multiplicity two and there are no other eigenvalues on the imaginary axis. To do this, we need some additional notation. Let

$$
-B_{1}(0) A^{-1} A_{2}(0) \equiv S=\left[\begin{array}{ll}
S_{11} & S_{12}  \tag{2.14}\\
S_{21} & S_{22}
\end{array}\right], \quad-B_{1}(0) A^{-2} A_{2}(0) \equiv W=\left[\begin{array}{ll}
W_{11} & W_{12} \\
W_{21} & W_{22}
\end{array}\right]
$$

where $S_{i j}, W_{i j}$ are matrices with $S_{11}$ and $W_{11}$ scalars and the other matrices of obvious dimension. We suppose that

$$
\begin{equation*}
R_{0} \equiv-S_{11}^{2}+2\left[W_{11}+S_{12}\left(I+H_{0}(0)\right)^{-1} S_{21}\right] \neq 0, \quad S_{11} \neq 0 \tag{H4}
\end{equation*}
$$

and that

$$
\operatorname{det}\left[\begin{array}{cc}
i \omega I_{2 m}-S_{11} \hat{A} & S_{11} \hat{A} \hat{A}_{2}  \tag{H5}\\
-\hat{B}_{1} e^{-i \omega} & I_{2 n}-\hat{B}_{2} e^{-i \omega}
\end{array}\right] \neq 0 \quad \text { for } \omega \in \mathbb{R} \backslash\{0\}
$$

We make the final hypothesis that

$$
\begin{equation*}
\sigma\left(B_{2}(0)\right) \subset B_{\rho}, \quad \rho<1 \tag{H6}
\end{equation*}
$$

Lemma 2.2. If (H4)-(H6) are satisfied and $r_{0}=S_{11}$, then $\mu=0$ is an eigenvalue of (2.12) of multiplicity two and there is a $\delta>0$ such that the remaining eigenvalues satisfy $|\operatorname{Re} \mu| \geq \delta>0$ and there are only a finite number of eigenvalues with positive real parts.

Theorem 2.3. Suppose that $\mathcal{F}_{\lambda}(x)$ satisfies (2.5) and (H1)-(H6). Then there are a neighborhood $V$ of zero in $\mathbb{R}^{m} \times \mathbb{R}^{n}$ and a neighborhood $U$ of $(0,0)$ in the $(\lambda, \epsilon)$ plane and a sectorial region $S$ in $U$ such that, if $(\lambda, \epsilon) \in U$, then there is a periodic solution of (2.1) in $V$ given by a periodic function $\left(\tilde{x}_{\lambda, \epsilon}, \tilde{y}_{\lambda, \epsilon}\right)$ with period $2 \tau(\lambda, \epsilon)=2+2 S_{11} \epsilon+O(|\epsilon|(|\lambda|+|\epsilon|))$ as $(\lambda, \epsilon) \rightarrow(0,0)$ if and only if $(\lambda, \epsilon) \in S$. Furthermore, this solution is unique.

Of course, the sector $S$ must belong to the set $\epsilon>0$ in the $(\lambda, \epsilon)$ plane. We actually will show that, if $R_{1}>0$ (the supercritical case of period doubling of the map) and $R_{0}>0$, then the sector $S \subset\{(\lambda, \epsilon): \epsilon>0, \lambda>0\}$ and, for $\lambda=\lambda_{0}>0$, fixed, the set $\left\{\epsilon:\left(\epsilon, \lambda_{0}\right) \in S\right\}$ is an interval $\lambda_{0} \times\left(0, \epsilon_{0}\left(\lambda_{0}\right)\right)$, where

$$
\epsilon_{0}\left(\lambda_{0}\right)=\frac{1}{\pi}\left(\frac{2 \pi \lambda_{0}}{R_{0}}\right)^{\frac{1}{2}}+O\left(\lambda_{0}\right)
$$

as $\lambda_{0} \rightarrow 0$. For any $\epsilon \in\left(0, \epsilon_{0}\left(\lambda_{0}\right)\right.$ ), the periodic solution $\left(\tilde{x}_{\lambda_{0}, \epsilon}(t), \tilde{y}_{\lambda_{0}, \epsilon}\right)$ approaches a square wave as $\epsilon \rightarrow 0$; that is, the periodic solution $\left(\tilde{x}_{\lambda_{0}, \epsilon}(t), \tilde{y}_{\lambda_{0}, \epsilon}\right)$ has the property that $\left(\tilde{x}_{\lambda_{0}, \epsilon}(t), \tilde{y}_{\lambda_{0}, \epsilon}\right) \rightarrow d_{1 \lambda}$ (respectively, $\left.d_{2 \lambda}\right)$ as $\epsilon \rightarrow 0$ uniformly on compact sets of $(0,1)$ (respectively, (1, 2)).

If $R_{1}<0$ (the subcritical case of period doubling of the map) and $R_{0}>0$, the sector $S$ contains points $(\epsilon, \lambda)$ with $\lambda$ both negative and positive and the periodic orbits have a different structure as $\epsilon \rightarrow 0$. More precisely, for $\lambda=\lambda_{0}>0$, fixed, the set $\left\{\epsilon:\left(\epsilon, \lambda_{0}\right) \in S\right\}$ is an interval $\lambda_{0} \times\left(\epsilon_{0}\left(\lambda_{0}\right), \beta_{0}\left(\lambda_{0}\right)\right)$. For $\lambda=\lambda_{0}<0$, fixed, the set $\left\{\epsilon:\left(\epsilon, \lambda_{0}\right) \in S\right\}$ is an interval $\lambda_{0} \times\left(0, \alpha_{0}\left(\lambda_{0}\right)\right)$. For any $\epsilon \in\left(0, \alpha_{0}\left(\lambda_{0}\right)\right)$, the unique periodic solution $\left(\tilde{x}_{\lambda_{0}, \epsilon}(t), \tilde{y}_{\lambda_{0}, \epsilon}\right)$ becomes pulse like as $\epsilon \rightarrow 0$ in the following sense: the periodic solution $\left(\tilde{x}_{\lambda_{0}, \epsilon}(t), \tilde{y}_{\lambda_{0}, \epsilon}\right)$ has the property that $\left(\tilde{x}_{\lambda_{0}, \epsilon}(t), \tilde{y}_{\lambda_{0}, \epsilon}\right) \rightarrow 0$ as $\epsilon \rightarrow 0$ uniformly on compact sets of $(0,1) \cup(1,2)$. In the pulse like solution, the pulses in the solution occur near the integers and are opposite in direction. However, the magnitude of the pulse near the integers exceeds the magnitude of the corresponding period two point of the map.

The part of the boundary of the sector $S$ described by the curve $\Gamma=\{(\lambda, \epsilon(\lambda)\}$ corresponds to a Hopf bifurcation curve. In the supercritical case, the Hopf bifurcation at $\left(\lambda_{0}, \epsilon\left(\lambda_{0}\right)\right)$ is in the direction of $\epsilon<\epsilon\left(\lambda_{0}\right)$ and the periodic orbit has a unique continuation to the interval $\left\{\left(\lambda_{0}, \epsilon\right): 0<\epsilon<\epsilon\left(\lambda_{0}\right)\right.$. In the subcritical case, the Hopf bifurcation is in the direction of $\epsilon>\epsilon\left(\lambda_{0}\right)$ and the periodic orbit has a unique continuation in all of $S$ (see Figure 2.1).


Figure 2.1. Bifurcation Diagram in the $(\lambda, \epsilon)$-plane
Lemma 2.2 will be proved in Section 3 and Theorem 2.3 in Section 4. As remarked earlier, Sections 5 and 6 will be devoted to relating the hypotheses in Theorem 2.3 to those that are required to obtain a generic Hopf bifurcation on the curve $\Gamma$.
3. Proof of Lemma 2.2. For any integer $N$, Let $C_{N}=C\left([-1,0] ; \mathbb{R}^{N}\right)$. Let $\hat{\Omega}_{0}=$ $\left\{\left(\varphi_{1}, \varphi_{2}\right) \in C_{2 m} \times C_{2 n}: \varphi_{2}(0)=\hat{B}_{1} \varphi_{1}(-1)+\hat{B}_{2} \varphi_{2}(-1)\right\}$. Equation (2.12) generates a $C^{0}$-semigroup $\hat{S}(t)$ on $\hat{\Omega}_{0}$. If $\hat{\mathcal{A}}$ is the infinitesimal generator of $\hat{S}(t)$, then

$$
\begin{aligned}
\mathcal{D}(\hat{\mathcal{A}}) & =\left\{\left(\varphi_{1}, \varphi_{2}\right) \in \hat{\Omega}_{0}:\left(\dot{\varphi}_{1}, \dot{\varphi}_{2}\right) \in \hat{\Omega}_{0}, \dot{\varphi}_{1}(0)=S_{11} \hat{A} \varphi_{1}(0)-S_{11} \hat{A} \hat{A}_{2} \varphi_{2}(0)\right\} \\
\hat{\mathcal{A}}\left(\varphi_{1}, \varphi_{2}\right) & =\left(\dot{\varphi}_{1}, \dot{\varphi}_{2}\right) .
\end{aligned}
$$

In this section, we are going to prove a stronger statement than Lemma 2.2.
Lemma 3.1. If $r_{e}(\sigma(\hat{S}(t)))$ is the radius of the essential spectrum of $\hat{S}(t)$, then there is an $\alpha>0$ such that $r_{e}(\sigma(\hat{S}(t)))<e^{-\alpha t}$ for all $t \geq 0$. Also, if $e^{\mu t}$ is an eigenvalue of $\hat{S}(t)$, then $\mu$ is an eigenvalue of $\hat{\mathcal{A}}$. Furthermore, there is a $\delta>0$ such that 0 is an eigenvalue of $\hat{\mathcal{A}}$ of multiplicity two and no elements of $\sigma(\hat{\mathcal{A}}) \backslash\{0\}$ belongs to the set $|\operatorname{Re} \mu| \geq \delta$.
Proof. It is known (see [11]) that $r_{e}(\sigma(\hat{S}(t))) \leq e^{-\alpha t}$, where $\alpha$ is such that the eigenvalues of $B_{2}(0)$ have modulus $\leq e^{-\alpha}$. From (H6), we obtain the first statement of the theorem. It also is known that the eigenvalues of $\hat{S}(t)$ are given by $e^{\mu t}$, where $\mu$ is an eigenvalue of $\hat{\mathcal{A}}$, which coincide with the solutions of (2.13). The multiplicity of an eigenvalue of $\hat{\mathcal{A}}$ is the same as the multiplicity of the eigenvalues of (2.13) (see [11]). Thus, there is a $\delta>0$ such that, if $\mu$ is an eigenvalue of $\hat{\mathcal{A}}$ with $\operatorname{Re} \mu \in(-\delta, \delta)$, then $\mu$ is purely imaginary and there are only a finite number of eigenvalues with $\operatorname{Re} \mu \geq \delta$. From (H5), the only possible purely
imaginary eigenvalue is $\mu=0$. Therefore, we need only show that $\mu=0$ is an eigenvalue of multiplicity two.

Using (2.3), (H2), and (2.4), it is easy to verify that a basis for $\mathcal{N}(\hat{\mathcal{A}})$ is given by

$$
\varphi^{1}=\left[\begin{array}{l}
d_{1}  \tag{3.1}\\
d_{2}
\end{array}\right], \quad d_{2}=\left[\begin{array}{c}
e \\
-e
\end{array}\right], \quad d_{1}=\hat{A}_{2} d_{2}=\left[\begin{array}{c}
A_{2} e \\
-A_{2} e
\end{array}\right], \quad e=\left[\begin{array}{c}
1 \\
0 \\
\vdots \\
0
\end{array}\right] \in \mathbb{R}^{n}
$$

Also, a basis of $\mathcal{N}\left(\hat{\mathcal{A}}^{2}\right)$ is given by $\varphi^{1}, \varphi^{2}$, where $\varphi^{2}$ is defined by the following relations:

$$
\varphi^{2}(\theta)=\left[\begin{array}{c}
\theta d_{1}+d_{1}^{*}  \tag{3.2}\\
\theta d_{2}+d_{2}^{*}
\end{array}\right], \quad d_{2}^{*}=\frac{1}{S_{11}}\left[\begin{array}{c}
0 \\
\left(I+H_{0}(0)\right)^{-1} S_{21} \\
0 \\
-\left(I+H_{0}(0)\right)^{-1} S_{21}
\end{array}\right], \quad d_{1}^{*}=\hat{A}_{2} d_{2}^{*}+\frac{1}{S_{11}} d_{1}
$$

We now prove that $\mathcal{N}\left(\hat{\mathcal{A}}^{3}\right) \subset \mathcal{N}\left(\hat{\mathcal{A}}^{2}\right)$. If $\left(\varphi_{1}, \varphi_{2}\right) \in \mathcal{N}\left(\hat{\mathcal{A}}^{3}\right)$, then $\hat{\mathcal{A}}\left(\varphi_{1}, \varphi_{2}\right) \in \mathcal{N}\left(\hat{\mathcal{A}}^{2}\right)$ and there exist constants $\alpha, \beta$ such that $\hat{\mathcal{A}}\left(\varphi_{1}, \varphi_{2}\right)=\left(\dot{\varphi}_{1}, \dot{\varphi}_{2}\right)=\alpha \varphi^{1}+\beta \varphi^{2}$. Using (3.1), (3.2) and integrating, we deduce that there are vectors $h_{1} \in \mathbb{R}^{2 m}, h_{2} \in \mathbb{R}^{2 n}$ such that

$$
\varphi_{i}(\theta)=h_{i}+\theta\left(\alpha d_{i}+\beta d_{i}^{*}\right)+\frac{1}{2} \beta \theta^{2} d_{i}, \quad i=1,2
$$

Since $\left(\varphi_{1}, \varphi_{2}\right) \in \mathcal{D}(\hat{\mathcal{A}})$, we must have

$$
\begin{align*}
& \alpha d_{1}+\beta d_{1}^{*}=\dot{\varphi}_{1}(0)=S_{11} \hat{A}\left[\varphi_{1}(0)-\hat{A}_{2} \varphi_{2}(0)\right]=S_{11} \hat{A}\left[h_{1}-\hat{A}_{2} h_{2}\right]  \tag{3.3}\\
& h_{2}=\varphi_{2}(0)=\hat{B}_{1} \varphi_{1}(-1)+\hat{B}_{2} \varphi_{2}(-1)=\Sigma_{i=1}^{2} \hat{B}_{i}\left[h_{i}-\beta d_{i}^{*}+\left(\frac{1}{2} \beta-\alpha\right) d_{i}\right] \tag{3.4}
\end{align*}
$$

From (3.3), we have

$$
\begin{equation*}
h_{1}=\hat{A}_{2} h_{2}+\frac{1}{S_{11}} \hat{A}^{-1}\left(\alpha d_{1}+\beta d_{1}^{*}\right) \tag{3.5}
\end{equation*}
$$

If we observe that $d_{2}=\sum_{i=1}^{2} \hat{B}_{i} d_{i}, d_{2}^{*}=\sum_{i=1}^{2} \hat{B}_{i}\left(d_{i}^{*}-d_{i}\right)$, and substitute (3.5) into (3.4), we have

$$
\begin{equation*}
h_{2}=\left(\hat{B}_{1} \hat{A}_{2}+\hat{B}_{2}\right) h_{2}+\frac{1}{S_{11}} \hat{B}_{1} \hat{A}^{-1}\left(\alpha d_{1}+\beta d_{1}^{*}\right)-\beta d_{2}^{*}-\left(\frac{1}{2} \beta+\alpha\right) d_{2} \tag{3.6}
\end{equation*}
$$

If we let $h_{2}=\operatorname{col}\left(h_{21}, h_{22}\right) \in \mathbb{R}^{n} \times \mathbb{R}^{n}$, use (3.6) and the definitions of $d_{i}, d_{i}^{*}, \hat{A}_{i}, \hat{B}_{i}, i=$ 1,2 , we deduce that $h_{21}=-h_{22}$ and

$$
\begin{aligned}
h_{21} & =\left[\begin{array}{cc}
1 & 0 \\
0 & -H_{0}(0)
\end{array}\right] h_{21}+\frac{\alpha}{S_{11}} S e \\
& +\frac{\beta}{S_{11}^{2}} S\left[\begin{array}{c}
0 \\
\left(I+H_{0}(0)\right)^{-1} S_{21}
\end{array}\right]+\frac{\beta}{S_{11}^{2}} W e-\left[\begin{array}{c}
\frac{1}{2} \beta+\alpha \\
\frac{\beta}{S_{11}}\left(I+H_{0}(0)\right)^{-1} S_{21}
\end{array}\right],
\end{aligned}
$$

or

$$
\left[\begin{array}{c}
0 \\
I+H_{0}(0)
\end{array}\right] h_{21}=\left[\begin{array}{c}
\left(\frac{1}{S_{11}^{2}}\left[S_{12}\left(I+H_{0}(0)\right)^{-1} S_{21}+W_{11}\right]-\frac{1}{2}\right) \beta \\
\frac{\alpha S_{21}}{S_{11}}+\frac{\beta}{S_{11}^{2}} S_{22}\left(I+H_{0}(0)\right)^{-1} S_{21}+\frac{\beta}{S_{11}^{2}} W_{21}-\frac{\beta}{S_{11}}\left(I+H_{0}(0)\right)^{-1} S_{21}
\end{array}\right]
$$

From (H4), this implies that $\beta=0$ and, thus, $\hat{\mathcal{A}}\left(\varphi^{1}, \varphi^{2}\right)=\alpha \varphi^{1}$. Since $\varphi^{1}$ is a constant function, this implies that $\hat{\mathcal{A}}^{2}\left(\varphi^{1}, \varphi^{2}\right)=0$ and $\left(\varphi^{1}, \varphi^{2}\right) \in \mathcal{N}\left(\hat{\mathcal{A}}^{2}\right)$. This completes the proof of Lemma 3.1.
4. Proof of Theorem 2.3. Since 0 is an eigenvalue of (2.12) of multiplicity two and the spectrum of the corresponding semigroup $\hat{S}(t)$ for $t>0$ intersects the unit circle only at the point 1 , there should be a center manifold of dimension two of (2.11) which contains all of the periodic orbits. We show that such a manifold exists and we compute explicitly the vector field on this manifold up through terms of order 3. The method of computation is through the variation of constants formula which was developed in [9]. We first outline the ideas behind the variation of constants formula and then give the results of the computations without supplying the details which are rather long.

The idea is simple and consists in embedding the equation (2.8) into a neutral functional differential equation for which a variation of constants formula is well known. From the corresponding linear equation, we can determine a special decomposition of the space $C_{2 m} \times C_{2 n}$ which yields the variation of constants for the flow defined by (2.8). We consider the linear equation

$$
\begin{equation*}
\frac{d}{d t} D\left(u_{t}, v_{t}\right)=L\left(u_{t}, v_{t}\right) \tag{4.1}
\end{equation*}
$$

where

$$
D\left(\varphi_{1}, \varphi_{2}\right)=\left[\begin{array}{c}
\varphi_{1}(0) \\
\varphi_{2}(0)-\sum_{i=1}^{2} \hat{B}_{i} \varphi_{i}(-1)
\end{array}\right], \quad L\left(\varphi_{1}, \varphi_{2}\right)=\left[\begin{array}{c}
S_{11}\left(\hat{A} \varphi_{1}(0)-\hat{A} \hat{A}_{2} \varphi_{2}(0)\right) \\
0
\end{array}\right]
$$

Equation (4.1) defines a $C^{0}$-semigroup $S(t)$ on $C_{2 m} \times C_{2 n}$, whose infinitesimal generator $\mathcal{A}_{S}$ is given by

$$
\begin{aligned}
& \mathcal{D}\left(\mathcal{A}_{S}\right)=\left\{\left(\varphi_{1}, \varphi_{2}\right) \in C_{2 m} \times C_{2 n}:\left(\dot{\varphi}_{1}, \dot{\varphi}_{2}\right) \in C_{2 m} \times C_{2 n}, D\left(\dot{\varphi}_{1}, \dot{\varphi}_{2}\right)=L\left(\varphi_{1}, \varphi_{2}\right)\right\} \\
& \mathcal{A}_{S}\left(\varphi_{1}, \varphi_{2}\right)=\left(\dot{\varphi}_{1}, \dot{\varphi}_{2}\right)
\end{aligned}
$$

We know that $\mu=0$ is an eigenvalue of $\mathcal{A}_{S}$ of multiplicity $2+2 n$. If $\varphi^{1}, \varphi^{2}$ are defined in (3.1), (3.2), then we can obtain a basis of $\mathcal{N}\left(\mathcal{A}_{S}^{2+2 n}\right)$ given by

$$
\tilde{\Phi}=\left[\begin{array}{ll}
\tilde{\Phi}_{11} & \tilde{\Phi}_{12} \\
\tilde{\Phi}_{21} & \tilde{\Phi}_{22}
\end{array}\right]
$$

where, for $\theta \in[-1,0]$,

$$
\left.\left.\begin{array}{ll}
\tilde{\Phi}_{11}(\theta)=\left[\begin{array}{lll}
d_{1} & \theta d_{1}+d_{1}^{*}
\end{array}\right], & \tilde{\Phi}_{12}(\theta)=\left[\frac{1}{S_{11}} \hat{A}^{-1} d_{1}^{*}+\theta d_{1}^{*}+\frac{\theta^{2}}{2} d_{1}\right.
\end{array} \hat{A}_{2} Q\right] ~\right] ~\left[\begin{array}{lll}
d_{2} & \theta d_{2}+d_{2}^{*}
\end{array}\right], \quad \tilde{\Phi}_{22}(\theta)=\left[\begin{array}{lll}
\theta d_{2}^{*}+\frac{\theta^{2}}{2} d_{2} & Q \tag{4.2}
\end{array}\right],
$$

and $Q$ is the $2 n \times(2 n-1)$ matrix given by

$$
Q=\left[\begin{array}{ccc}
0 & 0 & 0  \tag{4.3}\\
\left(I+H_{0}(0)\right)^{-1}\left(I-H_{0}(0)\right)^{-1} & 0 & \left(I+H_{0}(0)\right)^{-1} H_{0}(0)\left(I-H_{0}(0)\right)^{-1} \\
0 & 1 & 0 \\
\left(I+H_{0}(0)\right)^{-1} H_{0}(0)\left(I-H_{0}(0)\right)^{-1} & 0 & \left(I+H_{0}(0)\right)^{-1}\left(I-H_{0}(0)\right)^{-1}
\end{array}\right]
$$

Furthermore, a basis for $\mathcal{N}\left(\left[\mathcal{A}_{S}^{*}\right]^{2+2 n}\right)$ is given by

$$
\tilde{\Psi}=\left[\begin{array}{cc}
\tilde{\Psi}_{11} & \tilde{\Psi}_{12} \\
0 & I_{2 n}
\end{array}\right]
$$

where, for $s \in[0,1]$,

$$
\tilde{\Psi}_{11}(s)=\left[\begin{array}{c}
c_{1}  \tag{4.4}\\
c_{1}^{*}+s c_{1}
\end{array}\right], \quad \tilde{\Psi}_{12}(s)=\left[\begin{array}{c}
s c_{2} \\
-s c_{2}^{*}+\frac{s^{2}}{2} c_{2}
\end{array}\right]
$$

and

$$
\begin{align*}
& c_{2}=\left[\begin{array}{ll}
e^{T} & -e^{T}
\end{array}\right], \quad c_{1}=\frac{1}{S_{11}} c_{2} \hat{B}_{1} \hat{A}^{-1} \\
& c_{2}^{*}=\frac{1}{S_{11}}\left[\begin{array}{llll}
0 & S_{12}\left(I+H_{0}(0)\right)^{-1} & 0 & -S_{12}\left(I+H_{0}(0)\right)^{-1}
\end{array}\right],  \tag{4.5}\\
& c_{1}^{*}=\frac{1}{S_{11}}\left(c_{2}-c_{2}^{*}\right) \hat{B}_{1} \hat{A}^{-1}-\frac{c_{2}}{S_{11}^{2}} \hat{B}_{1} \hat{A}^{-2} .
\end{align*}
$$

If we define the bilinear form $(\cdot, \cdot)$ for the neutral equation (4.1) 12 as

$$
\begin{aligned}
(\psi, \varphi)= & \psi_{1}(0) \varphi_{1}(0)+\psi_{2}(0)\left[\varphi_{2}(0)-\Sigma_{i=1}^{2} \hat{B}_{i} \varphi_{i}(-1)\right] \\
& -\int_{-1}^{0} \dot{\psi}_{2}(\xi+1) \Sigma_{i=1}^{2} \hat{B}_{i} \varphi_{i}(\xi) d \xi
\end{aligned}
$$

where $\psi=\left(\psi_{1}, \psi_{2}\right) \in C_{2 m}^{*} \times C_{2 n}^{*}, \varphi=\left(\varphi_{1}, \varphi_{2}\right) \in C_{2 m} \times C_{2 n}, C_{k}^{*}=C\left([0,1] ; \mathbb{R}^{k *}\right)$, then, with a few computations, we deduce that

$$
(\tilde{\Psi}, \tilde{\Phi})=\left[\begin{array}{cc}
P_{11} & P_{12}  \tag{4.6}\\
0 & P_{22}
\end{array}\right]
$$

where

$$
P_{11}=\left[\begin{array}{cc}
0 & \frac{R_{0}}{S_{11}^{2}}  \tag{4.7}\\
-\frac{R_{0}}{S_{11}^{2}} & c_{1}^{*} d_{1}^{*}+c_{2}^{*} d_{2}^{*}-\frac{2}{3}
\end{array}\right],
$$

$$
P_{12}=\left[\begin{array}{cccc}
p_{11} & \frac{1}{S_{11}} S_{12}(I+H(0))^{-1} & 0 & -\frac{1}{S_{11}} S_{12}(I+H(0))^{-1}  \tag{4.8}\\
p_{21} & \bar{R}_{0} & \frac{R_{0}}{2 S_{11}^{2}} & -\bar{R}_{0}
\end{array}\right],
$$

$$
P_{22}=\left[\begin{array}{cccc}
-\frac{R_{0}}{2 S_{2}^{2}} & 0 & 1 & 0  \tag{4.9}\\
R_{1}^{* 1} & I_{n-1} & 0 & 0 \\
\frac{R_{0}}{2 S_{2}^{2}} & 0 & 1 & 0 \\
-R_{1}^{*} & 0 & 0 & I_{n-1}
\end{array}\right],
$$

and

$$
\begin{align*}
p_{11} & =\frac{1}{S_{11}} c_{1} \hat{A}^{-1} \hat{A}_{2} d_{2}^{*}+\frac{1}{S_{11}^{2}} c_{1} \hat{A}^{-2} \hat{A}_{2} d_{2}-\frac{1}{S_{11}} c_{2} \hat{B}_{1} \hat{A}^{-1} d_{1}^{*}+\frac{2}{3}, \\
p_{21} & =\frac{1}{S_{11}}\left[c_{1}^{*} \hat{A}^{-1} d_{1}^{*}+c_{2}^{*} \hat{B}_{1} \hat{A}^{-1} d_{1}^{*}+c_{2}^{*} d_{2}^{*}+\frac{1}{4}\right], \\
\bar{R}_{0} & =\frac{1}{S_{11}^{2}}\left[S_{11} S_{12}-\left(S_{12}(I+H(0))^{-1} S_{22}+W_{21}\right)\right](I+H(0))^{-1},  \tag{4.10}\\
R_{1}^{*} & =\frac{1}{S_{11}^{2}}\left[(I+H(0))^{-1} S_{21}-W_{21}-S_{22}(I+H(0))^{-1} S_{21}\right] .
\end{align*}
$$

Since $R_{0} \neq 0$, the matrices $P_{11}, P_{22}$ are nonsingular. As a consequence, $(\tilde{\Psi}, \tilde{\Phi})$ is nonsingular. If we change the bases by the transformation

$$
\Psi=\left[\begin{array}{cc}
P_{11}^{-1} & -P_{11}^{-1} P_{12} P_{22}^{-1}  \tag{4.11}\\
0 & I_{2 n}
\end{array}\right] \tilde{\Psi}=\left[\begin{array}{cc}
P_{11}^{-1} \tilde{\Psi}_{11} & P_{11}^{-1} \tilde{\Psi}_{12}-P_{11}^{-1} P_{12} P_{22}^{-1} \\
0 & I_{2 n}
\end{array}\right]
$$

$$
\Phi=\tilde{\Phi}\left[\begin{array}{cc}
I_{2} & 0 \\
0 & P_{22}^{-1}
\end{array}\right]=\left[\begin{array}{cc}
\tilde{\Phi}_{11} & \tilde{\Phi}_{12} P_{22}^{-1} \\
\tilde{\Phi}_{21} & \tilde{\Phi}_{22} P_{22}^{-1}
\end{array}\right]
$$

then

$$
\begin{equation*}
(\Psi, \Phi)=I_{2+2 n} \tag{4.13}
\end{equation*}
$$

The explicit expressions for $P_{11}^{-1}$ and $P_{22}^{-1}$ are given by

$$
\begin{gathered}
P_{11}^{-1}=\left[\begin{array}{cc}
p^{*} & -\frac{S_{11}^{2}}{R_{0}} \\
\frac{S_{11}^{2}}{R_{0}} & 0
\end{array}\right], \quad p^{*}=\frac{S_{11}^{4}}{R_{0}^{2}}\left(c_{1}^{*} d_{1}^{*}+c_{2}^{*} d_{2}^{*}-\frac{2}{3}\right) \\
P_{22}^{-1}=\left[\begin{array}{cccc}
-\frac{S_{11}^{2}}{R_{0}} & 0 & \frac{S_{11}^{2}}{R_{0}} & 0 \\
\frac{S_{11}^{2} R_{1}^{*}}{R_{0}} & I_{n} & -\frac{S_{11}^{2} R_{1}^{*}}{R_{0}} & 0 \\
\frac{1}{2} & 0 & \frac{1}{2} & 0 \\
-\frac{S_{11}^{2} R_{1}^{*}}{R_{0}} & 0 & \frac{S_{11}^{2} R_{1}^{*}}{R_{0}} & I_{n}
\end{array}\right]
\end{gathered}
$$

A few elementary calculations also yield

$$
\begin{equation*}
\mathcal{A}_{S} \Phi=\Phi B \tag{4.14}
\end{equation*}
$$

where

$$
\begin{gather*}
B=\left[\begin{array}{cc}
B_{11} & B_{12} \\
0 & 0
\end{array}\right], \\
B_{11}=\left[\begin{array}{ll}
0 & 1 \\
0 & 0
\end{array}\right], \quad B_{12}=\left[\begin{array}{cccc}
0 & 0 & 0 & 0 \\
-\frac{S_{11}^{2}}{R_{0}} & 0 & \frac{S_{11}^{2}}{R_{0}} & 0
\end{array}\right] . \tag{4.15}
\end{gather*}
$$

We now decompose the solution space $\hat{\Omega}_{0}$ of (2.12) using the above bases $\Psi$ and $\Phi$. Let $Q^{+}$denote the union of the generalized eigenspaces of the eigenvalues of $\hat{\mathcal{A}}$ with real parts $>0$, let $Q^{+*}$ denote the union of the generalized eigenspaces of the eigenvalues of $\hat{\mathcal{A}}^{*}$ with real parts $>0$ and define

$$
Q^{-}=\left\{\varphi=\left(\varphi_{1}, \varphi_{2}\right) \in C_{2 m} \times C_{2 n}:(\Psi, \varphi)=0,(\psi, \varphi)=0 \text { for all } \psi \in Q^{+*}\right\}
$$

It is shown in [9] that

$$
\hat{\Omega}_{0}=Q^{-} \oplus\left[\Phi_{1}\right] \oplus Q^{+}
$$

where [•] denotes span and

$$
\Phi_{1}=\left[\begin{array}{l}
\tilde{\Phi}_{11}  \tag{4.16}\\
\tilde{\Phi}_{21}
\end{array}\right]
$$

and $\tilde{\Phi}_{11}, \tilde{\Phi}_{21}$ are defined in (4.2). Relations (4.14), (4.15) imply that

$$
\begin{equation*}
\tilde{\Phi}_{1}(\sigma)=\Phi_{1}(0) e^{B_{11} \sigma}, \sigma \in \mathbb{R} \tag{4.17}
\end{equation*}
$$

We introduce the notation

$$
\bar{F}_{h \lambda}(\varphi)=\left[\begin{array}{c}
0 \\
F_{h \lambda}\left(\varphi_{1}(0), \varphi_{2}(0)\right)
\end{array}\right], \quad \bar{G}_{h \lambda}(\varphi)=\left[\begin{array}{c}
0 \\
G_{h \lambda}\left(\varphi_{1}(-1), \varphi_{2}(-1)\right)
\end{array}\right],
$$

$X_{0}(\theta)=0$ for $\theta \in[-1,0), X_{0}(0)=I_{2 m+2 n}$ and

$$
\begin{equation*}
\Theta(\varphi)=\varphi+X_{0} \bar{G}_{h \lambda}(\varphi) \tag{4.18}
\end{equation*}
$$

where $\varphi=\operatorname{col}\left(\varphi_{1}, \varphi_{2}\right) \in C_{2 m} \times C_{2 n}$. Using the variation of constants formula in [9] and the standard methods in the theory of center manifolds, it is possible to show that a center manifold of (2.1) in a neighborhood of 0 and $h, \lambda$ sufficiently small, has the form

$$
\varphi=\Phi_{1} U+W_{h \lambda}(U)+X_{0} \bar{G}_{h \lambda}\left(\Phi_{1} U\right), \quad W_{h \lambda}(U) \in Q=Q^{-} \oplus Q^{+}, U \in \mathbb{R}^{2}
$$

Furthermore, if we make the transformation of variables

$$
\left[\begin{array}{l}
u_{t} \\
v_{t}
\end{array}\right]=\Phi_{1} U(t)+W_{h \lambda}(U(t))+X_{0} \bar{G}_{h \lambda}\left(\Phi_{1} U(t)\right)
$$

in (2.11), then the flow on the center manifold is given by

$$
\begin{align*}
\frac{d U(t)}{d t}= & \left.B_{11} U(t)\right) \Psi_{11}(0) \bar{F}_{h \lambda}\left(\Theta\left(\Phi_{1} U(t)+W_{h \lambda}(U(t))\right)\right)  \tag{4.19}\\
& +\left[B_{11} \Psi_{12}(0)+B_{12}\right] \bar{G}_{h \lambda}\left(\Phi_{1} U(t)+W_{h \lambda}(U(t))\right.
\end{align*}
$$

Furthermore, the function $W_{h \lambda}(U)$ satisfies the integral equation

$$
\begin{align*}
W_{h \lambda}(U(t))= & -\int_{-\infty}^{t} d_{\sigma} \hat{S}(t-\sigma) X_{0}^{Q^{-}} \bar{G}_{h \lambda}\left(\Phi_{1} U(\sigma)+W_{h \lambda}(U(\sigma))\right) \\
& +\int_{-\infty}^{t} \hat{S}(t-\sigma) X_{0}^{Q^{-}} \bar{F}_{h \lambda}\left(\Theta\left(\Phi_{1} U(\sigma)+W_{h \lambda}(U(\sigma))\right)\right) d \sigma  \tag{4.20}\\
& +\int_{t}^{\infty} d_{\sigma} \hat{S}(t-\sigma) X_{0}^{Q^{+}} \bar{G}_{h \lambda}\left(\Phi_{1} U(\sigma)+W_{h \lambda}(U(\sigma))\right) \\
& \left.-\int_{t}^{\infty} \hat{S}(t-\sigma) X_{0}^{Q^{+}} \bar{F}_{h \lambda}\left(\Phi_{1} U(\sigma)+W_{h \lambda}(U(\sigma))\right)\right) d \sigma
\end{align*}
$$

We now want to obtain an approximation of the vector field on the center manifold. It turns out that the only terms in $W_{h \lambda}(U)$ that are important are the cubic terms in $U$ of the Taylor series of the function $W_{00}(U)$. The necessary terms in $h, \lambda$ are obtained directly from the vector field in (2.1).

It is clear that

$$
W_{h \lambda}(U)=O\left((|h|+|\lambda|)|U|+|U|^{2}\right)
$$

as $(|h|,|\lambda|,|U|) \rightarrow(0,0,0)$. By using the same argument as in [], we deduce that

$$
\begin{align*}
W_{00}(U)= & -\int_{-\infty}^{0} d_{\sigma} \hat{S}(-\sigma) X_{0}^{Q^{-}}\left[\frac{1}{2} D_{\varphi}^{2} \bar{G}_{00}(0)\left(\Phi_{1}(\sigma) U\right)^{2}\right] \\
& +\int_{-\infty}^{0} \hat{S}(-\sigma) X_{0}^{Q^{-}}\left[\frac{1}{2} D_{\varphi}^{2} \bar{F}_{00}(0)\left(\Phi_{1}(\sigma) U\right)^{2}\right] d \sigma \\
& +\int_{0}^{\infty} d_{\sigma} \hat{S}(-\sigma) X_{0}^{Q^{+}}\left[\frac{1}{2} D_{\varphi}^{2} \bar{G}_{00}(0)\left(\Phi_{1}(\sigma) U\right)^{2}\right]  \tag{4.21}\\
& -\int_{0}^{\infty} \hat{S}(-\sigma) X_{0}^{Q^{+}}\left[\frac{1}{2} D_{\varphi}^{2} \bar{F}_{00}(0)\left(\Phi_{1}(\sigma) U\right)^{2}\right] d \sigma \\
& +O\left(|U|^{3}\right)
\end{align*}
$$

as $|U| \rightarrow 0$.

As in [7], we can show that, for any $\zeta \in \mathbb{R}^{2 n+2 m}$ and any integer $k \geq 0$,

$$
\begin{align*}
\int_{0}^{\infty} \hat{S}(\sigma) \sigma^{k} X_{0}^{Q^{-}} \zeta d \sigma & -\int_{-\infty}^{0} \hat{S}(\sigma) \sigma^{k} X_{0}^{Q^{+}} \zeta d \sigma  \tag{4.22}\\
& =k!(-1)^{k+1}\left(\mathcal{A}_{\hat{S}} \mid Q\right)^{-(k+1)} X_{0}^{Q} \zeta
\end{align*}
$$

If $d_{1}, d_{2}$ are defined as in (3.2) and we let

$$
\xi=\operatorname{col}\left(d_{1}, d_{2}\right), \quad U=\operatorname{col}\left(U_{1}, U_{2}\right)
$$

use (4.16), apply (4.22) for $k=0$, and perform several computations, we conclude that

$$
\begin{equation*}
W_{00}(U)=\operatorname{col}\left(W_{1}^{0}(U), W_{2}^{0}(U)\right)+O\left(\left|U_{2}\right||U|+|U|^{3}\right) \tag{4.23}
\end{equation*}
$$

as $|U| \rightarrow 0$, where

$$
\begin{align*}
W_{1}^{0}(U) & =\frac{1}{2} U_{1}^{2}\left[\begin{array}{l}
A_{2} \gamma+b \\
A_{2} \gamma+b
\end{array}\right], \quad W_{2}^{0}(U)=\frac{1}{2} U_{1}^{2}\left[\begin{array}{l}
-X_{0} a+\gamma \\
-X_{0} a+\gamma
\end{array}\right] \\
a & =D_{(x, y)}^{2} g_{0}(0)\left(A_{2} e, e\right), \quad b=D_{y}^{2} f_{0}(0)(e, e)  \tag{4.24}\\
\gamma & =M\left[a+B_{1} b\right], \quad M=\left[\begin{array}{cc}
\frac{1}{2} & 0 \\
0 & \left(I+H_{0}(0)\right)^{-1}
\end{array}\right] .
\end{align*}
$$

With the above estimate on $W_{00}(U)$, we can estimate the vector field in (4.19) for $(h, \lambda)=(0,0)$. Using the definitions of $\Phi_{1}$ in (4.16) and $\Theta$ in (4.18), we can verify that

$$
\left.\begin{array}{c}
\Psi_{11}(0) \bar{F}_{00}\left(\Theta\left(\Phi_{1} U+W_{00}(U)\right)\right) \\
=\left[\begin{array}{c}
O\left(|U|^{3}\right) \\
\frac{2 S_{11}^{2}}{R_{0}} e^{T} B_{1}\left[f^{*}\left(e, \frac{1}{2} \gamma\right)+\frac{1}{3!} D_{y}^{3} f_{0}(0)(e)^{3}\right] U_{1}^{3}
\end{array}\right] \\
\left(B_{11} \Psi_{12}(0)+B_{12}\right) \bar{G}_{00}\left(\Phi_{10} U+W_{00}(U)\right) \\
=\left[\begin{array}{c}
2 S_{11}^{2} \\
R_{0}
\end{array} e^{T}\left[g^{*}\left(\left(A_{2} e, e\right), \frac{1}{2}\left(A_{2} \gamma+b, \gamma\right)\right)+\frac{1}{3!} D_{(x, y)}^{3} g_{0}(0)\left(A_{2} e, e\right)^{3}\right] U_{1}^{3}\right] \tag{4.26}
\end{array}\right],
$$

where the second component in each of the formulas (4.25), (4.26) are $O\left(\left|U_{2}\right||U|^{2}+|U|^{4}\right)$ as $|U| \rightarrow 0$, and where $f^{*}, g^{*}$ are the quadratic forms defined by

$$
\begin{align*}
f^{*}(y, \bar{y}) & =y^{T} D_{y}^{2} f_{0}(0) \bar{y} \\
g^{*}((\bar{x}, \bar{y}),(\tilde{x}, \tilde{y})) & =\left[\begin{array}{ll}
\bar{x}^{T} & \bar{y}^{T}
\end{array}\right] D_{(x, y)}^{2} g_{0}(0)\left[\begin{array}{c}
\tilde{x} \\
\tilde{y}
\end{array}\right] . \tag{4.27}
\end{align*}
$$

For $(h, \lambda)=(0,0)$, equations (4.25)-(4.27) give us the relevant terms in the Taylor series expansion of the vector field in (4.19). In the next result, we show that the coefficient of $U_{1}^{3}$ is related to the hypothesis (H3) imposed for the generic period doubling of the map associated with the system (2.1). In fact, except for a nonzero constant factor, the sum of the terms in $U_{1}^{3}$ in (4.25), (4.26) is equal to the constant $R_{1}$ in (H3).

Lemma 4.1. If $R_{1}$ is defined as in (H3), then

$$
\begin{align*}
R_{1}= & e^{T} B_{1}\left[f^{*}\left(e, \frac{1}{2} \gamma\right)+\frac{1}{3!} D_{y}^{3} f_{0}(0)(e)^{3}\right]  \tag{4.28}\\
& +e^{T}\left[g^{*}\left(\left(A_{2} e, e\right), \frac{1}{2}\left(A_{2} \gamma+b, \gamma\right)\right)+\frac{1}{3!} D_{(x, y)}^{3} g_{0}(0)\left(A_{2} e, e\right)^{3}\right]
\end{align*}
$$

Proof. From the definition of $\mathcal{F}_{\lambda}(y)$ in (2.2), we deduce that

$$
\begin{align*}
& \mathcal{F}_{0}(y)=\left(B_{1} A_{2}+B_{2}\right) y+\frac{1}{2} B_{1} D_{y}^{2} f_{0}(0)(y)^{2}+\frac{1}{2} D_{(x, y)}^{2} g_{0}(0)\left(A_{2} y, y\right)^{2}  \tag{4.29}\\
& \quad+g^{*}\left(\left(A_{2} y, y\right),\left(\frac{1}{2} D_{y}^{2} f_{0}(0)(y)^{2}, 0\right)\right)+\frac{1}{3!} B_{1} D_{y}^{3} f_{0}(0)(y)^{3}+\frac{1}{3!} D_{(x, y)}^{3} g_{0}(0)\left(A_{2} y, y\right)^{3} .
\end{align*}
$$

From (4.29), we deduce that

$$
\begin{align*}
k_{3}(0)= & e^{T}\left\{g^{*}\left(\left(A_{2} e, e\right),\left(\frac{1}{2} D_{y}^{2} f_{0}(0)(e)^{2}, 0\right)\right)\right. \\
& \left.+\frac{1}{3!} B_{1} D_{y}^{3} f_{0}(0)(e)^{3}+\frac{1}{3!} D_{(x, y)}^{3} g_{0}(0)\left(A_{2} e, e\right)^{3}\right\},  \tag{4.30}\\
k_{2}(0) & =e^{T}\left[B_{1} f^{*}\left(e, I_{*}\right)+g^{*}\left(\left(A_{2} e, e\right),\left(A_{2} I_{*}, I_{*}\right)\right)\right], \tag{4.31}
\end{align*}
$$

where $I_{*}$ is the $m \times(m-1)$ matrix $\operatorname{col}\left(0, I_{m-1}\right)$,

$$
\begin{align*}
& k_{1}(0)=\frac{1}{2} e^{T}\left[B_{1} f^{*}(e, e)+g^{*}\left(\left(A_{2} e, e\right),\left(A_{2} e, e\right)\right)\right]  \tag{4.32}\\
& H_{1}(0)=\frac{1}{2} I_{*}^{T}\left[B_{1} f^{*}(e, e)+g^{*}\left(\left(A_{2} e, e\right),\left(A_{2} e, e\right)\right)\right] \tag{4.33}
\end{align*}
$$

From the definition of $M$ and $\gamma$ in (4.24), we see that

$$
\begin{aligned}
\gamma & =\frac{1}{2} e^{T}\left(a+B_{1} b\right) e+I_{*}\left(I-H_{0}(0)\right)^{-1} I_{*}^{T}\left(a+B_{1} b\right) \\
& =k_{1}(0) e+2 I_{*}\left(I-H_{0}(0)\right)^{-1} H_{1}(0)
\end{aligned}
$$

which together with (4.32) allows us to conclude that

$$
\begin{align*}
e^{T} B_{1} & {\left[f^{*}\left(e, \frac{1}{2} \gamma\right)+\frac{1}{3!} D_{y}^{3} f_{0}(0)(e)^{3}\right] } \\
= & \frac{1}{2} k_{1}(0) e^{T} B_{1} f^{*}(e, e)+e^{T} B_{1} f^{*}\left(e, I_{*}\right)(I-H(0))^{-1} H_{1}(0)  \tag{4.34}\\
& +\frac{1}{3!} e^{T} B_{1} D_{y}^{3} f_{0}(0)(e)^{3}
\end{align*}
$$

In a similar way, we verify that

$$
\begin{align*}
&\left.e^{T}\left[g^{*}\left(\left(A_{2} e, e\right), \frac{1}{2}\left(A_{2} \gamma+b, \gamma\right)\right)\right)+\frac{1}{3!} D_{(x, y)}^{3} g_{0}(0)\left(A_{2} e, e\right)^{3}\right] \\
&= \frac{1}{2} k_{1}(0) e^{T} g^{*}\left(\left(A_{2} e, e\right),\left(A_{2} e, e\right)\right)  \tag{4.35}\\
&+\frac{1}{2} e^{T} g^{*}\left(\left(A_{2} e, e\right),\left(A_{2} I_{*}, I_{*}\right)\right)\left(I-H_{0}(0)\right)^{-1} H_{1}(0) \\
&+e^{T} g^{*}\left(\left(\left(A_{2} e, e\right),\left(\frac{1}{2} D_{y}^{2} f_{0}(0)(e)^{2}, 0\right)\right)+\frac{1}{3!} D_{(x, y)}^{3} g_{0}(0)\left(A_{2} e, e\right)^{3} .\right.
\end{align*}
$$

If we now use (4.30)-(4.35), we obtain the relation (4.28) and the lemma is proved.
As remarked above, the relevant terms in $h, \lambda$ in the Taylor series of the vector field in (4.16) are easier to obtain since they will not depend upon any specific knowledge of the function $W_{h \lambda}(U)$. After several rather straightforward but lengthy computations, we deduce that

$$
\begin{gather*}
\Psi_{11}(0)\left[\bar{F}_{h \lambda}\left(\Theta\left(\Phi_{1} U+W_{h \lambda}(U)\right)\right)-\bar{F}_{00}\left(\Theta\left(\Phi_{1} U+W_{00}(U)\right)\right)\right. \\
+\left[B_{11} \Psi_{12}(0)+B_{12}\left[\bar{G}_{h \lambda}\left(\Phi_{1} U+W_{h \lambda}(U)\right)-\bar{G}_{00}\left(\Phi_{1} U+W_{00}(U)\right)\right.\right.  \tag{4.36}\\
=A_{h \lambda} U+O\left((|\lambda|+|h|)|U|^{2}+(|\lambda|+|h|)^{2}|U|\right)
\end{gather*}
$$

as $(h, \lambda, U) \rightarrow(0,0,0)$, where

$$
A_{h \lambda}=\left[\begin{array}{cc}
\alpha \lambda & \beta_{1} \lambda+\beta_{2} h  \tag{4.37}\\
-\frac{2 S_{11}^{2}}{R_{0}} \lambda & \frac{2 S_{11}}{R_{0}} h+\ell^{*} \lambda
\end{array}\right]
$$

and $\alpha, \beta_{1}, \beta_{2}, \ell^{*}$ are constants.
From Lemma 4.1, formulas (4.25)-(4.27) and (4.36), (4.37), we conclude that the vector field in (4.19) is given by

$$
\begin{align*}
\dot{U}_{1}= & \alpha \lambda U_{1}+\left(1+\beta_{1} \lambda+\beta_{2} h\right) U_{2}+O\left((|\lambda|+|h|)|U|^{2}+(|\lambda|+|h|)^{2}|U|+|U|^{3}\right) \\
\dot{U}_{2}=- & \frac{2 S_{11}^{2}}{R_{0}} \lambda U_{1}+\left(\frac{2 S_{11}}{R_{0}} h+\ell^{*} \lambda\right) U_{2}+\frac{2 S_{11}^{2} R_{1}}{R_{0}} U_{1}^{3}+\ell_{1}^{*} U_{1}^{2} U_{2}+\ell_{2}^{*} U_{1} U_{2}^{2}+\ell_{3}^{*} U_{2}^{3}  \tag{4.38}\\
& +O\left((|\lambda|+|h|)|U|^{2}+(|\lambda|+|h|)^{2}|U|+|U|^{4}\right),
\end{align*}
$$

where $\ell^{*}, \ell_{1}^{*}, \ell_{2}^{*}, \ell_{3}^{*}$ are constants.
To analyze the periodic solutons of (4.38), it is convenient to write them in a different form. Since $R_{0} \neq 0, S_{11} \neq 0$, we can make a linear change of variables of the form $U \mapsto$ $Z=(I+P(h, \lambda)) U$ and change the time scale $t \mapsto(1+\delta(h, \lambda)) t$, where $P(0)=0, \delta(0)=0$ such that the new equation for $Z$ is given by

$$
\begin{align*}
\dot{Z}_{1}= & Z_{2}+O\left((|\lambda|+|h|)|Z|^{2}+(|\lambda|+|h|)^{2}|Z|+|Z|^{3}\right) \\
\dot{Z}_{2}= & -\frac{2 S_{11}^{2}}{R_{0}} \lambda Z_{1}+\left(\frac{2 S_{11}}{R_{0}} h+\ell^{*} \lambda\right) Z_{2}+\frac{2 S_{11}^{2} R_{1}}{R_{0}} Z_{1}^{3}+\ell_{1}^{*} Z_{1}^{2} Z_{2}+\ell_{2}^{*} Z_{1} Z_{2}^{2}+\ell_{3}^{*} Z_{2}^{3}  \tag{4.39}\\
& +O\left((|\lambda|+|h|)|Z|^{2}+(|\lambda|+|h|)^{2}|Z|+|Z|^{4}\right) .
\end{align*}
$$

This equation now has the same form as the one considered in [9]. The discussion there for the existence of periodic solutions corresponding to periodic solutons of (2.1) can be repeated verbatum to complete the proof of Theorem 2.3.
5. Hopf bifurcation curve. For the periodic orbit whose existence is given by Theorem 2.3, we want to impose additional conditions on the vector field in (2.1) which will ensure that this orbit occured through a Hopf bifurcation and that supercritical (subcritical) period doubling of the map corresponds to stability (instability) of the periodic orbit of (2.1). In this section, we specify conditions which will imply the existence of such a Hopf bifurcation curve and, in the next section, we consider the stability properties.

We consider the linear system

$$
\begin{align*}
& \epsilon \dot{x}(t)+A x(t)=A A_{2}(\lambda) y(t)  \tag{5.1}\\
& y(t)=B_{1}(\lambda) x(t-1)+B_{2}(\lambda) y(t-1)
\end{align*}
$$

for which the characteristic matrix is

$$
\Delta(\lambda, \epsilon, \mu)=\left[\begin{array}{cc}
\epsilon \mu I_{m}+A & -A A_{2}(\lambda)  \tag{5.2}\\
-B_{1}(\lambda) e^{-\mu} & I_{n}-B_{2}(\lambda) e^{-\mu}
\end{array}\right]
$$

It is convenient to introduce the following definition. We say that a curve $\Gamma$ in the $(\lambda, \epsilon)$-plane, $\epsilon>0$, is a Hopf Bifurcation Curve of (5.1) if there is an $\epsilon^{*}>0$ and a continuous function $\lambda=\lambda(\epsilon), 0<\epsilon \leq \epsilon^{*}, \lambda(\epsilon) \rightarrow 0$ as $\epsilon \rightarrow 0$ such that, if $\Gamma_{\epsilon^{*}}=\{(\lambda, \epsilon):$ $\left.\lambda=\lambda(\epsilon), \epsilon \in\left(0, \epsilon^{*}\right]\right\}$, then, for any $\left(\lambda_{0}, \epsilon_{0}\right) \in \Gamma_{\epsilon^{*}}$, there are two purely imaginary solutions $\pm i \beta_{0}$ of the characteristic equation $\operatorname{det} \Delta\left(\lambda_{0}, \epsilon_{0}, \mu\right)=0$ and the remaining solutions $\mu$ satisfy $\operatorname{Re} \mu \neq 0$. We say that a Hopf Bifurcation Curve is a Generic Hopf Bifurcation Curve with respect to $\epsilon$ if, for fixed $\lambda_{0}$, the two eigenvalues $\mu\left(\lambda_{0}, \epsilon\right), \bar{\mu}\left(\lambda_{0}, \epsilon\right), \mu\left(\lambda_{0}, \epsilon_{0}\right)=i \beta_{0}$, satisfy $\operatorname{Re} d \mu\left(\lambda_{0}, \epsilon_{0}\right) / d \epsilon<0$. This type of transversal crossing of the imaginary axis of the eigenvalue $\mu\left(\lambda_{0}, \epsilon\right)$ implies that there will be a Hopf bifurcation with respect to $\epsilon$ at $\epsilon_{0}$. We say that a Hopf Bifurcation Curve is the First Hopf Bifurcation Curve if, for each fixed $\epsilon_{0} \in\left(0, \epsilon^{*}\right]$, all eigenvalues $\mu$ of (5.1) with $\lambda=\lambda\left(\epsilon_{0}\right)$ and $\epsilon>\epsilon_{0}$ have $\operatorname{Re} \mu<0$. The Generic First Hopf Bifurcation Curve is the most interesting because there is a transfer of stability of the origin at $\epsilon=\epsilon_{0}$; that is, the origin is stable for $\epsilon>\epsilon_{0}$ and unstable for $\epsilon<\epsilon_{0}$. From the physical origins of the problem, this is natural because we expect that the origin is stable for large $\epsilon$ (by a change of time scale, this is small delay) and eventually becomes unstable for small $\epsilon$ (large delay). The First Hopf Bifurcation Curve represents the first change in the stability properties of the origin.

We want to determine the First Hopf Bifurcation Curve for (5.1) which is generic with respect to $\epsilon$.

We retain the condition (H2) which corresponds to $\lambda=0$ being a period doubling for the map. To have all of the relevant hypotheses close at hand, we repeat this hypothesis.

$$
B_{1}(\lambda) A_{2}(\lambda)+B_{2}(\lambda) \equiv C(\lambda)=\left[\begin{array}{cc}
-(1+\lambda) & 0  \tag{H2}\\
0 & H_{0}(\lambda)
\end{array}\right], \quad \sigma\left(H_{0}(0) \subset\{z:|z|<1\}\right.
$$

With the notation,

$$
-B_{1}(0) A^{-1} A_{2}(0) \equiv S=\left[\begin{array}{ll}
S_{11} & S_{12}  \tag{5.3}\\
S_{21} & S_{22}
\end{array}\right], \quad-B_{1}(0) A^{-2} A_{2}(0) \equiv W=\left[\begin{array}{ll}
W_{11} & W_{12} \\
W_{21} & W_{22}
\end{array}\right]
$$

we also retain the hypothesis

$$
\begin{equation*}
R_{0} \equiv-S_{11}^{2}+2\left[W_{11}+S_{12}\left(I+H_{0}(0)\right)^{-1} S_{21}\right] \neq 0, \quad S_{11} \neq 0 \tag{H4}
\end{equation*}
$$

Of course, we must keep the hypothesis

$$
\begin{equation*}
\sigma\left(B_{2}(0)\right) \subset B_{\rho}, \quad \rho<1 \tag{H6}
\end{equation*}
$$

The additional hypotheses that we need are

$$
\begin{equation*}
\min \{\operatorname{Re} z: z \in \sigma(A)\}>0 \tag{H7}
\end{equation*}
$$

$$
\operatorname{Det}\left[\begin{array}{cc}
i \theta I_{m}+A & -A A_{2}(0)  \tag{H8}\\
-B_{1}(0) e^{-i v} & I_{n}-B_{2}(0) e^{-i v}
\end{array}\right] \neq 0, \text { for } \theta>0,0 \leq v \leq 2 \pi,(\theta, v) \neq(0, \pi),
$$

$$
\operatorname{Det}\left[\begin{array}{cc}
\mu I_{m}+A & -A A_{2}(0)  \tag{H9}\\
-B_{1}(0) & I_{n}-B_{2}(0)
\end{array}\right] \neq 0, \text { for all } \mu \in \mathbb{C}, \operatorname{Re} \mu \geq 0
$$

The main result of this section is
Theorem 5.1. Under the assumptions (H2), (H4), (H6) and (H7), the hypotheses (H8) and (H9) are necessary and sufficient for the existence of the First Hopf Bifurcation Curve which is generic with respect to $\epsilon$.

The proof will be given in terms of several lemmas which bring out the role of each hypothesis.

Lemma 5.2. Under the hypotheses (H2), (H4), (H6) and (H7), there is an $\epsilon^{*}>0$ and a $C^{2}$-curve $\Gamma_{\epsilon^{*}}=\left\{\left((\lambda(\epsilon), \epsilon), \epsilon \in\left(0, \epsilon^{*}\right\}\right.\right.$ such that (i) $\lambda(\epsilon) \rightarrow 0$ as $\epsilon \rightarrow 0$, (ii) for any $\left(\lambda_{0}, \epsilon_{0}\right) \in \Gamma_{\epsilon^{*}}$, there is a unique pair of purely imaginary eigenvalues $\pm i v\left(\lambda_{0}, \epsilon_{o}\right)$ with $v \in(0,2 \pi)$. Furthermore, $\Gamma_{\epsilon^{*}}$ is the unique curve satisfying (i) and (ii).

Proof. If

$$
G(\alpha, \lambda)=B_{1}(\lambda)\left(\alpha I_{m}+A\right)^{-1} A_{2}(\lambda),
$$

then, under assumption (H7), we can verify that det $\Delta(\lambda, \epsilon, i v)=0$ if and only if det $\left[e^{i v} I_{n}-C(\lambda)+i \epsilon v G(i \epsilon v, \lambda)\right]=0$; that is, if and only if there is a nonzero vector $h \in \mathbb{C}^{\boldsymbol{n}}$ such that

$$
\begin{equation*}
\left[e^{i v} I_{n}-C(\lambda)+i \epsilon v G(i \epsilon v, \lambda)\right] h=0 \tag{5.4}
\end{equation*}
$$

If we set $\lambda=\epsilon=0$, then (5.4) implies that $\left[e^{i v} I_{n}-C(0)\right] h=0$, which, from (H2), implies that $v=\pi$ and $h=\gamma \operatorname{col}[1,0]$, where $\gamma$ is a constant. For $\lambda, \epsilon$ small, it follows that the first component of h in (5.4) is not zero and we can take $h=\operatorname{col}[1, \xi]$. With this observation, we can determine the solutions of (5.4) near $(\lambda, \epsilon, h)=(0,0, \operatorname{col}[1,0])$ by determining the zeros of the function

$$
\begin{align*}
& \mathcal{F}: \hat{I} \times I \times \mathbb{C}^{n-1} \rightarrow \mathbb{C}^{n} \\
& \mathcal{F}(v, \lambda, \epsilon, \xi)=\left[e^{i v} I_{n}-C(\lambda)+i \epsilon v G(i \epsilon v, \lambda)\right]\left[\begin{array}{l}
1 \\
\xi
\end{array}\right] \tag{5.5}
\end{align*}
$$

where $\hat{I} \subset \mathbb{R}$ is a neighborhood of $\pi$ and $I \subset \mathbb{R}$ is a neighborhood of 0 .
We have $\mathcal{F}(\pi, 0,0,0)=0$. If we define

$$
\left.T \equiv \frac{\partial \mathcal{F}(v, \lambda, \epsilon, \xi)}{\partial(v, \lambda, \xi)}\right|_{(\pi, 0,0,0)}: \mathbb{R} \times \mathbb{R} \times \mathbb{C}^{n-1} \rightarrow \mathbb{C}^{n} l
$$

then

$$
T(v, \lambda, \xi)=\left[\begin{array}{c}
-i v+\lambda \\
-\left(I+H_{0}(0)\right) \xi
\end{array}\right] .
$$

Since $I+H_{0}(0)$ is invertible, the map $T$ is invertible. Therefore, the Implicit Function Theorem implies that there are $C^{2}$-functions $v(\epsilon), \lambda(\epsilon), \xi(\epsilon)$ defined for $\epsilon \in\left(-\epsilon^{*}, \epsilon^{*}\right), \epsilon^{*}>0$, such that $v(0)=\pi, \lambda(0)=0 \xi(0)=0$, and

$$
\begin{equation*}
\mathcal{F}(v(\epsilon), \lambda(\epsilon), \epsilon, \xi(\epsilon))=0, \quad \epsilon \in\left(-\epsilon^{*}, \epsilon^{*}\right) \tag{5.6}
\end{equation*}
$$

Differentiating (5.6) with respect to $\epsilon$, we obtain

$$
\begin{align*}
& \dot{v}\left[i e^{i v} I_{n}+i \epsilon G(i \epsilon v, \lambda)-\epsilon^{2} v \frac{\partial}{\partial \alpha} G(i \epsilon v, \lambda)\right]\left[\begin{array}{l}
1 \\
\xi
\end{array}\right]+\dot{\lambda}\left[-\dot{C}(\lambda)+i \epsilon v \frac{\partial}{\partial \lambda} G(i \epsilon v, \lambda)\right]\left[\begin{array}{l}
1 \\
\xi
\end{array}\right]  \tag{5.7}\\
& +\left[e^{i v} I_{n}-C(\lambda)+i \epsilon v G(i \epsilon v, \lambda)\right]\left[\begin{array}{l}
0 \\
\dot{\xi}
\end{array}\right]+\left[i v G(i \epsilon v, \lambda)-\epsilon v^{2} \frac{\partial}{\partial \alpha} G(i \epsilon v, \lambda)\right]\left[\begin{array}{l}
1 \\
\xi
\end{array}\right]=0 .
\end{align*}
$$

If we set $\epsilon=0$ and use the fact that $-G(0,0)=S$, we obtain

$$
\left[\begin{array}{c}
-i \ddot{v}(0)+\dot{\lambda}(0) \\
-\left(I+H_{0}(0)\right) \dot{\xi}(0)
\end{array}\right]=-i v(0) G(0,0)\left[\begin{array}{c}
1 \\
\xi(0)
\end{array}\right]=i \pi\left[\begin{array}{c}
S_{11} \\
S_{21}
\end{array}\right]
$$

Hence,

$$
\begin{equation*}
\dot{v}(0)=-\pi S_{11}, \quad \dot{\lambda}(0)=0, \quad \dot{\xi}(0)=-i \pi\left(I+H_{0}(0)\right)^{-1} S_{21} . \tag{5.8}
\end{equation*}
$$

Differentiating (5.7) with respect to $\epsilon$, setting $\epsilon=0$ and using (5.8), we deduce that

$$
\begin{aligned}
{\left[\left(\dot{v}^{2}(0)\right.\right.} & \left.-i \ddot{v}(0)) I_{n}+2 i \dot{v}(0) G(0,0)-\dot{\lambda}(0) \dot{C}(0)-2 v^{2}(0) \frac{\partial}{\partial \alpha} G(0,0)\right]\left[\begin{array}{l}
1 \\
0
\end{array}\right] \\
& =2\left[i \dot{v}(0) I_{n}-i v(0) G(0,0)\right]\left[\begin{array}{l}
0 \\
\dot{\xi}
\end{array}\right]+\left[I_{n}+C(0)\right]\left[\begin{array}{l}
0 \\
\ddot{\xi}
\end{array}\right]
\end{aligned}
$$

From this relation and the fact that $\frac{\partial}{\partial \alpha} G(0,0)=W$, we have

$$
\dot{v}^{2}(0)+\ddot{\lambda}(0)-2 v^{2}(0) W_{11}=2 i\left[\dot{v}(0)-v(0) S_{12} \dot{\xi}(0)\right] .
$$

From this relation and (5.8), we have

$$
\ddot{\lambda}(0)=2 \pi^{2} W_{11}-\pi^{2} S_{11}^{2}+2 \pi^{2} S_{12}\left(I+H_{0}(0)\right)^{-1} S_{21}=\pi^{2} R_{0} \neq 0 .
$$

As a consequence, for small $\epsilon$, we have the Hopf bifurcation curve $\Gamma_{\epsilon^{*}}=\{(\lambda(\epsilon), \epsilon), \epsilon \in$ $\left.\left.\left(0, \epsilon^{*}\right)\right]\right\}$, where

$$
\begin{equation*}
\lambda(\epsilon)=\frac{1}{2} \pi^{2} R_{0} \epsilon^{2}+o\left(\epsilon^{2}\right) \tag{5.9}
\end{equation*}
$$

as $\epsilon \rightarrow 0$. This completes the proof of the lemma.
Remark 5.1. The expression (5.9) shows that the graph of the curve $\Gamma_{\epsilon^{*}}$ is over the positive (resp. negative) $\lambda$-axis if $R_{0}>0$ (resp. $R_{0}<0$ ).
Lemma 5.3. Suppose that the hypotheses (H2), (H4), (H6), (H7) and (H8) are satisfied and let $\Gamma_{\epsilon^{*}}=\left\{(\lambda(\epsilon), \epsilon), \epsilon \in\left(0, \epsilon^{*}\right]\right\}$ be given by Lemma 5.2. Then there is an $\epsilon_{*}>0$ such that $\Gamma_{\epsilon_{*}}$ is a Generic Hopf Bifurcation Curve.
Proof. We first show that there is an $\epsilon_{*}>0$ such that, for each $\epsilon_{0} \in\left(0, \epsilon_{*}\right]$, we have

$$
\begin{equation*}
\operatorname{det} \Delta\left(\lambda\left(\epsilon_{0}\right), \epsilon, i \beta\right) \neq 0, \text { for } \epsilon>\epsilon_{0}, \beta \geq 0 \tag{5.10}
\end{equation*}
$$

If this is not the case, then we can find sequences $\left\{\epsilon_{n}\right\},\left\{\hat{\epsilon}_{n}\right\},\left\{\beta_{n}\right\}$ in $(0, \infty)$ such that $\epsilon_{n}<\hat{\epsilon}_{n}, \epsilon_{n} \rightarrow 0$ as $n \rightarrow \infty$, and

$$
0=\operatorname{det}\left(\lambda\left(\epsilon_{n}\right), \hat{\epsilon}_{n}, i \beta_{n}\right) \equiv\left[\begin{array}{cc}
i \hat{\epsilon}_{n} \beta_{n} I_{m}+A & -A A_{2}\left(\lambda\left(\epsilon_{n}\right)\right)  \tag{5.11}\\
-B_{1}\left(\lambda\left(\epsilon_{n}\right)\right) e^{-i \beta_{n}} & I_{n}-B_{2}\left(\lambda\left(\epsilon_{n}\right)\right) e^{-i \beta_{n}}
\end{array}\right]
$$

From (5.11), it follows that $\hat{\epsilon}_{n} \beta_{n}, n \geq 1$, is bounded. Since $\lambda\left(\epsilon_{n}\right) \rightarrow 0$ as $n \rightarrow \infty$ and we are assuming (H8), we must have $\hat{\epsilon}_{n} \beta_{n} \rightarrow 0, e^{-i \beta_{n}} \rightarrow-1$ as $n \rightarrow \infty$. If we let $k_{n} \geq 0$ be an integer and $v_{n} \in(0,2 \pi]$ be such that $\beta_{n}=2 k_{n} \pi+v_{n}$, then we have $v_{n} \rightarrow \pi$ as $n \rightarrow \infty$. Also, if $\epsilon_{n}^{\prime}=\hat{\epsilon}_{n} \beta_{n} / v_{n}$, then

$$
\begin{equation*}
\lim _{n \rightarrow \infty} \epsilon_{n}^{\prime}=0, \quad \epsilon_{n}^{\prime} \geq \hat{\epsilon}_{n}>\epsilon_{n} \tag{5.12}
\end{equation*}
$$

and

$$
\begin{equation*}
\operatorname{det} \Delta\left(\lambda\left(\epsilon_{n}\right), \epsilon_{n}^{\prime}, i v_{n}\right)=\operatorname{det} \Delta\left(\lambda\left(\epsilon_{n}\right), \hat{\epsilon}_{n}, i \beta_{n}\right)=0 \tag{5.13}
\end{equation*}
$$

From (5.13) and Lemma 5.2, it follows that $\lambda\left(\epsilon_{n}\right)=\lambda\left(\epsilon_{n}^{\prime}\right)$ for sufficiently large $n$. From (5.9), the function $\lambda(\epsilon)$ is monotone for $\epsilon$ near 0 and, thus, we must have $\epsilon_{n}=\epsilon_{n}^{\prime}$, which contradicts (5.12). This completes the proof of relation (5.10).

We next show that, if $\mu=i v\left(\epsilon_{0}\right)$ is a zero of $\Delta\left(\lambda\left(\epsilon_{0}\right), \epsilon_{0}, \mu\right)$, then it is a simple zero. For fixed $\epsilon_{0} \in\left(0, \epsilon_{*}\right]$, let $\lambda_{0}=\lambda\left(\epsilon_{0}\right), \mu_{0}=i v\left(\epsilon_{0}\right)$ and $\xi_{0}=\xi\left(\epsilon_{0}\right)$, where $\lambda(\epsilon), v(\epsilon), \xi(\epsilon)$ are defined in (5.6). We define the function

$$
\begin{align*}
& \mathcal{K}: \mathbb{C} \times \mathbb{C}^{n-1} \times \mathbb{R} \rightarrow \mathbb{C}^{n} \\
& \mathcal{K}(\mu, \eta, \epsilon)=\left[e^{\mu} I_{n}-C\left(\lambda_{0}\right)+\epsilon \mu G\left(\epsilon \mu, \lambda_{0}\right)\right]\left[\begin{array}{l}
1 \\
\eta
\end{array}\right] . \tag{5.14}
\end{align*}
$$

We have $\mathcal{K}\left(\mu_{0}, \xi_{0}, \epsilon_{0}\right)=0$. Using (5.8) and (5.9), we deduce that

$$
\frac{\partial \mathcal{K}\left(\mu_{0}, \xi_{0}, \epsilon_{0}\right)}{\partial(\mu, \eta)}=\left[\begin{array}{cc}
-1 & 0 \\
0 & -\left(I+H_{0}(0)\right)
\end{array}\right]+O\left(\epsilon_{0}\right)
$$

as $\epsilon_{0} \rightarrow 0$. Thus, if $\epsilon_{*}$ is sufficiently small, this matrix has nonzero determinant and the Implicit Function Theorem implies that there are $C^{2}$-functions $\mu(\epsilon), \eta(\epsilon)$ defined for $\epsilon$ in a neighborhood $U$ of 0 such that $\mu\left(\epsilon_{0}\right)=i v\left(\epsilon_{0}\right), \eta\left(\epsilon_{0}\right)=\xi_{0}$, and

$$
\begin{equation*}
\mathcal{K}(\mu(\epsilon), \eta(\epsilon))=0, \quad \epsilon \in U \tag{5.15}
\end{equation*}
$$

In particular, $\mu=i v_{0}$ is a simple zero of $\operatorname{det} \Delta\left(\lambda_{0}, \epsilon_{0}, \mu\right)$.
We next determine how $\mu(\epsilon)$ depends upon $\epsilon$. Differentiating (5.15) with respect to $\epsilon$, we obtain

$$
\begin{gather*}
\dot{\mu}\left(\epsilon_{0}\right)\left[e^{i v_{0}} I_{n}+\epsilon_{0} G\left(i \epsilon_{0} v_{0}, \lambda_{0}\right)\right]\left[\begin{array}{c}
1 \\
\xi_{0}
\end{array}\right]+\left[e^{i v_{0}} I_{n}-C\left(\lambda_{0}\right)+i \epsilon_{0} v_{0} G\left(i \epsilon_{0} v_{0}, \lambda_{0}\right)\right]\left[\begin{array}{c}
0 \\
\dot{\eta}\left(\epsilon_{0}\right)
\end{array}\right]  \tag{5.16}\\
=-\left[i v_{0} G\left(i \epsilon_{0} v_{0}, \lambda_{0}\right)-\epsilon_{0} v_{0}^{2} \frac{\partial}{\partial \alpha} G\left(i \epsilon_{0} v_{0}, \lambda_{0}\right)\right]\left[\begin{array}{c}
1 \\
\xi_{0}
\end{array}\right]+O\left(\epsilon_{0}^{2}\right) .
\end{gather*}
$$

If we use (5.8), the relations $G(0,0)=-S, \partial G(0,0) / \partial \alpha=W,(5.16)$ and Talyor's theorem around $\epsilon_{0}$, we conclude that

$$
\begin{align*}
\dot{\mu}\left(\epsilon_{0}\right)\left[-\left(1+\epsilon_{0} S_{11}\right)+i \epsilon_{0} \pi S_{11}\right]= & i \epsilon_{0} \pi S_{12} \dot{\eta}\left(\epsilon_{0}\right)+i \pi S_{11}+i \pi S_{12} \xi_{0}+2 \epsilon_{0} \pi^{2} W_{11} \\
& +O\left(\epsilon_{0}^{2}\right)  \tag{5.17}\\
-\left[I+H_{0}(0)\right] \dot{\eta}\left(\epsilon_{0}\right)= & i \pi S_{21}+O\left(\epsilon_{0}\right)
\end{align*}
$$

The second relation of (5.17) implies that

$$
\dot{\eta}\left(\epsilon_{0}\right)=-i \pi\left[I+H_{0}(0)\right]^{-1} S_{21}+O\left(\epsilon_{0}\right)
$$

If we substitute this expression and $\xi_{0}=-i \pi\left(1+H_{0}(0)\right)^{-1} S_{21}+O\left(\epsilon_{0}\right)$ into (5.17), we conclude that

$$
\begin{aligned}
\dot{\mu}\left(\epsilon_{0}\right)[-(1 & \left.\left.+\epsilon_{0} S_{11}\right)+i \epsilon_{0} \pi S_{11}\right] \\
& =i \pi S_{11}+2 \epsilon_{0} \pi^{2}\left[W_{11}+S_{12}\left(I+H_{0}(0)\right)^{-1} S_{21}\right]+O\left(\epsilon_{0}^{2}\right)
\end{aligned}
$$

Using this relation and the definition of $R_{0}$ in (H4), we have

$$
\begin{equation*}
\operatorname{Re} \dot{\mu}\left(\epsilon_{0}\right)=-\epsilon_{0} R_{o} \pi^{2}+O\left(\epsilon_{0}^{2}\right) \tag{5.18}
\end{equation*}
$$

as $\epsilon_{0} \rightarrow 0$. For $\epsilon_{*}$ sufficiently small, this shows that the Hopf Bifurcation Curve is Generic and the lemma is proved.

We now complete the proof of sufficiency in Theorem 5.1 using these two lemmas and hypothesis (H9). We only need to prove that, for each $\epsilon_{0} \in\left(0, \epsilon_{*}\right]$, and all $\epsilon>\epsilon_{0}$, there is no zero $\mu(\epsilon)$ of the characteristic matrix in (5.2) with $\operatorname{Re} \mu(\epsilon)>0$. If $\operatorname{det} \Delta\left(\lambda\left(\epsilon_{0}\right), \hat{\epsilon}, \mu\right)=0$ has a solution $\mu(\hat{\epsilon})$ with $\operatorname{Re} \mu(\hat{\epsilon})>0$ for some $\hat{\epsilon}>\epsilon_{0}$, then we can extend this function to a continuous function $\mu(\epsilon)$ for $\epsilon \in\left(\epsilon_{0}, \infty\right)$ and have $\operatorname{det} \Delta\left(\lambda\left(\epsilon_{0}\right), \epsilon, \mu(\epsilon)\right)=0$ for $\epsilon \in\left(\epsilon_{0}, \infty\right)$. From Lemma 5.3 and, in particular, the implication (5.10), we must have $\operatorname{Re} \mu(\epsilon)>0$ for all $\epsilon \in\left(\epsilon_{0}, \infty\right)$. If we let $z(\epsilon)=\epsilon \mu(\epsilon)$, then

$$
\operatorname{det}\left[\begin{array}{cc}
z(\epsilon) I_{m}+A & -A A_{2}\left(\lambda\left(\epsilon_{0}\right)\right)  \tag{5.19}\\
-B_{1}\left(\lambda\left(\epsilon_{0}\right)\right) e^{-\frac{z(\epsilon)}{\epsilon}} & I_{n}-B_{2}\left(\lambda\left(\epsilon_{0}\right)\right) e^{-\frac{z(\epsilon)}{\epsilon}}
\end{array}\right]=\operatorname{det} \Delta\left(\lambda\left(\epsilon_{0}\right), \epsilon, \mu(\epsilon)\right)=0 .
$$

Since $\operatorname{Re} \frac{z(\epsilon)}{\epsilon}>0$, relation (5.19) implies that $z(\epsilon)$ is bounded in $\left(\epsilon_{0}, \infty\right)$. Hence there is a subsequence $\epsilon_{j} \rightarrow \infty$ as $j \rightarrow \infty$ such that $z\left(\epsilon_{j}\right) \rightarrow z^{*}$. It is clear that $\operatorname{Re} z^{*} \geq 0$. By letting $j \rightarrow \infty$ in (5.19), we deduce that

$$
\operatorname{det}\left[\begin{array}{cc}
z^{*} I_{m}+A & -A A_{2}\left(\lambda\left(\epsilon_{0}\right)\right) \\
-B_{1}\left(\lambda\left(\epsilon_{0}\right)\right) & I_{n}-B_{2}\left(\lambda\left(\epsilon_{0}\right)\right)
\end{array}\right]=0 .
$$

This contradicts (H9) if $\epsilon_{*}$ is small enough. This completes the proof of the existence of the First Hopf Bifurcation Curve under the hypotheses of the theorem.

The proof of the necessity can be supplied by carefully considering the arguments that were used in the proof of sufficiency. We do not give the details.

Let us now consider some special cases. The equation (1.3) can be written as a system

$$
\begin{align*}
\epsilon \dot{x}(t)+x(t) & =f_{\lambda}(y(t))  \tag{5.20}\\
y(t) & =x(t-1) .
\end{align*}
$$

The hypotheses (H1)-(H3) for generic period doubling of the map $x \mapsto f_{\lambda}(x)$ are equivalent to saying that the function $f_{\lambda}(x)$ can be written as

$$
\begin{equation*}
f_{\lambda}(x)=-(1+\lambda) x+\alpha x^{2}+\beta x^{3} \tag{5.21}
\end{equation*}
$$

where $\alpha^{2}+\beta \neq 0$. The linear variational equation about $x=0$ for (5.21) is a special case of (5.1) with $A=1, A_{2}(\lambda)=-(1+\lambda), B_{1}(\lambda)=1, B_{2}(\lambda)=0$. It is now obvious that the hypotheses (H2), (H4), (H6), (H7) are satisfied. A simple computation shows that hypotheses (H8) is equivalent to $i(\theta-\sin v)+1+\cos v \neq 0$ for $\theta>0,0 \leq v \leq 2 \pi, v \neq \pi$. This is clearly satisfied. Also, (H9) is equivalent to $\mu+2 \neq 0$ for $\operatorname{Re} \mu \geq 0$, which is true. Therefore, there is a Generic First Bifurcation Curve with respect to $\epsilon$.

Let us next discuss equation (1.4) with $\epsilon_{j}=\epsilon \alpha_{j}, j=1,2, \ldots, m$, and the function $f_{\lambda}(x)$ satisfying (5.21). It is possible to write this equation in a matrix form (5.1) and verify directly the hypotheses for the existence of a Generic Hopf Bifurcation Curve are satisfied. Rather than do this, we will show this curve exists by analyzing the characteristic equation. The proof is not difficult and also is a simple illustration of the ideas used in the proof of Theorem 2.3. The characteristic equation for the linearization of (1.4) about $x=0$ is given by

$$
\begin{equation*}
E(\epsilon, \lambda, \mu) \equiv\left(\epsilon \alpha_{m} \mu+1\right) \cdots\left(\epsilon \alpha_{1} \mu+1\right)+(1+\lambda) e^{-\mu}=0 \tag{5.22}
\end{equation*}
$$

In this particular situation, the determinant in (H8) is given by

$$
E_{1}(\theta, v) \equiv E\left(\frac{\theta}{v}, 0, i v\right)=\prod_{j=1}^{m}\left(i \theta \alpha_{j}+1\right)+e^{-i v}
$$

If $\theta \geq 0,0 \leq v<2 \pi$, are such that $E_{1}(\theta, v)=0$, then we have $\left|\prod_{j=1}^{m}\left(i \theta \alpha_{j}+1\right)\right|=1$; that is, $\prod_{j=1}^{m}\left(\theta^{2} \alpha_{j}^{2}+1\right)=1$. Since $\alpha_{j}>0$, we must have $\theta=0$ and hence $1+e^{-i v}=0$; that is, $v=\pi$. Therefore, the condition (H8) is satisfied.

Next, the function in (H9) has the form

$$
C_{2}(\mu) \equiv \prod_{j=1}^{m}\left(\mu \alpha_{j}+1\right)+1
$$

If there is a $\mu=u+i v, u \geq 0$, such that $C_{2}(u+i v)=0$, then it follows that $1=$ $\prod_{j=1}^{m}\left(\left(u \alpha_{j}+1\right)^{2}+v^{2}\right)$. Therefore, we must have $u=v=0$, which leads to the assertion that $0=C_{2}(0)=1+1$, which is a contradiction. As a consequence, $C_{2}(\mu) \neq 0$ for all $\mu \in \mathbb{C}$ with $\operatorname{Re} \mu \geq 0$; that is, (H9) is satisfied.

We can now show that, for fixed $\epsilon_{0}>0$, there is a unique $\lambda^{*}\left(\epsilon_{0}\right)>0$ such that (5.22) has exactly two purely imaginary roots and the remaining ones have negative real parts for $(\lambda, \epsilon)=\left(\lambda\left(\epsilon_{0}\right), \epsilon_{0}\right)$. For $\epsilon>\epsilon_{0}$, the origin is asymptotically stable and, for $0<\epsilon<\epsilon_{0}$, the origin is unstable. In this way, we obtain the existence of a Generic First Hopf Bifurcation Curve.

In the general matrix case, the problem becomes more complicated. More specifically, we can have the map associated with the differential equation undergo generic period doubling and have the linearization about the origin in the differential equation not possess a Generic First Hopf Bifurcation Curve. For example, let us consider the system

$$
\begin{align*}
\epsilon \dot{x}_{1}(t)+a x_{1}(t) & =a f_{\lambda}\left(x_{1}(t-1)\right) \\
\epsilon \dot{x}_{2}(t)+B x_{2}(t) & \left.=B \nu x_{2}(t-1)\right), \tag{5.23}
\end{align*}
$$

where $f_{\lambda}$ is the same function as before, $x_{2} \in \mathbb{R}, B, \nu$ are constants, $B<0,|\nu|<1$. The map obtained for $\epsilon=0$ is given by $\left(x_{1}, x_{2}\right) \mapsto\left(f_{\lambda}\left(x_{1}\right), \nu x_{2}\right)$ and it undergoes a generic period doubling at $\left(x_{1}, x_{2}, \lambda\right)=(0,0,0)$. The hypotheses (H2), (H4) and (H6) are satisfied, but we have not satisfied hypothesis (H7).

If we want to relate the period doubling bifurcation to a Generic First Hopf Bifurcation Curve, then we must have the solutions of the linear variational equation for $\lambda=0$ approach zero as $t \rightarrow \infty$; that is, the eigenvalues must have negative real parts. This is equivalent to having the solutions of the equation

$$
\begin{equation*}
\Delta(\mu) \equiv \mu+B\left(1-\nu e^{-\mu}\right)=0 \tag{5.24}
\end{equation*}
$$

having negative real parts. Since $\Delta(0)=B(1-\mu)<0$ and $\Delta(\mu) \rightarrow \infty$ as $\mu \rightarrow \infty$, there is a positive zero for $\Delta(\mu)$.
6. The Hopf bifurcation. Let us suppose that the hypotheses of Theorem 2.3 and Theorem 5.1 are satisfied. Then, as shown in Section 5, there is a first Hopf Bifurcation Curve $\Gamma_{\epsilon^{*}}=\left\{\left((\lambda(\epsilon), \epsilon), \epsilon \in\left(0, \epsilon^{*}\right\}\right.\right.$. Furthermore, if $R_{0}>0$, then, for each fixed $\lambda$, the zero solution is stable if $\epsilon<\epsilon^{*}$ and unstable if $\epsilon>\epsilon^{*}$. At $\epsilon=\epsilon^{*}$, the linearization of (2.4) at the origin has a pair of purely imaginary simple eigenvalues. As a consequence, we expect that there is a generic Hopf bifurcation arising at $\epsilon$. This is true and, in fact, we have

Theorem 6.1. Suppose that the hypotheses of Theorem 2.3 and Theorem 5.1 are satisfied. If, in addition, $R_{0}>0$, then there exists an $\epsilon^{*}>0$ such that, for each fixed $\left(\lambda_{0}, \epsilon_{0}\right) \in \Gamma_{\epsilon^{*}}$, $\lambda_{0}=\lambda\left(\epsilon_{0}\right)$, where $\Gamma_{\epsilon^{*}}$ is the First Hopf Bifurcation Curve, system (2.1) has a generic Hopf bifurcation at $(\lambda, \epsilon, x, y)=\left(\lambda_{0}, \epsilon_{0}, 0,0\right)$ which is supercritical (resp. subcritical) if $R_{1}>0$ (resp. $R_{1}<0$ ).

We remark that this result implies that the generic period doubling point $\lambda=0$ for the mapping $\mathcal{F}_{\lambda}$ is supercritical (resp. subcritical) if and only if the generic Hopf bifurcation at ( $\lambda_{0}, \epsilon_{0}$ ) is supercritical (resp. subcritical).

To prove Theorem 6.1, we are going to use some results from [10] for more general hybrid systems in which they gave a specific formula for a constant $\alpha^{*}$ depending upon the second and third derivatives of the nonlinearities in (2.1) evaluated at zero which determines the direction of the Hopf bifurcation. More specifically, the bifurcation is supercritical (resp. subcritical) if $\alpha^{*}>0$ (resp. $\alpha^{*}<0$ ). We now describe this constant, attempting to motivate each step and refer the reader to [10] for details.

Let

$$
\begin{gathered}
\tilde{f}(y, \lambda)=f(y, \lambda)-A_{2}(\lambda) y \\
\tilde{g}(x, y, \lambda)=g(x, y, \lambda)-B_{1}(\lambda) x-B_{2}(\lambda) y
\end{gathered}
$$

The linearized system of (2.1) is

$$
\begin{align*}
\epsilon \dot{x}(t) & =-A x(t)+A_{2}(\lambda) y(t) \\
y(t) & =B_{1}(\lambda) x(t-1)+B_{2}(\lambda) y(t-1) \tag{6.1}
\end{align*}
$$

From our hypotheses on $f$ and $g$, there is an $\epsilon^{*}>0$ such that, for each $0<\epsilon_{0}<\epsilon^{*}$, there is a unique $\lambda_{\epsilon_{0}}>0$ such that the linear system (6.1) for $\epsilon=\epsilon_{0}, \lambda=\lambda_{\epsilon_{0}}$ has an imaginary eigenvalue $\mu_{0}=i \nu_{0}=i \pi+O\left(\epsilon_{0}\right)$ and it is simple. For fixed $\lambda_{\epsilon_{0}}$, the eigenvalue $\mu_{0}(\epsilon)$ of (6.1) with $\mu_{0}\left(\epsilon_{0}\right)=\mu_{0}$ satisfies the inequality $\dot{\mu}_{0}\left(\epsilon_{0}\right)<0$.

Along with (6.1), we consider the linear neutral equation

$$
\begin{align*}
\epsilon \dot{x}(t) & =-A x(t)+A A_{2}(\lambda) y(t) \\
\frac{d}{d t}[y(t) & \left.-B_{1}(\lambda) x(t-1)-B_{2}(\lambda) y(t-1)\right]=0 . \tag{6.2}
\end{align*}
$$

This equation plays an important role in the variation of constants formula for a perturbation of (6.1) and, thus, an important role in the explicit computation of approximation of periodic orbits. The nonzero eigenvalues of (6.1) and (6.2) as well as their multiplicies coincide.

The characteristic matrix for (6.2) is

$$
\Delta\left(\epsilon_{0}, \mu\right)=\left[\begin{array}{cc}
\mu I+\frac{1}{\epsilon_{0}} A & -\frac{1}{\epsilon_{0}} A A_{2}\left(\lambda_{\epsilon_{0}}\right) \\
-\mu B_{1}\left(\lambda_{\epsilon_{0}}\right) e^{-\mu} & \mu\left(I-B_{2}\left(\lambda_{\epsilon_{0}}\right) e^{-\mu}\right)
\end{array}\right] .
$$

It is not difficult to show that there is a vector $\eta_{\epsilon_{0}}$ which satisfies $\Delta\left(\epsilon_{0}, \mu_{0}\right) \eta_{\epsilon_{0}}=0$ and has the form

$$
\eta_{\epsilon_{0}}=\operatorname{col}\left[\eta_{\epsilon_{0}}^{1}, \eta_{\epsilon_{0}}^{2}\right], \quad \eta_{\epsilon_{0}}^{1}=A_{2}\left(\lambda_{\epsilon_{0}}\right) e_{1}+O\left(\epsilon_{o}\right), \quad \eta_{\epsilon_{0}}^{2}=e_{1}+O\left(\epsilon_{0}\right), \quad e_{1}=\operatorname{col}(1,0) \in \mathbb{R}^{n} .
$$

As a consequence,

$$
\Phi^{\epsilon_{0}}(\theta)=\left[\eta_{\epsilon_{0}} e^{\mu_{0} \theta}, \bar{\eta}_{\epsilon_{0}} e^{-\mu_{0} \theta}\right] E
$$

is a basis for the eigenfunction space for $\mu_{0}$ and $\bar{\mu}_{0}$ considered as an eigenvalue of either (6.1) or (6.2). If $\Phi^{1, \epsilon_{0}}$ is the first column of $\Phi^{\epsilon_{0}}$, we have

$$
\lim _{\epsilon_{0} \rightarrow 0} \Phi^{1, \epsilon_{0}}(\theta)=\lim _{\epsilon_{0} \rightarrow 0}\left[\begin{array}{c}
\Phi^{11, \epsilon_{0}}(\theta)  \tag{6.3}\\
\Phi^{21, \epsilon_{0}}(\theta)
\end{array}\right]=\frac{1}{2}\left(e^{i \pi \theta}+e^{-i \pi \theta}\right)\left[\begin{array}{c}
A_{2}(0) e_{1} \\
e_{1}
\end{array}\right] .
$$

Now, if $\left[b_{0}^{1}, b_{0}^{2}\right] \in \mathbb{C}^{m+n^{*}}$ is such that $\left[b_{0}^{1}, b_{0}^{2}\right] \Delta\left(\epsilon_{0}, \mu_{0}\right)=0$, then we have

$$
\begin{gathered}
b_{0}^{1}=\epsilon_{0} \mu_{0} b_{0}^{2} B_{1}\left(\lambda_{\epsilon_{0}}\right) e^{-\mu_{0}}\left(\mu_{0} \epsilon_{0}+A\right)^{-1} \\
b_{0}^{2}\left[B_{0}^{1} e^{-\mu_{0}}\left(\mu_{0} \epsilon_{0}+A\right)^{-1} A A_{2}\left(\lambda_{\epsilon_{0}}\right)-\left(I-B_{2}\left(\lambda_{\epsilon_{0}}\right) e^{-\mu_{0}}\right)\right]=0 .
\end{gathered}
$$

So we can choose $b_{0}^{2}$ such that

$$
b_{0}^{2}=e_{1}^{*}+O\left(\epsilon_{0}\right), \quad e_{0}^{*}=[1,0] \in \mathbb{R}^{n *}
$$

If we let

$$
\gamma_{\epsilon_{0}}=\left[\gamma_{\epsilon_{0}}^{1}, \gamma_{\epsilon_{0}}^{2}\right], \quad \gamma_{\epsilon_{0}}^{i}=\frac{b_{0}^{i}}{b_{0}^{1} \eta_{\epsilon_{0}}^{1}+\mu_{0} b_{0}^{2} \eta_{\epsilon_{0}}^{2}}, \quad i=1,2
$$

then one is able to verify that

$$
\Psi^{\epsilon_{0}}(s)=E^{-1}\left[\begin{array}{c}
\gamma_{\epsilon_{0}} e^{-\mu_{0} s} \\
\bar{\gamma}_{\epsilon_{0}} e^{\mu_{0} s}
\end{array}\right]
$$

is a basis for the eigenfunction space corresponding to $\mu_{0}$ and $\bar{\mu}_{0}$ for the formal adjoint operator of (6.2) (not (6.1)). Moreover, we have

$$
\begin{align*}
\lim _{\epsilon_{0} \rightarrow 0} \Psi^{11, \epsilon_{0}}(s) \frac{A}{\epsilon_{0}} & =-\left(e^{i \pi s}+e^{-i \pi s}\right) e_{1}^{*} B_{1}(0)  \tag{6.4}\\
\lim _{\epsilon_{0} \rightarrow 0} \delta \Psi^{12, \epsilon_{0}}(s) & =-\left(e^{i \pi s}+e^{-i \pi s}\right) e_{1}^{*}
\end{align*}
$$

where $\Psi^{1, \epsilon_{0}}=\left[\Psi^{11, \epsilon_{0}}, \Psi^{12, \epsilon_{0}}\right]$ and $\Psi^{1, \epsilon_{0}}$ is the first row lof $\Psi^{e_{0}}$.
To obtain the direction of bifurcation, it is first necessary to determine the influence of the second order terms in (2.1) on the periodic orbit. These second order terms give rise to corrections in the solution of the form $\xi\left(\epsilon_{0}\right) e^{2 i \nu_{0} t}+\bar{\xi}\left(\epsilon_{0}\right) e^{-2 i \nu_{0} t}+\xi_{0}\left(\epsilon_{0}\right)$. It is shown in [10] that the quantity $\xi$ is given by

$$
\xi\left(\epsilon_{0}\right)=\frac{1}{8}\left[\begin{array}{cc}
2 \mu_{0}+\frac{A}{\epsilon_{0}} & -\frac{A A_{2}\left(\lambda\left(\epsilon_{0}\right)\right)}{\epsilon_{0}} \\
-B_{1}\left(\lambda\left(\epsilon_{0}\right)\right) e^{-2 \mu_{0}} & I-B_{2}\left(\lambda\left(\epsilon_{0}\right)\right) e^{-2 \mu_{0}}
\end{array}\right]^{-1}\left[\begin{array}{c}
\frac{A}{\epsilon_{0}} D_{y}^{2} f\left(0, \lambda\left(\epsilon_{0}\right)\right)\left\langle\eta_{\epsilon_{0}}^{2}\right\rangle^{2} \\
D_{z}^{2} g\left(0, \lambda\left(\epsilon_{0}\right)\right)\left\langle\eta_{\epsilon_{0}}\right\rangle^{2}
\end{array}\right]
$$

where $D_{z}=D_{(x, y)}$. By using (6.4) and letting $\xi\left(\epsilon_{0}\right)=\operatorname{col}\left[\xi^{1}\left(\epsilon_{0}\right), \xi^{2}\left(\epsilon_{0}\right)\right]$, we have

$$
\begin{align*}
& \lim _{\epsilon_{0} \rightarrow 0} \xi^{2}\left(\epsilon_{0}\right)=\xi_{0}^{2} \\
& \lim _{\epsilon_{0} \rightarrow 0} \xi^{1}\left(\epsilon_{0}\right)=\xi_{0}^{1} \tag{6.6}
\end{align*}
$$

where

$$
\begin{aligned}
& \xi_{0}^{1} \equiv A_{2}(0) \xi_{0}^{2}, \\
& \xi_{0}^{2} \equiv \frac{1}{8}\left[\begin{array}{c}
\frac{1}{2} e_{1}^{*}\left[B_{1}(0) D_{y}^{2} f(0,0)\left\langle e_{1}\right\rangle^{2}+D_{z}^{2} g(0,0)\left\langle\left[\begin{array}{c}
A_{2}(0) e_{1} \\
e_{1}
\end{array}\right]\right\rangle^{2}\right] \\
\left(I-H_{0}(0)\right)^{-1} E_{2}\left[B_{1}(0) D_{y}^{2} f(0,0)\left\langle e_{1}\right\rangle^{2}+D_{z}^{2} g(0,0)\left\langle\left[\begin{array}{c}
A_{2}(0) e_{1} \\
e_{1}
\end{array}\right]\right\rangle^{2}\right]
\end{array}\right]
\end{aligned}
$$

and $E_{2}$ is the $(n-1) \times n$ matrix:

$$
E_{2}=\left[\begin{array}{ccccccc}
0 & 1 & 0 & . & . & \cdot & 0 \\
0 & 0 & 1 & \cdot & . & \cdot & 0 \\
0 & 0 & 0 & . & . & \cdot & 0 \\
0 & 0 & 0 & . & . & \cdot & 0 \\
0 & 0 & 0 & . & . & \cdot & 0 \\
0 & 0 & 0 & \cdot & . & \cdot & 1
\end{array}\right]
$$

Similary, the constant $\xi_{0}\left(\epsilon_{0}\right)$ satisfies

$$
\lim _{\epsilon_{0} \rightarrow 0} \xi_{0}\left(\epsilon_{0}\right)=\lim _{\epsilon_{0} \rightarrow 0}\left[\begin{array}{c}
\xi_{0}^{1}\left(\epsilon_{0}\right)  \tag{6.7}\\
\xi_{0}^{2}\left(\epsilon_{0}\right)
\end{array}\right]=2\left[\begin{array}{l}
\xi_{0}^{1} \\
\xi_{0}^{2}
\end{array}\right] .
$$

From [10], the number $\alpha^{*}\left(\epsilon_{0}\right)$ for system (2.1) at $\left(\lambda\left(\epsilon_{0}\right), \epsilon_{0}\right)$ is given by

$$
\begin{aligned}
\alpha^{*}\left(\epsilon_{0}\right)= & -\int_{0}^{\frac{2 \pi}{\mu_{0}}} \Psi^{11, \epsilon_{0}}(s) \frac{A}{\epsilon_{0}} D_{y}^{2} f\left(0, \lambda\left(\epsilon_{0}\right)\right)\left\langle\Phi^{21, \epsilon_{0}}(s), \xi\left(\epsilon_{0}\right) e^{2 \mu_{0} s}+\bar{\xi}\left(\epsilon_{0}\right) e^{-2 \mu_{0} s}+\xi_{0}\left(\epsilon_{0}\right)\right\rangle d s \\
& -\frac{1}{3!} \int_{0}^{\frac{2 \pi}{\mu_{0}}} \Psi^{11, \epsilon_{0}}(s) \frac{A}{\epsilon_{0}} D_{y}^{3} f\left(0, \lambda\left(\epsilon_{0}\right)\right)\left\langle\Phi^{21, \epsilon_{0}}(s)\right\rangle^{3} d s \\
& +\int_{0}^{\frac{2 \pi}{\mu_{0}}} \dot{\Psi}^{12, \epsilon_{0}}(s) D_{z}^{2} g\left(0, \lambda\left(\epsilon_{0}\right)\right) \times \\
& \left\langle\Phi^{1, \epsilon_{0}}(s-1), \xi\left(\epsilon_{0}\right) e^{2 \mu_{0}(s-1)}+\bar{\xi}\left(\epsilon_{0}\right) e^{-2 \mu_{0}(s-1)}+\xi_{0}\left(\epsilon_{0}\right)\right\rangle d s \\
+ & \frac{1}{3!} \int_{0}^{\frac{2 \pi}{\mu_{0}}} \dot{\Psi}^{12, \epsilon_{0}}(s) D_{z}^{3} g\left(0, \lambda\left(\epsilon_{0}\right)\right)\left\langle\Phi^{1, \epsilon_{0}}(s-1)\right\rangle^{3} d s .
\end{aligned}
$$

By applying (4.4)-(4.7), we obtain

$$
\begin{aligned}
& \lim _{\epsilon_{0} \rightarrow 0} \alpha^{*}\left(\epsilon_{0}\right) \\
&= e_{1}^{*} \int_{0}^{2}\left(e^{i \pi s}+e^{-i \pi s}\right)^{2}\left(2+e^{2 i \pi s}+e^{-2 i \pi s}\right) D_{y}^{2} f(0,0)\left\langle\frac{1}{2} e_{1}, \xi_{0}^{2}\right\rangle d s \\
&+\frac{e_{1}^{*}}{3!} \int_{0}^{2}\left(e^{i \pi s}+e^{-i \pi s}\right)^{4} D_{y}^{3} f(0,0)\left\langle\frac{1}{2} e_{1}\right\rangle^{3} d s \\
&+e_{1}^{*} \int_{0}^{2}\left(e^{i \pi s}+e^{-i \pi s}\right)^{2}\left(2+e^{2 i \pi s}+e^{-2 i \pi s}\right) D_{z}^{2} g(0,0)\left\langle\frac{1}{2}\left[\begin{array}{c}
A_{2}(0) e_{1} \\
e_{1}
\end{array}\right],\left[\begin{array}{l}
\xi_{0}^{1} \\
\xi_{0}^{2}
\end{array}\right]\right\rangle d s \\
&+\frac{e_{1}^{*}}{3!} \int_{0}^{2}\left(e^{i \pi s}+e^{-i \pi s}\right)^{4} D_{z}^{3} f(0,0)\left\langle\frac{1}{2}\left[\begin{array}{c}
A_{2}(0) e_{1} \\
e_{1}
\end{array}\right]\right\rangle^{3} d s \\
&= 12 e_{1}^{*}\left[D_{y}^{2} f(0,0)\left\langle\frac{1}{2} e_{1}, \xi_{0}^{2}\right\rangle+D_{z}^{2} g(0,0)\left\langle\frac{1}{2}\left[\begin{array}{c}
A_{2}(0) e_{1} \\
e_{1}
\end{array}\right],\left[\begin{array}{c}
A_{2}(0) \xi_{0}^{2} \\
\xi_{0}^{2}
\end{array}\right]\right\rangle\right. \\
&\left.+\frac{1}{3!} D_{y}^{3} f(0,0)\left\langle\frac{1}{2} e_{1}\right\rangle^{3}+\frac{1}{3!} D_{z}^{3} g(0,0)\left\langle\frac{1}{2}\left[\begin{array}{c}
A_{2}(0) e_{1} \\
e_{1}
\end{array}\right]\right\rangle^{3}\right] \\
&= 6 R_{1},
\end{aligned}
$$

where the last equality follows from relation (2.5) and (H3).
The last equality shows that, for $\epsilon_{0}$ small, the quantity $\alpha^{*}\left(\epsilon_{0}\right)$ has the same sign as $R_{1}$. Therefore, from Lemma 2.1 and the above mentioned result from [10], we see that the Hopf bifurcation is supercritical if $R_{1}>0$ and subcritical if $R_{1}<0$. This completes the proof of Theorem 6.1.

## References

[1] M. L. Berre, E. Ressayre, A. Tallet and H. M. Gibbs, High-dimension chaotic attractors of a nonlinear ring cavity, Phys. Rev. Lett. 56 (1986), 274-277.
[2] S.-N. Chow, J. K. Hale and W. Huang, From sine waves to square waves in delay equations. Proc. Royal Soc. Edinburgh 120A (1992), 223-229.
[3] S.-N. Chow and W. Huang, Singular perturbation for delay differential equations and applications. J. Differential Equations. To appear.
[4] S.-N. Chow and J. Mallet-Paret, Singularly perturbed delay differential equations. In Coupled Nonlinear Oscillators. (Eds. J. Chandra and A. Scott) North Holland (1983).
[5] H. M. Gibbs, F. A. Hopf, D. D. L. Kaplan and R. L. Shoemaker, Observation of chaos in optical bistability, Phys. Rev. Lett. 46 (1981), 474-477.
[6] L. Glass and M. Mackey, Oscillation and chaos in physiological control systems, Science 191 (1977), 287-289.
[7] J.K. Hale and W. Huang, Period doubling in singularly perturbed delay equations. J. Differential Equations. To appear.
[8] J.K. Hale and W. Huang, Square and pulse waves in matrix delay differential equations. Dyn. Sys. Appl. 1(1992), 51-70.
[9] J.K. Hale and W. Huang, Variation of constants for hybrid systems of FDE. Proc. Royal Soc. Edinburgh. To appear.
[10] J.K. Hale and W. Huang, Hopf bifurcation analysis for hybrid systems. Proc. Int. Conf. Diff. Equations in memory of Stavros Busenburg. To appear.
[11] J.K. Hale and S. M. Verduyn-Lunel, Introduction to Functional Differential Equations, Springer-Verlag, New York, 1993.
[12] K. Ikeda, Multiple-valued stationary state and its instability of the transmitted light by a ring cavity system, Opt. Commun. 30 (1979), 257-261.
[13] K. Ikeda, H. Daido and O. Akimoto, Optical turbulence: Chaotic behavior of transmitted light from a ring cavity, Phys. Rev. Lett. 45 (1980), 709-712.
[14] M. Mackey and M. Milton, Comments on Theoretical Biology 1 (1990), 299327.
[15] J. Mallet-Paret and R. Nussbaum, Global continuation and asymptotic behavior for periodic solutions of a differential delay equation, Annali Mat. Pura Appl. 145 (1986), 33-128.
[16] J. Mallet-Paret and R. Nussbaum, Global continuation and complicated trajectories for periodic solutions for a differential delay equation, Proc. Symp. Pure Math., Am. Math. Soc. 45 (1986), Part 2, 155-167.
[17] M. Shimura, Analysis of some nonlinear phenomena in a transmission line. IEEE Trans. Circuit Theory 14 (1967), 60-69.
[18] R. Vallee, P. Dubois, M. Coté and C. Delisle, Phys. Rev. A36 (1987), 1327.
[19] R. Vallee and C. Marriott, Analysis of an $\mathrm{N}^{t h}$-order nonlinear differential delay equation, Phys. Rev. A39(1989), 197-205.
[20] M. Ważewska-Czyewska and A. Lasota, Mate. Stosowana 6 (1976), 25-40.

# CENTER FOR DYNAMICAL SYSTEMS AND NONLINEAR STUDIES Report Series 

## September 1994 -

## CDSNS94-

157 Rates of eigenvalues on a dumbbell domain. Simple eigenvalue case, J.M. Arrieta.
158 Existence of standing waves for competition-diffusion equations, Y. Kan-On.
159 Fisher-type property of travelling waves for competition-diffusion equations, Y. Kan-On.

160 Uniform ultimate boundedness and synchronization, H.M. Rodrigues.
161 Complete families of pseudotrajectories and shape of attractors, S.Yu. Pilyugin.
162 Aubry-Mather Theorem and quasiperiodic orbits for time dependent reversible systems, S.-N. Chow and M.-L. Pei.

163 Slowly-migrating transition layers for the discrete Allen-Cahn and Cahn-Hilliard equations, C.P. Grant and E.S. Van Vleck.

164 Partially hyperbolic sets from a co-dimension one bifurcation, T. Young.
165 Asymptotically autonomous semiflows: Chain recurrence and Lyapunov functions, K. Mischaikow, H. Smith and H.R. Thieme.

166 On the problem of stability in lattice dynamical systems, L.A. Bunimovich and E.A. Carlen.

167 Remark on continuous dependence of attractors on the shape of domain, A.V. Babin and S.Yu. Pilyugin.

168 Special pseudotrajectories for lattice dynamical systems, V.S. Afraimovich and S.Yu. Pilyugin.

169 Proof and generalization of Kaplan-Yorke conjecture on periodic solution of differential delay equations, J. Li and X-Z. He.

170 On the construction of periodic solutions of Kaplan-Yorke type for some differential delay equations, K. Gopalsamy, J. Li and X-Z. He.

171 Uniformly accurate schemes for hyperbolic systems with relaxation, R.E. Caflisch, S. Jin and G. Russo.

172 Numerical integrations of systems of conservation laws of mixed type, S. Jin.

173 A complex algorithm for computing Lyapunov values, R. Mao and D. Wang.
174 On minimal sets of scalar parabolic equations with skew-product structures, W. Shen and Y. Yi.

175 Lorenz type attractors from codimensional-one bifurcation, V. Afraimovich, S.-N. Chow and W. Liu.

176 Dynamics in a discrete Nagumo equation - spatial chaos, S.-N. Chow and W. Shen.

177 Stability and bifurcation of traveling wave solutions in coupled map lattices, S.-N. Chow and W. Shen.

178 Density of defects and spatial entropy in extended systems, V.S. Afraimovich and L.A. Bunimovich.

179 On the second eigenvalue of the Laplace operator penalized by curvature, E.M. Harrell II.

180 Singular limits for travelling waves for a pair of equations, V. Hutson and K. Mischaikow.

181 Conley Index Theory: Some recent developments, K. Mischaikow.
182 Variational principle for periodic trajectories of hyperbolic billiards, L.A. Bunimovich.

183 Epidemic waves: A diffusion model for fox rabies, W.M. Rivera.
184 High complexity of spatial patterns in gradient reaction-diffusion systems, V. Afraimovich, A. Babin and S.-N. Chow.

185 Spatial chaotic structure of attractors of reaction-diffusion systems, V. Afraimovich, A. Babin and S.-N. Chow.

186 Explosive non-negative solutions to two point boundary value problems, V. Anuradha, C. Brown and R. Shivaji.

187 Existence results for non-autonomous elliptic boundary value problems, V. Anuradha, S. Dickens, and R. Shivaji.

188 Existence results for superlinear semipositone BVP's, V. Anuradha, D.D. Hai and R. Shivaji.

189 Hopf bifurcation analysis for hybrid systems, J.K. Hale and W. Huang.
190 Dynamics of cyclic feedback systems, T. Gedeon and K. Mischaikow.
191 Cyclic feedback systems, T. Gedeon.

192 Traveling waves as limits of solutions on bounded domains, G. Fusco, J.K. Hale and J. Xun.

193 Equivalence of topological and singular transition matrices in the Conley index theory, C.K. McCord and K. Mischaikow.

194 Coupled oscillators on a circle, J.K. Hale.
195 Conventional multipliers for homoclinic orbits, V. Afraimovich, W. Liu and T. Young.

196 Wave train of monotone travelling waves for generic competition-diffusion systems, Y. Kan-on.

197 Instability of stationary solutions for Lotka-Volterra competition model with diffusion, Y. Kan-on.

198 Hyperbolic homoclinic points of $\mathbb{Z}^{d}$-actions in lattice dynamical systems, V.S. Afraimovich, S-N. Chow and W. Shen.

199 Least energy solutions of semilinear Neumann problems and asymptotics, X-B. Pan and X. Xu.

200 Effects of delays on dynamics, J.K. Hale.
201 Periodic solutions of singularly perturbed delay equations, J.K. Hale and W. Huang.

202 Uniform ultimate boundedness and synchronization for nonautonomous equations, V. Afraimovich and H.M. Rodrigues.

203 Algebraic transition matrices, R. Franzosa and K. Mischaikow.

## CDSNS95-

204 Ergodicity of minimal sets in scalar parabolic equations, W. Shen and Y. Yi.
205 Existence of infinitely many connecting orbits in a singularly perturbed ordinary differential equation, H. Kokubu, K. Mischaikow and H. Oka.

206 Asymptotic periodicity of positive solutions of reaction diffusion equations on a ball, X-Y. Chen and P. Poláčik.

207 Asymptotic behaviour of ground states, J. Hulshof and R. C.A.M. van der Vorst.
208 Synchronizations in lattices of nonlinear Duffing's oscillators, V.S. Afraimovich and S-N. Chow.

209 Population patterns for ecological models, H.I. Freedman and W.M. Rivera.
210 Generalization of a theorem of Malta and Palis, V. Afraimovich and T. Young.
211 Singular index pairs, K. Mischaikow, M. Mrozek and J.F. Reineck.

# Synchronizations in lattices of nonlinear 

Duffing's oscillators<br>by<br>V.S. Afraimovich and S.N. Chow

CDSNS95-208

# Synchronizations in Lattices of Nonlinear Duffing's Oscillators 

V. S. Afraimovich ${ }^{1}$ and S.-N. Chow ${ }^{2}$


#### Abstract

In this work we prove the possibility of stochastic synchronization in two dimensional lattice of coupled Duffing's oscillators with external periodic forces. The synchronization occurs provided coupling is dissipative and the coefficients of coupling is greater than some critical values. These values depend on parameters of individual subsystems and on the size of lattice.


## 1. Introduction

The phenomenon of stochastic synchronization was first observed in [1] for identical coupled subsystems and in [2] for different systems. In [2] a rigorous definition of synchronization was introduced and a theory of synchronization of parametrically excited, diffusively coupled nonlinear oscillators was suggested. Later, this phenomenon was observed in different fields - electrical engineering, biology, lasers systems, etc. (see, for example, [3-11]). Roughly speaking, synchronization occurs when two or more different systems exhibit similar behavior provided some dissipative coupling exists between them and a coefficient of this coupling is greater than some critical value.

A mathematical foundation of synchronization of two slightly different subsystems was first introduced in [9]. It seems that there are no rigorous mathematical results related to this except for $[1,2,9,10]$, even though this phenomenon is important in applications and a mathematically interesting topic.

Our work deals with synchronization in large lattices. For this case, some features of continuous nonequilibrium media become evident: for example, boundary conditions play an essential role. Of course, it is not surprising that a large lattice of coupled oscillators is a model of nonequilibrium media and must reflect its features. But mathematical difficulties increase in comparison to two coupled subsystems. Therefore, we restrict ourself in this

[^3]paper to the case of specific subsystems - Duffing's oscillator with periodic forces. It is well-known that each of them may manifest chaotic behavior for a wide range of parameters.

## 2. Formulation of the Problem

Let us consider the following system

$$
\left\{\begin{array}{l}
\dot{x}_{i}=y_{i}  \tag{1}\\
\dot{y}_{i}=-k_{i} y_{i}+\alpha_{i} x_{i}-\beta_{i} x_{i}^{3}+a_{i}+b_{i} \cos \omega t+c_{1} \Delta y_{i}+c_{2} \Delta x_{i}
\end{array}\right.
$$

where $i$ is a $p$-dimensional integer vector, i.e., $i \in \mathbb{Z}^{p}, k_{i}>0, \alpha_{i}>0, \beta_{i}>0, a_{i}, b_{i}, c_{1}>0$ and $c_{2}>0$ are constants. In this paper, for the sake of definiteness and in order to avoid cumbersome calculations and formulas, we consider only the case $p=2$. But our technique works and all results hold for an arbitrary value of $p \geq 2$. The symbol $\Delta$ denotes a discrete version of the Laplace operator:

$$
\Delta \xi_{i}=\xi_{i_{1}+1, i_{2}}+\xi_{i_{1}, i_{2}+1}+\xi_{i_{1}-1, i_{2}}+\xi_{i_{1}, i_{2}-1}-4 \xi_{i_{1}, i_{2}}, \quad i=\left(i_{1}, i_{2}\right) .
$$

Our assumptions are:
(i) Let $N>0$ be a fixed and large integer. We consider an $N \times N$-lattice, $1 \leq i_{1} \leq N$, $1 \leq i_{2} \leq N$, with the Dirichlet boundary conditions: $x_{i} \equiv 0, y_{i} \equiv 0$ if $i=\left(i_{1}, i_{2}\right)$ satisfies at least one of the following inequalities:

$$
i_{1} \leq 0, \quad i_{2} \leq 0, \quad i_{1} \geq N+1, \quad i_{2} \geq N+1
$$

(ii) The subsystems are almost identical, i.e.,

$$
\begin{gathered}
\left|k_{i}-k_{j}\right| \leq \epsilon \quad\left|\alpha_{i}-\alpha_{j}\right| \leq \epsilon, \quad\left|\beta_{i}-\beta_{j}\right| \leq \epsilon \\
\left|a_{i}-a_{j}\right| \leq \epsilon, \quad\left|b_{i}-b_{j}\right| \leq \epsilon
\end{gathered}
$$

where $1 \leq i_{1} \leq N, 1 \leq i_{2} \leq N$; the parameter $\epsilon$ characterizes scattering of the oscillators and is small; the parameters $c_{1}$ and $c_{2}$ are responsible for the magnitude of the coupling and play crucial roles in our study.

It is possible to show (see Section 5) that system (1) is dissipative in the following sense. Let $D>0$ be sufficiently large and $0<\omega<k_{i}$. Let

$$
U(x, y)=\sum_{i}\left(\frac{y_{i}^{2}}{2}-\frac{\alpha_{i} x_{i}^{2}}{2}+\frac{\beta_{i} x_{i}^{4}}{y}+\omega x_{i} y_{i}\right) .
$$

Define

$$
S=\{(x, y) \mid U(x, y)=D\}
$$

Then

$$
\left.\frac{d U}{d t}\right|_{S}=\left.(\operatorname{grad} U, Z)\right|_{S}<0
$$

where $Z$ is the vector field defined by the right hand side of (1). Thus, there exists a global attractor of (1) which belongs to a bounded region independent of $c_{1}$ and $c_{2}$. We note that there are many results related to existence of global attractors for infinite or finite dimensional dissipative systems (see, for example, [20], [21]). In our case, the proof of dissipativeness is nontrivial.

For any solution $x_{i}=x_{i}(t), y_{i}=y_{i}(t)$ belonging to the global attractor the following inequalities hold

$$
\begin{equation*}
\left|x_{i}(t)\right| \leq M, \quad\left|y_{i}(t)\right| \leq M \tag{2}
\end{equation*}
$$

where $M$ is a constant independent of $c_{1}, c_{2}$. Fix an arbitrary solution $\left\{x=x_{i}(t), y=y_{i}(t)\right\}$ belonging to the global attractor. Set

$$
u_{i j}(t)=x_{i}(t)-x_{j}(t), \quad v_{i j}(t)=y_{i}(t)-y_{j}(t)
$$

The main result of our paper is the following.
Theorem. Let $\underline{k}=\min _{i} k_{i}$. There exist constants $c_{1}^{*}>0, c_{2}^{*}>0$ and $0<q<1$ such that if

$$
c_{2}>c_{2}^{*}, \quad c_{1}>c_{1}^{*} \quad c_{1}>\frac{c_{2}}{q \underline{k}}
$$

then

$$
\varlimsup_{t \rightarrow \infty}\left|u_{i j}(t)\right| \leq K \epsilon, \quad \varlimsup_{t \rightarrow \infty}\left|v_{i j}(t)\right| \leq K \epsilon
$$

where $K$ is some constant independent of $\epsilon, c_{1}$ and $c_{2}$.
It follows from the Theorem that the projections of solutions on the coordinate subspaces are close to each other for an infinite interval of time. That is, the synchronization occurs.

## 3. Proof of Theorem

Introduce the mean values of the parameters:

$$
k=\frac{1}{N^{2}} \sum_{i} k_{i}, \quad \alpha=\frac{1}{N_{2}} \sum_{i} \alpha_{i}, \quad \beta=\frac{1}{N^{2}} \sum_{i} \beta_{i}, \quad a=\frac{1}{N^{2}} \sum_{i} a_{i}, \quad b=\frac{1}{N^{2}} \sum_{i} b_{i}
$$

and rewrite the system (1) in the following form

$$
\left\{\begin{array}{l}
\dot{x}_{i}=y_{i}  \tag{3}\\
\dot{y}_{i}=-k y_{i}+\alpha x_{i}-\beta x_{i}^{3}+a+b \cos \omega t+c_{1} \Delta y_{i}+c_{2} \Delta x_{i}+\delta_{i}
\end{array}\right.
$$

where

$$
\delta_{i}=-\left(k_{i}-k\right) y_{i}(t)+\left(\alpha_{i}-\alpha\right) x_{i}(t)-\left(\beta_{i}-\beta\right) x_{i}^{3}(t)+\left(a_{i}-a\right)+\left(b_{i}-b\right) \cos \omega t
$$

Taking (2) into account, we have $\left|\delta_{i}\right| \leq K_{1} \cdot \epsilon$, for some constant $K_{1}>0$.
The system for $u_{i j}(t), v_{i j}(t)$ takes the following form:

$$
\left\{\begin{array}{l}
\dot{u}_{i j}=v_{i j}  \tag{4}\\
\dot{v}_{i j}=-k v_{i j}+\alpha u_{i j}-\beta \gamma_{i j}(t) u_{i j}+c_{1} \mathcal{L} v_{i j}+c_{2} \mathcal{L} u_{i j}+\epsilon_{i j}(t)
\end{array}\right.
$$

where

$$
\begin{gathered}
\gamma_{i j}(t)=x_{i}^{2}(t)-x_{i} x_{j}+x_{j}^{2}(t), \quad \epsilon_{i j}(t)=\delta_{i}(t)-\delta_{j}(t) \\
\mathcal{L} v_{i j}=\Delta y_{i}-\Delta y_{j}, \quad \mathcal{L} u_{i j}=\Delta x_{i}-\Delta x_{j}
\end{gathered}
$$

Note that

$$
\begin{aligned}
\mathcal{L} u_{i j}= & \Delta x_{i}-\Delta x_{j} \\
= & \left(x_{i_{1}+1, i_{2}}-x_{j_{1}+1, j_{2}}\right)+\left(x_{i_{1}-1, i_{2}}-x_{j_{1}-1, j_{2}}\right)+\left(x_{i_{1}, i_{2}+1}-x_{j_{1}, j_{2}+1}\right) \\
& \quad+\left(x_{i_{1}, i_{2}-1}-x_{j_{1}, j_{2}-1}\right)-4\left(x_{i_{1}, i_{2}}-x_{j_{1}, j_{2}}\right) \\
= & u_{i_{1}+1, i_{2}, j_{1}+1, j_{2}}+u_{i_{1}-1, i_{2}, j_{1}-1, j_{2}}+u_{i_{1}, i_{2}+1, j_{1}, j_{2}+1}+u_{i_{1}, i_{2}-1, j_{1}, j_{2}-1}-4 u_{i_{1}, i_{2}, j_{1}, j_{2}} .
\end{aligned}
$$

Denoting $\left\{u_{i j}\right\}$ by $u,\left\{v_{i j}\right\}$ by $v,\left\{\gamma_{i j}\right\}$ by $\gamma,\left\{\epsilon_{i j}\right\}$ by $\underline{\epsilon}$, we can write

$$
\left\{\begin{array}{l}
\dot{u}=v  \tag{5}\\
\dot{v}=-k v+\alpha u-\beta \gamma * u+c_{1} \mathcal{L} v+c_{2} \mathcal{L} u+\underline{\epsilon}
\end{array}\right.
$$

where $(\gamma * u)_{i j}=\gamma_{i j} \cdot u_{i j}$.
We consider the system (5) for the values of $(i, j)$ belonging to the index set $J=$ $\bigcup_{\left|k_{2}\right| \leq N,\left|k_{2}\right| \leq N} J_{k_{1}} \times J_{k_{2}}$ such that the set $J_{k_{\alpha}}, \alpha=1,2$ is a discrete interval. $J_{k_{\alpha}}=\left\{\left(i_{\alpha}, j_{\alpha}\right) \mid\right.$ $j_{\alpha}=i_{\alpha}+k_{\alpha}$ where $-k_{\alpha} \leq i_{\alpha} \leq N$ if $k_{\alpha} \geq 0$ and $0 \leq i_{\alpha} \leq N-k_{\alpha}$ if $\left.k_{\alpha} \leq 0\right\}$. It is simple to see that

$$
d^{2} \stackrel{\text { def }}{=} \operatorname{card}\left(\bigcup_{\left|k_{\alpha}\right| \leq N} J_{k_{\alpha}}\right)=(N+1)+2[(N+2)+(2 N+1)] \cdot \frac{N}{2}=(N+1)(3 N+1)
$$

Therefore, $\left\{u_{i j}\right\}$ and $\left\{v_{i j}\right\}$ belong to $\mathbb{R}^{d^{2}} \times \mathbb{R}^{d^{2}}$, and $u_{i j} \equiv 0, v_{i j} \equiv 0$ if $(i, j) \notin J$. Let us remark that $\epsilon_{i j} \equiv 0$ if $(i, j) \notin J$.

We prove below that $u_{i j}$ is small for any $(i, j) \in J$. In particular, it will be so for $i$ and $j$ belonging to the original $N \times N$ lattice. We will essentially use below the following result.

Main Lemma. (i) The operator $\mathcal{L}: \mathbb{R}^{d^{2}} \times \mathbb{R}^{d^{2}} \rightarrow \mathbb{R}^{d^{2}} \times \mathbb{R}^{d^{2}}$ is self-adjoint. (ii) Let $\sigma(\mathcal{L})$ be the spectrum of $\mathcal{L}$ and

$$
\lambda_{0}=-8 \sin ^{2} \frac{\pi}{2(2 N+2)} .
$$

If $\lambda_{s} \in \sigma(\mathcal{L})$, then $\lambda_{s} \leq \lambda_{0}$.
The main Lemma will be proved in Section 4. The main Theorem will be proved by using several changes of variables. The first change of variables is linear.

Proposition 1. We assume that

$$
\begin{equation*}
\alpha-8 c_{2} \sin ^{2} \frac{\pi}{2(2 N+2)}<0 \tag{6}
\end{equation*}
$$

There exists a nonsingular self-adjoint linear map $A: \mathbb{R}^{d^{2}} \times \mathbb{R}^{d^{2}} \rightarrow \mathbb{R}^{d^{2}} \times \mathbb{R}^{d^{2}}$ such that the following change of variables

$$
\begin{equation*}
v=w+A u \tag{7}
\end{equation*}
$$

applied to the system (5) gives the following

$$
\left\{\begin{array}{l}
\dot{w}=-k w-A w+\beta \gamma * u+c_{1} \mathcal{L} w+\underline{\epsilon}  \tag{8}\\
\dot{u}=w+A u
\end{array} .\right.
$$

Furthermore, for any $\rho_{s} \in \sigma(A)$ there exists $\lambda_{s} \in \sigma(\mathcal{L})$ such that

$$
\begin{equation*}
\rho_{s}=\frac{1}{2}\left(-k+c_{1} \lambda_{s}+\sqrt{\left(-k+c_{1} \lambda_{s}\right)^{2}+4\left(\alpha+c_{2} \lambda_{s}\right)}\right)<0 . \tag{9}
\end{equation*}
$$

Proof. Differentiating (7), we obtain $\dot{v}=\dot{w}+A \dot{u}$, and it follows from (5) and (7) that

$$
\dot{w}+A w+A^{2} u=-k w-k A u+\alpha u-\beta \gamma * u+c_{1} \mathcal{L} w+c_{1} \mathcal{L} A u+c_{2} \mathcal{L} u+\underline{\epsilon}
$$

Claim that

$$
\begin{equation*}
A^{2} u=-k A u+\alpha u+c_{1} \mathcal{L} A u+c_{2} \mathcal{L} u \tag{10}
\end{equation*}
$$

If (10) is satisfied, then (8) will also be satisfied. We will show that there is an operator $A$ satisfying (10). Let $h^{s}$ be an eigenvector corresponding to the eigenvalue $\lambda_{s}$ of the operator $\mathcal{L}$. Hence, $\mathcal{L} h^{s}=\lambda_{s} h^{s}$, and the set $\left\{h^{s}\right\}$ forms a system of coordinates in $\mathbb{R}^{d^{2}} \times \mathbb{R}^{d^{2}}$. Each element $u$ can be represented as $u=. \sum_{s} u^{s} h^{s}$. Define $A$ as follows

$$
\begin{equation*}
A u=\sum_{s} u_{s} \rho_{s} h^{s} \tag{11}
\end{equation*}
$$

where $\rho_{s}$ is as in (9). Then

$$
\begin{align*}
& A^{2} u=\sum_{s} u_{s} \rho_{s}^{2} h^{s}, \quad-k A u=\sum_{s}-k \rho_{s} u_{s} h^{s}, \quad \alpha u=\sum_{s} \alpha u_{s} h^{s}, \\
& c_{2} \mathcal{L} u=\sum_{s} c_{2} \lambda_{s} u_{s} h^{s}, \quad c_{2} \mathcal{L} A u=\sum_{s} c_{1} \lambda_{s} \rho_{s} u_{s} h^{s} \tag{12}
\end{align*}
$$

Note that $\rho_{s}^{2}=-k \rho_{s}+\alpha+c_{2} \lambda_{s}+c_{1} \lambda_{s} \rho_{s}$. Substituting (12) into (10) we obtain the desired result. It is simple to check that $A$ is self-adjoint.

Now, we work with the system (8).
The second change of variables. Note that $\|\epsilon\|<K_{2} \cdot \epsilon$ where $K_{2}>0$ is a constant. If we show that

$$
\limsup _{t \rightarrow \infty}\|w(t)\|=O(\epsilon)
$$

then it will imply the desired result. We are looking for a change of variables in the form

$$
\begin{equation*}
w_{i j}=\eta_{i j}+g_{i j}(t) \cdot u_{i j} \tag{13}
\end{equation*}
$$

where $\eta=\left\{\eta_{i j}\right\}$ are new variables and $g(t)=\left\{g_{i j}(t)\right\}$ is some function to be determined. Differentiating (13),

$$
\dot{w}_{i j}=\dot{\eta}_{i j}+\dot{g}_{i j} u_{i j}+g_{i j} \cdot \dot{u}_{i j} .
$$

Substituting (8) and (13) into the above expression, we obtain

$$
\begin{align*}
\dot{\eta}_{i j} & +\dot{g}_{i j} u_{i j}+g_{i j}\left(\eta_{i j}+g_{i j} u_{i j}\right)+g_{i j}(A u)_{i j} \\
& =-k \eta_{i j}-k g_{i j} u_{i j}-\beta \gamma_{i j} u_{i j}+c_{1} \mathcal{L} \eta_{i j}+c_{1} \mathcal{L} g_{i j} u_{i j}+\epsilon_{i j} \tag{14}
\end{align*}
$$

If $g_{i j}(t)$ satisfies the equation

$$
\begin{equation*}
\dot{g}_{i j} u_{i j}+g_{i j}^{2} u_{i j}=-k g_{i j} u_{i j}-\beta \gamma_{i j} u_{i j}-g_{i j}(A u)_{i j}+c_{1} \mathcal{L} g_{i j} u_{i j} \tag{15}
\end{equation*}
$$

then the equation (14) becomes

$$
\begin{equation*}
\dot{\eta}_{i j}=-k \eta_{i j}-g_{i j}(t) \eta_{i j}+c_{1} \mathcal{L} \eta_{i j}+\epsilon_{i j} \tag{16}
\end{equation*}
$$

and

$$
\begin{equation*}
\dot{u}_{i j}=\eta_{i j}+g_{i j}(t) u_{i j}+(A u)_{i j} \tag{17}
\end{equation*}
$$

Solution of the equation (15). Rewrite the equation (15) in an operator form. For this, take into account that an arbitrary element $p=\left\{p_{i j}\right\}$ corresponds to the operator $(\operatorname{diag} p)$ which acts on the space $\{u\}$ as follows: $((\operatorname{diag} p)(u))_{i j}=p_{i j} u_{i j}$. So, the equation (15) can be written in the form

$$
\begin{align*}
(\operatorname{diag} \dot{g})(u) & +\left(\operatorname{diag} g^{2}\right)(u) \\
& =-k(\operatorname{diag} g)(u)-\beta(\operatorname{diag} \gamma)(u)-(\operatorname{diag} g)(A u)+c_{1} \mathcal{L}(\operatorname{diag} g)(u) \tag{18}
\end{align*}
$$

For the sake of convenience set $\bullet=(\operatorname{diag} \bullet)\left(\bullet\right.$ will be $g, \gamma, \dot{g}$, or $\left.g^{2}\right)$, then (18) becomes

$$
\begin{equation*}
\underline{\dot{g}} u+\underline{g}^{2} u=-k \underline{g} u-\beta \underline{\gamma} u-\underline{g} A u+c_{1} \mathcal{L} \underline{g} u \tag{19}
\end{equation*}
$$

The equation (19) will be satisfied if the operator Riccati equation

$$
\underline{\dot{g}}+\underline{g}^{2}=-k \underline{g}-\beta \underline{\gamma}-\underline{g} A+c_{1} \mathcal{L} \underline{g}
$$

is satisfied. Rewrite it as follows

$$
\begin{equation*}
\underline{\dot{g}}=-\underline{g}^{2}-\beta \underline{\gamma}+\mathcal{R} \underline{g} \tag{20}
\end{equation*}
$$

where $\mathcal{R}$ is a linear operator:

$$
\begin{equation*}
\mathcal{R} \underline{g}=-k \underline{g}-\underline{g} A+c_{1} \mathcal{L} \underline{g} \tag{21}
\end{equation*}
$$

Proposition 2. We have that $\mathcal{R}$ is self-adjoint and the spectrum $\sigma(\mathcal{R})=\left\{r_{s}\right\}$, where

$$
r_{s}=\frac{1}{2}\left(-k+c_{1} \lambda_{s}-\sqrt{\left(-k+c_{1} \lambda_{s}\right)^{2}+4\left(\alpha+c_{2} \lambda_{s}\right)}\right), \quad \lambda_{s} \in \sigma(\mathcal{L})
$$

Proof. By definition, the operator $\mathcal{R}$ acts on the space $\mathbb{R}^{d^{2}} \times \mathbb{R}^{d^{2}}$. We will show that the eigenvectors $\left\{h^{s}\right\}$ of $\mathcal{L}$ are also the eigenvectors of $\mathcal{R}$. By Proposition 1 and (21), we have

$$
\begin{aligned}
\mathcal{R} h^{s} & =-k h^{s}-h^{s} \rho_{s}+c_{1} \lambda_{s} h^{s}=\left(-k-\rho_{s}+c_{1} \lambda_{s}\right) h^{s} \\
& =\left(-k+c_{1} \lambda_{s}-\frac{1}{2}\left(-k+c_{1} \lambda_{s}+\sqrt{\left(-k+c_{1} \lambda_{s}\right)^{2}+4\left(\alpha+c_{2} \lambda_{s}\right)}\right) h^{s}=r_{s} h^{s} .\right.
\end{aligned}
$$

Corollary. (i)

$$
r=\max _{s} r_{s} \leq-\frac{1}{2}\left(k+8 c_{1} \sin ^{2} \frac{\pi}{2(2 N+2)}\right)<0
$$

(ii) If $S(t)$ is the fundamental matrix of the system $\dot{\xi}=\mathcal{R} \xi$ then

$$
\begin{equation*}
\left\|S\left(t-t_{0}\right)\right\| \leq e^{q r\left(t-t_{0}\right)} \tag{22}
\end{equation*}
$$

where $0<q<1$ is some constant independent of $c_{1}, c_{2}$.
We will look for a solution of (20) which is a solution of the following integral equation

$$
\begin{equation*}
\underline{g}(t)=-\int_{0}^{t} S(t-\tau)\left(\underline{g}^{2}(\tau)+\beta \underline{\gamma}(\tau)\right) d \tau \tag{23}
\end{equation*}
$$

Evidently, any continuous solution of (23), $t \in \mathbb{R}^{+}$, satisfies the equation (20). We will solve the equation (23) using the contraction mapping principle. For any $T>0$ and $m>0$, consider a metric space $\mathcal{H}_{T, m}$ of continuous vector functions

$$
\underline{h}:[0, T] \rightarrow \mathbb{R}^{d^{2}} \times \mathbb{R}^{d^{2}}
$$

satisfying

$$
\mathcal{H}_{T, m}=\{\underline{h}(\cdot) \mid \underline{h}(0)=0,\|\underline{h}(t)\| \leq m, 0 \leq t \leq T\}
$$

The space $\mathcal{H}_{T, m}$ is endowed with the $C^{0}$-metrics:

$$
\operatorname{dist}\left(\underline{h}^{\prime}(t), \underline{h}^{\prime \prime}(t)\right)=\sup _{0 \leq t \leq T}\left\|h^{\prime}(t)-h^{\prime \prime}(t)\right\| .
$$

Denote by $\Gamma$ the maximal norm: $\Gamma=\sup _{t}\|\gamma(t)\|$. By the uniform boundedness of the global attractor, $\Gamma$ is independent of $c_{1}, c_{2}$.

Proposition 3. Suppose

$$
\begin{equation*}
q^{2} r^{2}>4 \beta \Gamma \tag{24}
\end{equation*}
$$

and

$$
\begin{equation*}
m=m_{0} \equiv \frac{1}{2}\left(q|r|-\sqrt{q^{2} r^{2}-4 \beta \Gamma}\right) . \tag{25}
\end{equation*}
$$

Let $Q$ be an operator defined by:

$$
\begin{equation*}
Q \underline{h}(t)=-\int_{0}^{t} S(t-\tau)\left(\underline{h}^{2}(\tau)+\beta \underline{\gamma}(\tau)\right) d \tau \tag{26}
\end{equation*}
$$

Then the following statements hold
(i) $Q \mathcal{H}_{T, m_{0}} \subset \mathcal{H}_{T, m_{0}}$;
(ii) $Q$ is a contraction.

Proof. (i) It follows from the definition (26) that

$$
\|Q \underline{h}(t)\| \leq \int_{0}^{t} \epsilon^{q r(t-\tau)} \cdot\left(\left\|\underline{h}^{2}(\tau)\right\|+\beta \Gamma\right) d \tau
$$

By the definition of the norm $\|\cdot\|,\left\|\underline{h}^{2}(\tau)\right\| \leq\|\underline{h}(\tau)\|^{2}$. Thus,

$$
\begin{aligned}
\|Q \underline{h}(t)\| & \leq \int_{0}^{t} e^{q r(t-\tau)} \cdot\left(m_{0}^{2}+\beta \Gamma\right) d \tau \\
& =\frac{1}{q r}\left(1-e^{q r t}\right) \cdot\left(m_{0}^{2}+\beta \Gamma\right) \\
& <\frac{m_{0}^{2}+\beta \Gamma}{q|r|}
\end{aligned}
$$

If (24) and (25) are satisfied, then $\frac{m_{0}^{2}+\beta \Gamma}{q|r|}=m_{0}$ and $\|Q \underline{h}(t)\| \leq m_{0}$ for any $t \in[0, T]$.
(ii) Let $\underline{h}_{1}(t), \underline{h}_{2}(t) \in \mathcal{H}_{T, m_{0}}$ be arbitrary. By using the scalar product in $\mathbb{R}^{d^{2}} \times \mathbb{R}^{d^{2}}$ and the Cauchy-Schwarz inequality, we have

$$
\begin{aligned}
\left\|Q \underline{h}_{1}(t)-Q h_{2}(t)\right\| & \leq \int_{0}^{t} e^{q r(t-\tau)} \cdot\left\|h_{1}^{2}(\tau)-h_{2}^{2}(\tau)\right\| d \tau \\
& \leq \int_{0}^{t} e^{q r(t-\tau)} \cdot\left\|\underline{h}_{1}(\tau)-\underline{h}_{2}(\tau)\right\| \cdot 2 m_{0} d \tau \\
& \leq \frac{1}{q|r|}\left(1-e^{q r t}\right) \cdot 2 m_{0} \cdot \operatorname{dist}\left(\underline{h}_{1}(t), \underline{h}_{2}(t)\right) \\
& <\frac{2 m_{0}}{q|r|} \cdot \operatorname{dist}\left(\underline{h}_{1}(t), \underline{h}_{2}(t) \mid\right.
\end{aligned}
$$

Since

$$
\frac{2 m_{0}}{q|r|}=p \equiv \frac{q|r|-\sqrt{q^{2} r^{2}-4 \beta \Gamma}}{q|r|}<1
$$

we have

$$
\operatorname{dist}\left(Q \underline{h}_{1}(t), Q \underline{h}_{2}(t)\right) \leq p \operatorname{dist}\left(\underline{h}_{1}(t), \underline{h}_{2}(t)\right)
$$

then statement (ii) is proved.
The unique fixed point of $Q$, say $\underline{h}_{T}(t)$, is the desired solution. Evidently, if $Q \underline{h}_{T_{1}}(t)=$ $\underline{h}_{T_{\mathbf{1}}}(t), Q \underline{h}_{T_{2}}(t)=\underline{h}_{T_{2}}(t)$ and $T_{2}>T_{1}$, then $\underline{h}_{T_{1}}(t) \equiv \underline{h}_{T_{2}}(t)$ for all $0 \leq t \leq T_{1}$. Denote this fixed point by $\underline{g}(t)=\left\{g_{i j}(t)\right\}$. It satisfies the equation (15). Therefore the original system is represented in the form-(16), (17). It follows from Proposition 3 that

$$
\|\underline{g}(t)\| \leq m_{0}=\frac{2 \beta \Gamma}{q|r|+\sqrt{q^{2} r^{2}-4 \beta \Gamma}} \rightarrow 0 \quad \text { as } c_{1} \rightarrow \infty
$$

Moreover, since $g_{i j}(t)$ satisfies the equation

$$
\begin{equation*}
g_{i j}(t)=-\int_{0}^{t}\left(S(t-\tau)\left(\underline{g}^{2}(\tau)+\beta \underline{\gamma}(\ddot{\tau})\right)_{i j} d \tau\right. \tag{27}
\end{equation*}
$$

and $\mathcal{R}$ is self-adjoint, we have that $0 \leq-g_{i j}(t) \leq m_{0}$.
Solution of equations (16), (17). Denote by $T_{16}\left(t_{0}, t\right)$ the fundamental matrix of the linear system

$$
\dot{\eta}_{i j}=-k \eta_{i j}-g_{i j}(t) \eta_{i j}+c_{1} \mathcal{L} \eta_{i j}
$$

It follows from (27) then

$$
\left\|T_{16}\left(t_{0}, t\right)\right\| \leq a \cdot e^{\left(-k+m_{0}-c_{1} 8 \sin ^{2} \frac{x}{2(2 N+2)}\right)\left(t-t_{0}\right)}
$$

where $a$ is some constant. Therefore, if

$$
\begin{equation*}
k+c_{1} 8 \sin ^{2} \frac{\pi}{2(2 N+2)}>m_{0}=\frac{2 \beta \Gamma}{q|r|+\sqrt{q^{2} r^{2}-4 \beta \Gamma}} \tag{28}
\end{equation*}
$$

then the solution of (16) satisfies the inequality

$$
\begin{equation*}
\|\eta(t)\| \leq\|\eta(0)\| \cdot e^{\left(-k+m_{0}+c_{1} \lambda_{0}\right) t}+\frac{b}{\left|-k+c_{1} \lambda_{0}+m_{0}\right|} \cdot \epsilon_{1} \tag{29}
\end{equation*}
$$

where $\epsilon_{1}=\sup \|\epsilon(t)\|$ and $b$ is some constant.
Denote by $T_{17}\left(t_{0}, t\right)$ the fundamental matrix of the linear system

$$
\dot{u}_{i j}=g_{i j}(t) u_{i j}+(A u)_{i j} .
$$

By Proposition 1 and (27), we have

$$
\begin{equation*}
\left\|T_{17}\left(t, t_{0}\right)\right\| \leq a_{1} e^{q \rho_{0}\left(t-t_{0}\right)} \tag{30}
\end{equation*}
$$

where

$$
0>\rho_{0}=\max _{s} \rho_{s}=\frac{1}{2}\left(-k+c_{1} \lambda_{0}+\sqrt{\left(-k+c_{1} \lambda_{0}\right)^{2}+4\left(\alpha+c_{2} \lambda_{0}\right)}\right)
$$

and $a_{1}$ is some constant. Therefore, we obtain from (17) and (29) the following estimate

$$
\begin{align*}
\|u(t)\| \leq & e^{\rho_{0} t}\|u(0)\| \\
& +\|\eta(0)\| \cdot a_{1}\left|e^{\left(-k+m_{0}+c_{1} \lambda_{0}\right) t}-e^{q \rho_{0} t}\right| \cdot \frac{1}{\left(-q \rho_{0}+k-c \lambda_{0}-m_{0}\right)}  \tag{31}\\
& +\frac{a_{1} b}{q\left|\rho_{0} \|-k+c_{1} \lambda_{0}+m_{0}\right|} \cdot \epsilon_{1} .
\end{align*}
$$

Since $\epsilon_{1} \leq$ const $\cdot \epsilon$, (29) and (31) imply the validity of Theorem. The constant $c_{2}^{*}$ can be found from (6):

$$
\begin{equation*}
c_{2}^{*}=\frac{\alpha}{8 \sin ^{2} \frac{\pi}{2(2 N+2)}} \tag{32}
\end{equation*}
$$

and the constant $c_{1}^{*}$ can be found from (24), (28):

$$
\begin{equation*}
c_{1}^{*} \leq \frac{2 \sqrt{\beta \Gamma}-k}{8 \sin ^{2} \frac{\pi}{2(2 N+2)}} \tag{33}
\end{equation*}
$$

## 4. Spectra of Some Operators and Proof of Main Lemma

Since the operator $\mathcal{L}$ is determined by the discrete Laplace operator $\Delta$, we consider, first, the spectrum of $\Delta$ denoted by $\sigma(\Delta)$.
(i) For the case $p=1$, under Dirichlet boundary conditions the operator $\Delta$ can be represented by the following $M \times M$ matrix

$$
\left(\begin{array}{ccccc}
-2 & 1 & 0 & \cdots & \\
1 & -2 & 1 & 0 & \\
\cdots & \cdots & \cdots & & \\
0 & & 1 & -2 & 1 \\
0 & & 0 & 1 & -2
\end{array}\right)
$$

where $M$ is the number of points in the one-dimensional lattice. It is well-known (see, for example, [19]) that

$$
\sigma \quad\left\{\left.\lambda_{s}=-2-2 \cos \left(\frac{A}{M}\right)^{\pi s} \frac{\overline{=}}{+1} \right\rvert\, s=1, \ldots, M\right\} .
$$

Thus,

$$
\begin{equation*}
\max _{s} \lambda_{s}=-4 \sin ^{2} \frac{\pi}{2(M+1)} \tag{34}
\end{equation*}
$$

Furthermore, this operator is self-adjoint and its eigenvectors are mutually orthogonal in the corresponding Euclidean space.
(ii) For the case $p=2$, we will find the eigenvectors and eigenvalues of $\Delta$ by using separation of variables. This method was used in [12] for the case of periodic boundary conditions.

Let $M_{1}, M_{2}>1$ be fixed integers. We consider $\Delta$ acting on the $M_{1} \times M_{2}$ lattice. Let

$$
\Delta y_{i_{1}, i_{2}}=\lambda y_{i_{1} i_{2}}, \quad \text { where } \quad 1 \leq i_{1} \leq M_{1}, \quad 1 \leq i_{2} \leq M_{2} .
$$

Setting $y_{i_{1} i_{2}}=Y_{i_{1}} \cdot Y_{i_{2}}$, we have

$$
Y_{i_{2}} \Delta_{1} Y_{i_{1}}+Y_{i_{1}} \Delta_{1} Y_{2}=\lambda Y_{i_{1}} Y_{i_{2}}
$$

where $\Delta_{1}$ is the one-dimensional Laplace operator. If $\Delta_{1} Y_{i_{1}}=\mu_{1} Y_{i_{1}}$ and $\Delta_{1} Y_{i_{2}}=\mu_{2} Y_{i_{2}}$, then $\lambda=\mu_{s}+\mu_{t}$, where

$$
\begin{equation*}
\mu_{1}=-2-2 \cos \frac{\pi s}{M_{1}+1}, \quad \mu_{2}=-2-2 \cos \frac{\pi t}{M_{2}+1} \tag{35}
\end{equation*}
$$

Let

$$
\lambda_{s, t}=-4-2\left(\cos \frac{\pi s}{M_{1}+1}+\cos \frac{\pi t}{M_{2}+1}\right), \quad s=1, \ldots, M_{1}, \quad t=1, \ldots, M_{2}
$$

It is clear that $\Delta$ is self-adjoint for $d=2$, as well. If follows from (35) that

$$
\begin{equation*}
\max _{s, t} \lambda_{s, t} \leq-8 \sin ^{2} \frac{\pi}{2(M+1)} \tag{36}
\end{equation*}
$$

where $M=\max \left\{M_{1}, M_{2}\right\}$. Let $\mu_{s} Y^{(s)}=\Delta_{1} Y^{(s)}, s=1, \ldots, M_{1}, \mu_{t} Y^{(t)}=\Delta_{1} Y^{(t)}$, $t=1, \ldots, M_{2}$, and the eigenvectors $\left\{Y^{(s)}\right\},\left(\left\{Y^{(t)}\right\}\right)$ form an orthonormal basis in $\mathbb{R}^{M_{1}}$ $\left(\mathbb{R}^{M_{2}}\right)$. Let $Z_{i_{1} i_{2}}^{s, t}=Y_{i_{1}}^{s} \cdot Y_{i_{2}}^{t}, Z^{s, t}=\left\{Y_{i_{1} i_{2}}^{s, t}\right\}, i_{1}=1, \ldots, M_{1}, i_{2}=1, \ldots, M_{2}$. We proved that $\Delta_{2} Z^{s, t}=\lambda_{s, t} Z^{s, t}$. It remains to show that $M_{1} \cdot M_{2}$ vectors $Z^{s, t}$ form a basis in $\mathbb{R}^{M_{1}} \times \mathbb{R}^{M_{2}}$. It is simple to check that $Z^{s, t} \neq 0$. Let us check that they are linearly independent. It will be so if they are orthogonal. Show it. Fix two pairs $(s, t)$ and $\left(s_{0}, t_{0}\right)$. Then

$$
\sum_{i_{1}=1}^{M_{1}} \sum_{i_{2}=1}^{M_{2}} Z_{i_{1} i_{2}}^{s, t} \cdot Z_{i_{1} i_{2}}^{s_{0}, t_{0}}=\sum_{i_{1}=1}^{M_{1}} Y_{i_{1}}^{s} \cdot Y_{i_{1}}^{s_{0}} \sum_{i_{2}=1}^{M_{2}} Y_{i_{2}}^{t} \cdot Y_{i_{2}}^{t_{0}}=1
$$

if $s=s_{0}, t=t_{0}$ and 0 otherwise. Thus, $Z^{s, t}$ form an orthonormal basis in the Euclidean space $\mathbb{R}^{M_{1}} \times \mathbb{R}^{M_{2}}$.
(iii) Let us now look at the operator $\mathcal{L}$ :

$$
\begin{aligned}
\mathcal{L} v_{i_{1} i_{2} j_{1} j_{2}}= & v_{i_{1}-1 i_{2} j_{1}-1 j_{2}}+v_{i_{1}+1 i_{2} j_{1}+1, j_{2}}+v_{i_{1} i_{2}-1 j_{1} j_{2}-1} \\
& +v_{i_{1} i_{2}+1 j_{1} j_{2}+1}-4 v_{i_{1} i_{2} j_{1} j_{2}}, \quad i_{k}=1, \ldots, N, \quad k=1,2,3,4
\end{aligned}
$$

with the boundary conditions described above: $v_{i j} \equiv 0$ if $(i, j) \notin J$. The simple observation shows that for $\left(i_{1}, j_{1}\right) \in J_{k_{1}},\left(i_{2}, j_{2}\right) \in J_{k_{2}},-N \leq k_{1,2} \leq N$, if we introduce the notations $w_{i_{1} i_{2} k_{1} k_{2}}=v_{i_{1} i_{2} j_{1} j_{2}}$ then
$\mathcal{L} v_{i_{1} i_{2} j_{2} j_{2}}=\Delta_{2} w_{i_{1} i_{2} k_{1} k_{2}}=w_{i_{1}-1 i_{2} k_{1} k_{2}}+w_{i_{1}+1 i_{2} k_{1} k_{2}}+w_{i_{1} i_{2}+1 k_{1} k_{2}}+w_{i_{1} i_{2}-1 k_{1} k_{2}}-4 w_{i_{1} i_{2} k_{1} k_{2}}$
where
(i) $0 \leq i_{1} \leq N-k_{1}$ if $-N \leq k_{1} \leq 0$,
(ii) $-k_{1} \leq i_{1} \leq N$ if $0 \leq k_{1} \leq N$,
(iii) $0 \leq i_{2} \leq N-k_{2}$ if $-N \leq k_{2} \leq 0$,
(iv) $-k_{2} \leq i_{2} \leq N$ if $0 \leq k_{2} \leq N$.

Therefore, the eigenvalues of $\mathcal{L}$ are:

$$
\begin{aligned}
&-4-2\left(\cos \frac{s \pi}{N-k_{1}+2}+\cos \frac{t \pi}{N-k_{2}+2}\right), s=1, \ldots, N-k_{1}+1, \\
& t=1, \ldots, N-k_{2}+1, \text { in the case (i), (iii); } \\
&-4-2\left(\cos \frac{s \pi}{N-k_{1}+2}+\cos \frac{t \pi}{N+k_{2}+2}\right), s=1, \ldots, N-k_{1}+1, \\
& t=1, \ldots, N+k_{2}+1, \text { in the case (i), (iv); } \\
&-4-2\left(\cos \frac{s \pi}{N+k_{1}+2}+\cos \frac{t \pi}{N-k_{2}+2}\right), s=1, \ldots, N+k_{1}+1, \\
& t=1, \ldots, N-k_{2}+1, \text { in the case (ii), (iii); } \\
&-4-2\left(\cos \frac{s \pi}{N+k_{1}+2}+\cos \frac{t \pi}{N+k_{2}+2}\right), \quad s=1, \ldots, N+k_{1}+1, \\
& t=1, \ldots, N+k_{2}+1, \text { in the case (ii), (iv). }
\end{aligned}
$$

In the "worse" case ( $k_{1,2}= \pm N$ )

$$
\begin{equation*}
\max \{\lambda \mid \lambda \in \sigma(\mathcal{L})\}=-4\left(1+\cos \frac{(1+N) \pi}{2 N+2}\right)=-8 \sin ^{2} \frac{\pi}{2(2 N+2)} \tag{37}
\end{equation*}
$$

The support of the eigenvector, for example, corresponding to the eigenvalue

$$
-4-2\left(\cos \frac{s \pi}{N-k_{1}+2}+\cos \frac{t \pi}{N-k_{2}+2}\right)
$$

is a piece of the two dimensional discrete plane:

$$
\left\{\left(i_{1}, i_{2}, j_{1}, j_{2}\right) \mid j_{1}-i_{1}=k_{1}, j_{2}-i_{2}=k_{2}\right\}
$$

in the four-dimensional discrete set $J$. On this piece the operator $\mathcal{L}$ coincides with the Laplacian $\Delta$ with the Dirichlet boundary conditions. Thus, all eigenvectors with the support on this piece are mutually orthogonal. Two-dimensional planes corresponding to different pairs ( $k_{1}, k_{2}$ ) do not intersect. So, the eigenvectors corresponding to different pairs $\left(k_{1}, k_{2}\right)$ are mutually orthogonal. Therefore, all eigenvectors are mutually orthogonal and form a basis in $\mathbb{R}^{d^{2}} \times \mathbb{R}^{d^{2}}$. It follows from this that $\mathcal{L}$ is self-adjoint. (Though, it is possible to prove directly from the definition of $\mathcal{L}$.) This proves the Main Lemma. We denoted the spectrum $\sigma(\mathcal{L})$ by $\left\{\lambda_{s}\right\}$ and the corresponding eigenvectors by $\left\{h^{s}\right\}$.

## 5. Dissipativeness of the System

Denote by $\left\{\xi^{s}\right\}$ the eigenvectors of the discrete Laplace operator $\Delta$ for $d=2$, and by $\left\{\lambda_{s}\right\}$ the corresponding eigenvalues, see (35). It was shown in Section 4 that $\xi^{s_{1}} \perp \xi^{s_{2}}$. We assume that $\left\{\xi^{s}\right\}$ are chosen so that they form an orthonormal basis in $\mathbb{R}^{N^{2}} \times \mathbb{R}^{N^{2}}$ : $\sum_{i} \xi_{i}^{s_{1}} \xi_{i}^{s_{2}}=\delta_{s_{1} s_{2}}$. Let $x=\left\{x_{i}\right\}, y=\left\{y_{i}\right\}$ be elements of $\mathbb{R}^{N^{2}} \times \mathbb{R}^{N^{2}}$ and $\|x\|$ be the Euclidean norm, and the scalar product $(x, y)=\sum_{i} x_{i} y_{i}$. Consider the following representation of $x, y$ :

$$
x_{i}=\sum_{s} X_{s} \xi_{i}^{s}, \quad y_{i}=\sum_{s} Y_{s} \xi_{i}^{s} .
$$

Then, the following formulas hold:

$$
\begin{gather*}
\sum_{s} X_{s}^{2}=\|x\|^{2}, \quad \sum_{i} x_{i} y_{i}=\sum_{s} X_{s} Y_{s}, \quad \sum_{s} Y_{s}^{2}=\|y\|^{2}  \tag{38}\\
\sum_{i} y_{i} \Delta y_{i}=\sum_{s} \lambda_{s} Y_{s}^{2}, \quad \sum_{i} x_{i} \Delta x_{i}=\sum_{s} \lambda_{s} X_{s}^{2}  \tag{39}\\
\sum_{i} x_{i} \Delta y_{i}=\sum_{s} \lambda_{s} X_{s} Y_{s}, \quad \sum_{i} y_{i} \Delta x_{i}=\sum_{s} \lambda_{s} X_{s} Y_{s} \tag{40}
\end{gather*}
$$

Moreover, we have

$$
\begin{aligned}
\sum_{i} x_{i} y_{i} & =\sum_{i}\left(\sum_{s_{1}} X_{s_{1}} \xi_{i}^{s_{1}}\right)\left(\sum_{s_{2}} Y_{s_{2}} \xi_{i}^{s_{2}}\right)=\sum_{s_{1}, s_{2}} \sum_{i} X_{s_{1}} Y_{s_{2}} \xi_{i}^{s_{1}} \xi_{i}^{s_{2}} \\
& =\sum_{s_{1}, s_{2}} X_{s_{1}} Y_{s_{2}} \delta_{s_{1} s_{2}}=\sum_{s} X_{s} Y_{s} \\
\sum_{i} y_{i} \Delta y_{i} & =\sum_{i}\left(\sum_{s_{1}} Y_{s_{1}} \xi_{i}^{s_{1}}\right) \sum_{s_{2}} Y_{s_{2}} \lambda_{s_{2}} \xi_{i}^{s_{2}}=\sum_{s_{1}, s_{2}} \lambda_{s_{2}} \sum_{i} Y_{s_{1}} \xi_{i}^{s_{1}} Y_{s_{2}} \xi_{i}^{s_{2}} \\
& =\sum_{s_{1}, s_{2}} \lambda_{s_{2}} Y_{s_{1}} Y_{s_{2}} \delta_{s_{1} s_{2}}=\sum_{s} \lambda_{s} Y_{s}^{2} .
\end{aligned}
$$

Assume that the conditions of the Theorem are satisfied, i.e.,

$$
c_{1}>\frac{c_{2}}{q \underline{k}} .
$$

Thus,

$$
\begin{equation*}
c_{2}-\omega c_{1}=0 \tag{41}
\end{equation*}
$$

for some positive $\omega$ such that

$$
\begin{equation*}
0<\omega<q \underline{k}<k_{i} . \tag{42}
\end{equation*}
$$

Consider now the function

$$
U(x, y)=\sum_{i} \frac{y_{i}^{2}}{2}-\frac{\alpha_{i} x_{i}^{2}}{2}+\frac{\beta_{i} x_{i}^{4}}{4}+\omega x_{i} y_{i}
$$

We will show that $\frac{d U}{d t}<0$ on the surface $U(x, y)=D$ if $D$ is large enough. A direct calculation shows that

$$
\begin{align*}
& \frac{d U}{d t}=\sum_{i}\left(\omega-k_{i}\right) y_{i}^{2}+\omega \alpha_{i} x_{i}^{2}-\omega \beta_{i} x_{i}^{4}-\omega k_{i} x_{i} y_{i}+\left(y_{i}+\omega x_{i}\right)\left(a_{i}+b_{i} \cos \omega t\right)  \tag{43}\\
& \quad+\left(y_{i}+\omega x_{i}\right)\left(c_{1} \Delta y_{i}+c_{2} \Delta x_{i}\right)
\end{align*}
$$

It follows from (39) and (40) that

$$
\begin{equation*}
\sum_{i}\left(y_{i}+\omega x_{i}\right)\left(c_{1} \Delta y_{i}+c_{2} \Delta x_{i}\right)=\sum_{s} \lambda_{s}\left(c_{1} Y_{s}^{2}+\omega c_{2} X_{s}^{2}+\left(c_{2}+c_{1} \omega\right) X_{s} Y_{s}\right) \tag{44}
\end{equation*}
$$

By (41), the right hand part of (44) can be written as

$$
\sum_{s} \lambda_{s} c_{1}\left(Y_{s}+\frac{c_{2}+c_{1} \omega}{2 c_{1}} X_{s}\right)^{2}=\sum_{s} \lambda_{s} c_{1}\left(Y_{s}+\omega X_{s}\right)^{2} \leq 0
$$

Therefore, it follows from (43) that

$$
\begin{aligned}
\frac{d U}{d t} & \leq \sum_{i}\left(\omega-k_{i}\right) y_{i}^{2}+\omega \alpha_{i} x_{i}^{2}-\omega \beta_{i} x_{i}^{4}-\omega k_{i} x_{i} y_{i}+\left(y_{i}+\omega x_{i}\right)\left(a_{i}+b_{i} \cos \omega t\right) \\
& \leq-\sum_{i}\left(k_{i}-\omega\right) y_{i}^{2}+\omega \beta_{i} x_{i}^{4}-\omega \alpha_{i} x_{i}^{2}+\omega k_{i} x_{i} y_{i}-\left(\left|a_{i}\right|+\left|b_{i}\right|\right)\left(\left|y_{i}\right|+\omega\left|x_{i}\right|\right) \\
& =V(x, y), \quad \text { say. }
\end{aligned}
$$

We need to show that $V(x, y)<0$ provided $U(x, y)=D, D \gg 1$. Set $\mu=D^{-1 / 4}$, $\eta_{i}=y_{i} D^{-1 / 2}, \xi_{i}=x_{i} D^{-1 / 4}$. Then the equality $U(x, y)=D$ can be rewritten as

$$
\begin{equation*}
\sum_{i} \frac{1}{2} \eta_{i}^{2}+\frac{1}{4} \beta_{i} \xi_{i}^{4}-\frac{1}{2} \mu^{2} \alpha_{i} \xi_{i}^{2}+\mu \omega \xi_{i} \eta_{i}=1 \tag{45}
\end{equation*}
$$

Proposition 4. There exist $\mu_{0}>0$ and $C>0$ such that for any $0<\mu \leq \mu_{0}$ and any solution $(\xi, \eta)$ of (45), we have

$$
\|\xi\| \leq C, \quad\|\eta\| \leq C
$$

Proof. Suppose not. Then there exist sequences $\left\{\mu_{n}\right\},\left\{\xi^{(n)}\right\},\left\{\eta^{(n)}\right\}$ such that $\mu_{n} \rightarrow 0$ and $\left\|\xi^{(n)}\right\|+\left\|\eta^{(n)}\right\| \rightarrow \infty$ as $n \rightarrow \infty$ and $\left(\xi^{(n)}, \eta^{(n)}\right)$ is a solution of (45) for $\mu=\mu_{n}$.

One of the following holds:
(i) $\left\|\eta^{(n)}\right\| \leq$ constant $\stackrel{\text { def }}{=} C$ for all $n$;
(ii) $\left\|\eta^{(n)}\right\| \rightarrow \infty$ as $n \rightarrow \infty$.

Suppose (i) is true. Then, $\left\|\xi^{(n)}\right\| \rightarrow \infty$ as $n \rightarrow \infty$. Let $p(n)=\max _{i}\left|\xi_{i}^{(n)}\right|, \underline{\beta}=\min _{i} \beta_{i}$, $\bar{\alpha}=\max _{i} \alpha_{i}$. Thus, $p(n) \rightarrow \infty$ as $n \rightarrow \infty$. Represent (45) as

$$
\begin{equation*}
\sum_{i} \frac{1}{4} \beta_{i} \xi_{i}^{(n) 4}+\sum_{i} \frac{1}{2} \eta_{i}^{(n) 2}=1+\sum_{i} \frac{1}{2} \mu_{n}^{2} \alpha_{i} \xi_{i}^{(n) 2}+\sum_{i} \mu_{n} \omega \xi_{i}^{(n)} \eta_{i}^{(n)} \tag{46}
\end{equation*}
$$

It follows from (46) that

$$
\begin{equation*}
\frac{1}{4} \underline{\beta} p^{4}(n)+\frac{1}{2}\left\|\eta^{(n)}\right\|^{2}<1+\frac{1}{2} \mu_{n}^{2} \bar{\alpha}\left\|\xi^{(n)}\right\|^{2}+\mu_{n} \omega\left\|\xi^{(n)}\right\| \cdot\left\|\eta^{(n)}\right\| \tag{47}
\end{equation*}
$$

Dividing (47) by $p^{4}(n)$ and taking into account the inequality $\left\|\xi^{(n)}\right\| \leq N^{2} p(n)$, we obtain

$$
\begin{equation*}
\frac{1}{4} \underline{\beta}+\frac{1}{2 p^{4}(n)}\left\|\eta^{(n)}\right\|^{2}<\frac{1}{p^{4}(n)}+\frac{N^{4}}{2 p^{2}(n)} \mu_{n}^{2} \bar{\alpha}+\frac{\mu_{n} N^{2}}{p^{3}(n)} \cdot C \tag{48}
\end{equation*}
$$

The inequality (48) for sufficiently large $n$ contradicts (i).
Consider now the case (ii). Dividing (46) by $\left\|\eta^{(n)}\right\|^{2}$ and using notations above, we have

$$
\begin{equation*}
\frac{1}{2}+\frac{1}{4} \underline{\beta} \frac{p^{4}(n)}{\left\|\eta^{(n)}\right\|^{2}}<\frac{1}{\left\|\eta^{(n)}\right\|^{2}}+\frac{1}{2} \mu_{n}^{2} \bar{\alpha} \cdot \frac{N^{4} p^{2}(n)}{\left\|\eta^{(n)}\right\|^{2}}+\mu_{n} \omega \frac{N^{2} p(n)}{\left\|\eta^{(n)}\right\|} \tag{49}
\end{equation*}
$$

If $p(n)$ is a bounded sequence, then (49) cannot be satisfied for $n$ large enough, and we have a contradiction. Assume now that $p(n) \rightarrow \infty$ as $n \rightarrow \infty$. Then denoting by $\zeta_{n}$ the ratio $\frac{p(n)}{\left\|\eta^{(n)}\right\|}$, we obtain from (49)

$$
\begin{equation*}
\frac{1}{2}-\frac{1}{\left\|\eta^{(n)}\right\|^{2}}-\mu_{n} \omega N^{2} \zeta_{n}+\left(\frac{1}{4} \underline{\beta} p^{2}(n)-\frac{1}{2} \mu_{n}^{2} \bar{\alpha} N^{4}\right) \zeta_{n}^{2}<0 \tag{50}
\end{equation*}
$$

If $n$ is large, then the discriminant of the quadratic polynomial (with respect to $\zeta_{n}$ ) on the left hand side of (50) is negative and the leading coefficients are positive. Thus, (50) cannot be satisfied.

Corollary. There exists $\mu_{1}, 0<\mu_{1} \leq \mu_{0}$ and a constant $c_{0}>0$ such that for any $0<\mu \leq \mu_{1}$ and any solution ( $\xi, \eta$ ) of the equation (45) the following inequality holds

$$
\begin{equation*}
\sum_{i=1}^{\infty} \eta_{i}^{2}+\xi_{i}^{4}>c_{0} \tag{51}
\end{equation*}
$$

Proof. Let $\widetilde{\beta}=\min \left\{\frac{1}{2}, \min _{i} \frac{1}{4} \beta_{i}\right\}$. Then it follows from (45) and Proposition 4 that for $\mu<\mu_{0}$

$$
\sum_{i} \eta_{i}^{2}+\xi_{i}^{4} \geq \frac{1}{\widetilde{\beta}}\left(1-\mu \omega C^{2}\right)
$$

Setting $\mu_{1} \leq 1 / 2 \omega C^{2}$ we obtain the desired result with $c_{0} \geq 1 / 2 \widetilde{\beta}$.
It remains to show that $V(x, y)<0$. We have,

$$
\begin{aligned}
V(x, y) \leq- & {\left[(\underline{k}-\omega) \sum_{i} y_{i}^{2}+\omega \underline{\beta} \sum_{i} x_{i}^{4}-\omega \bar{\alpha}\|x\|^{2}\right.} \\
& \left.+\omega \sum_{i} k_{i} x_{i} y_{i}-\left|a_{i}\right|-\left|b_{i}\right|\left(\left|y_{i}\right|+\omega\left|x_{i}\right|\right)\right] \\
=- & D\left[(\underline{k}-\omega)\|\eta\|^{2}+\omega \underline{\beta} \sum_{i} \xi_{i}^{4}+\mu F(\xi, \eta, \mu)\right]
\end{aligned}
$$

where $F$ is a bounded function (as $\mu \searrow 0$ ). Therefore

$$
V(x, y) \leq-D\left[\min \{(\underline{k}-\omega), \omega \underline{\beta}\} c_{0}+\mu F(\xi, \eta, \mu)\right]<0 .
$$

So, we prove the dissipativeness of the system and the inequality (2).

## 6. Concluding Remarks

1. Our method works for any nonlinearity $f(x)$ satisfying the following condition

$$
f\left(x_{1}\right)-f\left(x_{2}\right)=\left(x_{1}-x_{2}\right) g\left(x_{1}, x_{2}\right)
$$

where $g\left(x_{1}, x_{2}\right)>0$ for $x_{1} \neq x_{2}$.
2. In the work of T. L. Carrol and L. M. Pecora (see [3, 4, 11] and bibliography in [11]), one approach for stochastic synchronization was suggested. It is based on the
following two points: (i) in the simplest case, all subsystems are identical and coupling admits the existence of spatially homogeneous solutions; (ii) stability of solutions is based on the "transversal" Lyapunov exponent technique (see also [16, 17]). We note that for (i), the synchronization phenomenon becomes the stability property of spatially-homogeneous solutions. For non-identical subsystems, synchronization means stability of some invariant submanifolds. In case (ii), one can give only "local" synchronization. In general, synchronization means that we have the realization of a synchronous regime for a wide region of initial conditions. It seems that the approach of these authors is probably only an interstitial step in the study of synchronization phenomena.
3. As a criterion of synchronization, it is suggested in [11] to find negative transversal Lyapunov exponents. But it is known from the theory of stability in ordinary differential equations (see, for example, [22]) that an invariant set may be destroyed under small perturbations if negative Lyapunov exponents in transversal directions are larger than negative Lyapunov exponents in the tangent directions of the invariant set. This indicates that "structurally stable" synchronization comes when the transversal Lyapunov exponents. turn out to be smaller than the tangent ones.
4. It is shown in [13-18] and other works that chains (or lattices) of coupled oscillators may exhibit a complex spatio-temporal behavior. In particular, they may have attractors of different dimensions. It seems to us, the dimension of attractors depends directly on the number of excited modes. We hope to study this elsewhere.

## References

1. H. Fujisaka, T. Yamada, Stability theory of synchronized motion in coupled oscillator systems, Progress of Theoretical Physics 69 No. 1 (1983) 32-47.
2. V. S. Afraimovich, N. N. Verichev, M. I., Rabinovich, Stochastic synchronization of oscillations in dissipative systems, Izvestigya Vysshikh Uchebnych Zavedenii, Radiofizika 29 No. 9 (1986) 1050-1060 [Sov. Radiophys. 29 (1986) 795].
3. T. L. Carrol and L. M. Pecora, Synchronization nonautonomous chaotic circuits, IEEE Trans. Circuits and Systems 38 (1991) 453-456.
4. L. M. Pecora, T. L. Carrol, Synchronization in chaotic systems, Phys. Rev. Letters 64 (1990) 821-824.
5. L. O. Chua, M. Itoh, L. Kosarev and K. Eckert, Chaos synchronization in Chua's circuits, J. of Circuits, Systems and Computers 3 No. 1 (1993) 93-108.
6. L. Fabiny, P. Colet and R. Roy, Coherence and phase dynamics of spatially coupled solid-state lasers, Phys. Rev. A 47 No. 5 (1993).
7. N. F. Rul'kov, A. R. Volkovsky, Threshold synchronization of chaotic relaxation oscillations, Phys. Letters A 179 No. 9 (1993) 332-336.
8. R. Brown, N. Rul'kov, N. Tufillaro, Sychronization in chaotic systems: The effects of additive noise and drift in the dynamics of the driving, Phys. Rev. E 50 No. 3 (1994) 4488-4508.
9. H. M. Rodrigues, Uniform ultimate boundedness and synchronization, Preprint CDSNS 94-160, 1994.
10. V. Afraimovich, H. Rodrigues, Uniform ultimate boundedness and synchronization for nonautonomous equations, Preprint CDSNS 94-202, 1994.
11. J. F. Heagy, T. L. Carrol and L. M. Pecora, Synchronous chaos in coupled oscillator systems, Phys. Rev. E 50, No. 3 .(1994) 1874-1885.
12. Y. Yan, Dimensions of attractors for discretizations for Navier-Stokes equations, Preprint AHPCRC 91-01, 1991.
13. D. K. Umberger, C. G. Grebogi, E. Ott, B. Afeyan, Spatiotemporal dynamics in a dispersively coupled chain of nonlinear oscillators, Phys. Rev. A. 39 No. 9 (1989) 4835-4842.
14. D. G. Aranson, M. Golubitsky, J. Mallet-Paret, Ponies on a merry-go-round in large arrays of Josephson junctions, Nonlinearity 4 (1991) 903.
15. S. H. Strogatz, C. M. Marcus, R. M. Westervelt, R. E. Mirollo, Collective dynamics of coupled oscillators with random pinning, Physica D $\mathbf{3 6}$ (1989) 23-50.
16. J. C. Alexander, I. Kan, J. A. Yorke, Z. You, Riddled basins, Int. J. Bif. Chaos 2 (1992) 795-813.
17. E. Ott, J. C. Sommerer, J. C. Alexander, I. Kan, J. A. Yorke, Scaling behavior of chaotic systems with riddled basins, Phys. Rev. Letters 71 (1993) 4134-4137.
18. V. S. Afraimovich, I. S. Aranson, M. I. Rabinovich, Multidimensional strange attractors and turbulence, Sov. Sci. Rev. C. Math./Phys. (ed by S. P. Novikov) Harwood Academic Publishers, vol. 8 (1989).
19. V. S. Afraimovich, On the Lyapunov dimension of invariant sets in a model of an active medium, Selecta Mathematica Sovietica 10, No. 3 (1991) 91-98.
20. J. K. Hale, Asymptotic behavior of dissipative systems, Math. Surveys and Monographs 25, AMS 1988.
21. A. V. Babin, M. I. Vishik, Attractors of Evolution Equations, North-Holland, 1991.
22. N. Fenichel, Persistence and smoothness of invariant manifolds for flows. Indiana Univ. Math. J. 21 (1971), 193-226.

## CENTER FOR DYNAMICAL SYSTEMS AND NONLINEAR STUDIES

## Report Series

September 1994-

## CDSNS94-

157 Rates of eigenvalues on a dumbbell domain. Simple eigenvalue case, J.M. Arrieta.
158 Existence of standing waves for competition-diffusion equations, Y. Kan-On.
159 Fisher-type property of travelling waves for competition-diffusion equations, Y. Kan-On.

160 Uniform ultimate boundedness and synchronization, H.M. Rodrigues.
161 Complete families of pseudotrajectories and shape of attractors, S.Yu. Pilyugin.
162 Aubry-Mather Theorem and quasiperiodic orbits for time dependent reversible systems, S.-N. Chow and M.-L. Pei.

163 Slowly-migrating transition layers for the discrete Allen-Cahn and Cahn-Hilliard equations, C.P. Grant and E.S. Van Vleck.

164 Partially hyperbolic sets from a co-dimension one bifurcation, T. Young.
165 Asymptotically autonomous semiflows: Chain recurrence and Lyapunov functions, K. Mischaikow, H. Smith and H.R. Thieme.

166 On the problem of stability in lattice dynamical systems, L.A. Bunimovich and E.A. Carlen.

167 Remark on continuous dependence of attractors on the shape of domain, A.V. Babin and S.Yu. Pilyugin.

168 Special pseudotrajectories for lattice dynamical systems, V.S. Afraimovich and S.Yu. Pilyugin.

169 Proof and generalization of Kaplan-Yorke conjecture on periodic solution of differential delay equations, J. Li and X-Z. He.

170 On the construction of periodic solutions of Kaplan-Yorke type for some differential delay equations, K. Gopalsamy, J. Li and X-Z. He.

171 Uniformly accurate schemes for hyperbolic systems with relaxation, R.E. Caflisch, S. Jin and G. Russo.

172 Numerical integrations of systems of conservation laws of mixed type, S. Jin.

173 A complex algorithm for computing Lyapunov values, R. Mao and D. Wang.
174 On minimal sets of scalar parabolic equations with skew-product structures, W. Shen and Y. Yi.

175 Lorenz type attractors from codimensional-one bifurcation, V. Afraimovich, S.-N. Chow and W. Liu.

176 Dynamics in a discrete Nagumo equation - spatial chaos, S.-N. Chow and W. Shen.

177 Stability and bifurcation of traveling wave solutions in coupled map lattices, S.-N. Chow and W. Shen.

178 Density of defects and spatial entropy in extended systems, V.S. Afraimovich and L.A. Bunimovich.

179 On the second eigenvalue of the Laplace operator penalized by curvature, E.M. Harrell II.

180 Singular limits for travelling waves for a pair of equations, V. Hutson and K. Mischaikow.

181 Conley Index Theory: Some recent developments, K. Mischaikow.
182 Variational principle for periodic trajectories of hyperbolic billiards, L.A. Bunimovich.

183 Epidemic waves: A diffusion model for fox rabies, W.M. Rivera.
184 High complexity of spatial patterns in gradient reaction-diffusion systems, V. Afraimovich, A. Babin and S.-N. Chow.

185 Spatial chaotic structure of attractors of reaction-diffusion systems, V. Afraimovich, A. Babin and S.-N. Chow.

186 Explosive non-negative solutions to two point boundary value problems, V. Anuradha, C. Brown and R. Shivaji.

187 Existence results for non-autonomous elliptic boundary value problems, V. Anuradha, S. Dickens, and R. Shivaji.

188 Existence results for superlinear semipositone BVP's, V. Anuradha, D.D. Hai and R. Shivaji.

189 Hopf bifurcation analysis for hybrid systems, J.K. Hale and W. Huang.
190 Dynamics of cyclic feedback systems, T. Gedeon and K. Mischaikow.
191 Cyclic feedback systems, T. Gedeon.

192 Traveling waves as limits of solutions on bounded domains, G. Fusco, J.K. Hale and J. Xun.

193 Equivalence of topological and singular transition matrices in the Conley index theory, C.K. McCord and K. Mischaikow.

194 Coupled oscillators on a circle, J.K. Hale.
195 Conventional multipliers for homoclinic orbits, V. Afraimovich, W. Liu and T. Young.

196 Wave train of monotone travelling waves for generic competition-diffusion systems, Y. Kan-on.

197 Instability of stationary solutions for Lotka-Volterra competition model with diffusion, Y. Kan-on.

198 Hyperbolic homoclinic points of $\mathbb{Z}^{d}$-actions in lattice dynamical systems, V.S. Afraimovich, S-N. Chow and W. Shen.

199 Least energy solutions of semilinear Neumann problems and asymptotics, X-B. Pan and X. Xu.

200 Effects of delays on dynamics, J.K. Hale.
201 Periodic solutions of singularly perturbed delay equations, J.K. Hale and W. Huang.

202 Uniform ultimate boundedness and synchronization for nonautonomous equations, V. Afraimovich and H.M. Rodrigues.

203 Algebraic transition matrices, R. Franzosa and K. Mischaikow.
CDSNS95-
204 Ergodicity of minimal sets in scalar parabolic equations, W. Shen and Y. Yi.
205 Existence of infinitely many connecting orbits in a singularly perturbed ordinary differential equation, H. Kokubu, K. Mischaikow and H. Oka.

206 Asymptotic periodicity of positive solutions of reaction diffusion equations on a ball, X-Y. Chen and P. Polácik.

207 Asymptotic behaviour of ground states, J. Hulshof and R. C.A.M. van der Vorst.
208 Synchronizations in lattices of nonlinear Duffing's oscillators, V.S. Afraimovich and S-N. Chow.

209 Population patterns for ecological models, H.I. Freedman and W.M. Rivera.
210 Generalization of a theorem of Malta and Palis, V. Afraimovich and T. Young.
211 Singular index pairs, K. Mischaikow, M. Mrozek and J.F. Reineck.

## OMB Number 3145-0058

NATIONAL SCIENCE FOUNDATION
4201 Wilson Blvd.
Arlington, VA 22230

PVPD Name and Address
Hale, Jack K

Atlanta, GA 303320420
United States

## NATIONAL SCIENCE FOUNDATION FINAL PROJECT REPORT

## PART I - PROJECT IDENTIFICATION INFORMATION



NSF Grant Conditions (Article 17, GC-1, and Article 9, FDP-11) require submission of a Final Project Report (NSF Form 98A) to the NSF program officer no later than 90 days alter the expiration of the award. Final Project Reports for expired awards must be received before new awards can be made (NSF Grants Policy Manual Section 677).

Below, or on a separate page attached to this form, provide a summary of the completed projects and technical information. Be sure to include your name and award number on each separate page. See below for more instructions.

## PART II - SUMMARY OF COMPLETED PROJECT (for public use)

The summary (about 200 words) must be self-contained and intelligible to a scientifically or technically literate reader. Without restating the project title, it should begin with a topic sentence stating the project's major thesis. The summary should include, if pertinent to the project being described, the following items:

- The primary objectives and scope of the project
- The techniques or approaches used only to the degree necessary for comprehension
- The findings and implications stated as concisely and informatively as possible

The publications under this grant centered on the qualitative properties of the flows defined by partial differential equations and functional differential equations - attractors, periodic solutions, Hopf blfurcation, traveling waves, coupled oscillators, synchronization, the effects of delays on dynamics and control and the relationship between parabolic regularization of hyperbolic conservation laws and the singular IImit.

PART III - TECHNICAL INFORMATION (for program management use)

List references to publications resulting from this award and briefly describe primary data, samples, physical collections, inventions, software, etc. created or gathered in the course of the research and, if appropriate, how they are being made available to the research community. Provide the NSF Invention Disclosure number for any invention.
(See enclosed TECHNICAL INFORMATION)

I certity to the best of my knowledge (1) the statements herein (excluding scientific hypotheses and scientific opinion) are true and complete, and (2) the text and graphics in this report as well as any accompanying publications or other documents, unless othenwise indicated, are the original work of the signatories or of individuals working under their supervision. I understand that willuly making a fasse statement or concealing a material fact in this report or any other communication suyritted to NSF is a criminal offenee (U.S. Code, Thie 18, Section 1001).


## IMPORTANT: <br> MAILING INSTRUCTIONS

Return this entire packet plus all attachments in the envelope attached to the back of this form. Please copy the information from Part I, Block I to the Attention block on the envelope.

## PART IV -- FINAL PROJECT REPORT -- SUMMARY DATA ON PROJECT PERSONNEL <br> (To be submitted to cognizant Program Officer upon completion of profect)

The data requested below are important for the development of a statistical profile on the personnel supported by Federal grants. The information on this part is solicited in response to Public Law 99-383 and 42 USC 1885C. All information provided will be treated as confidential and will be safeguarded in accordance with the provisions of the Privacy Act of 1974. You should submit a single copy of this part with each final project report. However, submission of the requested information is not mandatory and is not a precondition of future award(s). Check the "Decline to Provide Information" box below if you do not wish to provide the information.

Please enter the numbers of individuals supported under this grant.
Do not enter information for individuals working less than 40 hours in any calendar year.

|  | Senior Staff |  | PostDoctorals |  | Graduate Students |  | UnderGraduates |  | Other Participants ${ }^{\prime}$ |  |
| :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: |
|  | Male | Fem. | Male | Fem. | Male | Fem. | Male | Fern. | Male | Fem. |
| A. Total, U.S. Citizens | 2 |  |  |  |  |  |  |  |  |  |
| B. Total, Permanent Residents |  |  |  |  |  |  |  |  |  |  |
| U.S. Citizens or Permanent Residents ${ }^{2}$ : <br> American Indian or Alaskan Native . |  |  |  |  |  | $\ldots$ |  |  |  |  |
| Asian. |  |  |  |  |  |  |  |  |  |  |
| Black, Not of Hispanic Origin. |  |  |  |  |  |  |  |  |  |  |
| Hispanic. |  |  |  |  |  |  |  |  |  |  |
| Pacific Islander |  |  |  |  |  |  |  |  |  |  |
| White, Not of Hispanic Origin |  |  |  |  |  |  |  |  |  |  |
| C. Total, Other Non-U.S. Citizens |  |  |  |  | 3 |  |  |  |  |  |
| Specify Country <br> 1. |  |  |  |  |  |  |  |  |  |  |
| 2. |  |  |  |  |  |  |  |  |  |  |
| 3. |  |  |  |  |  |  |  |  |  |  |
| D. Total, All participants $(A+B+C)$ | 2 |  |  |  | 3 |  |  |  |  |  |
| Disabled ${ }^{3}$ |  |  |  |  |  |  |  |  |  |  |

$\square$ Dectine to Provide information: Check box if you do not wish to provide this information (you are still required to return this page along with Parts I-III).
${ }^{1}$ Category includes, for example, college and precollege teachers, conference and workshop participants.
${ }^{2}$ Use the category that best describes the ethnic/racial status for all U.S. Chizens and Non-citizens with Permanent Residency. (If more than one category applies, use the one category that most closely reflects the person's recogntion in the community.)
${ }^{3}$ A person having a physical or mental impairment that substantially limits one or more major life activities; who has a record of such impaiment; or who is regarded as having such impeimment. (Disabled individuals also should be counted under the appropriate ethnic/racial group unless they are classified as "Other Non-U.S. Citizens.")

AMERICAN INDAAN OR ALASKAN NATIVE: A person having origins in any of the original peoples of North America and who maintains cultural identification through tribal afflliation or community recognition.
ASIAN: A person having origins in any of the original peoples of East Asia, Southeast Asia or the Indian subcontinent. This area inctudes, for example, China, Incia, Indonesia, Japan, Korea and Vietnam.
BLACK, NOT OF HISPANIC ORIGIN: A person having origins in any of the black racial groups of Africa.
HiSPANIC: A person of Mexican, Puerto Rican, Cuben, Central or South American or other Spenish culure or origin, regardless of race.
PACIFIC ISLANDER: A person having origins in any of the orignal peoples of Hawal; he U.S. Pacific ferritories of Guam, American Samoa, and the Northem Marinas; the U.S. Trust Terriory of Paleu; the islands of Micronesia and Metanesia; or the Philippines.
WHITE, NOT OF HISPANIC ORIGIN: A person having origins in any of the original peoples of Europe, North Africa, or the Middle East.

1. Attractors in inhomogeneous conservation laws and parabolic regularizations. Trans. Am. Math. Soc 347(1995), 1239-1254.
2. Limits of semigroups depending on a parameter. Resenhas 1 (1993), 1-45.
3. Attractors and dynamics in partial differential equations. Nato-MSI Lectures, Cambridge, 1995.
4. (with Fusco and Xun) Traveling waves as limits of solutions of bounded domains. SIAM J. Math. Anal. 27(1996), 1554-1558.
5. (with Huang) Hopf bifurcation analysis for hybrid systems. In Differential Equations and Applications to Biology and Industry (Eds. Martell et al) World Scientific (1996), 157-172.
6. Coupled oscillators on a circle. Resenhas 1(1994), 441-457.
7. (with Huang) Periodic solutions of singularly perturbed delay equations. ZAMP 47(1996), 57-88.
8. Effects of delays on dynamics. In Topological Methods in Differential Equations and Inclusions (EDS. Granas and Frigon), Kluwer(1995), 191-238.
9. (with Afraimovich and Chow) Synchronization in lattices of coupled oscillatiors. Physcia D 103(1997), 442-451.
10. Diffusive coupline, dissipation and synchronization. J. Dyn. Diff. Eqns. 9(1997),1-52.
11. Dynamics of a scalar parabolic equation. Canadian Appl. Math. Quart. 5(3)(1997), 209-305.
12. 'Long-Time Behavior as A-0 of Parabolic Equations depending on At', by Anatoli Babin and Shui-Nee Chow.
13. 'Two definitions of Exponential Dichotomy for Skew-Product Semiflow in Banach Spaces' by Shui-Nee Chow and Hugo Leiva.
14. 'Floquet Bundles for Scalar Parabolic Equations' by Shui-Nee Chow, Kening Lu, and John Mallet-Paret.

[^0]:    *This research was partially supported by NSF grant DMS-9306265
    ${ }^{\dagger}$ CDSNS Georgia Tech, Atlanta, Georgia 30332
    ${ }^{\ddagger}$ CDSNS Georgia Tech and ULA-Venezuela: e-mail leiva@math.gatech.edu

[^1]:    ${ }^{1} K e y$ words. time dependent linear differential equations, skew-product flow, dynamical spectrum, Lyapunov exponents. This research was partially sopported by NSF grant DMS-9005420 and ULA AMS(MOS) subject classifications. primary 34C35, 34D05, 34G10

[^2]:    * Partially supported by NSF and DARPA

[^3]:    ${ }^{1}$ Partially supported by ARO Grant DAAH0493G0199 and NSF Grant DMS9404199.
    ${ }^{2}$ Partially supported by ARO, NSF and NIST.

