# Topological Orbital Equivalence with Asymptotic Phase for a Two Time-Scales Discrete-Time System* 


#### Abstract

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Abstract. Existence, smoothness, approximation, and attractiveness of a locally integral manifold are established for a two time-scales discrete-time system. This manifold contains all the solutions remaining in a specific compact subset. It allows us to define locally a triangular system which is topologically orbitally equivalent with asymptotic phase. It follows that (in)stability properties and existence of solutions of the original system, remaining after some time instant in the abovementioned compact subset, can be established from the study of a reduced-order system. We study this reduced-order system for a weakly nonstationary case, applying the stroboscopic method to approximate it by a practically meaningful, slowly time-varying system.


Key words. Discrete-time system, Two time-scales system, Topological equivalence, Integral manifold, Averaging.

## 1. Introduction

The objective of this paper is to provide a tool for local analysis of the following discrete-time system:

$$
\begin{align*}
X_{k+1} & =A\left(\theta_{k}\right) X_{k}+B\left(\theta_{k}\right) u_{k} \\
\theta_{k} & =\theta_{k-1}+\varepsilon C\left(X_{k}, \theta_{k-1}, k\right)
\end{align*}
$$

For $\varepsilon$ small enough and when $A(\theta)$ has no eigenvalue on the unit circle, this system exhibits a two time-scales property. The state of the fast subsystem is $X$ in $\mathbb{R}^{n}$ and the state of the slow subsystem is $\theta$ in $\mathbb{R}^{p}$. The $u_{k}$ represents a (closed-loop) forcing term.

To determine çonditions for which $S_{\varepsilon}$ has solutions bounded on $\mathbb{Z}$ or a semiinfinite time interval and when these solutions are stable, we consider $S_{\varepsilon}$ as a small pertubation of the "frozen" system

$$
\begin{align*}
X_{k+1} & =A\left(\theta_{k}\right) X_{k}+B\left(\theta_{k}\right) u_{k} \\
\theta_{k} & =\theta_{k-1} \tag{0}
\end{align*}
$$

which can be viewed as a family of linear systems indexed by $\theta_{k}$.

[^0]System $S_{\varepsilon}$ arises from the study of linear time-varying systems where the time variations, although small, depend on the state itself. Also, as noticed by Ljung and Söderström [LS] for example, $S_{z}$ describes most adaptive linear systems studied in the literature.

When the sequence $u_{k}$ satisfies the so-called test input assumption (see [RPK] and [PR]), i.e., when there exists a solution of $S_{\varepsilon}$ on $\mathbb{Z}$ or on a semi-infinite time interval whose $\theta$-component is constant, Riedle and Kokotovic [RK1] have performed stability analysis by linearization and by invoking the Krylov-BogoliubovMitropolski averaging theory (see [H1]). This theory has long been used in the case when $u_{k}$ is a stationary stochastic process. In particular, it has been used by Ljung to derive the "ordinary differential equation" technique for the case when $C(X, \theta, k)$ decreases to zero as $k$ tends to infinity [L1], [LS], [KC]. This technique has also been justified for the case where $C(X, \theta, k)$ is small but not decreasing [BMP], [DF]. In the deterministic case and for the specific system $S_{\varepsilon}$, its use as a heuristic in the nondecreasing case has been considered by Åström [A1], [A2]. Based on a linear averaging technique, but incorporating total stability arguments, some relaxation of the test input assumption has been obtained [ABJ].

When $u_{k}$ has only $l_{2}$-stationary properties (periodic, almost periodic, ...), existence of a particular $\mathbb{Z}$-bounded solution has been established using nonlinear averaging theory [BSA] or the Poincaré expansion method [P4], [PP]. It relies on the existence of a solution of a bifurcation equation. Again stability properties are established by linearization.

In this paper, up to Section 4, no other assumption, besides boundedness, are needed on $u_{k}$. Moreover, we are interested in a complete description of all the bounded solutions. Such a description is easily obtained for $S_{0}$ for any $\theta$-set where $A(\theta)$ is noncritical with respect to $u_{k}$. For example, all the $\mathbb{Z}$-bounded solutions of $S_{0}$ whose $\theta$-component lies in S (subset of $\mathbb{R}^{p}$ for which the eigenvalues of $A(\theta)$ are in the open unit disk), are given by the graph of the function $M_{0}: \mathbb{R}^{p} \times \mathbb{Z} \rightarrow \mathbb{R}^{n}$ :

$$
\begin{equation*}
M_{0}(\theta, k)=\sum_{i=-\infty}^{k-1} A(\theta)^{k-1-i} B(\theta) u_{i} \tag{1}
\end{equation*}
$$

This graph is an integral manifold which, restricted to $S$ for example, is normally attractive. A general theory is available proving the persistence of normally hyperbolic invariant manifolds under small perturbations (see [F2], [H1], [HPS], and [O]). Therefore we expect the existence of a function $M_{z}$ whose graph is an integral manifold of $S_{\varepsilon}$. And like $M_{0}$ for $S_{0}, M_{\varepsilon}$ should describe (at least locally) all the $\mathbb{Z}$-bounded solutions of $S_{\varepsilon}$.

After Section 2 where we introduce our assumptions and notations, in Section 3 we apply this general theory to establish existence and regularity of this function $M_{\varepsilon}$ and we prove that any solution ( $X_{k+1}, \theta_{k}$ ) of $S_{\varepsilon}$ which remains for ever in $B(O, x) \times S$ lies in the graph of $M_{\varepsilon}(B(O, x)$ is the closed ball of center $O$, radius $x)$.

Our main result is given in Section 4. We prove that the so-called "reduction principle" applies with an asymptotic phase. This "principle" has been introduced by Pliss [P1] and generalized by Aulbach [A3], Henry [H2], and Kelley [K]. It is of practical importance since it shows that stability properties and existence of solutions of $S_{\varepsilon}$ remaining in $B(O, x) \times S$ after time $k_{0}$ can be established from the
reduced-order system given by the restriction of $S_{\varepsilon}$ to the graph of $M_{\varepsilon}$, i.e., from the system $S M_{\varepsilon}^{\theta}$ obtained from the $\theta$-equation of $S_{\varepsilon}$ by replacing $X_{k}$ by $M_{\varepsilon}\left(\theta_{k-1}, k\right)$. We will see that this "principle" or more precisely the existence of a topological orbital equivalence is a constquence of Bowen's Shadowing Lemma (see [HPS] or Proposition 8.19 of [S]) which is again an aspect of normal hyperbolicity.

In Section 5 we study the system $S M_{\varepsilon}^{\theta}$ for the case where $u_{k}$ is weakly nonstationary. We apply Minorsky's stroboscopic method ideas [M]. This allows us to define a subset of $\mathbb{R}^{p}$ where the solutions of $S M_{\varepsilon}^{\theta}$ are hyperbollically attractive (and therefore, the same property holds for solutions of $S_{\varepsilon}$ ). This proves in particular that Åström's heuristic approach is technically sound when restricted to this subset.

Some of the results presented here have been established for the continuous-time case by Riedle and Kokotovic [RK2]. For the discrete-time case, they have been announced in [P3].

## 2. Assumptions and Notations

In the following, $\|\cdot\|$ denotes the Euclidian norm. The case when $u_{k}$ and $C(X, \theta, k)$ are periodic in $k$ is called the periodic case. Throughout this paper, $\varepsilon$ is positive.

The following assumptions are used:
A1. The sequence $u_{k}$ is bounded: $\left\|u_{k}\right\| \leq u$ for all $k \in \mathbb{Z}$.
Let $S$ be a compact set in $\mathbb{R}^{p}$, with a nonempty interior.
A2. If $\theta$ lies in $\mathbf{S}$, the eigenvalues of $A(\theta)$ are strictly inside the unit circle: $|\lambda\{A(\theta)\}| \leq \lambda_{0}$ for all $\theta \in \mathrm{S}, \lambda_{0}<1$.

A3. The functions $A(\theta), B(\theta)$ are Lipschitz continuous on S with $a_{1}, b_{1}$ the respective Lipschitz constants.

A direct consequence of assumptions A 2 and A 3 is the existence of $a, b, \lambda$ such that (see [F3]) $\|B(\theta)\| \leq b$ and $\left\|A(\theta)^{i}\right\| \leq a \lambda^{i}$, where $\lambda_{0}<\lambda<1$, for all $\theta \in \mathrm{S}$, for all $i \in \mathbb{N}$.

Let $B(O, x)$ be the closed ball of $\mathbb{R}^{n}$ with center $O$ and radius $x$.
A4. The function $C$ is Lipschitz continuous uniformly on $B(O, x) \times \mathbf{S} \times \mathbb{Z}$ :

$$
\begin{aligned}
\|C(X, \theta, k)\| & \leq c(x) \\
\left\|C\left(X^{0}, \theta^{0}, k\right)-C\left(X^{1}, \theta^{1}, k\right)\right\| & \leq c_{1}(x)\left(\left\|X^{0}-X^{1}\right\|+\left\|\theta^{0}-\theta^{1}\right\|\right)
\end{aligned}
$$

where $c, c_{1}$ are positive nondecreasing functions whose argument is omitted when no confusion is possible. In the following, the function $C$ could also depend on $\varepsilon$.

To study the smoothness properties of $M_{\ell}(\theta, k)$, we need:
A5. The functions $A, B, C$ are Lipschitz continuously differentiable: There exist linear maps $\partial A / \partial \theta, \partial B / \partial \theta, \partial C / \partial X$, and $\partial C / \partial \theta$ such that, with $(X, \theta)$ an interior point of $B(O, x) \times \mathbf{S}$ and $h_{x}, h_{\theta}$ two (sufficiently small) vectors in $\mathbb{R}^{n}, \mathbb{R}^{p}$ respectively, the
functions $\delta_{A}, \delta_{B}, \delta_{C}$ defined as (similarly for $\delta_{A}, \delta_{B}$ )

$$
\begin{aligned}
& \delta_{C}\left(X, \theta, k, h_{x}, h_{\theta}\right) \\
& \quad=C\left(X+h_{x}, \theta+h_{\theta}, k\right)-C(X \theta, k)-\frac{\partial C}{\partial X}(X, \theta, k) h_{x}-\frac{\partial C}{\partial \theta}(X, \theta, k) h_{\theta}
\end{aligned}
$$

satisfy (similarly for $\delta_{A}, \delta_{B}$ )

$$
\limsup _{h_{\theta} \rightarrow 0, h_{x} \rightarrow 0} \frac{\left\|\delta_{C}\left(X, \theta, k, h_{x}, h_{\theta}\right)\right\|}{\left\|h_{x}\right\|+\left\|h_{\theta}\right\|}=0
$$

Moreover, the functions $\partial A / \partial \theta(\theta), \partial B / \partial \theta(\theta), \partial C / \partial X(X, \theta, k)$, and $\partial C / \partial \theta(X, \theta, k)$ are Lipschitz continuous uniformly on $B(O, x) \times S \times \mathbb{Z}$ with $a_{2}, b_{2}, c_{2}(x)$, and $c_{2}(x)$ the respective Lipschitz constants.

Assumptions A1 and A3-A5 are smoothness conditions. Assumption A2 concerns the possibility of finding a constant $\theta$ such that the (linear) $X$-subsystem is exponentially stable. These assumptions are very similar to conditions C 1 and C 2 on p. 158 of [LS], introduced to derive the associated differential equation technique. One of the consequence of the results presented in this paper is to provide a geometric justification of the substitution of $X_{k+1}$ by $M_{0}\left(\theta_{k}, k\right)$ in the $\theta$-equation of $S_{\varepsilon}$, used to derive this associated differential equation. However, notice that for the time being no stationarity or decaying $\varepsilon$ is needed compared with C3-C6 used in the above reference. These extra assumptions are discussed in Section 5, when studying a reduced-order system.

To obtain this geometric insight, we consider $S_{\varepsilon}$ as a map from $\mathbb{R}^{n} \times \mathbb{R}^{p} \times \mathbb{Z}$ to $\mathbb{R}^{n} \times \mathbb{R}^{p} \times \mathbb{Z}$ taking

$$
S_{\varepsilon}:\left[\begin{array}{c}
X \\
\theta \\
k
\end{array}\right] \mapsto\left[\begin{array}{c}
Y=A(\psi) X+B(\psi) u_{k} \\
\psi=\theta+\varepsilon C(X, \theta, k) \\
k+1
\end{array}\right]
$$

Similarly, given the function $M_{\varepsilon}: S \times \mathbb{Z} \rightarrow \mathbb{R}^{n}$, we consider the system $S M_{\varepsilon}$ defined by

$$
S M_{\varepsilon}:\left[\begin{array}{c}
X^{M} \\
\theta^{M} \\
k
\end{array}\right] \mapsto\left[\begin{array}{c}
Y^{M}=A\left(\psi^{M}\right) X^{M}+B\left(\psi^{M}\right) u_{k} \\
\psi^{M}=\theta^{M}+\varepsilon C\left(M_{\varepsilon}\left(\theta^{M}, k\right), \theta^{M}, k\right) \\
k+1
\end{array}\right]
$$

The difference between $S_{\varepsilon}$ and $S M_{\varepsilon}$ is that, in the $\theta$-equation, $X$ is replaced by $M_{\varepsilon}$. Consequently, $S M_{\varepsilon}$, being "lower triangular," is structurally much simpler than $S_{e}$.

To study boundedness and stability properties of these systems, we consider the following sets: Let $\Theta$ be a compact subset of $\mathbf{S}$, we define $\mathbf{B}\left(\left[k_{0}, k_{1}\right), \Theta\right.$ ) (resp. $\mathbf{B}^{M}\left(\left[k_{0}, k_{1}\right), \Theta\right)$ ) as the set of sequences $\left(X_{k+1}, \theta_{k}\right)\left(\operatorname{resp} .\left(X_{k+1}^{M}, \theta_{k}^{M}\right)\right)$ solutions of $S_{\varepsilon}$ (resp. $S M_{\varepsilon}$ ) for $k$ in $\left[k_{0}, k_{1}\right.$ ) and satisfying

$$
\begin{aligned}
& X_{k_{0}+1}-M_{\varepsilon}\left(\theta_{k_{0}}, k_{0}+1\right) \in B\left(0, \frac{x-m}{\alpha}\right) \\
& \quad\left(\text { resp. } X_{k_{0}+1}^{M}-M_{\varepsilon}\left(\theta_{k_{0}}^{M}, k_{0}+1\right) \in B\left(0, \frac{x \pm m}{\alpha}\right)\right), \\
& \theta_{k} \in \Theta \quad\left(\text { resp. } \theta_{k}^{M} \in \Theta\right) \quad \text { for all } \quad k \in\left[k_{0}, k_{1}\right)
\end{aligned}
$$

where $m, \alpha$ are specified later. For sequences in these sets we define the elevation above the manifold as (similarly for $E^{M}$ ) $E_{k+1}=X_{k+1}-M_{\varepsilon}\left(\theta_{k}, k+1\right.$ ).
In the following, we say that the graph of $M_{\varepsilon}$ is an integral manifold of $S_{\varepsilon}$ locally on $\Theta$ if, for any $\left(\theta_{k}^{M}\right)$ solution of

$$
\theta_{k+1}^{M}=\theta_{k}^{M}+\varepsilon C\left(M_{\varepsilon}\left(\theta_{k}^{M}, k+1\right), \theta_{k}^{M}, k+1\right)
$$

such that $\theta_{k}^{M}$ lies in $\Theta$ for all $k$ in $\left[k_{0}, k_{1}\right.$ ), the sequence ( $X_{k+1}, \theta_{k}$ ) defined by $X_{k+1}=M_{\varepsilon}\left(\theta_{k}^{M}, k+1\right), \theta_{k}=\theta_{k}^{M}$, is a solution of $S_{\varepsilon}$ on $\left[k_{0}, k_{1}\right)$. Note that this implies that the graph of $M_{\varepsilon}$ is an integral manifold of $S M_{\varepsilon}$ locally on $\Theta$.
Our focus in this paper concerns the properties of solutions of $S_{\varepsilon}$ remaining in $B(O, x) \times S$. To facilitate our study, we use the classical trick consisting of modifying the system $S_{\varepsilon}$ into $\bar{S}_{\varepsilon}$, such that $S_{\varepsilon}$ and $\bar{S}_{\varepsilon}$ coincide on a compact subset of $B(O, x) \times \mathbf{S}$. For a compact set $\Theta$ in $\mathbb{R}^{p}$ we denote by $\Theta+\eta$, the " $\eta$-augmented" compact set of $\Theta$ :

$$
\Theta+\eta=\left\{\theta \in \mathbb{R}^{p} \mid \exists \psi \in \Theta:\|\psi-\theta\| \leq \eta\right\} .
$$

Since $\mathbf{S}$ has a nonempty interior we can find $\eta$ and compact sets $\mathbf{S}_{0}, \mathbf{S}_{1}$ with nonempty interiors, such that $\mathbf{S}_{0}+\eta \subset \mathbf{S}_{1}$ and $\mathbf{S}_{1}+\eta \subset \mathbf{S}$. We call a stopping function $s: \mathbb{R}^{p} \rightarrow$ $[0,1]$, a $C^{r}$-function $(0 \leq r \leq \infty)$ given by the Urysohn theorem satisfying $s(\theta)=1$ if and only if $\theta \in \mathbf{S}_{0}$, and $\theta \notin \mathbf{S}_{1}$ implies $s(\theta)=0$.

Let $s_{1}, s_{2}$ denote the Lipschitz constants of $s$ and its differential, respectively. We can take $s_{1}=1 / \eta$ if $s$ is only Lipschitz continuous. We define the function $\bar{C}: \mathbb{R}^{n} \times \mathbb{R}^{p} \times \mathbb{Z} \rightarrow \mathbb{R}^{p}$ by $\bar{C}(X, \theta, k)=s(\theta) C(X, \theta, k)$. The function $\bar{C}$ has the same properties as $C$, with, in particular, $\bar{c} \leq c, \bar{c}_{1} \leq s_{1} c+c_{1}$, and $\bar{c}_{2} \leq s_{2} c+2 s_{1} c_{1}+c_{2}$.

In this paper we always assume $\varepsilon \leq \varepsilon_{0}$ with $\varepsilon_{0} \leq \operatorname{Min}\left(\eta / c, 1 / \bar{c}_{1}\right)$. In this condition, if $\psi$ is defined as $\psi=\theta+\varepsilon \bar{C}(X, \theta, k)$, we have (i) $\theta \in \mathbf{S}_{1}$ implies $\psi \in \mathbf{S}$, (ii) $\psi \in \mathbf{S}-\mathbf{S}_{1}$ implies $\theta=\psi$, and (iii) $\theta \in \mathbf{S}$ implies the segment $[\theta, \psi] \subset \mathbf{S}$. For $\bar{C}$, we define the modified system $\bar{S}_{\varepsilon}$ (similarly for $\bar{S} M_{\varepsilon}$ ):

$$
\bar{S}_{\varepsilon}:\left[\begin{array}{c}
X \\
\theta \\
k
\end{array}\right] \rightarrow\left[\begin{array}{c}
Y=A(\psi) X+B(\psi) u_{k} \\
\psi=\theta+\varepsilon \bar{C}(X, \theta, k) \\
k+1
\end{array}\right] .
$$

In $\bar{S}_{c}$ we smoothly stop the $\theta$-component of any solution trying to leave $\mathrm{S}_{1}$. Clearly, a solution of $\bar{S}_{\varepsilon}$ is a solution of $S_{\varepsilon}$ on $\left[k_{0}, k_{1}\right)$ if, for all $k$ in $\left[k_{0}, k_{1}\right), \theta_{k}$ lies in $\mathbf{S}_{0}$. The idea of preventing the $\theta$-vector from leaving an admissible region is also used in practice [E], [LS].
To end this section, let us mention that the nonexplicit form of $S_{\varepsilon}$ is chosen to simplify the forthcoming derivations. However, this is done with no loss of generality since the system

$$
\begin{aligned}
& Y_{k+1}=A\left(\theta_{k}\right) Y_{k}+B\left(\theta_{k}\right) u_{k}, \\
& \theta_{k+1}=\theta_{k}+\varepsilon C\left(Y_{k}, \theta_{k}, k+1\right)
\end{aligned}
$$

can also be written as

$$
\begin{gathered}
\binom{Y_{k+1}}{Y_{k}}=\left(\begin{array}{cc}
A\left(\theta_{k}\right) & 0 \\
I & 0
\end{array}\right)\binom{Y_{k}}{Y_{k-1}}+\binom{B\left(\theta_{k}\right)}{0} u_{k}, \\
\theta_{k}=\theta_{k-1}+\varepsilon C\left(Y_{k-1}, \theta_{k-1}, k\right) .
\end{gathered}
$$

## 3. Existence of a Normally Attractive Locally Integral Manifold

Theorem 1. Under assumptions $\mathrm{A} 1-\mathrm{A} 4$, for any compact set $\mathrm{S}_{0}$ strictly contained in S , there exists $\varepsilon^{*}$ such that for any $\varepsilon, \varepsilon \leq \varepsilon^{*}$, we can find constants $\alpha, f^{0}, f^{1}, f_{1}^{0}$, a function $\tau_{0}(\varepsilon, x), 0 \leq \tau_{0}<1$, and a (possibly nonunique) function $M_{\varepsilon}: S \times \mathbb{Z} \rightarrow \mathbb{R}^{n}$ which is periodic in $k$ in the periodic case, such that:
(i) The graph of $M_{\varepsilon}$ is an integral manifold of $S_{\varepsilon}$ locally on $\mathrm{S}_{0}$.
(ii) Smoothness: $M_{\varepsilon}$ is bounded and Lipschitz continuous uniformly on $\mathbf{S} \times \mathbb{Z}$ with bounds $m$ and $m_{1}$, respectively. Moreover, if assumption A5 holds, $M_{\varepsilon}$ is Lipschitz continuously differentiable in $\theta$ with $m_{2}$ as a Lipschitz constant of its differential.
(iii) Approximation: let $M_{0}$ be the function defined in(1), then, uniformly on $\mathrm{S}_{0} \times \mathbb{Z}$, $\left\|M_{\mathrm{e}}(\theta, k)-M_{0}(\theta, k)\right\| \leq \varepsilon f^{0}$. Moreover, if assumption A5 holds, $M_{0}$ is continuously differentiable in $\theta$ and letting $\left(M_{0}^{1}(\theta, k)\right)$ be the unique $\mathbb{Z}$-bounded solution of

$$
X_{k+1}=A(\theta) X_{k}+B(\theta) u_{k}-\varepsilon A(\theta) \frac{\partial M_{0}}{\partial \theta}(\theta, k) C\left(M_{0}(\theta, k), \theta, k\right)
$$

we have, uniformly on $\mathrm{S}_{0} \times \mathbb{Z}$,

$$
\left\|M_{\varepsilon}(\theta, k)-M_{0}^{1}(\theta, k)\right\| \leq \varepsilon^{2} f^{1}, \quad\left\|\frac{\partial M_{\varepsilon}}{\partial \theta}(\theta, k)-\frac{\partial M_{0}}{\partial \theta}(\theta, k)\right\| \leq \varepsilon f_{1}^{0}
$$

(iv) Attractiveness: let $\left(X_{k+1}, \theta_{k}\right)$ be a sequence of $\mathbf{B}\left(\left[k_{0}, k_{1}\right), \mathrm{S}_{0}\right)$, with $\left(E_{k}\right)$ its corresponding sequence of elevations above the manifold; we have for all $k^{\prime}, k$, $k_{0} \leq k^{\prime} \leq k<k_{1},\left\|E_{k+1}\right\| \leq \alpha \tau_{0}(\varepsilon, x)^{k-k^{\prime}}\left\|E_{k^{\prime}+1}\right\|$. Moreover, if $\left(X_{k+1}, \theta_{k}\right)$ is a solution of $S_{\varepsilon}$ which lies in $B(O, x) \times S_{0}$ for all $k$ in $\mathbb{Z}$, then this solution lies in the graph of $M_{e}$, namely, for all $k, X_{k+1}=M_{\varepsilon}\left(\theta_{k}, k+1\right)$.
(v) All these properties hold for $S M_{e}$.

All the constants appearing in this statement are clarified in the proof.
Remarks. For the continuous-time case, existence of $M_{e}$, approximation by $M_{0}$, and exponential decaying of the elevation above the manifold are established in [RK2].

This theorem is a technical step toward our main result of Section 4. However, it gives us a first important geometric property of $S_{\varepsilon}$ :

Any solution of $S_{\varepsilon}$ remaining in the compact set $B(O, x) \times S_{0}$ lies in the graph of $M_{\varepsilon}$.
The end of this section is devoted to the proof of this theorem. It is sufficient to establish this result for the modified system $\bar{S}_{\varepsilon}$. One possible proof would call upon general theorems on persistence of normally hyperbolic integral manifolds. We prefer a less technical direct proof. It is an adaptation of the proof of Theorem 5.2 of [S] and is based on the graph transform technique.

Consider the image by $\bar{S}_{z}$ of the graph $\left\{\left(X=M_{\varepsilon}(\theta, k), \theta, k\right) \mid \theta \in S, k \in \mathbb{Z}\right\}$. We obtain a set of $(Y, \psi, k+1)$ contained in $\mathbb{R}^{n} \times S \times \mathbb{Z}$. If the graph is an integral
manifold, this set is contained in the graph itself, i.e., $Y=M_{\varepsilon}(\psi, k+1)$. This means that the following diagram commutes:

$$
\left[\begin{array}{cc}
{\left[\begin{array}{c}
X=M_{z}(\theta, k) \\
\theta \\
k
\end{array}\right]} \\
\uparrow \\
M_{t}
\end{array}\right] \xrightarrow{\bar{s}_{\varepsilon}}\left[\begin{array}{c}
Y=M_{\varepsilon}(\psi, k+1) \\
\psi \\
k+1
\end{array}\right]
$$

This also means that $M_{\varepsilon}$ is a fixed point of the operator $T$, called the graph transform and defined as $T(M)=\bar{S}_{\varepsilon} \circ M \circ\left(\bar{S} M_{\varepsilon}^{\theta}\right)^{-1}$. Our problem is reduced to studying the properties of this operator. Let us first introduce an

Adapted Metric. Given $\theta$ in S , for any vector $X$ in $\mathbb{R}^{n}$, we define its norm $\|X\|_{\theta}$ by $\|X\|_{\theta}=\sum_{i=0}^{\infty} \mu^{-i}\left\|A(\theta)^{i} X\right\|, \lambda<\mu<1$. From assumptions A2 and A3 it can be seen that:
(i) $\|X\| \leq\|X\|_{\theta} \leq \alpha\|X\|, \alpha=a \mu /(\mu-\lambda)$.
(ii) $\|A(\theta) X\|_{\theta} \leq \gamma\|X\|_{\theta}, \gamma=\mu[1-(\mu-\lambda) / a \mu]$.
(iii) $\|X\|_{\psi} \leq(1+\beta\|\theta-\psi\|)\|X\|_{\theta}, \beta=a_{1} a /(\mu-\lambda)$.

Remark. This metric allows us to exhibit the normal hyperbolicity property. Let $\left(X^{0}, \theta^{0}, k\right),\left(X^{1}, \theta^{1}, k\right)$ be two points in $B(O, x) \times \mathbf{S} \times \mathbb{Z}$ and let ( $Y^{0}, \psi^{0}, k+1$ ), ( $Y^{1}, \psi^{1}, k+1$ ) be their respective images by $\bar{S}_{\varepsilon}$. We have the following inequalities (compare with (2.21) and (2.22) of [P2]):

$$
\begin{aligned}
&\left\|Y^{0}-Y^{1}\right\|_{\psi^{0}} \leq \gamma(1+\varepsilon \beta \bar{c})\left\|X^{0}-X^{1}\right\|_{\theta^{0}}+\alpha\left(a_{1} x+b_{1} u\right)\left\|\psi^{0}-\psi^{1}\right\| \\
&\left\|\psi^{0}-\psi^{1}\right\| \geq-\varepsilon \bar{c}_{1}\left\|X^{0}-X^{1}\right\|_{\theta^{0}}+\left(1-\varepsilon \bar{c}_{1}\right)\left\|\theta^{0}-\theta^{1}\right\|
\end{aligned}
$$

Introducing the positive function $l(x)$ satisfying

$$
\frac{\gamma(1+\varepsilon \beta \bar{c}(x)) l(x)}{1-\varepsilon \bar{c}_{1}(x)(1+l(x))}+\alpha\left(a_{1} x+b_{1} u\right) \leq l(x)
$$

we obtain the following key technical triangular system:

$$
\begin{align*}
\left\|Y^{0}-Y^{1}\right\|_{\psi^{0}}-l\left\|\psi^{0}-\psi^{1}\right\| \leq & \tau_{0}(\varepsilon, x)\left(\left\|X^{0}-X^{1}\right\|_{\theta}-l\left\|\theta^{0}-\theta^{1}\right\|\right) \\
\left\|\psi^{0}-\psi^{1}\right\| \geq & -\varepsilon \bar{c}_{1}\left(\left\|X^{0}-X^{1}\right\|_{\theta}-l\left\|\theta^{0}-\theta^{1}\right\|\right)  \tag{2}\\
& +\left(1-\varepsilon \bar{c}_{1}(1+l)\right)\left\|\theta^{0}-\theta^{1}\right\|
\end{align*}
$$

with

$$
\tau_{0}(\varepsilon, x)=\gamma(1+\varepsilon \beta \bar{c}(x))\left(1+\varepsilon \frac{\bar{c}_{1}(x) l(x)}{1-\varepsilon \bar{c}_{1}(x)(1+l(x))}\right) .
$$

The normal attractivity property appears here. In particular, $\tau_{0}(\varepsilon, x)$ characterizes the contraction property of $S_{\varepsilon}$ in the direction normal to the integral manifold (see attractiveness). The term $\left(1-\varepsilon \bar{c}_{1}(1+l)\right)^{-1}$ characterizes the possible expansion along the manifold (see Lemma 1). From these characterizations and following [F2] or
[HPS] we expect the existence of an integral manifold which can be up to $r$ times continuously differentiable if $r$ is the largest integer such that

$$
\begin{equation*}
\tau_{0}(\varepsilon, x)\left(\frac{1}{1-\varepsilon \bar{c}_{1}(x)(1+l(x))}\right)^{r}<1 . \tag{3}
\end{equation*}
$$

This means that existence and smoothness of this integral manifold depend directly on how much more sharply $S_{\varepsilon}$ contracts in the normal direction than expands in the tangent direction or else on how much faster the $X$-component converges to its "steady-state" value than the $\theta$-component is changed.

Let $\mathbf{M}$ be the set of functions $M: \mathbf{S} \times \mathbb{Z} \rightarrow \mathbb{R}^{n}$ satisfying uniformly

$$
\begin{aligned}
\|M(\theta, k)\|_{\theta} & \leq m<x, \\
\left\|M\left(\theta^{0}, k\right)-M\left(\theta^{1}, k\right)\right\|_{\theta^{i}} & \leq m_{2}\left\|\theta^{0}-\theta^{1}\right\|,
\end{aligned} \quad i=0,1,
$$

where $m, m_{1}$ are as specified later. In the periodic case, $M(\theta, k)$ is chosen periodic in $k$. Equipped with the distance associated with the norm $|M|=$ $\operatorname{Sup}_{\theta \in \mathrm{S}, k \in \mathbb{Z}}\| \| M(\theta, k) \|_{\theta}, \mathbf{M}$ is a complete metric space.

In our above definition of the graph transform, we used the inverse function $\left(\bar{S} M_{\varepsilon}^{\theta}\right)^{-1}$. Let us prove that this makes sense.

Lemma 1. There exists $\varepsilon^{*}$ such that for any $\varepsilon, \varepsilon \leq \varepsilon^{*}, M$ in $\mathbf{M}$, and $k$ in $\mathbb{Z}$ we can find a unique function $D(M, k): S \rightarrow \mathbb{R}^{p}$ such that, uniformly on $\mathbf{M} \times \mathbb{Z} \times \mathbf{S}$ :
(i) $\|D(M, k)(\psi)\| \leq \bar{c}$.
(ii) $\left\|D(M, k)\left(\psi^{0}\right)-D(M, k)\left(\psi^{1}\right)\right\| \leq d_{1}\left\|\psi^{0}-\psi^{1}\right\|$.
(iii) $\left\|D\left(M^{0}, k\right)(\psi)-D\left(M^{1}, k\right)(\psi)\right\| \leq d_{1}^{m}\left|M^{0}-M^{1}\right|$.
(iv) $D(M, k)(\psi)=\bar{C}(M(\psi-\varepsilon D(M, k)(\psi), k), \psi-\varepsilon D(M, k)(\psi), k)$.

Proof. Given $M, k, \psi$ in $\mathbf{M} \times \mathbb{Z} \times \mathbf{S}_{1}$, we consider the complete metric space of vectors $\mathbf{D}$ of $\mathbb{R}^{p}$ such that $\|D\| \leq \bar{c}$. With our assumption on $\varepsilon$, if $\psi$ is in $\mathbf{S}_{1}$, $\psi-\varepsilon D$ is in S . On D , we define an operator $T(M, \psi, k)$ as $T(M, \psi, k)(D)=$ $\bar{C}(M(\psi-\varepsilon D, k), \psi-\varepsilon D, k)$. We have

$$
\begin{aligned}
& \left\|T\left(M^{0}, \psi^{0}, k\right)\left(D^{0}\right)-T\left(M^{1}, \psi^{1}, k\right)\left(D^{1}\right)\right\| \\
& \quad \leq \bar{c}_{1}\left[\left(1+m_{1}\right)\left\|\psi^{0}-\psi^{1}\right\|+\varepsilon\left(1+m_{1}\right)\left\|D^{0}-D^{1}\right\|+\left|M^{0}-M^{1}\right|\right] .
\end{aligned}
$$

From the Uniform Contraction Mapping Theorem [H1], the result follows with

$$
d_{1}=\frac{\bar{c}_{1}\left(1+m_{1}\right)}{1-\varepsilon \bar{c}_{1}\left(1+m_{1}\right)}, \quad d_{1}^{m}=\frac{1}{1-\varepsilon \bar{c}_{1}\left(1+m_{1}\right)}
$$

However, up to now $\psi$ was restricted to lie in $S_{1}$. The extension to $S$ is obtained by taking $D(M, k)(\psi)=0$ for all $\psi \in S-S_{1}$. Using the properties of $\bar{C}$, we can easily check that this is a valid extension satisfying (i)-(iv).

This lemma proves that $\psi=\theta+\varepsilon \bar{C}(M(\theta, k), \theta, k)$ if and only if $\theta=$ $\psi-\varepsilon D(M, k)(\psi)$ and $\psi \in S-S_{1}$ implies $D(M, k)(\psi)=0$. We also remark that $D(M, k)(\psi)$ is periodic in the periodic case.

Having established that the graph transform is well defined, we study its proper-
ties in a slightly more general context. Let us consider a graph transform $T$ defined on $\mathbf{M}$ by

$$
\begin{aligned}
T(M)(\psi, k+1) & =A(\psi) M(\theta, k)+B^{\prime}(\psi, M(\theta, k), k), \\
* \quad \theta & =\psi-\varepsilon D(M, k)(\psi),
\end{aligned}
$$

where $B^{\prime}: \mathbf{S} \times B(O, m) \times \mathbb{Z} \rightarrow \mathbb{R}^{n}$ satisfies uniformly:
(i) $\left\|B^{\prime}(\psi, X, k)\right\|_{\psi} \leq b^{\prime}$.
(ii) $\left\|B^{\prime}\left(\psi^{0}, X^{0}, k\right)-B^{\prime}\left(\psi^{1}, X^{1}, k\right)\right\|_{\psi^{1}} \leq b_{1}^{\psi}\left\|\psi^{0}-\psi^{1}\right\|+\varepsilon b_{1}^{x}\left\|X^{0}-X^{1}\right\|_{\psi^{\prime}}, i=0,1$.

## Lemma 2. Function $T(M)$ is contained in $\mathbf{M}$ and $T$ is a contraction.

Proof. (a) Function $T(M)$ is in $\mathbf{M}$ :
(i) With our adapted metric, we have

$$
\|T(M)(\psi, k+1)\|_{\psi} \leq \gamma(1+\varepsilon \beta \bar{c})\|M(\theta, k)\|_{\theta}+b^{\prime} .
$$

Hence $m$ should satisfy, with $\varepsilon$ sufficiently small, $\gamma(1+\varepsilon \beta \bar{c}) m+b^{\prime} \leq m$.
(ii) Using the properties of $M, D$, for $\psi^{0}, \psi^{1}$ in S , we have

$$
\begin{aligned}
& \left\|T(M)\left(\psi^{0}, k+1\right)-T(M)\left(\psi^{1}, k+1\right)\right\|_{\psi 1} \\
& \quad \leq\left[\left(\gamma+\varepsilon b_{1}^{x}\right)(1+\varepsilon \beta \bar{c}) m_{1}\left(1+\varepsilon d_{1}\right)+\alpha a_{1} m+b_{1}^{\psi}\right]\left\|\psi^{0}-\psi^{1}\right\|, \\
& \quad i=0,1 .
\end{aligned}
$$

Hence $m_{1}$ should satisfy, with $\varepsilon$ sufficiently small,

$$
\left(1+\varepsilon d_{1}\right)\left(\gamma+\varepsilon b_{1}^{x}\right)(1+\varepsilon \beta \bar{c}) m_{1}+\alpha a_{1} m+b_{1}^{\psi} \leq m_{1} .
$$

(iii) Function $T(M)$ is periodic in $k$ in the periodic case.
(b) The graph transform $T$ is a contraction: Let $M^{0}, M^{1}$ be two elements of $\mathbf{M}$, we have

$$
\begin{aligned}
& \left\|T\left(M^{0}\right)(\psi, k+1)-T\left(M^{1}\right)(\psi, k+1)\right\|_{\psi} \\
& \quad \leq\left(\gamma+\varepsilon b_{1}^{x}\right)(1+\varepsilon \beta \bar{c})\left\|M^{0}\left(\theta^{0}, k\right)-M^{1}\left(\theta^{1}, k\right)\right\|_{\theta^{0}} \\
& \quad \leq\left(\gamma+\varepsilon b_{1}^{x}\right)(1+\varepsilon \beta \bar{c})\left[\left|M^{0}-M^{1}\right|+\varepsilon m_{1}\left\|D\left(M^{0}, k\right)(\psi)-D\left(M^{1}, k\right)(\psi)\right\|\right] \\
& \quad \leq\left(\gamma+\varepsilon b_{1}^{x}\right)(1+\varepsilon \beta \bar{c})\left(1+\varepsilon m_{1} d_{1}^{m}\right)\left|M^{0}-M^{1}\right| .
\end{aligned}
$$

Since $\gamma<1$, taking the supremum on $\mathbf{S} \times \mathbb{Z}$ gives the result for $\varepsilon$ sufficiently small.

To prove existence of $M_{\varepsilon}(\theta, k)$, we apply this lemma with $B^{\prime}(\psi, X, k)=B(\psi) u_{k}$. It follows that the graph transform has a fixed point in $\mathbf{M}$ if

$$
\begin{gathered}
\gamma(1+\varepsilon \beta \bar{c}(m)) m+b u \alpha \leq m \leq x, \\
\frac{m_{1}}{1-\varepsilon \bar{c}_{1}(m)\left(1+m_{1}\right.} \gamma(1+\varepsilon \beta \bar{c}(m))+\alpha\left(b_{1} u+a_{1} m\right) \leq m_{1}, \\
\tau_{0}(\varepsilon, m)=\gamma(1+\varepsilon \beta \bar{c}(m))\left(1+\varepsilon \frac{\bar{c}_{1}(m) m_{1}}{1-\varepsilon \bar{c}_{1}(m)\left(1+m_{1}\right)}\right)<1 .
\end{gathered}
$$

In particular, this confirms condition (3) of the previous Remark.

We have proved the existence of a Lipschitz continuous function $M_{\varepsilon}: S \times \mathbb{Z} \rightarrow \mathbb{R}^{n}$ whose graph is an integral manifold of $S_{\varepsilon}$ locally on $S_{0}$ and which is periodic in the periodic case. We now establish that, under assumption A5, $M_{\varepsilon}$ is Lipschitz continuously differentiable in $\theta$. Assuming for the time being that this property holds and denoting by $L_{\varepsilon}(\theta, k)$ what the differential of $M_{e}(\theta, k)$ should be, formal derivations show that $L_{\varepsilon}(\theta, k)$ satisfies

$$
\begin{aligned}
L_{\varepsilon}(\psi, k+1) & =A(\psi) L_{\varepsilon}(\theta, k)+B^{\prime}\left(\psi, L_{\varepsilon}(\theta, k), k\right) \\
\psi & =\theta+\varepsilon \bar{C}\left(M_{\varepsilon}(\theta, k), \theta, k\right)
\end{aligned}
$$

where, with $\otimes$ denoting the tensor product,

$$
\begin{aligned}
B^{\prime}(\psi, Y, k)= & -\varepsilon A(\psi) Y \nabla \bar{C}(Y, \theta, k)\left(I+\varepsilon \nabla \bar{C}\left(Y, \theta_{i} k\right)\right)^{-1} \\
& +\frac{\partial A}{\partial \theta}(\psi) \otimes M_{\imath}(\theta, k)+\frac{\partial B}{\partial \psi}(\psi) \otimes u_{k} \\
\nabla \bar{C}(Y, \theta, k)= & \frac{\partial \bar{C}}{\partial \theta}\left(M_{\varepsilon}(\theta, k), \theta, k\right)+\frac{\partial \bar{C}}{\partial X}\left(M_{\varepsilon}(\theta, k), \theta, k\right) Y
\end{aligned}
$$

Hence $L_{\varepsilon}$ should be the fixed point of the following operator:

$$
\begin{aligned}
T(L)(\psi, k+1) & =A(\psi) L(\theta, k)+B^{\prime}(\psi, L(\theta, k), k) \\
\theta & =\psi-\varepsilon D\left(M_{\varepsilon}, k\right)(\psi)
\end{aligned}
$$

Let us prove that this operator has a fixed point. From its definition and assumption A5, $B^{\prime}$ satisfies the assumptions of Lemma 2 with, in particular,

$$
\|X\|_{\theta} \leq m_{1} \Rightarrow\left\|B^{\prime}(\psi, X, k)\right\|_{\psi} \leq \varepsilon \gamma(1+\varepsilon \beta \bar{c}) m_{1} \frac{\bar{c}_{1}\left(1+m_{1}\right)}{1-\varepsilon \bar{c}_{1}\left(1+m_{1}\right)}+\alpha\left(a_{1} m+b_{1} u\right)
$$

Noting that $m_{1}$ plays the role of $m$ and $m_{2}$ (the Lipschitz constant of $L_{\varepsilon}$ ) the role of $m_{1}$, Lemma 2 applies if

$$
\begin{gathered}
\gamma(1+\varepsilon \beta \bar{c}(m)) \frac{m_{1}}{1-\varepsilon \bar{c}_{1}(m)\left(1+m_{1}\right)}+\alpha\left(a_{1} m+b_{1} u\right) \leq m_{1}, \\
(1+\varepsilon \beta \bar{c}(m)) \frac{\left(\gamma+\varepsilon b_{1}^{x}\right) m_{2}}{1-\varepsilon \bar{c}_{1}(m)\left(1+m_{1}\right)}+\alpha a_{1} m_{1}+b_{1}^{\psi} \leq m_{2}, \\
\tau_{1}^{\prime}(\varepsilon, m)=(1+\varepsilon \beta \bar{c}(m))\left(\gamma+\varepsilon b_{1}^{x}\right)\left(1+\varepsilon \frac{m_{2} \bar{c}_{1}(m)}{1-\varepsilon \bar{c}_{1}(m)\left(1+m_{1}\right)}\right)<1 .
\end{gathered}
$$

Note that the condition $\tau_{1}^{\prime}<1$ guarantees not only existence but also Lipschitz continuity of $L_{\varepsilon}$ (compare with (3)).

To complete our proof let us show that the candidate $L_{\varepsilon}$ is indeed the differential of $M_{\varepsilon}$, i.e., for any $\theta$ interior point of $S$ we have, uniformly in $k$,

$$
\limsup _{h \rightarrow 0} \frac{\left\|M_{e}(\theta+h, k)-M_{\varepsilon}(\theta, k)-L_{\varepsilon}(\theta, k) h\right\|}{\|h\|}=0 .
$$

To prove this equality, we consider $\psi$ an interior point of $S$ and $\bar{h}$ a (sufficiently
small) vector in $\mathbb{R}^{p}$ and we introduce the following notations (see assumption A5):

$$
\begin{aligned}
\theta & =\psi-\varepsilon D\left(M_{\varepsilon}, k\right)(\psi), \\
h & =\bar{h}+\varepsilon\left[D\left(M_{\varepsilon}, k\right)(\psi)-D\left(M_{\varepsilon}, k\right)(\psi+\bar{h})\right], \\
\delta_{M}(\theta, k, h) & =M_{\varepsilon}(\theta+h, k)-M_{\varepsilon}(\theta, k)-L_{\varepsilon}(\theta, k) h, \\
\bar{\delta}_{C}(\theta, k, h) & =\delta_{\bar{C}}\left(M_{\varepsilon}(\theta, k), \theta, k, M_{\varepsilon}(\theta+h)-M_{\varepsilon}(\theta, k), h\right), \\
\Delta(\theta, k) & =\limsup _{h \rightarrow 0} \frac{\left\|\delta_{M}(\theta, k, h)\right\|_{\theta}}{\|h\|} .
\end{aligned}
$$

Note that $\theta$ is also an interior point of $S$. From assumption A5 and since $M_{\varepsilon}$ is Lipschitz continuous, we have

$$
\limsup _{h \rightarrow 0} \frac{\left\|\bar{\delta}_{C}(\theta, k, h)\right\|}{\|h\|}=0
$$

Let us find a recurrence satisfied by $\delta_{M}$. The definition of $M_{z}$ yields

$$
\begin{aligned}
\delta_{M}(\psi, k+1, \bar{h})= & A(\psi+\bar{h}) \delta_{M}(\theta, k, h)+A(\psi+\bar{h}) L_{\varepsilon}(\theta, k) h-L_{\varepsilon}(\psi, k+1) \bar{h} \\
& +(A(\psi+\bar{h})-A(\psi)) M_{\varepsilon}(\theta, k)+(B(\psi+\bar{h})-B(\psi)) u_{k} .
\end{aligned}
$$

We also have

$$
\begin{aligned}
\bar{h}= & {\left[I+\varepsilon \nabla \bar{C}\left(L_{\varepsilon}(\theta, k), \theta, k\right)\right] h+\varepsilon \frac{\partial \bar{C}}{\partial X}\left(M_{\varepsilon}(\theta, k), \theta, k\right) \delta_{M}(\theta, k, h) } \\
& +\varepsilon \bar{\delta}_{c}(\theta, k, h), \\
L_{\varepsilon}(\psi, k+1)= & A(\psi) L_{\varepsilon}(\theta, k)\left[I+\varepsilon \nabla \bar{C}\left(L_{\varepsilon}(\theta, k), \theta, k\right)\right]^{-1} \\
& +\frac{\partial A}{\partial \psi}(\psi) \otimes M_{\varepsilon}(\theta, k)+\frac{\partial B}{\partial \psi}(\psi) \otimes u_{k} .
\end{aligned}
$$

This yields

$$
\begin{aligned}
& \delta_{M}(\psi, k+1, \bar{h}) \\
&= A(\psi+\bar{h}) \delta_{M}(\theta, k, h)+\delta_{A}(\psi, k, \bar{h}) M_{\varepsilon}(\theta, k)+\delta_{B}(\psi, k, \bar{h}) u_{k} \\
&+(A(\psi+\bar{h})-A(\psi)) L_{\varepsilon}(\theta, k) h-\varepsilon A(\psi) L_{\varepsilon}(\theta, k) \\
& ? \\
& \times\left(I+\varepsilon \nabla \bar{C}\left(M_{\varepsilon}(\theta, k), \theta, k\right)\right)^{-1}\left[\frac{\partial \bar{C}}{\partial X}\left(M_{\varepsilon}(\theta, k), \theta, k\right) \delta_{M}(\theta, k, h)+\bar{\delta}_{C}(\theta, k, h)\right] .
\end{aligned}
$$

Noticing that $\|h\| /\left(1+\varepsilon d_{1}\right) \leq\|\bar{h}\| \leq\left(1+\varepsilon \bar{c}_{1}\left(1+m_{1}\right)\right)\|h\|$, taking the $\|\cdot\| \|$-norm, dividing by $\|h\|$, and taking the lim sup for $h$ going to zero, we obtain

$$
\Delta(\psi, k+1) \frac{1}{1+\varepsilon d_{1}} \leq \gamma(1+\varepsilon \beta \bar{c})\left(1+\varepsilon \frac{m_{1} \bar{c}_{1}}{1-\varepsilon \bar{c}_{1}\left(1+m_{1}\right)}\right) \Delta(\theta, k) .
$$

On the other hand, for any interior point $\theta$ of S , the properties of $M_{\varepsilon}, L_{\varepsilon}$ give, for all $k, \Delta(\theta, k) \leq 2 m_{1}$. Given $k$ and $\theta$, an interior point of $S$, we construct a sequence
$\left(\theta_{j}\right), j \in \mathbb{N}$, as $\theta_{j+1}=\theta_{j}-\varepsilon D\left(M_{e}, k-j-1\right)\left(\theta_{j}\right), \theta_{0}=\theta$. By induction, $\theta_{j}$ is an interior point of $S$ and, from the previous inequality, we have $\Delta(\theta, k) \leq \tau_{1}(\varepsilon, m)^{j} \Delta\left(\theta_{j}, k-j\right) \leq$ $\tau_{1}(\varepsilon, m)^{j} 2 m_{1}$ with

$$
\begin{aligned}
\tau_{1}(\varepsilon, m) & =\gamma(1+\varepsilon \beta \bar{c}(m))\left(1+\varepsilon \frac{m_{1} \bar{c}_{1}(m)}{1-\varepsilon \bar{c}_{1}(m)\left(1+m_{1}\right)}\right) \frac{1}{1-\varepsilon \bar{c}_{1}(m)\left(1+m_{1}\right)} \\
& =\tau_{0}(\varepsilon, m) \frac{1}{1-\varepsilon \bar{c}_{1}(m)\left(1+m_{1}\right)} .
\end{aligned}
$$

Our relation holds for all positive $j$, and, as expected from the previous Remark, we have established that $L_{z}$ is the differential of $M_{\varepsilon}$ if $\tau_{1}(\varepsilon, m)<1$.

Approximation. We have the existence of $M_{z}$, the solution of

$$
M_{\varepsilon}(\psi, k+1)=A(\psi) M_{\varepsilon}\left(\psi-\varepsilon D\left(M_{\varepsilon}, k\right)(\psi), k\right)+B(\psi) u_{k} .
$$

Getting a solution of this system for each $\psi$, would be equivalent to getting all the solutions of $\bar{S}_{\varepsilon}$ bounded on $\mathbb{Z}$. However, we notice that this equation can be rewritten as

$$
\begin{aligned}
M_{\varepsilon}(\psi, k+1)= & A(\psi) M_{\varepsilon}(\psi, k)+B(\psi) u_{k} \\
& +A(\psi)\left[M_{\varepsilon}\left(\psi-\varepsilon D\left(M_{\varepsilon}, k\right)(\psi), k\right)-M_{\varepsilon}(\psi, k)\right] .
\end{aligned}
$$

This is a linear system with a nonlinear forcing term which disappears for $\varepsilon=0$. Using the Poincaré method of expansion with respect to a small parameter, we can obtain a family of approximations:
" 0 "-order approximation: taking $\varepsilon=0$ gives

$$
M_{0}(\psi, k+1)=A(\psi) M_{0}(\psi, k)+B(\psi) u_{k}
$$

This is the solution of $S_{0}$.
" 1 "-order approximation: retaining the first-order term in $\varepsilon$ gives

$$
M_{1}(\psi, k+1)=A(\psi) M_{1}(\psi, k)+B(\psi) u_{k}-\varepsilon A(\psi) L_{0}(\psi, k) \bar{C}\left(M_{0}(\psi, k), \psi, k\right)
$$

again a linear system.
The quality of these approximations can be characterized as follows:
(i) We notice that $M_{0}$ belongs to the set $\mathbf{M}$, hence the graph transform properties imply

$$
\left|M_{0}-M_{\varepsilon}\right| \leq \frac{\left|M_{0}-T\left(M_{0}\right)\right|}{1-\tau_{0}(\varepsilon, m)}
$$

But
$\left\|M_{0}(\psi, k+1)-T\left(M_{0}\right)(\psi, k+1)\right\|_{\psi} \leq \gamma\left\|M_{0}\left(\psi-\varepsilon D\left(M_{0}, k\right)(\psi), k\right)-M_{0}(\psi, k)\right\|_{\psi}$.
It follows that

$$
\left|M_{0}-M_{\varepsilon}\right| \leq \varepsilon \frac{\gamma m_{1} \bar{c}}{1-\tau_{0}(\varepsilon, m)}=\varepsilon f^{0}
$$

Similarly, for the gradient matrix, we get

$$
\left|L_{0}-L_{\varepsilon}\right| \leq \frac{\left|L_{0}-T\left(L_{0}\right)\right|}{1-\tau_{1}^{\prime}(\varepsilon, m)}
$$

and

$$
\begin{aligned}
& L_{0}(\psi, k+1)-T\left(L_{0}\right)(\psi, k+1) \\
& \quad=A(\psi)\left[L_{0}(\psi, k)-L_{0}(\theta, k)\left(I+\varepsilon \nabla \bar{C}\left(L_{0}(\theta, k), \theta, k\right)\right)^{-1}\right] \\
& \quad+\frac{\partial A}{\partial \psi}(\psi) \otimes\left(M_{0}(\psi, k)-M_{\varepsilon}(\theta, k)\right)
\end{aligned}
$$

with $\theta=\psi-\varepsilon D\left(M_{0}, k\right)(\psi)$. Hence

$$
\left|L_{0}-L_{\varepsilon}\right| \leq \frac{\gamma\left(\varepsilon m_{2} \bar{c}+\varepsilon \frac{m_{1} \bar{c}_{1}\left(1+m_{1}\right)(1+\varepsilon \beta \bar{c})}{1-\varepsilon \bar{c}_{1}\left(1+m_{1}\right)}\right)+\alpha a_{1} m_{1}\left(\varepsilon m_{1} \bar{c}+\varepsilon f^{0}\right)}{1-\tau_{1}^{\prime}(\varepsilon, m)}=\varepsilon f_{1}^{0}
$$

(ii) The function $M_{0}^{1}$, defined in (iii) of Theorem 1 , is given by a noncritical linear system with bounded input. Hence, it is uniquely defined and bounded. Denoting $\Delta(\psi, k)=M_{e}(\psi, k)-M_{0}^{1}(\psi, k)$, we have

$$
\begin{aligned}
\Delta(\psi, k+1) & =A(\psi)\left[\Delta(\psi, k)+M_{\varepsilon}(\theta, k)-M_{\varepsilon}(\psi, k)+\varepsilon L_{0}(\psi, k) \bar{C}\left(M_{0}(\psi, k), \psi, k\right)\right], \\
\theta & =\psi-\varepsilon D\left(M_{\varepsilon}, k\right)(\psi) .
\end{aligned}
$$

We know that the segment $[\theta, \psi]$ is contained in $\mathbf{S}$. The Mean Value Theorem gives for some $\xi, 0 \leq \xi \leq 1$,

$$
M_{\varepsilon}(\theta, k)-M_{\varepsilon}(\psi, k)=-\varepsilon L_{\varepsilon}(\psi+\xi(\theta-\psi), k) \bar{C}\left(M_{\varepsilon}(\theta, k), \theta, k\right) .
$$

Hence, using the properties of $L_{e}, L_{0}-L_{v}, M_{0}-M_{\varepsilon}, D$, we get

$$
\begin{aligned}
& \left\|M_{\varepsilon}(\theta, k)-M_{\varepsilon}(\psi, k)+\varepsilon L_{0}(\psi, k) \bar{C}\left(M_{0}(\psi, k), \psi, k\right)\right\|_{\psi} \\
& \quad \leq \varepsilon^{2} m_{2} \xi \bar{c}^{2}+\varepsilon^{2} f_{1}^{0} \bar{c}+\varepsilon m_{1} \bar{c}_{1}\left(\varepsilon\left(1+m_{1}\right) \bar{c}+\varepsilon f^{0}\right) .
\end{aligned}
$$

This yields

$$
\left|M_{\varepsilon}-M_{1}\right| \leq \varepsilon^{2} \gamma \frac{m_{2} \bar{c}^{2}+f_{1}^{0} \bar{c}+m_{1} \bar{c}_{1}\left(\bar{c}\left(1+m_{1}\right)+f^{0}\right)}{1-\gamma}=\varepsilon^{2} f^{1} .
$$

Attractiveness. Let $(X, \theta, k)$ be an element of $B(O, x) \times S_{0} \times \mathbb{Z}$, for $x \geq m$, and let $(Y, \psi, k+1)$ be its image by $S_{\varepsilon}$. Paying attention to the fact that we have $X$ only in $B(O, x)$ and not in $B(O, m)$, we can however proceed, as in Lemma 1, to get

$$
\begin{aligned}
\left\|D\left(M_{\varepsilon}, k\right)(\psi)-C(X, \theta, k)\right\| & \leq d_{1}^{m}(x)\left\|X-M_{\varepsilon}(\theta, k)\right\|_{\theta}, \\
d_{1}^{m}(x) & =\frac{\bar{c}_{1}(x)}{1-\varepsilon \bar{c}_{1}(x)\left(1+m_{1}\right)} .
\end{aligned}
$$

Then, as in Lemma 2, we obtain $\left\|Y-M_{\varepsilon}(\psi, k+i)\right\|_{\psi} \leq \tau_{0}(\varepsilon, x)\left\|X-M_{\varepsilon}(\theta, k)\right\|_{\theta}$.

Applying this inequality to a sequence $\left(X_{k+1}, \theta_{k}\right)$ of $\mathbf{B}\left(\left[k_{0}, k_{1}\right), \mathbf{S}_{0}\right)$, we obtain for all $k, k^{\prime}, k_{0} \leq k^{\prime} \leq k<k_{1}$,

$$
\left\|X_{k+1}-M_{\varepsilon}\left(\theta_{k}, k+1\right)\right\| \leq \tau_{0}(\varepsilon, x)^{k-k^{\prime}} \alpha\left\|X_{k^{\prime}+1}-M_{\varepsilon}\left(\theta_{k^{\prime}}, k^{\prime}+1\right)\right\| .
$$

Moreover, if $X_{k+1}, \theta_{k}$ lies in $B(O, x) \times \mathrm{S}_{0}$ for all $k$, then this inequality holds for all $k^{\prime}, k^{\prime} \leq k$. This implies $X_{k+1}=M_{e}\left(\theta_{k}, k+1\right)$.

This completes the proof of Theorem 1.
In the following we are interested in characterizing the set $\mathbf{B}\left(\left[k_{0}, \infty\right), \mathbf{S}_{0}\right)$. The attractiveness property implies that this set is a subset of $l_{\infty}\left(\mathbb{R}^{n} \times \mathbb{R}^{p}\right)$ from which it inherits a complete metric space property. We also see thatithe corresponding sequences of elevation above the manifold have a natural exponentially weighted supremum norm:

$$
|E|=\operatorname{Sup}_{k \geq k_{0}}\left\{\rho^{k_{0}-k}\left\|E_{k}\right\|\right\}, \quad \rho<1 .
$$

For $\rho=\tau_{0}(\varepsilon, x)$, we know that the map from $\mathbf{B}\left(\left[k_{0}, \infty\right), S_{0}\right)$ to $\mathbb{R}^{n}$, giving $(X, \theta) \mapsto E$ is bounded with this norm. We even have

Lemma 3. Under assumptions A1-A5, there exists $\varepsilon^{*}$ such that for any $\varepsilon, \varepsilon \leq \varepsilon^{*}$, we can find $\rho<1$ such that the above map is Lipschitz continuous.

A key technical lemma to prove this statement is the following consequence of Hadamard's lemma [AE]:

Lemma 4. Let $f$ be a Lipschitz continuously differentiable function $f: \mathbf{C} \subset \mathbb{R}^{n} \rightarrow \mathbb{R}^{n}$, with $f_{1}, f_{2}$ as Lipschitz constants. If the segments $\left[x^{0}, x^{0}+\delta^{0}\right],\left[x^{1}, x^{1}+\delta^{1}\right]$ are contained in C , then

$$
\begin{aligned}
& \left\|f\left(x^{0}+\delta^{0}\right)-f\left(x^{0}\right)+f\left(x^{1}\right)-f\left(x^{1}+\delta^{1}\right)\right\| \\
& \quad \leq\left(f_{1}+\frac{1}{2} f_{2}\left\|\delta^{0}\right\|\right)\left\|\delta^{0}-\delta^{1}\right\|+f_{2}\left\|\delta^{0}\right\|\left\|x^{0}-x^{1}\right\| .
\end{aligned}
$$

Proof of Lemma 3. Let $\left(X^{0}, \theta^{0}, k\right),\left(X^{1}, \theta^{1}, k\right)$ be two points in $B(O, x) \times \mathrm{S}_{0} \times \mathbb{Z}$ and $\left(Y^{0}, \psi^{0}, k\right),\left(Y^{1}, \psi^{1}, k\right)$ be the respective images by $S_{\varepsilon}$. For $i=0$, 1 , we define

$$
\begin{gathered}
E^{i}=X^{i}-M_{\varepsilon}\left(\theta^{i}, k\right), \quad F^{i}=Y^{i}-M_{\varepsilon}\left(\psi^{i}, k\right), \quad G^{i}=X^{i}-M_{\varepsilon}\left(\theta^{i}+\Delta^{i}, k\right) \\
\Delta^{i}=\psi^{i}-\theta^{i}-\varepsilon D\left(M_{\varepsilon}, k\right)\left(\psi^{i}\right)=\varepsilon\left[C\left(X^{i}, \theta^{i}, k\right)-C\left(M_{\varepsilon}\left(\theta^{i}+\Delta^{i}, k\right), \theta^{i}+\Delta^{i}, k\right)\right]
\end{gathered}
$$

We have $\left\|\Delta^{0}\right\| \leq \varepsilon c_{1}\left(\left\|G^{0}\right\|+\left\|\Delta^{0}\right\|\right)$ and $\left\|G^{0}\right\| \leq\left\|E^{0}\right\|+m_{1}\left\|\Delta^{0}\right\|$ which implies

$$
\left\|\Delta^{0}\right\| \leq \varepsilon \frac{c_{1}}{1-\varepsilon c_{1}\left(1+m_{1}\right)}\left\|E^{0}\right\|, \quad\left\|G^{0}\right\| \leq \frac{1-\varepsilon c_{1}}{1-\varepsilon c_{1}\left(1+m_{1}\right)}\left\|E^{0}\right\|
$$

On the other hand, by definition, $F^{0}-F^{1}=A\left(\psi^{0}\right)\left(G^{0}-G^{1}\right)+\left(A\left(\psi^{0}\right)-A\left(\psi^{1}\right)\right) G^{1}$. But, noticing that, for $i=0,1$, the segments $\left[X^{i}, M_{\varepsilon}\left(\theta^{i}+\Delta^{i}, k\right)\right]$ and $\left[\theta^{i}, \theta^{i}+\Delta^{i}\right]$ are
contained in $B(O, x)$ and $\mathrm{S}_{1}$, respectively, we apply Lemma 4 to $M_{\varepsilon}$ and $C$ to obtain

$$
\begin{aligned}
\left\|G^{0}-G^{1}\right\|_{\psi^{0}} \leq & \left\|E^{0}-E^{1}\right\| \psi^{0}+\left(m_{1}+\frac{1}{2} m_{2}\left\|\Delta^{0}\right\|\right)\left\|\Delta^{0}-\Delta^{1}\right\| \\
\quad & \quad+m_{2}\left\|\Delta^{0}\right\|\left\|\theta^{0}-\theta^{1}\right\|, \\
\left\|\Delta^{0}-\Delta^{1}\right\| \leq & \varepsilon\left[\left(c_{1}+\frac{1}{2} c_{2}\left(\left\|\Delta^{0}\right\|+\left\|G^{0}\right\|\right)\right)\left(\left\|G^{0}-G^{1}\right\|+\left\|\Delta^{0}-\Delta^{1}\right\|\right)\right. \\
& \left.+c_{2}\left(\left\|\Delta^{0}\right\|+\left\|G^{0}\right\|\right)\left(\left\|X^{0}-X^{1}\right\|+\left\|\theta^{0}-\theta^{1}\right\|\right)\right] .
\end{aligned}
$$

It follows that

$$
\begin{aligned}
& {\left[1-\varepsilon\left(c_{1}+\frac{1}{2} c_{2}\left(\left\|\Delta^{0}\right\|+\left\|G^{0}\right\|\right)\right)\left(1+m_{1}+\frac{1}{2} m_{2}\left\|\Delta^{0}\right\|\right)\right]\left\|G^{0}-G^{1}\right\| \psi^{0} } \\
& \leq {\left[1-\varepsilon\left(c_{1}+\frac{1}{2} c_{2}\left(\left\|\Delta^{0}\right\|+\left\|G^{0}\right\|\right)\right)\right]\left\|\left\|E^{0}-E^{1}\right\| \psi^{0}\right.} \\
&+\left[\left(1-\varepsilon c_{1}\right) m_{2}\left\|\Delta^{0}\right\|+\varepsilon c_{2} m_{1}\left(\left\|\Delta^{0}\right\|+\left\|G^{0}\right\|\right)\right]\left\|\theta^{0}-\theta^{1}\right\| \\
&+\varepsilon c_{2}\left(\left\|\Delta^{0}\right\|+\left\|G^{0}\right\|\right)\left(m_{1}+\frac{1}{2} m_{2}\left\|\Delta^{0}\right\|\right)\left\|X^{0}-X^{3}\right\| .
\end{aligned}
$$

This implies existence of a constant $q_{0}$ such that

$$
\left\|F^{0}-F^{1}\right\|_{\psi^{0}} \leq \rho\left\|E^{0}-E^{1}\right\|_{\theta_{0}}+q_{0}\left\|E^{0}\right\|\left(\left\|X^{0}-X^{1}\right\|+\left\|\theta^{0}-\theta^{1}\right\|\right)
$$

with

$$
\rho=\frac{\gamma(1+\varepsilon \beta c)\left[1-\varepsilon\left(c_{1}+\left(\frac{1}{2} c_{2}\left(1-\varepsilon c_{1}\left(1+m_{1}\right)\right)\right)\left\|E^{0}\right\|\right)\right]}{1-\varepsilon\left(c_{1}+\left(\frac{1}{2} c_{2}\left(1-\varepsilon c_{1}\left(1+m_{1}\right)\right)\right)\left\|E^{0}\right\|\right)\left(1+m_{1}+\varepsilon \frac{1}{2} m_{2} c_{1} /\left(1-\varepsilon c_{1}\left(1+m_{1}\right)\right)\left\|E^{0}\right\|\right)} .
$$

Now we apply these inequalities to two sequences $\left(X_{k+1}^{0}, \theta_{k}^{0}\right),\left(X_{k+1}^{1}, \theta_{k}^{1}\right)$ of $\mathbf{B}\left(\left[k_{0}, \infty\right), \mathbf{S}_{0}\right)$ : we have

$$
\left\|E_{k}^{0}\right\| \leq(x-m) \tau_{0}^{k-k_{0}}, \quad\left\|E_{k_{0}+1}^{0}-E_{k_{0}+1}^{1}\right\|\left\|_{k_{0}} \leq \alpha\right\| X_{k_{0}+1}^{0}-X_{k_{0}+1}^{1}\left\|+m_{1}\right\| \theta_{k_{0}}^{0}-\theta_{k_{0}}^{1} \|,
$$

hence

$$
\left\|E_{k+1}^{0}-E_{k+1}^{1}\right\|_{\rho_{k}^{0}} \leq \rho\left\|E_{k}^{0}-E_{k}^{1}\right\|_{\theta_{k-1}^{0}}+q_{0}(x-m) \tau_{0}^{k-k_{0}}\left(\left\|X_{k}^{0}-X_{k}^{1}\right\|+\left\|\theta_{k}^{0}-\theta_{k}^{1}\right\|\right),
$$

which yields

$$
\left\|E_{k+1}^{0}-E_{k+1}^{1}\right\| \leq \rho^{k-k_{0}}\left(1+\frac{\tau_{0}}{\rho-\tau_{0}}\right) q \operatorname{Sup}_{k \geq k_{0}}\left(\left\|X_{k+1}^{0}-X_{k+1}^{1}\right\|+\left\|\theta_{k}^{0}-\theta_{k}^{1}\right\|\right)
$$

with $q=\operatorname{Max}\left(\alpha, m_{1}, q_{0}(x-m)\right)$.

## 4. Topological Orbital Equivalence with Asymptotic Phase

In this section we establish the existence of a homeomorphism between subsets of $\mathbf{B}\left(\left[k_{0}, \infty\right), \mathbf{S}_{0}\right)$ and $\mathbf{B}^{M}\left(\left[k_{0}, \infty\right), \mathbf{S}_{0}\right)$. For such a strong property to hold, solutions of $S_{\varepsilon}$ must be very close to solutions of $S M_{\varepsilon}$. As a consequence of attractiveness, we know that for any sequence $\left(X_{k+1}, \theta_{k}\right)$ (resp. $\left(X_{k+1}^{M}, \theta_{k}^{M}\right)$ ) of $\mathbf{B}\left(\left[k_{0}, \infty\right), \mathbf{S}_{0}\right)$ (resp. ( $\left.\mathbf{B}^{M}\left(\left[k_{0}, \infty\right), \mathbf{S}_{0}\right)\right)$ we have for all $k, k \geq k_{0}$,

$$
\begin{aligned}
\rho^{k_{0}-k}\left\|\left(X_{k+1}, \theta_{k}, k+1\right)-S M_{\varepsilon}\left(X_{k}, \theta_{k-1}, k\right)\right\| & \leq \varepsilon\left(1+\left(a_{1} x+b_{1} u\right)\right) c_{1} \alpha(x-m), \\
\rho & =\tau_{0}(\varepsilon, x),
\end{aligned}
$$

and

$$
\begin{aligned}
\rho^{k_{0}-k}\left\|\left(X_{k+1}^{M}, \theta_{k}^{M}, k+1\right)-S_{\varepsilon}\left(X_{k}^{M}, \theta_{k-1}^{M}, k\right)\right\| & \leq \varepsilon\left(1+\left(a_{1} x+b_{1} u\right)\right) c_{1} \alpha(x-m) \\
\rho & =\gamma(1+\varepsilon \beta c)
\end{aligned}
$$

This means that $\left(X_{k+1}, \theta_{k}\right)$ (resp. $\left(X_{k+1}^{M}, \theta_{k}^{M}\right)$ ) is an $\varepsilon$-pseudo solution of $S M_{\varepsilon}$ (resp. $S_{\varepsilon}$ ) in the neighborhood of a normally attractive integral manifold. Invoking the Shadowing Lemma (Proposition 8.19 of [S]) we may expect the existence of a unique solution $\left(\tilde{X}_{k+1}^{M}, \tilde{\theta}_{k}^{M}\right)$ (resp. $\left(\tilde{X}_{k+1}, \tilde{\theta}_{k}\right)$ ) of $S M_{z}$ (resp. $S_{\varepsilon}$ ) $\varepsilon$-close to the corresponding $\varepsilon$-pseudo solution for the exponentially weighted supremum norm. This allows us to define a map $\Phi\left(k_{0}\right)\left(\mathrm{resp} . \Phi^{M}\left(k_{0}\right)\right)$ as $\Phi\left(k_{0}\right)(X, \theta)_{k}=\left(\widetilde{X}_{k+1}^{M}, \tilde{\theta}_{k}^{M}\right)$ (resp. $\Phi^{M}\left(k_{0}\right)\left(X^{M}, \theta^{M}\right)_{k}=\left(\tilde{X}_{k+1}, \tilde{\theta}_{k}\right)$. Clearly, from the $\varepsilon$-closeness property, this map is a very good candidate for the homecmorphism we are looking for.

To obtain the Shadowing Lemma in our context, we have to shrink the set $\mathrm{S}_{0}$ again: let $\mathbf{S}_{0}^{\prime}, \mathbf{S}_{1}^{\prime}$ be compact sets with nonempty interior such that $\mathbf{S}_{0}^{\prime}+\eta \subset \mathbf{S}_{1}^{\prime}$, $\mathbf{S}_{1}^{\prime}+\eta \subset \mathbf{S}_{0}$.

Theorem 2. Under assumptions A1-A4, for any compact set $\mathrm{S}_{1}^{\prime}$ strictly contained in $\mathbf{S}$, there exists $\varepsilon^{*}$ such that for any $\varepsilon, \varepsilon \leq \varepsilon^{*}$, we can find $g^{x}, g^{\theta}, g_{1}, \sigma, 0 \leq \sigma<1$, and, for any $k_{0}$, maps

$$
\begin{aligned}
\Phi\left(k_{0}\right): \mathbf{B}\left(\left[k_{0}, \infty\right), \mathbf{S}_{1}^{\prime}\right) \rightarrow \mathbf{B}^{M}\left(\left[k_{0}, \infty\right), \mathbf{S}_{0}\right) \text { and } \\
\Phi^{M}\left(k_{0}\right): \mathbf{B}^{M}\left(\left[k_{0}, \infty\right), \mathbf{S}_{1}^{\prime}\right) \rightarrow \mathbf{B}\left(\left[k_{0}, \infty\right), \mathbf{S}_{0}\right)
\end{aligned}
$$

such that (similarly for $\Phi^{M}$ ):
(i) For any sequence $\left(X_{k+1}, \theta_{k}\right)$ in $\mathbf{B}\left(\left[k_{0}, \infty\right), \mathbf{S}_{1}^{\prime}\right)$, its image $\left(\tilde{X}_{k+1}^{M}, \tilde{\theta}_{k}^{M}\right)$ by $\Phi\left(k_{0}\right)$ satisfies for all $k, k \geq k_{0}$,

$$
\begin{gathered}
X_{k_{0}+1}=\widetilde{X}_{k_{0}+1}^{M}, \quad \sigma^{k_{0}-k}\left\|X_{k+1}-\widetilde{X}_{k+1}^{M}\right\| \leq \varepsilon g^{x}\left(\left\|E_{k_{0}+1}\right\|\right), \\
\sigma^{k_{0}-k}\left\|\theta_{k}-\tilde{\theta}_{k}^{M}\right\| \leq \varepsilon g^{\theta}\left(\left\|E_{k_{0}+1}\right\|\right),
\end{gathered}
$$

where $g^{x}, g^{\theta}$ are positive nondecreasing functions of the norm of the elevation above the manifold at time $k_{0}$.
(ii) $\Phi\left(k_{0}\right)(X, \theta)_{k}=\Phi\left(k_{0}^{\prime}\right)(X, \theta)_{k}, k_{0} \leq k_{0}^{\prime} \leq k$.
(iii) If assumption A5 holds, then $\Phi\left(k_{0}\right)$ is Lipschitz continuous in the following strong sense: for all $k, k \geq k_{0}$,

$$
\begin{aligned}
& \left\|\tilde{X}_{k+1}^{0}-\tilde{X}_{k+1}^{1}\right\|+\left\|\tilde{\theta}_{k}^{0}-\tilde{\theta}_{k}^{1}\right\| \\
& \quad \leq \cdot\left\|X_{k+1}^{0}-X_{k+1}^{1}\right\|+\left\|\theta_{k}^{0}-\theta_{k}^{1}\right\|+\sigma^{k-k_{0}} g_{1} \operatorname{Sup}_{\kappa \geq k_{0}}\left(\left\|X_{\kappa+1}^{0}-X_{\kappa+1}^{1}\right\|+\left\|\theta_{\kappa}^{0}-\theta_{\kappa}^{1}\right\|\right) .
\end{aligned}
$$

Moreover, $\Phi\left(k_{0}\right)\left(\right.$ resp. $\left.\Phi^{M}\left(k_{0}\right)\right)$ is injective and the restriction of $\Phi^{M}\left(k_{0}\right) \circ$ $\Phi\left(k_{0}\right)\left(r e s p . \Phi\left(k_{0}\right) \circ \Phi^{M}\left(k_{0}\right)\right)$ to $\mathbf{B}\left(\left[k_{0}, \infty\right), \mathbf{S}_{0}^{\prime}\right)\left(r e s p . \mathbf{B}^{M}\left(\left[k_{0}, \infty\right), \mathbf{S}_{0}^{\prime}\right)\right)$ is the identity map.

Remarks. For the continuous time case, Riedle and Kokotovic [RK2] have shown that if there exists a uniformly stable solution $\left(\theta_{k}^{M}\right)$ of $S M_{\varepsilon}^{\theta}$, remaining in $\mathbf{S}_{1}^{\prime}$ after time $k_{0}$, then $\mathbf{B}\left(\left[k_{0}, \infty\right), S_{0}\right)$ is not empty and has a nonempty interior.

Since we have only $\sigma^{k_{0}-k}\left\|\theta_{k}-\tilde{\theta}_{k}\right\| \leq \eta$ we have to take $\theta_{k}$ in $\mathbf{S}_{0}^{\prime}$ (resp. $\mathbf{S}_{1}^{\prime}$ ) to guarantee that $\tilde{\theta}_{k}^{M}$ is in $\mathbf{S}_{1}^{\prime}$ (resp. $\mathrm{S}_{0}$ ).

Since we have an exponentially decaying distance between ( $X_{k+1}, \theta_{k}$ ) and $\left(\tilde{X}_{k+1}^{M}, \tilde{\theta}_{k}^{M}\right), \Phi\left(k_{0}\right)$ is an asymptotic phase.

We have established $\mathbf{B}\left(\left[k_{0}, \infty\right), \mathbf{S}_{0}^{\prime}\right) \subset \Phi^{M}\left(k_{0}\right)\left(\mathbf{B}^{M}\left(\left[k_{0}, \infty\right), \mathbf{S}_{1}^{\prime}\right)\right)$. This means that any solution ( $X_{k+1}, \theta_{k}$ ) of $S_{c}$ which remains in $B(O, x) \times S_{0}^{\prime}$ on a semi-infinite time interval (i.e., belongs to $\mathbf{B}\left(\left[k_{0}, \infty\right), \mathbf{S}_{0}^{\prime}\right)$ for some $k_{0}$ ) can be approximated with an exponentially decaying distance by a solution of $S M_{t}$ satisfying the same property. Moreover, these two sequences have the same $X$-values $X_{k_{0}+1}$ at time $k_{0}$ and their $\theta$-values at this time are at an $\varepsilon$-distance, the magnitude of this distance increasing with the norm of the elevation above the manifold at time $k_{0}$. These two solutions have the same type of Lyapunov stability in each of the following cases: stability, uniform stability, asymptotic stability, uniform asymptotic stability, and instability. Unformity follows from (ii) and the independence of $\sigma, g^{x}, g^{\theta}, g_{1}$ in $k_{0}$. Asymptotic or exponential property results from the $\sigma^{k-k_{0}}$ term in (iii). (In)stability is a consequence of the continuity property (iii) as shown in the following lemma:

Lemma 5. If the sequence $\left(X_{k+1}, \theta_{k}\right)$ (resp. $\left(X_{k+1}^{M}, \theta_{k}^{M}\right)$ ) in the interior (in the $l_{\infty}\left(\mathbb{R}^{n} \times \mathbb{R}^{p}\right)$ sense $)$ of $\mathbf{B}\left(\left[k_{0}, \infty\right), \mathbf{S}_{0}^{\prime}\right)\left(\right.$ resp. $\left.\mathbf{B}^{M}\left(\left[k_{0}, \infty\right), \mathbf{S}_{0}^{\prime}\right)\right)$ is a stable solution of $S_{\varepsilon}\left(\right.$ resp. $\left.S M_{\varepsilon}\right)$, its image by $\Phi\left(k_{0}\right)\left(\right.$ resp. $\left.\Phi^{M}\left(k_{0}\right)\right)$ is a stable solution of $S M_{\varepsilon}$ (resp. $S_{\varepsilon}$ ).

Proof. As above let ( $\tilde{X}_{k+1}^{M}, \tilde{\theta}_{k}^{M}$ ) denote the image of $\left(X_{k+1}, \theta_{k}\right)$ by $\Phi\left(k_{0}\right)$. Let ( $X_{k+1}^{0}, \theta_{k}^{0}$ ) be the stable solution of $S_{\varepsilon}$ in the interior of $\mathbf{B}\left(\left[k_{0}, \infty\right), \mathbf{S}_{0}^{\prime}\right)$. Stability implies existence of an open neighborhood $\mathbf{V}\left(X_{k_{0}+1}^{0}, \theta_{k_{0}}^{0}\right)$ in $\mathbb{R}^{n} \times \mathbb{R}^{p}$ on which the injective map associating a solution of $S_{\varepsilon}$ to its value at time $k_{0}$ is continuous. Composing by $\Phi\left(k_{0}\right)$, we obtain a continuous injective map:

$$
\begin{gathered}
\Psi: \mathbf{V}\left(X_{k_{0}+1}^{0}, \theta_{k_{0}}^{0}\right) \rightarrow \mathbf{B}^{M}\left(\left[k_{0}, \infty\right), \mathbf{S}_{1}^{\prime}\right), \\
\left(X_{k_{0}+1}, \theta_{k_{k_{0}}}\right) \mapsto\left(\tilde{X}_{k+1}^{M}, \tilde{\theta}_{k}^{M}\right) .
\end{gathered}
$$

Restricting $\Psi$ to time $k=k_{0}$, we have a continuous injective map:

$$
\begin{aligned}
& \Psi_{k_{0}}: V\left(X_{k_{0}+1}^{0}, \theta_{k_{0}}^{0}\right) \subset \mathbb{R}^{n} \times \mathbb{R}^{p} \rightarrow \mathbb{R}^{n} \times \mathbb{R}^{p}, \\
&\left(X_{k_{0}+1}, \theta_{k_{0}}\right) \mapsto\left(\tilde{X}_{k_{0}+1}^{M}, \tilde{\theta}_{k_{0}}^{M}\right) .
\end{aligned}
$$

It follows from Brouwer's Domain Invariance Theorem (Theorem 3.3.2 of [L2]) that this application is a homeomorphism from the open set $\mathbf{V}\left(X_{k_{0}+1}^{0}, \theta_{k_{0}}^{0}\right)$ into its open image $\overline{\mathbb{V}}\left(\widetilde{X}_{k_{0}+1}^{M 0}, \tilde{\theta}_{k_{0}}^{M 0}\right)$. We conclude that the following map is continuous and is defined on a nonempty open subset of $\mathbb{R}^{n} \times \mathbb{R}^{p}$ :

$$
\Psi \circ \Psi_{k_{0}}^{-1}: \mathbf{V}\left(\tilde{X}_{k_{0}+1}^{M 0}, \tilde{\theta}_{k_{0}}^{M 0}\right) \rightarrow \mathbf{B}^{M}\left(\left[k_{0}, \infty\right), \mathbf{S}_{1}^{\prime}\right) .
$$

This is nothing but the map associating the solution $\left(\tilde{X}_{k+1}^{M}, \tilde{\theta}_{k}^{M}\right)$ of $S M_{\varepsilon}$ to its value $\left(\widetilde{X}_{k_{0}+1}^{M}, \tilde{\theta}_{k_{0}}^{M}\right)$ at time $k_{0}$. Since this can be done for any $k_{0}^{\prime}, k_{0}^{\prime} \geq k_{0}$, we have proved stability of $\left(\tilde{X}_{k+1}^{M 0}, \tilde{\theta}_{k}^{M 0}\right)$.

From this lemma, instability is obtained by contradiction.
Noticing that if $\left(X_{k+1}^{M}, \theta_{k}^{M}\right)$ belongs to $\mathbf{B}^{M}\left(\left[k_{0}, \infty\right), \mathbf{S}_{0}\right)$, then this solution of $S M_{\varepsilon}$ and $\left(\theta_{k}^{M}\right)$ as a solution of $S M_{\varepsilon}^{\theta}$ have the same type of Lyapunov stability or instability, we have established the main result of this paper:

Main result. Under assumptions A1-A5, for each $x, x \geq m$, there exists $\varepsilon^{*}$ such that for any $\varepsilon, \varepsilon \leq \varepsilon^{*}$, the systems $S_{\varepsilon}$ and $S M_{\varepsilon}$, when restricted to $B(O, x) \times S_{0}^{\prime}$, are topologically orbitally equivalent with asymptotic phase. More precisely, for any $k_{0}$, each solution $\left(X_{k+1}, \theta_{k}\right)$ of $S_{\varepsilon}$ such that $\theta_{k} \in S_{0}^{\prime}$ for all $k, k \geq k_{0}$, and $\| X_{k_{0}+1}-$ $M_{\varepsilon}\left(\theta_{k_{0}}, k_{0}+1\right) \| \leq(x-m) / \alpha$, can be obtained from a solution $\left(\theta_{k}^{M}\right)$ of $S M_{\varepsilon}^{\theta}$ satisfying $\theta_{k}^{M} \in \mathbf{S}_{1}^{\prime}$ for all $k, k \geq k_{0}$, and $X_{k+1}-M_{z}\left(\theta_{k}^{M}, k+1\right)$ and $\theta_{k}^{M}-\theta_{k}$ decay exponentially. Moreover, $\left(\theta_{k}^{M}\right)$ and $\left(X_{k+1}, \theta_{k}\right)$ have the same type of Lyapunov stability or instability.

Proof of Theorem 2. We study an auxiliary system defined as follows. Let ( $\phi_{k}$ ) be a sequence whose elements lie in $\mathbf{S}_{1}^{\prime}$ and satisfy, uniformly on $\mathbb{N}$, the set of positive integers, $\left\|\phi_{k+1}-\phi_{k}\right\| \leq \varepsilon c$. We consider the following auxiliary system on $\mathbb{R}^{n} \times \mathbb{R}^{p}:$

$$
\begin{align*}
\partial X_{k+1} & =A\left(\phi_{k}+\partial \theta_{k}\right) \partial X_{k}+B^{\prime \prime}\left(\partial \theta_{k}, k\right), \quad \partial X_{1} \in B\left(0, \frac{x-m}{\alpha}\right) \\
\partial \theta_{k} & =\partial \theta_{k-1}+\varepsilon C^{\prime \prime}\left(\partial X_{k}, \partial \theta_{k-1}, k\right)
\end{align*}
$$

where $B^{\prime \prime}, C^{\prime \prime}$ are defined and Lipschitz continuous on $B(O, \eta) \times \mathbb{N}$ and $B(O, x-m) \times B(O, \eta) \times \mathbb{N}$, respectively, with $b_{1}^{\prime \prime}, c_{1}^{\prime \prime}$ as respective Lipschitz constants and such that, uniformly on $\mathbb{N}, B^{\prime \prime}(0, k)=0,\left\|C^{\prime \prime}(0,0, k)\right\| \leq c_{1}^{\prime \prime} x_{1} \rho^{k-1}, \rho<1$. For two sequences $\left(\phi_{k}^{0}\right),\left(\phi_{k}^{1}\right)$ satisfying, uniformly on $\mathbb{N},\left\|\phi_{k}^{0}-\phi_{k}^{1}\right\| \leq \nu$ we are interested in the relation between the solutions of the corresponding auxiliary systems $\partial S^{0}, \partial S^{1}$, assuming that

$$
\begin{gathered}
\left\|B^{\prime \prime}\left(\partial \theta^{0}, k\right)-B^{\prime \prime 1}\left(\partial \theta^{1}, k\right)\right\| \leq b_{2}^{\prime \prime}\left[\operatorname{Max}\left\{\left\|\partial \theta^{0}\right\|,\left\|\partial \theta^{1}\right\|\right\} v+\left\|\partial \theta^{0}-\partial \theta^{1}\right\|\right] \\
\left\|C^{\prime \prime 0}\left(\partial X^{0}, \partial \theta^{0}, k\right)-C^{\prime \prime}\left(\partial X^{1}, \partial \theta^{1}, k\right)\right\| \\
\leq c_{2}^{\prime \prime}\left[\left(\operatorname{Max}\left\{\left\|\partial X^{0}\right\|+\left\|\partial \theta^{0}\right\|,\left\|\partial X^{1}\right\|+\left\|\partial \theta^{1}\right\|\right\}+\rho^{k-1}\right) v\right. \\
\left.+\left\|\partial X^{0}-\partial X^{1}\right\|+\left\|\partial \theta^{0}-\partial \theta^{1}\right\|\right]
\end{gathered}
$$

As for describing the $\varepsilon$-pseudo solution property, we use an exponentially weighted supremum norm to study this system:

$$
|\partial X|=\operatorname{Sup}_{k \geq 0}\left(\sigma^{-k}\left\|\partial X_{k+1}\right\|\right), \quad|\partial \theta|=\operatorname{Sup}_{k \geq 0}\left(\sigma^{-k}\left\|\partial \theta_{k}\right\|\right), \quad \rho \leq \sigma<1
$$

The following lemma states the existence of a unique solution of $\partial S$ which is bounded for this norm.

Lemma 6. Under the above assumptions and assumptions A 2 and A 3 , there exists $\varepsilon^{*}$ such that for any $\varepsilon, \varepsilon \leq \varepsilon^{*}$, we can find constants $\partial^{z}, g_{1}, \sigma$ such that for any initial $\partial X$-condition $\partial X_{1}$ in $B(O,(x-m) / \alpha)$ :
(i) The system $\partial S$ has a unique bounded solution satifying

$$
|\partial X| \leq x_{1}^{\prime}=\operatorname{Max}\left(\alpha\left\|\partial X_{1}\right\|, \varepsilon \partial^{x}\right), \quad|\partial \theta| \leq \varepsilon \frac{c_{1}^{\prime \prime}\left(x_{1}+x_{1}^{\prime}\right)}{1-\sigma-\varepsilon c_{1}^{\prime \prime}} \leq \eta
$$

(ii) If $\left(\partial X_{k+1}^{0}, \partial \theta_{k}^{0}\right),\left(\partial X_{k+1}^{1}, \partial \theta_{k}^{1}\right)$ are these solutions for the systems $\partial S^{0}, \partial S^{1}$, respectively, we have $\left|\partial X^{0}-\partial X^{1}\right|+\left|\partial \theta^{0}-\partial \theta^{1}\right| \leq g_{1}\left(\left\|\partial X_{1}^{0}-\partial X_{1}^{1}\right\|+\nu\right)$. The $\sigma, \partial^{x}$, and $g_{1}$ are specified in the proof.

Let us explain the interest of this lemma for proving Theorem 2:
If $\left(X_{k+1}, \theta_{k}\right)$ is a sequence of $\mathbf{B}\left(\left[k_{0}, \infty\right), \mathbf{S}_{1}^{\prime}\right)$, we define its auxiliary system $\partial S$ by choosing

$$
\begin{aligned}
\phi_{k} & =\theta_{k+k_{0}}, \quad \partial X_{1}=0, \\
B^{\prime \prime}(\partial \theta, k) & =\left(A\left(\phi_{k}+\partial \theta\right)-\left(A\left(\phi_{k}\right)\right)\right) X_{k+k_{0}}+\left(B\left(\phi_{k}+\partial \theta\right)-B\left(\phi_{k}\right)\right) u_{k+k_{0}}, \\
C^{\prime \prime}(\partial X, \partial \theta, k) & =C\left(M_{\varepsilon}\left(\phi_{k-1}+\partial \theta, k+k_{0}\right), \phi_{k-1}+\partial \theta, k+k_{0}\right)-C\left(X_{k+k_{0}}, \phi_{k-1}, k+k_{0}\right) .
\end{aligned}
$$

Using the properties of the elevation above the manifold and Lemma 4, we can check that $B^{\prime \prime}, C^{\prime \prime}$ satisfy our assumptions with

$$
\begin{aligned}
b_{1}^{\prime \prime}= & a_{1} x+b_{1} u, \\
b_{2}^{\prime \prime}= & \operatorname{Max}\left\{x\left(a_{1}+\frac{a_{2}}{2} \eta\right)+u\left(b_{1}+\frac{b_{2}}{2} \eta\right), a_{2} x+b_{2} u, a_{1}\right\}, \\
c_{1}^{\prime \prime}= & c_{1}\left(1+m_{1}\right), \\
c_{2}^{\prime \prime}= & \operatorname{Max}\left\{\left(c_{1}+\frac{c_{2}}{2}(x-m)+\eta\left(1+m_{1}\right)\right)\left(m_{1}+\frac{m_{2}}{2} \eta\right),\right. \\
& \left(c_{1}+\frac{c_{2}}{2}(x-m)+\eta\left(1+m_{1}\right)\right) m_{2}+c_{2}\left(1+m_{1}\right), \\
& \left.\left(c_{1}+\frac{c_{2}}{2}(x-m)+\eta\left(1+m_{1}\right)\right) q\left(1+\frac{\tau_{0}}{\rho-\tau_{0}}\right)+c_{2}(x-m)\right\},
\end{aligned}
$$

where $\eta$ is the "distance" between the imbeded compact sets, $q$ and $\rho$ are given by Lemma 3, $x_{1}=\alpha\left\|E_{k_{0}+1}\right\|$ and $v=\operatorname{Sup}_{k \geq k_{0}}\left(\left\|X_{k+1}^{0}-X_{k+1}^{1}\right\|+\left\|\theta_{k}^{0}-\theta_{k}^{1}\right\|\right)$. Then if ( $\partial X_{k+1}, \partial \theta_{k}$ ) is the particular solution of $\partial S$ given by Lemma 6 , we define $\Phi\left(k_{0}\right)(X, \theta)_{k}=\left(X_{k+1+k_{0}}^{2}+\partial X_{k+1}, \theta_{k+k_{0}}+\partial \theta_{k}\right)$. From the properties of $\left(\partial X_{k+1}, \partial \theta_{k}\right)$, it follows that $\Phi\left(k_{0}\right)$ satisfies (i)-(iii) of Theorem 2.

Similarly, if $\left(X_{k+1}^{M}, \theta_{k}^{M}\right)$ is a sequence of $\mathbf{B}^{M}\left(\left[k_{0}, \infty\right), \mathbf{S}_{1}^{\prime}\right)$, we define $\partial S$ by

$$
\begin{aligned}
\phi_{k}= & \theta_{k+k_{0}}^{M}, \quad \partial X_{1}=X_{k_{0}+1}^{M}-M_{\varepsilon}\left(\theta_{k_{0}}^{M}, k_{0}+1\right), \\
B^{\prime \prime}(\partial \theta, k)= & \left(A\left(\phi_{k}+\partial \theta\right)-A\left(\phi_{k}\right)\right) M_{\varepsilon}\left(\phi_{k-1}, k+k_{0}\right) \\
& +\left(B\left(\phi_{k}+\partial \theta\right)-B\left(\phi_{k}\right)\right) u_{k+k_{0}}, \\
C^{\prime \prime}(\partial X, \partial \theta, k)= & C\left(M_{\varepsilon}\left(\phi_{k-1}, k+k_{0}\right)+\partial X, \phi_{k-1}+\partial \theta, k+k_{0}\right) \\
& -C\left(M_{\varepsilon}\left(\phi_{k-1}, k+k_{0}\right), \phi_{k-1}, k+k_{0}\right) .
\end{aligned}
$$

Again our assumptions are satisfied with

$$
\begin{aligned}
& b_{1}^{\prime \prime}=a_{1} m+b_{1} u, \\
& b_{2}^{\prime \prime}=\operatorname{Max}\left\{m\left(a_{1}+\frac{a_{2}}{2} \eta\right)+u\left(b_{1}+\frac{b_{2}}{2} \eta\right), a_{2} m+b_{2} u+a_{1} m_{1}\right\}, \\
& c_{1}^{\prime \prime}=c_{1} \\
& c_{2}^{\prime \prime}=\operatorname{Max}\left\{\left(c_{1}+\frac{c_{2}}{2}(x-m+\eta)\right), c_{2}\left(1+m_{1}\right)\right\},
\end{aligned}
$$

and

$$
\rho=0, \quad v=\operatorname{Sup}_{k \geq k_{0}}\left\|\theta_{k}^{M 0}-\theta_{k}^{M 1}\right\|, \quad x_{1}=0
$$

With ( $\partial X_{k+1}, \partial \theta_{k}$ ) the corresponding solution of $\partial S$, we define

$$
\Phi^{M}\left(k_{0}\right)\left(X^{M}, \theta^{M}\right)_{k}=\left(M_{\varepsilon}\left(\theta_{k+k_{0}}^{M}, k+1+k_{0}\right)+\partial X_{k+1}, \theta_{k+k_{0}}^{M}+\partial \theta_{k}\right) .
$$

Then $\Phi^{M}\left(k_{0}\right)$ satisfies (i)-(iii) of Theorem 2. Note that, in this case, we have

$$
g^{x}(e)=\frac{\alpha\left(a_{1} x+b_{1} u\right)}{\sigma-\gamma(1+\varepsilon \beta c)} \frac{\sigma \alpha c_{1}^{\prime \prime}}{1-\sigma-\varepsilon c_{1}^{\prime \prime}} e .
$$

This is obtained by studying the sequence $\left(\tilde{X}_{k+1}-X_{k+1}^{M}\right)$, knowing that (for $e$ larger than $\varepsilon \partial^{x}$ )

$$
\tilde{X}_{k_{0}+1}=X_{k_{0}+1}^{M}, \quad\left\|\tilde{\theta}_{k}-\theta_{k}^{M}\right\| \leq \varepsilon \frac{c_{1}^{\prime \prime}}{1-\sigma-\varepsilon c_{1}^{\prime \prime}} \alpha e \sigma^{k-k_{0}} .
$$

Proof of Lemma 6. (i) Let $\mathbf{D}$ be the set of sequences $\left(\partial \theta_{k}\right)$ in $\mathbb{R}^{p}$ such that, for all positive $k$, we have $|\partial \theta| \leq \varepsilon \hat{\partial}^{\theta} \leq \eta$. Here, $\mathbf{D}$ is a complete metric space for the distance associated with the norm $|\partial \theta|$.

For any sequence ( $\partial \theta_{k}$ ) in $\mathbf{D}$, we consider the sequence ( $\partial X_{k}$ ) uniquely defined as the solution of $\partial X_{k+1}=A\left(\tilde{\phi}_{k}\right) \partial X_{k}+B^{\prime \prime}\left(\partial \theta_{k}, k\right)$ with $\partial X_{1}$ as the initial condition and $\tilde{\phi}_{k}=\phi_{k}+\partial \theta_{k}$. Since $\tilde{\phi}_{k}$ is in $S$, we have

$$
\left\|\partial X_{k+1}\right\|\left\|_{\tilde{\phi}_{k}} \leq \gamma\left(1+\beta\left\|\tilde{\phi}_{k}-\tilde{\phi}_{k-1}\right\|\right)\right\| \partial X_{k}\left\|\tilde{\phi}_{k-1}+\alpha b_{1}^{\prime \prime}\right\| \partial \theta_{k} \|
$$

It follows that $|\partial X| \leq x_{1}^{\prime}=\operatorname{Max}\left(\alpha\left\|\partial X_{1}\right\|, \varepsilon \partial^{x}\right)$ if $\varepsilon, \sigma, \partial^{\theta}, \partial^{x}$ satisfy $\gamma(1+\beta(\varepsilon c+$ $\left.\left.\varepsilon \delta^{\theta} \sigma^{k-1}(1+\sigma)\right)\right) \partial^{x}+\alpha b_{1}^{\prime \prime} \partial^{\theta} \sigma \leq \sigma \partial^{x}$. Also for two sequences $\left(\partial \theta_{k}^{0}\right),\left(\partial \theta_{k}^{1}\right)$ in $\mathbf{D}$, we have

$$
\begin{aligned}
\left\|\partial X_{k+1}^{0}-\partial X_{k+1}^{1}\right\|_{\dot{\phi}_{k}^{0}} \leq & \gamma\left(1+\beta\left\|\tilde{\phi}_{k}^{0}-\tilde{\phi}_{k-1}^{0}\right\|\right)\left\|\partial X_{k}^{0}-\partial X_{k}^{1}\right\|_{\dot{\phi}_{k-1}^{0}} \\
& +\alpha\left(a_{1} x_{1}^{\prime} \sigma^{k-1}+b_{1}^{\prime \prime}\right)\left\|\partial \theta_{k}^{0}-\partial \theta_{k}^{1}\right\| .
\end{aligned}
$$

This implies $\left|\partial X^{0}-\partial X^{1}\right| \leq \partial_{1}^{x}\left|\partial \theta^{0}-\partial \theta^{1}\right|$ if

$$
\gamma\left(1+\beta\left(\varepsilon c+\varepsilon \partial^{\theta} \sigma^{k-1}(1+\sigma)\right)\right) \partial_{1}^{x}+\alpha\left(a_{1} x_{1}^{\prime} \sigma^{k}+b_{1}^{\prime \prime} \sigma\right) \leq \sigma \partial_{1}^{x}
$$

Let us now define an operator $T$ on $\mathbf{D}$ by $T(\partial \theta)_{k}=\partial \theta_{k+1}-\varepsilon C^{\prime \prime}\left(\partial X_{k+1}, \partial \theta_{k}, k+1\right)$. Clearly, from the above inequalities ( $\partial X_{k+1}, \partial \theta_{k}$ ), defined this way, is the solution mentioned in the lemma if and only if $\left(\partial \theta_{k}\right)$ is a unique fixed point of $T$ in $\mathbf{D}$.
(a) The operator $T(\partial \theta)$ is in D . We have

$$
\begin{aligned}
\left\|T(\partial \theta)_{k}\right\| & \leq|\partial \theta| \sigma^{k+1}+\varepsilon c_{1}^{\prime \prime}\left(|\partial X|+|\partial \theta|+x_{1}\right) \sigma^{k} \\
& \leq \sigma^{k} \varepsilon\left(\sigma \partial^{\theta}+c_{1}^{\prime \prime}\left(x_{1}^{\prime}+\varepsilon \partial^{\theta}+x_{1}\right)\right)
\end{aligned}
$$

which means that $\varepsilon, x_{1}^{\prime}, \partial^{\theta}, \sigma$ should satisfy $\rho \leq \sigma<1$ and $\sigma \partial^{\theta}+c_{1}^{\prime \prime}\left(x_{1}+x_{1}^{\prime}+\right.$ $\left.\varepsilon \partial^{\theta}\right) \leq \partial^{\theta}$.
(b) The operator $T$ is a contraction (uniformly in $\partial X_{1}$ ):

$$
\begin{aligned}
\left\|T\left(\partial \theta^{0}\right)_{k}-T\left(\partial \theta^{1}\right)_{k}\right\| & \leq \sigma^{k+1}\left|\partial \theta^{0}-\partial \theta^{1}\right|+\varepsilon c_{1}^{\prime \prime}\left(1+\partial_{1}^{x}\right) \sigma^{k}\left|\partial \theta^{0}-\partial \theta^{1}\right| \\
& \leq \sigma^{k}\left(\sigma+\varepsilon c_{1}^{\prime \prime}\left(1+\partial_{1}^{x}\right)\right)\left|\partial \theta^{0}-\partial \theta^{1}\right|
\end{aligned}
$$

Therefore the first part of our lemma holds if

$$
\begin{aligned}
\varepsilon \partial^{\theta} \leq \eta, \quad \rho & \leq \sigma<1, \quad x_{1}^{\prime}=\operatorname{Max}\left(\alpha\left\|\partial X_{1}\right\|, \varepsilon \partial^{x}\right) \\
\alpha b_{1}^{\prime \prime} \partial^{\theta} & \leq\left[\sigma-\gamma\left(1+\varepsilon \beta\left(c+\partial^{\theta}(1+\sigma)\right)\right)\right] \partial^{x}, \\
\alpha\left(a_{1} x_{1}^{\prime}+b_{1}^{\prime \prime} \sigma\right) & \leq\left[\sigma-\gamma\left(1+\varepsilon \beta\left(c+\partial^{\theta}(1+\sigma)\right)\right)\right] \partial_{1}^{x}, \\
c_{1}^{\prime \prime}\left(x_{1}+x_{1}^{\prime}\right) & \leq\left(1-\sigma-\varepsilon c_{1}^{\prime \prime}\right) \partial^{\theta}, \\
0 & \leq 1-\sigma-\varepsilon c_{1}^{\prime \prime}\left(1+\partial_{1}^{x}\right) .
\end{aligned}
$$

(ii) Proceeding as for (i), we obtain

$$
\begin{aligned}
& \left\|\partial X_{k+1}^{0}-\partial X_{k+1}^{1}\right\| \|_{\dot{\phi_{k}^{0}}} \\
& \quad \leq \gamma\left(1+\varepsilon \beta\left(c+\partial^{\theta} \sigma^{k-1}(1+\sigma)\right)\right)\left\|\partial X_{k}^{0}-\partial X_{k}^{1}\right\| \|_{\phi_{k-1}^{0}} \\
& \quad+\alpha\left[a_{1} x_{1}^{\prime} \sigma^{k-1}\left(v+\sigma^{k}\left|\partial \theta^{0}-\partial \theta^{1}\right|\right)+b_{2}^{\prime \prime}\left(\varepsilon \partial^{\theta} \sigma^{k} v+\sigma^{k}\left|\partial \theta^{0}-\partial \theta^{1}\right|\right)\right] .
\end{aligned}
$$

It follows that $\left|\partial X^{0}-\partial X^{1}\right| \leq \alpha\left\|\partial X_{1}^{0}-\partial X_{1}^{1}\right\|+g_{1}^{\nu} v+g_{1}^{\theta}\left|\partial \theta^{0}-\partial \theta^{1}\right|$ with

$$
\begin{aligned}
& \gamma\left(1+\varepsilon \beta\left(c+\partial^{\theta} \sigma^{k-1}(1+\sigma)\right)\right) g_{1}^{v}+\alpha\left(a_{1} x_{1}^{\prime}+\varepsilon \partial^{\theta} b_{2}^{\prime \prime} \sigma\right) \leq \sigma g_{1}^{v} \\
& \gamma\left(1+\varepsilon \beta\left(c+\partial^{\theta} \sigma^{k-1}(1+\sigma)\right)\right) g_{1}^{\theta}+\alpha\left(a_{1} x_{1}^{\prime} \sigma^{k}+b_{2}^{\prime \prime} \sigma\right) \leq \sigma g_{1}^{\theta}
\end{aligned}
$$

Now since $\left(\partial \theta_{k}^{0}\right),\left(\partial \theta_{k}^{1}\right)$ are fixed points of the operator $T$ defined in (i), we have $\left\|\partial \theta_{k}^{0}-\partial \theta_{k}^{1}\right\|$

$$
\leq \sigma^{k}\left[\sigma\left|\partial \theta^{0}-\partial \theta^{1}\right|+\varepsilon c_{2}^{\prime \prime}\left(\left(x_{1}^{\prime}+\varepsilon \partial^{\theta}+1\right) v+\left|\partial X^{0}-\partial X^{1}\right|+\left|\partial \theta^{0}-\partial \theta^{1}\right|\right)\right]
$$

Our conclusion follows with

$$
\begin{aligned}
\alpha\left(a_{1} x_{1}^{\prime}+\varepsilon b_{2}^{\prime \prime} \partial^{\theta} \sigma\right) & \leq\left[\sigma-\gamma\left(1+\varepsilon \beta\left(c+\partial^{\theta}(1+\sigma)\right)\right)\right] g_{1}^{\nu}, \\
\alpha\left(a_{1} x_{1}^{\prime}+b_{2}^{\prime \prime} \sigma\right) & \leq\left[\sigma-\gamma\left(1+\varepsilon \beta\left(c+\partial^{\theta}(1+\sigma)\right)\right)\right] g_{1}^{\theta}, \\
\operatorname{Max}\left(\alpha, g_{1}^{v}+\varepsilon \frac{c_{2}^{\prime \prime}\left(1+g_{1}^{\theta}\right)\left(x_{1}^{\prime}+\varepsilon \partial^{\theta}+1\right)}{1-\sigma}\right) & \leq\left(1-\varepsilon \frac{c_{2}^{\prime \prime}\left(1+g_{1}^{\theta}\right)}{1-\sigma}\right) g_{1} .
\end{aligned}
$$

To complete the proof of Theorem 2, we have to study the relation between $\Phi$ and $\Phi^{M}$. Let $\left(X_{k+1}, \theta_{k}\right)$ be a sequence in $\mathbf{B}\left(\left[k_{0}, \infty\right), \mathbf{S}_{0}^{\prime}\right)$. From part (i) of the theorem,
its image $\left(\tilde{X}_{k+1}^{M}, \tilde{\theta}_{k}^{M}\right)$ by $\Phi\left(k_{0}\right)$ lies in $\mathbf{B}^{M}\left(\left[k_{0}, \infty\right), \mathbf{S}_{1}^{\prime}\right)$. Similarly, the image $\left(\tilde{X}_{k+1}, \tilde{\tilde{\theta}}_{k}\right)$ by $\Phi^{M}\left(k_{0}\right)$ of $\left(\tilde{X}_{k+1}^{M}, \tilde{\theta}_{k}^{M}\right)$ lies in $\mathbf{B}\left(\left[k_{0}, \infty\right), S_{0}\right)$. Moreover, we have $\tilde{X}_{k_{0}+1}=\tilde{X}_{k_{0}+1}=$ $X_{k_{0}+1}$. By uniqueness, $\Phi^{M}\left(k_{0}\right) \circ \Phi\left(k_{0}\right)$ is the identity map if $\tilde{\theta}_{k_{g}}=\theta_{k_{0}}$. From (i), we have $\left\|\theta_{k}-\tilde{\theta}_{k}\right\| \leq \sigma^{k-k_{0}} \varepsilon g^{\theta}$ and $\left\|\tilde{\theta}_{k}-\tilde{\theta}_{k}\right\| \leq \sigma^{k-k_{0}} \varepsilon g^{\theta}$. Hence $\left\|\tilde{\theta}_{k}-\theta_{k}\right\| \leq 2 \sigma^{k-k_{0}} \varepsilon g^{\theta}$. But since $\left(X_{k+1}, \theta_{k}\right)$ and $\left(\tilde{X}_{k+1}, \tilde{\theta}_{k}\right)$ are in $B(O, x) \times \mathrm{S}_{0}$ for all $k, k \geq k_{0}$, we apply inequalities (2) iteratively to get for all $k, k \geq k_{0}$,

$$
\begin{aligned}
\left\|X_{k+1}-\tilde{X}_{k+1}\right\|\left\|_{\theta_{k}}-l\right\| \theta_{k}-\tilde{\tilde{\theta}}_{k} \| & \leq-\tau_{0}(\varepsilon, x)^{k-k_{0}} l\left\|\theta_{k_{0}}-\tilde{\theta}_{k_{0}}\right\| \\
\left\|\theta_{k}-\tilde{\theta}_{k}\right\| & \geq\left(1-\varepsilon c_{1}(1+l)\right)^{k-k_{0}}\left\|\theta_{k_{0}}-\tilde{\theta}_{k_{0}}\right\| .
\end{aligned}
$$

Therefore, for all $k, k \geq k_{0}$,

$$
\left\|\theta_{k_{0}}-\tilde{\theta}_{k_{0}}\right\| \leq 2 \varepsilon g^{\theta}\left(\frac{\sigma}{1-\varepsilon c_{1}(1+l)}\right)^{k-k_{0}}
$$

Hence our result holds if $0<1-\sigma-\varepsilon c_{1}(1+l)$. Proofs of $\Phi\left(k_{0}\right) \circ \Phi^{M}\left(k_{0}\right)=I_{d}$ and injectivity properties follow the same lines.

## 5. Study of the Reduced-Order System $S M_{e}^{\theta}$

With Theorem 2, we have established that stability and existence of solutions of $S_{\varepsilon}$ remaining in $B(O, x) \times \mathrm{S}_{0}^{\prime}$ after time $k_{0}$ can be obtained from similar properties of solutions of $S M_{\varepsilon}^{\theta}$. Therefore, we can concentrate our attention on this system. It is a nonlinear nonautonomous system and many approaches can be used to study stability and existence of bounded solutions.

In the $l_{2}$-stationary (stochastic process) case, averaging theory, leading to the associated differential equation technique, turns out to be a very appropriate tool to deal with the difficulty due to time dependence. This has been demonstrated by Ljung in the $C(X, \theta, k)$-vanishing case [L1], [LS], [KC] and by Anderson et al. [ABJ] and Benveniste et al. [BMP] in the nonvanishing case. Noticing that, as $\varepsilon$ is made smaller, $S M_{\varepsilon}^{\theta}$ becomes closer to a first-order approximation of an ordinary differential equation, our result allows us to decompose the associated differential equation technique into three steps:
(i) Application of the topological orbital equivalence to replace $S_{\varepsilon}$ by $S M_{\varepsilon}$.
(ii) Application of the averaging theory to $S M_{\varepsilon}^{\theta}$.
(iii) Approximation of a difference equation by a differential equation.

In the nonstationary case, one possibility to simplify the time dependence is to extend to the discrete-time case the stroboscopic method ideas according to Minorsky [M] (see also [F1]): if $u_{k}$ and $C(X, \theta, k)$ have a "slowly varying period" $p(K)$, then by flash illumination at times $t(K), t(K)+p(K), t(K)+p(K)+p(K+1)$, $\ldots$ we see a weakly nonstationary advance map, i.e., the solutions of $S M_{\varepsilon}^{\theta}$ observed on the time interval $[t(K), t(K)+p(K))$ are very similar to the solutions observed on the time interval $[t(K)+p(K), t(K)+p(K)+p(K+1)$ ). The idea is then to approximate this advance map. Let $(p(K))$ be a sequence of bounded positive integers, $p(K) \leq p$. Given $k_{0}$, we define flash illumination times after $k_{0}$ by
$t(K+1)=t(K)+p(K), t(0)=k_{0}$. The advance $\operatorname{map} S M_{\varepsilon K}^{\theta}$ from $t(K)$ to $t(K+1)$ of $S M_{\varepsilon}^{\theta}$, written as

$$
\begin{equation*}
\theta_{t(K+1)}^{M}=\theta_{t(K)}^{M}+\varepsilon C_{\varepsilon}^{\Sigma}\left(\theta_{t(K)}^{M}, K+1\right) \tag{k}
\end{equation*}
$$

is obtained from

$$
\theta_{k}^{M}=\theta_{t(K)}^{M}+\varepsilon \sum_{i=(\mathrm{K})}^{k-1} C\left(M_{\varepsilon}\left(\theta_{i}^{M}, i+1\right), \theta_{i}^{M}, i+1\right), \quad t(K)+1 \leq k \leq t(K+1) .
$$

To approximate $S M_{\varepsilon_{K}}^{\theta}$, we introduce the system $S M_{0_{k}}^{\theta}$ :

$$
\begin{equation*}
\psi_{K+1}=\psi_{K}+\varepsilon C_{0}^{\Sigma}\left(\psi_{K}, K+1\right) \tag{k}
\end{equation*}
$$

with

$$
C_{0}^{\Sigma}(\psi, K+1)=\sum_{i=t(K)}^{t(K+1)-1} C\left(M_{0}(\psi, i+1), \psi, i+1\right)
$$

Under assumptions A1-A5, $C_{0}^{\Sigma}$ is Lipschitz continuously differentiable on $\mathbf{S} \times \mathbb{N}$ with $c_{1}^{\sigma}, c_{2}^{\sigma}$ the respective Lipschitz constants: $c_{1}^{\sigma} \leq p c_{1}\left(1+m_{1}\right)$ and $c_{2}^{\sigma} \leq$ $p\left(c_{2}\left(1+m_{1}\right)^{2}+c_{1} m_{2}\right)$.

System $S M_{0 K}^{\theta}$ is simpler than $S M_{\varepsilon}^{\theta}$. In particular, for adaptive systems, the $C$-function is given by the controller designer and is typically the product of a gain vector times an adaptation error. In this case, $C_{0}^{\Sigma}$ and $\partial C_{0}^{\Sigma} / \partial \psi$ are correlations on the time interval $\left[t(K), t(K+1)\right.$ ) of components of $M_{0}$ ог $\partial M_{0} / \partial \theta$. And, from their definitions, $M_{0}$ and $\partial M_{0} / \partial \theta$ can be obtained by implementation of sensitivity filters and observation of the feedback system using a constant parameter vector $\theta$, i.e., in a classical linear feedback context. This latter aspect makes the assumptions on $C_{0}^{\Sigma}$ interpretable in terms of signals properties.

Although simpler, $S M_{0 K}^{\theta}$ is very helpful for understanding the behavior of solutions of $S M_{\varepsilon}^{\theta}$. This is possible since $S M_{0 K}^{\theta}$ is an $\varepsilon^{2}$-approximation of $S M_{\varepsilon K}^{\theta}$. Indeed, from Theorem 1, under assumptions A1-A5, we have, uniformly on $S_{0} \times \mathbb{N}$,

$$
\left\|C_{0}^{\Sigma}(\psi, K)-C_{\varepsilon}^{\Sigma}(\psi, K)\right\| \leq \varepsilon v_{0}, \quad\left\|\frac{\partial C_{0}^{\Sigma}}{\partial \psi}(\psi, K)-\frac{\partial C_{\varepsilon}^{\Sigma}}{\partial \psi}(\psi, K)\right\| \leq \varepsilon v_{1}
$$

where $v_{0}, v_{1}$ can be obtained by induction on $p$ as

$$
\begin{aligned}
& v_{0} \leq p c_{1} f^{0}+c_{1}^{\sigma} c \frac{p-1}{2} \\
& v_{1} \leq p\left(c_{2}\left(1+m_{1}\right) f^{0}+c_{1} f_{1}^{0}\right)+c_{2}^{\sigma} c \frac{p-1}{2}+\left(\left(1+\varepsilon c_{1}\left(1+m_{1}\right)\right)^{p-1}-1\right) c_{1}\left(1+m_{1}\right)
\end{aligned}
$$

Again invoking hyperbolicity properties, namely conservation of stable and unstable manifolds of hyperbolic solutions under small perturbations (see [H1] and [S]), we may expect that, to any hyperbolic solution of $S M_{0_{k}}^{\theta}$, there corresponds an $\varepsilon$-close solution of $S M_{\varepsilon}^{\theta}$ with the same hyperbolicity property.

However, a difficulty remaining in the study of stability of solutions of $S M_{0_{K}}^{\theta}$ is the time variations. They have two causes: the time variations of the system itself and the motion of the solutions studied.

To take care of the system time variations, we consider the case where, uniformly on $\mathrm{S}_{0} \times \mathbb{N},\left\|\left(\partial C_{0}^{\Sigma} / \partial \psi\right)(\psi, K)-\left(\partial C_{0}^{\Sigma} / \partial \psi\right)(\psi, K+1)\right\|$ is "small." This assumption concerns essentially the $k$-dependence of $C(X, \theta, k)$ and $u_{k}$. In particular, it is trivially satisfied in the periodic case by choosing $p(K)$ constant, equal to the period. It holds also in the almost-periodic case, i.e., if there exists a Lipschitz continuously differentiable function $\hat{C}_{0}^{\Sigma}(\psi)$ such that, for any $\varepsilon$, we can find $p$ for which we have, uniformly on $S \times \mathbb{N}$,

$$
\left\|C_{0}^{\Sigma}(\psi, K)-\hat{C}_{0}^{\Sigma}(\psi)\right\| \leq \varepsilon \hat{v}_{0}, \quad\left\|\frac{\partial C_{0}^{\Sigma}}{\partial \psi}(\psi, K)-\frac{\partial \hat{C}_{0}^{\Sigma}}{\partial \psi}(\psi)\right\| \leq \varepsilon \hat{v}_{1} .
$$

In this case, we replace $C_{0}^{\Sigma}(\psi, K)$ by $\mathcal{C}_{0}^{\Sigma}(\psi)$ in the definition of $S M_{0 K}^{\theta}$.
To take care of the motion of the solutions, we consider thơse evolving in a set where $C_{0}^{\Sigma}(\psi, K)$ is "small" uniformly in $K$. This condition is trivially satisfied if $S M_{0_{k}}^{\theta}$ has a fixed point, i.e., if there exits $\psi^{*}$ such that, for all $K, C_{0}^{\Sigma}\left(\psi^{*}, K\right)=0$. This equation is precisely the bifurcation equation obtained by averaging theory [BSA], [RK2] or critical systems theory [PP]. Existence of solutions for this equation has been studied for model reference adaptive systems with a fixed-point argument (Section 4.5 of [R]) or applying degree theory [PCP].

From this discussion, we introduce the following definition:
Definition. Given strictly positive constants $\zeta, \varepsilon$, we define the set $\mathbf{P}(\zeta, \varepsilon)$ as

$$
\mathbf{P}(\zeta, \varepsilon)=\left\{\psi \in \mathbf{S}_{0}^{\prime} \left\lvert\, \begin{array}{ll}
\mathrm{A} 6.1: & \left\|\frac{\partial C_{0}^{\Sigma}}{\partial \psi}(\psi, K)-\frac{\partial C_{0}^{\Sigma}}{\partial \psi}(\psi, K+1)\right\| \\
& +\varepsilon c_{2}^{\sigma}\left\|C_{0}^{\Sigma}(\psi, K)\right\| \leq \varepsilon w<\varepsilon \zeta^{2} \quad \text { for all } K, \\
\text { A6.2: } & \operatorname{Max}_{i \in\lfloor 1, p]} \operatorname{Re}\left(\lambda_{i}\left\{\frac{\partial C_{0}^{\Sigma}(\psi, K)}{\partial \psi}\right\}\right) \leq-\zeta \text { for all } K,
\end{array}\right.\right\}
$$

where $\operatorname{Re}(\cdot)$ and $\lambda_{i}\{\cdot\}$ denote the real part and the $i$ th eigenvalue, respectively.
Remark. In the adaptive linear systems context, the inequality A6.2, involved in the definition of $\mathbf{P}(\zeta, \varepsilon)$, is related to the so-called "signal dependent positivity condition." In the test input assumption case (see the Introduction), it can be interpreted as the positivity of an operator restricted to act on specific signals (see [ABJ], [RK1], and [RPK]).

Theorem 3. Under assumptions A1-A5, there exists $\varepsilon^{*}$ such that if we can find $\zeta$ and $\varepsilon, 0<\varepsilon \leq \varepsilon^{*}$, for which $\mathbf{P}(\zeta, \varepsilon)$ contains a solution $\left(\psi_{K}\right)$ of $S_{0 K}^{\theta}$, then $S M_{\varepsilon}^{\theta}$ has an exponentially stable solution remaining in $\mathbf{S}_{1}^{\prime}$ after time $k_{0}$ and $\varepsilon$-close to $\left(\psi_{K}\right)$ at times $t(K)$, i.e., for all $K,\left\|\theta_{t(K)}^{M}-\psi_{K}\right\| \leq \varepsilon \delta$ with $v_{0} / \zeta<\delta$.

The constant $\delta$, appearing in this statement, is clarified in the proof.
Remarks. This stroboscopic method approach extends, in mowe general situations, the local averaging technique proposed by Kosut et al. [KAM] (see also [ABJ]) for studying $S M_{0}^{\theta}$ after linearization under a relaxed test input assumption.

With Theorems 2 and 3, we have established the following result:
For $\varepsilon$ sufficiently small, if $S M_{0 K}^{\theta}$ has a solution remaining in a set $\mathbf{P}(\zeta, \varepsilon)$, then $S_{\varepsilon}$ has exponentially stablensolutions remaining in $B(O, x) \times \mathrm{S}_{0}$ after time $k_{0}$.

In other words, this proves that, as far as stability is concerned, the heuristic technique proposed by $\AA$ Aström [A1]; [A2] is theoretically sound when restricted to the set $\mathbf{P}(\zeta, \varepsilon)$.

Finally, $\left(\psi_{K}\right)$ is an $\varepsilon^{2}$-approximation of a solution of $S M_{\varepsilon_{K}}^{\theta}$ if $C_{0}^{\Sigma}$ is replaced by $C_{1}^{\Sigma}$, an $\varepsilon^{2}$-approximation of $C_{0}^{\Sigma}$ defined with $M_{1}$, the $\varepsilon^{2}$-approximation of $M_{\varepsilon}$.

As a key step to proving Theorem 3, we establish the following result:

Lemma 7. Let $F\left(K_{1}, K_{0}\right)$ be the transition matrix of

$$
\Delta_{K+1}=\left(I+\varepsilon \frac{\partial C_{0}^{\Sigma}}{\partial \psi}\left(\psi_{K}, K+1\right)\right) \Delta_{K}
$$

with $\left(\psi_{K}\right)$, the solution in $\mathbf{P}(\zeta, \varepsilon)$ given by Theorem 3. For any $\zeta^{\prime}, 0<\zeta^{\prime}<\zeta-$ $\varepsilon\left(c_{1}^{\sigma}\right)^{2} / 2$, there exists $v$ (independent of $\psi_{K}$ ) such that for any $K_{0}, K_{1}, 0 \leq K_{0}<K_{1}$,

$$
\left\|F\left(K_{1}, K_{0}\right)\right\| \leq \frac{1+\varepsilon v^{\prime}}{1-\varepsilon \zeta^{\prime \prime}}\left(1-\varepsilon \zeta^{\prime \prime}\right)^{K_{1}-K_{0}}
$$

with $\zeta^{\prime \prime}=\zeta^{\prime}-\left[(1+\varepsilon v)\left(1-\varepsilon \zeta^{\prime}\right) w\right]^{1 / 2}$ and $v^{\prime}=v\left(1-\varepsilon \zeta^{\prime}\right)-\zeta^{\prime \prime}$, where $w$ is obtained from A6.1.

Proof. Given $\mathrm{K}_{1}$, the transition matrix satisfies for any $N, K_{0} \leq N<K_{1}$,

$$
\begin{aligned}
F\left(N+1, K_{0}\right)= & \left(I+\varepsilon \frac{\partial C_{0}^{\Sigma}}{\partial \psi}\left(\psi_{K_{1}-1}, K_{1}\right)\right) F\left(N, K_{0}\right) \\
& +\varepsilon\left[\frac{\partial C_{0}^{\Sigma}}{\partial \psi}\left(\psi_{N}, N+1\right)-\frac{\partial C_{0}^{\Sigma}}{\partial \psi}\left(\psi_{K_{1}-1}, K_{1}\right)\right] F\left(N, K_{0}\right) .
\end{aligned}
$$

To derive the property of $F\left(K_{1}, K_{0}\right)$, we use the following three inequalities:
(i) From assumption A6.2 in the definition of the set $\mathbf{P}(\zeta, \varepsilon)$, we have, uniformly on $\mathbf{P}(\zeta, \varepsilon) \times \mathbb{N}$,

$$
\begin{aligned}
\left\|I+\varepsilon \frac{\partial C_{0}^{\Sigma}}{\partial \psi}(\psi, K)\right\| & \leq 1+\varepsilon c_{1}^{\sigma}, \\
\operatorname{Max}_{i \in[1, p]}\left|\lambda_{i}\left\{I+\varepsilon \frac{\partial C_{0}^{\Sigma}}{\partial \psi}(\psi, K)\right\}\right| & \leq\left[1-2 \varepsilon \zeta+\varepsilon^{2}\left\|\frac{\partial C_{0}^{\Sigma}}{\partial \psi}(\psi, K)\right\|^{2}\right]^{1 / 2} \\
& \leq 1-\varepsilon \zeta+\varepsilon^{2} \frac{\left(c_{1}^{\sigma}\right)^{2}}{2} .
\end{aligned}
$$

Then it follows from Theorem 5 of [F3] that for any $\zeta^{\prime}, 0<\zeta^{\prime}<\zeta-$
$\varepsilon\left(c_{1}^{\sigma}\right)^{2} / 2$, there exists $v\left(v \geq c_{1}^{\sigma}\right)$ such that, uniformly on $\mathbf{P}(\zeta, \varepsilon) \times \mathbb{N}$,

$$
\left\|\left(I+\varepsilon \frac{\partial C_{0}^{\Sigma}}{\partial \psi}(\psi, K)\right)^{i}\right\| \leq(1+\varepsilon v)\left(1-\varepsilon \zeta^{\prime}\right)^{i} \quad \text { for all } \quad i \in \mathbb{N} .
$$

(ii) From assumption A6.1 and the fact that $\left(\psi_{K}\right)$ is a solution of $S M_{0_{K}}^{\theta}$, we obtain

$$
\begin{aligned}
& \left\|\frac{\partial C_{0}^{\Sigma}}{\partial \psi}\left(\psi_{N}, N+1\right)-\frac{\partial C_{0}^{\Sigma}}{\partial \psi}\left(\psi_{K_{1}}, K_{1}\right)\right\| \\
& \quad \leq \sum_{K=N+1}^{K_{1}-1}\left\|\frac{\partial C_{0}^{\Sigma}}{\partial \psi}\left(\psi_{K-1}, K\right)-\frac{\partial C_{0}^{\Sigma}}{\partial \psi}\left(\psi_{K}, K+1\right)\right\| \\
& \quad \leq \varepsilon w\left(K_{1}-N-1\right)
\end{aligned}
$$

(iii) Let $u_{k}$ be a sequence of positive real numbers satisfying, for any $k, k \geq 1$, $u_{k} \leq a \lambda^{k} u_{0}+b \sum_{i=0}^{k-1} \lambda^{k-1-i}(k-1-i) u_{i}, a \geq 0, b \geq 0, \lambda \geq 0$, then we can check by induction that

$$
u_{k} \leq\left[\frac{\lambda a+\sqrt{\lambda b}}{2}(\lambda+\sqrt{\lambda b})^{k-1}+\frac{\lambda a-\sqrt{\lambda b}}{2}(\lambda-\sqrt{\lambda b})^{k-1}\right] u_{0}
$$

Now use the variation of constants formula; take the Euclidian norm and use the two first inequalities to obtain

$$
\begin{aligned}
\left\|F\left(K_{1}, K_{0}\right)\right\| \leq & (1+\varepsilon v)\left(1-\varepsilon \zeta^{\prime}\right)^{K_{1}-K_{0}} \\
& +\varepsilon^{2} w(1+\varepsilon v) \sum_{N=K_{0}}^{K_{1}-1}\left(1-\varepsilon \zeta^{\prime}\right)^{K_{1}-1-N}\left(K_{1}-1-N\right)\left\|F\left(N, K_{0}\right)\right\| .
\end{aligned}
$$

The result follows from the third inequality
Proof of Theorem 3. (a) Existence: The idea is to find a solution $\left(\theta_{t(K)}^{M}\right)$ of $S M_{\varepsilon K}^{\theta}$ satisfying $\left\|\theta_{t(K)}^{M}-\psi_{K}\right\| \leq \varepsilon \delta$. This will solve our problem, since the corresponding solution $\left(\theta_{k}^{M}\right)$ of $S M_{\varepsilon}^{\theta}$ satisfies, for all $k, t(K)-p / 2 \leq k \leq t(K)+p / 2$,

$$
\left\|\theta_{t(K)}^{M}-\theta_{k}^{M}\right\| \leq \varepsilon c \frac{p}{2}
$$

In particular, knowing that $\psi_{K}$ lies in $\mathbf{S}_{0}^{\prime}, \theta_{k}^{M}$ lies in $\mathbf{S}_{1}^{\prime}$ if $\delta$ and $\varepsilon$ satisfy $\delta \leq \eta / \varepsilon-$ $c(p / 2)$. Hence, let $\Delta_{K}$ be defined recursively by

$$
\begin{aligned}
\Delta_{0}= & 0 \\
\Delta_{K+1}= & \left(I+\varepsilon \frac{\partial C_{0}^{\Sigma}}{\partial \psi}\left(\psi_{K}, K+1\right)\right) \Delta_{K} \\
& +\varepsilon\left[C_{\varepsilon}^{\Sigma}\left(\psi_{K}+\Delta_{K}, K+1\right)-C_{0}^{\Sigma}\left(\psi_{K}, K+1\right)-\frac{\partial C_{0}^{\Sigma}}{\partial \psi}\left(\psi_{K}, K+1\right) \Delta_{K}\right]
\end{aligned}
$$

then $\left(\psi_{K}+\Delta_{K}\right)$ is a solution of $S M_{s K}^{\theta}$. We remark that for any $\Delta$ such that the segment $\left[\psi_{K}, \psi_{K}+\Delta\right.$ ] is contained in $S_{0}$, application of the Mean Value Theorem
yields for some $\xi, 0 \leq \xi \leq 1$,

$$
\begin{aligned}
& \left.\| C_{\varepsilon}^{\Sigma}\left(\psi_{K}+\Delta, K+1\right)-C_{0}^{\Sigma}\left(\psi_{K}, K+1\right)-\frac{\partial C_{0}^{\Sigma}}{\partial \psi}\left(\psi_{K}, K+1\right) \Delta\right) \| \\
& \quad \leq \varepsilon v_{0}+\left\|\frac{\partial \stackrel{\rightharpoonup}{C_{0}^{\Sigma}}}{\partial \psi}\left(\psi_{K}+\xi \Delta, K+1\right)-\frac{\partial C_{0}^{\Sigma}}{\partial \psi}\left(\psi_{K}, K+1\right)\right\|\|\Delta\| \leq \varepsilon v_{0}+c_{2}^{\sigma}\|\Delta\|^{2}
\end{aligned}
$$

Now apply the variation of constants formula and Lemma 7. Then

$$
\left\|\Delta_{K+1}\right\| \leq \varepsilon \frac{1+\varepsilon v^{\prime}}{1-\varepsilon \zeta^{\prime \prime}} \sum_{N=0}^{K}\left(1-\varepsilon \zeta^{\prime \prime}\right)^{K-N}\left(\varepsilon v_{0}+c_{2}^{\sigma}\left\|\Delta_{N}\right\|^{2}\right)
$$

Hence, we have established $\left\|\Delta_{K}\right\| \leq \varepsilon \delta$ and, therefore, the existence of a solution ( $\theta_{k}^{M}$ ) of $S M_{\varepsilon}^{\theta}$ remaining in $S_{1}^{\prime}$ after time $k_{0}$, if $\zeta^{\prime}, \delta$ satisfy

$$
\frac{1+\varepsilon v^{\prime}}{1-\varepsilon \zeta^{\prime \prime}} \frac{v_{0}+\varepsilon c_{2}^{\sigma} \delta^{2}}{\zeta^{\prime \prime}} \leq \delta \leq \frac{\eta}{\varepsilon}-c \frac{p}{2}
$$

(b) Exponential stability: For any two solutions $\left(\theta_{k}^{M 0}\right),\left(\theta_{k}^{M 1}\right)$ of $S M_{\varepsilon}^{\theta}$ remaining in $S_{0}$, we have, uniformly in $K$,

$$
\left\|\theta_{k}^{M 0}-\theta_{k}^{M 1}\right\| \leq\left(1+\varepsilon c_{1}\left(1+m_{1}\right)\right)^{p}\left\|\theta_{t(K)}^{M 0}-\theta_{t(K)}^{M 1}\right\|, \quad t(K) \leq k \leq T(K+1) .
$$

Hence exponential stability of the solution $\left(\theta_{k}^{M}\right)$ obtained in part (a) follows from exponential stability of $\left(\theta_{t(\mathrm{~K})}^{M}\right)$, the solution of $S M_{\varepsilon K}^{\theta}$. But, from Lyapunov's theorem, for this property to hold it is sufficient that the origin be an exponentially stable solution of the following linear system:

$$
\Delta_{K+1}=\left(I+\varepsilon \frac{\partial C_{\varepsilon}^{\Sigma}}{\partial \psi}\left(\theta_{t(K)}^{M}, K+1\right)\right) \Delta_{K} .
$$

From part (a), we know

$$
\left\|\frac{\partial C_{\varepsilon}^{\Sigma}}{\partial \psi}\left(\theta_{t(K)}^{M}, K+1\right)-\frac{\partial C_{0}^{\Sigma}}{\partial \psi}\left(\psi_{K}, K+1\right)\right\| \leq \varepsilon v_{1}+\varepsilon c_{1}^{\dot{\sigma}} \delta
$$

Therefore, with Lemma 7 and Lemma B. 5 on p. 118 of [A4], our property holds if

$$
1-\varepsilon \zeta^{\prime \prime}+\varepsilon^{2} \frac{1+\varepsilon v^{\prime}}{1-\varepsilon \zeta^{\prime \prime}}\left(v_{1}+c_{1}^{\sigma} \delta\right)<1
$$

In conclusion, our theorem holds if $\varepsilon$ and $\delta$ satisfy, for $v$ given by Theorem 5 of [F3],

$$
\begin{gathered}
0<\zeta^{\prime}<\zeta-\varepsilon \frac{\left(c_{1}^{\delta}\right)^{2}}{2}, \quad \zeta^{\prime \prime}=\zeta^{\prime}-\left[(1+\varepsilon v)\left(1-\varepsilon \zeta^{\prime}\right) w\right]^{1 / 2}, \quad v^{\prime}=v\left(1-\varepsilon \zeta^{\prime}\right)-\zeta^{\prime \prime} \\
w<\frac{\zeta^{\prime 2}}{(1+\varepsilon v)\left(1-\varepsilon \zeta^{\prime}\right)} \\
\\
\frac{1+\varepsilon v^{\prime}}{1-\varepsilon \zeta^{\prime \prime}} \frac{v_{0}+\varepsilon c_{2}^{\sigma} \delta^{2}}{\zeta^{\prime \prime}} \leq \delta \leq \frac{\eta}{\varepsilon}-c \frac{p}{2} \\
1-\varepsilon \zeta^{\prime \prime}+\varepsilon^{2} \frac{1+\varepsilon v^{\prime}}{1-\varepsilon \zeta^{\prime \prime}}\left(v_{1}+c_{1}^{\sigma} \delta\right)<1
\end{gathered}
$$

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