# Computer assisted proof of the existence of the Lorenz attractor in the Shimizu-Morioka system 

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#### Abstract

We prove that the Shimizu-Morioka system has a Lorenz attractor for an open set of parameter values. For the proof we employ a criterion proposed by Shilnikov, which allows to conclude the existence of the attractor by examination of the behavior of only one orbit. The needed properties of the orbit are established by using computer assisted numerics. Our result is also applied to the study of local bifurcations of triply degenerate periodic points of three-dimensional maps. It provides a formal proof of the birth of discrete Lorenz attractors at various global bifurcations.


## 1. Introduction and main results

In this paper we provide a solution to a long-standing open problem. We prove, by employing rigorous numerics, that Shimizu-Morioka system has a Lorenz attractor for an open set of parameter values. The fact itself is well-known, see [1, 2, 3]; however its rigorous proof was missing which impeded the further development of the mathematical theory of Lorenz-like attractors. Here, we report a computer assisted proof which formally closes this problem.
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Figure 1. Two kinds of Lorenz attracto in the Shimizu-Morioka systemr: standard (for $\alpha=0.45, \lambda=0.9$ ) on the left, and with a lacuna (for $\alpha=0.5, \lambda=0.85$ ) on the right.

The Shimizu-Morioka system

$$
\left\{\begin{array}{l}
\dot{x}=y,  \tag{1}\\
\dot{y}=(1-z) x-\lambda y, \\
\dot{z}=-\alpha z+x^{2}
\end{array}\right.
$$

where $(x, y, z)$ are coordinates in $\mathbb{R}^{3}$ and $\alpha>0, \lambda>0$ are parameters, was introduced by T.Shimizu and N.Morioka in [4] and extensively studied numerically by A. Shilnikov in [1, 2]. One of the main findings was that there is a large open region in the $(\alpha, \lambda)$ plane where this system has a strange attractor very similar to the classical attractor of the Lorenz model (5), see Fig. 1. This is more than just the similarity in shape: as one can infer from the pictures numerically obtained in [1], the Poincaré map on a two-dimensional cross-section is hyperbolic (very strongly contracting in one direction and expanding in the other direction), so the attractor in the Shimizu-Morioka system can be described by the Afraimovich-Bykov-Shilnikov geometric Lorenz model [5, 6], i.e., it is a Lorenz attractor.

We will describe the geometric Lorenz model and give the corresponding definition of the Lorenz attractors in a moment. However, we first want to stress that ShimizuMorioka model is special (and maybe more important than the classical Lorenz model). The reason is that system (1) is a truncated normal form for certain codimension3 bifurcations of equilibria and periodic orbits [3, 7]. Because Lorenz or Lorenz-like attractors persist at small perturbations, our result on the existence of the Lorenz attractor in the normal form system (1) implies that Lorenz-like attractors exist for any small perturbation of (1). This proves the emergence of the Lorenz attractor (or its discrete analogue) in the corresponding class of codimension-3 bifurcations. In particular, in this paper, using Theorem 3, we prove the existence of discrete Lorenz attractors in a class of three-dimensional polynomial maps (3D Hénon maps, see Theorem 4 and remarks after it).

### 1.1. Pseudohyperbolicity

In order to talk about Lorenz attractors, we need a proper definition of them. We use the Afraimovich-Bykov-Shilnikov geometric model [5, 6] which we describe using the notion of pseudo-hyperbolicity introduced in [8, 9]. A system (a smooth flow or a diffeomorphism) in $\mathbb{R}^{n}$ is called pseudo-hyperbolic in a strictly forward-invariant domain $\mathcal{D} \subseteq \mathbb{R}^{n}$ if:

1) there are directions in which the dynamical system (a flow or a diffeomorphism) is strongly contracting ("strongly" means that any possible contraction in transverse directions is always strictly weaker);
2) transverse to the contracting directions the system is volume-expanding (i.e. the volume is stretched exponentially).

In precise terms, condition 1 reads as follows. For each point of $\mathcal{D}$, we assume that there exists a pair of transversal subspaces $N_{1}$ and $N_{2}$ (with $\operatorname{dim}\left(N_{2}\right)=k \geq 1$ and $\left.\operatorname{dim}\left(N_{1}\right)=n-k\right)$, continuously depending on the point, such that the families of these subspaces are invariant with respect to the derivative $D X_{t}$ of the time- $t$ map $X_{t}$ of the system, i.e., $D X_{t} N_{1}(x)=N_{1}\left(X_{t}(x)\right)$ and $D X_{t} N_{2}(x)=N_{2}\left(X_{t}(x)\right)$ for all $t \geq 0$ (in the case of a diffeomorphism, $t$ runs all positive integer values). We also assume that there exist constants $C>0, \alpha>0$ and $\beta>0$ such that for each $x \in \mathcal{D}$ and all $t \geq 0$

$$
\begin{equation*}
\left\|\left.D X_{t}(x)\right|_{N_{2}}\right\| \leq C e^{-\alpha t} \tag{2}
\end{equation*}
$$

and

$$
\begin{equation*}
\left\|\left.D X_{t}(x)\right|_{N_{2}}\right\| \cdot\left\|\left(\left.D X_{t}(x)\right|_{N_{1}}\right)^{-1}\right\| \leq C e^{-\beta t} \tag{3}
\end{equation*}
$$

The volume-expansion condition 2 reads as follows: there exist constants $C>0$ and $\sigma>0$ such that for each $x \in \mathcal{D}$ and all $t \geq 0$

$$
\begin{equation*}
\operatorname{det}\left(\left.D X_{t}(x)\right|_{N_{1}}\right) \geq C e^{\sigma t} \tag{4}
\end{equation*}
$$

Inequality (3) (the so-called cone condition) ensures that the invariant families of subspaces $N_{1}$ and $N_{2}$ continuously persist for all $C^{1}$-small perturbations of the system, so the pseudohyperbolic structure is a robust property of the system. Inequality (4) guarantees that for every orbit in $\mathcal{D}$ its maximal Lyapunov exponent is positive. Therefore, the existence of the pseudohyperbolic structure in a bounded domain $\mathcal{D}$ ensures chaotic dynamics in $\mathcal{D}$, which cannot be destroyed by small smooth perturbations.

Conditions (22), (3) imply also the existence of a strong-stable invariant foliation $\mathcal{N}^{s s}$ in $\mathcal{D}$, whose field of tangents is the family $N_{2}$. For every two points in the same leaf of $\mathcal{N}^{s s}$ the distance between their forward orbits tends to zero exponentially, i.e., the forward dynamics of all points in the same leaf are the same.

### 1.2. Lorenz attractor

The Lorenz attractor of the Afraimovich-Bykov-Shilnikov model is the attractor of a pseudohyperbolic system of differential equations with $\operatorname{dim}\left(N_{1}\right)=2$. Specifically,
consider a system $X$ of differential equations with a saddle equilibrium state $O$. Assume that $O$ has a one-dimensional unstable manifold $W^{u}(O)$ and an $(n-1)$-dimensional stable manifold $W^{s}(O)$, i.e., if $\lambda_{1}, \ldots, \lambda_{n}$ are the eigenvalues of the linearized system at $O$, then $\lambda_{1}>0$ and $\operatorname{Re} \lambda_{j}<0$ for $j \geq 2$. Assume that

$$
\lambda_{2}>\operatorname{Re} \lambda_{j} \quad \text { for } j \geq 3
$$

and

$$
\lambda_{1}+\lambda_{2}>0 .
$$

This means that the pseudohyperbolicity conditions (2), (3) and (4) are fulfilled at the point $O$, with $N_{1}$ being the two-dimensional eigenspace corresponding to the eigenvalues $\lambda_{1}$ and $\lambda_{2}$ and $N_{2}$ being the eigenspace corresponding to the rest of eigenvalues. We assume that the pseudohyperbolicity property also holds in a sufficiently large, bounded neighborhood $\mathcal{D}$ of $O$. We assume that $\mathcal{D}$ is strictly forward-invariant, i.e., there exists $T>0$ such that the image of the closure of $\mathcal{D}$ by the time- $T$ map lies strictly inside $\mathcal{D}$.

Moreover, we assume that in $\mathcal{D}$ there exists an $(n-1)$-dimensional cross-section $\Pi$ to the flow, such that for every point in $\mathcal{D}$ its forward orbit either tends to $O$ (i.e., it lies in $W^{s}(O)$ ) or hits $\Pi$ at some moment of time. We assume that $\Pi$ is divided by a smooth ( $n-2$ )-dimensional surface $\Pi_{0}$ into 2 halves, $\Pi_{+}$and $\Pi_{-}$, such that the orbits starting in $\Pi_{+}$and $\Pi_{-}$return to $\Pi$ again, while the orbits starting in $\Pi_{0}$ tend to $O$ as $t \rightarrow+\infty$, i.e., $\Pi_{0} \subset W^{u}(O) \cap \Pi$. Thus, the orbits of $X$ define the Poincaré map $T: \Pi_{+} \cup \Pi_{-} \rightarrow \Pi$.

This map is smooth outside a discontinuity surface $\Pi_{0}$. The orbits starting close to $\Pi_{0}$ pass near $O$, so the return time tends to infinity as the initial point $M$ tends to $\Pi_{0}$; note that $T M$ tends to one of the two points, $M_{+}$and $M_{-}$, where $W^{u}(O)$ intersects $\Pi$. The unstable manifold $W^{u}(O)$ is one-dimensional, so $W^{u}(O) \backslash O$ consists of exactly two orbits, $\Gamma_{+}$and $\Gamma_{-}$, called separatrices, and the points $M_{+}$and $M_{-}$are the first points where the $\Gamma_{+}$and $\Gamma_{-}$intersect $\Pi$. We take the convention that $\lim _{M \rightarrow \Pi_{0}} T(M)$ equals to $M_{+}$if the initial point $M$ approaches $\Pi_{0}$ from $\Pi_{+}$and $M_{-}$if $M$ approaches $\Pi_{0}$ from $\Pi_{-}$.

The invariant foliation $\mathcal{N}^{s s}$ of the system $X$ also corresponds to a codimension- 1 strong-stable invariant foliation $\mathcal{N}_{\Pi}^{s s}$ for the map $T$ : the leaves of $\mathcal{N}_{\Pi}^{s s}$ are obtained as intersections with $\Pi$ of the orbits of the leaves of $\mathcal{N}^{s s}$ by the flow. The foliation is contracting, i.e., the iterations by $T$ of any two points in the same leaf of $\mathcal{N}_{\Pi}^{s s}$ converge exponentially to each other. However, the dynamics transverse to the strong-stable foliation are chaotic. The expansion of areas by the flow transverse to $\mathcal{N}^{s s}$ implies that the Poincaré map $T$ is uniformly expanding in the direction transverse to $\mathcal{N}_{\Pi}^{s s}$.

In other words, the pseudohyperbolicity of the flow in $\mathcal{D}$ implies that the Poincaré map $T$ is locally uniformly hyperbolic. However, one cannot directly apply the theory of uniform hyperbolicity to this map, as it has a singularity (at $\Pi_{0}$ ). Nevertheless, the (singular) hyperbolicity of the Poincaré map $T$ makes a comprehensive analysis of the structure of the Lorenz attractor in the above described Afraimovich-Bykov-Shilnikov model possible, as it was done in [6]. The results of these papers can be summarized as follows [10]:

Theorem 1. [6] Under conditions above, the set of non-wandering orbits of the system in $\mathcal{D}$ consists of a uniquely defined, two-dimensional closed invariant set $\mathcal{A} \subset \mathcal{D}$ (which is called the Lorenz attractor) and a (possibly empty) one-dimensional closed invariant set $\Sigma$ (which may intersect $\mathcal{A}$ but is not a subset of $\mathcal{A}$ ) such that
(1) the separatrices $\Gamma_{+}$and $\Gamma_{-}$and the saddle $O$ lie in $\mathcal{A}$;
(2) $\mathcal{A}$ is transitive, and saddle periodic orbits are dense in $\mathcal{A}$;
(3) $\mathcal{A}$ is the limit of a nested sequence of hyperbolic, transitive, compact invariant sets each of which is equivalent to a suspension over a finite Markov chain with positive topological entropy;
(4) $\mathcal{A}$ is structurally unstable: arbitrarily small smooth perturbations of the system lead to the creation of homoclinic loops to $O$ and to the subsequent birth and/or disappearance of saddle periodic orbits within $\mathcal{A}$;
(5) when $\Sigma=\emptyset$, the set $\mathcal{A}$ is the maximal attractor in $\mathcal{D}$;
(6) when the set $\Sigma$ is non-empty, it is a hyperbolic set equivalent to a suspension over a finite Markov chain (it may have zero entropy, e.g. be a single saddle periodic orbit);
(7) every forward orbit in $\mathcal{D}$ tends to $\mathcal{A} \cap \Sigma$ as $t \rightarrow+\infty$;
(8) when $\Sigma \neq \emptyset$, the maximal attractor in $\mathcal{D}$ is $\operatorname{cl}\left(W^{u}(\Sigma)\right)=\mathcal{A} \cap W^{u}(\Sigma)$;
(9) A attracts all orbits from its neighbourhood when $\mathcal{A} \cup \Sigma=\emptyset$.

In [6] the pseudohyperbolicity conditions were expressed in a different (equivalent) form, as explicitly verifiable conditions that ensure the hyperbolicity of the Poincaré map $T$. Similar hyperbolicity conditions were checked for the classical Lorenz model

$$
\begin{equation*}
\dot{x}=-10(x-y), \quad \dot{y}=x(28-z)-y, \quad \dot{z}=-\frac{8}{3} z+x y \tag{5}
\end{equation*}
$$

by W.Tucker, with the use of rigorous numerics. In this way, in [11, 12], a computer assisted proof of the existence of the Lorenz attractor in system (5) was done.

### 1.3. Shilnikov criterion

In our approach to the proof of the existence of the Lorenz attractor in the MoriokaShimizu system we also rely on the computer assistance, however we need much less computations. We employ a criterion proposed by Shilnikov in [13] that allows to conclude the existence of the Lorenz attractor by examination of the behavior of only one orbit (a separatrix) of the system of differential equation, instead of the direct check of the hyperbolicity of the Poincaré map $T$ which would require a high precision computation of a huge number of orbits.

In [13], several criteria for the birth of the Lorenz attractor were proposed. We use the following one. Consider a system $X$ of differential equations in $\mathbb{R}^{n}$, which have a saddle equilibrium state $O$ with one-dimensional unstable manifold. Namely, let $\lambda_{1}, \ldots, \lambda_{n}$ are the eigenvalues of the linearized system at $O$. We assume that

$$
\lambda_{1}>0>\lambda_{2}>\operatorname{Re} \lambda_{j} \quad \text { for } j \geq 3
$$

Let the system be symmetric with respect to a certain involution $\mathcal{R}$ and that $O$ is a symmetric equilibrium, i.e., $\mathcal{R} O=O$. The eigenvectors $e_{1}$ and $e_{2}$ corresponding to the
eigenvalues $\lambda_{1}$ and $\lambda_{2}$ must be $\mathcal{R}$-invariant. We assume that $\mathcal{R} e_{1}=-e_{1}$ and $\mathcal{R} e_{2}=e_{2}$. This, in particular, implies that the two unstable separatrices $\Gamma_{+}$and $\Gamma_{-}$, which are tangent at $O$ to $e_{1}$, are symmetric to each other, $\Gamma_{+}=\mathcal{R} \Gamma_{-}$.

Let the system satisfy the following 3 conditions:

1. Assume that both separatrices $\Gamma_{+}$and $\Gamma_{-}$return to O as $t \rightarrow+\infty$ and are tangent to the direction $e_{2}$ when entering $O$ (it follows from the symmetry that they are tangent to each other as they enter $O$ ).
2. Assume that the so-called saddle value $\sigma=\lambda_{1}+\lambda_{2}$ is zero.
3. Assume that the so-called separatrix value $A$ satisfies the condition

$$
\begin{equation*}
0<|A|<2 . \tag{6}
\end{equation*}
$$

The definition of the separatrix value for our case is given in Section 1.5, see (14), (15). One can describe $|A|$ as the maximal extent, to which infinitesimal two-dimensional areas can be expanded by the system along the homoclinic loop (the sign of $A$ determines whether the orientation of the areas for which this maximal expansion is achieved is changed during the propagation along the loop); more about the definition of $A$ can be seen in [14, 15].

According to [13], bifurcations of systems satisfying the above described conditions lead to the birth of a Lorenz attractor. To make this statement precise, we note that conditions 1 and 2 describe a codimension-2 bifurcation in the class of $\mathcal{R}$-symmetric systems. Suppose the system is embedded into a two-parameter family of systems $X_{\mu, \varepsilon}$ such that by changing the parameters $\mu$ and $\varepsilon$ we can independently vary the saddle value $\sigma$ near zero and split the homoclinic loop $\Gamma_{+}$(by the symmetry, the homoclinic loop $\Gamma_{-}$will be split as well). Then we have

Theorem 2. [13] If condition (6) holds at the bifurcation moment (when the system $X_{\mu, \varepsilon}$ has a homoclinic butterfly with a zero saddle value), then there exists an open region in the plane of parameters $\mu, \varepsilon$, for which system $X_{\mu, \varepsilon}$ has a Lorenz attractor of the Afraimovich-Bykov-Shilnikov model.

A proof of this result was given by Robinson in [14] under certain additional assumptions on the eigenvalues $\lambda$; in full generality this theorem was proven in [10].

We show in Theorems 5 and 6 that Shimizu-Morioka system (1) satisfies conditions 1-3 for some value of parameters $(\alpha, \lambda)$ and that the homoclinic loops can be split and the saddle value can be varied independently when $\alpha$ and $\lambda$ vary. This, as explained, implies our main result:

Theorem 3. There exists an open set in the plane of parameters $(\alpha, \lambda)$ for which the Shimizu-Morioka system has a Lorenz attractor.

### 1.4. Discrete Lorenz attractor in local and global bifurcations.

The main theorem can be applied to the study of local bifurcations of triply degenerate periodic points of three-dimensional maps. For example, consider a 3D Hénon-like map

$$
\begin{align*}
& (x, y, z) \mapsto(\bar{x}, \bar{y}, \bar{z}) \\
& \quad \bar{x}=y, \quad \bar{y}=z, \quad \bar{z}=M_{1}+B x+M_{2} y-z^{2} \tag{7}
\end{align*}
$$

where $M_{1}, M_{2}$ and $B$ are parameters and $(x, y, z) \in \mathbb{R}^{3}$. Pictures numerically obtained in [16] show that this map has a strange attractor which looks very similar to the Lorenz attractor, even though this is a discrete dynamical system and not a system of differential equations. The explanation to this fact (the emergence of a discrete analogue of a Lorenz attractor) can be obtained based on the following observation.

At $M_{1}=-1 / 4, M_{2}=1, B=1$ this map has a fixed point at $x=y=z=1 / 2$, and this point has all three multipliers (the eigenvalues of the linearization matrix at this point) on the unit circle: $(-1,-1,1)$. This is a codimension-3 bifurcation and, as shown in [16], the flow normal form for this bifurcation in this map is given by the Shimizu-Morioka system. As the normal form has a Lorenz attractor, the map itself has a discrete version of it.

More precisely, at $\left(M_{1}, M_{2}, B\right)$ close to $(-1 / 4,1,1)$ and $\left(B+M_{1}-1\right)^{2}+4 M_{1}>0$ we can shift the coordinate origin to the fixed point $x=y=z=x^{*}=\left(B+M_{2}-1+\right.$ $\left.\sqrt{\left(B+M_{2}-1\right)^{2}+4 M_{1}}\right) / 2$. Then the map will take the form

$$
\begin{equation*}
\bar{x}=y, \quad \bar{y}=z, \quad \bar{z}=B x+M_{2} y-2 x^{*} z-z^{2} . \tag{8}
\end{equation*}
$$

Introduce small parameters $\varepsilon_{1}=1-B, \varepsilon_{2}=1-M_{2}, \varepsilon_{3}=2 x^{*}-1$. It was shown in [16] that in the region $\varepsilon_{1}>0, \varepsilon_{1}+\varepsilon_{3}>\varepsilon_{2}$, the second iteration of (8) near the origin is $O\left(s^{2}\right)$-close, in appropriately chosen rescaled coordinates, to the time-s shift by the flow of a system, which is $O(s)$-close to the Shimizu-Morioka system (1) with $\alpha=\left(\varepsilon_{1}+\varepsilon_{2}+\varepsilon_{3}\right) /(4 s), \lambda=\left(\varepsilon_{1}-\varepsilon_{3}\right) /(2 s), s=\sqrt{\left(\varepsilon_{1}-\varepsilon_{2}+\varepsilon_{3}\right) / 2}$.

This means that if we make $N$ iterations of map (8), such that $N$ is even and of order $s^{-1}$, the result will be $O(s)$-close to the time- 1 map of the Shimizu-Morioka system. In other words, the square of map (8) is the Poincare map for a certain small, time-periodic perturbation of the Shimizu-Morioka system. According to [9], the pseudohyperbolicity persists at small time-periodic perturbations. Therefore, as by Theorem 3 the Shimizu-Morioka system has a pseudohyperbolic attractor for an open set of $(\alpha, \lambda)$ values, the same holds true for the map (8) for the corresponding set of values of $\varepsilon_{1,2,3}$. This attractor is called a discrete Lorenz attractor because it shape is similar to the Lorenz attractor in the Shimizu-Morioka system. However, its structure is much more complicated than that of the Lorenz attractor described in section 1.2 , In particular, the discrete Lorenz attractor may contain homoclinic tangencies [9] and heterodimensional cycles involving saddle periodic orbits with different dimensions of the unstable manifold [17] (more discussions can be found in [19]). This makes a complete description of the dynamics of the discrete Lorenz attractor impossible but, anyway, its pseudohyperbolicity allows to conclude that every orbit in this attractor behaves chaotically and this property persists for all small perturbations.

Since the map (8) is obtained from (7) just by a change of coordinates, we have the following

Theorem 4. There exists an open set in the space of parameters $\left(M_{1}, M_{2}, B\right)$ for which the 3D Hénon map (7) has a pseudohyperbolic discrete Lorenz attractor.

Note that the same conclusion holds for larger classes of three-dimensional maps. It was shown in [19] that the normal form for the bifurcations of the zero fixed point of any map of the type

$$
\begin{aligned}
& \bar{x}=y \\
& \bar{y}=z \\
& \bar{z}=\left(1-\varepsilon_{1}\right) x+\left(1-\varepsilon_{2}\right) y-\left(1+\varepsilon_{3}\right) z+a y^{2}+b y z+c z^{2}+O\left(\|x, y, z\|^{3}\right)
\end{aligned}
$$

is the Shimizu-Morioka system, provided the condition

$$
\begin{equation*}
(c-a)(a-b+c)>0 \tag{9}
\end{equation*}
$$

is fulfilled. Therefore, by the same arguments as for the map (8), Theorem 3 provides a formal justification for the claim of [19] (see Lemma 3.1 there) about the existence of pseudohyperbolic attractors for an open set of parameters $\varepsilon_{1,2,3}$ in maps (9) whose coefficients satisfy condition (9).

Map (8) is particularly important for the theory of global bifurcations because it appears as a normal form (i.e., a very close - as close as we want - approximation in appropriately chosen coordinates) for the first-return maps near many types of homoclinic tangencies and heteroclinic cycles with tangencies in three-dimensional and higher-dimensional maps [18]. Therefore, by the robustness of the pseudohyperbolicity property, Theorem 4 implies that pseudohyperbolic attractors emerge near these types of homoclinic and heteroclinic cycles. Thus, Theorem 4 provides, finally, a formal proof to the results of [20, 21, 22, 23, 24, 25] about the birth of pseudohyperbolic (discrete Lorenz) attractors at various global bifurcations, including the birth of infinitely many coexisting pseudohyperbolic attractors, see [21, 22].

### 1.5. Main results of the rigorous numerics

We establish Theorem 3 in the following steps. First, we find good bounds for the values of parameters $(\alpha, \lambda)$ for which Shimizu-Morioka system (1) conditions 1 and 2 of section 1.3 (the existence of a homoclinic butterfly with zero saddle value $\sigma$ ). It is easy to check that condition $\sigma=0$ reads as

$$
\begin{equation*}
\lambda=\frac{1}{\alpha}-\alpha . \tag{10}
\end{equation*}
$$

We will study how the separatrices $\Gamma_{ \pm}$move as $\alpha$ varies while $\lambda$ is given by 10 ; the moment a separatrix forms a homoclinic loop corresponds to the parameter values we are looking for.

It is convenient to scale time $t \rightarrow t / \alpha$ and variables $x \rightarrow X \sqrt{\alpha}, y \rightarrow Y \alpha \sqrt{\alpha}$. Then the system takes the form

$$
\begin{align*}
\dot{X} & =Y \\
\dot{Y} & =(a+1)(1-Z) X-a Y  \tag{11}\\
\dot{Z} & =-Z+X^{2}
\end{align*}
$$

where

$$
a=\frac{1}{\alpha^{2}}-1=\lambda / \alpha
$$

The eigenvalues of the linearization matrix at $(0,0,0)$ are $(-1,1,-(1+a))$ with corresponding eigenvectors given by $(0,0,1),(1,1,0)$ and $\left(-(1+a)^{-1}, 1,0\right)$, respectively. We will investigate (11) for $a \approx 1.72$, so the hyperbolic equilibrium state ( $0,0,0$ ) has a one dimensional unstable manifold, and a two dimensional stable manifold. When the manifolds intersect, we have a solution that tends to zero both as $t \rightarrow+\infty$ and $t \rightarrow-\infty$, i.e., a homoclinic loop to the equilibrium state.
Theorem 5. There exists $a=a_{0} \in\left[a_{l}, a_{r}\right]$, where

$$
\begin{aligned}
& a_{l}=1.72432329151541-10^{-13}, \\
& a_{r}=1.72432329151541+10^{-13},
\end{aligned}
$$

for which (11) has a homoclinic orbit $X_{0}(t), Y_{0}(t), Z_{0}(t)$ to the equilibrium state $(0,0,0)$. As a changes between $a_{l}$ and $a_{r}$, the loop splits. Taking

$$
\xi=1-10^{-4}, \quad c=3.5, \quad T=26
$$

we have the bound:

$$
\left|X_{0}(T+t)\right|,\left|Y_{0}(T+t)\right|,\left|Z_{0}(T+t)\right| \leq c e^{-\xi t}\left\|\left(X_{0}(T), Y_{0}(T), Z_{0}(T)\right)\right\|,
$$

for $t \geq 0$, and for $t \leq 0$ we have:

$$
\left|X_{0}(t)\right|,\left|Y_{0}(t)\right|,\left|Z_{0}(t)\right| \leq c e^{-\xi|t|}\left\|\left(X_{0}(0), Y_{0}(0), Z_{0}(0)\right)\right\|
$$

The proof of the above theorem is obtained with a computer assistance and is given in section 3.2. The mere existence of a homoclinic loop in system (11) can be obtained purely analytically [26]. However, we need good estimates on $a_{0}$ and the corresponding homoclinic solution, which the methods of [26] do not provide. Crucially, these estimates are used in the next theorem where the separatrix value $A$ is estimated.

The separatrix value can be defined as follows (cf. [14, 26]). Let a system of three differential equations have a homoclinic solution $\left(X_{0}(t), Y_{0}(t), Z_{0}(t)\right)$ to a hyperbolic equilibrium state at zero, so $\left(X_{0}(t), Y_{0}(t), Z_{0}(t)\right) \rightarrow 0$ as $t \rightarrow \pm \infty$. Let

$$
\frac{d}{d t}\left(\begin{array}{l}
x  \tag{12}\\
y \\
z
\end{array}\right)=B(t)\left(\begin{array}{l}
x \\
y \\
z
\end{array}\right)
$$

be the linearization of the system along the loop. In particular, in the case of system (11) we have

$$
B=\left(\begin{array}{ccc}
0 & 1 & 0 \\
\left(a_{0}+1\right)\left(1-Z_{0}(t)\right) & -a_{0} & -\left(a_{0}+1\right) X_{0}(t) \\
2 X_{0}(t) & 0 & -1
\end{array}\right)
$$

Let $\xi_{1}, \xi_{2}$ be any two vectors and let $\eta=\xi_{1} \times \xi_{2}$ be their vector product. If the evolution of $\xi_{1}$ and $\xi_{2}$ is defined by (12), then the evolution of $\eta$ is governed by

$$
\begin{equation*}
\frac{d \eta}{d t}=-\left(B^{\top}-\operatorname{tr}(B) I\right) \eta \tag{13}
\end{equation*}
$$

where $I$ is the $(3 \times 3)$ identity matrix. This equation describes the evolution of infinitesimal two-dimensional areas near the homoclinic loop.

Since $\left(X_{0}(t), Y_{0}(t), Z_{0}(t)\right)$ tends to zero exponentially, the asymptotic behavior of solutions of (13) as $t \rightarrow \pm \infty$ is determined by the limit matrix

$$
\hat{B}=-\left(B_{\infty}^{\top}-\operatorname{tr}\left(B_{\infty}\right) I\right)
$$

where $B_{\infty}=\lim _{t \rightarrow \pm \infty} B(t)$, which is the linearization of the original system at the hyperbolic equilibrium at zero. If $\lambda_{1}>0>\lambda_{2}>\lambda_{3}$ are the eigenvalues of $B_{\infty}$, then the eigenvalues of $\hat{B}$ are

$$
\sigma_{1}=\lambda_{1}+\lambda_{2}, \quad \sigma_{2}=\lambda_{1}+\lambda_{3}, \quad \sigma_{3}=\lambda_{2}+\lambda_{3}
$$

We are interested here in the case of zero saddle value, i.e., $\lambda_{1}+\lambda_{2}=0$. Then the eigenvalues of $\hat{B}$ are $0, \sigma_{2}<0$ and $\sigma_{3}<0$, so every solution of 13 tends, as $t \rightarrow+\infty$ to a constant times the eigenvector of $\hat{B}$ that corresponds to the zero eigenvalue (this is the vector $v_{0}$, which is the vector product of the eigenvectors of $B_{\infty}$ that correspond to the eigenvalues $\lambda_{1}$ and $\lambda_{2}$ ). It also follows that only one solution of (13) tends to $v_{0}$ in backward time, as $t \rightarrow-\infty$. We take this particular solution $\eta_{0}(t)$ and denote

$$
\begin{equation*}
\lim _{t \rightarrow+\infty} \eta_{0}(t)=A v_{0} \tag{14}
\end{equation*}
$$

The coefficient $A$ is the sought separatrix value. One can see that

$$
\begin{equation*}
|A|=\sup \lim _{t \rightarrow+\infty} \frac{\|\eta(t)\|}{\|\eta(-t)\|} \tag{15}
\end{equation*}
$$

where the supremum is taken over all the solutions of (13). Thus, $A$ determines the maximal expansion of infinitesimal areas along the homoclinic loop.

In the case of system (11) system (13) becomes

$$
\eta^{\prime}(t)=\left(\begin{array}{ccc}
-\left(a_{0}+1\right) & -\left(a_{0}+1\right)\left(1-Z_{0}(t)\right) & -2 X_{0}(t)  \tag{16}\\
-1 & -1 & 0 \\
0 & \left(a_{0}+1\right) X_{0}(t) & -a_{0}
\end{array}\right) \eta(t)
$$

This equation, in the limit $t \rightarrow \pm \infty$ becomes

$$
\eta^{\prime}(t)=\hat{B} \eta(t)=\left(\begin{array}{ccc}
-\left(a_{0}+1\right) & -\left(a_{0}+1\right) & 0  \tag{17}\\
-1 & -1 & 0 \\
0 & 0 & -a_{0}
\end{array}\right) \eta(t)
$$

The eigenvalues are $0,-(a+2),-a$ with corresponding eigenvectors $v_{0}=(1,-1,0)$, $v_{-(a+2)}=(1+a, 1,0)$ and $v_{-a}(0,0,1)$, respectively.
Theorem 6. There exists an orbit $\eta(t)$ of (16), for which

$$
\begin{aligned}
\lim _{t \rightarrow-\infty} \eta(t) & \in[0.99984336210766,1.0001566378923] v_{0} \\
\lim _{t \rightarrow+\infty} \eta(t) & \in[0.62606812264791,0.62663392848044] v_{0}
\end{aligned}
$$

meaning that

$$
\begin{equation*}
\lim _{t \rightarrow+\infty} \frac{\|\eta(t)\|}{\|\eta(-t)\|} \in[0.62597007201516,0.6267320984754] \tag{18}
\end{equation*}
$$

The proof of the above theorem is obtained with computer assistance and is given in Section 4.3. By definition (14) of the separatrix value, estimate (18) gives the following bounds for the separatrix value of the homoclinic loop in system (11):

$$
A \in[0.62597007201516,0.6267320984754] .
$$

Importantly $0<A<1$. So, by applying Shilnikov criterion we obtain our main result, Theorem 3, see section 1.3 .

The plot of the homoclinic trajectory from Theorem 5 can be seen on the left plot in Figure 6 on page 18. The plot of the heteroclinic orbit from Theorem 6 is given in Figure 12 on page 28 .

## 2. Some notations

In the subsequent sections we present a methodology for establishing homoclinic orbits and for the computation of the separatrix value. First we introduce the following notations.

We will write $B_{k}(R)$ for ball in $\mathbb{R}^{k}$ of radius $R$, centered at zero. For a matrix $A$ we define the logarithmic norm of $A$ as (see [27] and literature cited there)

$$
l(A)=\lim _{h \rightarrow 0^{+}} \frac{\|I+A h\|-\|I\|}{h} .
$$

We will also use the notation

$$
\begin{aligned}
& m(A)= \begin{cases}\frac{1}{\|A\|^{-1}} & \operatorname{det} A \neq 0 \\
0 & \text { otherwise }\end{cases} \\
& m_{l}(A)=-l(-A) .
\end{aligned}
$$

The $m(A)$ is a number with the property that $\|A v\| \geq m(A)\|v\|$. The $m_{l}(A)$ can be interpreted as a 'bound from below' of the logarithmic norm.

Let $u, s \in \mathbb{N}$. For a function $f: \mathbb{R}^{u} \times \mathbb{R}^{s} \rightarrow \mathbb{R}^{u} \times \mathbb{R}^{s}$, we will use the notations $(x, y) \in \mathbb{R}^{u} \times \mathbb{R}^{s}$, where $x \in \mathbb{R}^{u}$ and $y \in \mathbb{R}^{s}$. The $x$ will play the role of an "unstable" coordinate and $y$ will be the "stable" coordinate, hence the choice of the notation $u, s$. We will also write $f_{x}$, and $f_{y}$ for the projections of $f$ onto $\mathbb{R}^{u}$ and $\mathbb{R}^{s}$, respectively. For a set $D \subset \mathbb{R}^{u} \times \mathbb{R}^{s}$ we define

$$
\begin{aligned}
& {\left[\frac{\partial f_{x}}{\partial x}(D)\right]:=\left\{A=\left(a_{i j}\right) \in \mathbb{R}^{u \times u}: a_{i, j} \in\left[\inf _{p \in D} \frac{\partial f_{x_{i}}}{\partial x_{j}}(p), \sup _{p \in D} \frac{\partial f_{x_{i}}}{\partial x_{j}}(p)\right]\right\}} \\
& {\left[\frac{\partial f_{x}}{\partial y}(D)\right]:=\left\{A=\left(a_{i j}\right) \in \mathbb{R}^{u \times s}: a_{i, j} \in\left[\inf _{p \in D} \frac{\partial f_{x_{i}}}{\partial y_{j}}(p), \sup _{p \in D} \frac{\partial f_{x_{i}}}{\partial y_{j}}(p)\right]\right\}}
\end{aligned}
$$

We also define

$$
\begin{aligned}
m\left(\frac{\partial f_{x}}{\partial x}(D)\right) & :=\inf _{A \in\left[\frac{\partial f_{x}}{\partial x}(D)\right]} m(A) \\
m_{l}\left(\frac{\partial f_{x}}{\partial x}(D)\right) & :=\inf _{A \in\left[\frac{\partial f_{x}}{\partial x}(D)\right]} m_{l}(A)
\end{aligned}
$$



Figure 2. The local unstable manifold $W_{\theta}^{u}$ in red, and the local stable manifold $W_{\theta}^{s}$ in green.

$$
\left\|\frac{\partial f_{x}}{\partial y}(D)\right\|:=\sup _{A \in\left[\frac{\partial f_{x}}{\partial y}(D)\right]}\|A\| .
$$

We will use the notation $\operatorname{int}(D), \bar{D}$ and $\partial D$ for the interior, closure and boundary of a set $D$, respectively.

## 3. Establishing homoclinics

In this section we give an overview of the method for establishing the existence of homoclinic orbits to fixed points. The method is written for the case where we consider a parameter dependent ODE with the vector field $f: \mathbb{R}^{3} \times \mathbb{R} \rightarrow \mathbb{R}^{3}$,

$$
\begin{equation*}
p^{\prime}=f(p, a) \tag{19}
\end{equation*}
$$

and $a \in A$ is a parameter, with $A=\left[a_{l}, a_{r}\right] \subset \mathbb{R}$. We assume that for each $a \in A$ (19) has a hyperbolic fixed point $p_{a}^{*}$, with one dimensional unstable manifold and two dimensional stable manifolds. (If dimensions are the other way around we can change the sign of the vector field.)

Let $\Phi_{t}(p, a)$ be the flow induced by 19). Let $\bar{B}_{u}(R)=[-R, R] \subset \mathbb{R}, \bar{B}_{s}(R) \subset \mathbb{R}^{2}$ and let

$$
D=\bar{B}_{u}(R) \times \bar{B}_{s}(R) \subset \mathbb{R}^{3},
$$

be a neighborhood of the smooth family of fixed points, meaning that we assume $p_{a}^{*} \in \operatorname{int} D$ for any $a \in A$.

We denote by $W_{a}^{u}$ the local unstable manifold of $p_{a}^{*}$ in $D$ and by $W_{a}^{s}$ the local stable manifold of $p_{a}^{*}$ in $D$, i.e.

$$
\begin{align*}
& W_{a}^{u}=\left\{p \in D: \Phi_{t}(p, a) \in D \text { for } t \leq 0 \text { and } \lim _{t \rightarrow-\infty} \Phi_{t}(p, a)=p_{a}^{*}\right\},  \tag{20}\\
& W_{a}^{s}=\left\{p \in D: \Phi_{t}(p, a) \in D \text { for } t \geq 0 \text { and } \lim _{t \rightarrow+\infty} \Phi_{t}(p, a)=p_{a}^{*}\right\} . \tag{21}
\end{align*}
$$

We assume that $W_{a}^{u}$ and $W_{a}^{s}$ are graphs of $C^{1}$ functions

$$
\begin{aligned}
w_{a}^{u}: \bar{B}_{u}(R) & \rightarrow \bar{B}_{s}(R), \\
w_{a}^{s}: \bar{B}_{s}(R) & \rightarrow \bar{B}_{u}(R),
\end{aligned}
$$



Figure 3. We have the 1-dimensional unstable manifold of $p_{a}^{*}$ in red, and the 2dimensional local stable manifold $W_{a}^{s}$ in $D$ in green. The $h(a)$ is the signed distance along the $x$ coordinate between $W_{a}^{s}$ and $\Phi_{T}\left(p_{a}^{u}, a\right)$; this is the distance along the dotted line on the plot.
meaning that (see Figure 2)

$$
\begin{align*}
W_{a}^{u} & =\left\{\left(x, w_{a}^{u}(x)\right): x \in \bar{B}_{u}(R)\right\}, \\
W_{a}^{s} & =\left\{\left(w_{a}^{s}(y), y\right): y \in \bar{B}_{s}(R)\right\} . \tag{22}
\end{align*}
$$

Let

$$
\begin{equation*}
p_{a}^{u}:=\left(R, w_{a}^{u}(R)\right) \in \mathbb{R}^{3} . \tag{23}
\end{equation*}
$$

Consider $T>0$ and assume that for all $a \in A, \Phi_{T}\left(p_{a}^{u}, a\right) \in D$. Let us define

$$
h: A \rightarrow \mathbb{R}
$$

as

$$
\begin{equation*}
h(a)=\pi_{x} \Phi_{T}\left(p_{a}^{u}, a\right)-w_{a}^{s}\left(\pi_{y} \Phi_{T}\left(p_{a}^{u}, a\right)\right) \tag{24}
\end{equation*}
$$

We now state a natural result, that $h(a)=0$ implies an intersection of the stable and unstable manifolds of $p_{a}^{*}$. (See Figure 3.)
Theorem 7. [28]If

$$
\begin{equation*}
h\left(a_{l}\right)<0 \quad \text { and } \quad h\left(a_{r}\right)>0 \tag{25}
\end{equation*}
$$

then there exists a $\psi \in A$ for which we have a homoclinic orbit to $p_{\psi}^{*}$.

### 3.1. Computer assisted bounds for unstable manifolds of fixed points of ODEs

In order to apply Theorem 7, we need to be able to establish bounds for $h$ defined in (24). Using a rigorous, interval arithmetic based integrator|l|, it is possible to obtain rigorous enclosures for $\Phi_{T}$. In this section we discuss how to obtain bounds on parameterizations $w_{a}^{s}, w_{a}^{u}$ of stable and unstable manifolds of fixed points. For simplicity, we will skip the parameter $a$ and consider an ODE

$$
\begin{equation*}
p^{\prime}=f(p) . \tag{26}
\end{equation*}
$$

|| in our application we use the CAPD package: http://capd.ii.uj.edu.pl/

When combined with Theorem 7, we can apply below results for (19) with fixed $a$.
Let $D \subset \mathbb{R}^{u} \times \mathbb{R}^{s}$,

$$
D=\bar{B}_{u}(R) \times \bar{B}_{s}(R) .
$$

In this section we do not need to assume that $u=1$ and $d=2$. The result works for arbitrary dimensions.

We define

$$
\begin{align*}
\vec{\mu} & =\sup _{z \in D}\left\{l\left(\frac{\partial f_{y}}{\partial y}(z)\right)+\frac{1}{L}\left\|\frac{\partial f_{y}}{\partial x}(z)\right\|\right\}  \tag{27}\\
\vec{\xi} & =m_{l}\left(\frac{\partial f_{x}}{\partial x}(D)\right)-L\left\|\frac{\partial f_{x}}{\partial y}(D)\right\| \tag{28}
\end{align*}
$$

Definition 8. We say that the vector field $f$ satisfies rate conditions if

$$
\begin{equation*}
\vec{\mu}<0<\vec{\xi} \tag{29}
\end{equation*}
$$

Definition 9. We say that $D=\bar{B}_{u}(R) \times \bar{B}_{s}(R)$ is an isolating block for (26) if
(i) For any $q \in \partial \bar{B}_{u}(R) \times \bar{B}_{s}(R)$,

$$
\left(\pi_{x} f(q) \mid \pi_{x} q\right)>0
$$

(ii) For any $q \in \bar{B}_{u}(R) \times \partial \bar{B}_{s}(R)$,

$$
\left(\pi_{y} f(q) \mid \pi_{y} q\right)<0
$$

Definition 10. We define the unstable set in $D$ as

$$
W^{u}=\left\{z: \Phi_{t}(z) \in D \text { for all } t<0\right\}
$$

Below theorem is a simplified (adapted to fixed points) version of the results from [27, Theorem 30]. The paper [27] is in the setting of normally hyperbolic invariant manifolds, hence the Theorem 31 from that paper is more involved than below result.
Theorem 11. Let $k \geq 1$. Assume that $f$ is $C^{1}$ and satisfies the rate conditions. Assume also that $D=\bar{B}_{u}(R) \times \bar{B}_{s}(R)$ is an isolating segment for $f$. Then the set $W^{u}$ is a manifold, which is a graph over $\bar{B}_{u}(R)$. To be more precise, there exists a $C^{1}$ function

$$
w^{u}: \bar{B}_{u}(R) \rightarrow \bar{B}_{s}(R),
$$

such that

$$
\begin{equation*}
W^{u}=\left\{\left(x, w^{u}(x)\right): x \in \bar{B}_{u}(R)\right\} . \tag{30}
\end{equation*}
$$

Moreover, $w^{u}$ is Lipschitz with constant L.
Proof. The proof is given in Appendix 5.1. It is a modification of the argument from [27]. The main difference is that the results from [27] are for the setting in which in addition to the hyperbolic coordinates $x, y$ we have a centre coordinate. Due to this the Lipschitz bound from Theorem 11 is sharper compared with [27]. In the proof we focus on this issue.

We will now discuss the contraction rate along the stable manifold $W^{u}$ from Theorem 11. First we shall need an auxiliary result. Consider a flow $\Phi_{t}(z)$, for $\Phi: \mathbb{R} \times \mathbb{R}^{u} \times \mathbb{R}^{s} \rightarrow \mathbb{R}^{u} \times \mathbb{R}^{s}$,

$$
\Phi_{t}(x, y)=\left(\pi_{x} \Phi_{t}(x, y), \pi_{y} \Phi_{t}(x, y)\right)
$$

and define the following constants

$$
\begin{align*}
& \mu(h)=\sup _{z \in U}\left\{\left\|\frac{\partial \pi_{y} \Phi_{t}}{\partial y}(h, z)\right\|+\frac{1}{L}\left\|\frac{\partial \pi_{y} \Phi_{t}}{\partial x}(h, z)\right\|\right\}  \tag{31}\\
& \xi(h)=m\left[\frac{\partial \pi_{x} \Phi_{t}}{\partial x}(h, U)\right]-L \sup _{z \in U}\left\|\frac{\partial \pi_{x} \Phi_{t}}{\partial y}(h, z)\right\| \tag{32}
\end{align*}
$$

Theorem 12. [27, Theorem 31]Let $\vec{\xi}$ and $\vec{\mu}$ be the constants defined in (27) and (28). If $\Phi_{t}$ is the flow induced by $(\sqrt{26})$, then for $h>0$

$$
\begin{aligned}
& \mu(h)=1+h \vec{\mu}+O\left(h^{2}\right), \\
& \xi(h)=1+h \vec{\xi}+O\left(h^{2}\right) .
\end{aligned}
$$

Theorem 13. Let $W^{u}$ be the manifold established in Theorem 11. Then for any $p_{1}, p_{2} \in W^{u}$

$$
\left\|\Phi_{-t}\left(p_{1}\right)-\Phi_{-t}\left(p_{2}\right)\right\| \leq c e^{-\vec{\xi} t}\left\|\pi_{x}\left(p_{1}-p_{2}\right)\right\| \quad \text { for all } t \geq 0
$$

for $c=2 \sqrt{1+L^{2}}$.
Proof. Let $p_{1}, p_{2} \in W^{u}$. Since $w^{u}$ is Lipschitz with constant $L$,

$$
\begin{equation*}
\left\|\pi_{y}\left[p_{1}-p_{2}\right]\right\| \leq L\left\|\pi_{x}\left[p_{1}-p_{2}\right]\right\| \tag{33}
\end{equation*}
$$

Let $q_{1}, q_{2} \in W^{u}$. If $\Phi_{t}\left(q_{1}\right), \Phi_{t}\left(q_{2}\right) \in D$ for $t \in(0, T]$, then for $0<h \leq T$ holds

$$
\begin{aligned}
\left\|\pi_{x}\left(\Phi_{h}\left(q_{1}\right)-\Phi_{h}\left(q_{2}\right)\right)\right\| & \geq \xi(h)\left\|\pi_{x}\left(q_{1}-q_{2}\right)\right\| \\
& =\left(1+h \vec{\xi}+O\left(h^{2}\right)\right)\left\|\pi_{x}\left(q_{1}-q_{2}\right)\right\| .
\end{aligned}
$$

If we take $h=\frac{T}{N}$, then

$$
\begin{aligned}
\left\|\pi_{x}\left(\Phi_{T}\left(q_{1}\right)-\Phi_{T}\left(q_{2}\right)\right)\right\| & \geq\left(1+h \vec{\xi}+O\left(h^{2}\right)\right)^{N}\left\|\pi_{x}\left(q_{1}-q_{2}\right)\right\| \\
& \rightarrow e^{T \vec{\xi}}\left\|\pi_{x}\left(q_{1}-q_{2}\right)\right\|, \quad N \rightarrow \infty .
\end{aligned}
$$

Observe that from the above it follows that (we set $p_{i}=\Phi_{-T}\left(q_{i}\right)$ )

$$
\begin{equation*}
\left\|\pi_{x}\left(p_{1}-p_{2}\right)\right\| \geq e^{T} \overrightarrow{\vec{\xi}}\left\|\pi_{x}\left(\Phi_{-T}\left(p_{1}\right)-\Phi_{-T} p_{2}\right)\right\| \tag{34}
\end{equation*}
$$

for any $T>0$.
Using (33) in the third line and (34) in the last line,

$$
\begin{aligned}
& \left\|\Phi_{-T}\left(p_{1}\right)-\Phi_{-T}\left(p_{2}\right)\right\|^{2} \\
= & \left\|\pi_{x}\left[\Phi_{-T}\left(p_{1}\right)-\Phi_{-T}\left(p_{2}\right)\right]\right\|^{2}+\left\|\pi_{y}\left[\Phi_{-T}\left(p_{1}\right)-\Phi_{-T}\left(p_{2}\right)\right]\right\|^{2} \\
\leq & \left(1+L^{2}\right)\left\|\pi_{x}\left[\Phi_{-T}\left(p_{1}\right)-\Phi_{-T}\left(p_{2}\right)\right]\right\|^{2} \\
\leq & \left(1+L^{2}\right) e^{-2 T} \vec{\xi}\left\|\pi_{x}\left[p_{1}-p_{2}\right]\right\|^{2},
\end{aligned}
$$

which concludes the proof.

Remark 14. Theorems 11, 13 can also be applied to establish bounds on the stable manifold. In order to do so, it is enough to consider $p^{\prime}=-f(p)$ instead of 26), and to swap the roles of the coordinates $x, y$. The unstable manifold for the vector field $-f$ is the stable manifold for $f$.

### 3.2. Computer assisted proof of homoclinic intersection

In this section we give an overview of the computer assisted proof of Theorem 5 .
To apply the method from Sections 3, 3.1 to conduct a computer assisted proof we follow the steps, as outlined in [28]:

## Algorithm 1.

(i) In local coordinates around zero, using Theorem 11, establish the bounds on the unstable manifolds for the family of vector fields .
(ii) By changing sign of the vector field, using the same procedure as in stepi, establish bounds on the stable manifolds.
(iii) Propagate the bounds on the unstable manifold along the flow, and establish the homoclinic intersection using Theorem 7 .

To obtain bounds for the stable/unstable manifolds, we use the local coordinates $\left(x, y_{1}, y_{2}\right)$,

$$
\begin{equation*}
(X, Y, Z)=C\left(x, y_{1}, y_{2}\right) \tag{35}
\end{equation*}
$$

with,

$$
C=\left(\begin{array}{ccc}
1 & -0.36706363121968 & 0 \\
1 & 1 & 0 \\
0 & 0 & 1
\end{array}\right)
$$

Note that $\|C\| \approx 1.5278$ and $\left\|C^{-1}\right\| \approx 1.1175$.
Coordinates $x, y_{1}, y_{2}$ align the system (11) so that $x$ is the (rough) unstable direction, and $y_{1}, y_{2}$ are (roughly) stable. Note that we use the same local coordinates for all the parameters $a$.

To obtain the bound on $W^{u}$ we choose

$$
D=\bar{B}_{u}(R) \times \bar{B}_{s}(R),
$$

with $R=10^{-5}$, and use Theorem 11 to obtain an enclosure of the unstable manifold $W^{u}$. In our computer assisted proof, we have a Lipschitz bound $L=4 \cdot 10^{-5}$ for the slope of the manifold for all parameters $a \in A$. See Figure 4. (Note the scale on the axes. The enclosure is in fact quite sharp.) Applying Theorem 13 we obtain

$$
\begin{aligned}
& \xi_{1}=0.99999999967813 \\
& c=2 \sqrt{1+L^{2}}
\end{aligned}
$$



Figure 4. The projection onto $x, y_{1}$ coordinates of the bounds on $W_{a}^{u}$.

For any trajectory $(X(t), Y(t), Z(t))$ starting from a point $p=C\left(x, w^{u}(x)\right)$, for $t \leq 0$, holds

$$
|X(t)|,|Y(t)|,|Z(t)| \leq c\|C\|\left\|C^{-1}\right\| e^{-\xi_{1}|t|}\|p\|<3.5 e^{-\xi|t|}\|p\|
$$

To establish the bounds for the two dimensional stable manifold $W^{s}$, we consider the vector field (11) with reversed sign (this means that we also swap the roles of the coordinates $x, y_{1}, y_{2}$ ), and apply Theorem 11 once again. We use the same Lipschitz bound $L$. We consider all the parameters $a \in A$. We do so by subdividing $A$ into several intervals, and performing the interval arithmetic enclosure of $W^{s}$; we use the interval enclosure of the map $f$ with the intervals on $a$. In Figure 5 we see the bound on the enclosure. From Theorem 13 we obtain

$$
\xi_{2}=0.99998045374688
$$

Thus, for any trajectory $(X(t), Y(t), Z(t))$ starting from a point $p=C\left(w^{s}\left(y_{1}, y_{2}\right), y_{1}, y_{2}\right)$ holds

$$
|X(t)|,|Y(t)|,|Z(t)| \leq c\|C\|\left\|C^{-1}\right\| e^{-\xi_{2} t}\|p\|<3.5 e^{-\xi t}\|p\| \quad \text { for } t \geq 0
$$

We now take $T=26$. The two rectangles in Figure 5 are the $\Phi_{T}\left(C p_{a}^{u}, a\right)$ for $a=a_{l}$ and $a=a_{r}$ (see (23) for the definition of $p_{u}^{u}$ ). Note that Figure 5 corresponds to the sketch from Figure 3. In Figure 5 we have the projection onto $x, y_{1}$ coordinates of what happens inside of the set $D$, without plotting the trajectory along the unstable manifold. Figure 6 presents the bounds, plotted in the original coordinates of the system.

We use the rigorous estimates for $\Phi_{T}\left(C p_{a_{l}}^{u}, a_{l}\right)$ and $\Phi_{T}\left(C p_{a_{r}}^{u}, a_{r}\right)$ to compute the following bounds (see (24) for the definition of the function $h$, )

$$
\begin{aligned}
& h\left(a_{l}\right) \in[1.5093787863274 e-09,3.9653443827625 e-09], \\
& h\left(a_{r}\right) \in[-3.9570292285809 e-09,-1.51777586447 e-09] .
\end{aligned}
$$

We also make sure that $\Phi_{T}\left(C p_{a}^{u}, a\right) \in D$ for all $a \in A$. We see that assumption (25) of Theorem 7 is satisfied, which means that we have a homoclinic connection for at least one of the parameters $a \in A$.


Figure 5. The bound on $W_{a}^{s}$, for all parameters $a$. On the left we have a non-rigorous plot, to illustrate the shape of our bound in three dimensions. In the middle and on the right, we have a projection onto the $x, y_{1}$ coordinates of the rigorous, computer assisted enclosure. The two rectangles depicted on the right hand side plots are $\Phi_{T}\left(p_{a_{l}}^{u}, a_{l}\right)$ (on the left, in red) and $\Phi_{T}\left(p_{a_{r}}^{u}, a_{r}\right)$ (on the right, in blue). All these plots are in the local coordinates.


Figure 6. The bound on $W_{a}^{u}$, for all parameters $a \in\left[a_{l}, a_{r}\right]$ on the left (in black). On the right, we have the bound on $W_{a}^{u}$ in black, valid for all $a \in\left[a_{l}, a_{r}\right]$. A trajectory along $W_{a}^{u}$, which leaves the neighbourhood by the top right corner, returns through the lower edge of the plot. In red, we have the $W_{a_{l}}^{u}$, as it returns to the neighbourhood, and in blue we have $W_{a_{r}}^{u}$. In green is the bound on $W_{a}^{s}$, for all parameters $a \in\left[a_{l}, a_{r}\right]$. The bound on $W_{a}^{s}$ in local coordinates is in fact is much tighter than on this plot. It is a strip that passes through the intersections of the green rectangles. The plots are in the original coordinates of the system, projected onto the $X, Y$ coordinates.

We have thus established a homoclinic orbit

$$
\left(X_{0}(t), Y_{0}(t), Z_{0}(t)\right)=\Phi_{t}\left(C p_{a_{0}}^{u}, a_{0}\right),
$$

for some $a_{0} \in A$.
The computer assisted proof has been done entirely by using the CAPD package and took under a second on a single core 3 Ghz Intel i7 processor.
$\dagger$ computer assisted proofs in dynamics: http://capd.ii.uj.edu.pl/

## 4. Computation of the separatrix value

In this section we will show how we prove Theorem6. First we will describe the method. We shall investigate the following ODE

$$
\begin{equation*}
\gamma^{\prime}(t)=(A+B(t)) \gamma(t), \tag{36}
\end{equation*}
$$

and assume that $a_{22}, a_{33}>0$ and that

$$
A=\left(\begin{array}{lll}
0 & 0 & 0 \\
0 & -a_{22} & 0 \\
0 & 0 & -a_{33}
\end{array}\right), \quad B(t)=\left(\begin{array}{lll}
b_{11}(t) & b_{12}(t) & b_{13}(t) \\
b_{21}(t) & b_{22}(t) & b_{23}(t) \\
b_{31}(t) & b_{32}(t) & b_{33}(t)
\end{array}\right)
$$

We assume also that there exist constants $c_{b}, \lambda>0$ for which

$$
\begin{equation*}
\left|b_{i j}(t)\right| \leq c_{b} e^{-\lambda|t|} \tag{37}
\end{equation*}
$$

Remark 15. In section 4.3, where we will apply the method to the Shimizu-Marioka system (11), the constant $c_{b}$ from (37) will be associated with the constant $c$ from Theorem 5, with the size of the neighborhood of zero in which we investigate the rate of convergence along the homoclinic, and also on some coordinate changes.
Remark 16. The (36) can be seen as (16) considered in appropriate local coordinates. This will be our approach when applying the results from this section for the proof of Theorem 6.

In the sequel in our investigations of the problem we will use the variables $\left(x, y_{1}, y_{2}\right)$ and quite often we will also write $y=\left(y_{1}, y_{2}\right)$, so that our coordinates will be $(x, y)$, where $x \in \mathbb{R}$ and $y \in \mathbb{R}^{2}$.

Our objective will be to obtain a proof of an orbit $\gamma^{*}(t)$ of (36) for which

$$
\begin{align*}
\lim _{t \rightarrow-\infty} \gamma^{*}(t) & =\left(x_{-}^{*}, 0,0\right)  \tag{38}\\
\lim _{t \rightarrow+\infty} \gamma^{*}(t) & =\left(x_{+}^{*}, 0,0\right) \tag{39}
\end{align*}
$$

and also to obtain explicit bounds on the fraction $\frac{x_{+}^{*}}{x_{-}^{*}}$.
To understand the behavior of (36) for large $|t|$, we treat the line $\{(x, 0,0), x \in \mathbb{R}\}$ as the normally hyperbolic manifold, with two stable directions ( $y_{1}$ and $y_{2}$ ). This is clearly visible in (36) with $B \equiv 0$, but for $|t|$ large $B$ is just a small perturbation so the normally hyperbolic behavior will survive. A bit problematic in this picture is the fact that we have here an non-autonomous system, so the above scenario needs some adjustments. This is the idea, which underlines our approach.
Definition 17. Let $Z \subset \mathbb{R} \times \mathbb{R}^{2}$ and $J \subset \mathbb{R}$. Let $V: Z \rightarrow \mathbb{R}$ be a $C^{1}$ function. We say that $V$ is a Lyapunov function on $J \times Z$ for (36) if for any solution $z: J \rightarrow \mathbb{R} \times \mathbb{R}^{2}$ of (36), $V(z(t))$ is decreasing $\left(\frac{d}{d t} V(z(t))<0\right)$ for $z(t) \in Z$ and $t \in J$.

We have

$$
\begin{equation*}
\frac{d}{d t} V(z(t))=((A+B(t)) z(t) \mid \nabla V(z(t))) \tag{40}
\end{equation*}
$$

We start with two technical lemmas:


Figure 7. The set $D$ from Lemma 18 is the vertical cone. The curved lines are the level sets of $V$. The arrows indicate the vector field of (36).

Lemma 18. Let $D=\left\{\|y\|^{2} \geq x^{2}\right\} \backslash\{0\}$ (see Figure 7). Let $J \subset \mathbb{R}$ be a set of all $t \in \mathbb{R}$, for which the following condition holds

$$
\begin{equation*}
6 c_{b} e^{-\lambda|t|}<\min \left(a_{22}, a_{33}\right) \tag{41}
\end{equation*}
$$

Then

$$
\begin{array}{ll}
V & : \mathbb{R} \times \mathbb{R}^{2} \rightarrow \mathbb{R} \\
V(x, y) & =\|y\|^{2}-x^{2}
\end{array}
$$

is a Lyapunov function for (36) on $J \times D$.
Proof. For $p=\left(x, y_{1}, y_{2}\right)$ we compute (skipping the dependence of $b_{i j}(t)$ on $t$ to simplify the notation):

$$
\begin{aligned}
& \left((A+B(t)) p \left\lvert\, \frac{1}{2} \nabla V(p)\right.\right) \\
= & -x^{2} b_{11}+y_{1}^{2}\left(b_{22}-a_{22}\right)+y_{2}^{2}\left(b_{33}-a_{33}\right) \\
& +x y_{1}\left(b_{21}-b_{12}\right)+x y_{2}\left(b_{31}-b_{13}\right)+y_{1} y_{2}\left(b_{23}+b_{32}\right) \\
\leq & \|y\|^{2}\left|b_{11}\right|+\left(y_{1}^{2}+y_{2}^{2}\right)\left(\max \left(\left|b_{22}\right|,\left|b_{33}\right|\right)+\max \left(-a_{22},-a_{33}\right)\right) \\
& +\frac{1}{2}\left(x^{2}+y_{1}^{2}\right)\left(\left|b_{21}\right|+\left|b_{12}\right|\right)+\frac{1}{2}\left(x^{2}+y_{2}^{2}\right)\left(\left|b_{31}\right|+\left|b_{13}\right|\right) \\
& +\frac{1}{2}\left(y_{1}^{2}+y_{2}^{2}\right)\left(\left|b_{23}\right|+\left|b_{32}\right|\right) \\
\leq & {\left[2 c_{b} e^{-\lambda|t|}+\max \left(-a_{22},-a_{33}\right)\right]\|y\|^{2} } \\
& +c_{b} e^{-\lambda|t|}\left(\left(x^{2}+y_{1}^{2}\right)+\left(x^{2}+y_{2}^{2}\right)+\left(y_{1}^{2}+y_{2}^{2}\right)\right) \\
\leq & {\left[2 c_{b} e^{-\lambda|t|}+\max \left(-a_{22},-a_{33}\right)\right]\|y\|^{2}+c_{b} e^{-\lambda|t|}\left(4\|y\|^{2}\right) } \\
= & {\left[6 c_{b} e^{-\lambda|t|}-\min \left(a_{22}, a_{33}\right)\right]\|y\|^{2} } \\
< & 0,
\end{aligned}
$$

which concludes the proof.


Figure 8. The set $D$ from Lemma 19 is represented by the two shaded rectangular strips. The horizontal lines are the level sets of $V$. The arrows indicate the vector field of 36 .

Lemma 19. Let $x_{0}, \rho, r>0$ and let

$$
D=\left[x_{0}-2 \rho, x_{0}+2 \rho\right] \times\{\|y\| \geq r\} \subset \mathbb{R} \times \mathbb{R}^{2} .
$$

Let $J$ be a set of $t \in \mathbb{R}$, such that

$$
\begin{equation*}
2 c_{b} e^{-\lambda|t|}\left(1+\frac{x_{0}+2 \rho}{r}\right)<\min \left(a_{22}, a_{33}\right), \tag{42}
\end{equation*}
$$

then

$$
\begin{array}{ll}
V & : \mathbb{R} \times \mathbb{R}^{2} \rightarrow \mathbb{R}, \\
V(x, y) & =\|y\|^{2}
\end{array}
$$

is a Lyapunov function for (36) on the set $J \times D$.
Proof. Let $p=\left(x, y_{1}, y_{2}\right)$ and let $R^{2}=y_{1}^{2}+y_{2}^{2} \geq r$. Note that $\pm y_{1} y_{2} \leq \frac{1}{2} R^{2}$ and $\left|y_{1}\right|,\left|y_{2}\right| \leq R$. We can compute (skipping the dependence of $b_{i j}(t)$ on $t$ to simplify the notation)

$$
\begin{aligned}
& \left((A+B(t)) p \left\lvert\, \frac{1}{2} \nabla V(p)\right.\right) \\
& =\left(-a_{22}+b_{22}\right) y_{1}^{2}+\left(-a_{33}+b_{33}\right) y_{2}^{2} \\
& \quad+\left(b_{23}+b_{32}\right) y_{1} y_{2}+b_{21} x y_{1}+b_{31} x y_{2} \\
& \leq \max \left(-a_{22}+b_{22},-a_{33}+b_{33}\right) R^{2} \\
& \quad+\frac{1}{2}\left(\left|b_{23}\right|+\left|b_{32}\right|\right) R^{2}+\left|b_{21}\right||x| R+\left|b_{31}\right||x| R \\
& \leq\left(\max \left(-a_{22},-a_{33}\right)+2 c_{b} e^{-\lambda|t|}\right) R^{2}+2 c_{b} e^{-\lambda|t|}\left(x_{0}+2 \rho\right) R \\
& \leq\left[\max \left(-a_{22},-a_{33}\right)+2 c_{b} e^{-\lambda|t|}\left(1+\frac{x_{0}+2 \rho}{r}\right)\right] R^{2} \\
& <0,
\end{aligned}
$$

as required.


Figure 9. The sets $U_{x, a}, U_{x, a}^{+}$and $U_{x, a}^{-}$.

Let $\gamma(t)$ be the solution of (36) with initial condition $\gamma\left(t_{0}\right)=p$. We shall define

$$
\phi_{t_{0}, t}(p):=\gamma\left(t_{0}+t\right) .
$$

Note that

$$
\begin{equation*}
\frac{d}{d t} \phi_{t_{0}, t}(p)=\left(A+B\left(t_{0}+t\right)\right) \phi_{t_{0}, t}(p) \tag{43}
\end{equation*}
$$

Definition 20. Let $r>0, a \geq 0$ and $x \in \mathbb{R}$ we define sets $U_{x, a}, U_{x, a}^{+}, U_{x, a}^{-} \subset \mathbb{R} \times \mathbb{R}^{2}$ by (see Figure 9)

$$
\begin{aligned}
U_{x, a} & =[x-a, x+a] \times \bar{B}(0, r), \\
U_{x, a}^{+} & =[x-a, x+a] \times \partial \bar{B}(0, r), \\
U_{x, a}^{-} & =\partial[x-a, x+a] \times \bar{B}(0, r) .
\end{aligned}
$$

### 4.1. Behavior for $t \rightarrow-\infty$

Below lemma allows us to obtain a bound on $x_{-}^{*}$ from (38). The idea is to choose some $x^{*} \in \mathbb{R}$ and an initial condition in $U_{x^{*}, 0}$ at an initial time $t^{*}<0$. If $t^{*}$ is small enough, then the influence of the matrix $B(t)$ on (36) will be small.

Lemma 21. Let us fix $x^{*}>0, \rho>0, r>0$ and $t^{*}<0$, such that conditions (41), (42) hold true for $x_{0}=x^{*}$. If

$$
\begin{equation*}
\frac{1}{\lambda} c_{b} e^{-\lambda\left|t^{*}\right|}\left(x^{*}+2 \rho+2 r\right)<\rho, \tag{44}
\end{equation*}
$$

then there exists a unique point $p^{*} \in U_{x^{*}, 0}$ such that $\phi_{t^{*}, t}\left(p^{*}\right)$ is convergent as $t \rightarrow-\infty$. Moreover,

$$
\lim _{t \rightarrow-\infty} \phi_{t^{*}, t}\left(p^{*}\right)=\left(x_{-}^{*}, 0,0\right) \in U_{x^{*}, \rho}
$$

Proof. Since (42) holds for $t^{*}$, by Lemma 19 we see that $U_{x^{*}, 2 \rho}^{+}$is an "entry set" for the set $U_{x^{*}, 2 \rho}$ (see Figures 89). What we mean by this statement is that for $p \in U_{x^{*}, 2 \rho}$ and any $t_{0}<t^{*}$ the trajectory $\phi_{t_{0}, s}(p)$ can not leave the set $U_{x^{*}, 2 \rho}$ by passing through
$U_{x^{*}, 2 \rho}^{+}$going forwards in time for $s \in\left[0, t^{*}-t_{0}\right]$. The only way it can exit is by passing through $U_{x^{*}, 2 \rho}^{-}$.

We will now show that for $p \in U_{x^{*}, \rho}$ and any $t_{0}<t^{*}$, the trajectory $\phi_{t_{0}, s}(p)$ will never leave $U_{x^{*}, 2 \rho}$ for $s \in\left[0, t^{*}-t_{0}\right]$. Since we can not exit through $U_{x^{*}, 2 \rho}^{+}$, we need to take care so that for $s \in\left[0, t^{*}-t_{0}\right]$ the $\phi_{t_{0}, s}(p)$ does not reach $U_{x^{*}, 2 \rho}^{-}$. Let us use the notation

$$
\phi_{t_{0}, s}(p)=\left(x(s), y_{1}(s), y_{2}(s)\right)
$$

As long as $\left|\pi_{x} \phi_{t_{0}, t}(p)-x^{*}\right| \leq 2 \rho$ for $t \in[0, s]$, we have the following bound (using (44) in the last inequality)

$$
\begin{align*}
& \left|\pi_{x} \phi_{t_{0}, s}(p)-x^{*}\right| \\
& \leq\left|\pi_{x} p-x^{*}\right|+\left|\int_{0}^{s} b_{11}\left(t_{0}+t\right) x(t)+b_{12}\left(t_{0}+t\right) y_{1}(t)+b_{13}\left(t_{0}+t\right) y_{2}(t) d t\right| \\
& \leq \rho+\int_{-\infty}^{t^{*}-t_{0}} c_{b} e^{\lambda\left(t_{0}+t\right)}\left(x^{*}+2 \rho+2 r\right) d t \\
& =\rho+\frac{1}{\lambda} c_{b} e^{\lambda t^{*}}\left(x^{*}+2 \rho+2 r\right) \\
& <\hat{c} 2 \rho \tag{45}
\end{align*}
$$

for a fixed $\hat{c} \in\left(\frac{1}{2}, 1\right)$, which does not depend on $s$.

$$
\forall t_{0} \leq t^{*}, \forall s \in\left[0, t^{*}-t_{0}\right], \quad\left\|\pi_{x} \phi_{t_{0}, s}(p)-x^{*}\right\| \leq 2 \rho
$$

This shows that $\phi_{t_{0}, s}$ will not leave $U_{x^{*}, 2 \rho}$. Since trajectories from $U_{x^{*}, \rho}$ do not leave $U_{x^{*}, 2 \rho}$ and can not touch $U_{x^{*}, 2 \rho}^{+}$, we have established that

$$
\begin{equation*}
\forall t_{0}<t^{*}, \forall s \in\left(0, t^{*}-t_{0}\right], \quad \phi_{t_{0}, s}\left(U_{x^{*}, \rho}\right) \subset \operatorname{int} U_{x^{*}, 2 \rho} . \tag{46}
\end{equation*}
$$

Mirror estimates to (45) together with (44) lead to

$$
\left|\pi_{x} \phi_{t_{0}, s}(p)-\pi_{x} p\right| \leq \frac{1}{\lambda} c_{b} e^{\lambda t^{*}}\left(x^{*}+2 \rho+2 r\right)<\rho,
$$

hence we proved that

$$
\begin{equation*}
\left|\pi_{x} \phi_{t_{0}, s}(p)-\pi_{x} p\right|<\rho, \quad \forall t_{0}<t^{*} \forall p \in U_{x^{*}, \rho} \forall s \in\left[0, t^{*}-t_{0}\right] . \tag{47}
\end{equation*}
$$

In particular we obtain for any $y \in \overline{B(0, r)}$, for all $t_{0} \leq t^{*}$ and $s \in\left[0, t^{*}-t_{0}\right]$

$$
\begin{equation*}
\pi_{x} \phi_{t_{0}, s}\left(x^{*}-\rho, y\right)<x^{*}, \quad \pi_{x} \phi_{t_{0}, s}\left(x^{*}+\rho, y\right)>x^{*} \tag{48}
\end{equation*}
$$

From (46) and (48), due to the topological alignment (see Figure 10) of the sets $\phi_{t_{0}, s}\left(U_{x^{*}, \rho}\right)$ and $U_{x^{*}, 0}$, it follows that for any $t_{0}<t^{*}$ there exist points $q=q\left(t_{0}\right) \in U_{x^{*}, \rho}$ and $p=p\left(t_{0}\right) \in U_{x^{*}, 0}$, such that $\phi_{t_{0}, t^{*}-t_{0}}(q)=p$ and $\phi_{t_{0}, s}(q) \in U_{x^{*}, 2 \rho}$ for $s \in\left[0, t^{*}-t_{0}\right]$. This implies that

$$
\begin{align*}
& \phi_{t^{*}, t_{0}-t^{*}}(p) \in U_{x^{*}, \rho},  \tag{49}\\
& \phi_{t^{*}, s}(p) \in U_{x^{*}, 2 \rho}
\end{align*} \quad \text { for } s \in\left[t_{0}-t^{*}, 0\right] . .
$$

Let $p_{n}=p(-n)\left(\right.$ by taking $\left.t_{0}=-n\right)$.


## Figure 10.

Since $U_{x^{*}, 0}$ is compact, we can pass to a convergent subsequence $p_{n_{k}} \rightarrow p^{*} \in U_{x^{*}, 0}$, and obtain a point for which

$$
\begin{equation*}
\phi_{t^{*}, s}\left(p^{*}\right) \in U_{x^{*}, 2 \rho} \quad \text { for } s \in(-\infty, 0] . \tag{50}
\end{equation*}
$$

We now need to show that $\phi_{t^{*}, s}\left(p^{*}\right)$ is convergent as $s \rightarrow-\infty$.
We first focus on the $x$ coordinate. We will show that for any $\varepsilon$, there exists an $S<0$, such that for any $s_{1}, s_{2}<S$ holds (the Cauchy condition for the function $s \mapsto \pi_{x} \phi_{t^{*}, s_{1}}\left(p^{*}\right)$ for $\left.s \rightarrow-\infty\right)$

$$
\begin{equation*}
\left\|\pi_{x} \phi_{t^{*}, s_{1}}\left(p^{*}\right)-\pi_{x} \phi_{t^{*}, s_{2}}\left(p^{*}\right)\right\| \leq \varepsilon . \tag{51}
\end{equation*}
$$

Assume that $s_{1}<s_{2}$ and that $u=s_{1}-s_{2}$. Then, by estimates analogous to (45),

$$
\begin{aligned}
& \left\|\pi_{x} \phi_{t^{*}, s_{1}}\left(p^{*}\right)-\pi_{x} \phi_{t^{*}, s_{2}}\left(p^{*}\right)\right\| \\
& =\left\|\pi_{x} \phi_{t^{*}, s_{2}+u}\left(p^{*}\right)-\pi_{x} \phi_{t^{*}, s_{2}}\left(p^{*}\right)\right\| \\
& \leq \frac{1}{\lambda} c_{b} e^{s_{2}}\left(x^{*}+2 \rho+2 r\right) .
\end{aligned}
$$

We see that by choosing negative $s_{2}$, with $\left|s_{2}\right|$ sufficiently large, we obtain (51). This implies the existence of $\lim _{s \rightarrow-\infty} \pi_{x} \phi_{t^{*}, s_{1}}\left(p^{*}\right)$.

We will now show that

$$
\begin{equation*}
\lim _{s \rightarrow-\infty} \pi_{y} \phi_{t^{*}, s}\left(p^{*}\right)=0 \tag{52}
\end{equation*}
$$

Suppose that this is not the case. Thus, there would exist a $\delta$ and a sequence $s_{i} \rightarrow-\infty$ such that

$$
\left|\pi_{y} \phi_{t^{*}, s_{i}}\left(p^{*}\right)\right|>\delta .
$$

Let us choose $s_{i}$ negative enough so that the $V(x, y)=\|y\|^{2}$ would be a Lyapunov function on $J \times D$ with $J=\left(-\infty, s_{i}\right]$ and

$$
D=\left[x^{*}-2 \rho, x^{*}+2 \rho\right] \times\left\{\|y\| \geq \frac{\delta}{2}\right\} .
$$

This is possible, since we can apply Lemma 19 provided that condition (42) is satisfied. The (42) will hold for $r=\frac{\delta}{2}$, provided that $t<T$, for $T$ negative enough. Since $V$ is a Lyapunov function, by going backward in time, for $s<s_{i}<T, \phi_{t^{*}, s}$ will exit $U_{x^{*}, 2 \rho}$. This contradicts (50), hence

$$
\lim _{s \rightarrow-\infty} \pi_{y} \phi_{t^{*}, s}\left(p^{*}\right)=0
$$

We have shown that $\phi_{t^{*}, s}\left(p^{*}\right)$ is convergent as $s \rightarrow-\infty$, and that the limit lies in $U_{x^{*}, 2 \rho}$. The fact that the limit is in fact in $U_{x^{*}, \rho}$ follows from (47).

We will now show that the point $p^{*}$ is unique. Suppose that we have a second point $p^{* *} \in U_{x^{*}, 0}$ for which $\lim _{s \rightarrow-\infty} \phi_{t^{*}, s}\left(p^{* *}\right)$ exists. Consider the following time dependent change of coordinates $\left(t^{*}+t, \xi=\left(t^{*}+t, x, y\right)-\phi_{t^{*}, t}\left(p^{*}\right)\right)$, which means that we just continuously change the location of the origin, by subtracting some solution. Our equation is linear, hence in new coordinates the equation is the same, $\xi^{\prime}=(A+B(t)) \xi$.

Let us take an initial condition $\xi\left(t^{*}\right)=p^{* *}-p^{*}$. If $p^{* *} \neq p^{*}$, since by construction $p^{* *}, p^{*} \in U_{x^{*}, 0}$, we see that $\xi\left(t^{*}\right) \in D=\left\{\|y\|^{2}>|x|\right\}$. Since we assumed that for $t^{*}$ condition (41) holds, it also holds for all $t<t^{*}$. By Lemma 18, this means that for $s<0$, $\xi\left(t^{*}+s\right)$ stays inside $D$. Moreover, since $V(x, y)=\|y\|^{2}-x$ is a Lyapunov function, $\xi\left(t^{*}+s\right)$ must tend to infinity as $s \rightarrow-\infty$. Since $\xi\left(t^{*}+s\right)=\phi_{t^{*}, s}\left(p^{* *}\right)-\phi_{t^{*}, s}\left(p^{*}\right)$, this means that we can not have both $p^{* *} \neq p^{*}$ and $\phi_{t^{*}, s}\left(p^{* *}\right)$ convergent. We have thus shown that $p^{*}$ is the only point in $U_{x^{*}, 0}$ for which the limit exists.

Corollary 22. If we consider $t^{*}=0, x^{*}=1, r=\rho$, and when

$$
\begin{aligned}
& 0<\lambda-4 c_{b}, \\
& 0<\min \left(a_{22}, a_{33}\right)-6 c_{b}, \\
& r>\max \left(\frac{c_{b}}{\lambda-4 c_{b}}, \frac{2 c_{b}}{\min \left(a_{22}, a_{33}\right)-6 c_{b}}\right),
\end{aligned}
$$

then assumptions of Lemma 21 are satisfied.

### 4.2. Behavior for $t \rightarrow \infty$

We now consider a point in $U_{x^{* *}, 0}$ with initial time $t^{* *}>0$.
Lemma 23. Let us fix $x^{* *}>0, \rho>0, r>0$ and $t^{* *}>0$, such that conditions (41), (42) hold true for $x_{0}=x^{* *}$.

If

$$
\frac{1}{\lambda} c_{b} e^{-\lambda t^{* *}}\left(x^{* *}+2 \rho+2 r\right)<\rho,
$$

then for any $p^{* *} \in U_{x^{* *}, 0}$ the $\phi_{t^{* *}, t}\left(p^{* *}\right)$ is convergent as $t \rightarrow+\infty$. Moreover,

$$
\begin{equation*}
\lim _{t \rightarrow+\infty} \phi_{t^{* *}, t}\left(p^{* *}\right)=\left(x_{+}^{* *}, 0,0\right) \in U_{x^{* *}, \rho} . \tag{53}
\end{equation*}
$$

Proof. The proof follows analogous in steps as the proof of Lemma 21. Since (42) holds for $t^{*}$, by Lemma 19 we see that $U_{x^{* *}, 2 \rho}^{+}$is an "entry set" for the set $U_{x^{* *}, 2 \rho}$ (see Figures

8 (9). By a similar derivation to (45) we can establish that trajectories from $U_{x^{*}, 0}$ do not leave $U_{x^{*}, \rho}$,

$$
\phi_{t^{* *}, s}\left(U_{x^{*}, 0}\right) \subset U_{x^{* *}, \rho} \quad \text { for all } t>0
$$

Using the same method (but taking the limit to $+\infty$ instead of $-\infty$ ) as for the proof of (51) we can show that we have convergence of $\pi_{x} \phi_{t^{* *}, t}\left(p^{* *}\right)$ as $t$ goes to $+\infty$ for any $p^{* *} \in U_{x^{* *}, 0}$.

We need to show that

$$
\begin{equation*}
\lim _{s \rightarrow+\infty} \pi_{y} \phi_{t^{* *}, s}\left(p^{* *}\right)=0 \tag{54}
\end{equation*}
$$

Suppose that this is not the case. Then there would exist a $\delta$ and a sequence $s_{i} \rightarrow+\infty$ such that

$$
\begin{equation*}
\left|\pi_{y} \phi_{t^{* *}, s_{i}}\left(p^{* *}\right)\right|>\delta . \tag{55}
\end{equation*}
$$

Let us choose $T$ positive enough so that $V(x, y)=\|y\|^{2}$ would be a Lyapunov function on $[T,+\infty) \times D$, for

$$
D=\left[x^{* *}-2 \rho, x^{* *}+2 \rho\right] \times\left\{\|y\| \geq \frac{\delta}{2}\right\} .
$$

This is possible, since we can apply Lemma 19 provided that condition (42) is satisfied. The 42 will hold for $r=\frac{\delta}{2}$, provided that $T<t$, for $T$ large enough. Since $V$ is a Lyapunov function, by going forwards in time $\phi_{t^{* *}, s}\left(p^{* *}\right)$ will enter $\left[x^{* *}-2 \rho, x^{* *}+2 \rho\right] \times\left\{\|y\|<\frac{\delta}{2}\right\}$ and once it enters, it will remain there. This contradicts (55), hence we have established (54).

Corollary 24. Consider $t^{* *}=0$. If

$$
\begin{aligned}
& 0<\min \left(a_{22}, a_{33}\right)-6 c_{b}, \\
& 0<\lambda-2 c_{b},
\end{aligned}
$$

and

$$
\begin{aligned}
\rho & <\frac{1}{2}\left[r\left(\frac{1}{2 c_{b}} \min \left(a_{22}, a_{33}\right)-1\right)-x^{* *}\right], \\
\rho & >\frac{c_{b}\left(x^{* *}+2 r\right)}{\lambda-2 c_{b}},
\end{aligned}
$$

then assumptions of Lemma 23 are satisfied.
Lemmas 21, 23 lead to the following algorithm for establishing bounds for $x_{-}^{*}$ and $x_{+}^{*}$ :

## Algorithm 2:

(i) Fix $r, \rho>0$ and choose $t^{*}<0$ small enough to satisfy assumptions of Lemma 21 . The Lemma 21 then ensures that

$$
x_{-}^{*} \in\left[x^{*}-\rho, x^{*}+\rho\right] .
$$

(ii) By rigorous numerical integration, evaluate the bound for $\phi_{t^{*}, t}\left(U_{x^{*}, 0}\right)$. The $t>0$ needs to be chosen large enough so that $t^{* *}=t^{*}+t$ satisfies assumptions of Lemma 23. The $t$ also needs to be large enough so that $\phi_{t^{*}, t}\left(U_{x^{*}, 0}\right) \subset U_{x^{* *}, b}$, for some $b>0$.
(iii) By Lemma $23 \lim _{t \rightarrow+\infty} \phi_{t^{*}, t}\left(p^{*}\right) \in U_{x^{* *}, b+\rho}$. This ensures that

$$
x_{+}^{*} \in\left[x^{* *}-b-\rho, x^{* *}+b+\rho\right] .
$$

We then have the following bound:

$$
\frac{x_{+}^{*}}{x_{-}^{*}} \in\left[\frac{x^{* *}-b-\rho}{x^{*}+\rho}, \frac{x^{* *}+b+\rho}{x^{*}-\rho}\right] .
$$

### 4.3. Application to the Shimizu Morioka system

In this section we prove Theorem 6. We will consider a matrix $P$ consisting of the eigenvectors of (17), of the form

$$
P=\left(\begin{array}{lll}
1 & a_{0}+1 & 0  \tag{56}\\
-1 & 1 & 0 \\
0 & 0 & 1
\end{array}\right)
$$

Thus, $P$ is transition from the eigenvectors basis to the standard one. In coordinates $\gamma$ given by $\eta=P \gamma$, the ODE (16) takes form

$$
\begin{equation*}
\gamma^{\prime}=(A+B(t)) \gamma, \tag{57}
\end{equation*}
$$

for

$$
\begin{align*}
& A=\left(\begin{array}{ccc}
0 & 0 & 0 \\
0 & -\left(a_{0}+2\right) & 0 \\
0 & 0 & -a_{0}
\end{array}\right) . \\
& B(t)=\frac{a_{0}+1}{a_{0}+2}\left(\begin{array}{ccc}
-Z_{0}(t) & Z_{0}(t) & -2 \frac{X_{0}(t)}{a_{0}+1} \\
-Z_{0}(t) & Z_{0}(t) & -2 \frac{X_{0}(t)}{a_{0}+1} \\
-\left(a_{0}+2\right) X_{0}(t) & \left(a_{0}+2\right) X_{0}(t) & 0
\end{array}\right) . \tag{58}
\end{align*}
$$

Since $a_{0} \approx 1.72$ from the bound on $X_{0}(t), Y_{0}(t), Z_{0}(t)$ from Theorem 5 by (58), we see that each coefficient $b_{i j}(t)$ of the matrix $\mathcal{B}(t)$ is bounded by

$$
\begin{aligned}
& \left|b_{i j}(t)\right| \leq\left(a_{0}+1\right) 3.5 e^{-\xi|t|}\left\|\left(X_{0}(0), Y_{0}(0), Z_{0}(0)\right)\right\| \quad \text { for } t \leq 0 \\
& \left|b_{i j}(T+t)\right| \leq\left(a_{0}+1\right) 3.5 e^{-\xi t}\left\|\left(X_{0}(T), Y_{0}(T), Z_{0}(T)\right)\right\| \quad \text { for } t \geq 0
\end{aligned}
$$

This finishes establishing the needed ingredients for Algorithm 2.
Below is the bound on the set $U_{1, \rho}$ from Lemma 21,

$$
U_{1, \rho}=\left(\begin{array}{c}
{[0.99984336210766,1.0001566378923]} \\
{[-0.00015663789234007,0.00015663789234007]} \\
{[-0.00015663789234007,0.00015663789234007]}
\end{array}\right)
$$

By Lemma 21, there is a point in $U_{1,0}$, that converges to $\left(x_{-}^{*}, 0,0\right) \in U_{1, \rho}$. We therefore take the set $U_{1,0}$ as the initial point from which we integrate (57) forward in time. A bound on all trajectories that start in $U_{1,0}$ is depicted in Figure 11. (We make also the


Figure 11. The computer assisted bound on the heteroclinic trajectory of (57).


Figure 12. The computer assisted bound on the heteroclinic trajectory of 16 .
same plot in coordinates $\eta$ in Figure 12.) A trajectory from such enclosure makes a loop, to finish at time $T$ closer to the origin, in a cubical enclosure $\Gamma=\Gamma_{1} \times \Gamma_{2} \times \Gamma_{3} \subset \mathbb{R}^{3}$. We use $\Gamma$ to compute the bound on the set $U_{\Gamma_{1}, \rho}$ from Lemma 23 , obtaining

$$
U_{\Gamma_{1}, \rho} \subset\left(\begin{array}{c}
{[0.62606812264791,0.62663392848044]} \\
{[-5.3960913395776 e-05,5.3960913395776 e-05]} \\
{[-5.3960913395776 e-05,5.3960913395776 e-05]}
\end{array}\right)
$$

By Lemma 23, there exists a point $\left(x_{+}^{*}, 0,0\right) \in U_{\Gamma_{1}, \rho}$ and a trajectory $\gamma(t)$ of (57), for with

$$
\begin{aligned}
\lim _{t \rightarrow-\infty} \gamma(t) & =\left(x_{-}^{*}, 0,0\right), \\
\lim _{t \rightarrow+\infty} \gamma(t) & =\left(x_{+}^{*}, 0,0\right) .
\end{aligned}
$$

From the fact that $\lim _{t \rightarrow \pm \infty} \eta(t)=\lim _{t \rightarrow \pm \infty} P \gamma(t)$, we obtain the claim of Theorem 6.
The computer assisted proof has been done entirely by using the CAPD $\ddagger$ package and took under a second on a single core 3 Ghz Intel i7 processor.
$\ddagger$ computer assisted proofs in dynamics: http://capd.ii.uj.edu.pl/

## 5. Appendix

### 5.1. Proof of Theorem 11

Before we give the proof, we introduce some auxiliary tools.
Let

$$
J_{u}(z, L)=\left\{(x, y) \in \mathbb{R}^{u} \times \mathbb{R}^{s}:\left\|\pi_{y} z-y\right\| \leq L\left\|\pi_{x} z-y\right\|\right\}
$$

The set $J_{u}(z, L)$ defines a cone of slope $L$, centered at $z$. Below theorem establish cone alignment for a map $\phi$, which satisfies certain bounds on its derivative.

Theorem 25. [29, Theorem 27] Let $U \subset \mathbb{R}^{u} \times \mathbb{R}^{s}$ be a convex neighborhood of zero and assume that $\phi: U \rightarrow \mathbb{R}^{u} \times \mathbb{R}^{s}$ is a $C^{1}$. If for $M>0$

$$
\begin{aligned}
& v \geq \sup _{z \in U}\left\{\left\|\frac{\partial \phi_{y}}{\partial y}(z)\right\|+\frac{1}{L}\left\|\frac{\partial \phi_{y}}{\partial x}(z)\right\|\right\}, \\
& \zeta \leq m\left[\frac{\partial \phi_{x}}{\partial x}(U)\right]-L \sup _{z \in U}\left\|\frac{\partial \phi_{x}}{\partial y}(z)\right\|,
\end{aligned}
$$

and

$$
\frac{\zeta}{v}>1
$$

then for $z \in U$

$$
\begin{equation*}
\phi\left(J_{u}(z, L) \cap U\right) \subset \operatorname{int} J_{u}(\phi(z), L) \cup\{\phi(z)\} \tag{59}
\end{equation*}
$$

Note that if the vector field (26) satisfies condition $\vec{\mu}<0<\vec{\xi}$, then by Theorem 12, for sufficiently small $h>0$, and for $\mu(h)$ and $\xi(h)$ defined in 31 32), the

$$
v=\mu(h), \quad \zeta=\xi(h),
$$

will satisfy the assumptions of Theorem 25 for $M=L$. This will imply (59) for $\phi=\Phi_{h}$, for the flow $\Phi$ induced by (26).

We now give the sketch of the proof of Theorem 11 .
Proof of Theorem 11. The proof of the theorem follows from a mirror argument to the proof of Theorem 30 from [27]. The one important difference is that the result from [27] is written in the context where in addition to the hyperbolic directions $x, y$, we also have a center coordinate. Here such coordinate does not exist, which allows us to obtain better bounds on the slope of the established manifold. Also, due to the lack of the center coordinate, we have less inequalities in the assumptions of our theorem compared to [27]. Instead of repeating the proof of Theorem 30 from [27] we refer the reader to the source, and will focus here on the Lipschitz bounds of the manifold, which is the improvement of the current result over [27].

The proof of Theorem 30 from [27] follows from a graph transform method [29]. We start with a flat function

$$
\begin{equation*}
\bar{B}_{u}(R) \ni x \rightarrow 0 \in \bar{B}_{s}(R), \tag{60}
\end{equation*}
$$

and propagate it using the graph transform method. The manifold $W^{u}$ is obtained by passing to the limit. Theorem 30 from [27] ensures that $W^{u}$ is a graph of the function $w^{u}$, meaning that we have (30). (We refer the reader to [27] for the proof of this procedure that is just outline here.) We will focus here on the Lipschitz bounds obtained at the limit. From [27] we would obtain directly the Lipschitz bound $1 / L$, but here we will show that the bound (due to lack of the center coordinate) is in fact $L$.

The important issue is that (29) ensures that the assumptions of Theorem 25 are satisfied (with $M=L$,) for the map $z \rightarrow \Phi_{h}(z)$, by taking $v=\mu(h, L), \zeta=\xi(h, L)$ and sufficiently small $h>0$. This means that the cones $J_{u}$ are preserved along the flow in the sense of (59); i.e. that

$$
\Phi_{h}\left(J_{u}(z, M) \cap U\right) \subset \operatorname{int} J_{u}\left(\Phi_{h}(z), M\right) \cup\left\{\Phi_{h}(z)\right\}
$$

For any $x \in \bar{B}_{u}(R)$ the graph of 60 is inside of the cone $J_{u}((x, 0), L)$. The graph of (60) after propagating by $\Phi_{h}$ using the graph transform, will also be contained in cones $J_{u}$. This implies that, after passing to the limit with the graph transform, for any $x_{1}, x_{2} \in \bar{B}_{u}(R)$

$$
\left(x_{1}, w^{u}\left(x_{1}\right)\right) \in J_{u}\left(\left(x_{2}, w^{u}\left(x_{2}\right)\right), L\right),
$$

hence

$$
\left\|w^{u}\left(x_{2}\right)-w^{u}\left(x_{1}\right)\right\| \leq L\left\|x_{1}-x_{2}\right\| .
$$

This means that $w^{u}$ is Lipschitz with constant $L$, as required.
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