# Existence of periodic solutions for the Lotka-Volterra type systems 

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#### Abstract

In this paper we prove the existence of nonstationary periodic solutions of delay Lotka-Volterra equations. In the proofs we use the $S^{1}$-degree due to Dylawerski et al. [G. Dylawerski, K. Geba, J. Jodel, W. Marzantowicz, An $S^{1}$-equivariant degree and the Fuller index, Ann. Polon. Math. 63 (1991) 243-280]. © 2006 Elsevier Inc. All rights reserved.


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## 1. Introduction

The aim of this paper is to prove the existence of nonstationary periodic solutions of autonomous delay differential equations of Lotka-Volterra type

$$
\left\{\begin{array}{l}
\dot{u}_{1}(t)=u_{1}(t)\left(r_{1}-a_{11} u_{1}(t-\tau)-a_{12} u_{2}(t-\tau)-\cdots-a_{1 n} u_{n}(t-\tau)\right),  \tag{1.1}\\
\dot{u}_{2}(t)=u_{2}(t)\left(r_{2}-a_{21} u_{1}(t-\tau)-a_{22} u_{2}(t-\tau)-\cdots-a_{2 n} u_{n}(t-\tau)\right), \\
\quad \vdots \\
\dot{u}_{n}(t)=u_{n}(t)\left(r_{n}-a_{n 1} u_{1}(t-\tau)-a_{n 2} u_{2}(t-\tau)-\cdots-a_{n n} u_{n}(t-\tau)\right),
\end{array}\right.
$$

[^0]where $n \geqslant 1, \tau>0, r_{1}, \ldots, r_{n} \in \mathbb{R}, a_{i j} \in \mathbb{R}$, for $i, j=1, \ldots, n$.
It is known that a broad class of problems in mathematical biology, economics and mechanics are described in the form above with initial conditions
\[

$$
\begin{cases}u_{i}(s)=\varphi_{i}(s), & s \in[-\tau, 0], \varphi_{i}(0)>0  \tag{1.2}\\ \varphi_{i} \in C([-\tau, 0], \mathbb{R}), & i=1,2, \ldots, n\end{cases}
$$
\]

In case $n=1$, problem (1.1) is known as delay logistic equation. The existence and multiplicity of solutions of delay logistic equation has been investigated by many authors (cf. Goparlsamy [3] and Hale [4] and references therein). To compare with the method employed here with that for delay logistic equation, we illustrate the proof for the existence of periodic solutions of the logistic equation

$$
\begin{equation*}
\dot{u}(t)=\alpha u(t)(1-u(t-\tau)), \tag{1.3}
\end{equation*}
$$

where $\alpha>0$. For each initial function $\varphi \in C([-\tau, 0], \mathbb{R})$ with $\varphi(0)>0$, one can find a solution $u(\varphi)$ of (1.3) with initial value $u(\varphi)(s)=\varphi(s), s \in[-\tau, 0]$. We put

$$
z(\varphi, \alpha)=\min \{t>0: u(\varphi)(t)=0, \dot{u}(\varphi)(t)>0\}
$$

and $[A(\alpha) \varphi](t)=u(\varphi)(t-z(\varphi, \alpha)-\tau)$ for $t>0$. Then one can see that each fixed point $u$ of $A(\alpha)$ is a periodic solution of (1.3). The existence of the nonstationary fixed points of $A(\alpha)$ is proved by combination of the Hopf bifurcation theorem and fixed point theorems (cf. Hale [4, Section 11.4]). For $\tau=1$, it is known that $\alpha=\pi / 2$ is the bifurcation point of solutions to (1.3) and for each $\alpha>\pi / 2$, problem (1.3) has a nonstationary periodic solution.

On the other hand, it is natural to ask if there are multiple solutions of (1.3) for sufficiently large $\tau$. The multiple existence of periodic solution of (1.3) for sufficiently large $\tau$ also follows from the Hopf bifurcation. In general, the methods employed for delay logistic equation are not valid for (1.1) with $n>1$.

In this paper we work with the space of periodic functions instead of considering the initial value problem and make use of the $S^{1}$-degree, see [2], to prove the multiplicity of solutions of problem (1.1). Applications of the degree for equivariant maps to the study of periodic solutions of a van der Pol system one can find in $[1,5]$.

To avoid unnecessary complexity, we restrict ourselves to the case $n=2$, that is, we consider the coupled equations of the form

$$
\left\{\begin{array}{l}
\dot{u}(t)=u(t)\left(r_{1}-a_{11} u(t-\tau)-a_{12} v(t-\tau)\right),  \tag{1.4}\\
\dot{v}(t)=v(t)\left(r_{2}-a_{21} u(t-\tau)-a_{22} v(t-\tau)\right)
\end{array}\right.
$$

Our argument does not depend on any specific property of $n=2$. That is why our result is valid for $n \geqslant 2$ with modifications of assumptions for the case that $n \geqslant 2$.

We impose the following conditions on matrix $A=\left[\begin{array}{lll}a_{11} & a_{12} \\ a_{21} & a_{22}\end{array}\right]$ :
(A0) $a_{11}, a_{12}, a_{21}, a_{22}, b_{1}, b_{2}>0$, where $\left[\begin{array}{l}b_{1} \\ b_{2}\end{array}\right]=A^{-1}\left[\begin{array}{l}r_{1} \\ r_{2}\end{array}\right]$,
(A1) $\langle A x, x\rangle>0$ for all $x \in \mathbb{R}^{2} \backslash\{0\}$,
(A2) a matrix $\left[\begin{array}{lll}b_{1} a_{11} & b_{1} a_{12} \\ b_{2} a_{21} & b_{2} a_{22}\end{array}\right]$ possesses two real eigenvalues $\mu_{1}, \mu_{2}>0$.

Remark 1.1. Notice that for an arbitrary $n \in \mathbb{N}$ assumptions (A0)-(A2) can be reformulated in the following way:
(A0) $a_{i j}, b_{i}>0$ for $1 \leqslant i, j \leqslant n$,
(A1) $\langle A x, x\rangle>0$ for all $x \in \mathbb{R}^{n} \backslash\{0\}$,
(A2) a matrix $\operatorname{diag}\left(b_{1}, \ldots, b_{n}\right) \cdot A$ possesses only real, positive eigenvalues $\mu_{1}, \ldots, \mu_{n}$, each with multiplicity one.

We can now formulate the main result of this article.
Theorem 1.1. Fix $\tau>0$ such that $\min \left\{\frac{2 \pi}{\mu_{1}}, \frac{2 \pi}{\mu_{2}}\right\}<\tau<\infty$. Assume that there are $n_{1}, n_{2} \in \mathbb{N} \cup\{0\}$ such that $n_{1} \neq n_{2}$ and for $i=1,2$,

$$
\begin{equation*}
\frac{\pi}{2}+2 n_{i} \pi<\mu_{i} \tau<\frac{\pi}{2}+2\left(n_{i}+1\right) \pi \tag{1.5}
\end{equation*}
$$

Under the above assumptions there is at least one nonstationary $\tau$-periodic solution of (1.4).
After this introduction our paper is organized as follows.
For the convenience of the reader in Section 2 we have repeated the relevant material from [2] without proofs, thus making our exposition self-contained.

In Section 3 we have performed a functional setting for our problem. This section is of technical nature. Namely, applying transformation of functions and fixing the period we have obtained a parameterized problem (3.2) which is equivalent to the original problem (1.4). Next we have defined a Banach space $E$ which is an infinite-dimensional representation of the group $S^{1}$, an open $S^{1}$-invariant subset $\Theta_{0} \subset E$ and an $S^{1}$-equivariant compact operator $F:\left(E \times \mathbb{R}^{+}\right) \times[0,1] \rightarrow E$, see formula (3.5), such that solutions of equation $F\left(\left(\left(x_{1}, x_{2}\right), \lambda\right), 1\right)=\left(x_{1}, x_{2}\right)$ in $\Theta_{0} \times \mathbb{R}^{+}$are exactly periodic solutions of problem (3.2).

In Section 4 we have defined an open, bounded $S^{1}$-invariant subset $\Omega_{\lambda_{1}, \lambda_{2}} \subset \Theta_{0} \times \mathbb{R}^{+} \subset$ $E \times \mathbb{R}^{+}$such that the homotopy $Q-F(\cdot, \theta)$, defined by (3.5) does not vanish on $\partial \Omega_{\lambda_{1}, \lambda_{2}}$. This allow us to simplify computations of the $S^{1}$-degree of $Q-F(\cdot, 1)$ on $\Omega_{\lambda_{1}, \lambda_{2}}$, see Lemma 4.4.

In Section 5 we have proved Theorem 1.1.

## 2. $S^{1}$-degree

In this section we have compiled some basic facts on the $S^{1}$-degree defined in [2]. Let $S^{1}=$ $\{z \in \mathbb{C}:|z|=1\}=\left\{e^{i \cdot \theta}: \theta \in[0,2 \pi)\right\}$ be the group with an action given by the multiplication of complex numbers. For any fixed $m \in \mathbb{N}$ we denote by $\mathbb{Z}_{m}$ a cyclic group of order $m$ and define homomorphism $\rho_{m}: S^{1} \rightarrow G L(2, \mathbb{R})$ as follows

$$
\rho_{m}\left(e^{i \cdot \theta}\right)=\left[\begin{array}{cc}
\cos (m \theta) & -\sin (m \theta) \\
\sin (m \theta) & \cos (m \theta)
\end{array}\right]
$$

Let $E$ be a Banach space which is an $S^{1}$-representation. We denote by $Q: E \times \mathbb{R} \rightarrow E$ the projection. For each closed subgroup $H$ of $S^{1}$ and each $S^{1}$-invariant subset $\Omega \subset E$, we denote by $\Omega^{H}$ the subset of fixed points of the action of $H$ on $\Omega$. For given $a \in E, S_{a}^{1}=\left\{s \in S^{1}: s \cdot a=a\right\}$ is called the isotropy group of $a$ and the set $S^{1} \cdot a=\left\{s \cdot a: s \in S^{1}\right\}$ is called the orbit of $a$. Denote
by $\Gamma_{0}$ the free abelian group generated by $\mathbb{N}$ and let $\Gamma=\mathbb{Z}_{2} \oplus \Gamma_{0}$. Then $\gamma \in \Gamma$ means $\gamma=\left\{\gamma_{r}\right\}$, where $\gamma_{0} \in \mathbb{Z}_{2}$ and $\gamma_{r} \in \mathbb{Z}$ for $r \in \mathbb{N}$.

Fix an open, bounded $S^{1}$-invariant subset $\Omega \subset E \times \mathbb{R}$ and continuous $S^{1}$-equivariant compact mapping $\Phi: \operatorname{cl}(\Omega) \rightarrow E$ such that $(Q+\Phi)(\partial \Omega) \subset E \backslash\{0\}$. In this situation the $S^{1}$-degree $\operatorname{Deg}(Q+\Phi, \Omega)=\left\{\gamma_{r}\right\} \in \Gamma$, where $\gamma_{0}=\operatorname{deg}_{S^{1}}(Q+\Phi, \Omega)$ and $\gamma_{r}=\operatorname{deg}_{\mathbb{Z}_{r}}(Q+\Phi, \Omega), r \in \mathbb{N}$, has been defined in [2].

Theorem 2.1. (Cf. [2].) Let $E$ be a Banach space which is a representation of the group $S^{1}$, $\Omega_{0}, \Omega_{1}, \Omega_{2} \subset \Omega$ be open bounded, $S^{1}$-invariant subsets of $E \times \mathbb{R}$. Assume that $\Phi: c l(\Omega) \rightarrow E$ is a compact $S^{1}$-equivariant mapping such that $(Q+\Phi)(\partial \Omega) \subset E \backslash\{0\}$. Then there exists a $\Gamma$-valued function $\operatorname{Deg}(Q+\Phi, \Omega)$ called the $S^{1}$-degree, satisfying the following properties:
(a) if $\operatorname{deg}_{H}(Q+\Phi, \Omega) \neq 0$, then $(Q+\Phi)^{-1}(0) \cap \Omega^{H} \neq \emptyset$,
(b) if $(Q+\Phi)^{-1}(0) \cap \Omega \subset \Omega_{0}$, then $\operatorname{Deg}(Q+\Phi, \Omega)=\operatorname{Deg}\left(Q+\Phi, \Omega_{0}\right)$,
(c) if $\Omega_{1} \cap \Omega_{2}=\emptyset$ and $(Q+\Phi)^{-1}(0) \cap \Omega \subset \Omega_{1} \cup \Omega_{2}$, then

$$
\operatorname{Deg}(Q+\Phi, \Omega)=\operatorname{Deg}\left(Q+\Phi, \Omega_{1}\right)+\operatorname{Deg}\left(Q+\Phi, \Omega_{2}\right)
$$

(d) if $h: c l(\Omega) \times[0,1] \rightarrow E$ is an $S^{1}$-equivariant compact homotopy such that $(Q+h)(\partial \Omega \times[0,1]) \subset E \backslash\{0\}$, then $\operatorname{Deg}\left(Q+h_{0}, \Omega\right)=\operatorname{Deg}\left(Q+h_{1}, \Omega\right)$.

Let $E_{0}, E$ be real Banach spaces. We denote by $K\left(E_{0}, E\right)$ the set of compact operators $B: E_{0} \rightarrow E$. Fix $B_{0} \in K(E, E)$. A real number $\mu$ is called a characteristic value of $B_{0}$ if $\operatorname{dim} \operatorname{ker}\left(I-\mu B_{0}\right)>0$. Suppose now that $A_{0}=I-B_{0}$ is an invertible operator. Denote by $\left\{\mu_{1}<\mu_{2}<\cdots<\mu_{p}\right\}$ the set of all characteristic values of $B_{0}$ contained in [0, 1]. We set $\operatorname{sgn} A_{0}=(-1)^{d}$, where $d=\sum_{i=1}^{p} \operatorname{dim} \operatorname{ker}\left(I-\mu_{i} B_{0}\right)$.

Fix $B_{1} \in K(E \times \mathbb{R}, E)$ and define $A=Q+B_{1}: E \times \mathbb{R} \rightarrow E$. Assume that $A$ is surjective. Since $A$ is a Fredholm operator of index 1 and $A$ is surjective, $\operatorname{dim} \operatorname{ker} A=1$. Fix $v \in \operatorname{ker} A \backslash\{0\}$ and define a linear functional $\xi: E \times \mathbb{R} \rightarrow E$ such that $\xi(v)=1$. Finally, define an operator $A^{\sim}: E \times \mathbb{R} \rightarrow E \times \mathbb{R}$ by $A^{\sim}(w)=(A w, \xi(w))$ and $\operatorname{sgn}(A, v)=\operatorname{sgn}\left(A^{\sim}\right)$.

Assume additionally that $f=Q+\Phi \in C^{1}(c l(\Omega), E)$. Suppose that $a \in \Omega$ is such that there is $k \in \mathbb{N}$ such that $S_{a}^{1}=\mathbb{Z}_{k}, f^{-1}(0) \cap \Omega=S^{1} \cdot a \approx S^{1} / \mathbb{Z}_{k}$ and $D f(a): E \times \mathbb{R} \rightarrow E$ is surjective. Let $f^{\mathbb{Z}_{k}}: E^{\mathbb{Z}_{k}} \times \mathbb{R} \rightarrow E^{\mathbb{Z}_{k}}$ denotes the restriction of $f$. Then $D f^{\mathbb{Z}_{k}}(a): E^{\mathbb{Z}_{k}} \times \mathbb{R} \rightarrow E^{\mathbb{Z}_{k}}$ is also surjective. We denote by $v$ the tangent vector to the orbit $S^{1} \cdot a$ at $a$. Notice that $v \in \operatorname{ker} D f^{\mathbb{Z}_{k}}(a)$.

This theorem ensures the nontriviality of the $S^{1}$-index of the nondegenerate $S^{1}$-orbit $S^{1} \cdot a$.
Theorem 2.2. (Cf. [2].) Under the above assumptions $\operatorname{deg}_{\mathbb{Z}_{k}}(f, \Omega)=\operatorname{sgn}\left(D f^{\mathbb{Z}_{k}}(a), v\right)$. Moreover, $\operatorname{deg}_{\mathbb{Z}_{k^{\prime}}}(f, \Omega)=0$ for every $k^{\prime}>k$.

## 3. Functional setting

Throughout the rest of this article we assume that assumptions (A0)-(A2) are fulfilled. Moreover, we fix $\tau>0$ satisfying assumptions of Theorem 1.1.

In this section we convert problem (1.4) to an equivalent problem (3.2). Next we define spaces on which we will work and define a homotopy $F$ of $S^{1}$-invariant compact mappings. The study of periodic solutions of problem (3.2) is equivalent to the study of fixed $S^{1}$-orbits of the operator $F(\cdot, 1)$.

By a transformation of functions (cf. [4]), one can see that problem (1.4) is equivalent to the problem

$$
\left\{\begin{array}{l}
\dot{u}(t)=-\left(a_{11} u(t-\tau)+a_{12} v(t-\tau)\right)\left(b_{1}+u(t)\right),  \tag{3.1}\\
\dot{v}(t)=-\left(a_{21} u(t-\tau)+a_{22} v(t-\tau)\right)\left(b_{2}+v(t)\right)
\end{array}\right.
$$

Let $(u, v) \in C(\mathbb{R}, \mathbb{R}) \times C(\mathbb{R}, \mathbb{R})$ be a pair of periodic functions with period $T$. Then by putting $\lambda=\frac{T}{2 \pi}, x_{1}(t)=u(\lambda t)$ and $x_{2}(t)=v(\lambda t)$ for $t \in \mathbb{R}$, we have that $x=\left(x_{1}, x_{2}\right)$ is a $2 \pi$-periodic solution of problem

$$
\left\{\begin{array}{l}
\dot{x}_{1}(t)=-\lambda\left(a_{11} x_{1}(t-\tau / \lambda)+a_{12} x_{2}(t-\tau / \lambda)\right)\left(b_{1}+x_{1}(t)\right),  \tag{3.2}\\
\dot{x}_{2}(t)=-\lambda\left(a_{21} x_{1}(t-\tau / \lambda)+a_{22} x_{2}(t-\tau / \lambda)\right)\left(b_{2}+x_{2}(t)\right) .
\end{array}\right.
$$

We will study the existence of $2 \pi$-periodic solutions of (3.2) for some $\lambda>0$ instead of looking for periodic solutions of (3.1). We note that each function $u:[0,2 \pi] \rightarrow \mathbb{R}$ with $u(0)=u(2 \pi)$ is extended to a $2 \pi$-periodic function on $\mathbb{R}$. Therefore we identify a $2 \pi$-periodic function $u$ on $\mathbb{R}$ with a function on $[0,2 \pi]$ with $u(0)=u(2 \pi)$.

Define a Banach space

$$
\widehat{E}=\left\{x \in C([0,2 \pi], \mathbb{R}): \int_{0}^{2 \pi} \dot{x}(t)^{2} d t<\infty, x(0)=x(2 \pi) \text { and } \int_{0}^{2 \pi} x(t) d t=0\right\}
$$

with a norm $\|\cdot\|$ given by $\|x\|^{2}=\int_{0}^{2 \pi} \dot{x}(t)^{2}+x(t)^{2} d t$, for $x \in \widehat{E}$. Moreover, we put $|u|_{\infty}=$ $\sup \{|u(t)|: t \in[0,2 \pi]\}$ for $u \in \widehat{E}$. Put $E=\widehat{E} \times \widehat{E}$ and define an action $\rho: S^{1} \times E \rightarrow E$ of the group $S^{1}$ as follows

$$
\begin{equation*}
\rho\left(e^{i \phi},\left(x_{1}(t), x_{2}(t)\right)\right)=\left(x_{1}(t+\phi), x_{2}(t+\phi)\right) \quad \bmod 2 \pi \tag{3.3}
\end{equation*}
$$

Define an open $S^{1}$-invariant subset $\Theta_{0} \subset E$ as follows

$$
\Theta_{0}=\left\{\left(x_{1}, x_{2}\right) \in E:-b_{i}<x_{i}(t) \text { for } t \in[0,2 \pi], i=1,2\right\} .
$$

Next we define a homotopy of $S^{1}$-equivariant mappings $F:\left(E \times \mathbb{R}^{+}\right) \times[0,1] \rightarrow E$ associated to problem (3.2) such that if $\left(\left(x_{1}, x_{2}\right), \lambda\right) \in \Theta_{0} \times \mathbb{R}^{+}$satisfies $F\left(\left(x_{1}, x_{2}\right), \lambda, 1\right)=\left(x_{1}, x_{2}\right)$, then $\left(\left(x_{1}, x_{2}\right), \lambda\right)$ is a solution of problem (3.2).

Let $\beta: E \rightarrow[0,1]$ be a continuous mapping and $\mathcal{N}: E \times[0,1] \rightarrow E$ be a mapping defined by

$$
\mathcal{N}\left(\left(x_{1}, x_{2}\right), \theta\right)(t)=\binom{-\left(a_{11} x_{1}(t-\tau / \lambda)+a_{12} x_{2}(t-\tau / \lambda)\right)\left(b_{1}+\theta x_{1}(t)\right)}{-\left(a_{21} x_{1}(t-\tau / \lambda)+a_{22} x_{2}(t-\tau / \lambda)\right)\left(b_{2}+\theta x_{2}(t)\right)}
$$

for $\left(\left(x_{1}, x_{2}\right), \theta\right) \in E \times[0,1]$. We put

$$
c_{1}\left(\left(x_{1}, x_{2}\right), \theta\right)=\binom{c_{11}\left(\left(x_{1}, x_{2}\right), \theta\right)}{c_{12}\left(\left(x_{1}, x_{2}\right), \theta\right)}=-\frac{\lambda}{2 \pi} \int_{0}^{2 \pi} \beta\left(x_{1}, x_{2}\right) \mathcal{N}\left(\left(x_{1}, x_{2}\right), \theta\right)(s) d s \quad \text { and }
$$

$$
\begin{aligned}
c_{2}\left(\left(x_{1}, x_{2}\right), \theta\right) & =\binom{c_{21}\left(\left(x_{1}, x_{2}\right), \theta\right)}{c_{22}\left(\left(x_{1}, x_{2}\right), \theta\right)} \\
& =-\frac{\lambda}{2 \pi} \int_{0}^{2 \pi} \int_{0}^{t} \beta\left(x_{1}, x_{2}\right) \mathcal{N}\left(\left(x_{1}, x_{2}\right), \theta\right)(s) d s d t-\pi c_{1}\left(\left(x_{1}, x_{2}\right), \theta\right)
\end{aligned}
$$

We define a mapping $F:\left(E \times \mathbb{R}^{+}\right) \times[0,1] \rightarrow C(\mathbb{R}, \mathbb{R}) \times C(\mathbb{R}, \mathbb{R})$ by

$$
\begin{align*}
& F\left(\left(\left(x_{1}, x_{2}\right), \lambda\right), \theta\right)(t) \\
& \quad=-\lambda \int_{0}^{t} \beta\left(x_{1}, x_{2}\right) \mathcal{N}\left(\left(x_{1}, x_{2}\right), \theta\right)(s) d s-t c_{1}\left(\left(x_{1}, x_{2}\right), \theta\right)-c_{2}\left(\left(x_{1}, x_{2}\right), \theta\right) . \tag{3.4}
\end{align*}
$$

From the definitions of $c_{1}$ and $c_{2}$, we can see that for all $\left(\left(x_{1}, x_{2}, \lambda\right), \theta\right) \in\left(E \times \mathbb{R}^{+}\right) \times[0,1]$

$$
F\left(\left(\left(x_{1}, x_{2}\right), \lambda\right), \theta\right)(0)=F\left(\left(\left(x_{1}, x_{2}\right), \lambda\right), \theta\right)(2 \pi) \quad \text { and } \quad \int_{0}^{2 \pi} F\left(\left(\left(x_{1}, x_{2}\right), \lambda\right), \theta\right)(t) d t=0
$$

holds.
Summing up, $F\left(\left(\left(x_{1}, x_{2}\right), \lambda\right), \theta\right) \in E$ for all $\left(\left(\left(x_{1}, x_{2}\right), \lambda\right), \theta\right) \in\left(E \times \mathbb{R}^{+}\right) \times[0,1]$. It is also easy to see that $F:\left(E \times \mathbb{R}^{+}\right) \times[0,1] \rightarrow E$ is an $S^{1}$-equivariant compact mapping.

From the definition of $F$, we find that $\left(\left(\left(x_{1}, x_{2}\right), \lambda\right), \theta\right) \in\left(\Theta_{0} \times \mathbb{R}^{+}\right) \times[0,1]$ satisfies

$$
\begin{equation*}
F\left(\left(\left(x_{1}, x_{2}\right), \lambda\right), \theta\right)=Q\left(\left(x_{1}, x_{2}\right), \lambda\right) \tag{3.5}
\end{equation*}
$$

if and only if

$$
\left\{\begin{array}{l}
\dot{x}_{1}(t)=-\lambda \beta\left(x_{1}, x_{2}\right)\left(a_{11} x_{1}(t-\tau / \lambda)+a_{12} x_{2}(t-\tau / \lambda)\right)\left(b_{1}+\theta x_{1}(t)\right)-c_{11}\left(\left(x_{1}, x_{2}\right), \theta\right),  \tag{3.6}\\
\dot{x}_{2}(t)=-\lambda \beta\left(x_{1}, x_{2}\right)\left(a_{21} x_{1}(t-\tau / \lambda)+a_{22} x_{2}(t-\tau / \lambda)\right)\left(b_{2}+\theta x_{2}(t)\right)-c_{21}\left(\left(x_{1}, x_{2}\right), \theta\right) .
\end{array}\right.
$$

We claim that system (3.6) is equivalent to

$$
\left\{\begin{array}{l}
\dot{x}_{1}(t)=-\lambda \beta\left(x_{1}, x_{2}\right)\left(a_{11} x_{1}(t-\tau / \lambda)+a_{12} x_{2}(t-\tau / \lambda)\right)\left(b_{1}+\theta x_{1}(t)\right),  \tag{3.7}\\
\dot{x}_{2}(t)=-\lambda \beta\left(x_{1}, x_{2}\right)\left(a_{21} x_{1}(t-\tau / \lambda)+a_{22} x_{2}(t-\tau / \lambda)\right)\left(b_{2}+\theta x_{2}(t)\right) .
\end{array}\right.
$$

What is left is to show that $c_{i 1}\left(\left(x_{1}, x_{2}\right), \theta\right)=0$ for $i=1,2$.
Fix $i \in\{1,2\}$ and notice that

$$
\begin{aligned}
\frac{d}{d t} \ln \left(b_{i}+x_{i}(t)\right) & =\frac{\dot{x}_{i}(t)}{b_{i}+x_{i}(t)} \\
& =-\lambda \beta\left(x_{1}, x_{2}\right)\left(a_{i 1} x_{1}(t-\tau / \lambda)+a_{i 2} x_{2}(t-\tau / \lambda)\right)-\frac{c_{i 1}\left(\left(x_{1}, x_{2}\right), \theta\right)}{b_{i}+x_{i}(t)}
\end{aligned}
$$

Thus

$$
\begin{aligned}
& \ln \left(b_{i}+x_{i}(t)\right)-\ln \left(b_{i}+x_{i}(s)\right) \\
& \quad=-\lambda \int_{s}^{t} \beta\left(x_{1}, x_{2}\right)\left(a_{i 1} x_{1}(w-\tau / \lambda)+a_{i 2} x_{2}(w-\tau / \lambda)\right) d w-\int_{s}^{t} \frac{c_{i 1}\left(\left(x_{1}, x_{2}\right), \theta\right)}{b_{i}+x_{i}(w)} d w
\end{aligned}
$$

for $s, t \in \mathbb{R}$ with $s<t$. Since

$$
x_{i}(2 \pi)=x_{i}(0), \quad \int_{0}^{2 \pi} x_{1}(t) d t=\int_{0}^{2 \pi} x_{2}(t) d t=0, \quad \int_{0}^{2 \pi} \frac{c_{i 1}\left(\left(x_{1}, x_{2}\right), \theta\right)}{b_{i}+x_{i}(t)} d t=0 .
$$

Finally, condition $x_{i}(t)>-b_{i}$ for all $t \in[0,2 \pi]$, implies $c_{i 1}\left(\left(x_{1}, x_{2}\right), \theta\right)=0$, which completes the proof.

We finish this section with the following lemma which yields a priori estimates for periodic solutions of problem (3.8).

## Lemma 3.1.

(1) For $\lambda_{1}, \lambda_{2} \in \mathbb{R}^{+}$with $\lambda_{1}<\lambda_{2}$, there exist positive numbers $m_{0}$, $\left\{d_{i}\right\}_{1 \leqslant i \leqslant 4}$ such that for each $\lambda \in\left[\lambda_{1}, \lambda_{2}\right], \alpha \in[0,1]$ and $\tau \geqslant 1$, each solution $\left(x_{1}, x_{2}\right) \in \Theta_{0}$ of the following problem

$$
\left\{\begin{array}{l}
\dot{x}_{1}(t)=-\alpha \lambda\left(a_{11} x_{1}(t-\tau / \lambda)+a_{12} x_{2}(t-\tau / \lambda)\right)\left(b_{1}+x_{1}(t)\right),  \tag{3.8}\\
\dot{x}_{2}(t)=-\alpha \lambda\left(a_{21} x_{1}(t-\tau / \lambda)+a_{22} x_{2}(t-\tau / \lambda)\right)\left(b_{2}+x_{2}(t)\right),
\end{array}\right.
$$

satisfies $\left|\dot{x}_{i}\right|_{\infty}<m_{0},\left|\ddot{x}_{i}\right|_{\infty}<m_{0}$ for $i=1,2$, and

$$
-b_{1}<-d_{1}<x_{1}(t)<d_{3}, \quad-b_{2}<-d_{2}<x_{2}(t)<d_{4} \quad \text { on }[0,2 \pi] ;
$$

(2) there exists $\alpha_{0} \in(0,1)$ such that there is no nontrivial solution of (3.8) for $\alpha \in\left[0, \alpha_{0}\right]$.

Proof. (1) Let $\lambda_{1}, \lambda_{2} \in \mathbb{R}^{+}$with $\lambda_{1}<\lambda_{2}$ and $\lambda \in\left[\lambda_{1}, \lambda_{2}\right]$. Let $\left(x_{1}, x_{2}\right) \in \Theta_{0}$ be a solution of (3.8). Then for $i \in\{1,2\}$ we find that

$$
\ln \left(b_{i}+x_{i}(t)\right)-\ln \left(b_{i}+x_{i}(s)\right)=-\alpha \lambda \int_{s}^{t} \sum_{j=1}^{2} a_{i j} x_{j}(w-\tau / \lambda) d w
$$

for $t, s \in \mathbb{R}$ with $s \leqslant t$. Let $s \in \mathbb{R}$ such that $x_{i}(s)=0$. Then we have

$$
b_{i}+x_{i}(t)=b_{i} \exp \left(-\alpha \lambda \int_{s}^{t} \sum_{j=1}^{2} a_{i j} x_{j}(w-\tau / \lambda) d w\right) .
$$

Then noting that $x_{i}(t)>-b_{i}$ for $i=1,2$, and taking into account assumption (A0) we obtain

$$
\begin{equation*}
x_{1}(t)<d_{3}=b_{1} \exp \left(2 \pi \lambda\left(a_{11} b_{1}+a_{12} b_{2}\right)\right)-b_{1} \quad \text { for all } t \in[0,2 \pi] . \tag{3.9}
\end{equation*}
$$

Similarly as above we obtain

$$
\begin{equation*}
x_{2}(t)<d_{4}=b_{2} \exp \left(2 \pi \lambda\left(a_{21} b_{1}+a_{22} b_{2}\right)\right)-b_{2} \quad \text { for all } t \in[0,2 \pi] . \tag{3.10}
\end{equation*}
$$

On the other hand, for all $t \in[0,2 \pi]$ we have

$$
\begin{aligned}
& -b_{1}<-d_{1}=b_{1} \exp \left(-2 \pi \lambda\left(a_{11} d_{3}+a_{12} d_{4}\right)\right)-b_{1}<x_{1}(t) \quad \text { and } \\
& -b_{2}<-d_{2}=b_{2} \exp \left(-2 \pi \lambda\left(a_{21} d_{3}+a_{22} d_{4}\right)\right)-b_{2}<x_{2}(t)
\end{aligned}
$$

From (3.8) and the inequalities above we find that

$$
\begin{equation*}
\left|\dot{x}_{i}(t)\right| \leqslant \max _{i=1,2} C\left(\left|a_{i 1}\right| d_{3}+\left|a_{i 2}\right| d_{4}\right)\left(b_{i}+d_{i+2}\right) \quad \text { for } t \in[0,1] \text { and } i=1,2 . \tag{3.11}
\end{equation*}
$$

We also have by differentiating the both sides of (3.8) and using the inequalities above that $\left\{\left|\ddot{x}_{i}(t)\right|: t \in[0,2 \pi], i=1,2\right\}$ is bounded, which completes the proof of (1).
(2) Suppose that there exists a sequence $\left\{\alpha_{n}\right\} \subset \mathbb{R}^{+}$and $\left\{\left(x_{1 n}, x_{2 n}\right)\right\} \subset E$ such that $\lim _{n \rightarrow \infty} \alpha_{n}=0$ and each $\left(x_{1 n}, x_{2 n}\right)$ is a solution of (3.8) with $\alpha=\alpha_{n}$. Then (3.9) and (3.10) holds with $x_{1}$ and $x_{2}$ replaced by $x_{1 n}$ and $x_{2 n}$, respectively. Then we have that $\lim _{n \rightarrow \infty}\left|x_{i n}\right|_{\infty}=0$, $i=1,2$. By subtracting subsequences, we may assume, without any loss of generality, that $\left|x_{1 n}\right|_{\infty} \geqslant\left|x_{2 n}\right|_{\infty}$ for all $n \geqslant 1$. We put $u_{i n}(t)=x_{i n} /\left|x_{1 n}\right|_{\infty}$ for $i=1,2$. Then we have that $u_{\text {in }} \in H$ for $n \geqslant 1$ and $i=1,2$. We also have $\left|u_{1 n}\right|_{\infty}=1$ for all $n \geqslant 1$. Then it follows that for $n$ sufficiently large $\left|\dot{u}_{1 n}(t)\right| \leqslant 2 b_{1} \alpha_{n} \lambda\left(a_{11}+a_{12}\right)$ for all $t \in[0,2 \pi]$. That is $\lim _{n \rightarrow \infty}\left|\dot{u}_{1 n}\right|_{\infty}=0$. This contradicts to the fact that $\left|u_{1 n}\right|_{\infty}=1$ for all $n \geqslant 1$, which completes the proof of (2).

## 4. Homotopies of admissible $S^{1}$-equivariant mappings

The aim of this section is to define an open, bounded $S^{1}$-invariant subset $\Omega_{\lambda_{1}, \lambda_{2}} \subset \Theta_{0} \times \mathbb{R}^{+} \subset$ $E \times \mathbb{R}^{+}$such that the homotopy $Q-F(\cdot, \theta)$, defined by (3.5), does not vanish on $\partial \Omega_{\lambda_{1}, \lambda_{2}}$. We underline that solutions of equation $Q\left(\left(x_{1}, x_{2}\right), \lambda\right)=F\left(\left(\left(x_{1}, x_{2}\right), \lambda\right), 1\right)$ in $\Omega_{\lambda_{1}, \lambda_{2}}$ are exactly the periodic solutions of problem (3.2) in $\Omega_{\lambda_{1}, \lambda_{2}}$.

We finish this section with Lemma 4.4, where we reduce the computation of the $S^{1}$-degree of $Q-F(\cdot, 1)$ on $\Omega_{\lambda_{1}, \lambda_{2}}$ to the computation of the $S^{1}$-degree of $Q-F(\cdot, 0)$ on $\Omega_{\lambda_{1}, \lambda_{2}}$.

We first consider the following eigenvalue problem associated with problem (3.2)

$$
\left\{\begin{array}{l}
\dot{u}(t)=-\gamma \lambda b_{1}\left(a_{11} u(t-\tau / \lambda)+a_{12} v(t-\tau / \lambda)\right),  \tag{4.1}\\
\dot{v}(t)=-\gamma \lambda b_{2}\left(a_{21} u(t-\tau / \lambda)+a_{22} v(t-\tau / \lambda)\right),
\end{array}\right.
$$

where $\lambda, \tau>0, \gamma \in \mathbb{R}$ and $(u, v) \in E$. By assumption (1.5) $\mu_{1} \neq \mu_{2}$. Hence from assumption (A2) it follows that there exists a nondegenerate matrix $P$ such that

$$
P\left[\begin{array}{ll}
a_{11} b_{1} & a_{12} b_{1}  \tag{4.2}\\
a_{21} b_{2} & a_{22} b_{2}
\end{array}\right] P^{-1}=\left[\begin{array}{cc}
\mu_{1} & 0 \\
0 & \mu_{2}
\end{array}\right]
$$

Then linear problem (4.1) is transformed into the form

$$
\left\{\begin{array}{l}
\dot{u_{1}}(t)=-\gamma \lambda \mu_{1} u_{1}(t-\tau / \lambda),  \tag{4.3}\\
\dot{u_{2}}(t)=-\gamma \lambda \mu_{2} u_{2}(t-\tau / \lambda),
\end{array}\right.
$$

by putting $\left(u_{1}, u_{2}\right)=P^{-1}(u, v)$.
Let $i=1$, 2. We put $u_{i}(t)=\sum_{k=1}^{\infty}\left(c_{k} \cos k t+s_{k} \sin k t\right)$, where $\left\{c_{k}\right\},\left\{s_{k}\right\} \subset \mathbb{R}$. Then $\dot{u}_{i}(t)=$ $\sum_{k=1}^{\infty}\left(\left(k s_{k}\right) \cos k t+\left(-k c_{k}\right) \sin k t\right)$ and

$$
\begin{aligned}
& -\gamma \lambda \mu_{i} u\left(t-\tau^{\prime}\right) \\
& =-\gamma \lambda \mu_{i} \sum_{k=1}^{\infty}\left\{c_{k}\left(\cos k t \cos k \tau^{\prime}+\sin k \tau^{\prime} \sin k t\right)+s_{k}\left(\cos k \tau^{\prime} \sin k t-\sin k \tau^{\prime} \cos k t\right)\right\} \\
& =-\gamma \lambda \mu_{i} \sum_{k=1}^{\infty}\left\{\left(c_{k} \cos k \tau^{\prime}-s_{k} \sin k \tau^{\prime}\right) \cos k t+\left(c_{k} \sin k \tau^{\prime}+s_{k} \cos k \tau^{\prime}\right) \sin k t\right\},
\end{aligned}
$$

where $\tau^{\prime}=\tau / \lambda$. If $u_{1}$ is a nontrivial solution of (4.3), then

$$
\begin{aligned}
& k s_{k}=-\gamma \lambda \mu_{1}\left(c_{k} \cos k \tau^{\prime}-s_{k} \sin k \tau^{\prime}\right), \\
& -k c_{k}=-\gamma \lambda \mu_{1}\left(c_{k} \sin k \tau^{\prime}+s_{k} \cos k \tau^{\prime}\right)
\end{aligned}
$$

for all $k \in \mathbb{N}$. That is we obtain the following system of linear equations

$$
c_{k} \cos k \tau^{\prime}-s_{k}\left(\sin k \tau^{\prime}-\frac{k}{\gamma \lambda \mu_{1}}\right)=0, \quad c_{k}\left(\sin k \tau^{\prime}-\frac{k}{\gamma \lambda \mu_{1}}\right)+s_{k} \cos k \tau^{\prime}=0
$$

for all $k \in \mathbb{N}$. Then $\cos ^{2} k \tau^{\prime}+\left(\sin k \tau^{\prime}-\frac{k}{\gamma \lambda \mu_{1}}\right)^{2}=0$ and therefore we find that

$$
\begin{equation*}
\frac{k}{\gamma \lambda \mu_{1}}=1, \quad \frac{k \tau}{\lambda}=\frac{\pi}{2}+2 n \pi, \quad \text { for some } n \in \mathbb{N} \cup\{0\} . \tag{4.4}
\end{equation*}
$$

If $u_{2}$ is a nontrivial solution of (4.3), then by the same argument as above we obtain that

$$
\begin{equation*}
\frac{k}{\gamma \lambda \mu_{2}}=1, \quad \frac{k \tau}{\lambda}=\frac{\pi}{2}+2 n \pi, \quad \text { for some } n \in \mathbb{N} \cup\{0\} . \tag{4.5}
\end{equation*}
$$

Consequently, we have that the eigenvalue $\gamma$ of problem (4.3) is of the form

$$
\begin{equation*}
\gamma=\frac{1}{\mu_{i} \tau}\left(\frac{\pi}{2}+2 n \pi\right), \quad i=1,2 \text { and } n \in \mathbb{N} . \tag{4.6}
\end{equation*}
$$

Based on the observation above we obtain the following lemma.
Lemma 4.1. Let $\lambda_{1}, \lambda_{2} \in \mathbb{R}^{+}$with $\lambda_{1}<\lambda_{2}$. Then the origin is an isolated solution of (3.2), i.e., there exists $m_{1}>0$ such that if $\left(\left(x_{1}, x_{2}\right), \lambda\right) \in \Theta_{0} \times\left[\lambda_{1}, \lambda_{2}\right]$ is a nontrivial solution of (3.2) then $\left(x_{1}, x_{2}\right) \notin\left\{\left(x_{1}, x_{2}\right) \in \Theta_{0}:\left\|x_{1}\right\| \leqslant m_{1},\left\|x_{2}\right\| \leqslant m_{1}\right\}$.

Proof. Let $\lambda_{1}, \lambda_{2} \in \mathbb{R}^{+}$with $\lambda_{1}<\lambda_{2}$. Suppose that there exists a sequence $\left\{\left(\left(x_{1 n}, x_{2 n}\right), \lambda_{n}\right)\right\} \subset$ $E \times \mathbb{R}^{+}$such that each $\left(\left(x_{1 n}, x_{2 n}\right), \lambda_{n}\right)$ is a solution of (3.2) and

$$
\lim _{n \rightarrow \infty}\left\|x_{1 n}\right\|=\lim _{n \rightarrow \infty}\left\|x_{2 n}\right\|=0
$$

We may assume that $\lim _{n \rightarrow \infty} \lambda_{n}=\lambda_{0} \in\left[\lambda_{1}, \lambda_{2}\right]$ and $\left|x_{2 n}\right|_{\infty} \leqslant\left|x_{1 n}\right|_{\infty}$ for all $n \geqslant 1$. We put $u_{i n}(t)=x_{i n} /\left|x_{1 n}\right|_{\infty}$ for each $n \geqslant 1$ and $i=1,2$. Then we have

$$
\left\{\begin{array}{l}
\dot{u}_{1 n}(t)=-\lambda_{n}\left(a_{11} u_{1 n}\left(t-\tau / \lambda_{n}\right)+a_{12} u_{2 n}\left(t-\tau / \lambda_{n}\right)\right)\left(b_{1}+x_{1 n}(t)\right), \\
\dot{u}_{2 n}(t)=-\lambda_{n}\left(a_{21} u_{1 n}\left(t-\tau / \lambda_{n}\right)+a_{22} u_{2 n}\left(t-\tau / \lambda_{n}\right)\right)\left(b_{2}+x_{2 n}(t)\right) .
\end{array}\right.
$$

Then one can see that $\sup \left\{\left\|\dot{u}_{i n}\right\|: n \geqslant 1, i=1,2\right\}<\infty$. By differentiating the equalities above, we also have that $\sup \left\{\left\|\ddot{u}_{i n}\right\|: n \geqslant 1, i=1,2\right\}<\infty$. Therefore we may assume that for each $i$, $u_{i n} \rightarrow u_{i}$ and $\dot{u}_{i n} \rightarrow \dot{u}_{i}$ strongly in $E$. Then we have

$$
\left\{\begin{array}{l}
\dot{u}_{1}(t)=-\lambda_{0} b_{1}\left(a_{11} u_{1}\left(t-\tau / \lambda_{0}\right)+a_{12} u_{2}\left(t-\tau / \lambda_{0}\right)\right),  \tag{4.7}\\
\dot{u}_{2}(t)=-\lambda_{0} b_{2}\left(a_{21} u_{1}\left(t-\tau / \lambda_{0}\right)+a_{22} u_{2}\left(t-\tau / \lambda_{0}\right)\right)
\end{array}\right.
$$

That is, (4.1) holds with $\gamma=1$. By assumption (1.5), we have that (4.6) does not hold with $\gamma=1$. Therefore problem (4.7) has no nontrivial solution. Then $u_{1} \equiv 0$. This contradicts the definition of $u_{1}$.

Now fix $\lambda_{1}, \lambda_{2} \in \mathbb{R}^{+}$with $\lambda_{1}<\lambda_{2}$ and $\alpha_{0}, m_{0}, m_{1},\left\{d_{i}\right\}_{1 \leqslant i \leqslant 4}$ be the positive numbers satisfying the assertion of Lemmas 3.1 and 4.1. We may assume without any loss of generality that $d_{i+2}>\left(b_{i}+d_{i}\right) / 2$ for $i=1,2$. Define open bounded $S^{1}$-invariant subsets

$$
\begin{gathered}
\Theta_{M}=\left\{\left(x_{1}, x_{2}\right) \in E:-\frac{b_{i}+d_{i}}{2}<x_{i}(t)<2 d_{i+2}, \text { for } i=1,2, t \in[0,2 \pi]\right\}, \\
\widetilde{\Theta}_{M}=\left\{\left(x_{1}, x_{2}\right) \in E:-d_{i}<x_{i}(t)<d_{i+2}, \text { for } i=1,2, t \in[0,2 \pi]\right\}
\end{gathered}
$$

and closed $S^{1}$-invariant subset as follows

$$
\Theta_{m_{1}}=\left\{\left(x_{1}, x_{2}\right) \in E:\left\|x_{i}\right\| \leqslant m_{1} / 2, i=1,2\right\}
$$

Since $b_{i}>d_{i}, i=1,2$, and $m_{1}>0$ can be chosen sufficiently small

$$
\Theta_{m_{1}} \subsetneq \widetilde{\Theta}_{M} \subsetneq \Theta_{M} \subsetneq \Theta_{0}
$$

Moreover, define open bounded $S^{1}$-invariant subsets in the following way

$$
\begin{gather*}
\Theta_{1}=\left\{\left(x_{1}, x_{2}\right) \in E: \int_{0}^{2 \pi} \dot{x}_{i}(t)^{2} d t<2 \pi m_{0}^{2}, i=1,2\right\}, \quad \widetilde{\Theta}=\left(\widetilde{\Theta}_{M} \cap \Theta_{1}\right) \backslash \Theta_{m_{1}} \\
\Theta=\left(\Theta_{M} \cap \Theta_{1}\right) \backslash \Theta_{m_{1}} \tag{4.8}
\end{gather*}
$$

and notice that $\widetilde{\Theta}_{M} \subset \operatorname{cl}\left(\widetilde{\Theta}_{M}\right) \subset \Theta$. Then we can choose $\delta_{0}>0$ such that $\operatorname{dist}^{2}\left(\widetilde{\Theta}_{M}, \partial \Theta\right) \geqslant \delta_{0}$. Let $\xi:[0,+\infty) \rightarrow(0,1]$ be a smooth function such that

$$
\xi(t)= \begin{cases}\alpha_{0} & \text { for } t=0  \tag{4.9}\\ \text { strictly increasing } & \text { for } 0<t<\delta_{0} \\ 1 & \text { for } t \geqslant \delta_{0}\end{cases}
$$

We put that

$$
\begin{equation*}
\beta\left(x_{1}, x_{2}\right)=\xi\left(\operatorname{dist}^{2}\left(\left(x_{1}, x_{2}\right), \partial \Theta\right)\right) \quad \text { for }\left(x_{1}, x_{2}\right) \in E . \tag{4.10}
\end{equation*}
$$

Then $\beta \in C^{1}(E ; \mathbb{R})$ and we have

$$
\beta\left(x_{1}, x_{2}\right)= \begin{cases}1 & \text { for }\left(x_{1}, x_{2}\right) \in \Theta_{m_{1}},  \tag{4.11}\\ \alpha_{0} & \text { for }\left(x_{1}, x_{2}\right) \in \operatorname{cl}\left(E \backslash\left(\Theta_{M} \cap \Theta_{1}\right)\right) .\end{cases}
$$

Put $\Omega_{\lambda_{1}, \lambda_{2}}=\Theta \times\left(\lambda_{1}, \lambda_{2}\right)$ and

$$
\mathcal{S}=\left\{\left(\left(\left(x_{1}, x_{2}\right), \lambda\right), \theta\right) \in \Omega_{\lambda_{1}, \lambda_{2}} \times[0,1]: F\left(\left(\left(x_{1}, x_{2}\right), \lambda\right), \theta\right)=Q\left(\left(x_{1}, x_{2}\right), \lambda\right)\right\} .
$$

Lemma 4.2. Under the above assumptions $\mathcal{S} \cap\left(\left(\partial \Theta \times\left(\lambda_{1}, \lambda_{2}\right)\right) \times[0,1]\right)=\emptyset$.
Proof. Suppose that $\left(\left(\left(x_{1}, x_{2}\right), \lambda\right), \theta\right) \in \mathcal{S}$. Then (3.6) holds for $\left(\left(\left(x_{1}, x_{2}\right), \lambda\right), \theta\right)$. Multiplying (3.6) by $\theta$, and denoting $u_{i}(t)=\theta x_{i}(t)$ for $i=1,2$, we find that

$$
\left\{\begin{array}{l}
\dot{u}_{1}(t)=-\lambda \beta\left(x_{1}, x_{2}\right)\left(a_{11} u_{1}(t-\tau / \lambda)+a_{12} u_{2}(t-\tau / \lambda)\right)\left(b_{1}+u_{1}(t)\right), \\
\dot{u}_{2}(t)=-\lambda \beta\left(x_{1}, x_{2}\right)\left(a_{21} u_{1}(t-\tau / \lambda)+a_{22} u_{2}(t-\tau / \lambda)\right)\left(b_{2}+u_{2}(t)\right),
\end{array}\right.
$$

holds. By Lemma 3.1 we obtain $\int_{0}^{2 \pi} \dot{x}_{i}(t)^{2} d t<2 \pi m_{0}^{2}, i=1,2$. Then we have $\left(x_{1}, x_{2}\right) \notin \partial \Theta_{1}$.
If $\left(x_{1}, x_{2}\right) \in \partial \Theta_{M}$, then we have that $\beta\left(x_{1}, x_{2}\right)=\alpha_{0}$. Then by (2) of Lemma 3.1, we have that $u_{1} \equiv u_{2} \equiv 0$. This contradicts to $\left(x_{1}, x_{2}\right) \in \operatorname{cl}(\Theta)$. If $\left(x_{1}, x_{2}\right) \in \partial \Theta_{m_{1}}$, then $\beta\left(x_{1}, x_{2}\right)=1$ and $\left\|u_{i}\right\|<m_{1} / 2$, for $i=1,2$. Then by Lemma 4.1, we have that $x_{1} \equiv x_{2} \equiv 0$. Thus we have that $\mathcal{S} \cap \partial \Omega_{\lambda_{1}, \lambda_{2}}=\emptyset$, which completes the proof.

The following result is known. For completeness, we give a proof.
Lemma 4.3. Suppose that $\tau=0$. Then problem (3.6) does not have nonstationary periodic solution $\left(\left(x_{1}, x_{2}\right), \lambda\right) \in E \times \mathbb{R}^{+}$for any $\theta \in[0,1]$.

Proof. Let $\theta \in[0,1]$ and $\left(\left(x_{1}, x_{2}\right), \lambda\right) \in \Theta_{0} \times \mathbb{R}^{+}$satisfy (3.6). We first consider the case that $\theta>0$. Since $\tau=0$, problem (3.6) reduces to the problem

$$
\left\{\begin{array}{l}
\dot{x}_{1}(t)=-\lambda\left(a_{11} x_{1}(t)+a_{12} x_{2}(t)\right)\left(b_{1}+\theta x_{1}(t)\right),  \tag{4.12}\\
\dot{x}_{2}(t)=-\lambda\left(a_{21} x_{1}(t)+a_{22} x_{2}(t)\right)\left(b_{2}+\theta x_{2}(t)\right) .
\end{array}\right.
$$

We integrate the both sides of (4.12) from 0 to $2 \pi$. Then by the periodicity, we have

$$
\left\{\begin{array}{l}
a_{11} \int_{0}^{2 \pi} x_{1}(t)^{2} d t+a_{12} \int_{0}^{2 \pi} x_{1}(t) x_{2}(t) d t=0 \\
a_{21} \int_{0}^{2 \pi} x_{1}(t) x_{2}(t) d t+a_{22} \int_{0}^{2 \pi} x_{2}(t)^{2} d t=0
\end{array}\right.
$$

Then one can see that $x_{1}=x_{2} \equiv 0$ from condition (A1). We next consider the case that $\theta=0$. In this case we multiply Eqs. (4.12) by $x_{i}$ and integrate over $[0,2 \pi]$. Then we have the equalities above. This completes the proof.

Lemma 4.4. Suppose that $\lambda_{1}=\frac{\tau}{2 j_{1} \pi}<\lambda_{2}=\frac{\tau}{2 j_{2} \pi}$, where $j_{1}, j_{2} \in \mathbb{N}$. Then

$$
\operatorname{Deg}\left(Q-F(\cdot, 0), \Omega_{\lambda_{1}, \lambda_{2}}\right)=\operatorname{Deg}\left(Q-F(\cdot, 1), \Omega_{\lambda_{1}, \lambda_{2}}\right)
$$

Proof. To prove the assertion, it is sufficient to show that there exists no solution of (3.6) in $\partial \Omega_{\lambda_{1}, \lambda_{2}}=\partial\left(\Theta \times\left(\lambda_{1}, \lambda_{2}\right)\right)=\operatorname{cl}(\Theta) \times\left\{\lambda_{1}, \lambda_{2}\right\} \cup \partial \Theta \times\left(\lambda_{1}, \lambda_{2}\right)$. We first see that there exists no solution on $\operatorname{cl}(\Theta) \times\left\{\lambda_{1}, \lambda_{2}\right\}$. From the definitions of $\lambda_{1}, \lambda_{2}$, we have that problem (3.6) is equivalent to (4.12) with $\lambda=\lambda_{1}$ or $\lambda=\lambda_{2}$. Then by Lemma 4.3, we find that $x_{1}=x_{2}=0$. This contradicts to the assumption that $\left(x_{1}, x_{2}\right) \in \operatorname{cl}(\Theta)$. We also have by Lemma 4.2 that there exists no solution of (3.6) in $\partial \Theta \times\left(\lambda_{1}, \lambda_{2}\right)$, which completes the proof.

## 5. Proof of Theorem 1.1

Throughout this section we assume that $\lambda_{1}=\frac{\tau}{2 j_{1} \pi}<\lambda_{2}=\frac{\tau}{2 j_{2} \pi}$, where $j_{1}, j_{2} \in \mathbb{N}$ and put $\Omega_{\lambda_{1}, \lambda_{2}}=\Theta \times\left(\lambda_{1}, \lambda_{2}\right)$. From Theorem 2.1 and Lemma 4.4 it follows that to finish the proof of Theorem 1.1, it is sufficient to show that $\operatorname{Deg}\left(Q-F(\cdot, 0), \Omega_{\lambda_{1}, \lambda_{2}}\right) \neq \Theta \in \Gamma$. But the mapping $F(\cdot, 0)$ is still too complicated to calculate the $S^{1}$-degree. Therefore we will provide another homotopy $G$ of $S^{1}$-equivariant compact mappings such that $F(\cdot, 0)=G(\cdot, 0)$ and the $S^{1}$-degree of $Q-G(\cdot, 1)$ on $\Omega_{\lambda_{1}, \lambda_{2}}$ is easy to compute.

We fix a $C^{1}$-mapping $\sigma:\left[\lambda_{1}, \lambda_{2}\right] \rightarrow\left[\lambda_{1}, \lambda_{2}\right]$ such that $\sigma$ is increasing on $\left[\lambda_{1}, \lambda_{2}\right]$ with $\sigma\left(\lambda_{1}\right)=\lambda_{1}$ and $\sigma\left(\lambda_{2}\right)=\lambda_{2}$, and

$$
\sigma\left(\lambda_{k, n}\right)=\lambda_{k, n} \quad \text { and } \quad \dot{\sigma}\left(\lambda_{k, n}\right)=0 \quad \text { for each } \lambda_{k, n}=\frac{k \tau}{\frac{\pi}{2}+2 n \pi} \in\left[\lambda_{1}, \lambda_{2}\right], k, n \in \mathbb{N} .
$$

We now define a homotopy of $S^{1}$-equivariant mappings $G: \Omega_{\lambda_{1}, \lambda_{2}} \times[0,1] \rightarrow E$ by

$$
\begin{equation*}
G\left(\left(\left(x_{1}, x_{2}\right), \lambda\right), \theta\right)=-(\theta \sigma(\lambda)+(1-\theta) \lambda) \int_{0}^{t} \beta\left(x_{1}, x_{2}\right) \mathcal{N}\left(\left(x_{1}, x_{2}\right), 0\right) d s \tag{5.1}
\end{equation*}
$$

By definition of $\mathcal{N}(\cdot, 0)$, we have $G\left(\left(\left(x_{1}, x_{2}\right), \lambda\right), \theta\right) \in E$ for $\left(\left(\left(x_{1}, x_{2}\right), \lambda\right), \theta\right) \in\left(E \times \mathbb{R}^{+}\right) \times$ $[0,1]$. If $\left(\left(\left(x_{1}, x_{2}\right), \lambda\right), \theta\right) \in \Omega_{\lambda_{1}, \lambda_{2}} \times[0,1]$ satisfies $Q\left(\left(x_{1}, x_{2}\right), \lambda\right)=G\left(\left(\left(x_{1}, x_{2}\right), \lambda\right), \theta\right)$ then

$$
\left\{\begin{array}{l}
\dot{x}_{1}(t)=-(\theta \sigma(\lambda)+(1-\theta) \lambda) b_{1} \beta\left(x_{1}, x_{2}\right)\left(a_{11} x_{1}(t-\tau / \lambda)+a_{12} x_{2}(t-\tau / \lambda)\right)  \tag{5.2}\\
\dot{x}_{2}(t)=-(\theta \sigma(\lambda)+(1-\theta) \lambda) b_{2} \beta\left(x_{1}, x_{2}\right)\left(a_{21} x_{1}(t-\tau / \lambda)+a_{22} x_{2}(t-\tau / \lambda)\right) .
\end{array}\right.
$$

Lemma 5.1. Under the above assumptions:

$$
\operatorname{Deg}\left(Q-F(\cdot, 0), \Omega_{\lambda_{1}, \lambda_{2}}\right)=\operatorname{Deg}\left(Q-G(\cdot, 1), \Omega_{\lambda_{1}, \lambda_{2}}\right)
$$

Proof. By the same argument as in the proof of Lemma 4.4, we see that

$$
\operatorname{Deg}\left(Q-G(\cdot, 0), \Omega_{\lambda_{1}, \lambda_{2}}\right)=\operatorname{Deg}\left(Q-G(\cdot, 1), \Omega_{\lambda_{1}, \lambda_{2}}\right)
$$

Then since $G(\cdot, 0)=F(\cdot, 0)$, we have by Lemma 4.4 that the assertion holds.

For $n, m \in \mathbb{N}$ define $\Phi(n, m)=\left\{j \in \mathbb{N}:\left[\frac{n}{j}\right]<\left[\frac{m}{j}\right]=\frac{m}{j}\right\}$. Notice that if $n<m$ then $\Phi(n, m) \neq \emptyset$.

The following lemma plays crucial role in our article.

Lemma 5.2. Let assumptions of Theorem 1.1 be fulfilled. If $n_{1}<n_{2}, j \in \Phi\left(n_{1}, n_{2}\right)$ and $\lambda_{1}=$ $\frac{\tau}{2(j+1) \pi}, \lambda_{2}=\frac{\tau}{2 j \pi}$ then $\operatorname{Deg}\left(Q-G(\cdot, 1), \Omega_{\lambda_{1}, \lambda_{2}}\right) \neq \Theta \in \Gamma$.

Proof. Before we prove this lemma, we outline the main steps of the proof. Namely, we will prove that $(Q-G(\cdot, 1))^{-1}(0) \cap \Omega_{\lambda_{1} \lambda_{2}}$ consists of a finite number of nondegenerate orbits $S^{1} \cdot a_{1}, \ldots, S^{1} \cdot a_{p}$. Since these orbits are nondegenerate, there are open bounded $S^{1}$-invariant subsets $U_{i} \subset \operatorname{cl}\left(U_{i}\right) \subset \Omega_{\lambda_{1} \lambda_{2}}, i=1, \ldots, p$, such that $(Q-G(\cdot, 1))^{-1}(0) \cap U_{i}=S^{1} \cdot a_{i}, i=$ $1, \ldots, p$. Moreover, we will prove that there are $k_{0} \in \mathbb{N}$ and $1 \leqslant i_{0} \leqslant p$ such that $S_{a_{i_{0}}}^{1}=\mathbb{Z}_{k_{0}}$ and $S_{a_{i}}^{1} \neq \mathbb{Z}_{k_{0}}$ for every $i \neq i_{0}$.

By Theorem 2.1 we obtain

$$
\operatorname{Deg}\left(Q-G(\cdot, 1), \Omega_{\lambda_{1}, \lambda_{2}}\right)=\operatorname{Deg}\left(Q-G(\cdot, 1), U_{1}\right)+\cdots+\operatorname{Deg}\left(Q-G(\cdot, 1), U_{p}\right) \in \Gamma
$$

From the above and Theorem 2.2 we obtain that

$$
\begin{aligned}
\operatorname{deg}_{\mathbb{Z}_{k_{0}}}\left(Q-G(\cdot, 1), \Omega_{\lambda_{1}, \lambda_{2}}\right) & =\operatorname{deg}_{\mathbb{Z}_{k_{0}}}\left(Q-G(\cdot, 1), U_{1}\right)+\cdots+\operatorname{deg}_{\mathbb{Z}_{k_{0}}}\left(Q-G(\cdot, 1), U_{p}\right) \\
& =\operatorname{deg}_{\mathbb{Z}_{k_{0}}}\left(Q-G(\cdot, 1), U_{i_{0}}\right) \neq 0 \in \mathbb{Z}
\end{aligned}
$$

Let us begin the proof. First of all notice that since $\mu_{1} \neq \mu_{2}$ applying change of coordinates (4.2) to system (5.2) we obtain the following equivalent system:

$$
\left\{\begin{array}{l}
\dot{x}_{1}(t)=-\sigma(\lambda) \beta\left(P^{-1}\left(x_{1}, x_{2}\right)\right) \mu_{1} x_{1}(t-\tau / \lambda),  \tag{5.3}\\
\dot{x}_{2}(t)=-\sigma(\lambda) \beta\left(P^{-1}\left(x_{1}, x_{2}\right)\right) \mu_{2} x_{2}(t-\tau / \lambda) .
\end{array}\right.
$$

Notice that system (5.3) does not have solutions on $\partial\left(P \Omega_{\lambda_{1}, \lambda_{2}}\right)=\partial\left(P \Theta \times\left(\lambda_{1}, \lambda_{2}\right)\right)$. Since $\beta\left(P^{-1}\left(x_{1}, x_{2}\right)\right)=\alpha_{0}$ for any $\left(x_{1}, x_{2}\right) \in P\left(\operatorname{cl}\left(E \backslash\left(\Theta_{M} \cap \Theta_{1}\right)\right)\right)$, we find that (5.3) does not have solutions on $P\left(E \backslash\left(\Theta_{M} \cap \Theta_{1}\right)\right) \times\left(\lambda_{1}, \lambda_{2}\right)$.

Therefore we can choose $R \gg r>0$ such that

$$
\begin{aligned}
& \operatorname{Deg}\left(Q-P G(\cdot, 1) P^{-1}, P \Theta \times\left(\lambda_{1}, \lambda_{2}\right)\right) \\
& \quad=\operatorname{Deg}\left(Q-P G(\cdot, 1) P^{-1},\left(D_{R}(E) \backslash \operatorname{cl}\left(D_{r}(E)\right) \times\left(\lambda_{1}, \lambda_{2}\right)\right)\right)
\end{aligned}
$$

where

$$
\begin{gathered}
P\left(c l\left(\Theta_{M} \cap \Theta_{1}\right)\right) \subset D_{R}(E)=\{x \in E:\|x\|<R\} \\
c l\left(D_{r}(E)\right)=\{x \in E:\|x\| \leqslant r\} \subset P\left(\Theta_{m_{1}}\right)
\end{gathered}
$$

Here we replace $\beta$ by a function for which the calculation of degree is easier.
Let $\tilde{\xi} \in C^{\infty}([0,+\infty),[0,1])$ with $\tilde{\xi}(t)=1$ for $0 \leqslant t \leqslant \sqrt{r}, \tilde{\xi}(t)=\alpha_{0}$ for $t \in[\sqrt{R},+\infty)$ and $\tilde{\xi}$ is strictly monotone decreasing on $[\sqrt{r}, \sqrt{R}]$. Define $\tilde{\beta}: E \rightarrow\left[\alpha_{0}, 1\right]$ as follows

$$
\begin{equation*}
\tilde{\beta}\left(x_{1}, x_{2}\right):=\tilde{\xi}\left(\left\|\left(x_{1}, x_{2}\right)\right\|^{2}\right) . \tag{5.4}
\end{equation*}
$$

We denote by $\widetilde{G}(\cdot, 1)$ the mapping $G(\cdot, 1)$ with $\beta$ replaced by $\tilde{\beta}$. Since maps $\beta, \tilde{\beta}$ coincide on $\partial\left(D_{R}(E) \backslash \operatorname{cl}\left(D_{r}(E)\right)\right),\left.\widetilde{G}(\cdot, 1)\right|_{\partial\left(D_{R}(E) \backslash c l\left(D_{r}(E)\right)\right)}=\left.G(\cdot, 1)\right|_{\partial\left(D_{R}(E) \backslash c l\left(D_{r}(E)\right)\right)}$ and by the homotopy invariance of the $S^{1}$-degree we have

$$
\begin{aligned}
& \operatorname{Deg}\left(Q-P G(\cdot, 1) P^{-1}, P \Theta \times\left(\lambda_{1}, \lambda_{2}\right)\right) \\
& \quad=\operatorname{Deg}\left(Q-P \widetilde{G}(\cdot, 1) P^{-1},\left(D_{R}(E) \backslash c l\left(D_{r}(E)\right) \times\left(\lambda_{1}, \lambda_{2}\right)\right)\right)
\end{aligned}
$$

For the simplicity of notation we will denote $\widetilde{G}(\cdot, 1), \tilde{\beta}$ and $\tilde{\xi}$ by $G(\cdot, 1), \beta$ and $\xi$, respectively.
If $a(t)=\left(a_{1}(t), a_{2}(t)\right)$ satisfies (5.3), then we have by (4.6) that $\lambda=\lambda_{k, n}=\frac{k \tau}{\pi / 2+2 n \pi}$ for some $k, n \in \mathbb{N}$ and $a_{i}(t) \in \operatorname{span}\{\cos k t, \sin k t\}$ for $i=1,2$. Then by the definition of $\sigma$, we find that $\sigma\left(\lambda_{k, n}\right)=\lambda_{k, n}=\frac{k \tau}{\pi / 2+2 n \pi}$ for some $k, n \geqslant 1$.

Now suppose that $a_{1} \not \equiv 0$. Then since $\sigma\left(\lambda_{k, n}\right) \beta\left(a_{1}, a_{2}\right) \mu_{1}=k$ and $\tau \mu_{1} \neq \pi / 2+2 n \pi$ for $n \geqslant 0$, we find that $\beta\left(a_{1}, a_{2}\right)<1$. Then taking into account (1.5) we obtain

$$
k<\lambda_{k, n} \mu_{1}=\frac{k \tau \mu_{1}}{\frac{\pi}{2}+2 n \pi}<k \frac{\frac{\pi}{2}+2\left(n_{1}+1\right) \pi}{\frac{\pi}{2}+2 n \pi} .
$$

Therefore we have $n \leqslant n_{1}$. On the other hand, we have by the definition that

$$
\begin{equation*}
\frac{\tau}{2(j+1) \pi} \leqslant \lambda_{k, n}=\frac{k \tau}{\left(\frac{\pi}{2}+2 n \pi\right)} \leqslant \frac{\tau}{2 j \pi} \tag{5.5}
\end{equation*}
$$

which is equivalent to $k j \leqslant n<k(j+1)$.
Therefore $1 \leqslant k \leqslant\left[n_{1} / j\right]$. Then noting that $\mu_{1} \neq \mu_{2}$, we have that

$$
(a(t), \lambda)=\left(\left(a_{1}(t), a_{2}(t)\right), \lambda\right)=\left(\left(c_{1, k} \cos k t, 0\right), \lambda_{k, n}\right) \quad \text { for some } 1 \leqslant k \leqslant\left[n_{1} / j\right], 1 \leqslant n \leqslant n_{1},
$$

where $c_{1, k}>0$ is such that $\beta\left(c_{1, k} \cos k t, 0\right) \lambda_{k, n} \mu_{1}=k$.
Similarly, we have that if $a_{2} \not \equiv 0$,

$$
(a(t), \lambda)=\left(\left(a_{1}(t), a_{2}(t)\right), \lambda\right)=\left(\left(0, c_{2, k} \cos k t\right), \lambda_{k, n}\right) \quad \text { for some } 1 \leqslant k \leqslant\left[n_{2} / j\right], 1 \leqslant n \leqslant n_{2},
$$

where $c_{2, k}>0$ is such that $\beta\left(0, c_{2, k} \cos k t\right) \lambda_{k, n} \mu_{2}=k$.
It is clear, that the map $s \rightarrow \beta(s u)$ is decreasing for any $u \in D_{R}(E) \backslash \operatorname{cl}\left(D_{r}(E)\right)$. Then since $\beta\left(a_{1}, a_{2}\right)<1$, the map $s \rightarrow \beta\left(s a_{1}, s a_{2}\right)$ is strictly decreasing in $[1-\varepsilon, 1+\varepsilon]$.

This implies that each $\left\{\left(\rho\left(e^{i \theta},\left(a_{1}(t), a_{2}(t)\right)\right), \lambda\right): \theta \in[0,2 \pi)\right\}$ is an isolated orbit satisfying (5.3).

Now fix $\left(a_{0}(t), \lambda_{k, n}\right)=\left(\left(a_{1}(t), a_{2}(t)\right), \lambda_{k, n}\right)=\left(\left(c_{1, k} \cos k t, 0\right), \lambda_{k, n}\right)$, where $1 \leqslant k \leqslant\left[n_{1} / j\right]$ and $1 \leqslant n<n_{1}$. Then $\left(\dot{a}_{0}(t), 0\right)=\left(\left(\dot{a}_{1}(t), \dot{a}_{2}(t)\right), 0\right)=\left(\left(-c_{1, k} k \sin k t, 0\right), 0\right)$ is the tangent vector to the orbit $S^{1} \cdot\left(a_{0}, \lambda_{k, n}\right)$ at $\left(a_{0}, \lambda_{k, n}\right)$.

Summing up, we have proved that $(Q-G(\cdot, 1))^{-1}(0)$ consists of a finite number of $S^{1}$-orbits. Below we prove that these orbits are nondegenerate.

For simplicity of notation we put $x=\left(x_{1}(t), x_{2}(t)\right)$. Then

$$
\begin{gather*}
f(x, \lambda)=\binom{f_{1}(x, \lambda)}{f_{2}(x, \lambda)}=\binom{\int_{0}^{t}-\sigma(\lambda) \beta(x) \mu_{1} x_{1}\left(s-\frac{\tau}{\lambda}\right) d s}{\int_{0}^{t}-\sigma(\lambda) \beta(x) \mu_{2} x_{2}\left(s-\frac{\tau}{\lambda}\right) d s}, \\
D_{x} f=\binom{D_{x} f_{1}}{D_{x} f_{2}}, \quad D_{\lambda} f=\binom{D_{\lambda} f_{1}}{D_{\lambda} f_{2}} . \tag{5.6}
\end{gather*}
$$

Then we obtain

$$
\begin{array}{rl}
D_{x} & f\left(a_{0}, \lambda_{k, n}\right)(v) \\
& =-\int_{0}^{t}\left(\begin{array}{c}
\left.k v_{1}\left(s-\frac{\tau}{\lambda_{k, n}}\right)+\begin{array}{l}
2 \lambda_{k, n} \mu_{1} a_{1}\left(s-\frac{\tau}{\lambda_{k, n}}\right) \xi^{\prime}\left(a_{0}\right)\left\langle a_{0}, v\right\rangle \\
k \frac{\mu_{2}}{\mu_{1}} v_{2}\left(s-\frac{\tau}{\lambda_{k, n}}\right)
\end{array}\right) d s \\
=-\int_{0}^{t}\binom{k v_{1}\left(s-\frac{\tau}{\lambda_{k, n}}\right)+2 \lambda_{k, n} \mu_{1} a_{1}\left(s-\frac{\pi}{2 k}\right) \xi^{\prime}\left(a_{0}\right)\left\langle a_{1}, v_{1}\right\rangle}{ k \frac{\mu_{2}}{\mu_{1}} v_{2}\left(s-\frac{\tau}{\lambda_{k, n}}\right)} d s \\
=\left(\begin{array}{cc}
T_{1} & 0 \\
0 & T_{2}
\end{array}\right)\binom{v_{1}}{v_{2}},
\end{array}\right.
\end{array}
$$

where $\xi^{\prime}\left(a_{0}\right)<0$. On the other hand,

$$
D_{\lambda} f(x, \lambda)=\binom{\int_{0}^{t}-\beta(x) \mu_{1}\left(\sigma^{\prime}(\lambda) x_{1}\left(s-\frac{\tau}{\lambda}\right)+\frac{\sigma(\lambda) \tau}{\lambda^{2}} \dot{x}_{1}\left(s-\frac{\tau}{\lambda}\right)\right) d s}{\int_{0}^{t}-\beta(x) \mu_{2}\left(\sigma^{\prime}(\lambda) x_{2}\left(s-\frac{\tau}{\lambda}\right)+\frac{\sigma(\lambda) \tau}{\lambda^{2}} \dot{x}_{2}\left(s-\frac{\tau}{\lambda}\right)\right) d s}
$$

and noting that $\sigma\left(\lambda_{k, n}\right)=\lambda_{k, n}, \sigma^{\prime}\left(\lambda_{k, n}\right)=0$ and that $a_{0}=\left(a_{1}, 0\right)$ we obtain

$$
\begin{align*}
D_{\lambda} f\left(a_{0}, \lambda_{k, n}\right) & =\binom{-\int_{0}^{t} \beta\left(a_{0}\right) \mu_{1} \lambda_{k, n} \frac{\tau}{\lambda_{k, n}^{2}} \dot{a}_{1}\left(s-\frac{\tau}{\lambda_{k, n}}\right) d s}{0} \\
& =\binom{-\int_{0}^{t} \frac{k \tau}{\lambda_{k, n}^{2}} \dot{a}_{1}\left(s-\frac{\pi}{2 k}\right) d s}{0}=\binom{-\frac{k \tau}{\lambda_{k, n}^{2}} a_{1}\left(t-\frac{\pi}{2 k}\right)}{0}=\binom{T_{3}}{0} . \tag{5.8}
\end{align*}
$$

Let us consider the following eigenvalue problem $v=\mu D_{x} f\left(a_{0}, \lambda_{k, n}\right) v$, i.e.,

$$
\binom{v_{1}}{v_{2}} v=\mu\binom{T_{1} v_{1}}{T_{2} v_{2}}
$$

for $v=\left(v_{1}, v_{2}\right) \in E$, which is equivalent to the following system

$$
\left\{\begin{array}{l}
\dot{v}_{1}(t)=-\mu\left(k v_{1}\left(t-\frac{\tau}{\lambda_{k, n}}\right)+2 \lambda_{k, n} \mu_{1} a_{1}\left(t-\frac{\pi}{2 k}\right) \xi^{\prime}\left(a_{0}\right)\left\langle a_{1}, v_{1}\right\rangle\right),  \tag{5.9}\\
\dot{v}_{2}(t)=-\mu \frac{k \mu_{2}}{\mu_{1}} v_{2}\left(t-\frac{\tau}{\lambda_{k, n}}\right) .
\end{array}\right.
$$

Since $\left\langle a_{1}, \dot{a}_{1}\right\rangle=0$, it is easy to verify that $\mu=1$ is the eigenvalue with corresponding eigenvector $\left(\dot{a}_{1}, 0\right)$.

Summing up, we obtain

$$
\begin{aligned}
Q-D f\left(a_{0}, \lambda_{k, n}\right) & =Q-\left(D_{x} f\left(a_{0}, \lambda_{k, n}\right), D_{\lambda} f\left(a_{0}, \lambda_{k, n}\right)\right)=Q-T \\
& =\left(\begin{array}{ccc}
I d & 0 & 0 \\
0 & I d & 0
\end{array}\right)-\left(\begin{array}{ccc}
T_{1} & 0 & T_{3} \\
0 & T_{2} & 0
\end{array}\right): E \times \mathbb{R} \rightarrow E
\end{aligned}
$$

is a surjection such that $\operatorname{ker}(Q-T)=\operatorname{span}\left\{\left(\left(\dot{a}_{1}, 0\right), 0\right)\right\}$. Notice that we have just proved that $S^{1}$-orbits of $(Q-G(\cdot, 1))^{-1}(0)$ are nondegenerate. In other words the assumptions of Theorem 2.2 are satisfied.

Since $\left(a_{0}(t), \lambda_{k, n}\right)$ is an isolated nondegenerate solution of (5.3) and $S_{a_{0}}^{1}=\mathbb{Z}_{k}$, applying Theorem 2.2 we obtain

$$
\begin{equation*}
\operatorname{deg}_{\mathbb{Z}_{k}}(Q-f, \Omega)= \pm 1, \quad \operatorname{deg}_{\mathbb{Z}_{k^{\prime}}}(Q-f, \Omega)=0 \quad \text { for } k^{\prime}>k \tag{5.10}
\end{equation*}
$$

for an open, bounded $S^{1}$-invariant subset $\Omega \subset \operatorname{cl}(\Omega) \subset \Omega_{\lambda_{1}, \lambda_{2}}$ such that $(Q-f)^{-1}(0) \cap \Omega=$ $S^{1} \cdot a_{0} \times\left\{\lambda_{k, n}\right\}$.

The same computation one can perform for $\left(a_{0}(t), \lambda\right)=\left(\left(a_{1}(t), a_{2}(t)\right), \lambda\right)=\left(\left(0, c_{2, k} \cos k t\right)\right.$, $\lambda_{k, n}$ ) for some $1 \leqslant k, n \leqslant n_{2}$, satisfying (5.5).

Summing up, we have proved that $(Q-f)^{-1}(0) \cap \Omega_{\lambda_{1} \lambda_{2}}=S^{1} \cdot a_{1} \cup \cdots \cup S^{1} \cdot a_{p}$, i.e., it consist of a finite number of nondegenerate $S^{1}$-orbits $S^{1} \cdot a_{1}, \ldots, S^{1} \cdot a_{p}$, each with nontrivial $S^{1}$-index, see formula (5.10). Since these orbits are nondegenerate, there are open bounded $S^{1}$-invariant subsets $U_{i} \subset c l\left(U_{i}\right) \subset \Omega_{\lambda_{1} \lambda_{2}}, i=1, \ldots, p$, such that $(Q-G(\cdot, 1))^{-1}(0) \cap U_{i}=S^{1} \cdot a_{i}, i=$ $1, \ldots, p$. And consequently by the properties of $S^{1}$-degree we obtain

$$
\operatorname{Deg}\left(Q-f, \Omega_{\lambda_{1}, \lambda_{2}}\right)=\operatorname{Deg}\left(Q-f, U_{1}\right)+\cdots+\operatorname{Deg}\left(Q-f, U_{p}\right) \in \Gamma
$$

Notice that for $k_{0}=\left[\frac{n_{2}}{j}\right]$ only $n=n_{2}$ satisfies (5.5). Therefore, there is exactly one solution of (5.3) in $\Omega_{\lambda_{1}, \lambda_{2}}$ of the form ( $\left.\left(0, c_{2, k_{0}} \cos k_{0} t\right), \lambda_{k_{0}, n_{2}}\right)$. Moreover, other solutions of (5.3) are of the form $\left(\left(c_{1, k} \cos k t, 0\right), \lambda_{k, n}\right)$ or $\left(\left(0, c_{2, k} \cos k t\right), \lambda_{k, n}\right)$, where $k<k_{0}$. In other words there is exactly one orbit with the isotropy group $\mathbb{Z}_{k_{0}}$.

Finally, combining Theorem 2.2 with (5.10) we obtain $\operatorname{deg}_{\mathbb{Z}_{k_{0}}}\left(Q-f, \Omega_{\lambda_{1}, \lambda_{2}}\right) \neq 0$, which completes the proof.

Proof of Theorem 1.1. Without loss of generality, we can assume that $n_{1}<n_{2}$. Fix $j \in$ $\Phi\left(n_{1}, n_{2}\right)$ and define $\lambda_{1}=\frac{\tau}{2(j+1) \pi}, \lambda_{2}=\frac{\tau}{2 j \pi}, \Omega_{\lambda_{1} \lambda_{2}}=\Theta \times\left(\lambda_{1}, \lambda_{2}\right)$, where $\Theta \subset E$ is an open bounded $S^{1}$-invariant subset defined by (4.8). In other words $\Omega_{\lambda_{1} \lambda_{2}}$ is a Cartesian product of an "annulus" $\Theta$ and an open interval $\left(\lambda_{1}, \lambda_{2}\right)$.

To complete the proof it is enough to show that $(Q-F(\cdot, 1))^{-1}(0) \cap \Omega_{\lambda_{1} \lambda_{2}} \neq \emptyset$, where the operator $F$ is defined by formula (3.4). By Theorem 2.1 it is enough to show that either $(Q-F(\cdot, 1))^{-1}(0) \cap \partial \Omega_{\lambda_{1} \lambda_{2}} \neq \emptyset$, or $\operatorname{Deg}\left(Q-F(\cdot, 1), \Omega_{\lambda_{1}, \lambda_{2}}\right) \neq \Theta \in \Gamma$.

By Lemma 4.4 we obtain that $(Q-F)^{-1}(0) \cap\left(\partial \Omega_{\lambda_{1} \lambda_{2}} \times[0,1]\right)=\emptyset$. Therefore by the homotopy property of the $S^{1}$-degree we obtain that

$$
\operatorname{Deg}\left(Q-F(\cdot, 1), \Omega_{\lambda_{1}, \lambda_{2}}\right)=\operatorname{Deg}\left(Q-F(\cdot, 0), \Omega_{\lambda_{1}, \lambda_{2}}\right)
$$

By Theorem 2.1, what is left is to show that $\operatorname{Deg}\left(Q-F(\cdot, 0), \Omega_{\lambda_{1}, \lambda_{2}}\right) \neq \Theta \in \Gamma$.
From Lemma 5.1 it follows that

$$
\operatorname{Deg}\left(Q-F(\cdot, 0), \Omega_{\lambda_{1}, \lambda_{2}}\right)=\operatorname{Deg}\left(Q-G(\cdot, 1), \Omega_{\lambda_{1}, \lambda_{2}}\right)
$$

Finally, by Lemma 5.2 we obtain $\operatorname{Deg}\left(Q-G(\cdot, 1), \Omega_{\lambda_{1}, \lambda_{2}}\right) \neq \Theta \in \Gamma$. Notice that we have just proved that $\operatorname{Deg}\left(Q-F(\cdot, 1), \Omega_{\lambda_{1}, \lambda_{2}}\right) \neq \Theta \in \Gamma$. The rest of the proof is a direct consequence of Theorem 2.1.

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