PERIODIC PERTURBATIONS WITH DELAY OF COUPLED DIFFERENTIAL EQUATIONS ON MANIFOLDS WITH APPLICATION TO A SUNFLOWER-LIKE EQUATION

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Abstract. We investigate the structure of the set of $T$-periodic solutions to periodically perturbed coupled delay differential equations on differentiable manifolds. By using fixed point index and degree-theoretic methods we prove the existence of branches of $T$-periodic solutions to the considered equations. As main application of our methods, we study a generalized version of the well-known sunflower equation, which is a second order delay differential equation used to model the helical movement of the tip of a sunflower plant (see, e.g., [20, 2]).

In the sense that the main results of those papers can be deduced from ours (see Remark 2 below). As a further motivation we will present an application of our methods to a generalized version of the well-known sunflower equation (cf. (13)), which is a second order delay differential equation used to model the helical movement of the tip of a sunflower plant (see, e.g., [20, 4]).

Let us describe more precisely our setting. Let $M \subseteq \mathbb{R}^k$ and $N \subseteq \mathbb{R}^s$ be boundaryless smooth manifolds, let $f: \mathbb{R} \times M \times N \times M \times N \to \mathbb{R}^k$ be tangent to $M$, and let $g: M \times N \to \mathbb{R}^s$ and $h: \mathbb{R} \times M \times N \times M \times N \to \mathbb{R}^s$ be tangent to $N$: This means that, for any $(t, p, q, v, w) \in \mathbb{R} \times M \times N \times M \times N$, then $g(p, q, v, w)$ and $h(t, p, q, v, w)$ belong to the tangent space $T_pN$, and $f(t, p, q, v, w)$ is in $T_pM$, respectively. Let also $a: \mathbb{R} \to \mathbb{R}$ be continuous. Given $T > 0$, we assume that $f$, $h$ and $a$ are $T$-periodic in the $t$ variable. Consider the following system of delay differential equations for $\lambda \geq 0$:

$$
\begin{align*}
\dot{x}(t) &= \lambda f(t, x(t), y(t), x(t-r), y(t-r)), \\
\dot{y}(t) &= a(t) g(x(t), y(t)) + \lambda h(t, x(t), y(t), x(t-r), y(t-r)),
\end{align*}
$$

where the time lag $r > 0$ is given. This system is equivalent to a single parameter-dependent delay differential equation on the product manifold $M \times N \subseteq \mathbb{R}^{k+s}$.

Denote by $C_T(M)$ and $C_T(N)$ the spaces of $T$-periodic continuous functions from $\mathbb{R}$ to $M$ and $N$, respectively, with the topology of uniform convergence. We investigate the properties of the set of the $T$-periodic triples (or briefly $T$-triples) of (1), i.e. of those triples $(\mu, x, y) \in [0, \infty) \times C_T(M) \times C_T(N)$, where $(x, y)$ is a solution to (1) when $\lambda = \mu$. In particular, we shall give conditions for the existence of a noncompact connected component of nontrivial $T$-triples (which we call a “branch”) emanating from the set $\nu^{-1}(0)$, where $\nu: M \times N \to \mathbb{R}^{k+s}$ is the vector field, tangent to $M \times N \subseteq \mathbb{R}^{k+s}$, given by

$$
\nu(p, q) = (w_1(p, q), g(p, q)),
$$

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with \( w[p, q] := \frac{1}{T} \int_0^T f(t, p, q, p, q) \, dt \). In the present setting, a \( T \)-triple \((\lambda, x, y)\) of (11) is said to be trivial if \( \lambda = 0 \) and \((x, y)\) is constant.

It is tempting to try to achieve the desired generalization of [21] by simply using a time-transformation as in [9][22] to get rid of the factor \( a(t) \) in (11) and then adapt the argument of [3][12] to the present case. Nevertheless, this simple procedure does not work because the transformed perturbing term would result in a form inappropriate for our methods. In fact, the time-transformation used in [22] does not preserve the fixed-delay structure. Instead, to prove our result, we follow the lead of [7] and combine the techniques of [22] and [21].

As main application of our methods we study the set of periodic solutions of a parametrized second order delayed differential equation (the one called sunflower-like equation). This parametrized equation is derived in Section 4 in a fairly direct way starting from the sunflower equation, and reads as follows:

\[
\dot{y}(t) = a(t)\dot{y}(t) + \lambda\phi(y(t), y(t-r)), \quad \lambda \geq 0,
\]

where \( a: \mathbb{R} \to \mathbb{R} \) and \( \phi: \mathbb{R}^2 \to \mathbb{R} \) are continuous and \( a \) is \( T \)-periodic with \( \phi \neq 0 \). By using elementary transformations we will show how the above equation can be equivalently rewritten as systems of delay differential equations of type (11) to which our methods apply. Thus, we find in a natural way a result about the structure of the set of \( T \)-periodic solutions of the sunflower-like equation above.

2. **Poincaré–type translation operator**

Consider the system of equations (11). We are interested in its \( T \)-periodic solutions. Without loss of generality, as suggested in [3], we will assume that \( T \geq r \). In fact, for \( n \in \mathbb{N} \), the system (11) and

\[
\begin{align*}
\dot{x}(t) &= \lambda f(t, x(t), y(t), x(t-r-nT), y(t-(r-nT))) \\
\dot{y}(t) &= a(t)g(x(t), y(t)) + \lambda h(t, x(t), y(t), x(t-(r-nT)), y(t-(r-nT)))
\end{align*}
\]

have the same \( T \)-periodic solutions. Thus, if necessary, one can replace \( r \) with \( r-nT \), where \( n \in \mathbb{N} \) is such that \( 0 < r-nT \leq T \).

Let us now introduce some notation. Given any \( X \subseteq \mathbb{R}^k \), \( \tilde{X} \) denotes the metric space \( C([-r, 0], X) \) with the distance inherited from the Banach space \( \tilde{R}^k = C([-r, 0], \mathbb{R}^k) \) with the usual supremum norm.

Given any \((p, q) \in M \times N\), denote by \( p^\# \in \tilde{M} \) and \( q^\# \in \tilde{N} \) the constant functions \( p^\#(t) \equiv p \) and \( q^\#(t) \equiv q \), \( t \in [-r, 0] \), respectively. Thus, \((p^\#, q^\#) \in \tilde{M} \times \tilde{N} \cong \tilde{M} \times \tilde{N} \). For any \( U \subseteq M \times N \), define \( U^\# = \{(p^\#, q^\#) \in \tilde{M} \times \tilde{N} : (p, q) \in U\} \). Also, given \( W \subseteq \tilde{M} \times \tilde{N} \), we put \( W^\# = \{(p, q) \in M \times N : (p^\#, q^\#) \in W\} \). Finally, we will denote by \( C_T(X) \) the metric subspace of the Banach space \((C_T(\mathbb{R}^k), \| \cdot \|)\) of all the \( T \)-periodic continuous maps \( x: \mathbb{R} \to X \) (as above, with the usual \( C^0 \) norm). Observe that \( C_T(X) \) is complete if and only if \( X \) is complete (or, equivalently, closed as a subset of \( \mathbb{R}^k \)). Nevertheless, since \( M \) and \( N \) are locally compact, \( C_T(M \times N) \simeq C_T(M) \times C_T(N) \) is always locally complete.

Assume now, unless differently stated, that \( a, f, g \) and \( h \) are \( C^1 \). Consider the map \( H \) with domain \( \mathcal{D}_H \subseteq \mathbb{R} \times \tilde{M} \times \tilde{N} \times \mathbb{R} \) in \( \tilde{M} \times \tilde{N} \) defined by

\[
H(\lambda, \varphi, \psi, \mu)(\theta) = \left( x_{\lambda, \mu}(\varphi, \psi, T + \theta), y_{\lambda, \mu}(\varphi, \psi, T + \theta) \right), \quad \theta \in [-r, 0],
\]
where \( t \mapsto (x_{\lambda,\mu}(\varphi, \psi, t), y_{\lambda,\mu}(\varphi, \psi, t)) \) denotes the unique maximal solution of the initial-value problem

\[
\begin{aligned}
\dot{x}(t) &= \lambda \mu f(t, x(t), y(t), x(t-r), y(t-r)) + (1-\mu) \frac{a(t)}{\delta} v(t), \\
\dot{y}(t) &= a(t) g(x(t), y(t)) + \lambda \mu b(t, x(t), y(t), x(t-r), y(t-r)), \\
x(t) &= \varphi(t), \quad y(t) = \psi(t),
\end{aligned}
\]

for \( t > 0 \),

\( t \in [-r, 0]. \)

Well known properties of differential equations imply that \( \mathcal{D}_H \) is an open subset of \( \mathbb{R} \times \tilde{M} \times \tilde{N} \). A similar argument shows that the set \( \mathcal{D}' := \{ (\varphi, \psi) \in \tilde{M} \times \tilde{N} : (0, \varphi, \psi, 1) \in \mathcal{D}_H \} \) is open as well. Also, since we are assuming \( T \geq r \) (see above), the Theorem of Ascoli-Arzela implies that \( H \) is a locally compact map (compare, e.g., [18] or [5]).

**Remark 1.** Consider the following equation:

\[
\begin{aligned}
\dot{x}(t) &= 0, \\
\dot{y}(t) &= a(t) g(x(t), y(t)),
\end{aligned}
\]

Given \( V \subseteq \tilde{M} \times \tilde{N} \) such that \( \overline{V} \subseteq \mathcal{D}' \) we have that all solutions of (3) starting at time \( t = 0 \) from \( V_{\#} \) are defined (at least) for \( t \in [0, T] \). An argument similar to, e.g., [22, Remark 2.3] or [7, Remark 2.1] shows that the same assertion holds for (3) when \( a(t) \) is replaced with its average \( \bar{a} = \frac{\int_{[0, T]} a(t) \, dt}{T} \):

\[
\begin{aligned}
\dot{x}(t) &= 0, \\
\dot{y}(t) &= \bar{a}(t) g(x(t), y(t)),
\end{aligned}
\]

In fact, one could prove that solutions of (3) and of (4), leaving at time \( t = 0 \) from the same point, coincide at time \( t = T \). Thus, \( T \)-periodic orbits (images of solutions) of (3) and (4) must coincide. More precisely, let \( \{ \Phi_t \}_{t \in \mathbb{R}} \) be the local flow associated to (4). That is, \( \Phi: U \to \tilde{M} \times \tilde{N} \) is defined on an open subset \( U \) of \( \tilde{M} \times \tilde{N} \), containing \( \{ 0 \} \times \tilde{M} \times \tilde{N} \), with the property that for any \( (p, q) \in \tilde{M} \times \tilde{N} \) the curve \( t \mapsto \Phi_t(p, q) \) is the maximal solution of (4) given the initial condition \( \Phi_0(p, q) = (p, q) \). Then, given \( \tau \in \mathbb{R} \), the domain of \( \Phi_{\tau} \) is the open set consisting of those points \( (p, q) \in \tilde{M} \times \tilde{N} \) for which the maximal solution of (4) starting from \( (p, q) \) at \( t = 0 \) is defined up to \( \tau \). (We are interested, in particular, to the case \( \tau = T \).) Let \( \{ \Psi_t \}_{t \in \mathbb{R}} \) be the analogous local flow associated to (3). The argument of the above cited remarks show that \( \Psi_{\tau}(p, q) = \Phi_{\tau}(p, q) \) whenever this relation makes sense, in particular for all \( (p, q) \in V_{\#} \).

The following definition is convenient:

**Definition 1.** We say that \( V \subseteq \tilde{M} \times \tilde{N} \) has the constant periodic property for (3) if any \( T \)-periodic solution \((x, y)\) of Equation (3) that intersects \( \partial V_{\#} \) is constant.

We have the following result:

**Lemma 1.** Let \( V \subseteq \tilde{M} \times \tilde{N} \) be open and such that

\[
Z_V := \{ (p\#, q\#) \in V : \nu(p, q) = 0 \}
\]

is compact. Then, there exists an open neighborhood \( W \subseteq V \) of \( Z_V \) and \( \eps > 0 \) s.t. \([0, \eps] \times W \times [0, 1] \subseteq \mathcal{D}_H \) and \( H([0, \eps] \times W \times [0, 1]) \) is compact.

Assume in addition that \( V_{\#} \) is relatively compact, \( \overline{V} \subseteq \mathcal{D}' \) and that \( V \) has the constant periodic property for (3) (Definition 3). Then \( W \) can be taken in such a way that it has the constant periodic property as well. That is, if \((x, y)\) is a \( T \)-periodic solution of (3) intersecting \( \partial W_{\#} \), then \((x, y)\) is constant.
Proof. One immediately checks that the set $Z_V$ consists of $T$-periodic solutions of (3). Thus, we have that $Z_V \subseteq \mathcal{D}'$ and the first part of the lemma follows from the local compactness of $H$.

Let us now prove the second part of the assertion. Let $\{\Phi_t\}_{t \in \mathbb{R}}$ be the local flow associated to (3) as in Remark 1. The map $(t,p,q) \mapsto \Phi_t(p,q)$ is continuous and, therefore, the “attainable set” $\mathcal{A}_T := \Phi_{[0,T]}(\mathcal{V}_{\#})$ is compact. Thus, the union $\mathcal{O}_T$ of all $T$-periodic orbits of (3) starting from points of $\mathcal{V}_{\#}$, being closed in $\mathcal{A}_T$, is compact as well. Clearly, since $\mathcal{O}_T$ is autonomous, $\mathcal{O}_T$ is actually the set of all $T$-periodic orbits of (3) that intersect $\mathcal{V}_{\#}$.

Remark 1 shows that $\mathcal{O}_T$ consists indeed of all $T$-periodic orbits of (3) that intersect $\mathcal{V}_{\#}$. Let us denote by $K$ the union of $Z_V$ with this set. Clearly $K$ is contained in $\mathcal{D}'$. The local compactness of $H$ implies the existence of an open neighborhood $W \subseteq V$ of $K$ and a positive $\varepsilon$ with the property that $[0,\varepsilon] \times \overline{W} \times [0,1] \subseteq \mathcal{D}_H$ and $H([0,\varepsilon] \times \overline{W} \times [0,1])$ is compact. The second part of the claim follows now from the fact a $T$-periodic solution of (3) whose image intersects the boundary $\partial W_{\#}$, of the set $W$ just constructed, necessarily intersects $\partial V_{\#}$ and thus must be constant.

It is convenient to set

$$Q_T^{\lambda} = H(\lambda,\cdot,\cdot,1), \quad \text{and} \quad \tilde{Q}_T^{\lambda} = H(\lambda,\cdot,\cdot,0).$$

We will denote the domain of $H(\cdot,\cdot,\cdot,1)$ by the letter $\mathcal{D}$. The following is the main result of this section (cf. [12, 17]). It relates the fixed point index of $Q_T^{\lambda}$ for small $\lambda > 0$ (see, e.g., [15, 17] for an introduction) with the degree of the tangent vector field $\nu$. Recall that this notion, roughly speaking, counts (algebraically) the zeros of a vector field; for an exposition of this topic we refer, e.g., to [15] or [17].

**Theorem 1.** Given $V \subseteq \overline{M} \times \overline{N}$ open and such that

(i) $V_{\#}$ is relatively compact;

(ii) There exists $s > 0$ such that $[0,s] \times \overline{V} \subseteq \mathcal{D}$;

(iii) $Z_V$ is compact;

(iv) If $(x,y)$ is a $T$-periodic solution of (3) whose image intersects $\partial V_{\#}$, then $(x,y)$ is constant.

Then there exists $\lambda_{*} \in (0,s]$ such that, for $\lambda \in (0,\lambda_{*})$, $\text{ind}(Q_T^{\lambda},V)$ is well defined and

$$\text{ind}(Q_T^{\lambda},V) = \text{sign}(\#)^{\dim N} \text{deg} \left(-\nu, V_{\#}\right).$$

The symbol “$\text{ind}(Q_T^{\lambda},V)$” in the above formula denotes the fixed point index of $Q_T^{\lambda}$ in the open set $V$, whereas “$\text{deg} \left(-\nu, V_{\#}\right)$” denotes the degree of the tangent vector field $-\nu$ in the open subset $V_{\#}$ of $M \times N$.

**Proof of Theorem 1**. Let $W$ and $\varepsilon$ be as in Lemma 1. Consider the sets

$$S = \{ (\lambda,\varphi,\psi) \in [0,\varepsilon] \times \overline{W} : H(\lambda,\varphi,\psi,1) = (\varphi,\psi) \},$$

$$S_0 = S \cap \{ (0) \times \overline{M} \times \overline{N} \}.$$

Clearly, $S$ is compact being a closed subset of the compact set $[0,\varepsilon] \times H([0,\varepsilon] \times \overline{W} \times [0,1])$. Thus $S_0$ is compact as well. Using the definition of $Q_T^{\lambda}$, we will prove the following fact:

**Claim 1.** There exists $\lambda_0 \in (0,\min\{\varepsilon,s\}]$ such that if $(\varphi,\psi) \in V$ is a fixed point of $Q_T^{\lambda}$ with $\lambda \in (0,\lambda_0]$ then $(\varphi,\psi) \in W$. That is, $Q_T^{\lambda}$ has no fixed points in $\overline{V} \setminus W$ for $\lambda \in (0,\lambda_0]$. 

To prove this claim we proceed by contradiction. If the claim is false there exist sequences \( \{\lambda_n\} \subseteq (0, \lambda_0] \), and \( \{(\varphi_n, \psi_n)\} \subseteq \overline{V} \setminus W \), with \( \lambda_n \to 0 \) and \( (\lambda_n, \varphi_n, \psi_n) \in \mathcal{S} \). By the compactness of \( S_0 \cap (\overline{V} \setminus W) \) we can assume that \( (\varphi_n, \psi_n) \to (\varphi_0, \psi_0) \in S_0 \cap (\overline{V} \setminus W) \). The continuous dependence on data shows that the solution of (3) with initial data \((\varphi_0, \psi_0)\) is \( T \)-periodic. Assumption (iv) shows that there exists \( p_0 \in M \) and \( q_0 \in N \) such that \((\varphi_0, \psi_0) = (p_0^#, q_0^#)\). Clearly, one has \( g(p_0, q_0) = 0 \).

Let \((x_n, y_n)\) be the unique maximal solution of

\[
\begin{aligned}
\dot{x}(t) &= \lambda_n f(t, x(t), y(t), x(t-r), y(t-r)), \\
\dot{y}(t) &= a(t) f(x(t), y(t)) + \lambda_n b(t, x(t), y(t), x(t-r), y(t-r)), \\
x(t) &= \varphi_n(t), \\
y(t) &= \psi_n(t),
\end{aligned}
\]

for \( t > 0 \), and \( t \in [-r, 0] \).

Then,

\[
0 = x_n(T) - x_n(0) = \lambda_n \int_0^T f(t, x_n(t), y_n(t), x_n(t-r), y_n(t-r)) \, dt.
\]

So that, in particular,

\[
0 = \int_0^T f(t, x_n(t), y_n(t), x_n(t-r), y_n(t-r)) \, dt
\]

and, passing to the limit, we get

\[
0 = \int_0^T f(t, p_0, q_0, p_0, q_0) \, dt = w_\lambda(p_0, q_0).
\]

Hence, \( \nu(p_0, q_0) = (w_\lambda(p_0, q_0), g(p_0, q_0)) = 0 \). This contradicts the choice of \( W \) and completes the proof of Claim 1.

Claim 1 shows that, for \( \lambda \in (0, \lambda_0] \), the set of the fixed points of \( Q^T_{\lambda} \) that lie in \( V \) is, in fact, contained in \( W \). Hence, by the compactness of \( \mathcal{S} \), it is compact. As a consequence, \( \text{ind}(Q^T_{\lambda}, V) \) and \( \text{ind}(Q^T_{\lambda}, W) \) are well-defined and, by the excision property,

\[
\text{ind}(Q^T_{\lambda}, V) = \text{ind}(Q^T_{\lambda}, W), \quad \text{for } \lambda \in (0, \lambda_0].
\]

In fact, when \( \lambda \) is sufficiently small, something more can be obtained:

**Claim 2.** There exists \( \lambda_* \in (0, \lambda_0] \), such that the homotopy \( H_\lambda : \overline{V} \times [0, 1] \to M \times N \) given by \( H_\lambda(\varphi, \psi, \mu) = H(\lambda, \varphi, \psi, \mu) \), is admissible for each \( \lambda \in (0, \lambda_0] \).

To prove the claim we ought to show that for each \( \lambda \in (0, \lambda_*] \), \( \lambda_* > 0 \) sufficiently small, the set of fixed points

\[
\mathcal{F}_\lambda = \{(\varphi, \psi) \in \overline{V} : H(\lambda, \varphi, \psi, \mu) = (\varphi, \psi), \text{ for some } \mu \in [0, 1]\},
\]

which is compact being a closed subset of \( H([0, \varepsilon] \times \overline{V} \times [0, 1]) \), is contained in \( W \). Suppose by contradiction that this is not the case, that is, that such a choice of \( \lambda_* \) cannot be done. Then there are sequences \( \{\lambda_n\} \subseteq (0, \lambda_0] \), \( \{\mu_n\} \subseteq [0, 1] \) and \( \{(\varphi_n, \psi_n)\} \subseteq \partial W \) with \( \lambda_n \to 0 \) and

\[
H(\lambda_n, \varphi_n, \psi_n, \mu_n) = (\varphi_n, \psi_n).
\]

As in the proof of Claim 1, by the compactness of \( H([0, \varepsilon] \times \overline{V} \times [0, 1]) \) we can assume that \((\varphi_n(0), \psi_n(0)) \to (\varphi_0, \psi_0) \in \partial W \). The continuous dependence on data shows that the solution of (3) with initial data \((\varphi_0, \psi_0)\) is \( T \)-periodic. Assumption (iv) shows that there exists \( p_0 \in M \) and \( q_0 \in N \) such that \((\varphi_0, \psi_0) = (p_0^#, q_0^#)\). Clearly,
we get \( g(p_0, q_0) = 0 \). From (3) it follows that if \((x_n, y_n)\) is the solution of
\[
\begin{aligned}
\dot{x}(t) &= \lambda_n \left[ \mu_n f(t, x(t), y(t), x(t-r), y(t-r)) \\
&\quad + (1 - \mu_n) \frac{a(t)}{\hat{q}} w_1(x(t), y(t)) \right], \\
\dot{y}(t) &= a(t) g(x(t), y(t)) \\
x(t) &= \varphi_n(t), \quad y(t) = \psi_n(t), \quad t \geq 0,
\end{aligned}
\]

Then,
\[
0 = x_n(T) - x_n(0) = \lambda_n \int_0^T \mu_n f(t, x_n(t), y_n(t), x_n(t-r), y_n(t-r)) \, \text{d}t \\
&\quad + \lambda_n \int_0^T (1 - \mu_n) \frac{a(t)}{\hat{q}} w_1(x_n(t), y_n(t)) \, \text{d}t.
\]

So that
\[
0 = \mu_n \int_0^T f(t, x_n(t), y_n(t), x_n(t-r), y_n(t-r)) \, \text{d}t \\
&\quad + (1 - \mu_n) \int_0^T \frac{a(t)}{\hat{q}} w_1(x_n(t), y_n(t)) \, \text{d}t.
\]

Passing to the limit we get
\[
0 = \mu_0 \int_0^T f(t, p_0, q_0, p_0, q_0) \, \text{d}t + (1 - \mu_0) \int_0^T \frac{a(t)}{\hat{q}} w_1(p_0, q_0) \, \text{d}t = \mu_0 w_1(p_0, q_0) + (1 - \mu_0) w_1(p_0, q_0) = w_1(p_0, q_0),
\]

which contradicts the choice of \( W \) and proves Claim 2.

Claim 2, along with the homotopy invariance property, imply that for \( \lambda \in (0, \lambda_*] \)
\[
\text{ind}(Q_{\lambda}, W) = \text{ind}(H(\lambda, \cdot, \cdot, 1), W) = \text{ind}(H(\lambda, \cdot, \cdot, 0), W) = \text{ind}(Q_{\lambda}, W).
\]

Consider the tangent vector field \( v_{\lambda} \) on \( M \times N \) given by
\[
v_{\lambda}(p, q) := \left( \frac{\lambda}{\hat{q}} w_1(p, q), \lambda g(p, q) \right).
\]

Theorem 3.2 of [21] imply that, for each fixed \( \lambda \in (0, \lambda_*] \)
\[
\text{ind}(Q_{\lambda}, W) = \text{sign}(\hat{q})^{\dim(M \times N)} \deg(-v_{\lambda}, W_{\#}).
\]

Since \( \lambda > 0 \), a well known property of the degree yields
\[
\text{deg}(-v_{\lambda}, W_{\#}) = \text{deg}(-v_1, W_{\#}).
\]

Lemma 1 of [21] shows that
\[
\text{deg}(-v_1, W_{\#}) = \text{sign}(\hat{q})^{\dim M} \deg(-\nu, W_{\#}),
\]

hence, by equalities (7) – (10), taking into account that \( \dim(M \times N) = \dim M + \dim N \), we get
\[
\text{ind}(Q_{\lambda}, W) = \text{ind}(Q_{\lambda}, W) = \text{sign}(\hat{q})^{\dim(M \times N)} \deg(-v_{\lambda}, W_{\#})
\]
\[
= \text{sign}(\hat{q})^{\dim(M \times N)} \deg(-v_1, W_{\#})
\]
\[
= \text{sign}(\hat{q})^{\dim(M \times N)} \text{sign}(\hat{q})^{\dim(M)} \deg(-\nu, W_{\#})
\]
\[
= \text{sign}(\hat{q})^{2 \dim M + \dim N} \deg(-\nu, W_{\#})
\]
\[
= \text{sign}(\hat{q})^{\dim N} \deg(-\nu, W_{\#}).
\]
Finally, by \((5), (11)\) and the excision property of the degree, we get
\[
\operatorname{ind}(Q^1_T, V) = \operatorname{ind}(Q^1_T, W) \\
= \text{sign}(\hat{g})^{\dim N} \deg(-\nu, W) \\
= \text{sign}(\hat{g})^{\dim N} \deg(-\nu, V),
\]
which proves the assertion. \(\square\)

3. Branches of \(T\)-periodic solutions

Let \(T > 0\) be given, by \(C_T(\mathbb{R}^d)\) we mean the Banach space of all the continuous \(T\)-periodic functions \(\zeta: \mathbb{R} \to \mathbb{R}^d\) whereas \(C_T(X)\) denotes the metric subspace of \(C_T(\mathbb{R}^d)\) consisting of all those \(\zeta \in C_T(\mathbb{R}^d)\) that take values in \(X\). It is not difficult to prove that \(C_T(X)\) is complete if and only if \(X\) is closed in \(\mathbb{R}^d\).

It is also convenient to introduce the following notation: Given \((p, q)\) in \(M \times N\), let \((\tilde{p}, \tilde{q}) \in C_T(M \times N) = C_T(\mathbb{R}, M \times N)\) be the constant maps \((\tilde{p}(t), \tilde{q}(t)) \equiv (p, q), t \in \mathbb{R}\).

We are now in a position to state our result concerning the “branches” of \(T\)-triples of \((1)\). Its proof follows closely the one of Th. 5.1 in \([12]\) (see also \([10, 7]\)), for this reason we only provide a sketch for the sake of completeness.

**Theorem 2.** Let \(\Omega\) be an open subset of \([0, \infty) \times C_T(M \times N)\), and let \(\Omega_{M \times N} := \{(p, q) \in \Omega : (0, \tilde{p}, \tilde{q}) \in \Omega\}\). Assume that \(\deg(\nu, \Omega_{M \times N})\) is well-defined and nonzero. Then there exists a connected set \(\Gamma\) of nontrivial \(T\)-triples for \((1)\) in \(\Omega\) whose closure in \([0, \infty) \times C_T(M \times N)\) meets \(\nu^{-1}(0) \cap \Omega_{M \times N}\) and is not contained in any compact subset of \(\Omega\). In particular, if \(M \times N\) is closed in \(\mathbb{R}^{k+s}\) and \(\Omega = [0, \infty) \times C_T(M \times N)\), then \(\Gamma\) is unbounded.

The proof of this theorem is based on the following global connection result (see \([10]\)), which will also be needed later.

**Lemma 2.** Let \(Y\) be a locally compact metric space and let \(Z\) be a compact subset of \(Y\). Assume that any compact subset of \(Y\) containing \(Z\) has nonempty boundary. Then \(Y \setminus Z\) contains a connected set whose closure (in \(Y\)) intersects \(Z\) and is not compact.

We are now ready to sketch the proof of the theorem.

**Sketch of the proof of Theorem 3** This proof can be roughly divided into three steps:

**Step 1.** We assume first that the maps \(a, \bar{f}, \bar{g}\) and \(\bar{h}\) are \(C^1\), so that uniqueness of solutions holds for \((1)\). Consider the following notion:

A triple \((\lambda, \varphi, \psi) \in [0, \infty) \times \tilde{M} \times \tilde{N}\) is said to be a starting triple for \((1)\) if the following initial value problem has a \(T\)-periodic solution:

\[
\begin{cases}
\dot{x}(t) = \lambda f(t, x(t), y(t)), \\
\dot{y}(t) = a(t)g(x(t), y(t)) + \lambda h(t, x(t), y(t)), \\
x(t) = \varphi(t), \\
y(t) = \psi(t),
\end{cases} \quad t > 0,
\]

\[
\begin{cases}
\dot{x}(t) = \varphi(t), \\
\dot{y}(t) = \psi(t),
\end{cases} \quad t \in [-r, 0],
\]

A triple of the type \((0, p^\#, q^\#)\) with \(g(p, q) = 0\) is clearly a starting triple and will be called a trivial starting triple. The set of all starting triples for \((1)\) will be denoted by \(S\). By known continuous dependence properties of delay differential equations the set \(\mathcal{V} \subseteq [0, \infty) \times \tilde{M} \times \tilde{N}\) of all triples \((\lambda, \varphi, \psi)\) such that the unique solution of \((12)\) is defined at least up to \(T\) is open (compare it to the set \(\mathcal{D}\) defined in section 2). Clearly \(\mathcal{V}\) contains the set \(S\) of all starting pairs for \((1)\).
Given an open set $W$ of $[0, \infty) \times \tilde{M} \times \tilde{N}$, let

$$W^0_\# := \left( W \cap \{0\} \times \tilde{M} \times \tilde{N} \right)_\# = \left\{ (p, q) \in M \times N : (0, p^#, q^#) \in W \right\}.$$  

Our first step consists of proving that, if $\deg (\nu, W^0_\#)$ is well-defined and nonzero, then there exists in $S \cap W$ a connected set $G$ of nontrivial starting triples whose closure in $S \cap W$ meets $\{ (0, p^#, q^#) \in W : g(p, q) = 0 \}$ and is not compact.

The proof of this fact follows closely the one of Proposition 4.1 in [12] using Theorem II in place of [12] Th. 3.2. Loosely speaking, this proof uses the properties of the fixed point index and of the degree of a tangent vector field to obtain a contradiction with Lemma 2. (Compare also [7, Th. 4.1].)

**Step 2.** As in Step 1 we assume that the maps $a$, $f$, $g$ and $h$ are $C^1$. Denote by $X$ the set of $T$-periodic triples of (1) and by $S$ the set of starting triples of the same equation, as above. Define the map $\Pi : X \to S$ by

$$\Pi(\lambda, x, y) = (\lambda, x|_{[-r,0]}, y|_{[-r,0]})$$

and observe that $\Pi$ is continuous, onto and, since $f$, $g$ and $h$ are smooth, it is also one to one. Furthermore, by the continuous dependence on data, $\Pi^{-1} : S \to X$ is continuous as well. Take

$$S_\Omega = \left\{ (\lambda, \varphi, \psi) \in S : \text{the solution of (12) is contained in } \Omega \right\},$$

so that $X \cap \Omega$ and $S_\Omega$ correspond under the homeomorphism $\Pi : X \to S$. Thus, $S_\Omega$ is an open subset of $S$ and, consequently, we can find an open subset $W$ of $[0, \infty) \times \tilde{M} \times \tilde{N}$ such that $S \cap W = S_\Omega$. This implies, as in [12] Th. 5.1, that

$$\{ (p, q) \in W^0_\# : g(p, q) = 0 \} = \{ (p, q) \in \Omega \cap N : g(p, q) = 0 \}.$$

The excision property of the degree of tangent vector fields yields

$$\deg (g, W^0_\#) = \deg (g, \Omega \cap N) \neq 0.$$

By Step 1 we deduce the existence of a connected set

$$\Sigma \subseteq (S \cap W) \setminus \{ (0, p^#, q^#) \in W : g(p, q) = 0 \}$$

whose closure in $S \cap W$ meets $\{ (0, p^#, q^#) \in W : g(p, q) = 0 \}$ and is not compact. Clearly, $\Gamma = \Pi^{-1}(\Sigma)$ satisfies the assertion.

**Step 3.** We now only need to remove the $C^1$-regularity assumption on the maps $a$, $f$, $g$ and $h$ replacing it with continuity. This is done by an approximation procedure that follows closely the one used in [12] Th. 5.1. For this reason we skip the details. \qed

**Remark 2.** One can easily check that Theorem II implies both [3] Lemma 4.5 and [4] Lemma 3.5, albeit in the less general case of boundaryless manifolds, which are valid for a single differential equation of the form

$$\dot{x}(t) = \lambda f(t, x(t), x(t-r)),$$

where $f : \mathbb{R} \times M \times M \to \mathbb{R}^k$ is tangent to $M$. At the same time, Theorem II extends [12] Thm. 5.1 (see also Th. 4.1 in [4]) that applies to the differential equations of the following type:

$$\dot{y}(t) = a(t)g(y) + \lambda h(t, y(t), y(t-r)),$$

where $g : M \to \mathbb{R}^k$ and $h : \mathbb{R} \times M \times M \to \mathbb{R}^k$ are tangent to $M$.\[\]
4. An application: A sunflower-like equation

In recent years there has been a growing interest around the dynamical behavior of delay differential equations and in their possible use in modeling biological and ecological systems. In particular, a certain amount of research has been dedicated to the sunflower equation, i.e.

\[ \ddot{y}(t) = -\frac{\alpha}{r} \dot{y}(t) - \frac{\beta}{r} \sin (y(t - r)), \]

where \( \alpha \) and \( \beta \) are experimental parameters and \( r > 0 \) is the finite time delay (see [13]). Somolinos in [20] showed the existence of periodic solutions to (13) for a certain range of values for the involved parameters \( \alpha, \beta, \) and \( r. \) This existence result covers both the cases of small and large amplitude limit cycles generated by Hopf bifurcation. More recently, Liu and Kalmár-Nagy [14] computed limit cycle amplitudes and frequencies for (13). Other meaningful results related to this equation can be found in [2, 8, 10].

Here, we are concerned with a parametrized differential equation which arise quite naturally from (13): It is obtained by assuming that the coefficient \(-\alpha/r\) of \( \ddot{y}(t) \) is actually a real valued function \( a: \mathbb{R} \to \mathbb{R} \) and setting \( \lambda = -\beta/r \) as the coefficient of the second term in the left-hand side of (13). The parametrized equation under consideration reads as follows:

\[ \ddot{y}(t) = a(t)\dot{y}(t) + \lambda \phi(y(t), y(t - r)), \quad \lambda \geq 0, \]

where \( a \) and \( \phi \) are continuous, \( a \) is \( T \)-periodic with average \( \bar{a} \neq 0. \) We point out that the assumption \( \bar{a} \neq 0 \) serves to generalize the constant coefficient of \( \ddot{y}(t) \) in (13).

Ideally, \( a(t) \) can be can be thought as a perturbation of the constant term \(-\alpha/r\), namely \( a(t) = -\alpha/r + \epsilon(t) \) where \( \epsilon(t) \) is continuous, \( T \)-periodic and sufficiently small so that \( \bar{a} = -\alpha/r + \bar{\epsilon} \neq 0. \)

We wish to look at how the results of the previous section apply to (14), especially in the case in which \( \phi(y(t), y(t - r)) = \sin (y(t - r)) \).

Let us now recall the notion of \( T \)-periodic pair (or \( T \)-pair for brevity) for Equation (14) and some related facts.

A pair \((\lambda, y) \in [0, \infty) \times C^1(\mathbb{R})\) is called a \( T \)-pair for (14) if \( y \) is a \( T \)-periodic solution of (14). A \( T \)-pair \((\lambda, y) \) is trivial if \( y \) is constant and \( \lambda = 0. \)

In order to study Equation (14), we introduce a transformation that allows us to rewrite this model in an equivalent but easier to handle form. We need the following technical lemma whose proof is a standard ODE argument which we provide for the sake of completeness.

**Lemma 3.** Let \( a: \mathbb{R} \to \mathbb{R} \) be as in (14) and such that \( \bar{a} \neq 0. \) Then, there exists a unique \( T \)-periodic \( C^1 \) function \( \sigma: \mathbb{R} \to \mathbb{R} \) for which

\[ a(t) = \frac{\dot{\sigma}(t)}{\sigma(t)} - \sigma(t), \quad \text{for all } t \in \mathbb{R}. \]

Clearly, \( \sigma \) has constant sign so that, in particular \( \bar{\sigma} \neq 0. \)

As a direct consequence, we have that Equation (14) can be rewritten as

\[ \ddot{y}(t) = \left( \frac{\dot{\sigma}(t)}{\sigma(t)} - \sigma(t) \right) \dot{y}(t) + \lambda \phi(y(t), y(t - r)), \quad \lambda \geq 0, \]

with \( \sigma \) chosen as in Lemma 3.

**Proof of Lemma 3.** It is easy to verify by inspection that, for any \( c \in \mathbb{R}, \)

\[ \zeta(t) := e^{-\int_0^t \bar{a}(s)ds} \left( c - \int_0^t e^{\int_0^s \bar{a}(r)dr} ds \right), \]
is a solution of equation

\[ (17) \quad \dot{\zeta}(t) = -\zeta(t)\alpha(t) - 1, \]

which corresponds to (15) under the transformation \( \zeta(t) = 1/\sigma(t) \). Clearly, \( \zeta(0) = c \). Taking

\[ c = c_0 := \frac{\int_0^T e^{\int_0^t \alpha(s)\,ds} \, ds}{e^{-T\dot{\vartheta}}} - 1, \]

(recall that \( \dot{\vartheta} \neq 0 \)) we get \( \zeta(0) = \zeta(T) \). Since the right-hand-side of (17) is \( T \)-periodic we obtain that \( \zeta \) is \( T \)-periodic as well. In fact, the above is the only choice of \( c \) for which \( \zeta(0) = \zeta(T) \); thus (17) has a unique \( T \)-periodic solution.

We need to prove that the function \( t \mapsto 1/\zeta(t) \) is a \( T \)-periodic solution of (15).

It is sufficient to show that \( \zeta(t) \neq 0 \) for all \( t \in \mathbb{R} \). We consider the two possibilities \( \dot{\vartheta} > 0 \) and \( \dot{\vartheta} < 0 \) separately:

**Case \( \dot{\vartheta} > 0 \).** Clearly, \( e^{-T\dot{\vartheta}} < 1 \) so that, since \( e^{-T\dot{\vartheta}} \int_0^T e^{\int_0^t \alpha(s)\,ds} \, ds > 0 \) we have \( c < 0 \). Now, being \( e^{-\int_0^t \alpha(s)\,ds} > 0 \) for all \( t \in \mathbb{R} \), we get \( \zeta(t) < 0 \) for all \( t \).

**Case \( \dot{\vartheta} < 0 \).** In this case one has \( e^{-T\dot{\vartheta}} > 1 \), thus

\[ 1 < \frac{e^{-T\dot{\vartheta}}}{e^{-T\dot{\vartheta}} - 1}. \]

Since \( t \mapsto e^{-\int_0^t \alpha(s)\,ds} \) is a positive function we have:

\[ \int_0^T e^{\int_0^t \alpha(s)\,ds} \, ds < \int_0^T e^{\int_0^t \alpha(s)\,ds} \, ds < \frac{e^{-T\dot{\vartheta}}}{e^{T\dot{\vartheta}} - 1} \int_0^T e^{\int_0^t \alpha(s)\,ds} \, ds, \]

so that \( \zeta(t) > 0 \).

Thus, in both cases, we find a \( T \)-periodic solution of (15). The uniqueness, follows from the fact that if \( t \mapsto \sigma(t) \) is a \( T \)-periodic solution of (15), hence defined for all \( t \in \mathbb{R} \), then \( \sigma(t) \neq 0 \). Then, \( t \mapsto 1/\sigma(t) \) is a \( T \)-periodic solution of (17) which, as discussed above, is unique. \( \square \)

**Remark 3.** From the proof of Lemma 3 it follows that \( \sigma \) has (constantly) the opposite sign of the average \( \dot{\vartheta} \). This fact has the obvious consequence that the signs of the averages of \( \sigma \) and of \( 1/\sigma \) coincide with \( -\text{sign}(\dot{\vartheta}) \). For the average \( \dot{\vartheta} \) of \( \sigma \) something more precise can be be deduced by the following simple argument:

\[ \dot{\vartheta} = \frac{1}{T} \int_0^T \alpha(t) \, dt = \frac{1}{T} \int_0^T \frac{\dot{\sigma}(t)}{\sigma(t)} \, dt - \frac{1}{T} \int_0^T \sigma(t) \, dt = -\dot{\vartheta}. \]

In fact,

\[ \int_0^T \frac{\dot{\sigma}(t)}{\sigma(t)} \, dt = \ln \left( \sigma(T) \right) - \ln \left( \sigma(0) \right) = 0 \]

because of the \( T \)-periodicity of \( \sigma \).

To investigate Equation (14) or, equivalently, (16) we follow the approach used in [21] §5. Along this path we find convenient to treat a more general class of equations, i.e.

\[ (18) \quad \dot{y}(t) = \left( \frac{\gamma(t)}{\gamma(t)} - \gamma(t) g(y(t)) \right) \dot{\gamma}(t) + \lambda f(t, y(t), g(t - r)), \quad \lambda \geq 0, \]

where \( f : \mathbb{R} \times \mathbb{R}^2 \to \mathbb{R} \) is continuous and \( T \)-periodic in \( t \), \( \gamma : \mathbb{R} \to \mathbb{R} \) is \( T \)-periodic and nonzero, and \( g : \mathbb{R} \to \mathbb{R} \) is \( C^1 \).
Introducing a new variable \( x \), Equation (18) can be equivalently rewritten in \( \mathbb{R}^2 \) (as in the so-called Liénard plane technique) as follows:

\[
\begin{aligned}
\dot{x}(t) &= \lambda f(t, y(t), y(t - r)) \gamma^{-1}(t), \\
\dot{y}(t) &= (x - G(y)) \gamma(t),
\end{aligned}
\]

(19)

where \( G(y) \) is a primitive of \( g(y) \) and \( \gamma \) is as in Lemma 3. Indeed, taking the derivative of the second equation in (19), we have

\[
\dot{y}(t) = \gamma(t) \left( x(t) - G(y(t)) \right) + \gamma(t) \dot{x}(t) - g(y(t)) \dot{y}(t)
\]

(20)

By this relation, one can easily see that (19) is equivalent to (18). Because of this equivalence, Theorem 2 can be applied to (18):

**Proposition 1.** Let \( f, g \) and \( \gamma \) be as in (18), and let \( \Omega \subseteq [0, \infty) \times C_T^1(\mathbb{R}) \) be open. Define the open subset of \( C_T(\mathbb{R} \times \mathbb{R}) \)

\[
\Omega := \{ (\lambda, \varphi, \psi) \in [0, \infty) \times C_T(\mathbb{R} \times \mathbb{R}) : (\lambda, \varphi) \in \Omega \},
\]

and, according to the notation of Theorem 2

\[
\Omega_{\mathcal{R}^2} = \{ (\lambda, p, q) \in [0, \infty) \times \mathbb{R} \times \mathbb{R} : (\lambda, \tilde{p}, \tilde{q}) \in \Omega \}.
\]

Consider the vector field \( \nu \) in \( \mathbb{R}^2 \), given by

\[
\nu(p, q) := (\tilde{w}(q), p - G(q)),
\]

with \( \tilde{w}(q) := \frac{1}{2} \int_0^T f(t, q, q) \gamma^{-1}(t) \, dt \). Assume that \( \nu \) is admissible in \( \Omega_{\mathcal{R}^2} \) for the degree and that \( \deg(\nu, \Omega_{\mathcal{R}^2}) \neq 0 \). Then, there exists a connected set of nontrivial \( T \)-pairs for (18) whose closure meets the set \( \{ (0, \tilde{p}) \in \Omega : \tilde{w}(p) = 0 \} \) and, is not compact.

**Proof.** By Theorem 2, there exists a connected set \( \Gamma \) of nontrivial \( T \)-triples for (19)

\[
\left\{ (0, \tilde{p}, \tilde{q}) \in \tilde{\Omega} : \tilde{w}(q) = 0, \quad p = G(q) \right\}
\]

and is not compact.

Observe that to any \( (\lambda, y, z) \in \Gamma \) one can associate the nontrivial \( T \)-pair \( (\lambda, y) \) for (18). In this way, one gets a connected set of nontrivial \( T \)-pairs for (18) whose closure meets the set \( \{ (0, \tilde{p}) \in \Omega : \tilde{w}(p) = 0 \} \) and is not compact. \( \square \)

**Example 3.** Consider Equation (18) with \( \gamma(t) = \sin(t) + 2 \) and \( g(y) \equiv 1 \); that is:

\[
\dot{x}(t) = \left( \frac{\cos(t)}{\sin(t) + 2} - (\sin(t) + 2) \right) \dot{x}(t) + \lambda x(t - r).
\]

(21)

Take \( T = 2\pi \). Clearly, the average \( \bar{x} = 2 \) and, for any \( q \in \mathbb{R} \),

\[
\tilde{w}(q) = \frac{1}{2\pi} \int_0^{2\pi} \frac{q}{\sin(t) + 2} \, dt = \frac{q}{\sqrt{3}}.
\]

Let \( \Omega = [0, \infty) \times C_T^1(\mathbb{R}) \). The vector field \( \nu(p, q) = (q/\sqrt{3}, p - q) \) is clearly admissible in \( \Omega_{\mathcal{R}^2} = \mathbb{R}^2 \) and has degree 1. Then, by Proposition 1, there exists a connected set of nontrivial \( 2\pi \)-pairs for (20) whose closure meets the set

\[
\{ (0, \tilde{p}) \in [0, \infty) \times C_T^1(\mathbb{R}) : \tilde{w}(p) = 0 \}
\]

and is not compact.
Remark 4. When \( \gamma(t) \equiv 1 \), the system of equations (19) reduces to
\[
\begin{aligned}
\dot{x}(t) &= \lambda f(t, y(t), y(t-r)), \\
\dot{y}(t) &= (x - G(y)),
\end{aligned}
\]
which is equivalent to the equation
\[
\dot{y}(t) = -g(y(t))\dot{y}(t) + \lambda f(t, y(t), y(t-r)).
\]
in the particular case when \( f(t, y(t), y(t-r)) = f(y(t-r)) \). Equation (21) gives a so-called delayed Liénard equation (or Liénard sunflower-type equation) see, e.g., [11, 12, 13, 21, 24, 25, 26]. Clearly, Proposition 1 applies (for the non-delayed case, see [21]).

When \( f \) does not depend on \( t \), Proposition 4 combined with Lemma 3 implies the main result of this section concerning Equation (14):

**Theorem 4.** Let \( \phi \) and \( a \) be as in (14) and let \( \Omega \subseteq [0, \infty) \times C_T^1(\mathbb{R}) \) be open. Take \( W_\Omega := \{ p \in \mathbb{R} : (0, \bar{p}) \in \Omega \} \), and let \( w(q) := \phi(q, q) \). Assume that \( \deg(w, W_\Omega) \) is well-defined and nonzero. Then, there exists a connected set of nontrivial \( T \)-pairs for (14), whose closure meets the set \( \{(0, \bar{p}) \in \Omega : w(p) = 0\} \) and is not compact.

**Proof.** By Lemma 3 there exists a unique \( T \)-periodic function of constant sign \( \sigma : \mathbb{R} \to \mathbb{R} \) such that \( a(t) = \dot{\sigma}(t)/\sigma(t) - \sigma(t) \). Therefore, (14) can be written in the form (15) with \( f(t, p, q) = \phi(p, q) \) for all \( (p, q) \in \mathbb{R}^2 \).

Take \( G(y) = y \). Then the maps \( \bar{w} \) and \( \nu \) of Proposition 4 become, respectively
\[
\bar{w}(q) = \frac{1}{T} \int_0^T \frac{\phi(q, q)}{\sigma(t)} \, dt = \Phi(q, q) \frac{1}{T} \int_0^T \frac{dt}{\sigma(t)}, \quad \nu(p, q) = (\bar{w}(q), p - q).
\]

Let \( \hat{\Omega}_{\mathbb{R}_2} \) as in Proposition 4 with \( \gamma = \sigma \). One easily checks that
\[
\deg(\nu, \hat{\Omega}_{\mathbb{R}_2}) = -\text{sign} \left( \frac{1}{T} \int_0^T \frac{dt}{\sigma(t)} \right) \deg(w, W_\Omega) = \text{sign} \dot{\sigma} \deg(w, W_\Omega),
\]
the last equality being a consequence of Remark 3. Thus, \( \deg(w, W_\Omega) \neq 0 \) implies \( \deg(\nu, \hat{\Omega}_{\mathbb{R}_2}) \neq 0 \). The assertion now follows from Proposition 4.

In the following example we consider the case of Equation (14) when the perturbing term \( \phi(y(t), y(t-r)) = \sin(y(t-r)) \).

**Example 5** (Sunflower-like equation). Consider the following scalar equation:
\[
\dot{x}(t) = a(t)x(t) + \lambda \sin(x(t-r)).
\]
where \( a : \mathbb{R} \to \mathbb{R} \) is as in (14). Let \( \Omega \) be the open subset of \([0, \infty) \times C_T^1(\mathbb{R}) \) given by \( \Omega = [0, \infty) \times C_T^1((-1, 1)) \), and let \( W_\Omega \) be as in Theorem 4. Let \( T = 2\pi \). One immediately checks that \( \deg(w, W_\Omega) = 1 \), where
\[
w(p) = \frac{1}{2\pi} \int_0^{2\pi} \sin(p) \, dt = \sin(p).
\]
By Theorem 4 there exists a connected set of nontrivial \( T \)-pairs for (14), whose closure meets the set \( \{(0, \bar{p}) \in \Omega : w(p) = 0\} \) and is not compact.

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