## Research Article

# Variational Methods for Almost Periodic Solutions of a Class of Neutral Delay Equations 

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Received 14 June 2007; Accepted 27 November 2007
Recommended by H. Bevan Thompson
We provide new variational settings to study the a.p. (almost periodic) solutions of a class of nonlinear neutral delay equations. We extend Shu and Xu (2006) variational setting for periodic solutions of nonlinear neutral delay equation to the almost periodic settings. We obtain results on the structure of the set of the a.p. solutions, results of existence of a.p. solutions, results of existence of a.p. solutions, and also a density result for the forced equations.

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## 1. Introduction

The aim of this paper is the study of a.p. (almost periodic) solutions of neutral delay equations of the following form:

$$
\begin{align*}
& D_{1} L\left(x(t-r), x(t-2 r), x^{\prime}(t-r), x^{\prime}(t-2 r), t-r\right)+D_{2} L\left(x(t), x(t-r), x^{\prime}(t), x^{\prime}(t-r), t\right) \\
& \quad=\frac{d}{d t}\left[D_{3} L\left(x(t-r), x(t-2 r), x^{\prime}(t-r), x^{\prime}(t-2 r), t-r\right)+D_{4} L\left(x(t), x(t-r), x^{\prime}(t), x^{\prime}(t-r), t\right)\right] \tag{1.1}
\end{align*}
$$

where $L:\left(\mathbb{R}^{n}\right)^{4} \times \mathbb{R} \rightarrow \mathbb{R}$ is a differentiable function $D_{j}$ denotes the partial differential with respect to the $j$ th vector variable, and $r \in(0, \infty)$ is fixed. We will consider the almost periodicity in the sense of Corduneau [1], and in the sense of Besicovitch [2].

A special case of (1.1) is the following forced neutral delay equation:

$$
\begin{align*}
& D_{1} K\left(x(t-r), x(t-2 r), x^{\prime}(t-r), x^{\prime}(t-2 r)\right)+D_{2} K\left(x(t), x(t-r), x^{\prime}(t), x^{\prime}(t-r)\right) \\
& -\frac{d}{d t}\left[D_{3} K\left(x(t-r), x(t-2 r), x^{\prime}(t-r), x^{\prime}(t-2 r)\right)+D_{4} K\left(x(t), x(t-r), x^{\prime}(t), x^{\prime}(t-r)\right)\right]=b(t) \tag{1.2}
\end{align*}
$$

where $K:\left(\mathbb{R}^{n}\right)^{4} \rightarrow \mathbb{R}$ is a differentiable function, and $b: \mathbb{R} \rightarrow \mathbb{R}^{n}$ is an a.p. forcing term. To see (1.2) as a special case of (1.1), it suffices to take

$$
\begin{equation*}
L\left(x_{1}, x_{2}, x_{3}, x_{4}, t\right):=K\left(x_{1}, x_{2}, x_{3}, x_{4}\right)-x_{1} \cdot b(t+r) \tag{1.3}
\end{equation*}
$$

where the point denotes the usual inner product in $\mathbb{R}^{n}$.
Another special case of (1.1) is the following forced second-order neutral delay equation:

$$
\begin{equation*}
x^{\prime \prime}(t-r)+D_{1} F(x(t-r), x(t-2 r))+D_{2} F(x(t), x(t-r))=b(t) \tag{1.4}
\end{equation*}
$$

where $b: \mathbb{R}^{n} \rightarrow \mathbb{R}$ and $F:\left(\mathbb{R}^{n}\right)^{2} \rightarrow \mathbb{R}$. To see that this last equation is a special case of (1.1), it suffices to take $L\left(x_{1}, x_{2}, x_{3}, x_{4}, t\right):=(1 / 2)\left\|x_{3}\right\|^{2}-F\left(x_{1}, x_{2}\right)+x_{1} \cdot b(t+r)$, where the norm is the usual Euclidian norm of $\mathbb{R}^{n}$. In their work [3], Shu and Xu study the periodic solutions of this last equation by using a variational method. We want to extend such a view point to the study of the a.p. solutions.

And so our approach to the study of the a.p. solutions of (1.1) consists to search critical points of a functional $\Phi$ defined on suitable Banach spaces of a.p. functions by

$$
\begin{equation*}
\Phi(x):=\lim _{T \rightarrow+\infty} \frac{1}{2 T} \int_{-T}^{T} L\left(x(t), x(t-r), x^{\prime}(t), x^{\prime}(t-r), t\right) d t \tag{1.5}
\end{equation*}
$$

At this time, we give some historical elements. Recall that the work [4] of Elsgolc treats the calculus of variations with a retarded argument on a bounded real interval. This work was followed by these ones of Hughes [5] and Sabbagh [6]. Since the variational problems can be seen as optimal control problems, recall also the existence of the theory of the periodic optimal control with retarded argument as developed by Colonius [7]. For instance, we consider a periodic optimal control problem with a criterion of the form $(1 / T) \int_{0}^{T} g(x(t), u(t), t) d t$ and with an equation of motion of the form $x^{\prime}(t)=f(x(t), x(t-r), u(t), t)$, where $x(t)$ is the state variable and $u(t)$ the control variable. In the special case, where $f(x(t), x(t-r), u(t), t)=$ $f_{1}(x(t), x(t-r), t)+u(t)$, the previous optimal control problem can be transformed into a calculus of variations problem with the criterion $(1 / T) \int_{0}^{T} g\left(x(t), f_{1}(x(t), x(t-r), t)-x^{\prime}(t), t\right) d t$, which is a special case of (1.5). Note that the Euler-Lagrange equation of such a variational problem is a special case of (1.1).

On another hand, calculus of variations in mean time was developed to study the a.p. solutions of some (nonretarded) differential equations [8-13]. Here, we extend this approach to treat equation like (1.1).

Now we describe the contents of this paper. In Section 2, we precise the notations about the function spaces used later. In Section 3, we establish a variational formalism suitable to the Bohr a.p. solutions; we give a variational principle and a result on the structure of the set of the a.p. solutions of (1.1) in the convex case. In Section 4, we establish a variational formalism suitable to the Besicovitch a.p. solutions, we give a variational principle, results of existence, and a result of density for the a.p. forced equations.

## 2. Notations

$\operatorname{AP}^{0}\left(\mathbb{R}^{n}\right)$ is the space of the Bohr almost periodic (Bohr a.p.) functions from $\mathbb{R}$ in $\mathbb{R}^{n}$; endowed with the supremum $\|\cdot\|_{\infty}$, it is a Banach space [1].
$\operatorname{AP}^{1}\left(\mathbb{R}^{n}\right):=\left\{x \in \mathcal{C}^{1}\left(\mathbb{R}, \mathbb{R}^{n}\right) \cap \mathrm{AP}^{0}\left(\mathbb{R}^{n}\right): x^{\prime} \in \operatorname{AP}^{0}\left(\mathbb{R}^{n}\right)\right\}$; endowed with the norm $\|x\|_{\mathcal{C}^{1}}:=$ $\|x\|_{\infty}+\left\|x^{\prime}\right\|_{\infty}$, it is a Banach space.

When $k \in \mathbb{N}^{*} \cup\{\infty\}, \operatorname{AP}^{k}\left(\mathbb{R}^{n}\right):=\left\{x \in \mathcal{C}^{k}\left(\mathbb{R}, \mathbb{R}^{n}\right): \forall j \leq k, d^{j} x / d t^{j} \in \operatorname{AP}^{0}\left(\mathbb{R}^{n}\right)\right\}$.
When $x \in \mathrm{AP}^{0}\left(\mathbb{R}^{n}\right)$, its mean value is

$$
\begin{equation*}
\mathcal{M}\{x(t)\}_{t}:=\lim _{T \rightarrow+\infty} \frac{1}{2 T} \int_{-T}^{T} x(t) d t \text { exists in } \mathbb{R}^{n} . \tag{2.1}
\end{equation*}
$$

The Fourier-Bohr coefficients of $x \in \mathrm{AP}^{0}\left(\mathbb{R}^{n}\right)$ are the complex vectors

$$
\begin{equation*}
a(x ; \lambda):=\lim _{T \rightarrow+\infty} \frac{1}{2 T} \int_{-T}^{T} e^{-i \lambda t} x(t) d t \tag{2.2}
\end{equation*}
$$

and $\Lambda(x):=\{\lambda \in \mathbb{R}: a(x, \lambda) \neq 0\}$.
When $p \in[1, \infty), B^{p}\left(\mathbb{R}^{n}\right)$ is the completion of $\mathrm{AP}^{0}\left(\mathbb{R}^{n}\right)\left(\right.$ in $\left.L_{\text {loc }}^{p}\left(\mathbb{R}, \mathbb{R}^{n}\right)\right)$ with respect to the norm $\|u\|_{p}:=\mathcal{M}\left\{|u|^{p}\right\}^{1 / p}$. When $p=2, B^{2}\left(\mathbb{R}^{n}\right)$ is a Hilbert spaces and its norm $\|\cdot\|_{2}$ is associated to the inner product $(u \mid v):=\mathcal{M}\{u \cdot v\}[2]$. The elements of these spaces $B^{p}\left(\mathbb{R}^{n}\right)$ are called Besicovitch almost periodic (Besicovitch a.p.) functions.

Recall the useful following fact: if $\left(u_{m}\right)_{m}$ is a sequence in $\mathrm{AP}^{0}\left(\mathbb{R}^{n}\right)$ and if $u \in L_{\text {loc }}^{p}\left(\mathbb{R}, \mathbb{R}^{n}\right)$ (Lebesgue space), which satisfy

$$
\begin{equation*}
\overline{\mathcal{M}}\left\{\left\|u_{m}-u\right\|^{p}\right\}^{1 / p}=\left(\limsup _{T \rightarrow \infty} \frac{1}{2 T} \int_{-T}^{T}\left\|u_{m}-u\right\|^{p} d t\right)^{1 / p} \rightarrow 0 \quad(m \longrightarrow 0) \tag{2.3}
\end{equation*}
$$

then $u \in B^{p}\left(\mathbb{R}^{n}\right)$ and we have $\left\|u_{m}-u\right\|_{p} \rightarrow 0(m \rightarrow 0)$.
We use the generalized derivative $\nabla u \in B^{2}\left(\mathbb{R}^{n}\right)$ of $u \in B^{2}\left(\mathbb{R}^{n}\right)$ (when it exists) defined by $\|\nabla u-(1 / s)(u(\cdot+s)-u)\|_{2} \rightarrow 0(s \rightarrow 0)$, and we define $B^{1,2}\left(\mathbb{R}^{n}\right):=\left\{u \in B^{2}\left(\mathbb{R}^{n}\right): \nabla u \in B^{2}\left(\mathbb{R}^{n}\right)\right\} ;$ endowed with the inner product $\langle u \mid v\rangle:=(u \mid v)+(\nabla u \mid \nabla v), B^{1,2}\left(\mathbb{R}^{n}\right)$ is a Hilbert space [11, 13].

If $E$ and $F$ are two finite-dimensional-normed spaces, $\operatorname{APU}(E \times \mathbb{R}, F)$ stands for the space of the functions $f: E \times \mathbb{R} \rightarrow F,(x, t) \mapsto f(x, t)$, which are almost periodic in $t$ uniformly with respect to $x$ in the classical sense given in [14].

To make the writing less heavy, we sometimes use the notations

$$
\begin{equation*}
\underline{u}(t):=(u(t), u(t-r), \nabla u(t), \nabla u(t-r)), \tag{2.4}
\end{equation*}
$$

when $u \in B^{1,2}\left(\mathbb{R}^{n}\right)$, and

$$
\begin{equation*}
\underline{x}(t):=\left(x(t), x(t-r), x^{\prime}(t), x^{\prime}(t-r)\right), \tag{2.5}
\end{equation*}
$$

when $x \in \mathrm{AP}^{1}\left(\mathbb{R}^{n}\right)$.

## 3. A variational setting for the Bohr a.p. functions

We consider the following condition:

$$
\begin{align*}
& L \in \operatorname{APU}\left(\left(\mathbb{R}^{n}\right)^{4} \times \mathbb{R}, \mathbb{R}\right) \text {, for all }(X, t) \in\left(\mathbb{R}^{n}\right)^{4} \times \mathbb{R}, \\
& \text { the partial differential } D_{X} L(X, t) \text { exists, and }  \tag{3.1}\\
& D_{X} L \in \operatorname{APU}\left(\left(\mathbb{R}^{n}\right)^{4} \times \mathbb{R}, \quad \mathscr{L}\left(\left(\mathbb{R}^{n}\right)^{4}, \mathbb{R}\right)\right) .
\end{align*}
$$

Lemma 3.1. Under [9] the functional $\Phi: A P^{1}\left(\mathbb{R}^{n}\right) \rightarrow \mathbb{R}$ defined by (1.5) which is of class $\mathcal{C}^{1}$, and for all $x, h \in A P^{1}\left(\mathbb{R}^{n}\right)$, then

$$
\begin{align*}
D \Phi(x) \cdot h=\mathcal{M} & \left\{D_{1} L(\underline{x}(t), t) \cdot h(t)+D_{2} L(\underline{x}(t), t) \cdot h(t-r)\right.  \tag{3.2}\\
& \left.+D_{3} L(\underline{x}(t), t) \cdot h^{\prime}(t)+D_{4} L(\underline{x}(t), t) \cdot h^{\prime}(t-r)\right\}_{t} .
\end{align*}
$$

Proof. We introduce the linear operator $\tau: \mathrm{AP}^{1}\left(\mathbb{R}^{n}\right) \rightarrow\left(\mathrm{AP}^{0}\left(\mathbb{R}^{n}\right)\right)^{4}$ by setting $\tau(x)(t):=\underline{x}(t)$. The four components of $\tau$ are continuous linear operators that imply the continuity of $\tau$, and therefore $\tau$ is of class $\mathcal{C}^{1}$, and for all $x, h \in \mathrm{AP}^{1}\left(\mathbb{R}^{n}\right)$ we have $D \tau(x) \cdot h=\tau(h)$.

Under (3.1), the Nemytski operator $\mathcal{N}_{L}:\left(\mathrm{AP}^{0}\left(\mathbb{R}^{n}\right)\right)^{4} \rightarrow \mathrm{AP}^{0}\left(\mathbb{R}^{n}\right)$, defined by $\mathcal{N}_{L}(X)(t):=$ $L(X(t), t)$, is of class $\mathcal{C}^{1}$ (cf. [15, Lemma 7]) and we have, for all $X, H \in \operatorname{AP}^{0}\left(\mathbb{R}^{n}\right)^{4}$, $\left(D \mathcal{N}_{L}(X) \cdot H\right)(t)=D_{X} L(X(t), t) \cdot H(t)$.

The linear functional $\mathcal{M}: \mathrm{AP}^{0}\left(\mathbb{R}^{n}\right) \rightarrow \mathbb{R}$ is continuous, therefore it is of class $\mathcal{C}^{1}$ and we have, for all $\phi, \psi \in \operatorname{AP}^{0}\left(\mathbb{R}^{n}\right), D \mathcal{M}\{\phi\} \cdot \psi=\mathcal{M}\{\psi\}$.

And so $\Phi=\mathscr{M} \circ \mathcal{N}_{L} \circ \tau$ is of class $\mathcal{C}^{1}$. Furthermore, we have

$$
\begin{align*}
D \Phi(x) \cdot h & =D \mathcal{M}\left(\mathcal{N}_{L} \circ \tau(x)\right) \circ D \mathcal{N}_{L}(\tau(x)) \circ D \tau(x) \cdot h \\
& =\mathscr{M}\left\{D \mathcal{N}_{L}(\tau(x)) \cdot \tau(h)\right\}  \tag{3.3}\\
& =\mathcal{M}\left\{D_{X} L(\underline{x}(t), t) \cdot \underline{h}(t)\right\}_{t}
\end{align*}
$$

and expressing $D_{X} L$ in terms of $D_{j} L$, we obtain the announced formula.
Note that in the case without delay, when $L$ is autonomous, that is, $L(X, t)=L(X)$, in [9], it is established that the functional $x \mapsto \mathcal{M}\left\{L\left(x, x^{\prime}\right)\right\}$ is of class $\mathcal{C}^{1}$ when $L$ is of class $\mathcal{C}^{1}$. In [16], we can find a proof of the differentiability of the Nemytski operator on $\operatorname{AP}^{0}\left(\mathbb{R}^{n}\right)$ which is different to this one of [9].

Theorem 3.2 (variational principle). Under (3.1), for $x \in A P^{1}\left(\mathbb{R}^{n}\right)$, the following two assertions are equivalent.
(i) $D \Phi(x)=0$, that is, $x$ is a critical point of $\Phi$ in $\operatorname{AP}^{1}\left(\mathbb{R}^{n}\right)$.
(ii) $x$ is a Bohr a.p. solution of (1.1).

Proof. First we assume (i). Since the mean value is translation invariant, we have

$$
\begin{align*}
& \mathcal{M}\left\{D_{2} L(\underline{x}(t), t) \cdot h(t-r)\right\}_{t}=\mathcal{M}\left\{D_{2} L(\underline{x}(t+r), t+r) \cdot h(t)\right\}_{t^{\prime}}  \tag{3.4}\\
& \mathcal{M}\left\{D_{4} L(\underline{x}(t), t) \cdot h^{\prime}(t-r)\right\}_{t}=\mathcal{M}\left\{D_{4} L(\underline{x}(t+r), t+r) \cdot h^{\prime}(t)\right\}_{t^{\prime}}
\end{align*}
$$

and so by using Lemma 3.1 we obtain, for all $h \in \operatorname{AP}^{1}\left(\mathbb{R}^{n}\right)$,

$$
\begin{align*}
0= & \mathcal{M}\left\{\left(D_{1} L(\underline{x}(t), t)+D_{2} L(\underline{x}(t+r), t+r)\right) \cdot h(t)\right\}_{t} \\
& +\mathcal{M}\left\{\left(D_{3} L(\underline{x}(t), t)+D_{4} L(\underline{x}(t+r), t+r)\right) \cdot h^{\prime}(t)\right\}_{t} . \tag{3.5}
\end{align*}
$$

Setting $q(t):=D_{1} L(\underline{x}(t), t)+D_{2} L(\underline{x}(t+r), t+r)$, denoting by $q_{k}(t)$ its coordinates for $k=1, \ldots, n$, setting $p(t):=D_{3} L(\underline{x}(t), t)+D_{4} L(\underline{x}(t+r), t+r)$, and denoting by $q_{k}(t)$ its coordinates for $k=1, \ldots, n$, we deduce from the previous equality that, for all $\phi \in \operatorname{AP}^{\infty}(\mathbb{R})$, we have $\mathcal{M}\left\{q_{k}(t) \cdot \phi(t)\right\}_{t}=-\mathcal{M}\left\{p_{k}(t) \cdot \phi^{\prime}(t)\right\}_{t}$. Then by reasoning like in the proof of Theorem 1 in [8], we
obtain that $D p_{k}=q_{k}$ in the sense of the a.p. distributions of Schwartz [17], and by using the proposition of the Fourier-Bohr series, we obtain that $p_{k}$ is $\mathbb{C}^{1}$ and that $p_{k}^{\prime}=q_{k}$ in the ordinary sense. From this, we obtain that $p(\cdot-r)$ is $\mathbb{C}^{1}$ and that $p^{\prime}(t-r)=q(t-r)$ which is exactly (ii).

Conversely by using the formula $\mathcal{M}\left\{l \cdot y^{\prime}\right\}=-\mathcal{M}\left\{l^{\prime} \cdot y\right\}$, for all $l \in \operatorname{AP}^{1}\left(\mathcal{L}\left(\mathbb{R}^{n}, \mathbb{R}\right)\right)$ and $y \in \operatorname{AP}^{1}\left(\mathbb{R}^{n}\right)$, and by translating the time, we obtain from (ii), for all $h \in \mathrm{AP}^{1}\left(\mathbb{R}^{n}\right)$, the following relation:

$$
\begin{align*}
0 & =\mathcal{M}\left\{\left(D_{1} L(\underline{x}(t), t)+D_{2} L(\underline{x}(t+r), t+r)\right) \cdot h(t)+\left(D_{3} L(\underline{x}(t), t)+D_{4} L(\underline{x}(t+r), t+r)\right) \cdot h^{\prime}(t)\right\}_{t} \\
& \left.=\mathcal{M}\left\{D_{1} L(\underline{x}(t), t) h(t)+D_{2} L(\underline{x}(t), t)\right) \cdot h(t-r)+D_{3} L(\underline{x}(t), t) h^{\prime}(t)+D_{4} L(\underline{x}(t), t) \cdot h^{\prime}(t-r)\right\}_{t} \\
& =D \Phi(x) \cdot h, \tag{3.6}
\end{align*}
$$

and so we have (i).
Theorem 3.2 is an extension to the nonautonomous case in presence of a delay of $[8$, Theorem 1]. Now we use Theorem 3.2 to provide some results on the structure of the set of the Bohr a.p. solutions of (1.1) in the case, where $L$ is autonomous and convex.

Theorem 3.3 (structure result). Assume that $L \in \mathcal{C}^{1}\left(\left(\mathbb{R}^{n}\right)^{4}, \mathbb{R}\right)$, and that $L$ is convex. Then the following assertions hold.
(i) The set of the Bohr a.p. solutions of (1.1) is a closed convex subset of $\operatorname{AP}^{1}\left(\mathbb{R}^{n}\right)$.
(ii) If $x^{1}$ is a $T^{1}$-periodic nonconstant solution of (1.1), if $x^{2}$ is a $T^{2}$-periodic nonconstant solution of (1.1), and if $T^{1} / T^{2}$ is no rational, then $(1-\theta) x^{1}+\theta x^{2}$ is a Bohr a.p. nonperiodic solution of (1.1), for all $\theta \in(0,1)$.
(iii) If $x$ is a Bohr a.p. solution of (1.1), then $\mathcal{M}\{x\}$ is a constant solution of (1.1).
(iv) If $x$ is a Bohr a.p. solution of $(1.1)$, if $T \in(0, \infty)$ is such that $a(x, 2 \pi / T) \neq 0$, then there exists a nonconstant $T$-periodic solution of (1.1).

Proof. Since $L$ is convex, the functional $\Phi$ of (1.5) is also convex on $\operatorname{AP}^{1}\left(\mathbb{R}^{n}\right)$. Since $L$ is autonomous and of class $\mathcal{C}^{1}, L$ satisfies (3.1), and so $\Phi$ is of class $\mathcal{C}^{1}$. Therefore, we have $\{x: \Phi(x)=\inf \Phi\}=\{x: D \Phi(x)=0\}$ which is closed and convex, and (i) becomes a consequence of Theorem 3.2. The assertion (ii) is a straightforward consequence of (i).

We introduce $C_{T, v}(x)(t):=(1 / v) \sum_{k=0}^{v-1} x(t+k T)$, when $x$ is a Bohr a.p. solution of (1.1), for all $v \in \mathbb{N}^{*}$. By using a theorem of Besicovitch (see [2, page 144]), there exists a $T$-periodic continuous function, denoted by $x^{T}$, such that $\lim _{v \rightarrow \infty}\left\|C_{T, v}(x)-x^{T}\right\|_{\infty}=0$.

We easily verify that $\lim _{v \rightarrow \infty}\left\|C_{T, v}(x)-x^{T}\right\|_{\mathbb{C}^{1}}=0$.
Since $L$ is autonomous, $t \mapsto x(t+k T)$ is a Bohr a.p. solution of (1.1). Since $C_{T, v}(x)$ is a convex combination of Bohr a.p. solutions of (1.1), $C_{T, v}(x)$ is a Bohr a.p. solution of (1.1), and $x^{T}$ also by using the closeness of the set of Bohr a.p. solutions. And so $x^{T}$ is a $T$-periodic solution of (1.1). By using a straightforward calculation, we see that $a\left(C_{T, v}(x), 2 \pi / T\right)=a(x, 2 \pi / T)$ and consequently $a\left(x^{T}, 2 \pi / T\right)=a(x, 2 \pi / T)$. When $a(x, 2 \pi / T) \neq 0$, then $x^{T}$ is not constant that proves (iv).

To prove (iii) it suffices to choose $T^{1} \in(0, \infty)$ such that $\left(2 \pi / T^{1}\right)(\mathbb{Z}-\{0\}) \cap \Lambda(x)=0$, and then all the Fourier-Bohr coefficients of $x^{T^{1}}$ are zero except (perhaps) the mean value of $x^{T^{1}}$ which is equal to $\mathcal{M}\{x\}$.

The assertions (i) and (ii) are extensions of [8, Theorems 3 and 4]; the assertions (iii) and (iv) are extension to neutral delay equations of [10, Theorem 2].

The space $\left(\mathrm{AP}^{1}\left(\mathbb{R}^{n}\right),\|\cdot\|_{\mathbb{C}^{1}}\right)$ does not possess good topological properties like to be a reflexive space. It is why in the following section we extend our variational formalism to the Hilbert space $B^{1,2}\left(\mathbb{R}^{n}\right)$.

## 4. A variational setting for the Besicovitch a.p. functions

$E$ and $F$ are Euclidean finite-dimensional spaces.
Lemma 4.1. Let $g \in \operatorname{APU}(E \times \mathbb{R}, F)$ be a function which satisfies the following Hölder condition:

$$
\begin{equation*}
\exists \alpha \in(0, \infty), \exists a \in[0, \infty), \forall t \in \mathbb{R}, \forall z, w \in E, \quad|g(z, t)-g(w, t)| \leq a \cdot|z-w|^{\alpha} \tag{4.1}
\end{equation*}
$$

Let $p, q \in[1, \infty)$ be such that $p=\alpha q$.
Then the following two assertions hold.
(i) If $u \in B^{p}(E)$, then $t \mapsto g(u(t), t) \in B^{q}(F)$.
(ii) The Nemytski operator on $g, \mathcal{N}_{g}: B^{p}(E) \rightarrow B^{q}(F)$ defined by $\mathcal{N}_{g} u(t):=g(u(t), t)$ satisfies $\left\|\Omega_{g} u-\Omega_{g} v\right\|_{q} \leq a \cdot\|u-v\|_{p}^{\alpha}$, for all $u, v \in B^{p}(E)$.

Proof. We set $b(t):=g(0, t)$, and so we have $b \in \operatorname{AP}^{0}(\mathbb{R})$ and the Hölder assumption implies $|g(x, t)| \leq a \cdot|z|^{\alpha}+b(t)$, for all $z \in E, t \in \mathbb{R}$. If $u \in B^{p}(E)$, then we have $|g(u(t), t)| \leq a \cdot|u(t)|^{\alpha}+b(t)$, for all $t \in \mathbb{R}$, and since $b$ is continuous, we have $b \in L_{\text {loc }}^{q}(\mathbb{R}, \mathbb{R})$ (the Lebesgue space), and since $\left(|u(t)|^{\alpha}\right)^{q}=|u(t)|^{p}$, we have $|u|^{\alpha} \in L_{\text {loc }}^{q}(\mathbb{R}, \mathbb{R})$. Since $u \in B^{p}(E)$, there exists a sequence $\left(u_{j}\right)_{j}$ in $\operatorname{AP}^{0}(E)$ such that $\lim _{j \rightarrow \infty}\left\|u-u_{j}\right\|_{p}=0$. By using [14, Theorem 2.7, page 16], setting $\varphi_{j}(t):=$ $g\left(u_{j}(t), t\right)$, we have $\varphi_{j} \in \mathrm{AP}^{0}(F)$, and a straightforward calculation gives us the following inequality:

$$
\begin{equation*}
\overline{\mathcal{M}}\left\{\left|g(u(t), t)-\varphi_{j}(t)\right|^{q}\right\}^{1 / q} \leq a \cdot \mathcal{M}\left\{\left|u-u_{j}\right|^{p}\right\}^{1 / q}=a\left\|u-u_{j}\right\|_{p^{\prime}}^{\alpha} \tag{4.2}
\end{equation*}
$$

and consequently we obtain

$$
\begin{equation*}
\lim _{j \rightarrow \infty} \overline{\mathcal{M}}\left\{\left|g(u(t), t)-\varphi_{j}(t)\right|^{p}\right\}^{1 / q}=0 \tag{4.3}
\end{equation*}
$$

that implies $t \mapsto g(u(t), t) \in B^{q}(F)$, and so (i) is proven; moreover the last inequality becomes the one of (ii) when we replace $\varphi_{j}(t)$ by $g(v(t), t)$.

This lemma is an extension to the nonautonomous case of [13, Theorem 1].
Lemma 4.2. Let $f \in \operatorname{APU}(E \times \mathbb{R}, F)$ be a function such that the partial differential $D_{1} f(z, t)$ exists, for all $(z, t) \in E \times \mathbb{R}$, such that $D_{1} f \in \operatorname{APU}(E \times \mathbb{R}, \Omega(E, F))$. We assume the following condition fulfilled.
(C) There exist $a_{1} \in[0, \infty)$, such that, for all $z, w \in E$, and for all $t \in \mathbb{R},\left|D_{1} f(z, t)-D_{1} f(w, t)\right| \leq$ $a_{1} \cdot|z-w|$.

Then the Nemytski operator $\mathcal{N}_{f}: B^{2}(E) \rightarrow B^{1}(F)$, defined by $\mathcal{N}_{f}(u)(t):=f(u(t), t)$, is of class $\mathcal{C}^{1}$ and, for all $u, h \in B^{2}(E)$,

$$
\begin{equation*}
\left(D \mathcal{N}_{f}(u) \cdot h\right)(t)=D_{1} f(u(t), t) \cdot h(t) . \tag{4.4}
\end{equation*}
$$

Proof
First step: we show that there exist $a_{0} \in[0, \infty), b \in B^{1}(E)$, such that, for all $(z, t) \in E \times \mathbb{R}$, $|f(z, t)| \leq a_{0}|z|^{2}+b(t)$. The following hold:

$$
\begin{align*}
\left|D_{1} f(z, t)-D_{1} f(0, t)\right| \leq a_{1} \cdot|z| & \\
\Rightarrow\left|D_{1} f(z, t)\right| & \leq\left|D_{1} f(z, t)-D_{1} f(0, t)\right|+\left|D_{1} f(0, t)\right|  \tag{4.5}\\
& \leq a_{1} \cdot|z|+\left|D_{1} f(0, t)\right| .
\end{align*}
$$

By using the mean value theorem (see [18, page 144]), we have, for all $(z, t) \in E \times \mathbb{R}$,

$$
\begin{align*}
|f(z, t)| & \leq|f(z, t)-f(0, t)|+|f(0, t)| \\
& \leq \sup _{\xi \in[0, z \mid}\left|D_{1} f(\xi, t)\right| \cdot|z-0|+|f(0, t)| \\
& \leq \sup _{\xi \in[0, z[ }\left(a_{1} \cdot|\xi|+\left|D_{1} f(0, t)\right|\right) \cdot|z|+|f(0, t)| \\
& =\left(a_{1} \cdot|z|+\left|D_{1} f(0, t)\right|\right) \cdot|z|+|f(0, t)|  \tag{4.6}\\
& =a_{1} \cdot|z|^{2}+\left|D_{1} f(0, t)\right| \cdot|z|+|f(0, t)| \\
& \leq a_{1} \cdot|z|^{2}+\frac{1}{2}\left|D_{1} f(0, t)\right|^{2}+\frac{1}{2}|z|^{2}+|f(0, t)| \\
& =\left(a_{1}+\frac{1}{2}\right) \cdot|z|^{2}+\frac{1}{2}\left|D_{1} f(0, t)\right|^{2}+|f(0, t)| .
\end{align*}
$$

Setting $b(t):=(1 / 2)\left|D_{1} f(0, t)\right|^{2}+|f(0, t)|$, then $a_{0}:=a_{1}+1 / 2$. Since $f \in \operatorname{APU}(E \times \mathbb{R}, F)$, and $D_{1} f \in \operatorname{APU}(E \times \mathbb{R}, \mathscr{L}(E, F))$, we have $b \in \operatorname{AP}^{0}(E) \subset B^{1}(E)$.
Second step: we show that $t \mapsto f(u(t), t) \in B^{1}(F)$ when $u \in B^{2}(E)$.
Let $u \in B^{2}(E)$. Then the inequality $|f(u(t), t)| \leq a_{0}|u(t)|^{2}+b(t)$ implies that

$$
\begin{equation*}
t \longmapsto f(u(t), t) \in L_{\mathrm{loc}}^{1}(\mathbb{R}, F) \tag{4.7}
\end{equation*}
$$

By using Lemma 4.1 with $p=2, q=2, \alpha=1$, and $g=D_{1} f$, we have $t \mapsto D_{1} f(u(t), t) \in$ $B^{2}(\mathcal{L}(E, F))$. Let $\left(u_{m}\right)_{m}$ be a sequence in $\mathrm{AP}^{0}(E)$ such that $\left\|u-u_{m}\right\|_{2} \rightarrow 0(m \rightarrow \infty)$. By using the mean value theorem (see [18, page 144]), we have, for all $t \in \mathbb{R}$,

$$
\begin{align*}
& \left|f\left(u_{m}(t), t\right)-f(u(t), t)-D_{1} f(u(t), t) \cdot\left(u_{m}(t)-u(t)\right)\right| \\
& \quad \leq\left(\sup _{\substack{\xi \in \in u(t), u_{m}(t) \mid}}\left|D_{1} f(\xi, t)-D_{1} f(u(t), t)\right|\right) \cdot\left|\left(u_{m}(t)-u(t)\right)\right|  \tag{4.8}\\
& \quad \leq a_{1} \cdot \sup _{\xi \in] u(t), u_{m}(t) \mid}|\xi-u(t)| \cdot\left|u_{m}(t)-u(t)\right| \leq a_{1} \cdot\left|u_{m}(t)-u(t)\right|^{2},
\end{align*}
$$

and consequently we obtain

$$
\begin{equation*}
\overline{\mathcal{M}}\left\{\left|f\left(u_{m}(t), t\right)-f(u(t), t)-D_{1} f(u(t), t) \cdot\left(u_{m}(t)-u(t)\right)\right|\right\}_{t} \leq a_{1} \cdot\left\|u_{m}-u\right\|_{2}^{2} . \tag{4.9}
\end{equation*}
$$

Since $t \mapsto D_{1} f(u(t), t) \in B^{2}(\mathcal{L}(E, F))$ and since $u_{m}-u \in B^{2}(E)$, we have

$$
\begin{equation*}
t \longmapsto D_{1} f(u(t), t) \cdot\left(u_{m}(t)-u(t)\right) \in B^{1}(F) . \tag{4.10}
\end{equation*}
$$

By using (see [14, Theorem 2.7, page 16]), we have

$$
\begin{equation*}
t \longmapsto f\left(u_{m}(t), t\right) \in \operatorname{AP}^{0}(F) \subset B^{1}(F), \tag{4.11}
\end{equation*}
$$

and so, by setting

$$
\begin{equation*}
\psi_{m}(t):=f\left(u_{m}(t), t\right)-D_{1} f\left(u_{m}(t), t\right) \cdot\left(u_{m}(t)-u(t)\right), \tag{4.12}
\end{equation*}
$$

we have $\psi_{m} \in B^{1}(F)$. The last inequality implies

$$
\begin{equation*}
\lim _{n \rightarrow \infty} \overline{\mathcal{M}}\left\{\left|f(u(t), t)-\psi_{m}(t)\right|\right\}_{t}=0 \tag{4.13}
\end{equation*}
$$

and therefore we have $t \mapsto f(u(t), t) \in B^{1}(F)$.
Third step: we show that, for all $u \in B^{2}(E)$, the operator $\mathcal{L}(u): B^{2} \rightarrow B^{1}(\mathbb{R})$, defined by $(\mathcal{L}(u)$. $h)(t):=D_{1} f(u(t), t) \cdot h(t)$, is linear continuous. We have yet seen that $t \mapsto D_{1} f(u(t), t) \cdot h(t) \in$ $B^{1}(F)$. The linearity of $\Omega(u)$ is easy to verify. By using a Cauchy-Schwartz-Bunyakovsky inequality, we have

$$
\begin{align*}
\mathcal{M}\left\{\left|D_{1} f(u(t), t) \cdot h(t)\right|\right\}_{t} & \leq \mathcal{M}\left\{\left|D_{1} f(u(t), t)\right| \cdot|h(t)|\right\}_{t} \\
& \leq \mathcal{M}\left\{\left|D_{1} f(u(t), t)\right|^{2}\right\}^{1 / 2} \cdot \mathcal{M}\left\{|h|^{2}\right\}^{1 / 2} \tag{4.14}
\end{align*}
$$

that proves the continuity of $\Omega(u)$.
Fourth step: we show the differentiability of $\Omega_{f}$.
Let $u \in B^{2}(E)$ and $h \in B^{2}(E)$. By using the mean value inequality (see [18, page 144$]$ ), we have, for all $t \in \mathbb{R}$,

$$
\begin{align*}
& \left|f(u(t)+h(t), t)-f(u(t), t)-D_{1} f(u(t), t) \cdot h(t)\right| \\
& \quad \leq \sup _{\xi \in] u(t), u(t)+h(t)[ }\left|D_{1} f(\xi, t)-D_{1} f(u(t), t)\right| \cdot|h(t)| \leq a_{1}\|h\|^{2} \tag{4.15}
\end{align*}
$$

and by using the monotonicity of $\Omega$, we obtain

$$
\begin{equation*}
\mathcal{M}\left\{\left|f(u(t)+h(t), t)-f(u(t), t)-D_{1} f(u(t), t) \cdot h(t)\right|\right\}_{t} \leq a_{1}\|h\|_{2}^{2} \tag{4.16}
\end{equation*}
$$

that is,

$$
\begin{equation*}
\left\|\Omega_{f}(u+h)-\mathcal{N}_{f}(u)-\mathcal{L}(u) \cdot h\right\|_{1} \leq a_{1}\|h\|_{2}^{2} \tag{4.17}
\end{equation*}
$$

that implies that $\Omega_{f}$ is differentiable at $u$ and that $D \mathcal{N}_{f}(u)=\Omega(u)$.
Fifth step: we show that $\Omega_{f}$ is of class $\mathcal{C}^{1}$.
Let $u, v \in B^{2}(E)$. By using (C), for all $h \in B^{2}(E)$, such that $\|h\|_{2} \leq 1$, for all $t \in \mathbb{R}$, we have

$$
\begin{align*}
\left|\left(D_{1} f(u(t), t)-D_{1} f(v(t), t)\right) \cdot h(t)\right| & \leq\left|D_{1} f(u(t), t)-D_{1} f(v(t), t)\right| \cdot|h(t)| \\
& \leq a_{1} \cdot|u(t)-v(t)| \cdot|h(t)| \tag{4.18}
\end{align*}
$$

That implies, by using the Cauchy-Schwartz-Bunyakovsky inequality, the following majorization holds:

$$
\begin{align*}
& \mathcal{M}\left\{\left|\left(D_{1} f(u(t), t)-D_{1} f(v(t), t)\right) \cdot h(t)\right|\right\}_{t} \\
& \quad \leq a_{1} \mathcal{M}\{|u(t)-v(t)| \cdot|h(t)|\}_{t} \leq a_{1}\|u-v\|_{2} \cdot\|h\|_{2} \leq a_{1}\|u-v\|_{2} \tag{4.19}
\end{align*}
$$

Therefore we have $\left\|D \mathcal{N}_{f}(u)-D \mathcal{N}_{f}(v)\right\|_{\perp} \leq a_{1}\|u-v\|_{2}$ that implies the continuity of $D \mathcal{N}_{f}$.

Note that Lemma 4.2 is an extension to the nonautonomous case of [13, Theorem 2].
Theorem 4.3 (variational principle). Let $L:\left(\mathbb{R}^{n}\right)^{4} \times \mathbb{R} \rightarrow \mathbb{R}$,

$$
\begin{equation*}
(X, t)=\left(x_{1}, x_{1}, x_{1}, x_{1}, t\right) \longmapsto L(X, t)=L\left(x_{1}, x_{1}, x_{1}, x_{1}, t\right) \tag{4.20}
\end{equation*}
$$

be a function and let $r \in(0, \infty)$. Assume the following conditions fulfilled:
$L \in \operatorname{APU}\left(\left(\mathbb{R}^{n}\right)^{4} \times \mathbb{R}, \mathbb{R}\right)$, the partial differentials $D_{k} L\left(x_{1}, x_{1}, x_{1}, x_{1}, t\right)$ exist, for all $\left(x_{1}, x_{1}, x_{1}, x_{1}, t\right) \in\left(\mathbb{R}^{n}\right)^{4} \times \mathbb{R}$
and for $k=1, \ldots, 4, D_{k} L \in \operatorname{APU}\left(\left(\mathbb{R}^{n}\right)^{4} \times \mathbb{R}, \Omega\left(\mathbb{R}^{n}, \mathbb{R}\right)\right)$;
there exists $a_{1} \in[0, \infty)$ such that $\left|L_{X}(X, t)-L_{X}(Y, t)\right| \leq a_{1}|X-Y|$,
for all $X, Y \in\left(\mathbb{R}^{n}\right)^{4}$, for all $t \in \mathbb{R}$, where $L_{X}$ is the partial differential with respect to $X \in\left(\mathbb{R}^{n}\right)^{4}$.

Then the functional $J: B^{1,2}\left(\mathbb{R}^{n}\right) \rightarrow \mathbb{R}$, defined by

$$
\begin{equation*}
J(u)=\mathcal{M}\{L(u(t), u(t-r), \nabla u(t), \nabla u(t-r), t)\}_{t^{\prime}} \tag{4.22}
\end{equation*}
$$

is of class $\mathcal{C}^{1}$, and the two following assertions are equivalent.
(i) $D J(u)=0$, that is, $u$ is a critical point of $J$.
(ii) $D_{1} L(u(t-r), u(t-2 r), \nabla u(t-r), \nabla u(t-2 r), t-r)+D_{2} L(u(t), u(t-r), \nabla u(t), \nabla u(t-r), t)=$ $\nabla\left[D_{3} L(u(t-r), u(t-2 r), \nabla u(t-r), \nabla u(t-2 r), t-r)+D_{4} L(u(t), u(t-r), \nabla u(t), \nabla u(t-\right.$ $r), t)]$ (equality in $\left.B^{2}\left(\mathcal{L}\left(\mathbb{R}^{n}, \mathbb{R}\right)\right)\right)$.

Definition 4.4. When $u \in B^{1,2}\left(\mathbb{R}^{n}\right)$ satisfies the equation of (ii) in Theorem 4.3 , we say that $u$ is a weak Besicovitch a.p. solution of (1.1).

Proof. We consider the operator $\mathcal{\perp}: B^{1,2}\left(\mathbb{R}^{n}\right) \rightarrow B^{2}\left(\mathbb{R}^{n}\right)^{4} \equiv B^{2}\left(\left(\mathbb{R}^{n}\right)^{4}\right)$, defined by $(\mathcal{L}(u))(t):=$ $(u(t), u(t-r), \nabla u(t), \nabla u(t-r))$. $\mathcal{L}$ is clearly linear continuous, therefore $\mathcal{L}$ is of class $\mathcal{C}^{1}$ and we have $D \perp(u) \cdot h=\perp(h)$.

We consider the Nemytski operator

$$
\begin{equation*}
\mathcal{N}_{L}: B^{2}\left(\left(\mathbb{R}^{n}\right)^{4}\right) \longrightarrow B^{1}(R), \quad\left(\mathcal{N}_{L}(u)\right)(t):=L(u(t), t) \tag{4.23}
\end{equation*}
$$

By using Lemma 4.2, $\mathcal{N}_{L}$ is of class $\mathcal{C}^{1}$ and, for all $U, H \in B^{1,2}\left(\left(\mathbb{R}^{n}\right)^{4}\right)$, we have

$$
\begin{align*}
\left(D \mathcal{N}_{L}(U) \cdot H\right)(t) & =L_{X}(U(t), t) \cdot H(t) \\
& =\sum_{k=1}^{4} D_{k} L\left(u_{1}(t), u_{2}(t), u_{3}(t), u_{4}(t), t\right) \cdot h_{k}(t) \tag{4.24}
\end{align*}
$$

The mean value $\mathcal{M}: B^{1}(\mathbb{R}) \rightarrow \mathbb{R}$ is linear continuous, therefore it is of class $\mathcal{C}^{1}$, and $D \mathcal{M}\{\phi\} \cdot \psi=$ $\mathcal{M}\{\psi\}$, for all $\phi, \psi \in B^{1}(\mathbb{R})$.

Consequently $J=\mathcal{M} \circ \mathcal{N}_{L} \circ \mathcal{L}$ is of class $\mathcal{C}^{1}$ as a composition of three mappings of class $\mathcal{C}^{1}$.

Let $u \in B^{1,2}\left(\mathbb{R}^{n}\right)$. If (i) is true then, for all $h \in B^{1,2}\left(\mathbb{R}^{n}\right)$, we have

$$
\begin{align*}
0= & D J(u) \cdot h=D \mathcal{M}\left\{\mathcal{N}_{L} \circ \Omega(u)\right\} \circ D \mathcal{N}_{L}(\perp(u)) \circ D \Omega(u) \cdot h \\
= & \mathcal{M}\left\{D \mathcal{N}_{L}(\mathscr{\Omega}(u)) \cdot \Omega(h)\right\} \\
= & \mathcal{M}\left\{D_{1} L(\underline{u}(t), t) \cdot h(t)+D_{2} L(\underline{u}(t), t) \cdot h(t-r)\right.  \tag{4.25}\\
& \left.+D_{3} L(\underline{u}(t), t) \cdot \nabla h(t)+D_{4} L(\underline{u}(t), t) \cdot \nabla h(t-r)\right\}_{t} \\
= & \mathcal{M}\left\{\left(D_{1} L(\underline{u}(t), t)+D_{2} L(\underline{u}(t+r), t+r)\right) \cdot h(t)\right\}_{t} \\
& +\mathcal{M}\left\{\left(D_{3} L(\underline{u}(t), t)+D 4 L(\underline{u}(t+r), t+r)\right) \cdot \nabla h(t)\right\}_{t}
\end{align*}
$$

and then we obtain (ii) by using [13, Proposition 10].
Conversely, if (ii) is true, then $t \mapsto D_{3} \mathcal{N}_{L}(\underline{u}(t), t)+D_{4} \mathcal{N}_{L}(\underline{u}(t+r), t+r) \in B^{1,2}\left(\mathbb{R}^{n}\right)$, and for all $h \in \operatorname{AP}^{1}\left(\mathbb{R}^{n}\right)$, we have

$$
\begin{align*}
& \mathcal{M}\left\{\left(D_{1} L(\underline{u}(t), t)+D_{2} L(\underline{u}(t+r), t+r)\right) \cdot h(t)\right\}_{t^{\prime}} \\
& -\mathcal{M}\left\{\nabla\left(D_{3} L(\underline{u}(t), t)+D_{4} L(\underline{u}(t+r), t+r)\right) \cdot h(t)\right\}_{t}=0, \tag{4.26}
\end{align*}
$$

therefore by using [13, Proposition 9], we obtain

$$
\begin{align*}
0= & \mathcal{M}\left\{\left(D_{1} L(\underline{u}(t), t)+D_{2} L(\underline{u}(t+r), t+r)\right) \cdot h(t)\right. \\
& \left.\quad+\left(D_{3} L(\underline{u}(t), t)+D_{4} L(\underline{u}(t+r), t+r)\right) \cdot h^{\prime}(t)\right\}_{t} \\
= & \mathcal{M}\left\{\left(D_{1} L(\underline{u}(t), t) \cdot h(t)+D_{2} L(\underline{u}(t), t)\right) \cdot h(t-r)\right.  \tag{4.27}\\
& \left.\quad+\left(D_{3} L(\underline{u}(t), t) h^{\prime}(t)+D_{4} L(\underline{u}(t), t)\right) \cdot h^{\prime}(t-r)\right\}_{t} \\
= & D J(u) \cdot h .
\end{align*}
$$

Since $\operatorname{AP}^{1}\left(\mathbb{R}^{n}\right)$ is dense in $B^{1,2}\left(\mathbb{R}^{n}\right)$, we have $D J(u) \cdot h=0$, for all $h \in B^{1,2}\left(\mathbb{R}^{n}\right)$, therefore $D J(u)=0$.

Note that the Theorem 4.3 is an extension to the nonautonomous case of [13, Theorem 4].
Theorem 4.5 (existence, uniqueness). Let $L:\left(\mathbb{R}^{n}\right)^{4} \times \mathbb{R} \rightarrow \mathbb{R}$ be a function which satisfies (4.21) and also satisfies the following two conditions:

$$
\begin{align*}
& L(\cdot, t):\left(\mathbb{R}^{n}\right)^{4} \longrightarrow \mathbb{R} \text { is convex, for all } t \in \mathbb{R} ;  \tag{4.28}\\
& \text { there exist } j \in\{1,2\}, k \in\{3,4\}, c \in(0, \infty) \\
& \text { such that, for all }\left(x_{1}, x_{2}, x_{3}, x_{4}, t\right) \in\left(\mathbb{R}^{n}\right)^{4} \times \mathbb{R} \text {, }  \tag{4.29}\\
& \text { we have } L\left(x_{1}, x_{2}, x_{3}, x_{4}, t\right) \geq c\left(\left|x_{j}\right|^{2}+\left|x_{k}\right|^{2}\right)
\end{align*}
$$

Then there exists a function $u \in B^{1,2}\left(\mathbb{R}^{n}\right)$ which is a weak Besicovitch a.p. solution of (1.1).
Moreover, if in addition the following condition fulfilled:
there exist $i \in\{1,2\}, l \in\{3,4\}, c_{1} \in(0, \infty)$
such that the function $M:\left(\mathbb{R}^{n}\right)^{4} \times \mathbb{R} \longrightarrow \mathbb{R}$, defined by
$M\left(x_{1}, x_{2}, x_{3}, x_{4}, t\right):=L\left(x_{1}, x_{2}, x_{3}, x_{4}, t\right)-\frac{c_{1}}{2}\left|x_{i}\right|^{2}-\frac{c_{1}}{2}\left|x_{l}\right|^{2}$,
is convex with respect to $\left(x_{1}, x_{2}, x_{3}, x_{4}, t\right)$, for all $t \in \mathbb{R}$,
then the weakBesicovitch a.p. solution of (1.1) is unique.

Proof. By using Theorem 4.3, the functional $J$ is of class $\mathcal{C}^{1}$ and, by using (4.28), $J$ is a convex functional. Assumption (4.29) ensures that, for all $u \in B^{1,2}\left(\mathbb{R}^{n}\right)$, we have

$$
\begin{equation*}
J(u) \geq c\left(\mathcal{M}\left\{|u|^{2}\right\}+\mathcal{M}\left\{|\nabla u|^{2}\right\}\right)=c \cdot\|u\|_{1,2}^{2} \tag{4.31}
\end{equation*}
$$

Since the mean value is translation invariant, consequently $J$ is coercive on $B^{1,2}\left(\mathbb{R}^{n}\right)$, and so (see [19, page 46]) there exists $u \in B^{1,2}\left(\mathbb{R}^{n}\right)$ such that $J(u)=\inf J$. Therefore we have $D J(u)=0$ and by using Theorem 4.3, $u$ is a weak Besicovitch a.p. solution of (1.1). The existence is proven.

To treat the uniqueness, we note that, under (4.30), the functional $I: B^{1,2}\left(\mathbb{R}^{n}\right) \rightarrow \mathbb{R}$, defined by $I(u):=J(u)-\left(c_{1} / 2\right) \mathcal{M}\left\{|u|^{2}\right\}-\left(c_{1} / 2\right) \mathcal{M}\left\{|\nabla u|^{2}\right\}$, is convex and since $J$ is of class $\mathcal{C}^{1}, I$ is also of class $\mathcal{C}^{1}$. Note that we have $D I(u)=D J(u)-c_{1}\langle u \mid \cdot\rangle$. By using the Minty monotonicity of the differential of a convex functional, for all $u, v \in B^{1,2}\left(\mathbb{R}^{n}\right)$, we have

$$
\begin{align*}
0 & \leq\langle D I(u)-D I(v), u-v\rangle=\langle D J(u)-D J(v), u-v\rangle-c_{1}\langle u-v \mid u-v\rangle \\
& \Longrightarrow\langle D J(u)-D J(v), u-v\rangle \geq c_{1}\|u-v\|_{1,2}^{2} . \tag{4.32}
\end{align*}
$$

Now if $u$ and $v$ are two weak Besicovitch a.p. solutions of (1.1), by using Theorem 4.3 we have $D J(u)=D J(v)=0$, and consequently $c_{1}\|u-v\|_{1,2}^{2}=0$, therefore $u=v$.

Theorem 4.6 (existence and density). Let $K \in \mathcal{C}^{2}\left(\left(\mathbb{R}^{n}\right)^{4}, \mathbb{R}\right)$ be a function which satisfies the following conditions:
there exists $a_{0} \in[0, \infty)$ such that $|K(x)| \leq a_{0}|X|^{2}$,
for all $X \in\left(\mathbb{R}^{n}\right)^{4}$, the following holds:
there exist $j \in\{1,2\}, k \in\{3,4\}, c \in(0, \infty)$
such that the function $G:\left(\mathbb{R}^{n}\right)^{4} \longrightarrow \mathbb{R}$, defined by
$G\left(x_{1}, x_{2}, x_{3}, x_{4}\right):=K\left(x_{1}, x_{2}, x_{3}, x_{4}\right)-\frac{c}{2}\left|x_{j}\right|^{2}-\frac{c}{2}\left|x_{k}\right|^{2}$,
is convex and nonnegative on $\left(\mathbb{R}^{n}\right)^{4}$.
The differential DK is Lipschitzian on $\left(\mathbb{R}^{n}\right)^{4}$.
Then the following conclusions hold.
(i) For all $b \in B^{2}\left(\mathbb{R}^{n}\right)$, there exists a unique $u \in B^{1,2}\left(\mathbb{R}^{n}\right)$ which is a weak Besicovitch a.p. solution of (1.2).
(ii) The set of the $b \in \operatorname{AP}^{0}\left(\mathbb{R}^{n}\right)$ for which there exists a Bohr a.p. solution of (1.2) is dense in $\operatorname{AP}^{0}\left(\mathbb{R}^{n}\right)$ with respect to the norm

$$
\begin{equation*}
\|b\|_{*}:=\sup \left\{\mathcal{M}\{b \cdot h\}: h \in B^{1,2}\left(\mathbb{R}^{n}\right),\|h\|_{1,2} \leq 1\right\} \tag{4.34}
\end{equation*}
$$

Proof. We introduce the functionals $E$ and $E_{1}$ from $B^{1,2}\left(\mathbb{R}^{n}\right)$ in $\mathbb{R}$ setting $E(u):=\mathcal{M}\{K(\underline{u}(t))\}_{t}$ and $E_{1}(u):=\mathcal{M}\{G(\underline{u}(t))\}_{t}$. They are special cases of the functional $J$ of the Theorem 4.3, and consequently they are of class $\mathcal{C}^{1}$. Note that $E_{1}(u)=E(u)-(c / 2)\|u\|_{1,2}^{2}$. By using the F. Riesz isomorphism $j: B^{1,2}\left(\mathbb{R}^{n}\right) \rightarrow B^{1,2}\left(\mathbb{R}^{n}\right)^{*},\langle j(u), v\rangle=\langle u, v\rangle$, for all $u, v \in B^{1,2}\left(\mathbb{R}^{n}\right)$, we can define
the gradients $\operatorname{grad} E(u):=j^{-1}(D E(u))$ and $\operatorname{grad} E_{1}(u):=j^{-1}\left(D E_{1}(u)\right)$. By using the Mintymonotonicity of grad $E_{1}$ (due to the convexity of $E_{1}$ ) we have, for all $u, v \in B^{1,2}\left(\mathbb{R}^{n}\right)$,

$$
\begin{equation*}
0 \leq\left\langle\operatorname{grad} E_{1}(u)-\operatorname{grad} E_{1}(v) \mid u-v\right\rangle=\langle\operatorname{grad} E(u)-\operatorname{grad} E(v) \mid u-v\rangle-c \cdot\|u-v\|_{1,2}^{2} \tag{4.35}
\end{equation*}
$$

that implies that grad $E$ is strongly monotone and consequently (see [20, page 100]) the following property holds:

$$
\begin{equation*}
\operatorname{grad} E \text { is a homeomorphism from } B^{1,2}\left(\mathbb{R}^{n}\right) \quad \text { on } B^{1,2}\left(\mathbb{R}^{n}\right) \tag{4.36}
\end{equation*}
$$

From each $b \in B^{2}\left(\mathbb{R}^{n}\right)$ we define the linear functional $b^{\#} \in B^{1,2}\left(\mathbb{R}^{n}\right)^{*}$ by setting $\left\langle b^{\#}, h\right\rangle:=$ $\mathcal{M}\{b(t+r) \cdot h(t)\}_{t}$.

Therefore we have $j^{-1}\left(b^{\#}\right) \in B^{1,2}\left(\mathbb{R}^{n}\right)$ and by using (4.36), there exists $u \in B^{1,2}\left(\mathbb{R}^{n}\right)$ such that $\operatorname{grad} E(u)=j^{-1}\left(b^{\#}\right)$, that is, $D E(u)=b^{\#}$ which means that, for all $h \in B^{1,2}\left(\mathbb{R}^{n}\right)$,

$$
\begin{align*}
& \mathcal{M}\{D K(\underline{u}(t)) \cdot \underline{h}(t)\}_{t}=\mathcal{M}\{b(t+r) \cdot h(t)\}_{t^{\prime}} \text { that is, } \\
& \mathcal{M}\left\{\left[D_{1} K(\underline{u}(t))+D_{2} K(\underline{u}(t+r))-b(t+r)\right] \cdot h(t)\right.  \tag{4.37}\\
& \left.\quad+\left[D_{3} K(\underline{u}(t))+D_{4} K(\underline{u}(t+r))-b(t+r)\right] \cdot \nabla h(t)\right\}_{t}=0
\end{align*}
$$

and by using [13, Proposition 10], we obtain that $u$ is a weak Besicovitch a.p. solution of (1.2).
About the uniqueness, note that if $v$ is a weak Besicovitch a.p. solution of (1.2), then we verify that $\mathcal{M}\{D K(\underline{v}(t)) \cdot \underline{h}(t)\}_{t}=\mathcal{M}\{b(t+r) \cdot h(t)\}_{t}$, for all $h \in B^{1,2}\left(\mathbb{R}^{n}\right)$, and consequently $D E(v)=b^{\#}$, that is, $\operatorname{grad} E(v)=j^{-1}\left(b^{\#}\right)=\operatorname{grad} E(u)$, and by using (4.36), we have $u=v$. And so (i) is proven.

Now we introduce the nonlinear unbounded operator

$$
\begin{equation*}
\mathcal{K}: \operatorname{Dom}(\mathcal{K}) \subset B^{1,2}\left(\mathbb{R}^{n}\right) \longrightarrow B^{2}\left(\mathbb{R}^{n}\right) \tag{4.38}
\end{equation*}
$$

defined by

$$
\begin{equation*}
(\nless K(u))(t):=D_{1} K(\underline{u}(t-r))+D_{2} K(\underline{u}(t))-\nabla\left[D_{3} K(\underline{u}(t-r))+D_{4} K(\underline{u}(t))\right] . \tag{4.39}
\end{equation*}
$$

And so $\nless(u)=b$ means that $u$ is a weak Besicovitch a.p. solution of (1.2). By using the assertion (i), $\mathcal{K}$ is bijective. We verify that

$$
\begin{equation*}
\|\mathcal{K}(u)-\not \subset(v)\|_{*}=\|\operatorname{grad} E(u)-\operatorname{grad} E(v)\|_{1,2^{\prime}} \tag{4.40}
\end{equation*}
$$

for all $u, v \in \operatorname{Dom}(\not \subset)$, and by using (4.36) we see that $\mathcal{K}$ is a homeomorphism from $\operatorname{Dom}(\not \subset)$ on $B^{2}\left(\mathbb{R}^{n}\right)$. Since $A^{2}\left(\mathbb{R}^{n}\right)$ is dense in $B^{1,2}\left(\mathbb{R}^{n}\right), \mathcal{K}\left(\operatorname{AP}^{2}\left(\mathbb{R}^{n}\right)\right)$ is dense in $B^{2}\left(\mathbb{R}^{n}\right)$ with respect to the norm $\|\cdot\|_{*}$, and since $\mathcal{K}\left(\operatorname{AP}^{2}\left(\mathbb{R}^{n}\right)\right) \subset A P^{0}\left(\mathbb{R}^{n}\right) \subset B^{2}\left(\mathbb{R}^{n}\right)$, we have proven (ii).

This result is an extension to the neutral delay equations of [13, Theorem 5].

## Acknowledgment

The authors thank an anonymous referee for several remarks which have permitted them to improve Theorems 4.3 and 4.5.

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