# INFINITELY MANY SOLUTIONS FOR A PERTURBED SCHRÖDINGER EQUATION 

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#### Abstract

We find multiple solutions for a nonlinear perturbed Schrödinger equation by means of the so-called Bolle's method.


1. Introduction. This note concerns with the elliptic equation

$$
\begin{equation*}
-\Delta u+V(x) u=g(x, u)+f(x) \quad \text { in } \mathbb{R}^{N} \tag{1}
\end{equation*}
$$

where $N \geq 2, V$ is a potential function on $\mathbb{R}^{N}, g: \mathbb{R}^{N} \times \mathbb{R} \rightarrow \mathbb{R}$ is a superlinear, but subcritical, nonlinearity (namely, it satisfies the Ambrosetti-Rabinowitz condition) and $f$ : $\mathbb{R}^{N} \rightarrow \mathbb{R}$ is a given function.

When $f=0$ the study of equation (1) begins with Rabinowitz's paper [15] and then it has been carried out by several authors (cf. [6] and references therein): even if it has a variational structure, the main problem with classical variational tools is the lack of compactness. Thus, in [15] the existence of a nontrivial solution is shown by using the Mountain Pass Theorem but assuming that $V \in C^{1}\left(\mathbb{R}^{N}, \mathbb{R}\right)$ is positive and coercive; later on in [6], by means of the Symmetric Mountain Pass Theorem (see [1, Theorem 2.8]), Bartsch and Wang find infinitely many solutions if $g$ is odd in $u$ and $V$ is a positive continuous function such that

$$
\text { meas }\left(\left\{x \in \mathbb{R}^{N}: V(x) \leq b\right\}\right)<+\infty \quad \text { for all } b>0
$$

Motivated by the fact that on bounded domains, starting with the pioneer papers $[2,3$, $16,19]$, it is shown that multiplicity results may persist when the symmetry is destroyed by a perturbation term (see also $[10,11,22]$ ), we study (1) for $f \neq 0$. Our approach is based on the so-called Bolle's method (cf. [9, 10]) and on some ideas in [22]. We are only aware of a few previous contributions in this direction: indeed, in [17, Theorem 1.1] (see also [18]) it is proved a multiplicity result for a problem related to ours, provided that the eigenvalues of the involved Schrödinger operator have a suitable growth; on the other hand, in [4] a sharp result is obtained under radial assumptions.

[^0]Hereafter, in order to have a variational formulation of the problem and to overcome the lack of compactness, we assume the following conditions:
$\left(H_{1}\right)$ the potential $V \in L_{\mathrm{loc}}^{2}\left(\mathbb{R}^{N}\right)$ is such that

$$
\begin{equation*}
\underset{x \in \mathbb{R}^{N}}{\operatorname{ess} \inf } V(x)>0 \tag{2}
\end{equation*}
$$

and

$$
\lim _{|x| \rightarrow+\infty} \int_{B_{1}(x)} \frac{1}{V(y)} \mathrm{d} y=0
$$

where $B_{1}(x)=\left\{y \in \mathbb{R}^{N}:|x-y|<1\right\}$;
$\left(H_{2}\right) g: \mathbb{R}^{N} \times \mathbb{R} \rightarrow \mathbb{R}$ is a Carathéodory function (i.e., $g(\cdot, t)$ is measurable in $\mathbb{R}^{N}$ for all $t \in \mathbb{R}$ and $g(x, \cdot)$ is continuous in $\mathbb{R}$ for a.e. $\left.x \in \mathbb{R}^{N}\right)$ such that there exist $a_{1}, a_{2}>0$, $\mu>2$ and $\delta>0$ small enough (cf. Remark 3.2) satisfying
$\left(g_{1}\right)|g(x, s)| \leq a_{1}|s|^{p-1}+\delta|s|$ for a.e. $x \in \mathbb{R}^{N}, s \in \mathbb{R}$, with $\left.p \in\right] 2,2^{*}[$;
$\left(g_{2}\right) g(x, s) s \geq \mu G(x, s)>0$ for a.e. $x \in \mathbb{R}^{N}, s \in \mathbb{R} \backslash\{0\} ;$
$\left(g_{3}\right) G(x, s) \geq a_{2}|s|^{\mu}$ for a.e. $x \in \mathbb{R}^{N}, s \in \mathbb{R}$;
$\left(g_{4}\right) g(x, \cdot)$ is odd for a.e. $x \in \mathbb{R}^{N}$,
with $G(x, s):=\int_{0}^{s} g(x, t) \mathrm{d} t$.
Remark 1.1. Assumption $\left(g_{3}\right)$ is somehow related to $\left(g_{2}\right)$ : indeed, by $\left(g_{2}\right)$ and direct computations it follows that for any $\varepsilon>0$ there exists a constant $a_{\varepsilon}>0$ such that

$$
G(x, s) \geq a_{\varepsilon}|s|^{\mu} \quad \text { if }|s| \geq \varepsilon, \text { for a.e. } x \in \mathbb{R}^{N}
$$

In what follows by a solution we mean a weak solution; classical solutions are found when all the involved functions are smooth enough (e.g, cf. [6]).

Our main result is the following.
Theorem 1.2. Assume that $\left(H_{1}\right)-\left(H_{2}\right)$ hold. Then, for all $f \in L^{\frac{\mu}{\mu-1}}\left(\mathbb{R}^{N}\right)$ problem (1) has infinitely many solutions, provided that

$$
\begin{equation*}
\frac{\mu}{\mu-1}<\frac{4}{N(p-2)} \tag{3}
\end{equation*}
$$

Clearly, by $\left(g_{1}\right)$ and $\left(g_{3}\right)$ it follows that $\mu \leq p$. If $\mu=p$, in particular when $g(x, s)$ is exactly a pure power, condition (3) can be rewritten as follows.
Corollary 1.3. Assume that $\left(H_{1}\right)$ holds and $g(x, u)=|u|^{p-2} u$, with $\left.p \in\right] 2,2^{*}[$. Then, for all $f \in L^{\frac{p}{p-1}}\left(\mathbb{R}^{N}\right)$ problem (1) has infinitely many solutions, provided that $\left.p \in\right] 2, p_{N}[$, where $p_{N}:=\frac{N+2+\sqrt{N^{2}+4}}{N}$.

Condition $\left(H_{1}\right)$ on function $V$ is weaker than those used in [6, 15], as shown in [17, Proposition 3.1]. On the other hand, our set of conditions $\left(H_{2}\right)$ is similar to the analogous in [17], even if a comparison between [17, Corollary 1.6] and our Theorem 1.2 can be carried out only when the spectrum of the Schrödinger operator is known (cf. Proposition 3.1). For example, if $N=3$ and $V(x)=|x|^{2}$, the corresponding operator in $L^{2}\left(\mathbb{R}^{3}\right)$ admits the sequence of eigenvalues $\left(\lambda_{k}\right)_{k}$, with $\lambda_{k}=2 k+3$ (cf. [13, p. 514]). Taking the model nonlinearity, Corollary 1.3 gives infinitely many solutions for (1) if $p$ varies in the range ]2, $\frac{5+\sqrt{13}}{3}$ [, while the range obtained in [17, Corollary 1.6] is smaller, being ] $2, \frac{19+\sqrt{213}}{14}[$. As usual, for results concerning with problems with broken symmetry, Theorem 1.2 is far from being optimal, since we do not cover the entire subcritical range $] 2,2^{*}$. In spite of this, when dealing with radial assumptions and $N \geq 3$, one finds almost optimal results (cf. [4, 5] for unbounded domains and [11, 20, 21] for bounded ones).

The paper is organized as follows: in Section 2 we recall Bolle's method, then in Section 3 we introduce the variational setting of our problem and prove some technical results; finally, in Section 4 we prove Theorem 1.2.

Notations. Throughout this paper we denote by

- $2^{*}=\frac{2 N}{N-2}$ if $N \geq 3,2^{*}=+\infty$ otherwise;
- $s^{\prime}$ the conjugate exponent of $s \geq 1$, namely $s^{\prime}=\frac{s}{s-1}$ if $s>1$ and $s^{\prime}=+\infty$ if $s=1$;
- $|\cdot|_{s}$ the standard norm in the Lebesgue space $L^{s}\left(\mathbb{R}^{N}\right), 1 \leq s \leq+\infty$;
- $m^{*}(\bar{x}, \Psi)$ the large Morse index of a $C^{2}$ functional $\Psi$ at a critical point $\bar{x}$;
- $d_{j}, C_{j}$ positive real numbers, for any $j \in \mathbb{N}$.

2. Bolle's perturbation method. In this section we introduce the Bolle's perturbation method firstly stated in [9] but in the version presented in [10] and improved in [12], as the involved functionals are $C^{1}$ instead of $C^{2}$. The key point of this approach is dealing with a continuous path of functionals $\left(I_{\theta}\right)_{\theta \in[0,1]}$ which starts at a symmetric functional $I_{0}$ and ends at the "true" non-even functional $I_{1}$ associated to the given perturbed problem, so that the critical points of mini-max type of the symmetric map $I_{0}$ "shift" into critical points of $I_{1}$.

Throughout this section, let $\left(\mathcal{H},\|\cdot\|_{\mathcal{H}}\right)$ be a Hilbert space with dual $\left(\mathcal{H}^{\prime},\|\cdot\|_{\mathcal{H}^{\prime}}\right)$ and $I:(\theta, v) \in[0,1] \times \mathcal{H} \mapsto I(\theta, v) \in \mathbb{R}$ a $C^{1}$ functional. For simplicity, let us set $I_{\theta}=I(\theta, \cdot):$ $\mathcal{H} \rightarrow \mathbb{R}$ and $I_{\theta}^{\prime}(\cdot)=\frac{\partial I}{\partial v}(\theta, \cdot): \mathcal{H} \rightarrow \mathcal{H}^{\prime}$, for each $\theta \in[0,1]$. Assume that $\mathcal{H}$ can be decomposed so that $\mathcal{H}=H_{-} \oplus H_{+}$, with $\operatorname{dim}\left(H_{-}\right)<+\infty$, and $\left(e_{k}\right)_{k \geq 1}$ is an orthonormal basis of $H_{+}$. Setting

$$
H_{0}=H_{-}, \quad H_{k+1}=H_{k} \oplus \mathbb{R} e_{k+1} \text { if } k \in \mathbb{N}
$$

we have that $\left(H_{k}\right)_{k}$ is an increasing sequence of finite dimensional subspaces of $\mathcal{H}$.
Furthermore, we define

$$
\begin{equation*}
\Gamma=\left\{\gamma \in C(\mathcal{H}, \mathcal{H}): \gamma \text { is odd and } \exists \rho>0 \text { s.t. } \gamma(v)=v \text { if }\|v\|_{\mathcal{H}} \geq \rho\right\} \tag{4}
\end{equation*}
$$

and

$$
c_{k}=\inf _{\gamma \in \Gamma} \sup _{v \in H_{k}} I_{0}(\gamma(v)) .
$$

Let us assume that:
$\left(A_{1}\right) I$ satisfies the following variant of the Palais-Smale condition: each sequence $\left(\left(\theta_{n}, v_{n}\right)\right)_{n} \subset[0,1] \times \mathcal{H}$ such that

$$
\begin{equation*}
\left(I\left(\theta_{n}, v_{n}\right)\right)_{n} \text { is bounded and } \lim _{n \rightarrow+\infty}\left\|I_{\theta_{n}}^{\prime}\left(v_{n}\right)\right\|_{\mathcal{H}^{\prime}}=0 \tag{5}
\end{equation*}
$$

converges, up to subsequences;
$\left(A_{2}\right)$ for all $b>0$ there exists $C_{b}>0$ such that, if $(\theta, v) \in[0,1] \times \mathcal{H}$, then

$$
\left|I_{\theta}(v)\right| \leq b \Longrightarrow\left|\frac{\partial I}{\partial \theta}(\theta, v)\right| \leq C_{b}\left(\left\|I_{\theta}^{\prime}(v)\right\|_{\mathcal{H}^{\prime}}+1\right)\left(\|v\|_{\mathcal{H}}+1\right)
$$

$\left(A_{3}\right)$ there exist two continuous maps $\eta_{1}, \eta_{2}:[0,1] \times \mathbb{R} \rightarrow \mathbb{R}$, with $\eta_{1}(\theta, \cdot) \leq \eta_{2}(\theta, \cdot)$ for all $\theta \in[0,1]$, which are Lipschitz continuous with respect to the second variable and such that, if $(\theta, v) \in[0,1] \times \mathcal{H}$, then

$$
\begin{equation*}
I_{\theta}^{\prime}(v)=0 \Longrightarrow \eta_{1}\left(\theta, I_{\theta}(v)\right) \leq \frac{\partial I}{\partial \theta}(\theta, v) \leq \eta_{2}\left(\theta, I_{\theta}(v)\right) \tag{6}
\end{equation*}
$$

$\left(A_{4}\right) I_{0}$ is even and for each finite dimensional subspace $\mathcal{V}$ of $\mathcal{H}$ it results

$$
\lim _{\substack{v \in \mathcal{V} \\\|v\|_{\mathcal{H}} \rightarrow+\infty}} \sup _{\theta \in[0,1]} I_{\theta}(v)=-\infty
$$

Now, for $i \in\{1,2\}$, let $\psi_{i}:[0,1] \times \mathbb{R} \rightarrow \mathbb{R}$ be the (unique, global) solution of the problem

$$
\left\{\begin{array}{l}
\frac{\partial \psi_{i}}{\partial \theta}(\theta, s)=\eta_{i}\left(\theta, \psi_{i}(\theta, s)\right) \\
\psi_{i}(0, s)=s
\end{array}\right.
$$

with $\eta_{i}$ as in $\left(A_{3}\right)$.

Note that $\psi_{i}(\theta, \cdot)$ is continuous, non-decreasing on $\mathbb{R}, i \in\{1,2\}$, and $\psi_{1}(\theta, \cdot) \leq \psi_{2}(\theta, \cdot)$. Moreover, we set

$$
\bar{\eta}_{1}(s)=\max _{\theta \in[0,1]}\left|\eta_{1}(\theta, s)\right|, \quad \bar{\eta}_{2}(s)=\max _{\theta \in[0,1]}\left|\eta_{2}(\theta, s)\right| .
$$

The following result holds (cf. [9, Theorem 3], [10, Theorem 2.2] and [12, Section 2]).
Theorem 2.1. Let $I:[0,1] \times \mathcal{H} \rightarrow \mathbb{R}$ be a $C^{1}$ path of functionals satisfying assumptions $\left(A_{1}\right)-\left(A_{4}\right)$. Then, there exists $C>0$ such that for all $k \in \mathbb{N}$ it results:
(a) either $I_{1}$ has a critical level $\widetilde{c}_{k}$ with $\psi_{2}\left(1, c_{k}\right)<\psi_{1}\left(1, c_{k+1}\right) \leq \widetilde{c}_{k}$,
(b) or $c_{k+1}-c_{k} \leq C\left(\bar{\eta}_{1}\left(c_{k+1}\right)+\bar{\eta}_{2}\left(c_{k}\right)+1\right)$.

Remark 2.2. We point out that, if $\eta_{2} \geq 0$ in $[0,1] \times \mathbb{R}$, the function $\psi_{2}(\cdot, s)$ is non-decreasing on $[0,1]$. Hence, $c_{k} \leq \widetilde{c}_{k}$ for all $c_{k}$ verifying case $(a)$.
3. Variational set-up. In this section we present the functional framework of our problem. Firstly, by (2) it makes sense to consider the weighted Sobolev space

$$
E_{V}:=H_{V}^{1}\left(\mathbb{R}^{N}\right)=\left\{u \in H^{1}\left(\mathbb{R}^{N}\right): \int_{\mathbb{R}^{N}} V(x) u^{2} \mathrm{~d} x<+\infty\right\}
$$

endowed with the norm

$$
\|u\|_{V}=\left(\int_{\mathbb{R}^{N}}\left(|\nabla u|^{2}+V(x) u^{2}\right) \mathrm{d} x\right)^{\frac{1}{2}}
$$

The following proposition (cf. [7, Theorems 3.1 and 4.1] and [17, Proposition 3.3]) is crucial to overcome the lack of compactness.

Proposition 3.1. Let $V: \mathbb{R}^{N} \rightarrow \mathbb{R}$ be such that $\left(H_{1}\right)$ holds. Then, for all $s \in\left[2,2^{*}[\right.$ it is $E_{V} \hookrightarrow \hookrightarrow L^{s}\left(\mathbb{R}^{N}\right)$, i.e. the embedding of $\left(E_{V},\|\cdot\|_{V}\right)$ in $\left(L^{s}\left(\mathbb{R}^{N}\right),|\cdot|_{s}\right)$ is compact. Moreover, the linear Schrödinger operator

$$
u \in C_{0}^{\infty}\left(\mathbb{R}^{N}\right) \mapsto-\Delta u+V(x) u \in L^{2}\left(\mathbb{R}^{N}\right)
$$

is essentially self-adjoint, the spectrum of its self-adjoint extension is an increasing sequence $\left(\lambda_{n}\right)_{n}$ of eigenvalues of finite multiplicity and

$$
L^{2}\left(\mathbb{R}^{N}\right)=\sum_{n=1}^{+\infty} M_{n} \quad \text { with } M_{n} \perp M_{m} \text { for } n \neq m
$$

where $M_{n}$ denotes the eigenspace corresponding to $\lambda_{n}$ for every $n \in \mathbb{N}$.
Remark 3.2. From $\left(g_{1}\right)$ and $\left(g_{3}\right)$ it follows $2<\mu \leq p<2^{*}$; hence, Proposition 3.1 implies $E_{V} \hookrightarrow \hookrightarrow L^{\mu}\left(\mathbb{R}^{N}\right)$ and $E_{V} \hookrightarrow \hookrightarrow L^{p}\left(\mathbb{R}^{N}\right)$. Moreover, as $E_{V} \hookrightarrow \hookrightarrow L^{2}\left(\mathbb{R}^{N}\right)$, for further use by $\alpha$ we denote the best embedding constant and in assumption $\left(g_{1}\right)$ we choose $\delta$ such that $\delta<\frac{1}{\alpha^{2}}$.

As direct consequence of Proposition 3.1 and [23, Theorem 1.22] we can state the following lemma.

Lemma 3.3. Assume that $\left(H_{1}\right)$ and $\left(g_{1}\right)$ hold. Then, setting $\Phi: E_{V} \rightarrow \mathbb{R}$ as

$$
\Phi(u)=\int_{\mathbb{R}^{N}} G(x, u) \mathrm{d} x \quad \text { for all } u \in E_{V}
$$

it results that $\Phi \in C^{1}\left(E_{V}, \mathbb{R}\right)$ with

$$
\Phi^{\prime}(u)[\varphi]=\int_{\mathbb{R}^{N}} g(x, u) \varphi \mathrm{d} x \quad \text { for all } \varphi \in E_{V}
$$

Moreover, $\Phi^{\prime}: E_{V} \rightarrow\left(E_{V}\right)^{\prime}$ is compact.

By Lemma 3.3 and standard arguments, the weak solutions of (1) are the critical points of the $C^{1}$ functional on $E_{V}$

$$
I_{1}(u)=\frac{1}{2} \int_{\mathbb{R}^{N}}\left(|\nabla u|^{2}+V(x) u^{2}\right) \mathrm{d} x-\int_{\mathbb{R}^{N}} G(x, u) \mathrm{d} x-\int_{\mathbb{R}^{N}} f u \mathrm{~d} x
$$

with

$$
I_{1}^{\prime}(u)[\varphi]=\int_{\mathbb{R}^{N}} \nabla u \cdot \nabla \varphi \mathrm{~d} x+\int_{\mathbb{R}^{N}} V(x) u \varphi \mathrm{~d} x-\int_{\mathbb{R}^{N}} g(x, u) \varphi \mathrm{d} x-\int_{\mathbb{R}^{N}} f \varphi \mathrm{~d} x
$$

for all $u, \varphi \in E_{V}$.
In order to apply the Bolle's perturbation method, we define the path of functionals $I:[0,1] \times E_{V} \rightarrow \mathbb{R}$ as follows:

$$
\begin{equation*}
I(\theta, u)=\frac{1}{2} \int_{\mathbb{R}^{N}}\left(|\nabla u|^{2}+V(x) u^{2}\right) \mathrm{d} x-\int_{\mathbb{R}^{N}} G(x, u) \mathrm{d} x-\theta \int_{\mathbb{R}^{N}} f v \mathrm{~d} x . \tag{7}
\end{equation*}
$$

Now we verify that, under our main assumptions, the path introduced in (7) satisfies conditions $\left(A_{1}\right)-\left(A_{4}\right)$ in Section 2.

Proposition 3.4. Assume that $\left(H_{1}\right)-\left(H_{2}\right)$ hold. Then, the family $\left(I_{\theta}\right)_{\theta \in[0,1]}$ verifies $\left(A_{1}\right)-\left(A_{4}\right)$.

Proof. The proof is organized in four steps.
Step 1. Let $\left(\left(\theta_{n}, u_{n}\right)\right)_{n} \subset[0,1] \times E_{V}$ be a sequence such that (5) holds; hence,

$$
I_{\theta_{n}}\left(u_{n}\right)=\frac{1}{2}\left\|u_{n}\right\|_{V}^{2}-\int_{\mathbb{R}^{N}} G\left(x, u_{n}\right) \mathrm{d} x-\theta_{n} \int_{\mathbb{R}^{N}} f u_{n} \mathrm{~d} x \leq d_{1}
$$

and

$$
\begin{aligned}
\left|I_{\theta_{n}}^{\prime}\left(u_{n}\right)\left[u_{n}\right]\right| & =\left|\left\|u_{n}\right\|_{V}^{2}-\int_{\mathbb{R}^{N}} g\left(x, u_{n}\right) u_{n} \mathrm{~d} x-\theta_{n} \int_{\mathbb{R}^{N}} f u_{n} \mathrm{~d} x\right| \\
& \leq \varepsilon_{n}\left\|u_{n}\right\|_{V}
\end{aligned}
$$

where $\varepsilon_{n} \searrow 0$ as $n \rightarrow+\infty$. Therefore, by $\left(g_{2}\right)$, Remark 3.2 and the Hölder inequality, it follows that

$$
\begin{aligned}
d_{1}+\frac{\varepsilon_{n}}{\mu}\left\|u_{n}\right\|_{V} & \geq I_{\theta_{n}}\left(u_{n}\right)-\frac{1}{\mu} I_{\theta_{n}}^{\prime}\left(u_{n}\right)\left[u_{n}\right] \\
& \geq\left(\frac{1}{2}-\frac{1}{\mu}\right)\left\|u_{n}\right\|_{V}^{2}-\left(1-\frac{1}{\mu}\right) \theta_{n} \int_{\mathbb{R}^{N}} f u_{n} \mathrm{~d} x \\
& \geq\left(\frac{1}{2}-\frac{1}{\mu}\right)\left\|u_{n}\right\|_{V}^{2}-\left(1-\frac{1}{\mu}\right) \theta_{n}|f|_{\mu^{\prime}}\left|u_{n}\right|_{\mu} \\
& \geq\left(\frac{1}{2}-\frac{1}{\mu}\right)\left\|u_{n}\right\|_{V}^{2}-d_{2}\left\|u_{n}\right\|_{V}
\end{aligned}
$$

thus the sequence $\left(u_{n}\right)_{n}$ is bounded in $E_{V}$ and $\left(A_{1}\right)$ follows by Proposition 3.1 and standard arguments.
Step 2. Since

$$
\frac{\partial I}{\partial \theta}(\theta, u)=-\int_{\mathbb{R}^{N}} f u \mathrm{~d} x
$$

by using again the Hölder inequality and Remark 3.2, we get that

$$
\left|\frac{\partial I}{\partial \theta}(\theta, u)\right| \leq d_{3}\|u\|_{V}
$$

hence $\left(A_{2}\right)$ holds.
Step 3. Taking $(\theta, u) \in[0,1] \times E_{V}$ such that $I_{\theta}^{\prime}(u)=0$, we have that

$$
\begin{aligned}
I_{\theta}(u) & =I_{\theta}(u)-\frac{1}{2} I_{\theta}^{\prime}(u)[u] \\
& =\frac{1}{2} \int_{\mathbb{R}^{N}} g(x, u) u \mathrm{~d} x-\int_{\mathbb{R}^{N}} G(x, u) \mathrm{d} x-\frac{\theta}{2} \int_{\mathbb{R}^{N}} f u \mathrm{~d} x
\end{aligned}
$$

then by the Hölder inequality, using $\left(g_{2}\right)$ and $\left(g_{3}\right)$ respectively, we get that

$$
\begin{aligned}
I_{\theta}(u) & \geq\left(\frac{\mu}{2}-1\right) \int_{\mathbb{R}^{N}} G(x, u) \mathrm{d} x-\frac{\theta}{2}|f|_{\mu^{\prime}}|u|_{\mu} \\
& \geq d_{4}|u|_{\mu}^{\mu}-d_{5}|u|_{\mu}
\end{aligned}
$$

Since $\mu>2$, direct computations and elementary inequalities give

$$
|u|_{\mu} \leq d_{6}\left(I_{\theta}^{2}(u)+1\right)^{\frac{1}{2 \mu}}
$$

Hence,

$$
\left|\frac{\partial I}{\partial \theta}(\theta, u)\right| \leq|f|_{\mu^{\prime}}|u|_{\mu} \leq C_{1}\left(I_{\theta}^{2}(u)+1\right)^{\frac{1}{2 \mu}}
$$

and inequality (6) holds with $\eta_{1}, \eta_{2}:[0,1] \times \mathbb{R} \rightarrow \mathbb{R}$ defined by

$$
\begin{equation*}
-\eta_{1}(\theta, s)=\eta_{2}(\theta, s)=C_{1}\left(s^{2}+1\right)^{\frac{1}{2 \mu}} \tag{8}
\end{equation*}
$$

therefore $\left(A_{3}\right)$ is proved.
Step 4. Finally, let us remark that by $\left(g_{4}\right)$ the functional $I_{0}$ is even on $E_{V}$ (see (7)); moreover, by $\left(g_{3}\right)$ and standard arguments we have that

$$
I(\theta, u) \leq \frac{1}{2}\|u\|_{V}^{2}-a_{2}|u|_{\mu}^{\mu}+|f|_{\mu^{\prime}}|u|_{\mu}
$$

Hence, taking any finite dimensional subspace $\mathcal{V}$ of $E_{V}$, as $\mu>2$ and all norms are equivalent on $\mathcal{V}$, property $\left(A_{4}\right)$ follows.
4. Proof of the main results. Our aim is to apply Theorem 2.1, therefore let us introduce a suitable class of mini- max values for the even functional $I_{0}$.

Denoting by $\left(e_{k}\right)_{k}$ the basis of eigenfunctions in $E_{V}$ found in Proposition 3.1, for any $k \geq 1$ let us set

$$
\begin{equation*}
E_{k}=\operatorname{span}\left\{e_{1}, \ldots, e_{k}\right\}, \quad E_{k}^{\perp}=\overline{\operatorname{span}\left\{e_{k+1}, \ldots\right\}} \tag{9}
\end{equation*}
$$

and

$$
\begin{equation*}
c_{k}=\inf _{\gamma \in \Gamma} \sup _{u \in E_{k}} I_{0}(\gamma(u)), \tag{10}
\end{equation*}
$$

where $\Gamma$ is as in (4) with $\mathcal{H}=E_{V}$.
In order to establish a lower estimate for the sequence $\left(c_{k}\right)_{k}$, we recall two lemmas, proved in [14, Corollary 2] and [4, Lemma 4.2] respectively.

Taking any $W: \mathbb{R}^{N} \rightarrow \mathbb{R}$, we denote by $N_{-}(-\Delta+W(x))$ the number of the negative eigenvalues of the operator $-\Delta+W(x)$ and set $W_{-}(x)=\min \{W(x), 0\}$.
Lemma 4.1. Let $N \geq 3$ and $W \in L^{\frac{N}{2}}\left(\mathbb{R}^{N}\right)$. Then, there exists $\bar{C}_{N}>0$ such that

$$
N_{-}(-\Delta+W(x)) \leq \bar{C}_{N}\left|W_{-}\right|_{\frac{N}{2}}^{\frac{N}{2}}
$$

If $N=2$ and $W \in L^{1+\varepsilon}\left(\mathbb{R}^{N}\right)$ for some $\varepsilon>0$, then there exists $\bar{C}_{\varepsilon}>0$ such that

$$
N_{-}(-\Delta+W(x)) \leq \bar{C}_{\varepsilon}\left|W_{-}\right|_{1+\varepsilon}^{1+\varepsilon}
$$

Lemma 4.2. Let $\left.p \in] 2,2+\frac{4}{N}\right]$. Then, for some $\left.\bar{p} \in\right] 2+\frac{4}{N}, 2^{*}[$, for all $\varepsilon>0$ there exists $D_{\varepsilon}>0$ such that

$$
\int_{\mathbb{R}^{N}}|u|^{p} \mathrm{~d} x \leq \varepsilon \int_{\mathbb{R}^{N}} u^{2} \mathrm{~d} x+D_{\varepsilon} \int_{\mathbb{R}^{N}}|u|^{\bar{p}} \mathrm{~d} x \quad \text { for all } u \in E_{V}
$$

Now, we are ready to prove our main result.
Proof of Theorem 1.2. By Proposition 3.4, Theorem 2.1 applies, so the proof of our result is complete if we rule out case ( $b$ ) for $k$ large enough or better, as by (8) condition (b) implies

$$
\begin{equation*}
c_{k+1}-c_{k} \leq C_{2}\left(\left(c_{k}\right)^{\frac{1}{\mu}}+\left(c_{k+1}\right)^{\frac{1}{\mu}}+1\right) \tag{11}
\end{equation*}
$$

with $c_{k}$ as in (10), it is enough to prove that (11) cannot hold for $k$ large enough.
In fact, if we assume that (11) holds for all $k \geq k_{0}$ for some $k_{0} \geq 1$, by [2, Lemma 5.3] it follows that there exist $\bar{C}>0$ and $\bar{k} \in \mathbb{N}$ such that

$$
\begin{equation*}
c_{k} \leq \bar{C} k^{\frac{\mu}{\mu-1}} \quad \text { for all } k \geq \bar{k} \tag{12}
\end{equation*}
$$

On the other hand, by $\left(g_{1}\right)$ it follows that

$$
|G(x, u)| \leq \frac{a_{1}}{p}|s|^{p}+\frac{\delta}{2}|s|^{2} \quad \text { for a.e. } x \in \mathbb{R}^{N}, s \in \mathbb{R}
$$

whence, by Remark 3.2 there exists $C^{*}>0$ such that

$$
I_{0}(u) \geq C^{*}\left(\frac{1}{2}\|u\|_{V}^{2}-C_{3}|u|_{p}^{p}\right) \quad \text { for all } u \in E_{V}
$$

From now on, we deal with the case $N \geq 3$, since the case $N=2$ follows by slight modifications. We claim that for any $p \in] 2,2^{*}[$ it is

$$
\begin{equation*}
c_{k} \geq C_{4} k^{\frac{4}{N(p-2)}} \quad \text { for all } k \geq 1 \tag{13}
\end{equation*}
$$

To this aim, two different cases occur.
Case 1. Let $2+\frac{4}{N} \leq p<2^{*}$. Setting

$$
K(u)=\frac{1}{2}\|u\|_{V}^{2}-C_{3}|u|_{p}^{p}
$$

and

$$
b_{k}=\inf _{\gamma \in \Gamma} \sup _{u \in E_{k}} C^{*} K(\gamma(u)),
$$

we have that

$$
\begin{equation*}
c_{k} \geq b_{k} \tag{14}
\end{equation*}
$$

Now, $\left[22\right.$, Theorem B] implies that for all $k \in \mathbb{N}$ there exists $u_{k} \in E_{V}$ such that

$$
\begin{equation*}
K^{\prime}\left(u_{k}\right)=0 \quad \text { and } \quad K\left(u_{k}\right) \leq b_{k}, \tag{15}
\end{equation*}
$$

with $m^{*}\left(u_{k}, K\right) \geq k$, i.e., the operator

$$
K^{\prime \prime}\left(u_{k}\right)=-\Delta+V(x)-C_{3} p(p-1)\left|u_{k}\right|^{p-2}
$$

has at least $k$ non-positive eigenvalues. Therefore, by [8, Proposition S1.3.1] and Lemma 4.1 with $W(x)=-C_{3} p(p-1)\left|u_{k}\right|^{p-2}$ we infer that

$$
k \leq N_{-}\left(K^{\prime \prime}\left(u_{k}\right)\right) \leq N_{-}\left(-\Delta-C_{3} p(p-1)\left|u_{k}\right|^{p-2}\right) \leq C_{5}\left|u_{k}\right|_{(p-2) \frac{N}{2}}^{(p-2) \frac{N}{2}}
$$

In this case, we have $(p-2) \frac{N}{2} \in\left[2,2^{*}[\right.$, then by Proposition 3.1 we get

$$
k \leq C_{6}\left\|u_{k}\right\|_{V}^{(p-2) \frac{N}{2}}
$$

As (15) implies $K^{\prime}\left(u_{k}\right)\left[u_{k}\right]=0$, then

$$
\begin{equation*}
\left\|u_{k}\right\|_{V}^{2}=C_{3} p\left|u_{k}\right|_{p}^{p} \tag{16}
\end{equation*}
$$

hence (13) follows by (14) - (16).
Case 2. Let $2<p<2+\frac{4}{N}$. By Lemma 4.2, for a suitable $\left.\bar{p} \in\right] 2+\frac{4}{N}, 2^{*}[$ and $\varepsilon>0$ small enough, there exist $b_{\varepsilon}, d_{\varepsilon}>0$ such that, setting

$$
K_{\varepsilon}(u)=b_{\varepsilon}\|u\|_{V}^{2}-d_{\varepsilon}|u|_{\bar{p}}^{\bar{p}},
$$

it results

$$
I_{0}(u) \geq K_{\varepsilon}(u) \quad \text { for all } u \in E_{V} .
$$

Then, let us define

$$
c_{k}^{\varepsilon}=\inf _{\gamma \in \Gamma} \sup _{u \in E_{k}} K_{\varepsilon}(\gamma(u))
$$

where $\Gamma$ is as in (4) with $\mathcal{H}=E_{V}$ and $E_{k}$ is as in (9). Plainly, $c_{k} \geq c_{k}^{\varepsilon}$. By applying the arguments developed in Case 1, but with $p$ replaced by $\bar{p}$ and $K$ by $K_{\varepsilon}$, also in this case (13) holds.

At last, by (3) inequality (13) yields to a contradiction with (12); therefore condition (a) in Theorem 2.1 holds for infinitely many $k \in \mathbb{N}$ and by Remark 2.2 the proof is complete.

Proof of Corollary 1.3. Proposition 3.4 follows by simpler arguments, with $\eta_{1}, \eta_{2}:[0,1] \times$ $\mathbb{R} \rightarrow \mathbb{R}$ defined by

$$
-\eta_{1}(\theta, s)=\eta_{2}(\theta, s)=C_{7}\left(s^{2}+1\right)^{\frac{1}{2 p}}
$$

this implies that (12) is now replaced by

$$
c_{k} \leq C_{8} k^{\frac{p}{p-1}} \quad \text { for } k \text { large enough. }
$$

Then, we can reason as in the proof of Theorem 1.2 , working directly on $I_{0}$.

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