## Università degli studi di Udine

## The homogeneous Hénon-Lane-Emden system

## Original

Availability:
This version is available http://hdl.handle.net/11390/1069605 since 2021-03-15T15:55:02Z

Publisher:

Published
DOI:10.1007/s00030-015-0330-5

Terms of use:
The institutional repository of the University of Udine (http://air.uniud.it) is provided by ARIC services. The aim is to enable open access to all the world.

Publisher copyright
(Article begins on next page)

# Nonlinear Differential Equations and Applications NoDea The homogeneous Hénon-Lane-Emden system <br> --Manuscript Draft-- 

| Manuscript Number: | NDEA-D-15-00054R1 |
| :--- | :--- |
| Full Title: | The homogeneous Hénon-Lane-Emden system |
| Article Type: | Original research |
| Keywords: | Eigenvalue problem; Hénon-Lane-Emden system; quasilinear elliptic system |
| Corresponding Author: | Roberta Musina, Ph.D. |
| Corresponding Author Secondary <br> Information: |  |
| Corresponding Author's Institution: |  |
| Corresponding Author's Secondary |  |
| Institution: | Andrea Carioli |
| First Author: | Andrea Carioli |
| First Author Secondary Information: | Roberta Musina, Ph.D. |
| Order of Authors: | We use variational methods to study the existence of a principal eigenvalue for the |
| homogeneous Hénon-Lane-Emden system on a bounded domain. Then we provide a |  |
| detailed insight into the problem in the linear case. |  |
| Order of Authors Secondary Information: | Points 1-2-3-4: we corrected the typos/inaccuracies <br> Points 5-6-7: we added Remark 4.1. |
| Abstract: |  |
| Response to Reviewers: |  |

# The homogeneous Hénon-Lane-Emden system 

Andrea Carioli and Roberta Musina


#### Abstract

We use variational methods to study the existence of a principal eigenvalue for the homogeneous Hénon-Lane-Emden system on a bounded domain. Then we provide a detailed insight into the problem in the linear case.


Mathematics Subject Classification (2010). Primary 35J47; Secondary $35 J 35$.
Keywords. Eigenvalue problem, Hénon-Lane-Emden system, quasilinear elliptic system.

## 1. Introduction

The Hénon-Lane-Emden system

$$
\begin{cases}-\Delta u=|x|^{a}|v|^{q-2} v & \text { in } \Omega  \tag{1.1}\\ -\Delta v=|x|^{b}|u|^{p-2} u & \text { in } \Omega \\ u=0=v & \text { on } \partial \Omega\end{cases}
$$

includes the second and fourth order Lane-Emden equations and the Hénon equation in astrophysics. Here $\Omega$ is a domain in $\mathbb{R}^{n}$ containing the origin, $p, q \in(1, \infty)$, and the weights are locally integrable, that is, $a, b>-n$.

Since the celebrated papers [15] by P.L. Lions and [17] by Mitidieri, where $a=b=0$ and $\Omega=\mathbb{R}^{n}$ are assumed, large efforts have been made in investigating (1.1) and related problems. It is difficult to give a complete list of references on this topic. We limit ourselves to cite $[2,3,4,5,8,11,12,13$, $18,22,25,27,28,29,30]$ and the references therein.

Most of the above mentioned papers require $n \geq 3$, deal with the socalled anticoercive case $(p-1)(q-1)>1$, and underline the role of the

[^0]"critical hyperbola"
$$
\frac{a+n}{q}+\frac{b+n}{p}=n-2
$$
in existence and nonexistence phenomena, as it separates the "subcritical case"
\[

$$
\begin{equation*}
\frac{a+n}{q}+\frac{b+n}{p}>n-2 \tag{1.2}
\end{equation*}
$$

\]

from the supercritical one.
Formally letting $q \searrow p^{\prime}=\frac{p}{p-1}$ in (1.1), and taking the homogeneities involved into account, one gets in the limit the eigenvalue problem

$$
\begin{cases}-\Delta u=\lambda_{1}|x|^{a}|v|^{p^{\prime}-2} v & \text { in } \Omega  \tag{P}\\ -\Delta v=\lambda_{2}|x|^{b}|u|^{p-2} u & \text { in } \Omega \\ u=0=v & \text { on } \partial \Omega\end{cases}
$$

and (1.2) reduces to

$$
\begin{equation*}
\frac{a}{p^{\prime}}+\frac{b}{p}+2>0 \tag{1.3}
\end{equation*}
$$

In the present paper we focus our attention on $\operatorname{Problem}(\mathcal{P})$. We emphasize the fact that we include the lower dimensional cases $n=1,2$, that actually present some remarkable peculiarities.

Nonexistence results have been obtained in [1], see also the recent papers [8, 10] for the case $\Omega=\mathbb{R}^{n}$. Montenegro [20] used degree theory to face Problem $(\mathcal{P})$ in a more general setting that includes non-self-adjoint elliptic operators. We adopt a variational approach that allows us to weaken the integrability assumptions on the coefficients from Montenegro's $L^{n}(\Omega)$ to $L^{1}(\Omega)$.

We look for finite energy solutions and for a principal eigenvalue to $(\mathcal{P})$, according to the next definitions.
Definition 1.1. The pair $(u, v)$ is a finite-energy solution to $(\mathcal{P})$ if:

- $u, v \in W^{2,1}(\Omega) \cap W_{0}^{1,1}(\Omega)$;
- $u \in L^{p}\left(\Omega,|x|^{b} d x\right), v \in L^{p^{\prime}}\left(\Omega,|x|^{a} d x\right)$, that is,

$$
\begin{equation*}
\int_{\Omega}|x|^{b}|u|^{p} d x<\infty, \quad \int_{\Omega}|x|^{a}|v|^{p^{\prime}} d x<\infty \tag{1.4}
\end{equation*}
$$

- $u, v$ are weak solutions to the elliptic equations in $(\mathcal{P})$. That is,
$\int_{\Omega} \nabla u \cdot \nabla \varphi d x=\lambda_{1} \int_{\Omega}|x|^{a}|v|^{p^{\prime}-2} v \varphi d x, \quad \int_{\Omega} \nabla v \cdot \nabla \varphi d x=\lambda_{2} \int_{\Omega}|x|^{a}|u|^{p-2} u \varphi d x$
for any test function $\varphi \in C_{c}^{\infty}(\Omega)$.
Definition 1.2. A real number $\mu$ is a principal eigenvalue for $(\mathcal{P})$ if for any pair of real numbers $\left(\lambda_{1}, \lambda_{2}\right)$ satisfying

$$
\begin{equation*}
\left|\lambda_{1}\right|^{p-1} \lambda_{1}\left|\lambda_{2}\right|^{p^{\prime}-1} \lambda_{2}=\mu^{p^{\prime}} \tag{1.5}
\end{equation*}
$$

Problem $(\mathcal{P})$ has a finite-energy solution $(u, v)$ such that $u, v>0$ in $\Omega$.

In Section 3 we prove the following result.
Theorem 1.3. Let $a, b>-n$ and let $\Omega \ni 0$ be a bounded and smooth domain in $\mathbb{R}^{n}$. If (1.3) holds, then Problem $(\mathcal{P})$ has a positive principal eigenvalue.
Notice that (1.3) is automatically satisfied if $n=1,2$ and $a, b>-n$. If $n \geq 3$ then assumption (1.3) can not be improved, see [7].

Our approach is based on the formal equivalence, already noticed for instance by Wang [31] and Calanchi-Ruf [6] in the anticoercive case, between $(\mathcal{P})$ and the fourth order eigenvalue problem

$$
\begin{cases}\Delta\left(|x|^{-a(p-1)}|\Delta u|^{p-2} \Delta u\right)=\mu|x|^{b}|u|^{p-2} u & \text { in } \Omega  \tag{1.6}\\ u=\Delta u=0 & \text { on } \partial \Omega\end{cases}
$$

where $\mu, \lambda_{1}$ and $\lambda_{2}$ satisfy (1.5). We will look for nontrivial solutions to (1.6) as constrained critical points for the energy
$\mathcal{E}(u)=\int_{\Omega}|x|^{-a(p-1)}|\Delta u|^{p} d x$ on the constraint $M=\left\{\int_{\Omega}|x|^{b}|u|^{p} d x=1\right\}$, in a suitably defined space $W_{N}^{2, p}\left(\Omega,|x|^{-a(p-1)} d x\right)$ of functions such that $u=0$ on $\partial \Omega$.

In spite of the apparent plainness of this program, its rigorous implementation needs a good understanding of some non trivial facts.

First of all one has to prove appropriate integral inequalities for smooth functions vanishing at the boundary of $\Omega$. This will be done in Lemma 2.4. If $n \geq 3$ we take advantage of the weighted Rellich-type inequality in [19] and [21, Lemma 2.14]. The lower dimensional cases $n=1,2$ require an ad hoc argument and a preliminary result, that can be found in the appendix.

Secondly, one needs to detect the "right" function space. The weighted space $W_{N}^{2, p}\left(\Omega,|x|^{-a(p-1)} d x\right)$ has to be "small enough" to be compactly embedded in $L^{p}\left(\Omega ;|x|^{b} d x\right)$. But the equivalence between weak solutions to (1.6) and finite-energy solutions to $(\mathcal{P})$ only holds if $W_{N}^{2, p}\left(\Omega,|x|^{-a(p-1)} d x\right)$ is "large enough".

The convenient definition of $W_{N}^{2, p}\left(\Omega,|x|^{-a(p-1)} d x\right)$ and details are given in Section 2. The above mentioned equivalence and Theorem 1.3 are proved in Section 3.
The last part of the paper is focused on the linear case $p=2$, so that (1.3) becomes

$$
\begin{equation*}
a+b+4>0 \tag{1.7}
\end{equation*}
$$

In Section 4 we prove that the linear system

$$
\begin{cases}-\Delta u=\lambda_{1}|x|^{a} v & \text { in } \Omega  \tag{1.8}\\ -\Delta v=\lambda_{2}|x|^{b} u & \text { in } \Omega \\ u=0=v & \text { on } \partial \Omega\end{cases}
$$

has a unique and simple principal eigenvalue $\mu_{1}>0$, and a discrete spectrum $\left\{\mu_{k}\right\}_{k \in \mathbb{N}}$. More precisely, the following facts hold.

Theorem 1.4. Let $a, b>-n$ and let $\Omega$ be a bounded and smooth domain in $\mathbb{R}^{n}$. If $n \geq 3$ assume also that (1.7) holds.
(a) There exists an increasing, unbounded sequence of eigenvalues $\left\{\mu_{k}\right\}_{k \in \mathbb{N}}$ such that Problem (1.8) has a nontrivial and finite-energy solution ( $u, v$ ) if and only if $\lambda_{1} \lambda_{2}=\mu_{k}$ for some integer $k \geq 1$.
(b) The first eigenvalue $\mu_{1}$ is the unique principal eigenvalue. In addition $\mu_{1}$ is simple, that is, if $(u, v)$ and $(\tilde{u}, \tilde{v})$ solve (1.8) and $\lambda_{1} \lambda_{2}=\mu_{1}$, then $\tilde{u}=\alpha u$ and $\tilde{v}=\beta v$ for some $\alpha, \beta \in \mathbb{R}$.

## 2. The functional setting

In this section we introduce and study certain second order weighted Sobolev spaces with Navier boundary conditions that are suitable for studying (1.6) via variational methods.

To simplify notation we set $s=a(p-1)$. Thus, from now on we assume that $s, b$ are given exponents such that

$$
s>n-n p, \quad b>-n
$$

even if not explicitly stated. In addition, $\Omega \subset \mathbb{R}^{n}$ will always denote a bounded and smooth domain. We denote by $c$ any universal positive constant.

We introduce the function space

$$
C_{N}^{2}(\bar{\Omega}):=\left\{u \in C^{2}(\bar{\Omega}) \mid u=0 \text { on } \partial \Omega\right\}
$$

Let $W_{N}^{2, p}\left(\Omega,|x|^{-s} d x\right)$ be the reflexive Banach space defined as the completion of the set

$$
D_{0}:=\left\{u \in C_{N}^{2}(\bar{\Omega}) \mid \Delta u \equiv 0 \text { on a neighborhood of the origin }\right\}
$$

with respect to the uniformly convex norm

$$
\|u\|_{s} \equiv\|u\|_{p, s}:=\left(\int_{\Omega}|x|^{-s}|\Delta u|^{p} d x\right)^{\frac{1}{p}}
$$

We begin to study the spaces $W_{N}^{2, p}\left(\Omega,|x|^{-s} d x\right)$ by pointing out some embedding results. Firstly, notice that the boundedness of the domain $\Omega$ implies

$$
\begin{equation*}
W_{N}^{2, p}\left(\Omega,|x|^{-s} d x\right) \hookrightarrow W_{N}^{2, p}\left(\Omega,|x|^{-s_{0}} d x\right) \quad \text { if } s_{0} \leq s \tag{2.1}
\end{equation*}
$$

Lemma 2.1. Assume $s>n-n p$. Then

$$
W_{N}^{2, p}\left(\Omega,|x|^{-s} d x\right) \hookrightarrow W^{2, \tau}(\Omega) \cap W_{0}^{1, \tau}(\Omega)
$$

where $\tau<\frac{n p}{n-s}<p$ if $s<0$, or $\tau=p$ otherwise.
Proof. If $s \geq 0$ the conclusion is immediate. Assume $s<0$. For any $u \in D_{0}$ and $\tau \in[1, n p /(n-s))$ we use elliptic regularity estimates, see for instance [14, Lemma 9.17], to get
$\|u\|_{W^{2, \tau}(\Omega)}^{\tau} \leq c \int_{\Omega}|\Delta u|^{\tau} d x \leq c\left(\int_{\Omega}|x|^{-s}|\Delta u|^{p} d x\right)^{\frac{\tau}{p}}\left(\int_{\Omega}|x|^{\frac{s \tau}{p-\tau}} d x\right)^{\frac{p-\tau}{p}}$.
The last integral is finite as $s>n-n p$, and the lemma is proved.

The next lemma will be used in Section 3 to rigorously prove the equivalence between the second order system $(\mathcal{P})$ and the fourth order equation (1.6).
Lemma 2.2. If $s>n-n p$, then $u \in W_{N}^{2, p}\left(\Omega,|x|^{-s} d x\right)$ if and only if

$$
\begin{equation*}
u \in W^{2,1} \cap W_{0}^{1,1}(\Omega) \quad \text { and } \quad-\Delta u \in L^{p}\left(\Omega,|x|^{-s} d x\right) \tag{2.2}
\end{equation*}
$$

Proof. Clearly, any $u \in W_{N}^{2, p}\left(\Omega,|x|^{-s} d x\right)$ satisfies (2.2) by Lemma 2.1.
Conversely, fix $u$ satisfying (2.2). Assume in addition that $-\Delta u=0$ almost everywhere on a ball $B_{r}$ about 0 , so that $-\Delta u \in L^{p}(\Omega)$. Hence, $u \in W^{2, p}(\Omega) \cap W_{0}^{1, p}(\Omega)$ by elliptic regularity theory. Extend $u$ to a function $u$ in $W^{2, p}\left(\mathbb{R}^{n}\right)$ with compact support and take a sequence of mollifiers $\left\{\rho_{k}\right\}_{k \in \mathbb{N}}$. Since, for $k$ large enough, $-\Delta\left(\rho_{k} * u\right) \equiv 0$ on $B_{r / 2}$ and $\rho_{k} * u \rightarrow u$ in $W^{2, p}(\Omega)$, then $-\Delta\left(\rho_{k} * u\right) \rightarrow-\Delta u$ in $L^{p}\left(\Omega,|x|^{-s} d x\right)$. Let $u_{k}$ be the solution to

$$
\begin{cases}-\Delta u_{k}=-\Delta\left(\rho_{k} * u\right) & \text { in } \Omega \\ u_{k}=0 & \text { on } \partial \Omega\end{cases}
$$

It turns out that $u_{k} \in D_{0} \cap W_{N}^{2, p}(\Omega)$, as $u_{k}$ is smooth up to the boundary of $\Omega$ by regularity theory, and $-\Delta u_{k} \equiv 0$ in $B_{r / 2}$. In addition, $u_{k} \rightarrow u$ in $W^{2, p}(\Omega)$ and $-\Delta u_{k} \rightarrow-\Delta u$ in $L^{p}\left(\Omega,|x|^{-s} d x\right)$, that is sufficient to conclude that $u \in W_{N}^{2, p}\left(\Omega,|x|^{-s} d x\right)$.

For a general $u$ satisfying (2.2) let $u_{k}$ be the unique solution to

$$
\begin{cases}-\Delta u_{k}=\chi_{\Omega_{k}}(-\Delta u) & \text { in } \Omega \\ u_{k}=0 & \text { on } \partial \Omega\end{cases}
$$

where $\Omega_{k}:=\Omega \backslash \bar{B}_{\varepsilon_{k}}$ and $\varepsilon_{k} \rightarrow 0$. Then $u_{k} \in W^{2, p} \cap W_{0}^{1, p}(\Omega)$ and $u_{k} \in$ $W_{N}^{2, p}\left(\Omega,|x|^{-s} d x\right)$ by the first part of the proof. Clearly, the sequence $\left\{u_{k}\right\}_{k \in \mathbb{N}}$ is bounded in $W_{N}^{2, p}\left(\Omega,|x|^{-s} d x\right)$, and we can assume that $u_{k} \rightarrow \bar{u}$ weakly in $W_{N}^{2, p}\left(\Omega,|x|^{-s} d x\right)$. On the other hand, the sequence $-\Delta u_{k}$ converges to $-\Delta u$ in $L^{p}\left(\Omega,|x|^{-s} d x\right)$ by Lebesgue's theorem. Thus $\bar{u}=u$, that is, $u \in$ $W_{N}^{2, p}\left(\Omega,|x|^{-s} d x\right)$.
The next corollary is an immediate consequence of Lemma 2.2.
Corollary 2.3. Assume $s>n-n p$. For any $f \in L^{p}\left(\Omega,|x|^{-s} d x\right)$, the unique solution u to

$$
\begin{cases}-\Delta u=f & \text { in } \Omega \\ u=0 & \text { on } \partial \Omega\end{cases}
$$

belongs to $W_{N}^{2, p}\left(\Omega,|x|^{-s} d x\right)$.
Next we deal with embeddings in weighted $L^{p}$ spaces.
Lemma 2.4. If $s+b+2 p \geq 0$, then

$$
\Lambda(s, b):=\inf _{\substack{u \in W_{N}^{2, p}\left(\Omega,|x|^{-s} d x\right) \\ u \neq 0}} \frac{\int_{\Omega}|x|^{-s}|\Delta u|^{p} d x}{\int_{\Omega}|x|^{b}|u|^{p} d x}>0
$$

Proof. First of all, notice that $L^{p}\left(\Omega,|x|^{b_{0}} d x\right) \hookrightarrow L^{p}\left(\Omega,|x|^{b} d x\right)$ if $b_{0} \leq b$, that together with (2.1) implies

$$
\begin{equation*}
\Lambda(s, b) \geq c \Lambda\left(s_{0}, b_{0}\right) \quad \text { if } s_{0} \leq s \text { and } b_{0} \leq b \tag{2.3}
\end{equation*}
$$

We start with the lowest dimensions $n=1,2$. Fix an exponent $s_{0} \leq s$, such that $n-n p<s_{0} \leq b(p-1)$. Then $\Lambda(s, b) \geq c \Lambda\left(s_{0}, \frac{s_{0}}{p-1}\right)>0$ by (2.3) and Lemma 4.2 in the Appendix.

Now assume $n \geq 3$. In addition, assume first that $s<n-2 p$. By the weighted Rellich inequality in [19] (see also [21, Lemma 2.14]) and using [21, Lemma 2.9], one has that there exists a positive and explicitly known constant $c=c(n, p, s)$, such that

$$
\begin{equation*}
c \int_{\Omega}|x|^{-s-2 p}|u|^{p} d x \leq \int_{\Omega}|x|^{-s}|\Delta u|^{p} d x \quad \text { for any } u \in C_{N}^{2}(\bar{\Omega}) \tag{2.4}
\end{equation*}
$$

that in particular gives $c=\Lambda(s,-s-2 p)>0$. Thus

$$
\Lambda(s, b) \geq c \Lambda(s,-s-2 p)>0
$$

by (2.3). Finally, if $s \geq n-2 p$, we fix a parameter $s_{0}$ such that

$$
\max \{n-n p,-2 p-b\}<s_{0}<n-2 p \leq s
$$

that is possible as $b>-n$ and $n \geq 3$. Then (2.3) and (2.4) with $s$ replaced by $s_{0}$ give $\Lambda(s, b) \geq c \Lambda\left(s_{0},-s_{0}-2 p\right)>0$, and the lemma is proved.

Remark 2.5. If $\Omega$ contains the origin and $s+b+2 p<0$, then $\Lambda(s, b)=0$. Indeed, fix a nontrivial $\psi \in C_{c}^{\infty}\left(B_{1} \backslash\{0\}\right)$. For $k$ large enough the function $\psi_{k}(x)=\psi(k x)$ belongs to $D_{0}$. Thus

$$
\Lambda(s, b) \leq \frac{\int_{\Omega}|x|^{-s}\left|\Delta \psi_{k}\right|^{p} d x}{\int_{\Omega}|x|^{b}\left|\psi_{k}\right|^{p} d x}=c k^{s+2 p+b}=o(1) \quad \text { as } k \rightarrow \infty
$$

Remark 2.6. If $n-n p<s<n-2 p$, then $C_{N}^{2}(\bar{\Omega}) \subset W_{N}^{2, p}\left(\Omega,|x|^{-s} d x\right)$ and the space

$$
C_{N}^{2}(\bar{\Omega} \backslash\{0\}):=\left\{u \in C_{N}^{2}(\bar{\Omega}) \mid u \equiv 0 \text { on a neighborhood of the origin }\right\}
$$

is dense in $W_{N}^{2, p}\left(\Omega,|x|^{-s} d x\right)$, see Lemma 2.14 in [21].
Remark 2.7. By Lemma 2.2, the set $D_{0}$ is dense in the standard Sobolev space $W_{N}^{2, p}(\Omega)=W^{2, p}(\Omega) \cap W_{0}^{1, p}(\Omega)$. The smaller set $C_{N}^{2}(\bar{\Omega} \backslash\{0\})$ is dense in $W_{N}^{2, p}(\Omega)$ if $n>2 p$, compare with Remark 2.6.

The next compactness result is a crucial point for studying the eigenvalue problem (1.6).

Lemma 2.8. If $s+b+2 p>0$ then $W_{N}^{2, p}\left(\Omega,|x|^{-s} d x\right)$ is compactly embedded into $L^{p}\left(\Omega,|x|^{b} d x\right)$.

Proof. It suffices to show that any sequence $\left\{u_{k}\right\}_{k \in \mathbb{N}}$ that converges weakly to the null function in $W_{N}^{2, p}\left(\Omega,|x|^{-s} d x\right)$ is compact in $L^{p}\left(\Omega,|x|^{b} d x\right)$. Fix such a sequence, and take $\varepsilon>0$ small. Since clearly $\left\{u_{k}\right\}_{k \in \mathbb{N}}$ is bounded in $W^{2, p}\left(\Omega \backslash \bar{B}_{\varepsilon}\right)$, then $|x|^{b}\left|u_{k}\right|^{p} \rightarrow 0$ in $L^{1}\left(\Omega \backslash \bar{B}_{\varepsilon}\right)$ by Rellich theorem. Therefore, for any $b_{0} \in(-n, b)$ we have that

$$
\int_{\Omega}|x|^{b}\left|u_{k}\right|^{p} d x=\int_{B_{\varepsilon}}|x|^{b}\left|u_{k}\right|^{p} d x+o(1) \leq \varepsilon^{b-b_{0}} \int_{\Omega}|x|^{b_{0}}\left|u_{k}\right|^{p} d x+o(1)
$$

Now, if $b_{0}$ is close enough to $b$, then $s+b_{0}+2 p>0$. Hence

$$
\int_{\Omega}|x|^{b}\left|u_{k}\right|^{p} d x \leq c \varepsilon^{b-b_{0}}+o(1)
$$

by Lemma 2.4. The conclusion follows, as $\varepsilon>0$ was arbitrarily chosen.

## 3. Two equivalent problems

In this section we provide the rigorous proof of the equivalence between the eigenvalue problems $(\mathcal{P})$ and (1.6). We start with a preliminary result.
Lemma 3.1. Assume that $a, b>-n$ and that (1.3) holds. For any $f \in$ $L^{p}\left(\Omega,|x|^{b} d x\right)$, the problem

$$
\left\{\begin{array}{l}
-\Delta u=|x|^{a}|v|^{p^{\prime}-2} v  \tag{3.1a}\\
-\Delta v=|x|^{b}|f|^{p-2} f \\
u \in W_{N}^{2, p}\left(\Omega,|x|^{-a(p-1)} d x\right), v \in W_{N}^{2, p^{\prime}}\left(\Omega,|x|^{-b\left(p^{\prime}-1\right)} d x\right)
\end{array}\right.
$$

admits a unique solution.
Proof. First of all, notice that $a(p-1)>n-n p, b\left(p^{\prime}-1\right)>n-n p^{\prime}$. Thus the results in Section 2 apply to the spaces $W_{N}^{2, p}\left(\Omega,|x|^{-a(p-1)} d x\right)$ and $W_{N}^{2, p^{\prime}}\left(\Omega,|x|^{-b\left(p^{\prime}-1\right)} d x\right)$.

Since $|x|^{b}|f|^{p-2} f \in L^{p^{\prime}}\left(\Omega,|x|^{-b\left(p^{\prime}-1\right)} d x\right)$, then Corollary 2.3 guarantees $(3.1 \mathrm{~b})$ has a unique solution $v \in W_{N}^{2, p^{\prime}}\left(\Omega,|x|^{-b\left(p^{\prime}-1\right)} d x\right)$. The embedding Lemma 2.4 gives that $|x|^{a}|v|^{p^{\prime}-2} v \in L^{p}\left(\Omega,|x|^{-a(p-1)} d x\right)$. Thus there exists a unique solution $u \in W_{N}^{2, p}\left(\Omega,|x|^{-a(p-1)} d x\right)$ to (3.1a), thanks again to Corollary 2.3 .
We are ready to prove the equivalence result we need.
Lemma 3.2. Assume that $a, b>-n$ and that (1.3) holds. Let $\mu, \lambda_{1}, \lambda_{2} \in \mathbb{R}$ satisfying (1.5). Then the following statements are equivalent.
(a) $u \in W_{N}^{2, p}\left(\Omega,|x|^{-a(p-1)} d x\right)$ is a weak solution to (1.6).
(b) The pair $u, v:=-|x|^{-a(p-1)}|\Delta u|^{p-2} \Delta u$ solves the system

$$
\left\{\begin{array}{l}
-\Delta u=|x|^{a}|v|^{p^{\prime}-2} v  \tag{3.2a}\\
-\Delta v=\mu|x|^{b}|u|^{p-2} u \\
u \in W_{N}^{2, p}\left(\Omega,|x|^{-a(p-1)} d x\right), v \in W_{N}^{2, p^{\prime}}\left(\Omega,|x|^{-b\left(p^{\prime}-1\right)} d x\right)
\end{array}\right.
$$

(c) The pair $u, v:=-|x|^{-a(p-1)}|\Delta u|^{p-2} \Delta u$ is a finite-energy solution to $(\mathcal{P})$, in the sense of Definition 1.1.

Proof. If $u \in W_{N}^{2, p}\left(\Omega,|x|^{-a(p-1)} d x\right)$, then $u \in L^{p}\left(\Omega,|x|^{b} d x\right)$ by Lemma 2.4. Thus we can apply Lemma 3.1 to find a unique pair $u_{0}, v_{0}$ such that

$$
\left\{\begin{array}{l}
-\Delta u_{0}=|x|^{a}\left|v_{0}\right|^{p^{\prime}-2} v_{0} \\
-\Delta v_{0}=\mu|x|^{b}|u|^{p-2} u \\
u_{0} \in W_{N}^{2, p}\left(\Omega,|x|^{-a(p-1)} d x\right), v_{0} \in W^{2, p^{\prime}}\left(\Omega,|x|^{-b\left(p^{\prime}-1\right)} d x\right)
\end{array}\right.
$$

Notice that $v_{0}=|x|^{-a(p-1)}\left|\Delta u_{0}\right|^{p-2}\left(-\Delta u_{0}\right)$ almost everywhere in $\Omega$. Therefore, if $u$ solves (1.6), then for any $\varphi \in D_{0}$ it holds that

$$
\begin{aligned}
\int_{\Omega}|x|^{-a(p-1)}|\Delta u|^{p-2} \Delta u \Delta \varphi d x & =\mu \int_{\Omega}|x|^{b}|u|^{p-2} u \varphi d x \\
& =\int_{\Omega}\left(-\Delta v_{0}\right) \varphi d x=\int_{\Omega} v_{0}(-\Delta \varphi) d x \\
& =\int_{\Omega}|x|^{-a(p-1)}\left|\Delta u_{0}\right|^{p-2} \Delta u_{0} \Delta \varphi d x
\end{aligned}
$$

that readily gives that $u=u_{0}$, since $u, u_{0} \in W_{N}^{2, p}\left(\Omega,|x|^{-a(p-1)} d x\right)$ and $D_{0}$ is dense in $W_{N}^{2, p}\left(\Omega,|x|^{-a(p-1)} d x\right)$. Hence also $v=v_{0}$, the pair $u, v$ solves (3.2), and the first implication is proved.

The equivalence between $(b)$ and $(c)$ is immediate, thanks to Lemma 2.1 and Corollary 2.3. It remains to be shown that $(b) \operatorname{implies}(a)$. If $(u, v)$ solves (3.2), then for every $\varphi \in D_{0}$ it holds that

$$
\begin{aligned}
\mu \int_{\Omega}|x|^{b}|u|^{p-2} u \varphi d x & =\int_{\Omega} v(-\Delta \varphi) d x \\
& =\int_{\Omega}|x|^{-a(p-1)}|\Delta u|^{p-2} \Delta u \Delta \varphi d x
\end{aligned}
$$

that is, $u$ solves (1.6).
Lemma 3.2 shows that finite energy solutions to $(\mathcal{P})$ are the stationary points of the functional

$$
u \mapsto \int_{\Omega}|x|^{-a(p-1)}|\Delta u|^{p} d x
$$

on the constraint

$$
M=\left\{\left.u \in W_{N}^{2, p}\left(\Omega,|x|^{-a(p-1)} d x\right)\left|\int_{\Omega}\right| x\right|^{b}|u|^{p} d x=1\right\}
$$

If $a, b>-n$ and (1.3) holds, then $M$ is weakly compact by Lemma 2.8. Thus the infimum

$$
\begin{equation*}
\mu:=\Lambda(a(p-1), b)=\inf _{\substack{u \in M \\ u \neq 0}} \int_{\Omega}|x|^{-a(p-1)}|\Delta u|^{p} d x \tag{3.3}
\end{equation*}
$$

is positive and attained.

Remark 3.3. Let $(u, v)$ be a finite energy solution to $(\mathcal{P})$. Clearly, if $n=1$ then $u, v$ are continuous on $\bar{\Omega}$. If $n \geq 2$ a standard bootstrap argument can be used to check that $u, v \in W_{\mathrm{loc}}^{2, q}(\Omega \backslash\{0\})$ for any $q \in[1, \infty]$. In particular, $u, v$ are continuous on $\Omega \backslash\{0\}$.

The next lemma deals with minimizers for $\mu$.
Lemma 3.4. Assume that $a, b>-n$ and that (1.3) holds.
If $u \in W_{N}^{2, p}\left(\Omega,|x|^{-a(p-1)} d x\right)$ achieves $\mu$, then, up to a change of sign, $u$ is superharmonic and positive on $\Omega$.
Proof. Let $v=|x|^{-a(p-1)}|\Delta u|^{p-2}(-\Delta u)$, so that the pair $(u, v)$ solves (3.2). Use Corollary 2.3 to introduce $u_{0}$ via

$$
\left\{\begin{array}{l}
-\Delta u_{0}=|x|^{a}|v|^{p^{\prime}-1} \\
u_{0} \in W_{N}^{2, p}\left(\Omega,|x|^{-a(p-1)} d x\right)
\end{array}\right.
$$

In particular, $u_{0}$ is superharmonic and positive on $\Omega$. Next, put

$$
g=|x|^{a}|v|^{p^{\prime}-2} v
$$

Thus $u$ and $u_{0}$ solve, for some $\tau \in[1, p)$,

$$
\left\{\begin{array} { l } 
{ - \Delta u = g } \\
{ u \in W ^ { 2 , \tau } \cap W _ { 0 } ^ { 1 , \tau } ( \Omega ) , }
\end{array} \quad \left\{\begin{array}{l}
-\Delta u_{0}=|g| \\
u_{0} \in W^{2, \tau} \cap W_{0}^{1, \tau}(\Omega)
\end{array}\right.\right.
$$

Since $-\Delta\left(u_{0} \pm u\right) \geq 0$ and $u_{0} \pm u=0$ on the boundary of $\Omega$, then $u_{0} \pm u \geq 0$, that is, $u_{0} \geq|u|$. On the other hand, $\left|\Delta u_{0}\right|=|g|=|\Delta u|$. Therefore

$$
\mu \leq \frac{\int_{\Omega}|x|^{-a(p-1)}\left|\Delta u_{0}\right|^{p} d x}{\int_{\Omega}|x|^{b}\left|u_{0}\right|^{p} d x} \leq \frac{\int_{\Omega}|x|^{-a(p-1)}|\Delta u|^{p} d x}{\int_{\Omega}|x|^{b}|u|^{p} d x}=\mu
$$

that is, $u_{0}$ attains $\mu$ and $u_{0}=|u|$. Since $u_{0}$ is positive in $\Omega$, then $u$ and $-\Delta u$ have constant sign (use Remark 3.3), as desired.

Proof of Theorem 1.3. Assume that (1.3) holds. The existence of a principal eigenvalue is immediate, thanks to Lemmata 2.8 and 3.4 and the equivalence given by Lemma 3.2.
Remark 3.5. Assume that $a, b>-n$ and

$$
\frac{a}{p^{\prime}}+\frac{b}{p}+2<0
$$

so that, in particular, $n \geq 3$. Then $\operatorname{Problem}(\mathcal{P})$ has no positive principal eigenvalue. This is a consequence of [1, Theorem 2.2].

We conclude the section pointing out a symmetry result about the infimum in (3.3). It is convenient to use the notation $\mu(a, b, p)$ to emphasize the dependence of $\mu$ on the exponents $a, b$ and $p$.
Proposition 3.6. If $a, b>-n$ and (1.3) holds, then $\mu\left(b, a, p^{\prime}\right)^{p}=\mu(a, b, p)^{p^{\prime}}$.

Proof. Let $u$ be an extremal for $\mu(a, b, p)$. By Lemma 3.2, the pair $u, v$, where

$$
v:=-|x|^{-a(p-1)}|\Delta u|^{p-2} \Delta u
$$

solves (3.2) with respect to the eigenvalue $\mu=\mu(a, b, p)$. Hence

$$
\begin{aligned}
& \int_{\Omega}|x|^{-b\left(p^{\prime}-1\right)}|\Delta v|^{p^{\prime}} d x=\mu(a, b, p)^{p^{\prime}} \int_{\Omega}|x|^{b}|u|^{p} d x \\
& \quad=\mu(a, b, p)^{p^{\prime}-1} \int_{\Omega}|x|^{-a(p-1)}|\Delta u|^{p} d x=\mu(a, b, p)^{p^{\prime}-1} \int_{\Omega}|x|^{a}|v|^{p^{\prime}} d x
\end{aligned}
$$

Thus $\mu\left(b, a, p^{\prime}\right) \leq \mu(a, b, p)^{p^{\prime}-1}$, or equivalently $\mu\left(b, a, p^{\prime}\right)^{p} \leq \mu(a, b, p)^{p^{\prime}}$. The opposite inequality follows by exchanging the roles of $u$ and $v$.

## 4. The linear problem

In this section we prove Theorem 1.4. First of all we notice that if $p=2$, then equation (1.6) reduces to

$$
\begin{equation*}
\Delta\left(|x|^{-a} \Delta u\right)=\mu|x|^{b} u \tag{4.1}
\end{equation*}
$$

We denote by $X_{a}$ the Hilbert space $W_{N}^{2,2}\left(\Omega,|x|^{-a} d x\right)$, endowed with norm $\|\cdot\|_{a}$ and scalar product $(\cdot \mid \cdot)_{a}$.
Proof of Theorem 1.4 (a). We formally introduce the "solution operator" to (4.1) under Navier boundary conditions. More precisely, we define the linear operator

$$
T: L^{2}\left(\Omega,|x|^{b} d x\right) \rightarrow X_{a}, \quad(T f \mid w)_{a}=\int_{\Omega}|x|^{b} f w d x \quad \text { for } w \in X_{a}
$$

Then $T$ is continuous, positive and self-adjoint. Let $j: X_{a} \rightarrow L^{2}\left(\Omega,|x|^{b} d x\right)$ be the embedding in Lemma 2.4. Then the operator

$$
T \equiv j \circ T, T: L^{2}\left(\Omega,|x|^{b} d x\right) \rightarrow L^{2}\left(\Omega,|x|^{b} d x\right)
$$

is compact. Thus the point spectrum $\sigma_{p}(T)$ of $T$ is a non-increasing sequence $\left\{\nu_{k}\right\}_{k \in \mathbb{N}}$ of positive numbers converging to 0 , and

$$
\frac{1}{\nu_{k}}=\min \left\{\left.\frac{\int_{\Omega}|x|^{-a}|\Delta u|^{2} d x}{\int_{\Omega}|x|^{b} u^{2} d x} \right\rvert\, u \in \Lambda_{i}^{\perp}, 1 \leq i \leq k-1\right\}
$$

where $\Lambda_{i}$ is the eigenspace relative to the eigenvalue $\nu_{i}$. Thus (a) readily follows by Lemma 3.2.

Proof of Theorem $1.4(b)$. We will use the theory of abstract positive operators on Banach lattices, for which we refer to the monograph [26]. Recall that $L^{2}\left(\Omega,|x|^{b} d x\right)$ has a natural Banach lattice structure induced by the cone $P_{+}$ of nonnegative functions. We will show that $T$ is positive and irreducible. Then, the conclusion will follow thanks to an adaptation of Theorem V.5.2 in [26], that guarantees that the following facts hold:

- the spectral radius $r(T) \in \mathbb{R}_{+}$is an eigenvalue;
- the eigenspace $\Lambda(r(T))$ has dimension one, and is spanned by a (unique, normalized) quasi-interior point of $P_{+}$;
- $r(T)$ is the unique eigenvalue of $T$ with a positive eigenvector.

To check that $T$ is irreducible we first recall that that the only closed ideals in $L^{2}\left(\Omega,|x|^{b} d x\right)$ are the ones of the form

$$
I_{A}=\left\{f \in L^{2}\left(\Omega,|x|^{b} d x\right) \mid f=0 \text { on } A\right\}
$$

where $A$ is a measurable set, see for instance [26, p. 157]. Therefore, we have to show that if $A$ satisfies

$$
0<\int_{\Omega}|x|^{b} \chi_{A} d x<\int_{\Omega}|x|^{b} d x
$$

then $I_{A}$ is not fixed by $T$.
Let $f \in I_{A}$ be a nonnegative fixed function. Then the problem

$$
\left\{\begin{array}{l}
-\Delta v=|x|^{b} f \\
v \in X_{b}
\end{array}\right.
$$

admits a solution by Corollary 2.3, and $v \in W^{2, \tau} \cap W_{0}^{1, \tau}(\Omega)$ for $\tau>1$ small enough. The minimum principle implies that $v$ is strictly positive in $\Omega$. For the same reason, the problem

$$
\left\{\begin{array}{l}
-\Delta u=|x|^{a} v \\
u \in X_{a},
\end{array}\right.
$$

defines a function $u$ that is strictly positive in $\Omega$. Hence $u \equiv T f \notin I_{A}$, and this proves the irreducibility property. The same argument proves also the positivity property.

Remark 4.1. There are of course a number of possible generalizations of the above results. For instance, the weights $|x|^{a},|x|^{b}$ may be replaced by more general measurable weights $\alpha(x), \beta(x)$, suitably pinched between power-type functions. More general differential operators might be considered as well. One could wonder what happens if $-\Delta$ is replaced by an $m$-order (possibly, fractional) differential operator.

Finally, one may wonder if a Faber-Krahn type inequality can hold, in particular in the non-linear case $p \neq 2$.

## Appendix. An inequality in lower dimensions

We sketch here the proof of a second order integral inequality in low dimensions by using, in essence, the Rellich-type identity in [19].

Lemma 4.2. Let $n \in\{1,2\}$ and let $\Omega$ be a a bounded domain in $\mathbb{R}^{n}$. If $n=2$, assume that $\Omega$ is of class $C^{2}$. For any $s>n-n p$, there exists a constant $c>0$ such that

$$
c \int_{\Omega}|x|^{\frac{s}{p-1}}|u|^{p} d x \leq \int_{\Omega}|x|^{-s}|\Delta u|^{p} d x
$$

for any $u \in C_{N}^{2}(\bar{\Omega})$ such that $\Delta u=0$ in a neighborhood of 0 .
Proof. We can assume that $\Omega$ is contained in the unit ball about the origin. Put

$$
a=\frac{s}{p-1}
$$

and notice that $a>-n$. We argue in a heuristic way. A more rigorous proof requires a suitable approximation of the weight $|x|^{a+2}$ by smooth functions. We omit details.

Fix $u \in C_{N}^{2}(\bar{\Omega})$ such that $\Delta u=0$ in a neighborhood of 0 . For $p \geq 2$ one clearly has

$$
(p-1) \int_{\Omega}|\nabla u|^{2}|u|^{p-2} d x=\int_{\Omega}(-\Delta u)|u|^{p-2} u d x .
$$

For general $p>1$, one can check that $|\nabla u|^{2}|u|^{p-2} \in L^{1}(\Omega)$ and

$$
\begin{equation*}
(p-1) \int_{\Omega}|\nabla u|^{2}|u|^{p-2} d x \leq \int_{\Omega}|\Delta u||u|^{p-1} d x \tag{4.2}
\end{equation*}
$$

Next, we are allowed to use integration by parts again and Hölder's inequality to estimate

$$
\begin{aligned}
& (a+2)(a+n) \int_{\Omega}|x|^{a}|u|^{p} d x=-\int_{\Omega}\left(\Delta|x|^{a+2}\right)|u|^{p} d x \\
& \quad=p \int_{\Omega}\left(\nabla|x|^{a+2} \cdot \nabla u\right)|u|^{p-2} u d x \leq p(a+2) \int_{\Omega}|x|^{a+1}|\nabla u||u|^{p-1} d x \\
& \quad \leq p(a+2)\left(\int_{\Omega}|\nabla u|^{2}|u|^{p-2} d x\right)^{\frac{1}{2}}\left(\int_{\Omega}|x|^{2 a+2}|u|^{p} d x\right)^{\frac{1}{2}} .
\end{aligned}
$$

Therefore, from $a+2 \geq a+n>0$ we infer

$$
\left(\frac{a+n}{p}\right)^{2} \int_{\Omega}|x|^{a}|u|^{p} d x \leq \int_{\Omega}|\nabla u|^{2}|u|^{p-2} d x \leq \frac{1}{p-1} \int_{\Omega}|u|^{p-1}|\Delta u| d x
$$

by (4.2). Then we use again Hölder's inequality to estimate

$$
c \int_{\Omega}|x|^{a}|u|^{p} d x \leq\left(\int_{\Omega}|x|^{a}|u|^{p} d x\right)^{\frac{1}{p^{\prime}}}\left(\int_{\Omega}|x|^{-a(p-1)}|\Delta u|^{p} d x\right)^{\frac{1}{p}}
$$

where $c=c(a, n, p)>0$. Thus

$$
c \int_{\Omega}|x|^{a}|u|^{p} d x \leq \int_{\Omega}|x|^{-a(p-1)}|\Delta u|^{p} d x
$$

as desired.

## Acknowledgements

The authors would like to thank the anonymous referee for comments, suggestions, and for having raised the open problems in Remark 4.1.

## References

[1] M.-F. Bidaut-Véron, Local behaviour of the solutions of a class of nonlinear elliptic systems. Adv. Differential Equations 5 (2000), no. 1-3, 147-192.
[2] M. F. Bidaut-Véron and H. Giacomini, A new dynamical approach of EmdenFowler equations and systems. Adv. Differential Equations 15 (2010), no. 11-12, 1033-1082.
[3] D. Bonheure, E. Moreira dos Santos and M. Ramos, Ground state and nonground state solutions of some strongly coupled elliptic systems. Trans. Amer. Math. Soc. 364 (2012), no. 1, 447-491.
[4] D. Bonheure, E. Moreira dos Santos and M. Ramos, Symmetry and symmetry breaking for ground state solutions of some strongly coupled elliptic systems. J. Funct. Anal. 264 (2013), no. 1, 62-96.
[5] J. Busca and R. Manásevich, A Liouville-type theorem for Lane-Emden systems. Indiana Univ. Math. J. 51 (2002), no. 1, 37-51.
[6] M. Calanchi and B. Ruf, Radial and non radial solutions for Hardy-Hénon type elliptic systems. Calc. Var. Partial Differential Equations 38 (2010), no. 1-2, 111-133.
[7] A. Carioli and R. Musina, The Hénon-Lane-Emden system: sharp nonexistence results. Preprint (2015).
[8] G. Caristi, L. D'Ambrosio and E. Mitidieri, Representation formulae for solutions to some classes of higher order systems and related Liouville theorems. Milan J. Math. 76 (2008), 27-67.
[9] Ph. Clément, D. G. de Figueiredo and E. Mitidieri, Positive solutions of semilinear elliptic systems. Comm. Partial Differential Equations 17 (1992), no. 5-6, 923-940.
[10] L. D'Ambrosio and E. Mitidieri, Hardy-Littlewood-Sobolev systems and related Liouville theorems. Discrete Contin. Dyn. Syst. Ser. S 7 (2014), no. 4, 653-671.
[11] D. G. de Figueiredo, I. Peral and J. D. Rossi, The critical hyperbola for a Hamiltonian elliptic system with weights. Ann. Mat. Pura Appl. (4) 187 (2008), no. 3, 531-545.
[12] M. Fazly, Liouville type theorems for stable solutions of certain elliptic systems. Adv. Nonlinear Stud. 12 (2012), no. 1, 1-17.
[13] M. Fazly and N. Ghoussoub, On the Hénon-Lane-Emden conjecture. Discrete Contin. Dyn. Syst. 34 (2014), no. 6, 2513-2533.
[14] D. Gilbarg and N. S. Trudinger, Elliptic partial differential equations of second order, second edition, Grundlehren der Mathematischen Wissenschaften, 224, Springer, Berlin, 1983.
[15] P.-L. Lions, The concentration-compactness principle in the calculus of variations. The limit case. I. Rev. Mat. Iberoamericana 1 (1985), no. 1, 145-201.
[16] E. Mitidieri, A Rellich identity and applications. Rapporti interni Università di Udine 25 (1990), 1-35.
[17] E. Mitidieri, A Rellich type identity and applications. Comm. Partial Differential Equations 18 (1993), no. 1-2, 125-151.
[18] E. Mitidieri, Nonexistence of positive solutions of semilinear elliptic systems in $\mathbf{R}^{N}$. Differential Integral Equations 9 (1996), no. 3, 465-479.
[19] E. Mitidieri, A simple approach to Hardy inequalities. Mat. Zametki 67 (2000), no. 4, 563-572; translation in Math. Notes 67 (2000), no. 3-4, 479-486.
[20] M. Montenegro, The construction of principal spectral curves for Lane-Emden systems and applications. Ann. Scuola Norm. Sup. Pisa Cl. Sci. (4) 29 (2000), no. 1, 193-229.
[21] R. Musina, Optimal Rellich-Sobolev constants and their extremals. Differential Integral Equations 27 (2014), no. 5-6, 579-600.
[22] R. Musina and K. Sreenadh, Radially symmetric solutions to the Hénon-LaneEmden system on the critical hyperbola. Commun. Contemp. Math. 16 (2014), no. 3, 1350030 (16 pages).
[23] L. A. Peletier and R. C. A. M. Van der Vorst, Existence and nonexistence of positive solutions of nonlinear elliptic systems and the biharmonic equation. Differential Integral Equations 5 (1992), no. 4, 747-767.
[24] Q. H. Phan, Liouville-type theorems and bounds of solutions for Hardy-Hénon elliptic systems. Adv. Differential Equations 17 (2012), no. 7-8, 605-634.
[25] P. Poláĉik, P. Quittner and P. Souplet, Singularity and decay estimates in superlinear problems via Liouville-type theorems, Part I: Elliptic systems. Duke Math. J. 139 (2007), 555-579.
[26] H. H. Schaefer, Banach lattices and positive operators. Springer, New York, 1974.
[27] J. Serrin and H. Zou, Non-existence of positive solutions of semilinear elliptic systems. in A tribute to Ilya Bakelman (College Station, TX, 1993), 55-68, Discourses Math. Appl., 3 Texas A \& M Univ., College Station, TX.
[28] J. Serrin and H. Zou, Non-existence of positive solutions of Hénon-Lane-Emden systems. Differential Integral Equations 9 (1996), 635-653.
[29] J. Serrin and H. Zou, Existence of positive solutions of the Hénon-Lane-Emden system. Atti Semin. Mat. Fis. Univ. Modena 46 (1998), 369-380.
[30] P. Souplet, The proof of the Lane-Emden conjecture in four space dimensions. Adv. Math. 221 (2009), no. 5, 1409-1427.
[31] X. J. Wang, Sharp constant in a Sobolev inequality. Nonlinear Anal. 20 (1993), no. 3, 261-268.

Andrea Carioli<br>Mathematics Area, S.I.S.S.A.<br>via Bonomea, 265<br>34136 Trieste, Italy<br>e-mail: acarioli@sissa.it<br>Roberta Musina<br>Dipartimento di Matematica e Informatica, Università di Udine<br>via delle Scienze, 206<br>33100 Udine, Italy<br>e-mail: roberta.musina@uniud.it


[^0]:    A. Carioli is partially supported by INdAM-GNAMPA.
    R. Musina is partially supported by MIUR-PRIN 201274FYK7_004.

