



On semilinear Δ_λ -Laplace equation

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ABSTRACT

We prove some existence, nonexistence and regularity results for the boundary value problem

$$\Delta_\lambda u + f(u) = 0 \quad \text{in } \Omega, \quad u|_{\partial\Omega} = 0,$$

where Ω is a bounded subset of \mathbb{R}^N , $N \geq 2$, and Δ_λ is a Δ_λ -Laplacian, i.e. a “degenerate” elliptic operator of the kind

$$\Delta_\lambda := \sum_{i=1}^N \partial_{x_i} (\lambda_i^2(x) \partial_{x_i}), \quad \lambda = (\lambda_1, \dots, \lambda_N).$$

Together with some assumptions made in Franchi and Lanconelli (1984) [1], the family λ is supposed to verify a condition making Δ_λ homogeneous of degree two with respect to a group of dilations in \mathbb{R}^N .

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1. Introduction

1.1. Homogeneous Δ_λ -Laplacians. The semilinear Dirichlet problem

In recent years a certain number of papers have been devoted to “degenerate” elliptic operators whose simplest prototype in \mathbb{R}^N is the following one.

$$\Delta_\alpha := \Delta_{(1)} + |x^{(1)}|^{2\alpha} \Delta_{(2)}, \quad \alpha > 0. \tag{1.1}$$

Here $x = (x^{(1)}, x^{(2)})$, $x^{(i)} \in \mathbb{R}^{N_i}$, $i = 1, 2$ denotes the point of \mathbb{R}^N , $N = N_1 + N_2$, and $\Delta_{(i)}$ stands for the classical Laplacian in \mathbb{R}^{N_i} .

Nowadays, Δ_α in (1.1) is usually quoted in the literature as Grushin’s operator. However, if α is a nonnegative integer, Δ_α falls into the general class of Hörmander’s operators, sum of squares of vector fields generating a Lie algebra of maximum rank at any point. If α is not an integer, then Δ_α is contained in a family of operators of the kind

$$\Delta_\lambda := \sum_{i=1}^N \partial_{x_i} (\lambda_i^2 \partial_{x_i}), \quad \partial_{x_i} = \frac{\partial}{\partial x_i} \tag{1.2}$$

first studied in [1–3] with a geometrical technique taking into account the property of the control (or Carnot–Carathéodory) metric $d := d_\lambda$ generated by the vector fields

$$X_i := \lambda_i \partial_{x_i}, \quad i = 1, \dots, N. \tag{1.3}$$

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The aim of this paper is to establish existence, nonexistence and regularity results for the problem

$$\begin{cases} \Delta_\lambda u + f(u) = 0 & \text{in } \Omega, \\ u|_{\partial\Omega} = 0 \end{cases} \quad (1.4)$$

where Ω is a bounded open subset of \mathbb{R}^N and Δ_λ is the operator in (1.2) related to a N -tuple $\lambda = (\lambda_1, \dots, \lambda_N)$ of continuous functions in \mathbb{R}^N verifying, together with the assumption in [1] (which will be recalled, for reader convenience, in Section 1.2), the following one.

(H1) There exists a group of dilations $(\delta_t)_{t>0}$,

$$\delta_t : \mathbb{R}^N \longrightarrow \mathbb{R}^N, \quad \delta_t(x) = \delta_t(x_1, \dots, x_N) = (t^{\varepsilon_1}x_1, \dots, t^{\varepsilon_N}x_N)$$

with $1 = \varepsilon_1 \leq \varepsilon_2 \leq \dots \leq \varepsilon_N$, such that λ_i is δ_t homogeneous of degree $\varepsilon_i - 1$, i.e.,

$$\lambda_i(\delta_t(x)) = t^{\varepsilon_i - 1} \lambda_i(x), \quad \forall x \in \mathbb{R}^N, \forall t > 0, i = 1, \dots, N. \quad (1.5)$$

The number

$$Q := \varepsilon_1 + \varepsilon_2 + \dots + \varepsilon_N \quad (1.6)$$

is the *homogeneous dimension* of \mathbb{R}^N with respect to $(\delta_t)_{t>0}$. It will play a crucial rôle both in the geometry and the functional setting naturally associated to Δ_λ . We explicitly remark that (1.5) is equivalent to the δ_t -homogeneity of degree one of the vector field X_i , that is, to the property

$$X_i(u(\delta_t(x))) = t(X_i u)(\delta_t(x)), \quad \forall x \in \mathbb{R}^N, \forall t > 0 \quad (1.7)$$

and for every $u \in C^\infty(\mathbb{R}^N)$.

Thus, (H1) implies that Δ_λ is δ_t -homogeneous of degree two. When the family λ satisfies all the hypotheses mentioned above, we call Δ_λ -Laplacian the operator in (1.2).

1.2. Hypotheses on $\lambda = (\lambda_1, \dots, \lambda_N)$

The function λ_i 's are continuous in \mathbb{R}^N , different from zero and of class C^1 in $\mathbb{R}^N \setminus \Pi$ where

$$\Pi = \left\{ (x_1, \dots, x_N) \in \mathbb{R}^N : \prod_{i=1}^N x_i = 0 \right\}.$$

Moreover, together with (H1), we assume the following properties.

(H2) $\lambda_1 = 1$, $\lambda_i(x) = \lambda_i(x_1, \dots, x_{i-1})$, $i = 2, \dots, N$.

(H3) There exists a constant $\rho \geq 0$ such that

$$0 \leq x_k \partial_{x_k} \lambda_i(x) \leq \rho \lambda_i(x) \quad \forall k \in \{1, \dots, i-1\}, \forall i = 2, \dots, N$$

and for every $x \in \mathbb{R}_+^N := \{(x_1, \dots, x_N) \in \mathbb{R}^N : x_i \geq 0 \forall i = 1, \dots, N\}$.

(H4) For every $x \in \mathbb{R}^N$, $\lambda_i(x) = \lambda_i(x^*)$ where

$$x^* = (|x_1|, \dots, |x_N|) \quad \text{if } x = (x_1, \dots, x_N).$$

Some remarks are in order.

Remark 1.1. If $\varepsilon_i = 1$ then $\lambda_i(x) = \lambda_i(0) > 0$ for every $x \in \mathbb{R}^N$. Indeed, if $\varepsilon_i = 1$ then λ_i is δ_t -homogeneous of degree zero and its continuity implies

$$\lambda_i(x) = \lambda_i(\delta_t(x)) = \lim_{t \searrow 0} \lambda_i(\delta_t(x)) = \lambda_i(0).$$

Moreover, $\lambda_i(0)$ has to be strictly positive since $\lambda_i > 0$ in $\mathbb{R}_+^N \setminus \Pi$. We want to stress that, vice versa, if $\lambda_i(0) > 0$ from (1.5) we get $\varepsilon_i = 1$ and $\lambda_i(x) = \lambda_i(0)$ for every $x \in \mathbb{R}^N$.

Note 1. Throughout the paper, without loss of generality, we assume $\lambda_i \equiv 1$ if $\varepsilon_i = 1$.

Remark 1.2. By condition (H2) the operator Δ_λ can be written as follows

$$\Delta_\lambda = \sum_{i=1}^N \lambda_i^2 \partial_{x_i}^2 = \sum_{i=1}^N X_i^2.$$

Remark 1.3. If the λ_i 's are smooth, then (H1) and (H2) imply the *hypoellipticity* of Δ_λ , i.e. the smoothness of the distributional solutions to $\Delta_\lambda u = f$ when f is smooth.

Indeed, since λ_i is smooth, δ_t -homogeneous and everywhere different from zero in $\mathbb{R}^N \setminus \Pi$, λ_i is a nonvanishing polynomial function (see, e.g., [4, Proposition 1.3.4]). Then, for every fixed $x \in \mathbb{R}^N$ there exists a multi-index $\alpha^{(i)} = (\alpha_1^{(i)}, \dots, \alpha_{i-1}^{(i)})$ such that

$$D^{\alpha^{(i)}} \lambda_i(x) \neq 0.$$

Using this property it is easy to recognize that

$$\text{Lie} \{X_1, \dots, X_N\}(x) \supseteq \{\partial_{x_1}, \dots, \partial_{x_N}\}.$$

Therefore

$$\text{rank Lie} \{X_1, \dots, X_N\} = N, \quad \forall x \in \mathbb{R}^N,$$

and, by the Hörmander Theorem in [5] and Remark 1.2, Δ_λ is hypoelliptic.

On the other hand, apart from the obvious case in which every λ_i is constant, Δ_λ is not elliptic at every point of Π (see Remark 1.1). Thus, if the λ_i 's are not constant

$$\dim(\text{span} \{\lambda_i(0)\partial_{x_i} : i = 1, \dots, N\}) < N = \dim(\text{span} \{\lambda_i(x)\partial_{x_i} : i = 1, \dots, N\}) \quad \forall x \notin \Pi.$$

As a consequence: there is no composition law \circ in \mathbb{R}^N making Δ_λ left translation invariant (see [4, Proposition 1.2.13]). So that, if the λ_i 's are not constant, there is no Lie group \mathbb{G} in \mathbb{R}^N making Δ_λ a sub-Laplacian on \mathbb{G} .

1.3. The Δ_λ -metric space

An absolutely continuous path $\gamma : [0, T] \rightarrow \mathbb{R}^N$ is called λ -subunit if there exist measurable functions $c_1, \dots, c_N : [0, T] \rightarrow \mathbb{R}$ such that, almost everywhere in $[0, T]$ one has

$$\dot{\gamma}(t) = \sum_{i=1}^N c_i(t)X_i(\gamma(t)), \quad \sum_{i=1}^N c_i^2(t) \leq 1.$$

Hereafter we agree to identify the vector fields

$$X_i(x) = \lambda_i(x)\partial_{x_i} \quad \text{with the function } \lambda_i(x)e_i, \quad e_i = \left(0, \dots, \underset{i}{1}, \dots, 0\right).$$

Assumption (H2) and the positivity of the λ_i 's in $\mathbb{R}^N \setminus \Pi$ imply the λ -connectivity of \mathbb{R}^N . Precisely, for every $x, y \in \mathbb{R}^N$ there exists a λ -subunit path $\gamma : [0, T] \rightarrow \mathbb{R}^N$ such that $\gamma(0) = x$ and $\gamma(T) = y$. We denote by $\Lambda(x, y)$ the family of the λ -subunit paths connecting x and y . Finally, if the subunit path γ is defined in the interval $[0, T]$, we let $l(\gamma) := T$. Then d_λ , what we call the λ -distance, is defined as follows: if $x, y \in \mathbb{R}^N$,

$$d_\lambda(x, y) := \inf\{l(\gamma) : \gamma \in \Lambda(x, y)\}.$$

It is quite trivial to verify that $(\mathbb{R}^N, d_\lambda)$ is a metric space. Sometime, in what follows, we will write d instead of d_λ . The d -ball of center $x \in \mathbb{R}^N$ and radius $r > 0$ will be denoted by $B_d(x, r)$. Hence

$$B_d(x, r) := \{y \in \mathbb{R}^N \mid d(x, y) < r\}.$$

A precise estimate of the d_λ -distance, and of the Lebesgue measure of the d_λ -balls, come from (H1)–(H4). Define a N -tuple of functions F_1, \dots, F_N on $\mathbb{R}_+^N \times [0, \infty[$ with the following recurrence law

$$\begin{cases} F_1(x, \tau) = \tau \\ F_i(x, \tau) = \tau \lambda_i(x_1 + F_1(x, \tau), \dots, x_{i-1} + F_{i-1}(x, \tau)), \quad i = 2, \dots, N. \end{cases} \tag{1.8}$$

Since λ_i is monotone increasing with respect to $x_j \in [0, \infty[$, $j = 1, \dots, i - 1$, and strictly positive in $\mathbb{R}_+^N \setminus \Pi$ the function $\tau \mapsto F_i(x, \tau)$ is strictly increasing in $[0, \infty[$, for every fixed $x \in \mathbb{R}_+^N$. We let

$$\phi_i(x, \cdot) = (F_i(x, \cdot))^{-1}, \quad i = 1, 2, \dots, N. \tag{1.9}$$

Then, by Theorems 2.6 and 2.7 in [1] we have the following. There exist two strictly positive constants c_1 and c_2 such that

$$c_1 \leq \frac{d(x, y)}{\sum_{i=1}^N \phi_i(x^*, |x_i - y_i|)} \leq c_2 \tag{1.10}$$

and

$$c_1 \leq \frac{|B_d(x, r)|}{\prod_{i=1}^N F_i(x^*, r)} \leq c_2 \tag{1.11}$$

for every $x = (x_1, \dots, x_N), y = (y_1, \dots, y_N) \in \mathbb{R}^N$ and for every $r > 0$. $|B_d(x, r)|$ denotes the Lebesgue measure of $B_d(x, r)$.

From (1.10) and (1.11), and from the properties of the functions F_i 's showed in the Appendix, we obtain some crucial result on the metric d_λ and on the measure of the d_λ -balls. Precisely:

(d1) there exist positive constants c_1^*, \dots, c_N^* such that

$$d_\lambda(x, y) \leq \sum_{i=1}^N c_i^* |x_i - y_i|^{\frac{1}{c_i}} \tag{1.12}$$

for every $x = (x_1, \dots, x_N)$ and $y = (y_1, \dots, y_N)$ in \mathbb{R}^N .
(See Corollary A.3.)

(d2) For every fixed compact subset $K \subseteq \mathbb{R}^N$ there exists $d^* = d_K^* > 0$ such that

$$d_\lambda(x, y) \geq d^* |x - y|, \quad \forall x, y \in K \tag{1.13}$$

(See (A.4).)

(d3) There exists a positive constant c_d such that

$$|B(x, 2r)| \leq c_d |B(x, r)| \quad \forall x \in \mathbb{R}^N, \forall r > 0, \tag{1.14}$$

where $c_d = \frac{c_2}{c_1}$. (See (A.1) and (1.11).)

Moreover, from (A.2),

$$|B(x, R)| \leq c_d 2^Q \left(\frac{R}{r}\right)^Q |B(x, r)| \tag{1.15}$$

for every $x \in \mathbb{R}^N$ and $0 < r < R$. Q is the number defined in (1.6), the homogeneous dimension of \mathbb{R}^N with respect to $(\delta_t)_{t>0}$.

1.4. The Δ_λ -functional setting

For a function u of class C^1 we let

$$|\nabla_\lambda u|^2 := \sum_{i=1}^N |\lambda_i \partial_{x_i} u|^2. \tag{1.16}$$

Given a bounded open set $\Omega \subseteq \mathbb{R}^N$, for every $p \in]1, \infty[$ we denote by

$$\dot{W}_\lambda^{1,p}(\Omega)$$

the closure of $C_0^1(\Omega)$ with respect to the norm

$$\|u\|_{1,p} := \left(\int_\Omega |\nabla_\lambda u|^p dx \right)^{\frac{1}{p}}. \tag{1.17}$$

In Section 3 we will directly recognize the continuous embedding

$$\dot{W}_\lambda^{1,p}(\Omega) \hookrightarrow L^{p_\lambda^*}(\Omega) \tag{1.18}$$

for every $p \in]1, Q[$ and

$$p_\lambda^* := \frac{pQ}{Q-p}. \tag{1.19}$$

As before, Q is the homogeneous dimension of \mathbb{R}^N with respect to $(\delta_t)_{t>0}$.

We would like to stress that the embedding (1.18) also follows from generalized doubling property (1.15) and the following Poincaré inequality

$$\int_{B_d(x,r)} |u - u_r| dy \leq Cr \int_{B_d(x,\theta r)} |\nabla_\lambda u| dy \quad \forall u \in C^1(\overline{B_d(x,\theta r)}) \tag{1.20}$$

where $C > 0$ and $\theta > 1$ are suitable constants independent of u, x and r (see [3] and [6, Section 3]) and u_r stands for the average of u on $B_d(x, r)$, i.e.,

$$u_r = \frac{1}{|B_d(x, r)|} \int_{B_d(x,r)} u(y) dy.$$

Indeed, nowadays it is well known that (1.15) and (1.20) imply the embedding (1.18); some important references are [7–10].

1.5. Examples of Δ_λ operators

Example 1.4. Let us split \mathbb{R}^N as follows

$$\mathbb{R}^N = \mathbb{R}^{N_1} \times \dots \times \mathbb{R}^{N_r}$$

and denote by $x = (x^{(1)}, \dots, x^{(r)})$, $x^{(i)} \in \mathbb{R}^{N_i}$, $i = 1, \dots, r$, the point of \mathbb{R}^N . Let $\Delta_{(i)}$ be the classical Laplace operator in \mathbb{R}^{N_i} . Given a multi-index $\alpha = (\alpha_1, \dots, \alpha_{r-1})$ $\alpha_j \geq 1, j = 1, \dots, r - 1$, define

$$\Delta_\alpha := \Delta_{(1)} + |x^{(1)}|^{2\alpha_1} \Delta_{(2)} + \dots + |x^{(r-1)}|^{2\alpha_{r-1}} \Delta_{(r)}. \tag{1.21}$$

Then $\Delta_\alpha = \Delta_\lambda$ with $\lambda = (\lambda^{(1)}, \dots, \lambda^{(r)})$ and $\lambda^{(i)} = |x^{(i-1)}|^{\alpha_{i-1}} 1^{(i)}$, $i = 1, \dots, r$. Here we agree to let $|x^{(0)}|^{\alpha_0} = 1$ and $1^{(i)} = (1, \dots, 1)$. A group of dilations for which λ satisfies (H1) is given by

$$\delta_t : \mathbb{R}^N \longrightarrow \mathbb{R}^N, \quad \delta_t(x^{(1)}, \dots, x^{(r)}) = (t^{\varepsilon_1} x^{(1)}, \dots, t^{\varepsilon_r} x^{(r)}) \tag{1.22}$$

with $\varepsilon_1 = 1$ and $\varepsilon_i = \alpha_{i-1} \varepsilon_{i-1} + 1, i = 2, \dots, r$. In particular, if $\alpha_1 = \dots = \alpha_{r-1} = 1$, the operator (1.21) and the dilation (1.22) becomes, respectively

$$\Delta_{(1)} + |x^{(1)}|^2 \Delta_{(2)} + \dots + |x^{(r-1)}|^2 \Delta_{(r)}$$

and

$$\delta_t(x^{(1)}, \dots, x^{(r)}) = (tx^{(1)}, t^2x^{(2)}, \dots, t^r x^{(r)}).$$

Example 1.5. Let $\mathbb{R}^N = \mathbb{R}^{N_1} \times \mathbb{R}^{N_2}$ and let $\mu : \mathbb{R}^{N_1} \rightarrow \mathbb{R}$ be continuous in its domain and of class C^1 and strictly positive outside the coordinate axis. Moreover, assume that $\mu(tx^{(1)}) = t^\alpha \mu(x^{(1)})$, for a suitable $\alpha > 0$, and for every $x^{(1)} \in \mathbb{R}^{N_1}$ and $t > 0$. Then, if we let $\lambda = (1^{(1)}, \mu 1^{(2)})$, we have

$$\Delta_\lambda = \Delta_{(1)} + (\mu(x^{(1)}))^2 \Delta_{(2)}. \tag{1.23}$$

This operator satisfies (H1) with respect to the dilations

$$\delta_t : \mathbb{R}^N \longrightarrow \mathbb{R}^N, \quad \delta_t(x^{(1)}, x^{(2)}) = (tx^{(1)}, t^{\alpha+1}x^{(2)}).$$

The class of the operators in (1.23) obviously contains Δ_α in (1.1), as well as

$$\Delta_z + |x|^{2\alpha_1} |y|^{2\alpha_2} \Delta_t, \quad z = (x, y) \in \mathbb{R}^m \times \mathbb{R}^n, \quad t \in \mathbb{R}^{N_2}.$$

Note 2. When $\mu(x^{(1)}) = \frac{1}{4}|x^{(1)}|$ the operator Δ_λ in (1.23) takes the form

$$\Delta_{x^{(1)}} + \frac{1}{4}|x^{(1)}|^2 \Delta_{x^{(2)}}.$$

On the other hand, if the dimension N_1 and N_2 verify the inequality $N_2 < \rho(N_1)$, where ρ is the so called Hurwitz–Radon function, then there exists a composition law \circ in \mathbb{R}^N making $\mathbb{H}_N := (\mathbb{R}^N, \circ, \delta_\lambda)$ a group of Heisenberg type (see [4, Remark 3.6.7]). Then if $\Delta_{\mathbb{H}_N}$ denotes the canonical sub-Laplacian on \mathbb{H}_N , for every smooth function $u : \mathbb{R}^N \rightarrow \mathbb{R}$ which is radially symmetric in the variable $x^{(1)}$, one has

$$\left(\Delta_{x^{(1)}} + \frac{1}{4}|x^{(1)}|^2 \Delta_{x^{(2)}} \right) u = \Delta_{\mathbb{H}_N} u,$$

see [4, page 251].

2. Some integral identities for Δ_λ : nonexistence results

In this section we prove some integral identities extending to the Δ_λ setting the classical Pohozaev identity for semilinear Poisson equation [11]. Pohozaev identity has been extended by several authors to general elliptic equations and systems, both in Riemannian and sub-Riemannian context, see, e.g., [12–14] and the references therein. To prove our identities we closely follow the original procedure of Pohozaev, just replacing the vector field $P = \sum_{i=1}^N x_i \partial_{x_i}$ in [11, page 1410], by

$$T := \sum_{i=1}^N \varepsilon_i x_i \partial_{x_i}, \tag{2.1}$$

the generator of the group of dilation $(\delta_t)_{t \geq 0}$ in (H1) (we say that T generates $(\delta_t)_{t \geq 0}$ since a function u is δ_t -homogeneous of degree m if and only if $Tu = mu$).

Throughout this section the λ_j 's are not supposed to verify assumptions (H3) and (H4) and $\Omega \subseteq \mathbb{R}^N$ will be assumed to be open, bounded, with C^1 -boundary and such that $\Omega = \text{int}(\overline{\Omega})$. We will denote by $\Lambda^2(\overline{\Omega})$ the linear space of the function $u \in C(\overline{\Omega})$ such that

$$X_j u, \quad X_j^2 u, \quad j = 1, \dots, N$$

exist in the weak sense of distributions in Ω and can be continuously extended to $\overline{\Omega}$ (as above, $X_j := \lambda_j \partial_{x_j}$).

Our first integral identity is the following one.

Theorem 2.1. For every $u \in \Lambda^2(\overline{\Omega})$ we have

$$\int_{\Omega} T(u) \Delta_{\lambda} u \, dx = \int_{\partial \Omega} \langle \nabla_{\lambda} u, \nu_{\lambda} \rangle T(u) \, ds - \frac{1}{2} \int_{\partial \Omega} |\nabla_{\lambda} u|^2 \langle T, \nu \rangle \, ds + \left(\frac{Q}{2} - 1 \right) \int_{\Omega} |\nabla_{\lambda} u|^2 \, dx, \quad (2.2)$$

where T is the vector field (2.1), $\langle \cdot, \cdot \rangle$ stands for the Euclidean inner product, ν is the outward normal to Ω , $\nu_{\lambda} = (\lambda_1 \nu_1, \dots, \lambda_N \nu_N)$ and

$$\nabla_{\lambda} u = (\lambda_1 \partial_{x_1} u, \dots, \lambda_N \partial_{x_N} u).$$

Proof. We first prove (2.2) assuming $u \in C^2(\overline{\Omega})$. An integration by parts gives¹

$$\int_{\Omega} T(u) \Delta_{\lambda} u \, dx = \int_{\partial \Omega} \varepsilon_i x_i \partial_{x_i} u \partial_{x_j} u \lambda_j^2 \nu_j \, ds - \int_{\Omega} \lambda_j^2 \partial_{x_j} u \partial_{x_j} (\varepsilon_i x_i \partial_{x_i} u) \, dx =: I_1 + I_2. \quad (2.3)$$

It is easily seen that

$$I_1 = \int_{\partial \Omega} T(u) \langle \nabla_{\lambda} u, \nu_{\lambda} \rangle \, ds, \quad (2.4)$$

while I_2 can be handled as follows

$$I_2 = - \int_{\Omega} \lambda_j^2 \partial_{x_j} u \delta_{ij} \varepsilon_i \partial_{x_i} u \, dx - \int_{\Omega} \lambda_j^2 \partial_{x_j} u \varepsilon_i x_i \partial_{x_j x_i} u \, dx =: I_{2,1} + I_{2,2}. \quad (2.5)$$

Obviously

$$I_{2,1} = - \int_{\Omega} \lambda_j^2 (\partial_{x_j} u)^2 \varepsilon_j \, dx.$$

Moreover, an integration by parts in $I_{2,2}$ gives

$$\begin{aligned} I_{2,2} &= - \int_{\partial \Omega} \lambda_j^2 \partial_{x_j} u \varepsilon_j x_i \partial_{x_j} u \nu_i \, ds + \int_{\Omega} \partial_{x_j} u \partial_{x_i} (\lambda_j^2 \partial_{x_j} u \varepsilon_i x_i) \, dx \\ &= - \int_{\partial \Omega} |\nabla_{\lambda} u|^2 \langle T, \nu \rangle \, ds + \int_{\Omega} |\nabla_{\lambda} u|^2 (\text{div } T) \, dx + \int_{\Omega} \lambda_j^2 \partial_{x_j} u \varepsilon_i x_i \partial_{x_j x_i} u \, dx + \int_{\Omega} (\partial_{x_j} u)^2 T \lambda_j^2 \, dx \\ &= (\text{keeping in mind that } \text{div } T = Q \text{ and that } \lambda_j \text{ is } \delta_t\text{-homogeneous} \\ &\quad \text{of degree } \varepsilon_j - 1, \text{ hence } T \lambda_j^2 = 2 \lambda_j T \lambda_j = 2(\varepsilon_j - 1) \lambda_j^2) \\ &\quad - \int_{\partial \Omega} |\nabla_{\lambda} u|^2 \langle T, \nu \rangle \, ds + Q \int_{\Omega} |\nabla_{\lambda} u|^2 \, dx - I_{2,2} + 2 \int_{\Omega} (\varepsilon_j - 1) (\lambda_j \partial_{x_j} u)^2 \, dx. \end{aligned}$$

Therefore,

$$I_{2,2} = - \frac{1}{2} \int_{\partial \Omega} |\nabla_{\lambda} u|^2 \langle T, \nu \rangle \, ds + \left(\frac{Q}{2} - 1 \right) \int_{\Omega} |\nabla_{\lambda} u|^2 \, dx - I_{2,1}$$

hence, by (2.5),

$$I_2 = - \frac{1}{2} \int_{\partial \Omega} |\nabla_{\lambda} u|^2 \langle T, \nu \rangle \, ds + \left(\frac{Q}{2} - 1 \right) \int_{\Omega} |\nabla_{\lambda} u|^2 \, dx.$$

¹ Repeated indices i and j are understood to be summed from 1 to N .

This identity, together with (2.3) and (2.4), implies (2.2) under the assumption $u \in C^2(\overline{\Omega})$. To complete the proof when $u \in \Lambda^2(\overline{\Omega})$ we argue as in [13, page 77]. Let $(\Omega_k)_{k \in \mathbb{N}}$ be a sequence of open set with smooth boundary such that

$$\overline{\Omega}_k \subseteq \Omega_{k+1} \quad \text{for every } k \in \mathbb{N} \text{ and } \cup_{k \in \mathbb{N}} \Omega_k = \Omega.$$

On every Ω_k we can approximate the function u with a sequence of $C^2(\overline{\Omega}_k)$ functions u_j for each of which identity (2.3) holds on Ω_k . Letting $j \rightarrow \infty$ we obtain (2.3) for u in Ω_k . Finally, as $k \rightarrow \infty$ we get our identity for $u \in \Lambda^2(\overline{\Omega})$. \square

From the previous theorem we easily obtain an integral identity for $\Lambda^2(\overline{\Omega})$ -solutions of the equation

$$\Delta_\lambda u + f(u) = 0, \tag{2.6}$$

where $f : \mathbb{R} \rightarrow \mathbb{R}$ is a continuous function. In what follows we let $F(u) := \int_0^u f(t)dt$ and we agree to call $\Lambda^2(\overline{\Omega})$ -solution of (2.6) every function $u \in \Lambda^2(\overline{\Omega})$ satisfying

$$\sum_{j=1}^N (X_j^2 u)(x) + f(u)(x) = 0 \quad \forall x \in \Omega.$$

Theorem 2.2. *Let $u \in \Lambda^2(\overline{\Omega})$ be a solution of (2.6). Then*

$$\begin{aligned} \int_{\Omega} \left(F(u) + \left(\frac{1}{Q} - \frac{1}{2} \right) uf(u) \right) dx &= \frac{1}{Q} \int_{\partial\Omega} \left(\langle T, \nu \rangle \left(F(u) - \frac{1}{2} |\nabla_\lambda u|^2 \right) \right. \\ &\quad \left. + \langle \nabla_\lambda u, \nu \rangle \left(T(u) + \left(\frac{Q}{2} - 1 \right) u \right) \right) ds. \end{aligned} \tag{2.7}$$

Moreover, if $u = 0$ on $\partial\Omega$

$$\int_{\Omega} \left(F(u) + \left(\frac{1}{Q} - \frac{1}{2} \right) uf(u) \right) dx = \frac{1}{2Q} \int_{\partial\Omega} \left(\frac{\partial u}{\partial \nu} \right)^2 |v_\lambda|^2 \langle T, \nu \rangle ds. \tag{2.8}$$

Proof. Since $Q = \text{div } T$, we have

$$Q \int_{\Omega} F(u) dx = \int_{\Omega} \text{div}(T)F(u) dx.$$

On the other hand, if $u \in C^2(\overline{\Omega})$, an integration by parts gives

$$\int_{\Omega} \text{div}(T)F(u) dx = \int_{\partial\Omega} \langle T, \nu \rangle F(u) ds - \int_{\Omega} T(u)f(u) dx$$

so that

$$Q \int_{\partial\Omega} F(u) dx = \int_{\partial\Omega} \langle T, \nu \rangle F(u) ds + \int_{\Omega} T(u)\Delta_\lambda u dx. \tag{2.9}$$

On the other hand,

$$\begin{aligned} \int_{\Omega} uf(u) dx &= - \int_{\Omega} (\lambda_i^2 \partial_{x_i}^2 u) u dx \\ &= - \int_{\partial\Omega} \lambda_i^2 \partial_{x_i} u \nu_i u ds + \int_{\Omega} \lambda_i^2 (\partial_{x_i} u)^2 dx, \end{aligned}$$

that is

$$\int_{\Omega} uf(u) dx = - \int_{\partial\Omega} u \langle \nabla_\lambda u, \nu_\lambda \rangle ds + \int_{\Omega} |\nabla_\lambda u|^2 dx. \tag{2.10}$$

Identities (2.9) and (2.10) can be extended to functions $u \in \Lambda^2(\overline{\Omega})$ by using an approximation argument as in the proof of the previous theorem. Using (2.9) and (2.10) in (2.2) one gets (2.7). When $u = 0$ on $\partial\Omega$ (2.7) becomes

$$\int_{\Omega} \left(F(u) + \left(\frac{1}{Q} - \frac{1}{2} \right) uf(u) \right) dx = \frac{1}{Q} \int_{\partial\Omega} \left(-\frac{1}{2} \langle T, \nu \rangle |\nabla_\lambda u|^2 + T(u) \langle \nabla_\lambda u, \nu_\lambda \rangle \right) ds. \tag{2.11}$$

We now remark that $u \in C^1(\overline{\Omega} \setminus \Pi)$, since $\lambda_j \partial_{x_j} u \in C^1(\overline{\Omega})$ and λ_j is C^1 and different from zero in $\mathbb{R}^N \setminus \Pi$. Thus, the condition $u = 0$ on $\partial\Omega$ implies

$$\partial_{x_j} u = \left(\frac{\partial u}{\partial \nu} \right) \nu_j \quad \text{at any point of } \partial\Omega \setminus \Pi$$

for every $j \in \{1, \dots, N\}$. As a consequence, on $\partial\Omega \setminus \Gamma$ we have

$$\begin{aligned} -\frac{1}{2}\langle T, \nu \rangle |\nabla_\lambda u|^2 + T(u) \langle \nabla_\lambda u, \nu_\lambda \rangle &= -\frac{1}{2} \varepsilon_i x_i \nu_i \lambda_i^2 (\partial_{x_j} u)^2 + \varepsilon_i x_i \partial_{x_i} u \lambda_j^2 \partial_{x_j} u \nu_j \\ &= -\frac{1}{2} \varepsilon_i x_i \nu_i \lambda_j^2 \left(\frac{\partial u}{\partial \nu}\right)^2 \nu_j^2 + \varepsilon_i x_i \left(\frac{\partial u}{\partial \nu}\right)^2 \nu_i \lambda_j^2 \nu_j^2 \\ &= \frac{1}{2} \langle T, \nu \rangle |\nu_\lambda|^2 \left(\frac{\partial u}{\partial \nu}\right)^2. \end{aligned}$$

Replacing this identity in (2.11) we obtain (2.8). \square

Our nonexistence results follow from the previous theorem. To proceed, we need the definition of δ_t -starshaped domain.

Definition 2.3. We say that Ω is δ_t -starshaped with respect to the origin if $0 \in \Omega$ and

$$\langle T, \nu \rangle \geq 0 \quad \text{at every point of } \partial\Omega.$$

Theorem 2.4. Let Ω be δ_t -starshaped with respect to the origin. Then the problem

$$\Delta_\lambda u + f(u) = 0 \quad \text{in } \Omega, \quad u|_{\partial\Omega} = 0 \tag{2.12}$$

has no solution $u \in \Lambda^2(\overline{\Omega})$, $u \not\equiv 0$, if

$$F(s) + \left(\frac{1}{Q} - \frac{1}{2}\right) sf(s) < 0 \tag{2.13}$$

for every $s \neq 0$. If (2.13) holds for every $s > 0$ then (2.12) has no nonnegative solution $u \not\equiv 0$.

Proof. Let $u \in \Lambda^2(\overline{\Omega})$ be a solution to (2.12). Since $\langle T, \nu \rangle \geq 0$ at any point of $\partial\Omega$ identity (2.8) implies

$$\int_\Omega \left(F(u) + \left(\frac{1}{Q} - \frac{1}{2}\right) uf(u) \right) dx \geq 0.$$

Using condition (2.13) in this inequality, one gets $u = 0$ a.e. in Ω , completing the proof of the theorem. \square

Condition (2.13) can be weakened to obtain a nonexistence result of nonnegative nontrivial solutions to (2.12). For this we need a unique continuation theorem that seems to have an interest in its own right.

Proposition 2.5. Let $u \in \dot{W}_\lambda^{1,2}(\Omega)$ be a weak nonnegative solution of $\Delta_\lambda u + cu = 0$ in Ω , with $c \in L_{\text{loc}}^{\frac{Q}{2}}(\Omega)$.² If there exists $x_0 \in \Omega$ such that

$$\int_{B_d(x_0,r)} u(x) dx = O(r^k) \quad \text{as } r \rightarrow 0, \text{ for every } k \in \mathbb{N},$$

then $u \equiv 0$ in the connected component of Ω containing x_0 .

Proof. The proof of this proposition follows exactly the same lines as the one of Corollary A.1 in [13]. We only stress that Jerison’s Poincaré inequality in [15], used in page 96 in [13], has to be replaced by the one related to our vector fields $\lambda_1 \partial_{x_1}, \dots, \lambda_N \partial_{x_N}$. We omit any further details. \square

We are ready to state our second nonexistence theorem.

Theorem 2.6. Let Ω be connected and δ_t -starshaped with respect to the origin. Then problem (2.12) has no nonnegative solution $u \in \Lambda^2(\overline{\Omega})$, $u \not\equiv 0$, if f is locally Lipschitz continuous, $f(0) = 0$ and

$$F(s) + \left(\frac{1}{Q} - \frac{1}{2}\right) sf(s) \leq 0 \quad \text{for every } s > 0. \tag{2.14}$$

² We refer to the next section for the definition of weak solutions.

Proof. Suppose $u \in A^2(\overline{\Omega})$, $u \geq 0$, is a solution of (2.12). Then (2.14) and identity (2.8) imply

$$\int_{\partial\Omega} \left(\frac{\partial u}{\partial \nu}\right)^2 |v_\lambda|^2 \langle T, \nu \rangle ds \leq 0.$$

The assumption $\langle T, \nu \rangle \geq 0$ on $\partial\Omega$ implies

$$\left(\frac{\partial u}{\partial \nu}\right)^2 \langle T, \nu \rangle = 0 \quad \text{at any point of } \partial\Omega.$$

On the other hand, since Ω is bounded, it must be $\langle T, \nu \rangle > 0$ on $V \cap (\partial\Omega \setminus \Gamma)$, for a suitable bounded and connected open set $V \subseteq \mathbb{R}^N \setminus \Gamma$. Thus

$$Du = \frac{\partial u}{\partial \nu} \equiv 0 \quad \text{in } V \cap (\partial\Omega \setminus \Gamma).$$

Setting $u \equiv 0$ in $(\mathbb{R}^N \setminus \Omega) \cap V$ we then obtain a weak solution to

$$\Delta_\lambda u + cu = 0 \quad \text{in } V,$$

where $c = \frac{f(u)}{u}$ where $u \neq 0$ and $c = 0$ otherwise. Thus, since f is locally Lipschitz, $c \in L^\infty(V)$ and, by Proposition 2.5 $u \equiv 0$ in V . A connectedness argument yields $u \equiv 0$ in Ω . \square

3. Existence and regularity results

In this section we aim to prove the existence of weak solutions to problem (1.4) under natural assumptions on the nonlinearity f . The second aim of the section is to provide some regularity results for the weak solutions.

Our existence result is the following one.

Theorem 3.1. *Let $\Omega \subseteq \mathbb{R}^N$ be a bounded open set and let $f : \overline{\Omega} \times \mathbb{R} \rightarrow \mathbb{R}$ be a continuous function satisfying*

$$f(x, u) = o(u) \quad \text{as } u \rightarrow 0 \quad \text{and} \quad f(x, u) = o\left(|u|^{\frac{Q+2}{Q-2}}\right) \quad \text{as } |u| \rightarrow \infty \tag{3.1}$$

uniformly with respect to $x \in \overline{\Omega}$. Let

$$F(x, u) := \int_0^u f(x, s) ds$$

and assume the existence of two constants $\mu > 2$ and $\mu_0 > 0$ such that

$$uf(x, u) \geq \mu F(x, u) \quad \text{for } |u| > \mu_0 \quad \text{and} \quad F(x, u) > 0 \quad \text{for } u > \mu_0. \tag{3.2}$$

Then, there exists a nonnegative weak solution $u \in \dot{W}_\lambda^{1,2}(\Omega)$, $u \neq 0$ to the equation

$$\Delta_\lambda u - \eta u + f(x, u) = 0 \tag{3.3}$$

for every $\eta \geq 0$.

In (3.1) the number Q stands for the homogeneous dimension of \mathbb{R}^N with respect to $(\delta_t)_{t>0}$, see (1.6). By a weak solution to (3.3) we mean a critical point of the functional

$$J(u) = \int_\Omega (|\nabla_\lambda u|^2 + \eta u^2 - F(x, u)) dx, \quad u \in \dot{W}_\lambda^{1,2}(\Omega). \tag{3.4}$$

The proof of Theorem 3.1 is a standard application of Ambrosetti–Rabinowitz’s Mountain Pass theorem and of the following compact embedding result.

Proposition 3.2. *Let $\Omega \subseteq \mathbb{R}^N$ be a bounded open set. Then the embedding*

$$\dot{W}_\lambda^{1,2}(\Omega) \hookrightarrow L^p(\Omega)$$

is compact for every $p \in [1, 2_\lambda^*[, 2_\lambda^* = \frac{2Q}{Q-2}$.

This proposition is a consequence of the Poincaré inequality (1.20), the doubling property of the d_λ -balls and a general result from analysis in metric spaces, see, e.g., the survey paper [16]. However, it also easily follows from next embedding theorem, which is of interest in its own right. Indeed, it can be seen as a regularity result, in terms of classical Sobolev spaces, for the functions $u \in \dot{W}_\lambda^{1,p}(\Omega)$.

Theorem 3.3. Let Ω be a bounded open subset of \mathbb{R}^N . Then, for every $p \in]1, \infty[$, we have

$$\dot{W}_\lambda^p(\Omega) \hookrightarrow \dot{W}_{\varepsilon_1, \dots, \varepsilon_N}^{\frac{1}{\varepsilon_1}, \dots, \frac{1}{\varepsilon_N}, p}(\Omega). \tag{3.5}$$

We denote by $\dot{W}_{\varepsilon_1, \dots, \varepsilon_N}^{\frac{1}{\varepsilon_1}, \dots, \frac{1}{\varepsilon_N}, p}(\Omega)$ the closure of $C_0^1(\Omega)$ with respect to the norm

$$\|u\|_{\frac{1}{\varepsilon}, p} := \sum_{j=1}^N \|u\|_{\frac{1}{\varepsilon_j}, p}$$

where

$$\|u\|_{\frac{1}{\varepsilon_j}, p} := \int_{\mathbb{R}^N} |\partial_{x_j} u|^p dx \quad \text{if } \varepsilon_j = 1$$

while

$$\|u\|_{\frac{1}{\varepsilon_j}, p} := \int_0^1 \left(\int_{\mathbb{R}^N} \frac{|u(x + se_j) - u(x)|^p}{s^{1+\frac{p}{\varepsilon_j}}} dx \right) ds \quad \text{if } \varepsilon_j > 1,$$

where, as usual, $e_j = (0, \dots, \underset{j}{1}, \dots, 0)$. From an embedding theorem for classical anisotropic Sobolev-type spaces of fractional orders, we know that

$$\dot{W}_{\varepsilon_1, \dots, \varepsilon_N}^{\frac{1}{\varepsilon_1}, \dots, \frac{1}{\varepsilon_N}, p}(\Omega) \hookrightarrow L^q(\Omega) \tag{3.6}$$

if

$$1 < p < \varepsilon_1 + \varepsilon_2 + \dots + \varepsilon_N = Q, \quad \frac{1}{p} - \frac{1}{q} = \frac{1}{Q} \quad \text{i.e., } q = \frac{pQ}{Q-p} := p_\lambda^*,$$

and the embedding is compact if

$$q < p_\lambda^*,$$

see, e.g., [17]. As a consequence, by Theorem 3.3, Proposition 3.2 immediately follows.

Proof of Theorem 3.3. To prove embedding (3.5) we use a family of inequalities proved in [2, Theorem 2.6]. To begin with, we remark that embedding (3.5) is equivalent to the inequality

$$\|u\|_{\frac{1}{\varepsilon_j}, p}^p \leq C \sum_{i=1}^N \int_{\Omega} |\lambda_i \partial_{x_i} u|^p dx, \quad j = 1, 2, \dots, p, \tag{3.7}$$

for every $u \in C_0^\infty(\Omega)$ and $C = C(\Omega, \lambda, p) > 0$. If $\varepsilon_j = 1$ inequality (3.7) is trivial because, in this case, $\lambda_j \equiv 1$ (see Remark 1.1) so that

$$\|u\|_{\frac{1}{\varepsilon_j}, p}^p = \int_{\Omega} |\lambda_j \partial_{x_j} u|^p dx.$$

Then, we can only consider the case $\varepsilon_j > 1$ and use inequality (2.6.b), page 1242, in [2]. From that inequality we obtain, for a suitable positive constant $C = C(\Omega, \lambda, p)$,

$$\int_0^1 \left(\int_{\mathbb{R}^N} \frac{|u(x + se_j) - u(x)|^p}{s(\varphi_j(s))^p} dx \right) ds \leq C \|u\|_{\lambda, p}, \quad \forall u \in C_0^\infty(\Omega), \tag{3.8}$$

where $\varphi_j(s) = \phi_j(0, s)$ and ϕ_j is defined in (1.9).³ Thus, to complete the proof of the theorem, it is enough to show that

$$\varphi_j(s) = c_j s^{\frac{1}{\varepsilon_j}} \tag{3.9}$$

for a suitable constant $c_j > 0$, and for every $j = 1, \dots, N$. Indeed, by Proposition A.2 in the Appendix we have

$$F_j(0, r) = a_j r^{\varepsilon_j}, \quad j = 1, \dots, N$$

³ We would like to explicitly remark that the assumption (2.6.a) of Theorem 2.6 in [2] is satisfied. Indeed, due to the δ_λ -homogeneity of λ_j , we have

$$(\varepsilon_j - 1)\lambda_j = T(\lambda_j) \leq \left(\max_{1 \leq i \leq N} \varepsilon_i \right) \langle \nabla \lambda_j, x \rangle.$$

for a suitable constant $a_i > 0$. Therefore

$$\varphi_j(s) = \left(\frac{1}{a_j} s\right)^{\frac{1}{\theta_j}}, \quad j \in \{1, \dots, N\}$$

and the proof is complete. \square

Proof of Theorem 3.1. Using assumption (3.1) and Proposition 3.2 we can show that the functional J satisfies the Palais–Smale condition. With an equally standard argument we can show, using (3.1) and (3.2) that J satisfies all the other assumptions of the Mountain Pass Theorem, see e.g. [18,19]. We explicitly remark that the argument based on the Maximum Principle for Elliptic Equation used in [18,19] in order to get nonnegative solutions, also works in the present context. Indeed, Δ_λ satisfies a weak and a strong maximum principle, thanks to Theorem 3.1 and the nonhomogeneous Harnack inequality of Theorem 5.5 in [20]. \square

We close this section by stating a regularity theorem which can be proved verbatim as Theorem 4.1 in [13].

Theorem 3.4. Let $\Omega \subseteq \mathbb{R}^N$ be a bounded open set and let $u \in \dot{W}_\lambda^{1,2}(\Omega)$ be a weak solution to (3.3), with f satisfying assumption (3.1). Then $u \in L^p(\Omega)$ for every $p \in [2, \infty[$.

This theorem, together with the nonhomogeneous Harnack inequality of Theorem 5.5 in [20], leads to the following corollary.

Corollary 3.5. Let the assumptions of Theorem 3.4 be satisfied and let $u \in \dot{W}_\lambda^{1,2}(\Omega)$ be a weak solution to (3.3). Then u is locally Hölder continuous in Ω . Precisely, $u \in L^\infty_{\text{loc}}(\Omega)$ and there exists $\theta \in]0, 1[$ and $C > 0$ such that

$$|u(x) - u(y)| \leq C \sup_{B_d(x_0, 2r)} |u| (d(x, y))^\theta$$

for every λ -ball $B_d(x_0, 2r) \subseteq \Omega$ and for every $x, y \in B_d(x_0, r)$. The constants θ and C are independent of u and of $B_d(x_0, 2r)$.

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Appendix

In this Appendix we prove some propositions regarding the functions F_i 's and ϕ_i 's from which important properties of d_λ and of $|B_d(x, r)|$ easily follow. First of all, we prove the following.

Proposition A.1. For every $x \in \mathbb{R}_+^N$ and for every $r > 0$ we have

$$F_i(x, 2r) \leq 2^{\varepsilon_i} F_i(x, r), \quad i = 1, \dots, N. \tag{A.1}$$

Proof. We argue by induction on the index $i \in \{1, \dots, N\}$. Inequality (A.1) is trivial if $i = 1$. Assume it holds for $j \leq i$ and prove it holds for $j = i + 1$. Indeed, since λ_{i+1} is increasing with respect to its argument

$$\begin{aligned} F_{i+1}(x, 2r) &= 2r\lambda_{i+1}(x_1 + F_1(x, 2r), \dots, x_i + F_i(x, 2r)) \\ &\leq 2r\lambda_{i+1}(x_1 + 2r, \dots, x_i + 2^{\varepsilon_i} F_i(x, r)) \\ &\leq 2r\lambda_{i+1}(2(x_1 + r), \dots, 2^{\varepsilon_i}(x_i + F_i(x, r))) \\ &= (\text{by the } \delta_t\text{-homogeneity of } \lambda_{i+1}) 2^{\varepsilon_i+1} r\lambda_{i+1}(x_1 + F_1(x, r), \dots, x_i + F_i(x, r)) \\ &= 2^{\varepsilon_i+1} F_{i+1}(x, r). \end{aligned}$$

This completes the proof. \square

From (A.1), with a standard argument, we get

$$F_i(x, R) \leq 2^{\varepsilon_i} \left(\frac{R}{r}\right)^{\varepsilon_i} F_i(x, r) \tag{A.2}$$

for $0 < r < R$.

The functions F_i 's take an explicit form at $x = 0$. Indeed, the following proposition holds.

Proposition A.2. For every $r > 0$ we have

$$F_i(0, r) = a_i r^{\varepsilon_i}, \quad i = 1, \dots, N$$

for suitable constant $a_i > 0$.

Proof. We know that $F_1(0, \tau) = \tau$, $F_j(0, \tau) = \tau \lambda_j(F_1(0, \tau), \dots, F_{j-1}(0, \tau))$ for $2 \leq j \leq N$. Using the δ_t -homogeneity of the λ_j 's, we obtain

$$F_2(0, \tau) = \tau \lambda_2(F_1(0, \tau)) = \tau \lambda_2(\tau) = \tau^{\varepsilon_2} \lambda_2(1) =: a_2 \tau^{\varepsilon_2}.$$

Moreover

$$F_3(0, \tau) = \tau \lambda_3(F_1(0, \tau), F_2(0, \tau)) = \tau \lambda_3(\tau, a_2 \tau^{\varepsilon_2}) = \tau^{\varepsilon_3} \lambda_3(1, a_2) =: a_3 \tau^{\varepsilon_3}.$$

An easy iteration argument shows that

$$F_j(0, \tau) = a_j \tau^{\varepsilon_j} \quad \text{for every } j \in \{1, \dots, N\}$$

with the a_j 's given by the recurrence formula:

$$a_1 = 1, \quad a_j = \lambda_j(a_1, \dots, a_{j-1}), \quad j \in \{2, \dots, N\}. \quad \square$$

From the previous proposition and inequalities (1.10) we immediately get the following corollary.

Corollary A.3. For every $x, y \in \mathbb{R}^N$ we have

$$d_\lambda(x, y) \leq \sum_{i=1}^N c_i^* |x_i - y_i|, \quad c_i^* = c_2 \left(\frac{1}{a_i} \right)^{\frac{1}{\varepsilon_i}}, \quad i = 1, \dots, N$$

for every $x = (x_1, \dots, x_N)$ and $y = (y_1, \dots, y_N) \in \mathbb{R}^N$.

To obtain a local lower estimate of d_λ we need one more proposition.

Proposition A.4. For every compact set $K \subseteq \mathbb{R}^N$ there exists a constant $d_i = d_i(K) > 0$ such that

$$F_i(x, r) \leq d_i r \quad \forall x \in K, \forall r \in [0, 1], \quad i = 1, \dots, N. \quad (\text{A.3})$$

Proof. The inequality (A.3) is obvious if $i = 1$. Suppose (A.3) is satisfied for $i \leq j$. Then, for $0 < r < 1$,

$$\begin{aligned} F_{j+1}(x, r) &= r \lambda_{j+1}(x_1 + F_1(x, r), \dots, x_j + F_j(x, r)) \\ &\leq r \lambda_{j+1}(x_1 + d_1 r, \dots, x_j + d_j r) \\ &\leq r \sup_{x \in K} \lambda_{j+1}(x_1 + d_1, \dots, x_j + d_j) =: d_{j+1} r. \end{aligned}$$

This completes the proof. \square

From (A.3) we get $F_i\left(x, \frac{r}{d_i}\right) \leq r$ for $r \leq d_i$. Hence $\frac{r}{d_i} \leq \phi_i(x, r)$ for $0 < r < d_i$. As a consequence, by (1.10),

$$d_\lambda(x, y) \geq c_1 \sum_{i=1}^N \phi_i(x^*, |x_i - y_i|) \geq c_1 \sum_{i=1}^N \frac{|x_i - y_i|}{d_i}.$$

Therefore, for a suitable constant d^* , we have

$$d_\lambda(x, y) \geq d^* |x - y|, \quad \forall x \in K, \forall y \in \mathbb{R}^N : |x - y| \leq d^*. \quad (\text{A.4})$$

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