# Moderate solutions of semilinear elliptic equations with Hardy potential 

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#### Abstract

Let $\Omega$ be a bounded smooth domain in $\mathbb{R}^{N}$. We study positive solutions of equation (E) $-L_{\mu} u+u^{q}=0$ in $\Omega$ where $L_{\mu}=\Delta+\frac{\mu}{\delta^{2}}$, $0<\mu, q>1$ and $\delta(x)=\operatorname{dist}(x, \partial \Omega)$. A positive solution of ( E ) is moderate if it is dominated by an $L_{\mu}$-harmonic function. If $\mu<C_{H}(\Omega)$ (the Hardy constant for $\Omega$ ) every positive $L_{\mu}$-harmonic function can be represented in terms of a finite measure on $\partial \Omega$ via the Martin representation theorem. However the classical measure boundary trace of any such solution is zero. We introduce a notion of normalized boundary trace by which we obtain a complete classification of the positive moderate solutions of ( E ) in the subcritical case, $1<q<q_{\mu, c}$. (The critical value depends only on $N$ and $\mu$.) For $q \geq q_{\mu, c}$ there exists no moderate solution with an isolated singularity on the boundary. The normalized boundary trace and associated boundary value problems are also discussed in detail for the linear operator $L_{\mu}$. These results form the basis for the study of the nonlinear problem.


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## 1. Introduction

In this paper, we investigate boundary value problem with measure data for the following equation

$$
\begin{equation*}
-\Delta u-\frac{\mu}{\delta^{2}} u+u^{q}=0 \tag{1.1}
\end{equation*}
$$

in a $C^{2}$ bounded domain $\Omega$, where $q>1, \mu \in \mathbb{R}$ and $\delta(x)=\operatorname{dist}(x, \partial \Omega)$. This problem is naturally linked to the theory of linear Schrödinger equations $-L^{V} u=0$ where $L^{V}:=\Delta+V$ and the potential $V$ satisfies $|V| \leq c \delta^{-2}$. Such equations have been studied in numerous papers (see [1,2] and the references therein).

[^0]Put

$$
\begin{equation*}
L_{\mu}:=\Delta+\frac{\mu}{\delta^{2}} . \tag{1.2}
\end{equation*}
$$

A solution $u \in L_{l o c}^{1}(\Omega)$ of the equation $-L_{\mu} u=0$ is called an $L_{\mu}$-harmonic function. Similarly, if

$$
-L_{\mu} u \geq 0 \quad \text { or } \quad-L_{\mu} u \leq 0
$$

we say that $u$ is $L_{\mu}$-superharmonic or $L_{\mu}$-subharmonic respectively. If $\mu=0$ we shall just use the terms harmonic, superharmonic, subharmonic.

Some problems involving equations (1.1) and (1.2) with $\mu<1 / 4$ were studied by Bandle, Moroz and Reichel [4]. They derived estimates of local $L_{\mu}$-subharmonic and superharmonic functions and applied these results to study conditions for existence or nonexistence of large solutions of (1.1). They also showed that the classical Keller-Osserman estimate [14,24] remains valid for (1.1).

The condition $\mu<\frac{1}{4}$ is related to Hardy's inequality. Denote by $C_{H}(\Omega)$ the best constant in Hardy's inequality, i.e.,

$$
\begin{equation*}
C_{H}(\Omega)=\inf _{H_{0}^{1}(\Omega)} \frac{\int_{\Omega}|\nabla u|^{2} d x}{\int_{\Omega}(u / \delta)^{2} d x} . \tag{1.3}
\end{equation*}
$$

By Marcus, Mizel and Pinchover [17], $C_{H}(\Omega) \in\left(0, \frac{1}{4}\right]$ and $C_{H}(\Omega)=\frac{1}{4}$ when $\Omega$ is convex. Furthermore the infimum is achieved if and only if $C_{H}(\Omega)<1 / 4$. By Brezis and Marcus [7], for every $\mu<1 / 4$ there exists a unique number $\lambda_{\mu, 1}$ such that

$$
\mu=\inf _{H_{0}^{1}(\Omega)} \frac{\int_{\Omega}\left(|\nabla u|^{2}-\lambda_{\mu, 1} u^{2}\right) d x}{\int_{\Omega}(u / \delta)^{2} d x}
$$

and the infimum is achieved. Thus $\lambda_{\mu, 1}$ is an eigenvalue of $-L_{\mu}$ and, by [7, Lemma 2.1], it is a simple eigenvalue. We denote by $\varphi_{\mu, 1}$ the corresponding positive eigenfunction normalized by $\int_{\Omega}\left(\varphi_{\mu, 1}^{2} / \delta^{2}\right) d x=1$.

The mapping $[1 / 4, \infty) \ni \mu \mapsto \lambda_{\mu, 1}$ is strictly decreasing. Therefore if $\mu<C_{H}(\Omega)$ then $\lambda_{\mu, 1}>0$. Consequently, in this case, $\varphi_{\mu, 1}$ is a positive supersolution of $-L_{\mu}$. This fact and a classical result of Ancona [2] imply that for every $y \in \partial \Omega$, there exists a positive $L_{\mu}$-harmonic function in $\Omega$ which vanishes on $\partial \Omega \backslash\{y\}$ and is unique up to a constant. Denote this function by $K_{\mu}^{\Omega}(\cdot, y)$, normalized by setting it equal to 1 at a fixed reference point $x_{0} \in \Omega$. The function $(x, y) \mapsto K_{\mu}^{\Omega}(x, y),(x, y) \in \Omega \times \partial \Omega$, is the $L_{\mu}$-Martin kernel in $\Omega$ relative to $x_{0}$. Further, by [2]:

Representation Theorem. For every $v \in \mathfrak{M}^{+}(\partial \Omega)$ the function

$$
\begin{equation*}
\mathbb{K}_{\mu}^{\Omega}[\nu](x):=\int_{\partial \Omega} K_{\mu}^{\Omega}(x, y) d \nu(y) \quad \forall x \in \Omega \tag{1.4}
\end{equation*}
$$

is $L_{\mu}$-harmonic, i.e., $L_{\mu} \mathbb{K}_{\mu}^{\Omega}[\nu]=0$. Conversely, for every positive $L_{\mu}$-harmonic function $u$ there exists a unique measure $v \in \mathfrak{M}^{+}(\partial \Omega)$ such that $u=\mathbb{K}_{\mu}^{\Omega}[\nu]$.

This theorem implies that - in the present case - the $L_{\mu}$-Martin boundary of $\Omega$ coincides with the Euclidean boundary. (For the general definition of Martin boundary see, e.g. [1]. However this notion will not be used here beyond the representation theorem stated above.) The measure $\nu$ such that $u=\mathbb{K}_{\mu}^{\Omega}[\nu]$ is called the $L_{\mu}$-boundary measure of $u$. If $\mu=0, v$ is equivalent to the classical measure boundary trace of $u$ (see Definition 1.1). But if $0<\mu<C_{H}(\Omega)$, it can be shown that, for every $v \in \mathfrak{M}^{+}(\partial \Omega)$, the measure boundary trace of $\mathbb{K}_{\mu}^{\Omega}[\nu]$ is zero (see Corollary 2.11 below).

In the case $\mu=0$, the boundary value problem

$$
\begin{align*}
-\Delta u+|u|^{q-1} u & =0 & & \text { in } \Omega \\
u & =v & & \text { on } \partial \Omega \tag{1.5}
\end{align*}
$$

where $q>1$ and $v$ is either a finite measure or a positive (possibly unbounded) measure, has been studied by numerous authors. Following Brezis [6], if $v$ is a finite measure, a weak solution of (1.5) is defined as follows: $u$ is a solution of the problem if $u$ and $\delta|u|^{q}$ are integrable in $\Omega$ and

$$
\begin{equation*}
\int_{\Omega}\left(-u \Delta \zeta+|u|^{q-1} u \zeta\right) d x=-\int_{\partial \Omega} \frac{\partial \zeta}{\partial \mathbf{n}} d \nu \quad \forall \zeta \in C_{0}^{2}(\bar{\Omega}) \tag{1.6}
\end{equation*}
$$

where $\mathbf{n}$ is the outer unit normal on $\partial \Omega$. Brezis proved that, if a solution exists then it is unique. Gmira and Véron [13] showed that there exists a critical exponent, $q_{c}:=\frac{N+1}{N-1}$, such that if $1<q<q_{c}$, (1.6) has a weak solution for every finite measure $v$ but, if $q \geq q_{c}$ there exists no positive solution with isolated point singularity.

Marcus and Véron [20] proved that every positive solution of the equation

$$
\begin{equation*}
-\Delta u+u^{q}=0 \tag{1.7}
\end{equation*}
$$

possesses a boundary trace given by a positive measure $v$, not necessarily bounded. In the subcritical case the blow-up set of the trace is a closed set. Furthermore they showed that, in this case, for every such positive measure $v$, the boundary value problem (1.5) has a unique solution.

In the case $q=2, N=2$ this result was previously proved by Le Gall [15] using a probabilistic definition of the boundary trace.

In the supercritical case the problem turned out to be much more challenging. It was studied by several authors using various techniques. The problem was studied by Le Gall, Dynkin, Kuznetsov, Mselati a.o. employing mainly probabilistic methods. Consequently the results applied only to $1<q \leq 2$. In parallel it was studied by Marcus and Veron employing purely analytic methods that were not subject to the restriction $q \leq 2$. A complete classification of the positive solutions of (1.5) in terms of their behavior at the boundary was provided by Mselati [18] for $q=2$, by Dynkin [11] for $q_{c} \leq q \leq 2$ and finally by Marcus [16] for every $q \geq q_{c}$. For details and related results we refer the reader to $[23,22,21,3,10]$ and the references therein.

In the case of equation (1.1) one is faced by the problem that, according to the classical definition of measure boundary trace, every positive $L_{\mu}$-harmonic function has measure boundary trace zero. Therefore, in order to classify the positive solutions of (1.1) in terms of their behavior at the boundary, it is necessary to introduce a different notion of trace. As in the study of (1.7), we first consider the question of boundary trace for positive $L_{\mu}$-harmonic or superharmonic functions.

We recall the classical definition of measure boundary trace.
Definition 1.1. (i) A sequence $\left\{D_{n}\right\}$ is a $C^{2}$ exhaustion of $\Omega$ if for every $n, D_{n}$ is of class $C^{2}, \bar{D}_{n} \subset D_{n+1}$ and $\cup_{n} D_{n}=\Omega$. If the domains are uniformly of class $C^{2}$ we say that $\left\{D_{n}\right\}$ is a uniform $C^{2}$ exhaustion.
(ii) Let $u \in W_{l o c}^{1, p}(\Omega)$ for some $p>1$. We say that $u$ possesses a measure boundary trace on $\partial \Omega$ if there exists a finite measure $v$ on $\partial \Omega$ such that, for every uniform $C^{2}$ exhaustion $\left\{D_{n}\right\}$ and every $\varphi \in C(\bar{\Omega})$,

$$
\left.\lim _{n \rightarrow \infty} \int_{\partial D_{n}} u\right|_{\partial D_{n}} \varphi d S=\int_{\partial \Omega} \varphi d \nu .
$$

Here $\left.u\right|_{D_{n}}$ denotes the Sobolev trace. The measure boundary trace of $u$ is denoted by $\operatorname{tr}(u)$.
For $\beta>0$, denote

$$
\Omega_{\beta}=\{x \in \Omega: \delta(x)<\beta\}, D_{\beta}=\{x \in \Omega: \delta(x)>\beta\}, \Sigma_{\beta}=\{x \in \Omega: \delta(x)=\beta\}
$$

Put

$$
\begin{equation*}
\alpha_{ \pm}:=\frac{1}{2} \pm \sqrt{\frac{1}{4}-\mu} \tag{1.8}
\end{equation*}
$$

It can be shown (see Corollary 2.11 below) that the classical measure boundary trace of $\mathbb{K}_{\mu}^{\Omega}[\nu]$ is zero but there exist constants $C_{1}, C_{2}$ such that, for every $v \in \mathfrak{M}(\partial \Omega)$,

$$
\begin{equation*}
C_{1}\|\nu\|_{\mathfrak{M}(\partial \Omega)} \leq \frac{1}{\beta^{\alpha_{-}}} \int_{\Sigma_{\beta}} \mathbb{K}_{\mu}^{\Omega}[\nu](x) d S_{x} \leq C_{2}\|\nu\|_{\mathfrak{M}(\partial \Omega)} \tag{1.9}
\end{equation*}
$$

for all $\beta \in\left(0, \beta_{0}\right)$ where $\beta_{0}>0$ depends only on $\Omega$. In view of this we introduce the following definition of trace.
Definition 1.2. A positive function $u$ possesses a normalized boundary trace if there exists a measure $v \in \mathfrak{M}^{+}(\partial \Omega)$ such that

$$
\begin{equation*}
\lim _{\beta \rightarrow 0} \frac{1}{\beta^{\alpha_{-}}} \int_{\Sigma_{\beta}}\left|u-\mathbb{K}_{\mu}^{\Omega}[\nu]\right| d S_{x}=0 \tag{1.10}
\end{equation*}
$$

The normalized boundary trace will be denoted by $\operatorname{tr}^{*}(u)$.
Remark. The notion of normalized boundary trace is well defined. Indeed, suppose that $v$ and $v^{\prime}$ satisfy (1.10). Put $v=$ $\left(\mathbb{K}_{\mu}^{\Omega}\left[\nu-v^{\prime}\right]\right)_{+}$then $v$ is a nonnegative $L_{\mu}$-subharmonic function, $v \leq \mathbb{K}\left[\nu+v^{\prime}\right]$ and $\operatorname{tr}^{*}(v)=0$. By Proposition 2.14, $v=0$, i.e., $\mathbb{K}_{\mu}^{\Omega}\left[\nu-\nu^{\prime}\right] \leq 0$. By interchanging the roles of $v$ and $\nu^{\prime}$, we deduce that $\mathbb{K}_{\mu}^{\Omega}\left[\nu^{\prime}-v\right] \leq 0$. Thus $v=v^{\prime}$.

Denote by $G_{\mu}^{\Omega}$ the Green function of $-L_{\mu}$ in $\Omega$ and, for every positive Radon measure $\tau$ in $\Omega$, put

$$
\mathbb{G}_{\mu}^{\Omega}[\tau](x):=\int_{\Omega} G_{\mu}^{\Omega}(x, y) d \tau(y)
$$

Denote by $\mathfrak{M}_{f}(\Omega), f$ a positive Borel function in $\Omega$, the space of Radon measures $\tau$ on $\Omega$ satisfying $\int_{\Omega} f d|\tau|<\infty$ and by $\mathfrak{M}_{f}^{+}(\Omega)$ the positive cone of this space.

If $\tau$ is a positive measure such that $\mathbb{G}_{\mu}^{\Omega}[\tau](x)<\infty$ for some $x \in \Omega$ then $\tau \in \mathfrak{M}_{\delta^{\alpha}}(\Omega)$ and $\mathbb{G}_{\mu}^{\Omega}[\tau]$ is finite everywhere in $\Omega$. The underlying reason for this is the behavior of the Green function at the boundary: for every $\beta>0$ there exists $c_{\beta}$ such that

$$
c_{\beta}^{-1} \delta(x)^{\alpha_{+}} \leq G_{\mu}^{\Omega}(x, y) \leq c_{\beta} \delta(x)^{\alpha_{+}} \quad \forall x \in \Omega_{\beta / 2}, y \in D_{\beta} .
$$

For details see Section 2.2 below.
We begin with the study of the linear boundary value problem,

$$
\begin{align*}
& -L_{\mu} u=\tau \quad \text { in } \Omega \\
& \operatorname{tr}^{*}(u)=v, \tag{1.11}
\end{align*}
$$

where $v \in \mathfrak{M}^{+}(\partial \Omega)$ and $\tau \in \mathfrak{M}_{\delta^{\alpha}}^{+}(\Omega)$. As usual we look for solutions $u \in L_{l o c}^{1}(\Omega)$ and the equation is understood in the sense of distributions. The representation theorem implies that if $\tau=0$ the problem has a unique solution, $u=\mathbb{K}_{\mu}^{\Omega}[\nu]$.

We list below our main results regarding this problem.

## Proposition I.

(i) If $u$ is a non-negative $L_{\mu}$-harmonic function and $\operatorname{tr}^{*}(u)=0$ then $u=0$.
(ii) If $\tau \in \mathfrak{M}_{\delta^{\alpha}+}^{+}(\Omega)$ then $\mathbb{G}_{\mu}^{\Omega}[\tau]$ has normalized trace zero. Thus $\mathbb{G}_{\mu}^{\Omega}[\tau]$ is a solution of (1.11) with $\nu=0$.
(iii) Let $u$ be a positive $L_{\mu}$-subharmonic function. If $u$ is dominated by an $L_{\mu}$-superharmonic function then $L_{\mu} u \in$ $\mathfrak{M}_{\delta^{\alpha}+}^{+}(\Omega)$ and $u$ has a normalized boundary trace. In this case $\operatorname{tr}^{*}(u)=0$ if and only if $u \equiv 0$.
(iv) Let $u$ be a positive $L_{\mu}$-superharmonic function. Then there exist $v \in \mathfrak{M}^{+}(\partial \Omega)$ and $\tau \in \mathfrak{M}_{\delta^{\alpha+}}^{+}(\Omega)$ such that

$$
\begin{equation*}
u=\mathbb{G}_{\mu}^{\Omega}[\tau]+\mathbb{K}_{\mu}^{\Omega}[\nu] . \tag{1.12}
\end{equation*}
$$

In particular, $u$ is an $L_{\mu}$-potential (i.e., $u$ does not dominate any positive $L_{\mu}$-harmonic function) if and only if $\operatorname{tr}^{*}(u)=0$.
(v) For every $\nu \in \mathfrak{M}^{+}(\partial \Omega)$ and $\tau \in \mathfrak{M}_{\delta^{\alpha}}^{+}(\Omega)$, problem (1.11) has a unique solution. The solution is given by (1.12).

Next we study the nonlinear boundary value problem,

$$
\begin{align*}
-L_{\mu} u+u^{q} & =0 \quad \text { in } \Omega \\
\operatorname{tr}^{*}(u) & =v \tag{1.13}
\end{align*}
$$

where $v \in \mathfrak{M}^{+}(\partial \Omega)$.
Definition 1.3. (i) A positive solution of (1.1) is $L_{\mu}$-moderate if it is dominated by an $L_{\mu}$-harmonic function.
(ii) A positive function $u \in L_{l o c}^{q}(\Omega)$ is a (weak) solution of (1.13) if it satisfies the equation (in the sense of distributions) and has normalized boundary trace $\nu$.

Definition 1.4. Put

$$
X(\Omega)=\left\{\zeta \in C^{2}(\Omega): \delta^{\alpha_{-}} L_{\mu} \zeta \in L^{\infty}(\Omega), \delta^{-\alpha_{+}} \zeta \in L^{\infty}(\Omega)\right\}
$$

A function $\zeta \in X(\Omega)$ is called an admissible test function for (1.13).
Following are our main results concerning the nonlinear problem (1.13). Theorems A-D apply to arbitrary exponent $q>1$.

Theorem A. Assume that $0<\mu<C_{H}(\Omega), q>1$. Let u be a positive solution of (1.1). Then the following statements are equivalent:
(i) $u$ is $L_{\mu}$-moderate.
(ii) $u$ admits a normalized boundary trace $v \in \mathfrak{M}^{+}(\partial \Omega)$. In other words, $u$ is a solution of (1.13).
(iii) $u \in L_{\delta^{\alpha+}}^{q}(\Omega)$ and

$$
\begin{equation*}
u+\mathbb{G}_{\mu}^{\Omega}\left[u^{q}\right]=\mathbb{K}_{\mu}^{\Omega}[\nu] \tag{1.14}
\end{equation*}
$$

where $v=\operatorname{tr}^{*}(u)$.
Furthermore, a positive function $u$ is a solution of (1.13) if and only if $u / \delta^{\alpha_{-}} \in L^{1}(\Omega), \delta^{\alpha_{+}} u^{q} \in L^{1}(\Omega)$ and

$$
\begin{equation*}
\int_{\Omega}\left(-u L_{\mu} \zeta+u^{q} \zeta\right) d x=-\int_{\Omega} \mathbb{K}_{\mu}^{\Omega}[\nu] L_{\mu} \zeta d x \quad \forall \zeta \in X(\Omega) \tag{1.15}
\end{equation*}
$$

Theorem B. Assume $0<\mu<C_{H}(\Omega), q>1$.
I. UniQueness. For every $v \in \mathfrak{M}^{+}(\partial \Omega)$, there exists at most one positive solution of (1.13).
II. Monotonicity. Assume $\nu_{i} \in \mathfrak{M}^{+}(\partial \Omega), i=1,2$. Let $u_{v_{i}}$ be the unique solution of (1.13) with $v$ replaced by $\nu_{i}$, $i=1$, 2 . If $\nu_{1} \leq \nu_{2}$ then $u_{\nu_{1}} \leq u_{\nu_{2}}$.
III. A-PRIORI ESTIMATE. There exists a positive constant $c=c(N, \mu, \Omega)$ such that every positive solution u of (1.13) satisfies,

$$
\begin{equation*}
\|u\|_{L_{\delta^{-\alpha_{-}}}^{1}(\Omega)}+\|u\|_{L_{\delta^{+}}^{q}}(\Omega) \leq c\|v\|_{\mathfrak{M}(\partial \Omega)} . \tag{1.16}
\end{equation*}
$$

Theorem C. Assume $0<\mu<C_{H}(\Omega), q>1$. If $v \in \mathfrak{M}^{+}(\partial \Omega)$ and $\mathbb{K}_{\mu}^{\Omega}[\nu] \in L_{\delta^{\alpha+}}^{q}(\Omega)$ then there exists a unique solution of the boundary value problem (1.13).

Corollary C1. For every positive function $f \in L^{1}(\partial \Omega)(1.13)$ with $v=f$ admits a unique positive solution.
Theorem D. Assume $0<\mu<C_{H}(\Omega), q>1$. If u is a positive solution of (1.13) then

$$
\begin{equation*}
\lim _{x \rightarrow y} \frac{u(x)}{\mathbb{K}_{\mu}^{\Omega}[\nu](x)}=1 \quad \text { non-tangentially, v-a.e. on } \partial \Omega . \tag{1.17}
\end{equation*}
$$

Let

$$
\begin{equation*}
q_{\mu, c}:=\frac{N+\alpha_{+}}{N-1-\alpha_{-}} . \tag{1.18}
\end{equation*}
$$

In the next two results we show, among other things, that $q_{\mu, c}$ is the critical exponent for (1.13). This means that, if $1<q<q_{\mu, c}$ then problem (1.13) has a unique solution for every measure $v \in \mathfrak{M}^{+}(\partial \Omega)$ but, if $q \geq q_{\mu, c}$ then the problem has no solution for some measures $v$, e.g. Dirac measure.

In Theorem E we consider the subcritical case $1<q<q_{\mu, c}$ and in Theorem F the supercritical case.
Theorem E. Assume $0<\mu<C_{H}(\Omega)$ and $1<q<q_{\mu, c}$. Then:
I. EXISTENCE AND UNIQUENESS. For every $\nu \in \mathfrak{M}^{+}(\partial \Omega)$ (1.13) admits a unique positive solution $u_{\nu}$.
II. Stability. If $\left\{v_{n}\right\}$ is a sequence of measures in $\mathfrak{M}^{+}(\partial \Omega)$ weakly convergent to $v \in \mathfrak{M}^{+}(\partial \Omega)$ then $u_{\nu_{n}} \rightarrow u_{\nu}$ in $L_{\delta^{-\alpha_{-}}}^{1}(\Omega)$ and in $L_{\delta^{\alpha_{+}}}^{q}(\Omega)$.
iiI. Local behavior. Let $v=k \delta_{y}$, where $k>0$ and $\delta_{y}$ is the Dirac measure concentrated at $y \in \partial \Omega$. Then, under the assumptions of Theorem $E$, the unique solution of (1.13), denoted by $u_{k \delta_{y}}$, satisfies

$$
\begin{equation*}
\lim _{x \rightarrow y} \frac{u_{k \delta_{y}}(x)}{K_{\mu}^{\Omega}(x, y)}=k \tag{1.19}
\end{equation*}
$$

Remark. Note that in part III we have 'uniform convergence' not just 'non-tangential convergence' as in Theorem D.
Theorem F. Assume $0<\mu<C_{H}(\Omega)$ and $q \geq q_{\mu, c}$. Then for every $k>0$ and $y \in \partial \Omega$, there is no positive solution of (1.1) with normalized boundary trace $k \delta_{y}$.

In the first part of the paper we study properties of positive $L_{\mu}$-harmonic functions and the boundary value problem (1.11). In the second part, these results are applied to a study of the corresponding boundary value problem for the nonlinear equation (1.1). These results yield a complete classification of the positive moderate solutions of (1.1) in the subcritical case. They also provide a framework for the study of positive solutions of (1.1) that may blow up at some parts of the boundary. The existence of such solutions in the subcritical case has been studied (by different methods) in [5]. The boundary trace for positive non-moderate solutions and corresponding boundary value problems will be treated in a forthcoming paper.

The main ingredients used in this paper are: the Representation Theorem previously stated and other basic results of potential theory (see [1]), a sharp estimate of the Green kernel of $-L_{\mu}$ due to Filippas, Moschini and Tertikas [9], estimates for convolutions in weak $L^{p}$ spaces (see [23, Section 2.3.2]) and the comparison principle obtained in [4].

## 2. The linear equation

Throughout this paper we assume that $0<\mu<C_{H}(\Omega)$.

### 2.1. Some potential theoretic results

We denote by $\mathfrak{M}_{\delta^{\alpha}}(\Omega), \alpha \in \mathbb{R}$, the space of Radon measures $\tau$ on $\Omega$ satisfying $\int_{\Omega} \delta^{\alpha}(x) d|\tau|<\infty$ and by $\mathfrak{M}_{\delta^{\alpha}}^{+}(\Omega)$ the positive cone of $\mathfrak{M}_{\delta^{\alpha}}(\Omega)$. When $\alpha=0$, we use the notation $\mathfrak{M}(\Omega)$ and $\mathfrak{M}^{+}(\Omega)$. We also denote by $\mathfrak{M}(\partial \Omega)$ the space of finite Radon measures on $\partial \Omega$ and by $\mathfrak{M}^{+}(\partial \Omega)$ the positive cone of $\mathfrak{M}(\partial \Omega)$.

Let $D$ be a $C^{2}$ domain such that $D \Subset \Omega$ and $h \in L^{1}(\partial \Omega)$. Denote by $\mathbb{S}_{\mu}(D, h)$ the solution of the problem

$$
\left\{\begin{align*}
-L_{\mu} u & =0 & & \text { in } D  \tag{2.1}\\
u & =h & & \text { on } \partial D .
\end{align*}\right.
$$

Lemma 2.1. Let u be $L_{\mu}$-superharmonic in $\Omega$ and $D$ be a $C^{2}$ domain such that $D \Subset \Omega$. Then $u \geq \mathbb{S}_{\mu}(D, u)$ a.e. in $D$.

Proof. Since $u$ is $L_{\mu}$-superharmonic in $\Omega$, there exists $\tau \in \mathfrak{M}^{+}(\Omega)$ such that $-L_{\mu} u=\tau$. Let $v$ be the solution of

$$
\left\{\begin{align*}
-L_{\mu} v & =\tau & & \text { in } D  \tag{2.2}\\
v & =0 & & \text { on } \partial D
\end{align*}\right.
$$

and put $w=\mathbb{S}_{\mu}(D, u)$. Then $w \geq 0$ and $\left.u\right|_{D}=v+w \geq v$.
Lemma 2.2. Let u be a nonnegative $L_{\mu}$-superharmonic and $\left\{D_{n}\right\}$ be a $C^{2}$ exhaustion of $\Omega$. Then

$$
\hat{u}:=\lim _{n \rightarrow \infty} \mathbb{S}_{\mu}\left(D_{n}, u\right)
$$

exists and is the largest $L_{\mu}$-harmonic function dominated by $u$.
Proof. By Lemma 2.1, $\mathbb{S}_{\mu}\left(D_{n}, u\right) \leq\left. u\right|_{D_{n}}$, hence the sequence $\left\{\mathbb{S}_{\mu}\left(D_{n}, u\right)\right\}$ is decreasing. Consequently, $\hat{u}$ exists and is an $L_{\mu}$-harmonic function dominated by $u$. Next, if $v$ is an $L_{\mu}$-harmonic function dominated by $u$ then $v \leq$ $\mathbb{S}_{\mu}\left(D_{n}, u\right)$ for every $n \in \mathbb{N}$. Letting $n \rightarrow \infty$ yields $v \leq \hat{u}$.

Definition 2.3. A nonnegative $L_{\mu}$-superharmonic function is called an $L_{\mu}$-potential if its largest $L_{\mu}$-harmonic minorant is zero.

As a consequence of Lemma 2.2, we obtain
Lemma 2.4. Let $u_{p}$ be a nonnegative $L_{\mu}$-superharmonic function in $\Omega$. If for some $C^{2}$ exhaustion $\left\{D_{n}\right\}$ of $\Omega$,

$$
\begin{equation*}
\lim _{n \rightarrow \infty} \mathbb{S}_{\mu}\left(D_{n}, u_{p}\right)=0 \tag{2.3}
\end{equation*}
$$

then $u_{p}$ is an $L_{\mu}$-potential in $\Omega$. Conversely, if $u_{p}$ is an $L_{\mu}$-potential, then (2.3) holds for every $C^{2}$ exhaustion $\left\{D_{n}\right\}$ of $\Omega$.

For easy reference we quote below the Riesz decomposition theorem (see [1]).
Theorem 2.5. Every nonnegative $L_{\mu}$-superharmonic function $u$ in $\Omega$ can be written in a unique way in the form $u=u_{p}+u_{h}$ where $u_{p}$ is an $L_{\mu}$-potential and $u_{h}$ is a nonnegative $L_{\mu}$-harmonic function in $\Omega$.

The next result is a consequence of the Fatou convergence theorem [1, Theorem 1.8] and the following well-known fact: if a function satisfies the local Harnack inequality, fine convergence at the boundary (in the sense of [1]) implies non-tangential convergence.

Theorem 2.6. Let $u_{p}$ be a positive $L_{\mu}$-potential and $u$ be a positive $L_{\mu}$-harmonic function. Assume that $\frac{u_{p}}{u}$ satisfies the Harnack inequality. Then

$$
\lim _{x \rightarrow y} \frac{u_{p}(x)}{u(x)}=0 \quad \text { non-tangentially, v-a.e. on } \partial \Omega
$$

where $v$ is the $L_{\mu}$-boundary measure of $u$.

### 2.2. The action of the Green and Martin kernels on spaces of measures

From [2], for every $y \in \partial \Omega$, there exists a positive $L_{\mu}$-harmonic function in $\Omega$ which vanishes on $\partial \Omega \backslash\{y\}$. When normalized, this function is unique. We choose a fixed reference point $x_{0}$ in $\Omega$ and denote by $K_{\mu, y}^{\Omega}$ this $L_{\mu}$-harmonic function, normalized by $K_{\mu, y}^{\Omega}\left(x_{0}\right)=1$. The function $K_{\mu}^{\Omega}(\cdot, y)=K_{\mu, y}^{\Omega}(\cdot)$ is the $L_{\mu}$-Martin kernel in $\Omega$, normalized at $x_{0}$.

For $v \in \mathfrak{M}(\partial \Omega)$ denote

$$
\mathbb{K}_{\mu}^{\Omega}[\nu](x)=\int_{\partial \Omega} K_{\mu}^{\Omega}(x, y) d \nu(y)
$$

In what follows the notation $f \sim g$ means: there exists a positive constant $c$ such that $c^{-1} f<g<c f$ in the domain of the two functions or in a specified subset of this domain. Of course, in the latter case, the constant depends on the subset.

Let $G_{\mu}^{\Omega}$ be the Green kernel for the operator $L_{\mu}$ in $\Omega \times \Omega$. Fix a point $x_{0} \in \Omega$. It is well known that the function $x \mapsto G_{\mu}^{\Omega}\left(x, x_{0}\right)$ behaves like the first eigenfunction $\varphi_{\mu, 1}(x)$ near the boundary, i.e., $G_{\mu}^{\Omega}\left(\cdot, x_{0}\right) \sim \varphi_{\mu, 1}$ in $\Omega_{\beta},(0<$ $\beta<\delta\left(x_{0}\right)$ ).

By [19, Lemmas 5,1, 5.2] (see also [8, Lemma 7] for an alternative proof)

$$
\begin{equation*}
c^{-1} \delta(x)^{\alpha_{+}} \leq \varphi_{\mu, 1}(x) \leq c \delta(x)^{\alpha_{+}} . \tag{2.4}
\end{equation*}
$$

Thus, if $0<\beta<\delta\left(x_{0}\right)$,

$$
\begin{equation*}
c_{\beta}^{-1} \delta(x)^{\alpha_{+}} \leq G_{\mu}^{\Omega}\left(x, x_{0}\right) \leq c_{\beta} \delta(x)^{\alpha_{+}} \quad \forall x \in \Omega_{\beta} . \tag{2.5}
\end{equation*}
$$

Therefore, if $\tau \in \mathfrak{M}_{\delta^{\alpha+}}(\Omega)$ then

$$
\mathbb{G}_{\mu}^{\Omega}[\tau](x):=\int_{\Omega} G_{\mu}^{\Omega}(x, y) d \tau(y)<\infty \quad \text { a.e. in } \Omega .
$$

Indeed, by (2.5) and the symmetry of the Green kernel, for every $x \in \Omega$, the integral over $\Omega_{\delta(x) / 2}$ is finite. For $y \in D_{\delta(x) / 4}, G_{\mu}^{\Omega}(x, y) \leq c|x-y|^{2-N}$. Therefore the integral is finite over this set as well. Inequality (2.5) also implies that, if $\tau$ is a positive Radon measure in $\Omega$ and $\mathbb{G}_{\mu}^{\Omega}[\tau](x)<\infty$ for some point $x \in \Omega$ then $\tau \in \mathfrak{M}_{\delta^{\alpha+}}(\Omega)$ and $\mathbb{G}_{\mu}^{\Omega}[\tau]$ is finite everywhere in $\Omega$.

By [9, Theorem 4.11], for every $x, y \in \Omega, x \neq y$,

$$
\begin{equation*}
G_{\mu}^{\Omega}(x, y) \sim \min \left\{|x-y|^{2-N}, \delta(x)^{\alpha_{+}} \delta(y)^{\alpha_{+}}|x-y|^{2 \alpha_{-}-N}\right\} \tag{2.6}
\end{equation*}
$$

Since

$$
K_{\mu}^{\Omega}(x, y):=\lim _{z \rightarrow y} \frac{G_{\mu}^{\Omega}(x, z)}{G_{\mu}^{\Omega}\left(x_{0}, z\right)} \quad \forall x \in \Omega
$$

it follows from (2.6) that

$$
\begin{equation*}
K_{\mu}^{\Omega}(x, y) \sim \delta(x)^{\alpha_{+}}|x-y|^{2 \alpha_{-}-N} \quad \forall x \in \Omega, y \in \partial \Omega . \tag{2.7}
\end{equation*}
$$

Let $G^{\Omega}=G_{0}^{\Omega}$ and $P^{\Omega}=P_{0}^{\Omega}$ denote the Green and Poisson kernels of $-\Delta$ in $\Omega$. Then, by (2.7)

$$
\begin{equation*}
\frac{K_{\mu}^{\Omega}(x, y)}{\delta(x)^{\alpha_{-}}} \sim \frac{\delta(x)}{|x-y|^{N}}\left(\frac{|x-y|}{\delta(x)}\right)^{2 \alpha_{-}} \sim P^{\Omega}(x, y)\left(\frac{|x-y|}{\delta(x)}\right)^{2 \alpha_{-}} . \tag{2.8}
\end{equation*}
$$

Denote $L_{w}^{p}(\Omega ; \tau), 1 \leq p<\infty, \tau \in \mathfrak{M}^{+}(\Omega)$, the weak $L^{p}$ space defined as follows: a measurable function $f$ in $\Omega$ belongs to this space if there exists a constant $c$ such that

$$
\begin{equation*}
\lambda_{f}(a ; \tau):=\tau(\{x \in \Omega:|f(x)|>a\}) \leq c a^{-p}, \quad \forall a>0 . \tag{2.9}
\end{equation*}
$$

The function $\lambda_{f}$ is called the distribution function of $f$ (relative to $\tau$ ). For $p \geq 1$, denote

$$
L_{w}^{p}(\Omega ; \tau)=\left\{f \text { Borel measurable }: \sup _{a>0} a^{p} \lambda_{f}(a ; \tau)<\infty\right\}
$$

and

$$
\begin{equation*}
\|f\|_{L_{w}^{p}(\Omega ; \tau)}^{*}=\left(\sup _{a>0} a^{p} \lambda_{f}(a ; \tau)\right)^{\frac{1}{p}} . \tag{2.10}
\end{equation*}
$$

This expression is not a norm, but for $p>1$, it is equivalent to the norm

$$
\begin{equation*}
\|f\|_{L_{w}^{p}(\Omega ; \tau)}=\sup \left\{\frac{\int_{\omega}|f| d \tau}{\tau(\omega)^{1 / p^{\prime}}}: \omega \subset \Omega, \omega \text { measurable }, 0<\tau(\omega)\right\} . \tag{2.11}
\end{equation*}
$$

More precisely,

$$
\begin{equation*}
\|f\|_{L_{w}^{p}(\Omega ; \tau)}^{*} \leq\|f\|_{L_{w}^{p}(\Omega ; \tau)} \leq \frac{p}{p-1}\|f\|_{L_{w}^{p}(\Omega ; \tau)}^{*} . \tag{2.12}
\end{equation*}
$$

Notice that, for every $\alpha>-1$,

$$
L_{w}^{p}\left(\Omega ; \delta^{\alpha} d x\right) \subset L_{\delta^{\alpha}}^{r}(\Omega) \quad \forall r \in[1, p) .
$$

For every $x \in \partial \Omega$, denote by $\mathbf{n}_{x}$ the outward unit normal vector to $\partial \Omega$ at $x$.
The following is a well-known geometric property of $C^{2}$ domains.
Proposition 2.7. There exists $\beta_{0}>0$ such that
(i) For every point $x \in \bar{\Omega}_{\beta_{0}}$, there exists a unique point $\sigma_{x} \in \partial \Omega$ such that $\left|x-\sigma_{x}\right|=\delta(x)$. This implies $x=$ $\sigma_{x}-\delta(x) \mathbf{n}_{\sigma_{x}}$.
(ii) The mappings $x \mapsto \delta(x)$ and $x \mapsto \sigma_{x}$ belong to $C^{2}\left(\bar{\Omega}_{\beta_{0}}\right)$ and $C^{1}\left(\bar{\Omega}_{\beta_{0}}\right)$ respectively. Furthermore, $\lim _{x \rightarrow \sigma(x)} \nabla \delta(x)=-\mathbf{n}_{x}$.

By combining (2.6), (2.7) and [23, Lemma 2.3.2], we obtain
Proposition 2.8. There exist constants $c_{i}$ depending only on $N, \mu, \beta, \Omega$ such that,

$$
\begin{align*}
& \left\|\mathbb{G}_{\mu}^{\Omega}[\tau]\right\|_{L_{w}^{N+\beta}\left(\Omega, \delta^{\beta}\right)}^{N-c_{1}}\|\tau\|_{\mathfrak{M}(\Omega)}, \quad \forall \tau \in \mathfrak{M}(\Omega), \beta>-1,  \tag{2.13}\\
& \left\|\mathbb{G}_{\mu}^{\Omega}[\tau]\right\|_{L_{w}^{N-2 \alpha_{-}}}^{\left(\Omega, \delta^{\left.\beta-\alpha_{+}\right)}\right.} \leq c_{1}\|\tau\|_{\mathfrak{M}_{\delta^{\alpha}+}(\Omega)}, \quad \forall \tau \in \mathfrak{M}_{\delta^{\alpha+}}(\Omega), \beta>-2 \alpha_{-},  \tag{2.14}\\
& \left\|\mathbb{K}_{\mu}^{\Omega}[\nu]\right\|_{L_{w}^{N-1-\alpha_{-}}}^{\left(\Omega, \delta^{\beta}\right)} \tag{2.15}
\end{align*} \leq c_{2}\|\nu\|_{\mathfrak{M}(\partial \Omega)}, \quad \forall v \in \mathfrak{M}(\partial \Omega), \beta>-1 .
$$

Proof. We assume that $\tau$ is positive; otherwise we replace $\tau$ by $|\tau|$. We consider $\tau$ as a positive measure in $\mathbb{R}^{N}$ by extending $\tau$ by zero outside of $\Omega$. For $a \in(0, N)$, denote $\Gamma_{a}(x)=|x|^{a-N}$. By [23, inequality (2.3.17)],

$$
\begin{equation*}
\left\|\Gamma_{a} * \tau\right\|_{L_{w}^{N+\beta}\left(\Omega, \delta^{\beta}\right)} \leq c\|\tau\|_{\mathfrak{M}(\Omega)} \quad \forall \beta>\max \{-1,-a\} \tag{2.16}
\end{equation*}
$$

where $c=c(N, a, \beta, \operatorname{diam}(\Omega))$. By (2.6),

$$
G_{\mu}^{\Omega}(x, y) \leq c \min \left\{\Gamma_{2}(x-y), \delta(x)^{\alpha_{+}} \delta(y)^{\alpha_{+}} \Gamma_{2 \alpha_{-}}(x-y)\right\} .
$$

Hence, by (2.16),

$$
\begin{aligned}
\left\|\mathbb{G}_{\mu}^{\Omega}[\tau]\right\|_{L_{w}^{N+2}\left(\Omega, \delta^{\beta}\right)}^{N+\beta} & \leq c\left\|\Gamma_{2} * \tau\right\|_{L_{w}^{N-2}}^{L^{N+\beta}}\left(\Omega, \delta^{\beta}\right) \\
& \leq c^{\prime}\|\tau\|_{\mathfrak{M}(\Omega)} \quad \forall \beta>-1, \\
\left\|\mathbb{G}_{\mu}^{\Omega}[\tau]\right\|_{L_{w}^{N-2 \alpha_{-}}}\left(\Omega, \delta^{\left.\beta-\alpha_{+}\right)}\right. & \leq c\left\|\Gamma_{2 \alpha_{-}} *\left(\delta^{\alpha_{+}} \tau\right)\right\|_{L_{w}^{N-2 \alpha_{-}}}\left(\Omega, \delta^{\beta}\right) \\
& \leq c\|\tau\|_{M_{\delta^{\alpha+}}(\Omega)} \quad \forall \beta>-2 \alpha_{-} .
\end{aligned}
$$

Next we extend $\nu$ by zero outside $\partial \Omega$ and observe that, by (2.7), $K_{\mu}^{\Omega}(x, y) \leq c \Gamma_{1+\alpha_{-}}(x-y)$. Hence $\mathbb{K}_{\mu}^{\Omega}[\nu] \leq c \Gamma_{1+\alpha_{-}} *$ $\nu$ and by (2.16),

$$
\left\|\mathbb{K}_{\mu}^{\Omega}[\nu]\right\|_{L_{w}^{N-1-\alpha_{-}}\left(\Omega, \delta^{\beta}\right)} \leq c\left\|\Gamma_{1+\alpha_{-}} * \nu\right\|_{L_{w}^{N-1-\alpha_{-}}\left(\Omega, \delta^{\beta}\right)}^{N+\beta} \leq v \|_{\mathfrak{M}(\partial \Omega)} \quad \forall \beta>-1
$$

Corollary 2.9. Let $\beta>-1$.
(i) If $\left\{\nu_{n}\right\} \subset \mathfrak{M}^{+}(\partial \Omega)$ converges weakly to $v \in \mathfrak{M}^{+}(\partial \Omega)$ then $\left\{\mathbb{K}_{\mu}^{\Omega}\left[\nu_{n}\right]\right\}$ converges to $\mathbb{K}_{\mu}^{\Omega}[\nu]$ in $L_{\delta^{\beta}}^{p}(\Omega)$ for every $p$ such that $1 \leq p<\frac{N+\beta}{N-1-\alpha_{-}}$.
(ii) If $\left\{\tau_{n}\right\} \subset \mathfrak{M}^{+}(\Omega)$ converges weakly (relative to $C_{0}(\bar{\Omega})$ ) to $\tau \in \mathfrak{M}^{+}(\Omega)$ then $\left\{\mathbb{G}_{\mu}^{\Omega}\left[\tau_{n}\right]\right\}$ converges to $\mathbb{G}_{\mu}^{\Omega}[\tau]$ in $L_{\delta \beta}^{p}(\Omega)$ for every $p$ such that $1 \leq p<\frac{N+\beta}{N-2}$.

Proof. We prove the first statement. The second is proved in a similar way.

Since $K_{\mu}^{\Omega}(x,.) \in C(\partial \Omega)$ for every $x \in \Omega,\left\{\mathbb{K}_{\mu}^{\Omega}\left[\nu_{n}\right]\right\}$ converges to $\mathbb{K}_{\mu}^{\Omega}[\nu]$ every where in $\Omega$. By Holder inequality and (2.15), we deduce that $\left\{\left(\mathbb{K}_{\mu}^{\Omega}\left[v_{n}\right]\right)^{p}\right\}$ is equi-integrable w.r.t. $\delta^{\beta} d x$ for any $1 \leq p<\frac{N+\beta}{N-1-\alpha_{-}}$. By Vitali's theorem, $\mathbb{K}_{\mu}^{\Omega}\left[\nu_{n}\right] \rightarrow \mathbb{K}_{\mu}^{\Omega}[\nu]$ in $L_{\delta^{\beta}}^{p}(\Omega)$.

### 2.3. Estimates related to the normalized trace

Proposition 2.10. There exist positive constants $C_{1}, C_{2}$ such that, for every $\beta \in\left(0, \beta_{0}\right)$,

$$
\begin{equation*}
C_{1} \beta^{\alpha_{-}} \leq \int_{\Sigma_{\beta}} K_{\mu}^{\Omega}(x, y) d S_{x} \leq C_{2} \beta^{\alpha_{-}} \quad \forall y \in \partial \Omega \tag{2.17}
\end{equation*}
$$

The constants $C_{1}, C_{2}$ depend on $N, \Omega, \mu$ but not on $y$.
Furthermore, for every $r_{0}>0$,

$$
\begin{equation*}
\lim _{\beta \rightarrow 0} \frac{1}{\beta^{\alpha_{-}}} \int_{\Sigma_{\beta} \backslash B r_{0}(y)} K_{\mu}^{\Omega}(x, y) d S_{x}=0 \quad \forall y \in \partial \Omega . \tag{2.18}
\end{equation*}
$$

For $r_{0}$ fixed, the rate of convergence is independent of $y$.
Proof. By (2.7),

$$
\begin{equation*}
\frac{1}{\beta^{\alpha_{-}}} \int_{\Sigma_{\beta} \backslash B_{r_{0}}(y)} K_{\mu}^{\Omega}(x, y) d S_{x} \leq c \beta^{\alpha_{+}-\alpha_{-}} . \tag{2.19}
\end{equation*}
$$

This implies (2.18).
For the next estimate it is convenient to assume that the coordinates are placed so that $y=0$ and the tangent hyperplane to $\partial \Omega$ at 0 is $x_{N}=0$ with the $x_{N}$ axis pointing into the domain. For $x \in \mathbb{R}^{N}$ put $x^{\prime}=\left(x_{1}, \cdots, x_{N-1}\right)$. Pick $r_{0} \in\left(0, \beta_{0}\right)$ sufficiently small (depending only on the $C^{2}$ characteristic of $\Omega$ ) so that

$$
\frac{1}{2}\left(\left|x^{\prime}\right|^{2}+\delta(x)^{2}\right) \leq|x|^{2} \quad \forall x \in \Omega \cap B_{r_{0}}(0) .
$$

Then, if $x \in \Sigma_{\beta} \cap B_{r_{0}}(0)=: \Sigma_{\beta, 0}$,

$$
\frac{1}{4}\left(\left|x^{\prime}\right|+\beta\right) \leq|x| .
$$

This inequality and (2.7) imply,

$$
\begin{aligned}
\int_{\Sigma_{\beta, 0}} K_{\mu}^{\Omega}(x, 0) d S_{x} & \leq c_{0} \beta^{\alpha_{+}} \int_{\Sigma_{\beta, 0}}\left(\left|x^{\prime}\right|+\beta\right)^{2 \alpha_{-}-N} d S_{x} \\
& \leq c_{1} \beta^{\alpha_{+}} \int_{\left|x^{\prime}\right|<r_{0}}\left(\left|x^{\prime}\right|+\beta\right)^{2 \alpha_{-}-N} d x^{\prime} \\
& \leq c_{2} \beta^{\alpha_{+}} \int_{0}^{r_{0}}(t+\beta)^{2 \alpha_{-}-2} d t \\
& <c_{2} \beta^{\alpha_{-}} \int_{1}^{\infty} \tau^{-2 \alpha_{+}} d \tau=\frac{c_{2}}{2 \alpha_{+}-1} \beta^{\alpha_{-}} .
\end{aligned}
$$

Thus, for $\beta<r_{0}$,

$$
\begin{equation*}
\frac{1}{\beta^{\alpha_{-}}} \int_{\Sigma_{\beta, 0}} K_{\mu}^{\Omega}(x, 0) d S_{x} \leq \frac{c_{2}}{2 \alpha_{+}-1} . \tag{2.20}
\end{equation*}
$$

Estimates (2.19) and (2.20) imply the second estimate in (2.17). The first estimate in (2.17) follows from (2.8).
Since (2.17) holds uniformly w.r. to $y \in \partial \Omega$, an application of Fubini's yields the following.
Corollary 2.11. For every $v \in \mathfrak{M}^{+}(\partial \Omega)$,

$$
\begin{align*}
C_{1}\|\nu\|_{\mathfrak{M}(\partial \Omega)} & \leq \liminf _{\beta \rightarrow 0} \int_{\Sigma_{\beta}} \frac{\mathbb{K}_{\mu}^{\Omega}[\nu]}{\delta(x)^{\alpha_{-}}} d S_{x} \\
& \leq \limsup _{\beta \rightarrow 0} \int_{\Sigma_{\beta}} \frac{\mathbb{K}_{\mu}^{\Omega}[\nu]}{\delta(x)^{\alpha_{-}}} d S_{x} \leq C_{2}\|\nu\|_{\mathfrak{M}(\partial \Omega)} \tag{2.21}
\end{align*}
$$

with $C_{1}, C_{2}$ as in (2.17).
Proposition 2.12. If $\tau \in \mathfrak{M}_{\delta^{\alpha+}}(\Omega)$ then

$$
\begin{equation*}
\operatorname{tr}^{*}\left(\mathbb{G}_{\mu}^{\Omega}[\tau]\right)=0 \tag{2.22}
\end{equation*}
$$

and, for $0<\beta<\beta_{0}$,

$$
\begin{equation*}
\frac{1}{\beta^{\alpha_{-}}} \int_{\Sigma_{\beta}} \mathbb{G}_{\mu}^{\Omega}[\tau] d S_{x} \leq c \int_{\Omega} \delta^{\alpha_{+}} d|\tau| \tag{2.23}
\end{equation*}
$$

where $c$ is a constant depending on $\mu, \Omega$.
Proof. We may assume that $\tau>0$. Denote $v:=\mathbb{G}_{\mu}^{\Omega}[\tau]$. We start with the proof of (2.23).
By Fubini's theorem and (2.6),

$$
\left.\begin{array}{rl}
\int_{\Sigma_{\beta}} v d S_{x} & \leq c\left(\int_{\Omega} \int_{\Sigma_{\beta} \cap B_{\frac{\beta}{2}}(y)}|x-y|^{2-N} d S_{x} d \tau(y)\right. \\
& \left.\left.+\beta^{\alpha_{+}} \int_{\Omega_{\Sigma_{\beta} \backslash B_{\frac{\beta}{2}}}} \int \right\rvert\, x\right)
\end{array}|x-y|^{2 \alpha_{-}-N} d S_{x} \delta^{\alpha_{+}}(y) d \tau(y)\right)=I_{1}(\beta)+I_{2}(\beta) . .
$$

Note that, if $x \in \Sigma_{\beta}$ and $|x-y| \leq \beta / 2$ then $\beta / 2 \leq \delta(y) \leq 3 \beta / 2$. Therefore

$$
\begin{aligned}
I_{1}(\beta) & \leq c_{1} \int_{\Sigma_{\beta} \cap B_{\frac{\beta}{2}}^{2}}|x-y|^{2-\alpha_{+}-N} d S_{x} \int_{\Omega} \delta(y)^{\alpha_{+}} d \tau(y) \\
& \leq c_{1}^{\prime} \int_{0}^{\beta / 2} r^{2-\alpha_{+}-N_{1}} r^{N-2} d r \int_{\Omega} \delta(y)^{\alpha_{+}} d \tau(y) \\
& \leq c_{1}^{\prime \prime} \beta^{\alpha_{-}} \int_{\Omega} \delta(y)^{\alpha_{+}} d \tau(y)
\end{aligned}
$$

and

$$
I_{2}(\beta) \leq c_{2} \beta^{\alpha_{+}} \int_{\beta / 2}^{\infty} r^{2 \alpha_{-}-N} r^{N-2} d r \int_{\Omega} \delta(y)^{\alpha_{+}} d \tau=c_{2}^{\prime} \beta^{\alpha_{-}} \int_{\Omega} \delta(y)^{\alpha_{+}} d \tau .
$$

This implies (2.23).
Given $\epsilon \in\left(0,\|\tau\|_{\mathfrak{M}_{\delta^{\alpha}+}(\Omega)}\right)$ and $\beta_{1} \in\left(0, \beta_{0}\right)$ put $\tau_{1}=\tau \chi_{\bar{D}_{\beta_{1}}}$ and $\tau_{2}=\tau-\tau_{1}$. Pick $\beta_{1}=\beta_{1}(\epsilon)$ such that

$$
\begin{equation*}
\int_{\Omega_{\beta_{1}}} \delta(y)^{\alpha+} d \tau \leq \epsilon \tag{2.24}
\end{equation*}
$$

Thus the choice of $\beta_{1}$ depends on the rate at which $\int_{\Omega_{\beta}} \delta^{\alpha_{+}} d \tau$ tends to zero as $\beta \rightarrow 0$.
Put $v_{i}=\mathbb{G}_{\mu}^{\Omega}\left[\tau_{i}\right]$. Then, for $0<\beta<\beta_{1} / 2$,

$$
\int_{\Sigma_{\beta}} v_{1} d S_{x} \leq c_{3} \beta^{\alpha_{+}} \beta_{1}^{2 \alpha_{-}-N} \int_{\Omega} \delta^{\alpha_{+}}(y) d \tau_{1}(y) .
$$

Thus,

$$
\begin{equation*}
\lim _{\beta \rightarrow 0} \frac{1}{\beta^{\alpha_{-}}} \int_{\Sigma_{\beta}} v_{1} d S_{x}=0 \tag{2.25}
\end{equation*}
$$

On the other hand, by (2.23) and (2.24),

$$
\begin{equation*}
\frac{1}{\beta^{\alpha_{-}}} \int_{\Sigma_{\beta}} v_{2} d S_{x} \leq c \epsilon \quad \forall \beta<\beta_{0} . \tag{2.26}
\end{equation*}
$$

This implies that $\operatorname{tr}^{*}(v)=0$.
It is well-known that $u$ is an $L_{\mu}$-potential if and only if there exists a positive measure $\tau$ in $\Omega$ such that $u=$ $\mathbb{G}_{\mu}^{\Omega}[\tau]$ (see e.g. [1, Theorem 12]). The estimate (2.6) implies that if $\mathbb{G}_{\mu}^{\Omega}[\tau] \not \equiv \infty$ then $\tau \in \mathfrak{M}_{\delta^{\alpha+}}(\Omega)$. Therefore as a consequence of the previous proposition:

Corollary 2.13. A positive $L_{\mu}$-superharmonic function $u$ is a potential if and only if $\mathrm{tr}^{*}(u)=0$.
Remark. Let $D \Subset \Omega$ be a $C^{2}$ domain and denote by $G_{\mu}^{D}$ and $P_{\mu}^{D}$ the Green and Poisson kernels of $L_{\mu}$ in $D$. (To avoid misunderstanding we point out that, in the formula defining $L_{\mu}, \delta(x)$ denotes, as before, the distance from $x$ to $\partial \Omega$, not to $\partial D$.) As every positive $L_{\mu}$ harmonic function has measure boundary trace zero, there is no Poisson kernel for $L_{\mu}$ in $\Omega$. However, $L_{\mu}$ has a Poisson kernel in every $C^{2}$ domain $D$ strictly contained in $\Omega$. This follows from the fact that the Green kernel $G_{\mu}^{D}$ exists and behaves like $G_{0}^{D}$.

Proposition 2.14. Let $w$ be a non-negative $L_{\mu}$-subharmonic function. If $w$ is dominated by an $L_{\mu}$-superharmonic function then $L_{\mu} w \in \mathfrak{M}_{\delta^{+}}^{+}(\Omega)$ and $w$ has a normalized boundary trace $\nu \in \mathfrak{M}^{+}(\partial \Omega)$. If, in addition, $\operatorname{tr}^{*}(w)=0$ then $w=0$.

Proof. The first assumption implies that there exists a positive Radon measure $\lambda$ in $\Omega$ such that $-L_{\mu} w=-\lambda$.
First assume that $\lambda \in \mathfrak{M}_{\delta^{\alpha}+}(\Omega)$. Then $v:=w+\mathbb{G}_{\mu}^{\Omega}[\lambda]$ is a non-negative $L_{\mu}$-harmonic function and consequently, by the representation theorem, $v=\mathbb{K}_{\mu}^{\Omega}[\nu]$ for some $v \in \mathfrak{M}^{+}(\partial \Omega)$. By Proposition 2.12, $\operatorname{tr}^{*}(w)=v$. If $v=0$ then $v=0$ and therefore $w=0$. Now let us drop the assumption on $\lambda$.

Let $v_{\beta}$ be the unique solution of the boundary value problem,

$$
-L_{\mu} v_{\beta}=-\lambda_{\beta} \text { in } D_{\beta}, \quad v_{\beta}=h_{\beta} \text { on } \partial D_{\beta}
$$

where $\lambda_{\beta}$ is the restriction of $\lambda$ to $D_{\beta}$ and $h_{\beta}$ is the restriction of $w$ to $\partial D_{\beta}$. (The uniqueness follows from [4, Lemma 2.3].) The uniqueness implies that $v_{\beta}=w L_{D_{\beta}}$. By assumption there exists a positive $L_{\mu}$-superharmonic function, say $V$, such that $w \leq V$. Hence

$$
w+\mathbb{G}_{\mu}^{D_{\beta}}\left[\lambda_{\beta}\right]=\mathbb{P}_{\mu}^{D_{\beta}}\left[h_{\beta}\right] \leq \mathbb{P}_{\mu}^{D_{\beta}}\left[V\left\lfloor{ }_{\partial D_{\beta}}\right] \leq V .\right.
$$

This implies that $\mathbb{G}_{\mu}^{\Omega}[\lambda]=\lim _{\beta \rightarrow 0} \mathbb{G}_{\mu}^{D_{\beta}}\left[\lambda_{\beta}\right]<\infty$. For fixed $x \in \Omega, G_{\mu}^{\Omega}(x, y) \sim \delta(y)^{\alpha_{+}}$. Therefore the finiteness of $\mathbb{G}_{\mu}^{\Omega}[\lambda]$ implies that $\lambda \in \mathfrak{M}_{\delta^{\alpha+}}(\Omega)$. By the first part of the proof $w$ has a normalized trace.

Remark. See Proposition 2.20 below for a complementary result.

### 2.4. Test functions

Denote

$$
X(\Omega)=\left\{\zeta \in C^{2}(\Omega): \delta^{\alpha-} L_{\mu} \zeta \in L^{\infty}(\Omega), \delta^{-\alpha} \zeta \in L^{\infty}(\Omega)\right\} .
$$

Proposition 2.15. For any $\zeta \in X(\Omega), \delta^{\alpha}-|\nabla \zeta| \in L^{\infty}(\Omega)$.
Proof. Let $\zeta \in X(\Omega)$ then there exist a positive constant $c_{1}$ and a function $f \in L^{\infty}(\Omega)$ such that $|\zeta| \leq c_{1} \delta^{\alpha+}$ and

$$
-L_{\mu} \zeta=\delta^{-\alpha_{-}} f
$$

Take arbitrary point $x_{*} \in \Omega_{\beta_{0}}$ and put $d_{*}=\frac{1}{2} \delta\left(x_{*}\right), y_{*}=\frac{1}{d_{*}} x_{*}, \zeta_{*}(y)=\zeta\left(d_{*} y\right)$ for $y \in \frac{1}{d_{*}} \Omega_{d_{*}}$. Note that if $x \in B_{d_{*}}\left(x_{*}\right)$ then $y=\frac{1}{d_{*}} x \in B_{1}\left(y_{*}\right)$ and $1 \leq \operatorname{dist}\left(y, \partial\left(\frac{1}{d_{*}} \Omega_{d_{*}}\right)\right) \leq 3$. In $B_{1}\left(y_{*}\right)$,

$$
-\Delta \zeta_{*}-\frac{\mu}{\operatorname{dist}\left(y, \partial\left(\frac{1}{d_{*}} \Omega_{d_{*}}\right)\right)^{2}} \zeta_{*}=d_{*}^{2-\alpha_{-}} \operatorname{dist}\left(y, \partial\left(\frac{1}{d_{*}} \Omega_{d_{*}}\right)\right)^{-\alpha_{-}} f\left(d_{*} y\right)
$$

By local estimate for elliptic equations [12, Theorem 8.32], there exists a positive constant $c_{2}=c_{2}(N, \mu)$ such that

$$
\max _{B_{\frac{1}{2}}\left(y_{*}\right)}\left|\nabla \zeta_{*}\right| \leq c_{2}\left[\max _{B_{1}\left(y_{*}\right)}\left|\zeta_{*}\right|+\max _{B_{1}\left(y_{*}\right)}\left(d_{*}^{2-\alpha_{-}} \operatorname{dist}\left(y, \partial\left(\frac{1}{d_{*}} \Omega_{d_{*}}\right)\right)^{-\alpha_{-}}\left|f\left(d_{*} y\right)\right|\right] .\right.
$$

This implies

$$
d_{*}\left|\nabla \zeta\left(x_{*}\right)\right| \leq c_{3}\left(\delta\left(x_{*}\right)^{\alpha_{+}}+\|f\|_{L^{\infty}(\Omega)} \delta\left(x_{*}\right)^{2-\alpha_{-}}\right),
$$

where $c_{3}=c_{3}\left(N, \mu, c_{1}\right)$. Therefore

$$
|\nabla \zeta(x)| \leq c_{4} \delta(x)^{\alpha_{+}-1} \quad \forall x \in \Omega_{\beta_{0}}
$$

where $c_{4}=c_{4}\left(N, \mu, c_{1},\|f\|_{L^{\infty}(\Omega)}\right)$. Thus $\delta^{-\alpha_{-}}|\nabla \zeta| \in L^{\infty}(\Omega)$.
Definition 2.16. Let $x_{0} \in \Omega$ and denote $\tilde{\beta}\left(x_{0}\right)=\min \left(\beta_{0}, \frac{1}{2} \delta\left(x_{0}\right)\right)$. We say that $\tilde{G}_{\mu}^{\Omega}$ is a proper regularization of $G_{\mu}^{\Omega}$ relative to $x_{0}$ if $\tilde{G}_{\mu}^{\Omega}(x)=G_{\mu}^{\Omega}\left(x_{0}, x\right)$ for $x \in \bar{\Omega}_{\tilde{\beta}\left(x_{0}\right)}, \tilde{G}_{\mu}^{\Omega} \in C^{2}(\Omega) \cap C(\bar{\Omega})$ and $\tilde{G}_{\mu}^{\Omega} \geq 0$ in $\Omega$. Similarly $\tilde{\delta}$ is a proper regularization of $\delta$ relative to $x_{0}$ if $\tilde{\delta}(x)=\delta(x)$ for $x \in \bar{\Omega}_{\tilde{\beta}\left(x_{0}\right)}, \tilde{\delta} \in C^{2}(\bar{\Omega})$ and $\tilde{\delta} \geq 0$ in $\Omega$.

Remark. Using (2.6) and (2.4), it is easily verified that the functions $\varphi_{\mu, 1}, \mathbb{G}_{\mu}^{\Omega}[\eta]$ (for $\eta \in L^{\infty}(\Omega)$ ), $\tilde{G}_{\mu}^{\Omega}$ and $\tilde{\delta}^{\alpha_{+}}$ belong to $X(\Omega)$. Moreover, using Proposition 2.15, one obtains,

$$
\zeta \in X(\Omega) \quad \text { and } \quad h \in C^{2}(\bar{\Omega}) \Longrightarrow h \zeta \in X(\Omega)
$$

In the proofs of the next two propositions we use the following construction. Let $D \Subset \Omega$ be a $C^{2}$ domain. The Green function for $-L_{\mu}$ in $D$ is denoted by $G_{\mu}^{D}$. (To avoid misunderstanding we point out that, in the formula defining $L_{\mu}$, $\delta(x)$ denotes, as before, the distance from $x$ to $\partial \Omega$, not to $\partial D$.) Given $x_{0} \in \Omega$ we construct a family of functions $\mathcal{G}\left(x_{0}\right)=\left\{\tilde{G}_{\mu}^{D_{\beta}}: 0<\beta<\frac{1}{2} \tilde{\beta}\left(x_{0}\right)\right\}$ such that, for each $\beta, \tilde{G}_{\mu}^{D_{\beta}}$ is a proper regularization of $G_{\mu}^{D_{\beta}}\left(x_{0}, \cdot\right)$ in $D_{\beta}$ and $\mathcal{G}\left(x_{0}\right)$ has the following properties:

- For every $\beta \in\left(0, \frac{1}{2} \tilde{\beta}\left(x_{0}\right)\right), \tilde{G}_{\mu}^{D_{\beta}} \in C^{2}\left(\bar{D}_{\beta}\right), \tilde{G}_{\mu}^{D_{\beta}} \geq 0$ and $\tilde{G}_{\mu}^{D_{\beta}}(x)=G_{\mu}^{D_{\beta}}\left(x_{0}, x\right)$ for $x \in D_{\beta} \backslash D_{\tilde{\beta}\left(x_{0}\right)}$.
- The sequences $\left\{\tilde{G}_{\mu}^{D_{\beta}}\right\}$ and $\left\{L_{\mu} \tilde{G}_{\mu}^{D_{\beta}}\right\}$ converge to $\tilde{G}_{\mu}^{\Omega}$ and $L_{\mu} \tilde{G}_{\mu}^{\Omega}$ respectively, as $\beta \rightarrow 0$, a.e. in $\Omega$.
- $\left\|\tilde{G}_{\mu}^{D_{\beta}}+\left|L_{\mu} \tilde{G}_{\mu}^{D_{\beta}}\right|\right\|_{L^{\infty}\left(D_{\beta}\right)} \leq M_{x_{0}}$ where $M_{x_{0}}$ is a positive constant independent of $\beta$.
$\mathcal{G}\left(x_{0}\right)$ will be called a uniform regularization of $\left\{G_{\mu}^{D_{\beta}}\right\}$.
For any function $h \in C^{2}(\partial \Omega)$, we say that $\tilde{h}$ is an admissible extension of $h$ relative to $x_{0}$ in $\bar{\Omega}$ if $\tilde{h}(x)=h(\sigma(x))$ for $x \in \Omega_{\tilde{\beta}\left(x_{0}\right)}$ and $\tilde{h} \in C^{2}(\bar{\Omega})$.


### 2.5. Nonhomogeneous linear equations

Here we discuss the boundary value problem (1.11) in $\Omega$.
Lemma 2.17. Let $u \in L_{\text {loc }}^{1}(\Omega)$ be a positive solution (in the sense of distributions) of equation

$$
\begin{equation*}
-L_{\mu} u=\tau \tag{2.27}
\end{equation*}
$$

in $\Omega$ where $\tau$ is a non-negative Radon measure.
If $\tau \in \mathfrak{M}_{\delta^{\alpha}+}(\Omega)$ then

$$
\begin{equation*}
-\int_{\Omega} \mathbb{G}_{\mu}^{\Omega}[\tau] L_{\mu} \zeta d x=\int_{\Omega} \zeta d \tau \quad \forall \zeta \in X(\Omega) \tag{2.28}
\end{equation*}
$$

Proof. We may assume that $\tau$ is positive. By Proposition 2.12, $\operatorname{tr}^{*}\left(\mathbb{G}_{\mu}^{\Omega}[\tau]\right)=0$. Therefore, given $\varepsilon>0$, there exists $\bar{\beta}=\bar{\beta}(\varepsilon)<\frac{1}{2} \beta_{0}$ such that,

$$
\begin{equation*}
\frac{1}{\beta^{\alpha_{-}}} \int_{\Sigma_{\beta}} \mathbb{G}_{\mu}^{\Omega}[\tau] d S_{x}<\varepsilon \quad \text { and } \quad \int_{\Omega_{\beta}} \delta^{\alpha_{+}} d \tau<\varepsilon \quad \forall \beta \in(0, \bar{\beta}] . \tag{2.29}
\end{equation*}
$$

Let

$$
I(\beta):=\int_{D_{\beta}} \mathbb{G}_{\mu}^{\Omega}[\tau] L_{\mu} \zeta d x+\int_{D_{\beta}} \zeta d \tau
$$

To prove (2.28) we show that

$$
\begin{equation*}
\lim _{\beta \rightarrow 0} I(\beta)=0 \tag{2.30}
\end{equation*}
$$

Put

$$
\tau_{1}:=\chi_{\bar{D}_{\bar{\beta}}} \tau, \quad \tau_{2}:=\chi_{\Omega_{\bar{\beta}}} \tau
$$

and, for $0<\beta<\bar{\beta}$,

$$
I_{k}(\beta):=\int_{D_{\beta}} \mathbb{G}_{\mu}^{\Omega}\left[\tau_{k}\right] L_{\mu} \zeta d x+\int_{D_{\beta}} \zeta d \tau_{k}, \quad k=1,2 .
$$

As $|\zeta| \leq c \delta^{\alpha_{+}}$and $\left|L_{\mu} \zeta\right| \leq \frac{c}{\delta^{\alpha-}},(2.29)$ implies,

$$
\begin{equation*}
\left|I_{2}(\beta)\right| \leq c \varepsilon \quad \forall \beta \in(0, \bar{\beta}) . \tag{2.31}
\end{equation*}
$$

For every $\beta \in(0, \bar{\beta})$,

$$
-\int_{D_{\beta}} \mathbb{G}_{\mu}^{\Omega}\left[\tau_{1}\right] L_{\mu} \zeta d x=\int_{D_{\beta}} \zeta d \tau_{1}+\int_{\Sigma_{\beta}} \frac{\partial \mathbb{G}_{\mu}^{\Omega}\left[\tau_{1}\right]}{\partial \mathbf{n}} \zeta d S_{x}-\int_{\Sigma_{\beta}} \mathbb{G}_{\mu}^{\Omega}\left[\tau_{1}\right] \frac{\partial \zeta}{\partial \mathbf{n}} d S_{x}
$$

Thus

$$
I_{1}(\beta)=-\int_{\Sigma_{\beta}} \frac{\partial \mathbb{G}_{\mu}^{\Omega}\left[\tau_{1}\right]}{\partial \mathbf{n}} \zeta d S_{x}+\int_{\Sigma_{\beta}} \mathbb{G}_{\mu}^{\Omega}\left[\tau_{1}\right] \frac{\partial \zeta}{\partial \mathbf{n}} d S_{x}=: I_{1,1}(\beta)+I_{1,2}(\beta)
$$

By Proposition 2.15 and (2.29),

$$
\begin{equation*}
\left|I_{1,2}(\beta)\right| \leq c \varepsilon \quad \forall \beta \in(0, \bar{\beta}) . \tag{2.32}
\end{equation*}
$$

Next we estimate $I_{1,1}(\beta)$ for $\beta \in(0, \bar{\beta} / 2)$. By Fubini,

$$
\begin{aligned}
I_{1,1}(\beta) & =-\int_{\Sigma_{\beta}} \frac{\partial}{\partial \mathbf{n}_{x}} \int_{D_{\bar{\beta}}} G_{\mu}^{\Omega}(x, y) d \tau_{1}(y) \zeta(x) d S_{x} \\
& =-\int_{D_{\bar{\beta}}} \int_{\Sigma_{\beta}} \frac{\partial G_{\mu}^{\Omega}(x, y)}{\partial \mathbf{n}_{x}} \zeta(x) d S_{x} d \tau_{1}(y)
\end{aligned}
$$

For every $y \in D_{\bar{\beta}}$ the function $x \mapsto G_{\mu}^{\Omega}(x, y)$ is $L_{\mu}$-harmonic in $\Omega_{\bar{\beta}}$. By local elliptic estimates, for every $\xi \in \Sigma_{\beta}$,

$$
\sup _{x \in B_{\beta / 4}(\xi)}\left|\nabla_{x} G_{\mu}^{\Omega}(x, y)\right| \leq c \beta^{-1} \sup _{x \in B_{\beta / 2}(\xi)} G_{\mu}^{\Omega}(x, y)
$$

By Harnack's inequality,

$$
\sup _{x \in B_{\beta / 2}(\xi)} G_{\mu}^{\Omega}(x, y) \leq c^{\prime} \inf _{x \in B_{\beta / 2}(\xi)} G_{\mu}^{\Omega}(x, y)
$$

The constants $c, c^{\prime}$ are independent of $\beta \in(0, \bar{\beta} / 2), y \in D_{\bar{\beta}}$ and $\xi \in \Sigma_{\beta}$. Therefore we obtain,

$$
\begin{equation*}
\left|\nabla_{x} G_{\mu}^{\Omega}(x, y)\right| \leq C \beta^{-1} G_{\mu}^{\Omega}(x, y) \quad \forall x \in \Sigma_{\beta}, \forall y \in D_{\bar{\beta}}, \forall \beta \in(0, \bar{\beta} / 2) \tag{2.33}
\end{equation*}
$$

Hence,

$$
\left|I_{1,1}(\beta)\right| \leq C \beta^{-1} \int_{\Sigma_{\beta}} \mathbb{G}_{\mu}^{\Omega}\left[\tau_{1}\right]|\zeta| d S_{x}
$$

As $|\zeta(x)| \leq c \delta(x)^{\alpha_{+}}$it follows that,

$$
\left|I_{1,1}(\beta)\right| \leq C \frac{1}{\beta^{\alpha_{-}}} \int_{\Sigma_{\beta}} \mathbb{G}_{\mu}^{\Omega}\left[\tau_{1}\right] d S_{x}
$$

Therefore, by (2.29),

$$
\begin{equation*}
\left|I_{1,1}(\beta)\right| \leq C^{\prime} \varepsilon \quad \forall \beta \in(0, \bar{\beta} / 2) \tag{2.34}
\end{equation*}
$$

Finally (2.30) follows from (2.31), (2.32) and (2.34).
Theorem 2.18. Let $v \in \mathfrak{M}^{+}(\partial \Omega)$ and $\tau \in \mathfrak{M}_{\delta^{\alpha+}}(\Omega)$. Then:
(i) Problem (1.11) has a unique solution. The solution is given by

$$
\begin{equation*}
u=\mathbb{G}_{\mu}^{\Omega}[\tau]+\mathbb{K}_{\mu}^{\Omega}[v] \tag{2.35}
\end{equation*}
$$

(ii) There exists a positive constant $c=c(N, \mu, \Omega)$ such that

$$
\begin{equation*}
\|u\|_{L_{\delta^{-\alpha_{-}}}^{1}(\Omega)} \leq c\left(\|\tau\|_{\mathfrak{M}_{\delta^{\alpha}+}(\Omega)}+\|\nu\|_{\mathfrak{M}(\partial \Omega)}\right) \tag{2.36}
\end{equation*}
$$

(iii) $u$ is a solution of (1.11) if and only if $u \in L_{\delta^{-\alpha_{-}}}^{1}(\Omega)$ and

$$
\begin{equation*}
-\int_{\Omega} u L_{\mu} \zeta d x=\int_{\Omega} \zeta d \tau-\int_{\Omega} \mathbb{K}_{\mu}^{\Omega}[v] L_{\mu} \zeta d x \quad \forall \zeta \in X(\Omega) \tag{2.37}
\end{equation*}
$$

Proof. (i) Proposition 2.12 implies that (2.35) is a solution of (1.11).
If $u$ and $u^{\prime}$ are two solutions of (1.11) then $v:=\left(u-u^{\prime}\right)_{+}$is a nonnegative $L_{\mu}$-subharmonic function such that $\operatorname{tr}^{*}(v)=0$ and $v \leq 2 \mathbb{G}_{\mu}^{\Omega}[|\tau|]$ which is a positive $L_{\mu}$-superharmonic function. By Proposition 2.14, $v \equiv 0$ and hence $u \leq u^{\prime}$ in $\Omega$. Similarly $u^{\prime} \leq u$, so that $u=u^{\prime}$.
(ii) In view of (2.14) and (2.15), (2.36) is an immediate consequence of (2.35).
(iii) Let $u$ be the solution of (1.11). By (2.36), $u \in L_{\delta^{-\alpha_{-}}}^{1}(\Omega)$ and by Lemma 2.17 and (2.35), $u$ satisfies (2.37).

Conversely, suppose that $u \in L_{\delta^{-\alpha_{-}}}^{1}(\Omega)$ and satisfies (2.37). We show that $u$ is a solution of (1.11) or, equivalently, of (2.35).

By (2.37) with $\zeta \in C_{c}^{\infty}(\Omega), u$ is a solution (in the sense of distributions) of the equation in (1.11). It remains to show that $\operatorname{tr}^{*}(u)=v$. Put $U=u-\mathbb{G}_{\mu}^{\Omega}[\tau]-\mathbb{K}_{\mu}^{\Omega}[\nu]$ and note that, as $-L_{\mu} u=\tau, U$ is $L_{\mu}$-harmonic.

Let $z \in \Omega$ and let $\mathcal{G}(z)$ be a uniform regularization of $\left\{G_{\mu}^{D_{\beta}}: 0<\beta<\frac{1}{2} \tilde{\beta}(z)\right\}$ (see Section 2.4). Then, for every $\beta \in\left(0, \frac{1}{2} \tilde{\beta}(z)\right), \tilde{G}_{\mu}^{D_{\beta}} \in C_{0}^{2}\left(\bar{D}_{\beta}\right)$. Recall that $\tilde{G}_{\mu}^{D_{\beta}}(x)=G_{\mu}^{D_{\beta}}(z, x)$. Therefore, as $\frac{\partial G_{\mu}^{D_{\beta}}(z, x)}{\partial \mathbf{n}_{x}}=P_{\mu}^{D_{\beta}}(z, x), x \in \Sigma_{\beta}$, we obtain

$$
\begin{equation*}
-\int_{D_{\beta}} U(x) L_{\mu} \tilde{G}_{\mu}^{D_{\beta}}(x) d x=\int_{\Sigma_{\beta}} U(x) P_{\mu}^{D_{\beta}}(z, x) d S_{x}=U(z) \tag{2.38}
\end{equation*}
$$

The second equality is a consequence of the fact that $U$ is $L_{\mu}$-harmonic. But $L_{\mu} \tilde{G}_{\mu}^{D_{\beta}}(x) \rightarrow L_{\mu} \tilde{G}_{\mu}^{\Omega}(z, x)$ pointwise and the sequence $\left\{L_{\mu} \tilde{G}_{\mu}^{D_{\beta}}\right\}$ is bounded by a constant $M_{z}$. We observe that $U \in L^{1}(\Omega)$; in fact by assumption $u \in L_{\delta^{-\alpha_{-}}}^{1}(\Omega)$ and therefore, by Proposition 2.8, $U \in L_{\delta^{-\alpha_{-}}}^{1}(\Omega)$. Consequently, by (2.38),

$$
U(z)=-\int_{\Omega} U(x) L_{\mu} \tilde{G}_{\mu}^{\Omega}(z, x) d x
$$

Since $G_{\mu}^{\Omega}(z, \cdot) \in X(\Omega)$, by (2.37) the right hand side vanishes. Thus $U$ vanishes in $\Omega$, i.e., $u$ satisfies (2.35).
Corollary 2.19. Let $u$ be a positive $L_{\mu}$ superharmonic function. Then there exist $v \in \mathfrak{M}^{+}(\partial \Omega)$ and $\tau \in \mathfrak{M}_{\delta^{\alpha}}^{+}(\Omega)$ such that (1.12) holds.

Proof. By the Riesz decomposition theorem $u$ can be written in the form $u=u_{p}+u_{h}$ where $u_{p}$ is an $L_{\mu}$-potential and $u_{h}$ is a non-negative $L_{\mu}$-harmonic function. Therefore there exists $v \in \mathfrak{M}^{+}(\partial \Omega)$ such that $u_{h}=\mathbb{K}_{\mu}^{\Omega}[\nu]$. Since $u_{p}$ is an $L_{\mu}$-potential there exists a positive Radon measure $\tau$ such that $u_{p}=\mathbb{G}_{\mu}^{\Omega}[\tau]$ (see e.g. [1, Theorem 12]). This necessarily implies that $\tau \in \mathfrak{M}_{\delta^{\alpha+}}(\Omega)$.

Proposition 2.20. Let $w$ be a non-negative $L_{\mu}$-subharmonic function. If $w$ has a normalized boundary trace then it is dominated by an $L_{\mu}$-harmonic function.

Proof. There exist a positive Radon measure $\tau$ in $\Omega$ and a measure $v \in \mathfrak{M}^{+}(\partial \Omega)$ such that

$$
-L_{\mu} w=-\tau \quad \text { in } \Omega, \quad \operatorname{tr}^{*}(w)=v .
$$

Let $u_{\beta}$ be the solution of

$$
-L_{\mu} u=-\tau_{\beta} \quad \text { in } D_{\beta}, \quad u=\mathbb{K}_{\mu}^{\Omega}[\nu] \quad \text { on } \Sigma_{\beta}
$$

where $\tau_{\beta}:=\tau \chi_{D_{\beta}}$. Then,

$$
u_{\beta}+\mathbb{G}_{\mu}^{D_{\beta}}\left[\tau_{\beta}\right]=\mathbb{K}_{\mu}^{\Omega}[\nu] .
$$

Letting $\beta \rightarrow 0$ we obtain,

$$
\mathbb{G}_{\mu}^{\Omega}[\tau] \leq \mathbb{K}_{\mu}^{\Omega}[\nu] .
$$

Hence $\tau \in \mathfrak{M}_{\delta^{+}}^{+}(\Omega)$ and consequently

$$
\begin{equation*}
w+\mathbb{G}_{\mu}^{\Omega}[\tau]=\mathbb{K}_{\mu}^{\Omega}[\nu] . \tag{2.39}
\end{equation*}
$$

## 3. The nonlinear equation

In this section, we consider the nonlinear equation

$$
\begin{equation*}
-L_{\mu} u+u^{q}=0 \tag{3.1}
\end{equation*}
$$

in $\Omega$ with $0<\mu<C_{H}(\Omega)$ and $q>1$.
Proof of Theorem A. Since $u$ is a positive solution of (1.1), $u$ is $L_{\mu}$-subharmonic. Assuming (i), $u$ is dominated by an $L_{\mu}$-harmonic function. Therefore, by Proposition 2.14, (i) $\Longrightarrow$ (ii) and $u \in L_{\delta^{\alpha+}}^{q}(\Omega)$. On the other hand, by Proposition 2.20 (ii) $\Longrightarrow$ (i).

As mentioned above, (i) implies that $u \in L_{\delta^{\alpha}}^{q}(\Omega)$ and that there exists $v \in \mathfrak{M}_{\delta^{\alpha+}}^{+}(\partial \Omega)$ such that $\operatorname{tr}^{*}(u)=v$. Therefore, by Theorem 2.18, (1.14) is a consequence of (2.37). Thus (i) $\Longrightarrow$ (iii).

Finally, the implication (iii) $\Longrightarrow$ (i) is obvious.
It remains to prove the last assertion. If $u$ is a positive solution of (1.13) then, by (iii), $u \in L_{\delta^{\alpha+}}^{q}(\Omega)$ and (1.15) follows from Theorem 2.18.

Conversely, assume that $\delta^{\alpha}+u^{q}, u / \delta^{\alpha-} \in L^{1}(\Omega)$ and (1.15) holds. Then, by (1.15) with $\zeta \in C_{c}^{\infty}(\Omega), u$ is a solution of (1.1). Taking $\zeta_{f}=\mathbb{G}_{\mu}^{\Omega}[f]$ where $f \in C_{c}(\Omega)$ and $f \geq 0$ we obtain

$$
\int_{\Omega}\left(\mathbb{K}_{\mu}^{\Omega}[\nu]-u\right) f d x=\int_{\Omega} u^{q} \zeta_{f} d x<\infty .
$$

This implies $u \leq \mathbb{K}_{\mu}^{\Omega}[\nu]$, i.e., $u$ is $L_{\mu}$-moderate. Therefore by (i), $u$ is a solution of (1.13).

## Proof of Theorem B.

Uniqueness. Let $u_{1}$ and $u_{2}$ be two positive solutions of (1.13). Then $v:=\left(u_{1}-u_{2}\right)_{+}$is a subsolution of (1.1) and therefore an $L_{\mu}$-subharmonic function. Furthermore, by (iii) in Theorem A, $u_{1}, u_{2} \in L_{\delta^{\alpha+}}^{q}(\Omega)$ and $v \leq \mathbb{G}_{\mu}^{\Omega}\left[u_{1}^{q}+\right.$ $\left.u_{2}^{q}\right]=: \bar{v}$. Obviously $\bar{v}$ is $L_{\mu}$ superharmonic and $\operatorname{tr}^{*}(v)=0$. Therefore, by Proposition 2.14, $v=0$. Thus $u_{1} \leq u_{2}$ and similarly $u_{2} \leq u_{1}$.

Monotonicity. As before, $v:=\left(u_{1}-u_{2}\right)_{+}$is $L_{\mu}$-subharmonic and it is dominated by an $L_{\mu}$-superharmonic function. Since $\nu_{1} \leq \nu_{2}, \operatorname{tr}^{*}(v)=0$. Hence by Proposition 2.14, $v=0$.
A-priori estimate. Suppose that $u$ is a positive solution of (1.13). Then (1.15) with $\zeta=\mathbb{G}_{\mu}^{\Omega}[1]$ implies (1.16). (Recall that $\mathbb{G}_{\mu}^{\Omega}[1] \sim \delta^{\alpha+}$.)

For the proof of the next theorem we need
Lemma 3.1. Let $D \Subset \Omega$ be a $C^{2}$ domain and $q>1$. If $h$ is a positive function in $L^{1}(\partial D)$ then there exists a unique solution of the boundary value problem,

$$
\begin{align*}
-L_{\mu} u+u^{q} & =0 & & \text { in } D \\
u & =h & & \text { on } \partial D . \tag{3.2}
\end{align*}
$$

Proof. First assume that $h$ is bounded. Let $P_{\mu}^{D}$ denote the Poisson kernel of $-L_{\mu}$ in $D$ and put $u_{0}:=\mathbb{P}_{\mu}^{D}[h]$. Thus $u_{0}$ is bounded. We show that there exists a non-increasing sequence of positive functions $\left\{u_{n}\right\}_{1}^{\infty}$, dominated by $u_{0}$, such that $u_{n}$ is the solution of the boundary value problem,

$$
\begin{align*}
-\Delta v+v^{q} & =\frac{\mu}{\delta^{2}} u_{n-1} \quad \text { in } D \\
v & =h \quad \text { on } \partial D \quad n=1,2, \ldots \tag{3.3}
\end{align*}
$$

As usual $\delta$ denotes the distance to $\partial \Omega$, not to $\partial D$. For $n=1, u_{0}$ is a supersolution of the problem and, obviously $v=0$ is a subsolution. Consequently there exists a unique solution $u_{1}$. By induction, for $n>1$,

$$
-\Delta u_{n-1}+u_{n-1}^{q}=\frac{\mu}{\delta^{2}} u_{n-2} \geq \frac{\mu}{\delta^{2}} u_{n-1}
$$

Thus $v=u_{n-1}$ is a supersolution of (3.3) and it is bounded. It follows that there exists $0 \leq u_{n} \leq u_{n-1}$ such that

$$
-\Delta u_{n}+u_{n}^{q}=\frac{\mu}{\delta^{2}} u_{n-1} \text { in } D, \quad u_{n}=h \text { on } \partial D .
$$

As the sequence is monotone we conclude that $u=\lim u_{n}$ is a solution of (3.2).
If $h \in L^{1}(\partial D)$, we approximate it by a monotone increasing sequence of non-negative bounded functions $\left\{h_{k}\right\}$. If $v_{k}$ is the solution of (3.2) with $h$ replaced by $h_{k}$ then $\left\{v_{k}\right\}$ increases (by the comparison principle [4, Lemma 3.2]) and $v=\lim v_{k}$ is a solution of (3.2).

Uniqueness follows by the comparison principle.
Proof of Theorem C. Put $u_{0}:=\mathbb{K}_{\mu}^{\Omega}[\nu]$ and $h_{\beta}:=u_{0}\left\llcorner_{\beta}\right.$. Let $u_{\beta}$ be the solution of (3.2) with $h$ replaced by $h_{\beta}$, $\beta \in\left(0, \beta_{0}\right)$. Since $u_{0}$ is a supersolution of (1.1) it follows that $\left\{u_{\beta}\right\}$ decreases as $\beta \downarrow 0$. Therefore $u:=\lim _{\beta \rightarrow 0} u_{\beta}$ is a solution of (1.1).

We claim that $\operatorname{tr}^{*}(u)=v$. Indeed,

$$
\begin{equation*}
u_{\beta}+\mathbb{G}_{\mu}^{D_{\beta}}\left[u_{\beta}^{q}\right]=\mathbb{P}_{\mu}^{D_{\beta}}\left[h_{\beta}\right]=u_{0} . \tag{3.4}
\end{equation*}
$$

Furthermore, in $D_{\beta}, u_{\beta} \leq u_{0} \in L_{\delta^{\alpha}}^{q}(\Omega)$. Therefore

$$
\mathbb{G}_{\mu}^{D_{\beta}}\left[u_{\beta}^{q}\right] \rightarrow \mathbb{G}_{\mu}^{\Omega}\left[u^{q}\right] .
$$

Hence, by (3.4),

$$
u+\mathbb{G}_{\mu}^{\Omega}\left[u^{q}\right]=u_{0}=\mathbb{K}_{\mu}^{\Omega}[\nu] .
$$

By Proposition 2.12, $\operatorname{tr}^{*}(u)=v$.
By Theorem B the solution is unique.
Proof of Corollary C1. By the previous theorem, if $v=f$ where $f$ is a positive bounded function then (1.13) has a solution. If $0 \leq f \in L^{1}(\Omega)$ then it is the limit of an increasing sequence of such functions. Therefore, once again problem (1.13) with $v=f$ has a solution.

Proof of Theorem D. Put $v=\mathbb{K}_{\mu}^{\Omega}[\nu]-u$. By the comparison principle $v \geq 0$. Clearly $v$ is $L_{\mu}$-superharmonic in $\Omega$ and, by definition $\operatorname{tr}^{*}(v)=0$. By Proposition I(iv) $v$ is an $L_{\mu}$ potential. Consequently, by Theorem 2.6,

$$
\lim _{x \rightarrow y} \frac{v(x)}{\mathbb{K}_{\mu}^{\Omega}[\nu]}=0 \quad \text { non-tangentially, } \nu \text { a.e. on } \partial \Omega .
$$

This implies (1.17).
Proof of Theorem E. By Proposition 2.8, specifically inequality (2.15), $\mathbb{K}_{\mu}^{\Omega}[\nu] \in L_{\delta^{\alpha+}}^{q}(\Omega)$ for every $q \in\left(1, q_{\mu, c}\right)$ and $v \in \mathfrak{M}^{+}(\partial \Omega)$. Therefore the first assertion of the theorem is a consequence of Theorem C.

We turn to the proof of stability. Put $v_{n}=\mathbb{K}_{\mu}^{\Omega}\left[v_{n}\right]$. By Proposition $2.8,\left\{v_{n}\right\}$ is bounded in $L_{\delta^{\alpha+}}^{q}(\Omega)$ for every $q \in\left(1, q_{\mu, c}\right)$ and in $L_{\delta^{-\alpha_{-}}}^{p}(\Omega)$ for every $p \in\left(1, \frac{N-\alpha_{-}}{N-1-\alpha_{-}}\right)$. In addition $v_{n} \rightarrow v$ pointwise in $\Omega$. This implies that $\left\{v_{n}^{q} \delta^{\alpha_{+}}\right\}$and $\left\{v_{n} / \delta^{\alpha_{-}}\right\}$are uniformly integrable in $\Omega$. Since $u_{v_{n}} \leq v_{n}$ it follows that this conclusion applies also to $\left\{u_{v_{n}}\right\}$.

By the extension of the Keller-Osserman inequality due to [4], the sequence $\left\{u_{v_{n}}\right\}$ is uniformly bounded in every compact subset of $\Omega$. Therefore, by a standard argument, we can extract a subsequence, still denoted by $\left\{u_{v_{n}}\right\}$ that converges pointwise to a solution $u$ of (1.1). In view of the uniform convergence mentioned above we conclude that

$$
u_{v_{n}} \rightarrow u \quad \text { in } L_{\delta^{\alpha_{+}}}^{q}(\Omega) \text { and in } L_{\delta^{-\alpha_{-}}}^{1}(\Omega)
$$

By Theorem A,

$$
u_{\nu_{n}}+\mathbb{G}_{\mu}^{\Omega}\left[u_{\nu_{n}}^{q}\right]=\mathbb{K}_{\mu}^{\Omega}\left[\nu_{n}\right] .
$$

In view of the previous observations, passing to the limit as $n \rightarrow \infty$, we obtain,

$$
u+\mathbb{G}_{\mu}^{\Omega}\left[u^{q}\right]=\mathbb{K}_{\mu}^{\Omega}[\nu] .
$$

Again by Theorem A it follows that $u$ is the (unique) solution of (1.13). Because of the uniqueness we conclude that the entire sequence $\left\{u_{v_{n}}\right\}$ (not just a subsequence) converges to $u$ as stated in assertion II. of the theorem.

Finally we prove assertion III. By Theorem A

$$
\begin{equation*}
u_{k \delta_{y}}+\mathbb{G}_{\mu}^{\Omega}\left[u_{k \delta_{y}}^{q}\right]=k K_{\mu}^{\Omega}(\cdot, y) . \tag{3.5}
\end{equation*}
$$

Combining (2.7), (2.6) and the fact $u_{k \delta_{y}} \leq k K_{\mu}^{\Omega}(\cdot, y)$, we obtain

$$
\frac{\mathbb{G}_{\mu}^{\Omega}\left[u_{k \delta_{y}}^{q}\right](x)}{K_{\mu}^{\Omega}(x, y)} \leq k^{q} \frac{\mathbb{G}_{\mu}^{\Omega}\left[\left(K_{\mu}^{\Omega}(., y)^{q}\right](x)\right.}{K_{\mu}^{\Omega}(x, y)} \leq c k^{q}|x-y|^{N+\alpha_{+}-q\left(N-1-\alpha_{-}\right)} .
$$

Since $1<q<q_{\mu, c}$, it follows that

$$
\lim _{x \rightarrow y} \frac{\mathbb{G}_{\mu}^{\Omega}\left[u_{k \delta_{y}}^{q}\right](x)}{K_{\mu}^{\Omega}(x, y)}=0 .
$$

Therefore, by (3.5), we obtain (1.19).
Proof of Theorem F. Let $y \in \partial \Omega$. By negation, assume that there exists a positive solution $u$ of (1.13) with $v=k \delta_{y}$ for some $k>0$. By Theorem A, $u \leq k \mathbb{K}_{\mu}^{\Omega}(., y)$ and $u \in L_{\delta^{\alpha+}}^{q}(\Omega)$. Let $\gamma \in(0,1)$ and denote $C_{\gamma}(y)=\{x \in \Omega: \gamma|x-y| \leq$ $\delta(x)\}$. By Theorem D,

$$
\lim _{x \in C_{\gamma}(y), x \rightarrow y} \frac{u(x)}{K_{\mu}^{\Omega}(x, y)}=k .
$$

This implies that there exist positive numbers $r_{0}, c$ such that

$$
\begin{equation*}
u(x) \geq c K_{\mu}^{\Omega}(x, y) \quad \forall x \in C_{\gamma}(y) \cap B_{r_{0}}(y) . \tag{3.6}
\end{equation*}
$$

By (2.7),

$$
\begin{aligned}
& J_{\gamma}:=\int_{C_{\gamma}(y) \cap B_{r_{0}}(y)}\left(K_{\mu}^{\Omega}(x, y)\right)^{q} \delta(x)^{\alpha_{+}} d x \\
& \geq c^{\prime} \int_{C_{\gamma}(y) \cap B_{r_{0}}(y)} \delta(x)^{\alpha_{+}(q+1)}|x-y|^{\left(2 \alpha_{-}-N\right) q} d x \\
& \geq c^{\prime} \gamma^{\alpha_{+}(q+1)} \int_{C_{\gamma}(y) \cap B_{r_{0}}(y)}|x-y|^{\alpha_{+}-q\left(N-1-\alpha_{-}\right)} d x .
\end{aligned}
$$

Since $q \geq q_{\mu, c}$ the last integral is divergent. But (3.6) and the fact that $u \in L_{\delta^{\alpha+}}^{q}(\Omega)$ imply that $J_{\gamma}<\infty$. We reached a contradiction.

## Conflict of interest statement

No conflict of interest.

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