

EXISTENCE OF SOLUTIONS AND REGULARITY NEAR THE CHARACTERISTIC BOUNDARY FOR SUB-LAPLACIAN EQUATIONS ON CARNOT GROUPS

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We prove that the best constant in the Folland–Stein embedding theorem on Carnot groups is achieved. This implies the existence of a positive solution of the Yamabe-type equation on Carnot groups. The second goal of the paper is to show a certain regularity of the Green’s function and solutions of the Yamabe equation involving the sub-Laplacian near the characteristic boundary of a domain in the considered groups.

1. Introduction

In this paper we consider problems of existence and regularity of solutions to a nonlinear Dirichlet problem involving sub-Laplacians on Carnot groups. The main motivation comes from the classical Yamabe problem, the question of determining the best constant in the Sobolev embedding inequality, and their CR counterparts.

A Carnot group is a simply connected and connected Lie group G , whose Lie algebra \mathfrak{g} admits a stratification $\mathfrak{g} = \bigoplus_{j=1}^r V_j$ with

$$[V_1, V_j] = V_{j+1} \quad \text{for } 1 \leq j < r, \quad [V_1, V_r] = \{0\}.$$

Let $Q = \sum_{j=1}^r j \dim V_j$ be the homogeneous dimension.

Our starting point is this embedding result of Folland and Stein, [Folland 1975]: For any $p \in (1, Q)$ there exists $S_p = S_p(G) > 0$ such that, for $u \in C_0^\infty(G)$,

$$(1-1) \quad \left(\int_G |u|^{p^*} dH \right)^{1/p^*} \leq S_p \left(\int_G |Xu|^p dH \right)^{1/p}.$$

Here, the horizontal gradient $|Xu|$ is defined as $|Xu| = (\sum_{j=1}^m (X_j u)^2)^{1/2}$, where $X = \{X_1, \dots, X_m\}$ is a basis of V_1 , $p^* = pQ/(Q - p)$, and dH is a fixed Haar measure on G .

Unlike the Euclidean case [Talenti 1976; Aubin 1976b], the value of the best possible constant or the nonnegative functions for which it is achieved is unknown.

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In the particular case when $p = 2$ and $G = \mathbb{R}^n$, this problem is related to the Yamabe problem ([Trudinger 1968; Aubin 1976a; Schoen 1984], see also the survey article [Lee and Parker 1987]); in the case of the Heisenberg group, it is related to the CR Yamabe problem. The case of odd-dimensional spheres is equivalent to the problem of determining the best constant in the L^2 Folland–Stein inequality on the Heisenberg group, and the functions for which it is achieved. The solution of this problem is instrumental for solving the general case, which is complete after the works of Jerison and Lee [1984; 1987; 1988; 1989], and Gamara and Yacoub [Gamara 2001; Gamara and Yacoub 2001].

We note that, when $p = 2$, the Yamabe equation is the Euler–Lagrange equation satisfied by the extremals of the naturally associated variational problem. It is also of interest to study the Yamabe equation on sets different from the whole group, especially in connection with certain blow-up arguments that appear, for example, in questions of existence of solutions on domains with a nontrivial topology [Bahri and Coron 1988; Brezis 1986; Citti and Uguzzoni 2001].

We consider similar problems in a general Carnot group. The method of concentration compactness of P. L. Lions can be used, and we show in Section 3 that there exists a best constant in the Folland–Stein inequality. It is achieved on the space $\mathring{D}^{1,p}(G)$, which is the closure of $C_0^\infty(G)$ with respect to the norm

$$\|u\|_{\mathring{D}^{1,p}(G)} = \left(\int_G |Xu|^p dH \right)^{1/p}.$$

This method does not allow an explicit determination of the best constant or the functions for which it is achieved. The problem can be formulated as a variational problem. We consider the case $p = 2$ in detail. The Euler–Lagrange equation of the nonnegative extremals leads to the Yamabe-type equation

$$(1-2) \quad \mathcal{L}u = -u^{2^*-1},$$

where $\mathcal{L}u = \sum_{j=1}^m X_j^2 u$.

While it is relatively easy to see that, in any domain, weak solutions to this equation are bounded, further regularity is based on intricate subelliptic estimates. In Section 4 we show that any weak solution of the equation

$$(1-3) \quad \mathcal{L}_p u = \sum_{j=1}^n X_j (|Xu|^{p-2} X_j u) = -|u|^{p^*-2} u \quad \text{in } G.$$

is a bounded function. We shall present the proof in a somewhat more general case. There are several similar results in the Riemannian case. Yamabe [1960], proved the boundedness for solutions of the Yamabe problem on a manifold without boundary; see also [Trudinger 1968]. For the ordinary Laplacian

in a bounded domain, Brézis and Kato [1979] established a result similar to part (1) of Theorem 4.1; for part (2), see [Brézis and Nirenberg 1983]. The proofs rely on a suitable modification of the test-function and truncation ideas introduced in Serrin’s seminal paper [1964] and in [Moser 1961]. We note that, in all these results, u is assumed from the space $L^p(\Omega)$, since this is part of the definition of the considered Sobolev spaces. This is not true in the Sobolev spaces $\mathring{\mathcal{D}}^{1,p}(\Omega)$ that we consider, since we only include the L^p norm of the horizontal gradient in our definition. Therefore, the results here are not exactly the same, besides our working on a Carnot group. Subsequently, Serrin’s ideas were generalized to the subelliptic setting in [Capogna et al. 1993] and, in different forms, also in [Holopainen 1992; Holopainen and Rickman 1992; Lanconelli and Uguzzoni 1998; Xu 1990].

The rest of the paper concerns the regularity of solutions on bounded domains. In Section 5 we show that, under certain geometric conditions on the boundary of the C^∞ connected bounded open set $\Omega \subset \mathbf{G}$, one can prove boundedness of the horizontal gradient of the Green’s function. We shall also obtain estimates along other vector fields. Similar results hold near the boundary of the domain for weak nonnegative solutions of the following Yamabe-type equation

$$(1-4) \quad \mathcal{L}u = -u^{2^*-1}, \quad u \in \mathring{\mathcal{D}}^{1,2}(\Omega), \quad u \geq 0.$$

In order to prove such estimates, we impose some geometrical conditions on the boundary. In particular, we point to the “convexity condition” [Garofalo and Vasiliev 2000], assumed to hold globally whenever we are working with nonlinear equations, unless we are on the Heisenberg group. The reason for such an assumption is that, at present, there is no proof of the boundary Schauder estimates in Lipschitz spaces, or of $\Gamma^{2,\alpha}$ regularity near the points from the noncharacteristic portion of the boundary. Assuming the global validity of the convexity condition allows us to avoid the use of Schauder estimates or extra a priori regularity assumptions. However, by using the Lipschitz Schauder theory near the noncharacteristic boundary for domains on the Heisenberg group [Jerison 1981a], we present in Section 5E the argument in the local setting, with the convexity assumption holding only near the characteristic boundary.

An example in [Jerison 1981b] implies that the estimates we obtain at the characteristic boundary fail without some assumption on the boundary. However, the convexity condition is clearly not necessary for estimates of the type considered. Boundedness of the horizontal gradient has been established for more general domains, satisfying a uniform outer-ball condition [Capogna et al. 1998]. Verifying such a condition is in general hard. A result as in the $C^{1,1}$ -boundary case in \mathbb{R}^n is not true. Also, the horizontal gradient involves only differentiation along vectors from the first layer. We give estimates along vector fields involving differentiation along vectors from other layers; in particular, we include the radial vector field.

Lastly, in [Garofalo and Vassilev 2000] it was required that the domain be uniformly star-like near the characteristic boundary, which is unnecessary, as we will show.

2. Preliminaries and notation

We introduce the relevant definitions and state some results which will be needed in the sequel. Consider the Lie algebra $\mathfrak{g} = \bigoplus_{j=1}^r V_j$ of \mathbf{G} . We assume that there is a scalar product on \mathfrak{g} , with respect to which the V_j 's are mutually orthogonal.

Let $X = \{X_1, \dots, X_m\}$ be a basis of V_1 , and continue to denote by X the corresponding system of sections on \mathbf{G} . The *sub-Laplacian* associated with X is the second-order partial differential operator on \mathbf{G} given by

$$\mathcal{L} = - \sum_{j=1}^m X_j^* X_j = \sum_{j=1}^m X_j^2$$

(recall that $X_j^* = -X_j$ in a Carnot group; see [Folland 1975]). From the assumption on the Lie algebra, one immediately sees that the system X satisfies the well-known finite-rank condition, and therefore that the operator \mathcal{L} is hypoelliptic, thanks to Hörmander's theorem [1967]. However, it fails to be elliptic, and the loss of regularity is measured by the step r of the stratification of \mathfrak{g} . For a function u on \mathbf{G} , set

$$|Xu| = \left(\sum_{j=1}^m (X_j u)^2 \right)^{1/2}.$$

For $1 \leq p < Q$, define $\mathcal{D}^{1,p}(\mathbf{G})$ as the closure of $C_0^\infty(\mathbf{G})$ with respect to the norm

$$(2-1) \quad \|u\|_{\mathcal{D}^{1,p}(\mathbf{G})} = \left(\int_{\mathbf{G}} |Xu|^p dH \right)^{1/p}.$$

We define the *Sobolev exponent* relative to p as the number $p^* = pQ/(Q - p)$, where Q is the homogeneous dimension defined below. The relevance of such a number is emphasized by the Folland–Stein inequality (1-1).

In any Carnot group, the exponential mapping $\exp : \mathfrak{g} \rightarrow \mathbf{G}$ is an analytic diffeomorphism. We use it to define analytic maps $\xi_i : \mathbf{G} \rightarrow V_i$, $i = 1, \dots, r$, through the equation $g = \exp \xi(g)$, where $\xi(g) = \xi_1(g) + \dots + \xi_r(g)$. With $m = \dim(V_1)$, the coordinates of ξ 's projection in the basis X_1, \dots, X_m will be denoted by $x_1 = x_1(g), \dots, x_m = x_m(g)$, that is,

$$(2-2) \quad x_j(g) = \langle \xi(g), X_j \rangle, \quad j = 1, \dots, m.$$

We set $x = x(g) = (x_1, \dots, x_m) \in \mathbb{R}^m$. Later we will need to exploit the properties of the exponential coordinates in the second layer of the stratification of \mathfrak{g} . We

thus fix an orthonormal basis Y_1, \dots, Y_k of V_2 and, similarly to (2-2), define the exponential coordinates in the second layer V_2 of a point $g \in \mathbf{G}$ by setting

$$(2-3) \quad y_i(g) = \langle \xi(g), Y_i \rangle, \quad i = 1 \dots k,$$

and $y = (y_1, \dots, y_k) \in \mathbb{R}^k$.

Every Carnot group is naturally equipped with a family of nonisotropic dilations defined by

$$(2-4) \quad \delta_\lambda(g) = \exp \circ \Delta_\lambda \circ \exp^{-1}(g), \quad g \in \mathbf{G},$$

where $\exp : \mathfrak{g} \rightarrow \mathbf{G}$ is the exponential map and $\Delta_\lambda : \mathfrak{g} \rightarrow \mathfrak{g}$ is defined by

$$\Delta_\lambda(\xi_1 + \dots + \xi_r) = \lambda \xi_1 + \dots + \lambda^r \xi_r.$$

The *topological dimension* of \mathbf{G} is $N = \sum_{j=1}^r \dim V_j$, whereas the *homogeneous dimension* of \mathbf{G} , attached to the automorphisms $\{\delta_\lambda\}_{\lambda>0}$, is given by

$$Q = \sum_{j=1}^r j \dim V_j.$$

One has $dH(\delta_\lambda(g)) = \lambda^Q dH(g)$, so that, with respect to the group dilations, the number Q plays the role of a dimension. Let Z be the infinitesimal generator of the one-parameter group of nonisotropic dilations $\{\delta_\lambda\}_{\lambda>0}$. Such a vector field is characterized by the property that a function $u : \mathbf{G} \rightarrow \mathbb{R}$ is homogeneous of degree s with respect to $\{\delta_\lambda\}_{\lambda>0}$ — that is, $u(\delta_\lambda(x)) = \lambda^s u(x)$ for every $x \in \mathbf{G}$ — if and only if $Zu = su$.

The Euclidean distance to the origin $|\cdot|$ on \mathfrak{g} induces a homogeneous norm $|\cdot|_{\mathfrak{g}}$ on \mathfrak{g} , and (via the exponential map) a norm on the group \mathbf{G} , in the following way (see also [Folland 1975]): First, for $\xi \in \mathfrak{g}$ with $\xi = \xi_1 + \dots + \xi_r$, $\xi_i \in V_i$, set

$$(2-5) \quad |\xi|_{\mathfrak{g}} = \left(\sum_{i=1}^r |\xi_i|^{2r!/i} \right)^{2r!};$$

then define $|g|_{\mathbf{G}} = |\xi|_{\mathfrak{g}}$ if $g = \exp \xi$. Such a norm on \mathbf{G} can be used to define a pseudodistance on \mathbf{G} :

$$(2-6) \quad \rho(g, h) = |h^{-1}g|_{\mathbf{G}}.$$

The pseudodistance ρ is equivalent to the Carnot–Carathéodory distance $d(\cdot, \cdot)$ generated by the system X , that is, there exists a constant $C = C(\mathbf{G}) > 0$ such that

$$(2-7) \quad C\rho(g, h) \leq d(g, h) \leq C^{-1}\rho(g, h), \quad g, h \in \mathbf{G},$$

see [Nagel et al. 1985]. We will almost exclusively work with the distance d , except in a few situations where we will find more convenient to use (2-6). The

Carnot–Carathéodory balls are defined in the obvious way,

$$B_R(x) \equiv B(x, R) = \{y \in \mathbf{G} \mid d(x, y) < R\}.$$

By left-translation and dilation it is easy to see that the Haar measure of $B(x, R)$ is proportional to R^Q , where $Q = \sum_{i=1}^r i \dim V_i$ is the homogeneous dimension of \mathbf{G} . One has, for every $f, g, h \in \mathbf{G}$ and any $\lambda > 0$,

$$(2-8) \quad d(gf, gh) = d(f, h), \quad d(\delta_\lambda(g), \delta_\lambda(h)) = \lambda d(g, h).$$

We also recall the Baker–Campbell–Hausdorff formula; see, for example, [Hörmander 1967]:

$$(2-9) \quad \exp \xi \exp \eta = \exp\left(\xi + \eta + \frac{1}{2}[\xi, \eta] + \dots\right), \quad \xi, \eta \in \mathfrak{g},$$

where the dots indicate a linear combination of terms of order three and higher, which is finite due to the nilpotency of \mathbf{G} . By definition, the order of an element in V_j is j .

We next list some results which play an important role in this paper. To state these, we recall that, given a bounded open set $D \subset \mathbf{G}$ and a function $\varphi \in C(\partial D)$, the Dirichlet problem for D and a sub-Laplacian \mathcal{L} consists in finding a solution to $\mathcal{L}w = 0$ in D which takes value φ on the boundary.

Theorem 2.1 (Bony’s maximum principle [Bony 1969]). *Let $D \subset \mathbf{G}$ be a connected bounded open set, and $\varphi \in C(\partial D)$. There exists a unique \mathcal{L} -harmonic function H_φ^D that solves the Dirichlet problem in the sense of Perron–Wiener–Brelot. Moreover, H_φ^D satisfies*

$$\sup_D |H_\varphi^D| \leq \sup_{\partial D} |\varphi|.$$

Theorem 2.2 (Schauder-type interior estimates [Danielli and Garofalo 1998]). *Let $D \subset \mathbf{G}$ be an open set, and suppose that w is \mathcal{L} -harmonic in D . For every $g \in D$ and $r > 0$ with $\bar{B}(g, r) \subset D$, one has*

$$|X_{j_1} X_{j_2} \dots X_{j_s} w(g)| \leq \frac{C}{r^s} \max_{\bar{B}(g, r)} |w|,$$

for $s \in \mathbb{N}$, $j_1, \dots, j_s \in \{1, \dots, m\}$, and some constant $C = C(\mathbf{G}, s) > 0$.

To state the next result, we introduce a definition. Given an open set $D \subset \mathbf{G}$, we denote by $\mathcal{L}^{1,\infty}(D)$ the space of those distributions $u \in L^\infty(D)$ such that $Xu \in L^\infty(D)$, endowed with the natural norm.

Theorem 2.3 (L^∞ Poincaré inequality [Garofalo and Nhieu 1998]). *For a Carnot group \mathbf{G} , there exists $C = C(\mathbf{G}) > 0$ such that, if $u \in \mathcal{L}^{1,\infty}(B(g_0, 3R))$, then u can be modified on a set of measure zero in $\bar{B}(g_0, R)$ so as to satisfy*

$$|u(g) - u(h)| \leq C d(g, h) \|u\|_{\mathcal{L}^{1,\infty}(B(g_0, 3R))}$$

for every $g, h \in \bar{B}(g_0, R)$. If furthermore $u \in C^\infty(B(g_0, 3R))$, then only the L^∞ norm of Xu suffices in the right-hand side of the previous inequality.

We note explicitly that the theorem asserts that every function $u \in \mathcal{L}^{1,\infty}(B(g_0, 3R))$ has a representative which is Lipschitz-continuous in $B(g_0, R)$ with respect to the Carnot–Carathéodory distance d . The reverse implication also holds; see [Garofalo and Nhieu 1998].

Finally, we recall a characterization of the nonisotropic Lipschitz spaces on Carnot groups [Krantz 1982]:

Theorem 2.4. $f \in \Gamma_\alpha(\mathbf{G})$ if and only if $f(g_t) \in \Lambda_\alpha$ for every horizontal curve g_t .

Here $\Gamma_\alpha(\mathbf{G})$ and Λ_α denote, correspondingly, the nonisotropic Lipschitz space of Folland and Stein on \mathbb{R}^N , and the isotropic Lipschitz space on \mathbb{R} . A curve is called horizontal if $\frac{d}{dt}g_t \in \text{span}\{X_1, \dots, X_m\}$.

3. Variational problems

We apply the concentration-compactness principle of Lions [1985a; 1985b] in the homogeneous setting of a Carnot group \mathbf{G} , to prove that, for any $1 < p < Q$, the best constant in the Folland–Stein embedding (1-1) is achieved. A simple argument shows that, without any loss of generality, we can consider only nonnegative functions; we shall do that throughout this section. Consequently, we show that, for any such p , the quasilinear equation with critical exponent,

$$(3-1) \quad \mathcal{L}_p u = \sum_{j=1}^n X_j(|Xu|^{p-2} X_j u) = -u^{p^*-1} \quad \text{in } \mathbf{G},$$

possesses a weak nonnegative solution that is also, up to a constant, an extremal for the following variational problem

$$(3-2) \quad I \stackrel{\text{def}}{=} \inf \left\{ \int_{\mathbf{G}} |Xu|^p \mid u \in C_0^\infty(\mathbf{G}), \int_{\mathbf{G}} |u|^{p^*} = 1 \right\}.$$

We used that $\mathcal{L}_p(cu) = c^{p-1} \mathcal{L}_p u$ to reduce the equation given by the Euler–Lagrange multiplier to (3-1). We shall consider a similar problem when restricting the test functions to those having with a certain symmetry. We prove that the corresponding infimum is achieved again. The purpose of this section is to record such basic results.

It is still an open question to find the norm of the Folland–Stein embedding. An interesting and more accessible problem is obtained by requiring that the group be of Heisenberg type. The precise value of the norm of the embedding when considering only functions with symmetries and $p = 2$ was found in [Garofalo and Vassilev 2001]. The proof required the existence result of this section. In

the particular case of the Heisenberg group, this problem was solved without the symmetry restriction by [Jerison and Lee \[1988\]](#).

3A. The best constant in the Folland–Stein inequality. A minimizing sequence $\{u_n\} \in C_0^\infty(\mathbf{G})$ of the variational problem (3-2) is characterized by the properties

$$(3-3) \quad \int_{\mathbf{G}} |u_n|^{p^*} = 1 \quad \text{and} \quad \int_{\mathbf{G}} |Xu_n|^p \xrightarrow[n \rightarrow \infty]{} I.$$

Two crucial aspects of equation (3-1) and of the above variational problem are their invariance with respect to group translations and dilations. The former is obvious, since the vector fields X_j are left-invariant. The latter must be suitably interpreted, and follows from the observation that

$$(3-4) \quad \mathcal{L}_p(u \circ \delta_\lambda) = \lambda^p \delta_\lambda \circ \mathcal{L}_p u.$$

If we thus define, for a solution u of (3-1) and for $\lambda > 0$, the rescaled function $u_\lambda = \lambda^\alpha u \circ \delta_\lambda$, then it is clear that u_λ satisfies (3-1) if and only if $\alpha = Q/p^* = (Q - p)/p$. These considerations lead to the introduction, for $u \in C_0^\infty(\mathbf{G})$, of two new functions:

$$(3-5) \quad \tau_h u \stackrel{\text{def}}{=} u \circ \tau_h, \quad h \in \mathbf{G},$$

where $\tau_h : \mathbf{G} \rightarrow \mathbf{G}$ is the left-translation operator $\tau_h(g) = hg$; and

$$(3-6) \quad u_\lambda \stackrel{\text{def}}{=} \lambda^{Q/p^*} u \circ \delta_\lambda, \quad \lambda > 0.$$

It is easy to see that the norms in the Folland–Stein inequality and the functionals in the variational problem (3-2) are invariant under the transformations (3-5) and (3-6). Only the second part requires a small computation, since dH is bi-invariant under translations.

$$\|\delta_\lambda u\|_{L^{p^*}(\mathbf{G})}^{p^*} = \int_{\mathbf{G}} |u(\delta_\lambda g)|^{p^*} dH(g) = \int_{\mathbf{G}} |u(g)|^{p^*} \lambda^{-Q} dH(g) = \lambda^{-Q} \|u\|_{L^{p^*}(\mathbf{G})}^{p^*}.$$

This shows that $\|u_\lambda\|_{L^{p^*}(\mathbf{G})} = \|u\|_{L^{p^*}(\mathbf{G})}$. Similarly,

$$\|\delta_\lambda u\|_{\mathfrak{D}^{1,p}(\mathbf{G})}^p = \|\lambda \delta_\lambda Xu\|_{L^{p^*}(\mathbf{G})}^p = \lambda^p \|\delta_\lambda Xu\|_{L^{p^*}(\mathbf{G})}^p = \lambda^{p-Q} \|u\|_{\mathfrak{D}^{1,p}(\mathbf{G})}^p.$$

Taking into account that $p^* = pQ/(Q - p)$, we obtain $\|Xu_\lambda\|_{\mathfrak{D}^{1,p}(\mathbf{G})} = \|Xu\|_{\mathfrak{D}^{1,p}(\mathbf{G})}$.

The main result about the existence of global minimizers is this:

Theorem 3.1. *Let \mathbf{G} be a Carnot group, and consider the minimization problem (3-2). Every minimizing sequence $\{u_n\}$ of (3-2) is relatively compact in $\mathfrak{D}^{1,p}(\mathbf{G})$, possibly after translating and dilating each of its elements using (3-5) and (3-6). In particular, there exists a minimum of (3-2), and the equation*

$$(3-7) \quad \mathcal{L}_p u = -u^{p^*-1}$$

admits a nontrivial nonnegative solution $u \in \mathring{D}^{1,p}(\mathbf{G})$.

The proof of [Theorem 3.1](#) is based on an adaptation of the concentration-compactness principle. In such an adaptation, the Euclidean space \mathbb{R}^n is replaced by a Carnot group \mathbf{G} , with its homogeneous structure and Carnot–Carathéodory distance. We mention that the implementation of Lions’ program relies, among others, on the Rellich–Kondrachov compact embedding. In the subelliptic setting, the proof of this result requires a substantial amount of work. A general version of it was proved in [[Garofalo and Nhieu 1996](#)]; it states that, if Ω denotes a bounded X -PS domain (Poincaré–Sobolev domain) in a Carnot–Carathéodory space, then the embedding

$$\mathcal{L}^{1,p}(\Omega) \subset L^q(\Omega)$$

is compact provided that $1 \leq q < p^* = pQ/(Q - p)$. Here, $\mathcal{L}^{1,p}(\Omega)$ indicates the Sobolev space of those functions $f \in L^p(\Omega)$ such that $Xf \in L^p(\Omega)$, endowed with the natural norm. Carnot groups are the basic models of Carnot–Carathéodory spaces. We shall apply such a result to an increasing sequence of bounded domains $\Omega_k \subset \Omega_{k+1} \subset \mathbf{G}$ such that $\Omega_k \nearrow \mathbf{G}$. As Ω_k , we can take the Carnot–Carathéodory ball B_k centered at the identity $e \in \mathbf{G}$ with radius k , since it was proved in [[Franchi et al. 1994](#); [Garofalo and Nhieu 1996](#)] that such sets are X -PS domains.

An important tool is the concentration function of a measure:

Definition 3.2. For a nonnegative measure $d\nu$ on \mathbf{G} , define the concentration function Q on $[0, \infty)$ by

$$(3-8) \quad Q(r) \stackrel{\text{def}}{=} \sup_{g \in \mathbf{G}} \left(\int_{B_r(g)} d\nu \right).$$

For a function f on \mathbf{G} , we will call the *concentration function of f* the concentration function of the measure $|f|^{p^*} dH$.

Similarly to Lions’ works, the crucial ingredients in the solution of the variational problem are the next lemmata. For the proofs, we refer to [[Lions 1984a](#); [1984b](#)]. As already mentioned, we have a suitable version of the Rellich–Kondrachov compact embedding to replace the usual embedding used in the proof of the next lemma:

Lemma 3.3. *Suppose ν_n is a sequence of probability measures on \mathbf{G} . There exists a subsequence, which we denote by $d\nu_n$, such that exactly one of the following three conditions holds:*

(1) *Compactness: There is a sequence $(g_n) \in \mathbf{G}$ such that, for every $\varepsilon > 0$, there exists $R > 0$ so that, for every n ,*

$$\int_{B(g_n, R)} d\nu_n \geq 1 - \varepsilon.$$

(2) *Vanishing:* For all $R > 0$, we have

$$\lim_{n \rightarrow \infty} \left(\sup_{g \in \mathbf{G}} \int_{B(g, R)} dv_n \right) = 0.$$

(3) *Dichotomy:* There exists a λ with $0 < \lambda < 1$, such that, for every $\varepsilon > 0$, there exist $R > 0$ and a sequence (g_n) with the property: Given $R' > R$, there exist nonnegative measures v_n^1 and v_n^2 for which

$$\begin{aligned} 0 \leq v_n^1 + v_n^2 \leq v_n, \\ \text{supp } v_n^1 \subset B(g_n, R) \quad \text{and} \quad \text{supp } v_n^2 \subset \mathbf{G} \setminus B(g_n, R'), \\ \left| \lambda - \int v_n^1 \right| + \left| (1 - \lambda) - \int v_n^2 \right| \leq \varepsilon. \end{aligned}$$

Remark 3.4. Using diagonal subsequences, we can also achieve

$$\begin{aligned} \text{supp } v_n^1 \subset B_{R_n}(g_n) \quad \text{and} \quad \text{supp } v_n^2 \subset \mathbf{G} \setminus B_{2R_n}(g_n), \\ \lim_{n \rightarrow \infty} \left| \lambda - \int_{\mathbf{G}} v_n^1 \right| + \left| (1 - \lambda) - \int_{\mathbf{G}} v_n^2 \right| = 0. \end{aligned}$$

Lemma 3.5. Suppose $u_n \rightharpoonup u$ weak-* in $\mathring{\mathcal{D}}^{1,p}(\mathbf{G})$, $\mu_n = |Xu_n|^p dH \rightharpoonup \mu$, and $v_n = |u_n|^{p^*} dH \rightharpoonup v$ weak-* in measure, where μ and v are bounded nonnegative measures on \mathbf{G} . There exist points $g_j \in \mathbf{G}$ and real numbers $d_j \geq 0$ and $e_j \geq 0$, at most countably many nonzero, such that

$$\begin{aligned} v &= |u|^{p^*} + \sum_j d_j \delta_{g_j} \\ \mu &\geq |Xu|^p dH + \sum_j e_j \delta_{g_j} \\ Id_j^{p/p^*} &\leq e_j, \end{aligned}$$

where I is the constant from (3-2). In particular,

$$\sum d_j^{p/p^*} < \infty.$$

Proof of Theorem 3.1. Since $p > 1$, we can assume that $u_n \rightharpoonup u$ weak-* in $\mathring{\mathcal{D}}^{1,p}(\mathbf{G})$ for some $u \in \mathring{\mathcal{D}}^{1,p}(\mathbf{G})$, by regarding $\mathring{\mathcal{D}}^{1,p}(\mathbf{G})$ as the dual of $\mathcal{D}^{1,p'}(\mathbf{G})$, where p' is the Hölder conjugate of p . From the Folland–Stein embedding theorem, this is also true for the weak-* convergence in $L^{p^*}(\mathbf{G})$. We can also assume the minimizing sequence is a.e. pointwise convergent on \mathbf{G} . This follows easily from Rellich’s theorem, applied successively to an exhaustion of \mathbf{G} by an increasing sequence of Carnot–Carathéodory balls.

The sequential lower semicontinuity of the norms shows that

$$\begin{aligned} \|u\|_{L^{p^*}(G)} &\leq \underline{\lim} \|u_n\|_{L^{p^*}(G)} = 1, \\ \|u\|_{\dot{G}^{1,p}(G)} &\leq \underline{\lim} \|u_n\|_{\dot{G}^{1,p}(G)} = I^{1/p}. \end{aligned}$$

Thus, it is enough to prove that $\int_G |u|^{p^*} = 1$, since, by the previous relations and the Folland–Stein inequality, and noticing that $I = 1/S_p^p$, we have

$$(3-9) \quad I \geq \int_G |Xu|^p \geq I \left(\int_G |u|^{p^*} \right)^{p/p^*} = I \quad \text{when} \quad \int_G |u|^{p^*} = 1,$$

which would give that u is a minimizer. In other words, we reduce the proof to showing that $u_n \rightarrow u$ in $L^{p^*}(G)$, as weak-* and norm convergence imply strong convergence.

Because of translation and dilation invariance, all the mentioned properties hold if we replace (u_n) with any translated and rescaled sequence, which we shall denote by (v_n) . We will consider the following measures,

$$dv_n \stackrel{\text{def}}{=} |v_n|^{p^*} dH, \quad d\mu_n \stackrel{\text{def}}{=} |Xv_n|^p dH.$$

where v_n is a suitable translation and dilation of u_n that is to be defined in a moment. Note that dv_n is also a sequence of probability measures. From the weak-* compactness of the unit ball, without loss of generality we can assume that $dv_n \rightharpoonup dv$ and $d\mu_n \rightharpoonup d\mu$ in the weak-* topology of all bounded nonnegative measures.

The desired convergence, that is, the fact that $\int_G |v|^{p^*} = 1$, will be obtained by applying the concentration-compactness principle, exactly as in [Lions 1985a] (see also [Struwe 1990]). We shall see that v is a probability measure, as well as that $dv = \int_G |v|^{p^*}$. Here v is the limit of the sequence (v_n) taken in various spaces, as we did for the sequence (u_n) .

Let $\hat{Q}_n(r)$ be the concentration function of u_n , that is,

$$(3-10) \quad \hat{Q}_n(r) \stackrel{\text{def}}{=} \sup_{h \in G} \left(\int_{B_r(h)} |u_n|^{p^*} dH \right).$$

Clearly, $\hat{Q}_n(0) = 0$, $\lim_{r \rightarrow \infty} \hat{Q}_n(r) = 1$, and \hat{Q}_n is a continuous nondecreasing function. Therefore, for every n we can find an $r_n > 0$ such that

$$(3-11) \quad \hat{Q}_n\left(\frac{1}{r_n}\right) = \frac{1}{2}.$$

Since the integral in (3-10) is absolutely continuous, it defines a continuous function of h , which, as $u_n \in L^{p^*}(G)$, tends to zero when $d(h, e) \rightarrow \infty$. Consequently,

the supremum is achieved, that is, for every n there exist a $h_n \in \mathbf{G}$ such that

$$(3-12) \quad \hat{Q}_n\left(\frac{1}{r_n}\right) = \int_{B_{1/r_n}(h_n)} |u_n(g)| dH(g).$$

The concentration functions Q_n of the dilated and translated sequence

$$(3-13) \quad v_n \stackrel{\text{def}}{=} r_n^{-Q/p^*} \delta_{r_n^{-1}} \tau_{g_n} u$$

satisfy

$$(3-14) \quad Q_n(1) = \int_{B_1(e)} dv_n \quad \text{and} \quad Q_n(1) = \frac{1}{2},$$

where we have set $g_n = \delta_{d_n} h_n^{-1}$.

The homogeneity properties of the metric are essential for proving (3-14). From the definition of v_n and (2-8),

$$(3-15) \quad \begin{aligned} \int_{B_r(g)} |v_n(h)|^{p^*} dH(h) &= \int_{\{d(g,h) < r\}} r_n^{-Q} |u_n(\delta_{r_n^{-1}} \tau_{g_n} h)|^{p^*} dH(h) \\ &= \int_{\{d(g, \tau_{g_n}^{-1} \delta_{r_n} z) < r\}} r_n^{-Q} |u_n(z)|^{p^*} r_n^Q dH(z) \\ &= \int_{\{d(\tau_{g_n} g, \delta_{r_n} z) < r\}} |u_n(z)|^{p^*} dH(z) \\ &= \int_{\{r_n d(\delta_{r_n^{-1}} \tau_{g_n} g, z) < r\}} |u_n(z)|^{p^*} dH(z) \\ &= \int_{\{d(h,z) < r/r_n\}} |u_n(z)|^{p^*} dH(z), \end{aligned}$$

where $h = \delta_{r_n^{-1}} \tau_{g_n} g$. By taking the supremum over $g \in \mathbf{G}$, we obtain

$$(3-16) \quad Q_n(1) = \hat{Q}_n\left(\frac{1}{r_n}\right) = \frac{1}{2}.$$

When $g = e$, we have

$$(3-17) \quad \int_{B_1(e)} |v_n(h)|^{p^*} dH(h) = \int_{B_{1/r_n}(h_n)} |u_n(z)|^{p^*} dH(z),$$

which shows that (3-14) holds.

At this point, we are ready to apply the lemmata and finish the proof. Notice that the vanishing case in Lemma 3.5 is ruled out by the normalization $Q_n(1) = 1/2$. We embed our variational problem in the family

$$(3-18) \quad I_\lambda \stackrel{\text{def}}{=} \inf \left\{ \int_{\mathbf{G}} |Xu|^p \mid u \in C_0^\infty(\mathbf{G}), \int_{\mathbf{G}} |u|^{p^*} = \lambda \right\}.$$

Since $I_\lambda = \lambda^{p/p^*} I$, we see that I_λ is strictly subadditive, that is,

$$(3-19) \quad I_1 < I_\alpha + I_{1-\alpha} \quad \text{for every } \alpha \in (0, 1).$$

Assume first that the compactness case in [Lemma 3.3](#) holds when applied to the sequence dv_n . Let g_n be points as in the lemma and, for $\varepsilon > 0$, choose $R = R(\varepsilon)$, such that

$$(3-20) \quad \int_{B_R(g_n)} dv_n \geq 1 - \varepsilon \quad \text{for every } n.$$

If $\varepsilon < 1/2$, then $\int_{B_R(g_n)} dv_n > 1/2$. Since by construction $\int_{B_1(e)} dv_n = 1/2$ while $\int_G dv_n = 1$, we see that $B_1(e)$ and $B_R(g_n)$ have a nonempty intersection, and thus there exist a ball $B_{2R+1}(e) \supset B_R(g_n)$. This implies that

$$(3-21) \quad \int_{B_{2R+1}(e)} dv_n \geq \int_{B_R(g_n)} dv_n \geq 1 - \varepsilon \quad \text{for every } n,$$

and therefore the condition (1) of [Lemma 3.3](#) holds with $g_n \equiv e$ for every n . By taking $\varepsilon \rightarrow 0$, we proved that

$$(3-22) \quad \int_G dv = 1.$$

If we look at the sequence $(d\mu_n)$ now, we have $d\mu_n \rightarrow d\mu$ and $\int_G d\mu_n \rightarrow I$, and thus $\int_G d\mu \leq I$. On the other hand, [Lemma 3.5](#) gives

$$\begin{aligned} dv_n \rightarrow dv &= |v|^{p^*} + \sum_j d_j \delta_{g_j}, \\ d\mu_n \rightarrow d\mu &\geq |Xv|^p dH + \sum_j e_j \delta_{g_j}, \end{aligned}$$

for certain $d_j, e_j \geq 0$ satisfying

$$(3-23) \quad I d_j^{p/p^*} \leq e_j.$$

We shall prove that all d_j 's are zero, and thus $\int_G |v|^{p^*} dH = 1$. Let

$$\alpha \stackrel{\text{def}}{=} \int_G |v|^{p^*} dH < 1.$$

Since $\int_G dv = 1$, we have $\sum d_j = 1 - \alpha$. From $\int_G d\mu \leq I$ we have

$$\int_G |Xv|^p dH \leq I - \sum e_j.$$

Now [\(3-23\)](#) gives

$$I = I_1 \geq \int_G |Xv|^p dH + \sum e_j \geq I_\alpha + \sum I d_j^{p/p^*} \geq I_\alpha + \sum I d_j.$$

From the strict subadditivity (3-19) of I_λ , we conclude that exactly one of the numbers α and d_j is nonzero. We claim that $\alpha = 1$ (and thus that all d_j 's are zero). Suppose that there is a d_j with $d_j = 1$ and $d\nu = \delta_{g_j}$. From the normalization, $Q_n(1) = 1/2$, and hence

$$(3-24) \quad \frac{1}{2} \geq \int_{B_1(g_j)} |v_n|^{p^*} dH \rightarrow \int_{B_1(g_j)} d\nu = 1,$$

which is a contradiction. Thus, we proved that $\|v\|_{L^{p^*}(G)} = \alpha = 1$ and $v_n \rightarrow v$ in $L^{p^*}(G)$, which shows that v is a solution of the variational problem; see (3-9).

Suppose we have dichotomy. We have to show that this leads to a contradiction. A simple process of taking a diagonal subsequence (see the remark after Lemma 3.3) shows that we can find a sequence $R_n > 0$ such that

$$(3-25) \quad \text{supp } v_n^1 \subset B_{R_n}(g_n), \quad \text{supp } v_n^2 \subset G \setminus B_{2R_n}(g_n),$$

$$(3-26) \quad \lim_{n \rightarrow \infty} \left| \lambda - \int_G v_n^1 \right| + \left| (1 - \lambda) - \int_G v_n^2 \right| = 0.$$

We fix a number ε such that

$$(3-27) \quad 0 < \varepsilon < \lambda^{p/p^*} + (1 - \lambda)^{p/p^*} - 1.$$

Such a choice of ε is possible, as for $0 < \lambda < 1$ and $p/p^* < 1$ we have $\lambda^{p/p^*} + (1 - \lambda)^{p/p^*} - 1 > 0$.

Let φ be a cut-off function $0 \leq \varphi \in C_0^\infty(B_2(e))$ with $\varphi \equiv 1$ on $B_1(e)$, and set $\varphi_n = \delta_{R_n^{-1}} \tau_{g_n} \varphi$. We have

$$(3-28) \quad \int_G |Xv_n|^p dH = \int_G |X(\varphi_n v_n)|^p dH + \int_G |X((1 - \varphi_n)v_n)|^p dH + \varepsilon_n.$$

Note that the remainder term ε_n is expressed by an integral over an annulus

$$(3-29) \quad A_n = B_{2R_n}(g_n) \setminus B_{R_n}(g_n).$$

Furthermore, we claim that

$$(3-30) \quad \varepsilon_n \geq o(1) - \varepsilon \int_G |Xv_n|^p dH, \quad \text{where } o(1) \rightarrow 0 \text{ as } n \rightarrow \infty.$$

Indeed, using the inequality

$$(|a| + |b|)^p \leq (1 + \varepsilon)|a|^p + C_{\varepsilon,p}|b|^p,$$

which holds for any $0 < \varepsilon < 1$, $p \geq 1$, and a suitable constant $C_{\varepsilon, p}$ depending on ε and p , we have

$$\begin{aligned} \varepsilon_n &= \int_{A_n} |Xv_n|^p dH - \int_{A_n} |X(\varphi_n v_n)|^p dH - \int_{A_n} |X((1 - \varphi_n)v_n)|^p dH \\ &\geq \int_{A_n} |Xv_n|^p (1 - \varphi_n^p - (1 - \varphi_n)^p) dH - 2C_{\varepsilon, p} \int_{A_n} |v_n|^p |X\varphi_n|^p dH \\ &\quad - \varepsilon \int_{A_n} |Xv_n|^p dH. \end{aligned}$$

Since $p > 1$ and $0 \leq \varphi \leq 1$, it follows that $1 \geq \varphi_n^p + (1 - \varphi_n)^p$, and thus

$$(3-31) \quad \varepsilon_n \geq -C \int_{A_n} |v_n|^p |X\varphi_n|^p dH - \varepsilon \int_G |Xv_n|^p dH.$$

First we use $|X\varphi_n| \leq C/R_n$, and then we apply Hölder's inequality on A_n ,

$$R_n^{-1} \|v_n\|_{L^p(A_n)} \leq R_n^{-1} |A_n|^{1/p - 1/p^*} \|v_n\|_{L^{p^*}(A_n)}.$$

Since $1/p - 1/p^* = 1/Q$ and, from the paragraph above (2-8),

$$(3-32) \quad |A_n| \sim R_n^Q,$$

we obtain

$$(3-33) \quad R_n^{-1} \|v_n\|_{L^p(A_n)} \leq C \|v_n\|_{L^{p^*}(A_n)}.$$

The last term in the above inequality can be estimated as follows:

$$\begin{aligned} (3-34) \quad \|v_n\|_{L^{p^*}(A_n)}^{p^*} &= \int_{A_n} dv_n = \int_G dv_n - \int_{G \setminus A_n} dv_n \\ &\leq \int_G dv_n - \int_{G \setminus A_n} dv_n^1 - \int_{G \setminus A_n} dv_n^2 \\ &= \int_G dv_n - \int_G dv_n^1 - \int_G dv_n^2. \end{aligned}$$

Hence, the claim (3-30) follows from

$$(3-35) \quad R_n^{-1} \|v_n\|_{L^p(A_n)} \leq C \left(\int_G dv_n - \int_G dv_n^1 - \int_G dv_n^2 \right)^{1/p^*} \rightarrow 0 \quad \text{as } n \rightarrow \infty.$$

Continuing with the proof that the dichotomy case does not occur, we use the definition of I , see (2-1) and (3-2), together with the above inequalities, to get

$$\begin{aligned} \|v_n\|_{\mathring{\mathcal{D}}^{1,p}(\mathbf{G})}^p &= \|\varphi_m v_n\|_{\mathring{\mathcal{D}}^{1,p}(\mathbf{G})}^p + \|(1 - \varphi_n)v_n\|_{\mathring{\mathcal{D}}^{1,p}(\mathbf{G})}^p + \varepsilon_n \\ &\geq I \left(\|\varphi_m v_n\|_{L^{p^*}(\mathbf{G})}^p + \|(1 - \varphi_n)v_n\|_{L^{p^*}(\mathbf{G})}^p \right) + \varepsilon_n \\ &\geq I \left(\left(\int_{B_{R_n}(g_n)} dv_n \right)^{p/p^*} + \left(\int_{\mathbf{G} \setminus B_{R_n}(g_n)} dv_n \right)^{p/p^*} \right) + \varepsilon_n \\ &\geq I \left(\left(\int_{\mathbf{G}} dv_n^1 \right)^{p/p^*} + \left(\int_{\mathbf{G}} dv_n^2 \right)^{p/p^*} \right) + \varepsilon_n. \end{aligned}$$

Letting $n \rightarrow \infty$, we obtain

$$I = \lim_{n \rightarrow \infty} \|v_n\|_{\mathring{\mathcal{D}}^{1,p}(\mathbf{G})}^p \geq I(\lambda^{p/p^*} + (1-\lambda)^{p/p^*}) - \varepsilon I,$$

which is a contradiction with the choice of ε in (3-27), and hence the dichotomy case of Lemma 3.3 cannot occur. The proof of the theorem is finished. \square

3B. The best constant in the presence of symmetries. We consider here the same problem as before, but we restrict the class of test functions.

Definition 3.6. Let \mathbf{G} be a Carnot group with Lie algebra $\mathfrak{g} = V_1 \oplus V_2 \cdots \oplus V_n$. We say that a function $U : \mathbf{G} \rightarrow \mathbb{R}$ has *partial symmetry* with respect to g_0 if there exist an element $g_0 \in \mathbf{G}$ such that for every $g = \exp(\xi_1 + \xi_2 + \cdots + \xi_n) \in \mathbf{G}$ one has

$$U(g_0 g) = u(|\xi_1(g)|, \dots, |\xi_{n-1}(g)|, \xi_n(g)),$$

for some function $u : [0, \infty) \times \cdots \times [0, \infty) \times V_n \rightarrow \mathbb{R}$.

A function U is said to have *cylindrical symmetry* if there exist $g_0 \in \mathbf{G}$ and $\varphi : [0, \infty) \times \cdots \times [0, \infty) \rightarrow \mathbb{R}$ for which

$$U(g_0 g) = \varphi(|\xi_1(g)|, |\xi_2(g)|, \dots, |\xi_n(g)|),$$

for every $g \in \mathbf{G}$.

We also define the spaces $\mathring{\mathcal{D}}_{\text{ps}}^{1,p}(\mathbf{G})$ and $\mathring{\mathcal{D}}_{\text{cyl}}^{1,p}(\mathbf{G})$ by

$$(3-36) \quad \mathring{\mathcal{D}}_{\text{ps}}^{1,p}(\mathbf{G}) \stackrel{\text{def}}{=} \left\{ u \in \mathring{\mathcal{D}}^{1,p}(\Omega) \mid u(g) = u(|\xi_1(g)|, \dots, |\xi_{n-1}(g)|, \xi_n(g)) \right\},$$

$$(3-37) \quad \mathring{\mathcal{D}}_{\text{cyl}}^{1,p}(\mathbf{G}) \stackrel{\text{def}}{=} \left\{ u \in \mathring{\mathcal{D}}^{1,p}(\Omega) \mid u(g) = u(|\xi_1(g)|, |\xi_2(g)|, \dots, |\xi_n(g)|) \right\}.$$

The effect of the symmetries (see also [Lions 1985b]) is manifested in the fact that, if the limit measure given by Lemma 3.5 concentrates at a point, then it must concentrate on the whole orbit of the group of symmetries. Therefore, in the cylindrical case there could be no points of concentration except at the origin,

while in the partially symmetric case the points of concentration lie in the center of the group.

Theorem 3.7. (1) *The norm of the embedding $\mathring{\mathcal{D}}_{\text{ps}}^{1,p}(\mathbf{G}) \subset L^{p^*}(\mathbf{G})$ is achieved.*

(2) *The norm of the embedding $\mathring{\mathcal{D}}_{\text{cyl}}^{1,p}(\mathbf{G}) \subset L^{p^*}(\mathbf{G})$ is achieved.*

To rule out the dichotomy case in the first part of the theorem, we prove:

Lemma 3.8. *Under the conditions of Lemma 3.3, the points g_n in the dichotomy part can be taken from the center of the group.*

Proof. Define the concentration function of v_n by

$$(3-38) \quad Q_n^{\text{ps}} \stackrel{\text{def}}{=} \sup_{h \in C(\mathbf{G})} \left(\int_{B_r(h)} dv_n \right).$$

The rest of the proof is identical to the proof of Lemma 3.3, with the remark that, in the dichotomy part, the definition of Q_n^{ps} shows that the points g_n can be taken to belong to the center. □

Proof of Theorem 3.7. We argue as for Theorem 3.1.

(1) Finding the norm of the embedding $\mathring{\mathcal{D}}_{\text{ps}}^{1,p}(\mathbf{G}) \subset \mathring{\mathcal{D}}^{1,p}(\Omega)$ leads to the variational problem

$$(3-39) \quad I^{\text{ps}} \equiv I_1^{\text{ps}} \stackrel{\text{def}}{=} \inf \left(\int_{\mathbf{G}} |Xu|^p : u \in \mathring{\mathcal{D}}_{\text{ps}}^{1,p}(\mathbf{G}), \int_{\mathbf{G}} |u|^{p^*} = 1 \right).$$

We take a minimizing sequence (u_n) , that is,

$$(3-40) \quad \int_{\mathbf{G}} |u_n|^{p^*} = 1 \quad \text{and} \quad \int_{\mathbf{G}} |Xu_n|^p \xrightarrow{n \rightarrow \infty} I^{\text{ps}}.$$

It is clear that $\mathring{\mathcal{D}}_{\text{ps}}^{1,p}(\mathbf{G})$ is invariant under the dilations (3-6). Using the Baker–Campbell–Hausdorff formula, it is easy to see that $\mathring{\mathcal{D}}_{\text{ps}}^{1,p}(\mathbf{G})$ is also invariant under the translations (3-5) by elements in the center $C(\mathbf{G})$ of \mathbf{G} . In order to extract a suitable dilated and translated subsequence of $\{u_n\}$, we have to make sure that we translate always by elements belonging to $C(\mathbf{G})$. For this, we define the concentration function of u_n as

$$(3-41) \quad \hat{Q}_n^{\text{ps}}(r) \stackrel{\text{def}}{=} \sup_{h \in C(\mathbf{G})} \left(\int_{B_r(h)} |u_n|^{p^*} dH \right).$$

We can fix $r_n > 0$ and $h_n \in C(\mathbf{G})$ such that (3-11) and (3-12) hold. Define the sequence $\{v_n\}$ as in (3-13). As mentioned above, $v_n \in \mathring{\mathcal{D}}_{\text{ps}}^{1,p}(\mathbf{G})$ as well. Equation (3-15) holds without any changes. By taking the supremum over $g \in C(\mathbf{G})$ we obtain (3-16), and by using (3-17) we obtain (3-14). At this point, we apply Lemma 3.8. The case of vanishing is ruled out from the normalization (3-14) of

the sequence $\{v_n\}$. Suppose we have dichotomy. As before, we take a sequence $R_n > 0$ such that (3-25) and (3-26) hold. We choose a cut-off function φ from the space $\mathcal{D}_{\text{ps}}^{1,p}(\mathbf{G})$, also satisfying

$$(3-42) \quad \text{supp } \varphi \subset \Omega_2(e) \quad \text{and} \quad \varphi \equiv 1 \quad \text{on } \Omega_1(e),$$

where $\Omega_r(g)$ denotes a gauge ball centered at g and of radius r , that is,

$$(3-43) \quad \Omega_r(g) = \{h \in \mathbf{G} \mid N(h^{-1}g) < r\}.$$

This can be done by setting $\varphi = \eta(N(g))$, where $\eta(t)$ is a smooth function on the real line, supported where $|t| < 2$ and with $\eta \equiv 1$ on $|t| \leq 1$. We define the cut-off functions φ_n as before, by $\varphi_n = \delta_{R_n^{-1}\tau_{g_n}}\varphi$. From Baker–Campbell–Hausdorff, they have partial symmetry with respect to the identity, since $g_n \in C(\mathbf{G})$ (see Lemma 3.8), and the gauge is a function with partial symmetry \mathbf{G} . By setting

$$(3-44) \quad A_n = \Omega_{2R_n}(g_n) \setminus \Omega_{R_n}(g_n)$$

and noting that

$$(3-45) \quad |A_n| \sim R_n^Q,$$

we see that (3-30) holds. Now, from the definition of I_{ps} , keeping in mind that φ_n and v_n have partial symmetry with respect to the identity, we obtain

$$\begin{aligned} \|v_n\|_{\mathcal{D}_{\text{ps}}^{1,p}(\mathbf{G})}^p &= \|\varphi_n v_n\|_{\mathcal{D}_{\text{ps}}^{1,p}(\mathbf{G})}^p + \|(1 - \varphi_n)v_n\|_{\mathcal{D}_{\text{ps}}^{1,p}(\mathbf{G})}^p + \varepsilon_n \\ &\geq I_{\text{ps}}(\|\varphi_n v_n\|_{L^{p^*}(\mathbf{G})}^p + \|(1 - \varphi_n)v_n\|_{L^{p^*}(\mathbf{G})}^p) + \varepsilon_n \\ &\geq I_{\text{ps}}\left(\left(\int_{B_{R_n}(g_n)} dv_n\right)^{p/p^*} + \left(\int_{\mathbf{G} \setminus B_{R_n}(g_n)} dv_n\right)^{p/p^*}\right) + \varepsilon_n \\ &\geq I_{\text{ps}}\left(\left(\int_{\mathbf{G}} dv_n^1\right)^{p/p^*} + \left(\int_{\mathbf{G}} dv_n^2\right)^{p/p^*}\right) + \varepsilon_n \\ &\geq I_{\text{ps}}(\lambda^{p/p^*} + (1-\lambda)^{p/p^*}) + \varepsilon_n. \end{aligned}$$

Letting $n \rightarrow \infty$, we come, as in (3-30), to

$$(3-46) \quad \lim_{n \rightarrow \infty} \|v_n\|_{\mathcal{D}_{\text{ps}}^{1,p}(\mathbf{G})}^p \geq I_{\text{ps}}(\lambda^{p/p^*} + (1-\lambda)^{p/p^*}) > I_{\text{ps}},$$

since $0 < \lambda < 1$ and $p/p^* < 1$. This contradicts that $\|v_n\|_{\mathcal{D}_{\text{ps}}^{1,p}(\mathbf{G})}^p \rightarrow I^{\text{ps}}$ as $n \rightarrow \infty$, which shows that the dichotomy case of Lemma 3.8 cannot occur. Hence, the compactness case holds. As in Theorem 3.1, we see that

$$\int_{\mathbf{G}} dv = 1.$$

Next, we apply [Lemma 3.5](#), with I replaced by I_{ps} . The important fact here is that the partial symmetry of the sequence $\{v_n\}$ implies that the points of concentration of $d\nu$, if they occur, must be in the center of the group. Having this in mind together with the definitions of the concentration functions, we can justify the validity of [\(3-24\)](#), and finish the proof of part (1).

(2) The vanishing case is ruled out by using the dilation (but not translation, because of the symmetries) invariance, and by normalizing the minimizing sequence with the condition $Q_n(1) = 1/2$; see [\(3-14\)](#).

Suppose that the dichotomy case occurs. We shall see that this leads to a contradiction. The points $\{g_n\}$ in the dichotomy part of [Lemma 3.3](#) must be a bounded sequence. If not, let $\varepsilon = \lambda/2$ and R as in the lemma. Because of the invariance under rotations in the layers of the functions v_n and the Haar measure dH (which is just the Lebesgue measure), for any arbitrarily fixed natural number N_0 we can find a point g_n and N_0 points on the orbit of g_n under rotations in one of the layers, such that the balls with radius R centered at all these points do not intersect. This leads to a contradiction, since the integral of the probability measure $d\nu_n$ over each of these balls is greater than $\lambda/2$. Thus, $\{g_n\}$ is a bounded sequence. This is, however, impossible since $d\nu_n$ are probability measures.

Therefore, the compactness case holds. As in the dichotomy part, we see that the sequence $\{g_n\} d\nu$ can concentrate only at the origin e . This is impossible. \square

4. Global boundedness of weak solutions

Let $p \in (1, Q)$ and denote by p^* the Sobolev conjugate $p^* = pQ/(Q - p)$, and by p' the Hölder conjugate $p' = p/(p - 1)$. Let $u \in \mathring{\mathcal{D}}^{1,p}(\Omega)$ be a weak solution, not necessarily bounded, of the equation [\(1-3\)](#) in an open set $\Omega \subset G$. “Weak solution” means that, for every $\varphi \in C_0^\infty(\Omega)$, we have

$$(4-1) \quad \int_{\Omega} |Xu|^{p-2} \langle Xu, X\varphi \rangle dH = \int_{\Omega} |u|^{p^*-2} u \varphi dH.$$

Note that $u^{p^*-1} \in L^{p^*/(p^*-1)}(\Omega) = L^{(p^*)'}$. From the definition of $\mathring{\mathcal{D}}^{1,p}(\Omega)$, we obtain that [\(4-1\)](#) holds for every $\varphi \in \mathring{\mathcal{D}}^{1,p}(\Omega)$. The main result of this section is that weak solutions as above are bounded functions. In the next theorem we prove a more general result.

Theorem 4.1. *Let $u \in \mathring{\mathcal{D}}^{1,p}(\Omega)$ be a weak solution to the equation*

$$(4-2) \quad \sum_{i=1}^m X_i(|Xu|^{p-2} X_i u) = -V|u|^{p-2} u \quad \text{in } \Omega,$$

that is,

$$(4-3) \quad \int_{\Omega} |Xu|^{p-2} \langle Xu, X\varphi \rangle dH = \int_{\Omega} V |u|^{p-2} u \varphi dH,$$

for every $\varphi \in C_0^\infty(\Omega)$.

(1) If $V \in L^{Q/p}(\Omega)$, then $u \in L^q(\Omega)$ for every $p^* \leq q < \infty$.

(2) If $V \in L^t(\Omega) \cap L^{Q/p}(\Omega)$ for some $t > Q/p$, then $u \in L^\infty(\Omega)$.

Proof. The assumption that $V \in L^{Q/p}(\Omega)$, together with the Folland–Stein inequality, shows that (4-3) holds true for any $\varphi \in \mathring{\mathcal{D}}^{1,p}(\Omega)$. This can be seen by approximating in the space $\mathring{\mathcal{D}}^{1,p}(\Omega)$ by a sequence of test functions $\varphi_n \in C_0^\infty(\Omega)$, which will allow us to put the limit function in the left-hand side of (4-3). On the other hand, the Folland–Stein inequality implies that $\varphi_n \rightarrow \varphi$ in $L^{p^*}(\Omega)$. Set $t_0 = Q/p$, and its Hölder conjugate $t_0' = t_0/(t_0 - 1)$. An easy computation gives

$$\frac{1}{t_0} + \frac{p-1}{p^*} = 1 - \frac{1}{p^*} = \frac{1}{(p^*)'}$$

Hölder's inequality then shows that $V|u|^{p-2}u \in L^{(p^*)'}(\Omega)$, which allows us to pass to the limit in the right-hand side of (4-3). We turn to the proofs of (1) and (2).

(1) It is enough to prove that, if $u \in \mathring{\mathcal{D}}^{1,p}(\Omega) \cap L^q(\Omega)$ with $q \geq p^*$, then $u \in L^{\kappa q}$ with $\kappa = p^*/p > 1$. Let $G(t)$ be a continuous, piecewise-smooth, and globally Lipschitz function on the real line, and set

$$(4-4) \quad F(u) = \int_0^u |G'(t)|^p dt.$$

Clearly, except at finitely many points, F is a differentiable function with a bounded and continuous derivative. From the chain rule (see [Garofalo and Nhieu 1996], for example) there follows that $F(u) \in \mathring{\mathcal{D}}^{1,p}(\Omega)$ is a legitimate test function in (4-3). The left-hand side, taking into account that $F'(u) = |G'(u)|^p$, can be rewritten as

$$\int_{\Omega} |Xu|^{p-2} \langle Xu, XF(u) \rangle dH = \int_{\Omega} |XG(u)|^p.$$

The Folland–Stein inequality (1-1) gives

$$(4-5) \quad \int_{\Omega} |Xu|^{p-2} \langle Xu, XF(u) \rangle dH \geq S_p \left(\int_{\Omega} |G(u)|^{p^*} \right)^{p/p^*}.$$

We choose $G(t)$ as

$$G(t) = \begin{cases} (\text{sign } t) |t|^{q/p} & \text{if } 0 \leq |t| \leq l, \\ l^{(q/p)-1} t & \text{if } l < |t|. \end{cases}$$

From the power growth of G , besides the preceding properties, this function also satisfies

$$(4-6) \quad |u|^{p-1}|F(u)| \leq C(q)|G(u)|^p \leq C(q)|u|^q.$$

The constant $C(q)$ depends also on p , but for us this is a fixed quantity. At this moment, the value of $C(q)$ is not important, but an easy calculation shows that $C(q) \leq Cq^{p-1}$, with C depending on p ; we will use this in part (2). Note that $pt'_0 = p^*$. Leave $M > 0$ to be fixed in a moment, and estimate the integral in the right-hand side of (4-3):

$$(4-7) \quad \int_{\Omega} V|u|^{p-2}uF(u) dH \\ = \int_{(|V| \leq M)} V|u|^{p-2}uF(u) dH + \int_{(|V| > M)} V|u|^{p-2}uF(u) dH \\ \leq M \int_{(|V| \leq M)} |u|^{p-1}F(u) dH \\ \quad + \left(\int_{(|V| > M)} |V|^{t_0} dH \right)^{1/t_0} \left(\int_{\Omega} (|u|^{p-1}F(u))^{t'_0} dH \right)^{1/t'_0} \\ \leq C(q)M \int_{\Omega} |G(u)|^p dH \\ \quad + C(q) \left(\int_{(|V| > M)} |V|^{t_0} dH \right)^{1/t_0} \left(\int_{\Omega} |G(u)|^{p^*} dH \right)^{p/p^*}.$$

At this point, we fix an M sufficiently large so that

$$C(q) \left(\int_{(|V| > M)} |V|^{t_0} dH \right)^{1/t_0} \leq \frac{S_p}{2},$$

which can be done because $V \in L^{t_0}$. Putting together (4-5) and (4-7) we come to our main inequality,

$$\frac{S_p}{2} \left(\int_{\Omega} |G(u)|^{p^*} dH \right)^{p/p^*} \leq C(q)M \int_{\Omega} |G(u)|^p dH \leq C(q)M \int_{\Omega} |u|^q dH.$$

By the Fatou and Lebesgue dominated convergence theorem, we can let l in the definition of G go to infinity, and obtain

$$\frac{S_p}{2} \left(\int_{\Omega} |u|^{qp^*/p} dH \right)^{p/p^*} \leq C(q)M \int_{\Omega} |u|^q dH.$$

The proof of (1) is finished.

(2) It is enough to prove that the $L^q(\Omega)$ norms of u are uniformly bounded by some sufficiently large but fixed L^{q_0} norm of u , $q_0 \geq p^*$, which is finite from (1). We shall do this by iteration. We use the function $F(u)$ from part (1) in the weak form (4-3) of our equation. The left-hand side is estimated from below as before, in (4-5). This time, though, we use Hölder's inequality to estimate from above the right-hand side:

$$(4-8) \quad \int_{\Omega} V|u|^{p-2}uF(u) dH \leq \|V\|_t \left\| |u|^{p-1}F(u) \right\|_{t'} \\ \leq \|V\|_t \left\| C(q)|G(u)|^p \right\|_{t'} \leq C(q) \|V\|_t \|u\|_{q_{t'}}^q.$$

With the estimate from below, we come to

$$S_p \|G(u)\|_{p^*}^p \leq C(q) \|V\|_t \|u\|_{q_{t'}}^q.$$

Letting $l \rightarrow \infty$, we obtain

$$(4-9) \quad \left\| |u|^{q/p} \right\|_{p^*}^p \leq \frac{C(q)}{S_p} \|V\|_t \|u\|_{q_{t'}}^q.$$

Set $\delta = p^*/(pt')$. The assumption $t > Q/p$ implies $\delta > 1$, since the latter is equivalent to $t' < p^*/p = t'_0$, as $t_0 = Q/p$. With this notation we can rewrite (4-9) as

$$(4-10) \quad \|u\|_{\delta q_{t'}} \leq \left(\frac{C(q)}{S_p} \right)^{1/q} \|V\|_t^{1/q} \|u\|_{q_{t'}}.$$

Recall that $C(q) \leq Cq^{p-1}$. At this point, we define $q_0 = p^*t'$ and $q_k = \delta^k q_0$. After a simple induction, we obtain

$$(4-11) \quad \|u\|_{q_k} \leq \prod_{j=0}^{k-1} (Cq_j^{p-1})^{t'/q_j} \|V\|_t^{t' \sum_{j=0}^{k-1} 1/q_j} \|u\|_{q_0}.$$

We observe that the right-hand side is finite,

$$(4-12) \quad \sum_{j=0}^{\infty} \frac{1}{q_j} = \frac{1}{q_0} \sum_{j=1}^{\infty} \frac{1}{\delta^j} < \infty \quad \text{and} \quad \sum_{j=1}^{\infty} \frac{\log q_j}{q_j} < \infty,$$

because $\delta > 1$. Letting $j \rightarrow \infty$, we obtain

$$\|u\|_{\infty} \leq C \|u\|_{q_0}. \quad \square$$

Remark 4.2. When Ω is a bounded open set, we trivially have $V \in L^{Q/p}(\Omega)$ whenever $V \in L^t(\Omega)$ with $t > Q/p$. Also, in this case one can obtain a uniform estimate of the $L^\infty(\Omega)$ norm of u by its $L^{p^*}(\Omega)$ norm that does not depend on the distribution function of V , as we had in the preceding theorem. This can even be

achieved in the unbounded case, assuming only that $V \in L^{Q/p}(\Omega)$, but requiring $u \in L^p(\Omega)$.

With the previous theorem proved, we turn to our original equation (1-3):

Theorem 4.3. *Take $p \in (1, Q)$ and let $\Omega \subset \mathbf{G}$ be an open set. If $u \in \mathring{D}^{1,p}(\Omega)$ is a weak solution to the equation*

$$\mathcal{L}_p u = \sum_{j=1}^n X_j (|Xu|^{p-2} X_j u) = -|u|^{p^*-2} u \quad \text{in } \Omega,$$

then $u \in L^\infty(\Omega)$.

Proof. We define $V = |u|^{p^*-p}$. From the Folland–Stein inequality we have $u \in L^{p^*}(\Omega)$, and thus $V \in L^{p^*/(p^*-p)}(\Omega)$. Since $p^*/(p^* - p) = Q/p$, part (1) of Theorem 4.1 shows that $u \in L^q(\Omega)$ for $p^* \leq q < \infty$. Therefore, $V \in L^{q/(p^*-p)}(\Omega)$ for any such q and thus, by part (2) of the same theorem, we conclude that $u \in L^\infty(\Omega)$. □

5. Regularity near the characteristic boundary

We start by introducing the geometric assumptions on the domain, and describe the regularity of the weak solutions to the Yamabe equation (1-4), which can be obtained from well-known results.

Let $\Omega \subset \mathbf{G}$ be a C^∞ domain whose boundary $\partial\Omega$ is an *orientable* hypersurface. We assume the existence of $\rho \in C^\infty(\mathbf{G})$ and $\gamma_\Omega > 0$ such that, for some $R \in \mathbb{R}$,

$$(5-1) \quad \Omega = \{g \in \mathbf{G} \mid \rho(g) < R\},$$

and such that $|D\rho(g)| \geq \gamma_\Omega > 0$ for every g in some relatively compact neighborhood K of $\partial\Omega$. We shall denote by U a bounded open set such that $\partial\Omega \subset U$, and set $\omega = \Omega \cap U$. We stress that, with this assumption, $\partial\omega$ is a compact set inside Ω . However, in Sections 5B, 5C, and 5E we shall make a different assumption, requiring that U contain only the characteristic points of the boundary, and hence $\partial\omega$ will reach $\partial\Omega$.

The assumptions on the domains are:

- *A-condition:* There exist $A, r_0 > 0$ such that, for every $g \in \partial\Omega$ and $r \in (0, r_0)$,

$$(5-2) \quad |(\mathbf{G} \setminus \Omega) \cap B(g, r)| \geq A |B(g, r)|.$$

- *Convexity:* There exists $M_1 > 0$ such that the defining function ρ of Ω satisfies the differential inequality

$$(5-3) \quad \mathcal{L}\rho \geq \frac{2}{M_1} \langle X\rho, X\psi \rangle \quad \text{in } \omega,$$

where $\psi(g) \stackrel{\text{def}}{=} |x(g)|^2$.

We remark that the A -condition is fulfilled if, for example, Ω satisfies the uniform corkscrew condition, see [Capogna and Garofalo 1998; Capogna et al. 1998]. These papers contain an extensive study of examples of domains that, in particular, satisfy (5-2). The A -condition allows us to adapt to the present setting Moser’s [1961] classical iteration arguments. On the other hand, the convexity condition allows the construction of useful barriers. It is satisfied if, for example, ρ is a strictly \mathcal{L} -superharmonic function.

From Theorem 4.3, we know that $u \in L^\infty(\Omega)$. This crucial information allows us to implement the local regularity theory of Folland and Stein [1974; Folland 1975] to conclude that $u \in C^\infty(\Omega)$.

The A -condition (5-2) permits us to adapt to the present setting the classical arguments that lead, via Moser’s iteration, to $u \in \Gamma^{0,\alpha}(\bar{\Omega})$ for some $0 < \alpha < 1$; see for example [Gilbarg and Trudinger 1983, Section 8.10]. Extending u by zero outside Ω , we can assume henceforth that

$$(5-4) \quad u \in \Gamma^{0,\alpha}(\mathbf{G}) \cap C^\infty(\Omega).$$

Given the domain Ω and an open neighborhood U , fixed as in the preceding discussion, we assume that $M_0 > 0$ has been chosen so that it fulfill the condition

$$(5-5) \quad \bar{U} \subset \{g \in \mathbf{G} \mid |x(g)|^2 \leq mM_0/4\}.$$

We shall consider the Riemannian distance $d_R(\cdot, \cdot)$ on \mathbf{G} , defined using the Euclidean metric on \mathfrak{g} via the exponential map; that is, if $g = \exp \xi$ and $h = \exp \eta$, we set

$$(5-6) \quad d_R(g, h) = |\xi - \eta|.$$

It is straightforward to estimate the Riemannian distance by the Carnot–Carathéodory distance:

$$(5-7) \quad d_R(g, h) \leq d(g, h).$$

The estimate in the other direction was proven by Tanaka [1975]. He works in a more general situation than ours, so we state the result as relevant to our setting.

Theorem 5.1. *Let \mathbf{G} be a Carnot group of step r , and take $g_0 \in \mathbf{G}$. For every ball $B(g_0, R)$ there exists a constant $C = C(\mathbf{G}, R)$ such that, if $g, h \in B(g_0, R)$, then*

$$(5-8) \quad d(g, h) \leq C d_R(g, h)^{1/r}.$$

5A. Barrier functions. The barrier functions defined in the next theorem are crucial to the rest of the section. The convexity property of the boundary is essentially what makes these functions useful.

Theorem 5.2. *Let $\Omega \subset G$ be a smooth connected bounded open set that satisfies the convexity condition (5-3). Let $M \geq \max\{M_0, M_1\}$. For $0 < \alpha \leq 1$, we define*

$$\Psi_\alpha = (R - \rho)^\alpha e^{-\psi/M}.$$

Under the stated hypothesis, for every $g \in \omega$,

$$(5-9) \quad \mathcal{L}\Psi_\alpha(g) \leq -\frac{m}{M}\Psi_\alpha(g).$$

We remark that, when working with \mathcal{L} -harmonic functions, one takes $\alpha = 1$. We also consider $\alpha < 1$ because of its use in the case of nonlinear equations, for example, Yamabe-type equations. The proof of the above theorem can be found in [Garofalo and Vassilev 2000].

Let η be a smooth vector field defined in U that is transversal to the boundary, that is,

$$(5-10) \quad \eta\rho(g_0) \equiv \left(\frac{d}{dt}\rho(g_t)\right)\Big|_{t=0} \neq 0, \quad g_0 \in \partial\Omega,$$

where $\frac{d}{dt}g_t = \eta(g_t)$.

Using the compactness of $\partial\Omega$, we can assume, possibly after taking a smaller U , that there exists a constant $\delta > 0$ such that

$$(5-11) \quad \eta\rho(g_0) \geq \delta > 0 \quad \text{for } g_0 \in \partial\Omega.$$

We note that the transversality condition (5-11) implies that the trajectories of η that start from points of $\partial\Omega$ fill a full open set ω , interior to Ω . Possibly by shrinking the set U , we can assume that $\omega = \Omega \cap U$. To fix the notation, we suppose that there exists $t_0 \in (0, 1)$ such that

$$g_t \in \omega \quad \text{for } 0 < t < t_0 \text{ and } g_0 \in \partial\Omega.$$

We shall hereafter use this generic transversal vector field and notation.

Lemma 5.3. *There exist $C_1, C_2 > 0$ such that, for every $g_0 \in \partial\Omega$ and $0 \leq t \leq t_0$,*

$$(5-12) \quad C_1 t^\alpha \leq \Psi_\alpha(g_t) \leq C_2 t^\alpha.$$

Proof. Under the assumptions we made, the proof follows from the Taylor formula. Indeed,

$$(5-13) \quad R - \rho(g_t) = t\eta\rho(g_0)(1 + O(t)),$$

where $O(t)$ denotes a function bounded by Ct , uniformly in $g_0 \in \partial\Omega$ and $0 \leq t \leq t_0$.

It is clear from the latter identity that (5-12) follows by using (5-11), the smoothness of ρ , and the boundedness of Ω . □

5B. Lipschitz estimates near the boundary. We start with a theorem about the Green’s function, since its proof is very simple because of the existence of barriers.

An important observation is that, thanks to the results of [Derridj \[1971\]](#), the Green’s function is a smooth function up to the noncharacteristic boundary; indeed, its \mathcal{L} -Laplacian vanishes near the boundary, and hence it is a smooth function. Therefore, when working with equations with a smooth right-hand side, we shall consider $\omega = U \cap \Omega$, where U is a sufficiently small open neighborhood of the characteristic set Σ defined in [\(5-37\)](#), rather than $\partial\Omega \subset U$ as we do in the case of the Yamabe equation; see the definitions of the convexity and A -conditions in [Section 5A](#).

Theorem 5.4. *Let $u \in C(\bar{\omega})$ satisfy*

$$\mathcal{L}u = 0 \quad \text{in } \omega, \quad u = 0 \quad \text{on } \partial\Omega.$$

If the convexity condition is satisfied in a neighborhood ω of the characteristic boundary Σ , then there exists a constant $C = C(\mathbf{G}, \Omega, u)$ such that

$$(5-14) \quad |u(g)| \leq C d_R(g, \partial\Omega) \quad \text{for every } g \in \omega.$$

Remark 5.5. Note that in this theorem the right-hand side uses the Riemannian distance. This is important for the estimates on derivatives that involve vectors not only from the first layer. Clearly, the same inequality holds for the Carnot–Carathéodory distance.

Proof. The proof uses Ψ as a barrier, and Bony’s maximum principle. As we saw before the statement of the theorem, we have

$$u \in C^\infty(\bar{\omega} \setminus \Sigma).$$

Equation [\(5-9\)](#) with $\Psi = \Psi_1$ shows that

$$\pm \mathcal{L}u(g) = 0 \geq \mathcal{L}C\Psi(g), \quad g \in \omega.$$

On the other hand, for a sufficiently large constant C , we have the estimate

$$(5-15) \quad C\Psi(g) \geq |u(g)| \quad \text{on } \partial\omega.$$

To see that [\(5-15\)](#) holds, we argue as follows. Its validity is clear on $\partial\Omega$ since both u and Ψ vanish there. On the other hand, the compact set $\partial\omega \cap \Omega$ is at a fixed distance away from the characteristic boundary, and the estimate [\(5-15\)](#) follows since $\Psi, u \in C^\infty(\partial\omega \cap \Omega)$ and from [Lemma 5.3](#); see the proof of [Theorem 5.14](#) for further details.

We can now apply Bony’s maximum principle ([Theorem 2.1](#)) and conclude that the bound holds inside as well:

$$|u(g)| \leq C\Psi(g), \quad g \in \omega.$$

This implies that u satisfies the Riemannian Lipschitz estimate

$$(5-16) \quad |u(g)| \leq C d_R(g, \partial\Omega), \quad g \in \omega,$$

which completes the proof. \square

The next theorem was proved in [Garofalo and Vassilev 2000] by also requiring that Ω be uniformly starlike along Σ . A minor modification shows that this is unnecessary. We shall use g_t as in (5-10).

Theorem 5.6. *Let Ω be a C^∞ open bounded set in a Carnot group \mathbf{G} , and u a solution of the Yamabe equation (1-4). If the convexity and A-conditions are satisfied, then there exists a constant $C = C(\mathbf{G}, \Omega, u)$ such that*

$$u(g) \leq C d_R(g, \partial\Omega) \quad \text{for every } g \in \omega.$$

Proof. We begin by observing that, thanks to (5-4) and $u = 0$ on $\partial\Omega$, we have, for any $g_0 \in \partial\Omega$,

$$(5-17) \quad u(g_t) \leq C d(g_t, g_0)^\alpha.$$

Theorem 5.1 now gives, for every $g_0 \in \partial\Omega$ and $0 \leq t \leq t_0$,

$$(5-18) \quad d(g_t, g_0) \leq C d_R(g_t, g_0)^{1/r} \leq C t^{1/r}$$

for some constant $C = C(\Omega) > 0$. Using (5-17) and (5-18), and setting $\alpha_0 = \alpha/r$, we infer that

$$(5-19) \quad u(g_t) \leq C t^{\alpha_0}$$

for every $g_0 \in \partial\Omega$ and $0 \leq t \leq t_0 < 1$.

We now let $\sigma = 2^* - 1 = (Q + 2)/(Q - 2)$. In the sequel it will be important that $\sigma > 1$. Since it is clear that (5-19) continues to hold if in the right-hand side we raise t to any exponent smaller than α_0 , we assume in what follows that $\sigma\alpha_0 < 1$ and that there is some $n \in \mathbb{N}$ for which

$$(5-20) \quad \sigma^n \alpha_0 = 1.$$

Next, we use the barriers constructed in Theorem 5.2. For any point $g_t \in \omega$ we have, from (1-4), (5-19), (5-12), and (5-9), that

$$\begin{aligned} -\mathcal{L}u(g_t) &= u(g_t)^\sigma \leq C t^{\sigma\alpha_0} \\ &\leq C C_1^{-1} \Psi_{\sigma\alpha_0}(g_t) \leq -C C_1^{-1} M m^{-1} \mathcal{L} \Psi_{\sigma\alpha_0}(g_t) \\ &= -\mathcal{L}(C^* \Psi_{\sigma\alpha_0})(g_t). \end{aligned}$$

Keeping in mind that, as g_0 varies in $\partial\Omega$ and t in $(0, t_0)$, the point g_t covers ω , we have proved that

$$\mathcal{L}(C^* \Psi_{\sigma\alpha_0} - u) \leq 0 \quad \text{in } \omega.$$

At this point, we observe that (possibly using a constant larger than C^*) we also have the estimate

$$(5-21) \quad C^* \Psi_{\sigma\alpha_0} \geq u \quad \text{on } \partial\omega.$$

To see that (5-21) holds, we argue as follows. It is clear that (5-21) holds on $\partial\omega \cap \partial\Omega$, since both u and $\Psi_{\sigma\alpha_0}$ vanish there. On the other hand, the compact set $\partial\omega \cap \Omega$ is at a fixed distance away from the boundary, and therefore, using that u and $\Psi_{\sigma\alpha_0}$ are smooth there and $\Psi_{\sigma\alpha_0}$ does not vanish, one trivially obtains the estimate by using the maximum of u .

We can now apply to ω Bony's maximum principle from [Theorem 2.1](#) to infer that a similar estimate also holds in ω . From this result and from the right-hand side of (5-12), we conclude that, for every $g_0 \in \omega$ and $0 < t < t_0$,

$$(5-22) \quad u(g_t) \leq Ct^{\sigma\alpha_0},$$

which shows that we have improved on (5-19). It is now clear that, by repeating the above arguments n times, where n is as in (5-20), we shall reach the desired conclusion (5-16). \square

5C. Estimates for the Green's function.

Theorem 5.7. *Let $u \in C^\infty(\bar{\omega} \setminus \partial\Omega) \cap C(\bar{\omega})$ satisfy*

$$\mathcal{L}u = 0 \quad \text{in } \omega, \quad u = 0 \quad \text{on } \partial\Omega.$$

If the convexity condition (5-3) is satisfied, then

$$Xu \in L^\infty(\omega).$$

The proof of [Theorem 5.7](#) is an immediate consequence of [Theorem 5.4](#), with the help of the Schauder-type estimates proved in [[Danielli and Garofalo 1998](#)].

Proof. Fix an arbitrary $g \in \omega$. With $r = \text{dist}(g, \partial\omega)/2$, consider the ball $B(g, r) \subset \bar{B}(g, r) \subset \omega$. Applying the interior Schauder estimates to the \mathcal{L} -harmonic function u , one has

$$(5-23) \quad |Xu(g)| \leq \frac{C}{r} \sup_{B(g,r)} |u|,$$

for some constant $C = C(\mathbf{G})$. At this point, we invoke [Theorem 5.4](#), which implies, for any $g' \in B(g, r)$, the inequality

$$|u(g')| \leq C'd(g', \partial\Omega) \leq C'(d(g', g) + d(g, \partial\Omega)) \leq 2C'r,$$

with a constant $C' = C'(\mathbf{G}, \omega, u)$. Substitution in (5-23) finishes the proof. \square

In the next theorems we show the boundedness of derivatives along some transversal to the boundary vector fields. We denote by Σ the characteristic boundary, which is defined in (5-37). Also, as in Theorem 5.4, $\omega = U \cap \Omega$ where U is an open neighborhood of Σ . We further set $\Delta = \partial\Omega \cap U$.

Theorem 5.8. *Suppose η is a vector field from the center of the Lie algebra of \mathbf{G} . Let $\tilde{\eta}$ be the corresponding left-invariant vector field. Assume further that $\Omega \subset \mathbf{G}$ satisfies the convexity condition (5-3) and $\tilde{\eta}$ is transversal to Σ . If u satisfies*

$$\mathcal{L}u = 0 \quad \text{in } \omega, \quad u = 0 \quad \text{on } \partial\Omega,$$

then

$$\tilde{\eta}u \in L^\infty(\omega).$$

Proof. For t very close to 0, we define $\omega_t = \{g \exp t\eta \mid 0 < t < t_0, g \in \omega\} \cap \omega$ and $\partial\Delta_t = \{g_0 \exp t\eta \mid g_0 \in \Delta\}$, and consider the difference quotients

$$(5-24) \quad \varphi_t = \frac{1}{t}(R_{\exp t\eta}u - u), \quad g \in \omega_t,$$

where $R_h u(g) = u(gh)$ is the right-translations operator. We claim that there exists a constant $C > 0$ such that, for all t sufficiently close to 0 and $g \in \omega_t$, one has

$$(5-25) \quad |\varphi_t(g)| \leq C.$$

Suppose the claim (5-25) is true. Passing to the limit as $t \rightarrow 0$, we conclude that $|\tilde{\eta}u(g)| \leq C$ for every $g \in \omega$, thus establishing the theorem.

We turn to the proof of (5-25). Note that

$$(5-26) \quad \mathcal{L}\varphi_t(g) = 0.$$

That is, each of the functions φ_t is \mathcal{L} -harmonic when u is harmonic. This follows since the considered \mathcal{L} -Laplacian is left invariant and, when η is in the center of \mathfrak{g} , we have that $h = \exp t\eta$ is in the center of the group, and hence the right and left translations by h coincide.

From Bony's maximum principle (Theorem 2.1), it is therefore enough to prove that, for some t_1 close to 0, (5-25) holds for $g \in \partial\omega_t$ and $t \in (0, t_1)$. Note that $\partial\omega_t = \Delta_t \cup (\partial\omega \cap \Omega)$. We analyze the two portions separately. Since any point $g \in \Delta_t$ can be written as $g = g_0 \exp t\eta$ for some $g_0 \in \Delta$, we have

$$(5-27) \quad \varphi_t(g) = \frac{u(g_0 \exp t\eta) - u(g_0)}{t}.$$

We recall Theorem 5.4, which gives

$$(5-28) \quad |u(g_0 \exp t\eta)| \leq C d_R(g_0 \exp t\eta, g_0) \leq Ct,$$

with a fixed constant C when g_0 belongs to a fixed compact set. This estimate shows that

$$(5-29) \quad |\varphi_t(g)| \leq C, \quad g \in \partial\Omega_t.$$

Finally, the same inequality on $\partial\omega \cap \Omega$ follows from the C^∞ regularity of u in Ω . In conclusion, we have proved the claim (5-25), and therefore the theorem. \square

Remark 5.9. The same proof can be used to show boundedness of the derivative along the radial vector field, assuming that it is transversal to $\partial\Omega$. The reason is that the corresponding difference quotients are \mathcal{L} -harmonic as well.

5D. Regularity of solutions of the Yamabe equation. Let u be a weak nonnegative solution of the Yamabe-type equation (1-4). We note again that, because of the present lack of boundary Schauder estimates (except in the case of the Heisenberg group: see 5.13), in the next theorem the convexity assumption is required to hold globally. Accordingly, for domains in the Heisenberg group, convexity will be assumed only near the characteristic boundary; see Theorem 5.14.

Theorem 5.10. *Let Ω be a C^∞ open bounded set in a Carnot group \mathbf{G} , and u a solution of (1-4). If the convexity assumption (5-3) and the A-condition (5-2) are satisfied, then*

$$Xu \in L^\infty(\Omega).$$

Proof. Since u is smooth away from $\partial\Omega$, in order to prove the theorem it will be enough to show that

$$(5-30) \quad Xu \in L^\infty(\omega),$$

where ω is fixed as before.

We begin by introducing $v = u^{2^*-1} * \Gamma$, where Γ is the positive fundamental solution of \mathcal{L} , that is, $\mathcal{L}\Gamma = -\delta$. According to [Folland 1975, Corollary 2.8], v satisfies the equation $\mathcal{L}v = -u^{2^*-1}$. Since by (5-4) u^{2^*-1} is in $\Gamma^{0,\beta}(\mathbf{G})$ for some $\beta \in (0, 1)$ (and u is compactly supported in \mathbf{G}), from [Folland 1975, Theorem 6.1] we have

$$(5-31) \quad v \in \Gamma_{\text{loc}}^{2,\beta}(\mathbf{G}).$$

Therefore, if we let $w \stackrel{\text{def}}{=} u - v$, to prove (5-30) it is enough to show it for w , that is, to prove that $Xw \in L^\infty(\omega)$. We notice that w is \mathcal{L} -harmonic, that is, $\mathcal{L}w = 0$ in Ω .

From Theorem 5.6,

$$u(g) \leq Cd(g, \partial\Omega) \quad \text{for every } g \in \omega.$$

Since we know that $u \in C^\infty(\bar{\Omega} \setminus \omega)$, we conclude that

$$(5-32) \quad u(g) \leq Cd(g, \partial\Omega) \quad \text{for every } g \in \Omega;$$

see [Theorem 2.3](#). Fix a point $g \in \omega$. With $r = \text{dist}(g, \partial\Omega)/2$, consider the ball $B(g, r) \subset \bar{B}(g, r) \subset \Omega$. Applying the interior Schauder estimates in [Theorem 2.2](#) to the \mathcal{L} -harmonic function $w - w(g)$, one has

$$(5-33) \quad |Xw(g)| \leq \frac{C}{r} \sup_{B(g,r)} (w - w(g)).$$

Note that (5-32) gives, for $g' \in B(g, r)$,

$$(5-34) \quad u(g') \leq C \text{dist}(g', \partial\Omega) \leq C(d(g', g) + \text{dist}(g, \partial\Omega)) \leq Cr.$$

Since $w = u - v$, in view of (5-34) and (5-32) one has, for $g' \in B(g, r)$,

$$(5-35) \quad |w(g') - w(g)| \leq (u(g') + u(g)) + |v(g') - v(g)| \leq C(r + |v(g') - v(g)|).$$

Finally, we observe that (5-31) implies that $v \in \mathcal{L}^{1,\infty}(\Omega)$. Applying [Theorem 2.3](#) once more, we conclude that

$$|v(g) - v(g')| \leq Cd(g, g') \quad \text{for } g, g' \in \Omega.$$

Substitution of this information in (5-35) gives

$$\sup_{B(g,r)} (w - w(g)) \leq Cr.$$

Combining the latter inequality with (5-33) brings the sought-for conclusion $Xw \in L^\infty(\omega)$. This finishes the proof of [Theorem 5.10](#). \square

To end this section, we note that one can show the boundedness of the radial derivative of solutions of the Yamabe equation. This was done in [[Garofalo and Vassilev 2000](#), Theorem 4.7], which we state below. The proof requires that \mathbf{G} be of step 2 and the considered domain be C^∞ , bounded, connected, uniformly starlike with respect to one of its points, and satisfying (5-2) and (5-3).

Theorem 5.11. *Let \mathbf{G} be a Carnot group of step 2. Consider a C^∞ connected, uniformly starlike, bounded open set $\Omega \subset \mathbf{G}$ satisfying (5-2) and (5-3). If u is a weak solution of (1-4), then*

$$(5-36) \quad Zu \in L^\infty(\Omega).$$

5E. The estimates near the characteristic boundary for a domain in the Heisenberg group. Let Ω be a smooth domain in a Carnot group. Denote by Σ its characteristic set with respect to the system $X = \{X_1, \dots, X_m\}$, that is,

$$(5-37) \quad \Sigma = \{g \in \partial\Omega \mid X_j(g) \in T_g(\partial\Omega), j = 1, \dots, m\}.$$

Clearly, when the domain is bounded this is a compact subset of the boundary. In the next well-known theorem we denote by $N = \sum_{j=1}^r \dim V_j$ the topological dimension of \mathbf{G} . The symbol H_{N-1} denotes the $(N-1)$ -dimensional Hausdorff measure constructed using the Riemannian distance on \mathbf{G} .

Theorem 5.12. *Let $\Omega \subset \mathbf{G}$ be a C^∞ domain. One has*

$$H_{N-1}(\Sigma) = 0.$$

This result is due to [Derridj \[1972\]](#). A more refined version has been recently proved by [Franchi and Wheeden \[1997\]](#).

We now consider the Heisenberg group \mathbb{H}^n and take $\Omega \subset \mathbb{H}^n$. We shall use the notation used until now, except that U will be a sufficiently small neighborhood of Σ — that is, $\Sigma \Subset U$ — and $\omega = U \cap \Omega$ as before. We also set $\Delta = \partial\Omega \cap U$. We recall the following Schauder estimates [[Jerison 1981a](#)]:

Theorem 5.13. *Let Ω be a bounded C^∞ domain in the Heisenberg group \mathbb{H}^n , and let $\varphi \in C_0^\infty(\mathbb{H}^n)$ be supported in a small neighborhood of a noncharacteristic point $g_0 \in \partial\Omega$. Given $f \in \Gamma^{k,\alpha}(\overline{\Omega})$, $k \in \mathbb{N} \cup \{0\}$, $0 < \alpha < 1$, for the unique solution u to the Dirichlet problem for the Kohn sub-Laplacian*

$$(5-38) \quad \mathcal{L}u = f \quad \text{in } \Omega, \quad u = 0 \quad \text{on } \partial\Omega,$$

one has $\varphi u \in \Gamma^{k+2,\alpha}(\overline{\Omega})$.

Following the arguments in the preceding sections, we can prove the following:

Theorem 5.14. *Let Ω be a smooth connected bounded open set in \mathbb{H}^n , satisfying the A-condition. Suppose that the convexity condition is satisfied in an interior neighborhood ω of its characteristic set. If u is a weak nonnegative solution of the CR Yamabe equation, then the horizontal gradient of u is bounded in Ω . If in addition $\partial\Omega$ is uniformly starlike along Σ , then the radial derivative Zu is bounded in Ω .*

Proof. As already mentioned, $u \in \Gamma^{0,\alpha}(\mathbb{H}^n) \cap C^\infty(\Omega)$. From [Theorem 5.13](#), by taking into account that $(Q+2)/(Q-2) > 1$, we have $u \in \Gamma_{\text{loc}}^{2,\alpha}(\overline{\Omega} \setminus \Sigma)$. Now we can argue as in [Theorem 5.10](#), which holds as long as we have [Theorem 5.6](#) and [Theorem 5.11](#), with the only difference that now $\partial\Omega$ should be replaced with Δ . However, by doing this we see that $\partial\omega = \Delta \cup (\partial\omega \setminus \Delta)$, with $\partial\omega \setminus \Delta$ reaching $\partial\Omega$. There are two places where this is important:

- (1) in [Theorem 5.6](#), for (5-21);
- (2) in [Theorem 5.11](#), for the bound on $\partial\omega_\lambda \setminus \partial\Omega_\lambda$, in the proof's last paragraph.

We start with (1). The set $\Lambda = (\partial\omega \cap \overline{\Omega}) \setminus \Delta$ is at a fixed distance away from the characteristic set Σ . Therefore, for every $g_0 \in \Lambda$, there exists $j \in \{1, \dots, m\}$ such

that $X_j \rho(g_0) \neq 0$. By continuity, the trajectories of X_j fill a (sufficiently small) full neighborhood V_{g_0} of g_0 . This means that there exists $t_0 = t(g_0) > 0$ such that every $g \in \Omega \cap V_{g_0}$ can be written as $g_1 \exp t X_j$ for some $g_1 \in \partial\Omega \cap V_{g_0}$ and some $t \in (0, t_0)$. Using the uniform transversality of X_j to $\partial\Omega$ in $\Omega \cap V_{g_0}$ and Taylor's formula, we infer the existence of $C = C(g_0) > 0$ such that

$$(5-39) \quad |R - \rho(g_1 \exp t X_j)| \geq C|t|$$

for every $g_1 \in \partial\Omega \cap V_{g_0}$ and $0 < t < t_0$. We now use that $u \in \Gamma_{\text{loc}}^{2,\alpha}(\bar{\Omega} \setminus \Sigma)$ to deduce the existence of a constant $C^* = C^*(u, g_0) > 0$ such that

$$u(g_1 \exp t X_j) \leq C^*|t| \leq C^*|t|^{\alpha\sigma_0}$$

for every $g_1 \in \partial\Omega \cap V_{g_0}$ and $0 < t < t_0$. The latter inequality and (5-39) allow us to conclude that (5-21) does hold in the set $\Omega \cap V_{g_0}$, for a constant depending on u and g_0 . By a finite-covering argument, we see that (5-21) continues to hold in the intersection of a small neighborhood of $\partial\Omega$ with Λ . We can thus separate from $\partial\Omega$. Once inside Ω , we can use the C^∞ smoothness of u to conclude that (5-21) holds on the remaining portion of $\partial\omega \cap \Omega$ as well. This proves (1).

The embedding theorem 5.25 in [Folland 1975] implies that

$$(5-40) \quad \Gamma_{\text{loc}}^{2,\beta}(\mathbf{G}) \subset \Lambda_{\text{loc}}^{1,\beta/2}(\mathbf{G}) = C_{\text{loc}}^{1,\beta/2}(\mathbf{G}),$$

where the latter space denotes the standard Hölder class with respect to the Riemannian distance $d_R(\cdot, \cdot)$ on \mathbf{G} . The proof of (2) follows from the above embedding. □

Remark 5.15. The previous theorem can be proved for any group of step 2, by requiring also that $u \in \Gamma_{\text{loc}}^{2,\alpha}(\bar{\Omega} \setminus \Sigma)$. With this assumption, we can also prove that

$$\tilde{\eta}u \in L^\infty(\Omega),$$

provided that $\tilde{\eta}$ is transversal to Σ . Here, $\eta \in V_2$ and $\tilde{\eta}$ is the corresponding left-invariant vector field on a group \mathbf{G} of step 2.

We want to give an example of domains in groups of step 2 satisfying the assumptions of the preceding sections. Let \mathbf{G} be a Carnot group of step 2. Using exponential coordinates, we define the function

$$(5-41) \quad f_\varepsilon(g) = ((\varepsilon^2 + |x(g)|^2)^2 + 16|y(g)|^2)^{1/4}, \quad \varepsilon \in \mathbb{R}.$$

For $R > 0$ and $\varepsilon \in \mathbb{R}$ with $\varepsilon^2 < R^2$, consider the C^∞ bounded open set

$$(5-42) \quad \Omega_{R,\varepsilon} = \{g \in \mathbf{G} \mid f_\varepsilon(g) < R\}.$$

When $\varepsilon = 0$, it is clear that $\Omega_{R,\varepsilon}$ is nothing but a gauge pseudo-ball centered at the group identity e , except that the natural gauge was defined in (2-5) without the

factor 16. Here, we have introduced such a factor for the purpose of keeping a consistent definition with the case of groups of Heisenberg type. For all practical purposes, the reader can ignore it and identify f_0 in (5-41) with (2-5). For $g \in \mathbf{G}$, we set $\Omega_{R,\varepsilon}(g) = \{h \in \mathbf{G} \mid f_\varepsilon(g^{-1}h) < R\} = g \Omega_{R,\varepsilon}$.

Theorem 5.16. *Let \mathbf{G} be a Carnot group of step 2. For every $\varepsilon \in \mathbb{R}$ with $\varepsilon^2 < R^2$, the domain $\Omega_{R,\varepsilon}(g)$ satisfies the A-condition, the convexity condition (globally), and is uniformly starlike.*

For the proof of this theorem see [Garofalo and Vassilev 2000].

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