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Regularity of minimizers of the calculus of variations in Carnot groups via hypoellipticity of systems of Hörmander type

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Abstract. We prove the hypoellipticity for systems of Hörmander type with constant coefficients in Carnot groups of step 2. This result is used to implement blow-up methods and prove partial regularity for local minimizers of non-convex functionals, and for solutions of non-linear systems which appear in the study of non-isotropic metric structures with scalings. We also establish estimates of the Hausdorff dimension of the singular set.

1. Introduction

The present paper constitutes the first part of a project aimed at developing the regularity theory for local minimizers of non-convex variational functionals and for weak solutions of quasi-linear systems in a class of non-isotropic spaces known as Carnot groups. Here, we focus on the non-evolutive case and on the model situation of groups of step 2. In a series of forthcoming articles we will study the technically more involved case of Carnot groups of arbitrary step, and we will also prove partial regularity results for quasi-linear parabolic systems which arise in the study of the harmonic maps flow. We will also consider the regularity theory for variational functionals and systems with special structure.

A Carnot group of step $r \geq 1$ is a simply connected Lie group \mathbf{G} whose Lie algebra admits a vector space decomposition in r layers $\mathfrak{g} = V^1 \oplus \dots \oplus V^r$ having the two properties: (i) \mathfrak{g} is stratified, i.e., $[V^1, V^j] = V^{j+1}$, $j = 1, \dots, r-1$; (ii) \mathfrak{g} is r -nilpotent, that is, $[V^j, V^r] = \{0\}$, $j = 1, \dots, r$. We note that the exponential map $\exp : \mathfrak{g} \rightarrow \mathbf{G}$ defines a global diffeomorphism. Such a group is equipped with a natural family of non-isotropic dilations. One starts with defining dilations Δ_λ on the Lie algebra as follows. If $X = X_1 + \dots + X_r \in \mathfrak{g}$, with $X_j \in V^j$, $j = 1, \dots, r$, then $\Delta_\lambda(X) = \lambda X_1 + \dots + \lambda^r X_r$. The group dilations $\delta_\lambda : \mathbf{G} \rightarrow \mathbf{G}$ are then introduced by lifting Δ_λ via the exponential map $\exp : \mathfrak{g} \rightarrow \mathbf{G}$, i.e.,

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$\delta_\lambda(x) = \exp \circ \Delta_\lambda \circ \exp^{-1}(x)$. We denote by dx the bi-invariant Haar measure on \mathbf{G} obtained by pushing forward the standard Lebesgue measure on \mathfrak{g} via the exponential map. When we say that a set has measure zero in \mathbf{G} we mean this with respect to dx . One easily recognizes that $(d \circ \delta_\lambda)x = \lambda^Q dx$, where the number

$$Q = \sum_{j=1}^r j \dim(V^j)$$

is the so-called *homogeneous dimension* of \mathbf{G} . Such number is larger than the topological dimension $\sum_{j=1}^r \dim(V^j)$ of the group.

The homogeneous structure of \mathbf{G} generated by the group dilations $\{\delta_\lambda\}_{\lambda>0}$ plays an important role in the analysis developed in this paper. Another aspect which will be equally important is the pseudo-metric structure which one associates to the stratification of \mathfrak{g} in the following natural way. We assume that \mathfrak{g} is equipped with an inner product $\langle \cdot, \cdot \rangle$ with respect to which the V^j 's are mutually orthogonal, and let $|\cdot| = (\langle \cdot, \cdot \rangle)^{1/2}$. For $X = X_1 + \dots + X_r \in \mathfrak{g}$ one defines

$$\|X\|_{\mathfrak{g}} = \left(\sum_{i=1}^r |X_i|^{2r/i} \right)^{2r},$$

and

$$\|x\|_{\mathbf{G}} = \|X\|_{\mathfrak{g}} \quad \text{if } x = \exp X.$$

Henceforth, we will ordinarily omit the subscripts and simply write $\|x\|$ and $\|X\|$. The homogeneous pseudo-norm on \mathbf{G} induces the following pseudo-distance $d(x, y) = \|x^{-1}y\|$. One has for every $x, y, z \in \mathbf{G}$ and for any $\lambda > 0$

$$d(zx, zy) = d(x, y), \quad d(\delta_\lambda(x), \delta_\lambda(y)) = \lambda d(x, y).$$

Using these properties it is easy to recognize that the Haar measure of $B(x, R) = \{y \in \mathbf{G} \mid d(x, y) < R\}$ is given by $|B(x, R)| = \omega_{\mathbf{G}} R^Q$, where $\omega_{\mathbf{G}} = |B(e, 1)|$ and e is the group identity. This growth of the pseudo-balls accounts for the role of a critical dimension played by the number Q in Theorems 1.1 and 1.2.

Carnot groups indirectly owe their name to the foundational paper of Carathéodory [Ca] on Carnot thermodynamics. Such Lie groups are similar to Euclidean spaces (which are themselves abelian Carnot groups with trivial step 1) in that they possess translations and dilations, and come with a measure which is invariant with respect to these operations. They also have a pseudo-metric structure, but when the step $r > 1$, then the non-abelian group law produces a strong distortion of the underlying geometry. As a consequence, the basic properties of the latter are extremely different, and much harder to grasp, than those of Euclidean spaces. There exists a very wide variety of situations in which Carnot groups appear as the natural setting, and it is impossible to attempt here even a brief description. Besides Carathéodory's cited paper, we mention the famous accessibility paper of Chow [Cho], which in many respects paved the way to the modern developments. In connection with optimal control theory and sub-Riemannian geometry one

should also see the articles [Br] and [Su], and the recent book [BR]. In partial differential equations Carnot groups first implicitly appeared in Kolmogorov's short note [Ko]. Thirty years later Hörmander [H] proved his fundamental theorem stating that, given a system X_o, X_1, \dots, X_m of smooth vector fields in \mathbb{R}^n , the condition $\text{rank Lie}\{X_o, X_1, \dots, X_m\} \equiv n$ is sufficient for the hypoellipticity of the partial differential operator $L = \sum_{j=1}^m X_j^2 + X_o$. These operators include trivially all elliptic and parabolic operators (with smooth coefficients), but in the important case when $m < n - 1$, they define a class which fails to be either elliptic or parabolic. For instance, Kolmogorov's equation for the probability density of a system with $2n$ degrees of freedom is precisely of this type. Here is the basic model in \mathbb{R}^3 : $X_1^2 u + X_o u = 0$, where $X_1 = \partial/\partial x$ and $X_o = -x\partial/\partial y - \partial/\partial t$, and one has $[X_1, X_o] = -\partial/\partial y$. Another operator, which is of paramount importance, is the real part of the Kohn sub-Laplacian on the Heisenberg group \mathbb{H}^n . In the three-dimensional model \mathbb{H}^1 such operator is given by $L = X_1^2 + X_2^2$, where $X_1 = \partial/\partial x + 2y\partial/\partial t$, $X_2 = \partial/\partial y - 2x\partial/\partial t$, and one has the commutation relation $[X_1, X_2] = -4\partial/\partial t$. The Heisenberg group is the most important Carnot group of step 2. For its description we refer the reader to section two. Carnot groups are the natural osculating spaces (or blow-up limits) for Cauchy-Riemann [FS], [S], and Carnot-Carathéodory manifolds [Pa], [Gr], [GN], [HK], [JX], [Che], [AK1], [AK2], [Ga]. In continuum mechanics they occur as limits of crystalline structures [Chr], but they also play an important role in other applied sciences, such as mechanical engineering or robotics [CK]. These Lie groups are thus the natural testing ground for problems in Calculus of Variations which arise in the study of non-isotropic metric structures with scalings. In this paper we focus on the more elementary case of groups of step 2. The reason for this choice is that the study of these models incorporates most of the essential ideas. Many of the geometric obstructions of the more general situation are already present for groups of step 2, but the technical aspect is considerably less complicated.

Given a Carnot group of step 2 with Lie algebra $\mathfrak{g} = V^1 \oplus V^2$, we fix an orthonormal basis $X = \{X_1, \dots, X_m\}$ of the first layer V^1 of \mathfrak{g} . Thanks to the stratification assumption, the elements of X and their brackets generate the whole algebra. We continue to use the same symbols X_j for the corresponding left-invariant vector fields on \mathbf{G} . For a function $u = (u^1, \dots, u^N) : \mathbf{G} \rightarrow \mathbb{R}^N$, the matrix $\{Xu\}_{i,j} = X_i u^j$ denotes the Jacobian matrix of u with respect to such basis. Since $\{Xu\}$ represents differentiation only with respect to a basis of the first layer V^1 , we will call it the *horizontal Jacobian matrix* of u . The main objective of this paper is to establish regularity properties in the non-horizontal (commutators) directions starting from structural assumptions which involve *only* control of the horizontal ones.

One of our main purposes is to study the partial regularity of local minimizers of functionals of the type

$$(1.1) \quad \mathcal{F}(u, \Omega) = \int_{\Omega} F(Xu) dp,$$

where $F : \mathbb{R}^{mN} \rightarrow \mathbb{R}^+$ is a C^2 function satisfying the following structural assumptions:

For some $2 \leq q < \infty$, one has for every $\xi, \eta \in \mathbb{R}^{mN}$,

$$F(\xi) \leq C(1 + |\xi|^q),$$

$$\sum_{\alpha, \beta=1}^N \sum_{i, j=1}^m \partial_{\xi_i^\alpha} \partial_{\xi_j^\beta} F(\xi) \eta_i^\alpha \eta_j^\beta \geq \lambda(1 + |\xi|)^{q-2} |\eta|^2 > 0.$$

Observe that while F is strictly convex as a function of Xu , it may fail to be so when viewed as a function of ∇u (the full Riemannian gradient of u). The prototype of the functionals we consider is the regularized horizontal q -energy

$$(1.2) \quad E_q^\varepsilon(u, \Omega) = \int_{\Omega} (\varepsilon + |Xu|^2)^{\frac{q}{2}} dp,$$

with $2 \leq q < \infty$. Here, Ω denotes a given open subset of \mathbf{G} .

Another objective of this paper is to prove the partial regularity of weak solutions of quasi-linear systems of the form

$$(1.3) \quad - \sum_{\beta=1}^N \sum_{i, j=1}^m X_i (A_{ij}^{\alpha\beta}(x, u) X_j u^\beta) = f^\alpha(x, u, Xu), \alpha = 1, \dots, N$$

where $A_{ij}^{\alpha\beta}(x, u)$ are Hölder continuous functions satisfying together with $f = (f^1, \dots, f^N)$ the hypothesis: *For every $M > 0$ there exist $\lambda(M), a(M) > 0$, such that for $x \in \Omega$ and $|u| \leq M$ one has*

$$(1.4) \quad \sum_{\alpha, \beta=1}^N \sum_{i, j=1}^m A_{ij}^{\alpha\beta}(x, u) \xi_i^\alpha \xi_j^\beta \geq \lambda(M) |\xi|^2,$$

and

$$(1.5) \quad |f(x, u, \xi)| \leq a(M) |\xi|^2 + b.$$

The model for (1.3) is the sub-elliptic harmonic map system for mappings $u : \Omega \rightarrow N$, where N is a Riemannian manifold. In this case the principal part of the operator involved is the *sub-Laplacian* on \mathbf{G}

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

The Dirichlet problem for such maps with suitable conditions on the target has been recently studied by Jost and Xu [JX] (see also [HS]). The parabolic version of (1.3) has clearly a geometric interest in view of the applications to the harmonic flow. This case will be studied in a sub-sequent article.

The regularity of local minimizers of (1.1) and of weak solutions to (1.3) will be measured in terms of the horizontal Hölder and Sobolev spaces introduced by Folland and Stein [FS], [F]. Denote by H^s the s -dimensional Hausdorff measure generated by the pseudo-metric d . We prove the following.

Theorem 1.1 *Let \mathbf{G} be a Carnot group of step two, and let $\Omega \subset \mathbf{G}$ be a bounded open set. If u is a local minimizer of (1.1) in Ω , then there exists a closed set $S \subset \Omega$, having measure zero, such that Xu is Hölder continuous in $\Omega \setminus S$. Moreover, if $q = 2$ then $H^{Q-2+\sigma}(S) = 0$, for every $\sigma > 0$.*

Theorem 1.2 *Let \mathbf{G} be a Carnot group of step two, and let $\Omega \subset \mathbf{G}$ be a bounded open set. If u is a weak solution of (1.3) in Ω , with the assumptions (1.4), (1.5), and if $a(M)M \leq \lambda(M)/2$, then there exists a closed set $S \subset \Omega$ such that Xu is Hölder continuous in $\Omega \setminus S$, and $H^{Q-2+\sigma}(S) = 0$, for every $\sigma > 0$.*

In the classical elliptic case corresponding results have been proved by various authors using either the blow-up technique or the approach of Morrey-Campanato (see [Mo], [GM], [GG1], [GG2], [Gi], [E]). The central result, on which both approaches rest, are the $W^{2,2}$ -estimates for constant coefficient systems. When we tried to extend the classical theory to the non-elliptic setting of Carnot groups we were confronted with a major obstruction: The $W^{2,2}$ -estimates for the corresponding constant coefficient systems of Hörmander type were lacking. The main contribution of the present paper consists in developing a direct proof of such estimates in the setting of Carnot groups of step two. After completing this part we will suitably modify and adapt several parts of either of the above mentioned approaches to the non-abelian setting.

The constant coefficient systems of interest for us are of the form

$$(1.6) \quad - \sum_{\beta=1}^N \sum_{i,j=1}^m X_i (A_{ij}^{\alpha\beta} X_j u^\beta + f_i^\alpha) = f^\alpha, \quad \alpha = 1, \dots, N,$$

where $A_{ij}^{\alpha\beta}$ is a tensor satisfying the coercivity condition on the horizontal directions

$$(1.7) \quad \sum_{\alpha,\beta=1}^N \sum_{i,j=1}^m A_{ij}^{\alpha\beta} \xi_i^\alpha \xi_j^\beta \geq \lambda |\xi|^2, \quad \xi \in \mathbb{R}^{mN},$$

and f_i^α, f^α are L^2 functions defined in Ω . In the following we will let $f = (f^1, \dots, f^N)$, and $f_i = (f_i^1, \dots, f_i^N)$. Systems such as (1.6) are the most basic examples of non-diagonal systems of Hörmander type. In the context of the present paper they either arise as limits of the blow-up's of the Euler-Lagrange equation of (1.1), or as the "frozen" version of the homogeneous part of (1.3). We will prove that for smooth f_i^α , and f^α , the system (1.6) satisfies the following estimate: For each $n > 0$ there exist $C_n > 0$, and multi-indices a, b, c, d , depending on n , such that for any multi-index k of order less than or equal to n

$$(1.8) \quad \sup_{B(0,1)} |X^k u| \leq C_n \left\{ \int_{B(0,2)} |u|^2 dx + \int_{B(0,2)} [|T^a X^b f^\alpha|^2 + |T^c X^d f_i^\alpha|^2] dp \right\}.$$

Here, T denotes a derivative along the direction of commutators. As a consequence of (1.8) we find the following basic result.

Theorem 1.3 *Let \mathbf{G} be a Carnot group of step two, then the system (1.6) is hypoelliptic in the sense that, given an open set $\Omega \subset \mathbf{G}$ and $f^\alpha, f_i^\alpha \in C^\infty(\Omega)$, a weak solution to (1.6) is $C^\infty(\Omega)$.*

After one has the hypoellipticity result, then standard arguments lead to L^p , $1 < p < \infty$, and Hölder estimates for higher order derivatives of the solutions.

It is important to note that for diagonal systems the $W^{2,2}$ -estimates and the ensuing hypoellipticity follow from the fundamental works of J.J. Kohn [K], Hörmander [H], Folland and Stein [FS], and Folland [F]. In this connection, one should also see the interesting recent paper of Xu and Zuily [XZ]. However, the results in these papers do not cover non-diagonal systems and this is why we needed to develop a new proof of these basic results. After the completion of the present paper E. Stein, M. Christ and D. Geller brought to our attention that the $W^{2,2}$ -estimates could also be derived from the analysis of pseudo-differential operators on homogeneous groups developed in [G], [CGGP], and [T]. Such approach is very general but also highly technical and complex. We feel that our proof, which is close in spirit to the works [C1], [C2], has perhaps the advantage of being more easily accessible to workers in pde's who are not familiar with the theory of pseudo-differential operators. Moreover, the method of fractional order difference quotients used here can be applied also to quasi-linear systems such as (1.3), or to those arising as Euler-Lagrange equations for (1.1). While for (1.3) the parametrix method can still be used (by freezing the coefficients), this is not the case for the Euler-Lagrange systems of (1.1), where in order to use the method of freezing the coefficients one needs to have *a-priori* regularity assumptions on the minimizers.

We use the method of fractional order difference quotients to prove the second part of Theorem 1.1. In order to estimate the size of the singular set one needs to establish higher differentiability of the minimizers not just in terms of the horizontal derivatives, but with respect to the underlying Riemannian structure. This is achieved by the following result.

Theorem 1.4 *If u is as in Theorem 1.1, with $q = 2$, then its Riemannian gradient ∇u has locally square integrable horizontal derivatives.*

In closing, we give a short description of the organization of the paper. Section two is devoted to some preliminary definitions and lemmas concerning Sobolev and Hölder spaces. We also collect various basic results concerning difference quotients. In section three we establish the L^2 estimates, the hypoellipticity for (1.6), and prove Theorem 1.4. In section four we present for the special case $f^\alpha = 0$ the detailed blow-up analysis of (1.3), and we prove Theorem 1.1. In section five we prove Theorem 1.2.

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2. Preliminaries

As we mentioned in the previous section this paper focuses on Carnot groups of step two. We will thus henceforth assume that \mathbf{G} represents such a group with Lie algebra $\mathfrak{g} = V^1 \oplus V^2$, with $[V^1, V^1] = V^2$, $[\mathfrak{g}, V^2] = 0$. We will denote by m the dimension of V^1 , and indicate with k the dimension of V^2 . The homogeneous dimension of \mathbf{G} will be thus given by $Q = m + 2k$. Let X_1, \dots, X_m be an orthonormal basis of V^1 and fix an orthonormal basis T_1, \dots, T_k of V^2 . The same symbols X_j and T_l will be used to indicate the corresponding left-invariant vector fields on \mathbf{G} . One can easily see that the Carnot Carathéodory metric – the control distance associated to the vector fields X_j via Chow's accessibility theorem [Cho] – is equivalent to the pseudo-distance d . Since we will exclusively work with d , we will not make further reference to the control metric. If $B(p, r)$ denotes the pseudo-ball centered at p with radius r , then as we said $|B(p, r)| = \omega_{\mathbf{G}} r^Q$.

In this paper we will always use the globally defined exponential coordinates on \mathbf{G} . It is possible to describe the group structure only in terms of the algebra by means of the formula of Baker-Campbell-Hausdorff (see Sect. 2.15 in [V]): If $X, Y \in \mathfrak{g}$, then

$$(2.1) \quad e^X e^Y = e^{X+Y+\frac{1}{2}[X, Y]}.$$

Using (2.1) one can write explicitly the dependence of the vector fields on the exponential coordinates. There exists constants $c_{i,j,l}$, $i, j = 1, \dots, m$, and $l = 1, \dots, k$, such that if $(x, t) \in \mathbb{R}^m \oplus \mathbb{R}^k$ represent the exponential coordinates of $p \in \mathbf{G}$, then

$$X_i(p) = \frac{\partial}{\partial x_i} + \sum_{l=1}^k \sum_{j \neq i, j=1}^m c_{i,j,l} x_j \frac{\partial}{\partial t_l}, \quad i = 1, \dots, m,$$

$$T_l(p) = \frac{\partial}{\partial t_l}, \quad l = 1, \dots, k.$$

We respectively denote by $Xu = (X_1u, \dots, X_mu)$ and $Tu = (T_1u, \dots, T_ku)$ the *horizontal* and *non-horizontal* gradient of a function u .

The simplest non-trivial example of Carnot group of step two is the *Heisenberg group* \mathbb{H}^n . The latter is the Lie group whose underlying manifold is $\mathbb{R}^{2n} \times \mathbb{R}$ endowed with the non-commutative group law

$$(2.2) \quad (z, t)(\zeta, s) = (z + \zeta, t + s + 2[z, \zeta]),$$

where $z, \zeta \in \mathbb{R}^{2n}$, $t, s \in \mathbb{R}$, and $[z, \zeta] = \sum_{i=1}^n (\zeta_i z_{n+i} - z_i \zeta_{n+i})$. The vector fields $X_i = \partial_{z_i} + 2z_{n+i} \partial_t$, $X_{i+n} = \partial_{z_{n+i}} - 2z_i \partial_t$, for $i = 1, \dots, n$, and $T = -4 \partial_t$ form a left-invariant basis for the Lie algebra of \mathbb{H}^n . The real part of the Kohn sub-Laplacian on \mathbb{H}^n is $\mathcal{L} = \sum_{j=1}^{2n} X_j^2$. The stratification of the Heisenberg algebra is given by $V^1 = \text{span}\{X_1, \dots, X_{2n}\}$, and $V^2 = \text{span}\{T\}$. The associated dilations are parabolic $\delta_\lambda(z, t) = (\lambda z, \lambda^2 t)$, with homogeneous dimension $Q = 2n + 2$. If V is a left-invariant vector field and $V = \sum_{l=1}^{2n} v_l X_l + v_{2n+1} T$, then we will write the exponential map as $\exp(V) = e^V = (v_1, v_2, \dots, v_{2n+1}) = (v', v_{2n+1})$.

In this case, the Baker-Campbell-Hausdorff formula follows immediately from the group law (2.2). Since the group is of step two the homogeneous gauge is given by $\|p\| = (|z|^4 + t^2)^{1/4}$, where $p = (z, t) \in \mathbb{H}^n$, with relative pseudo-distance $d(p, p') = \|p^{-1}p'\|$. It was proved by Cygan [Cy] that this is in fact a distance.

Let \mathbf{G} be a Carnot group of step two and $\Omega \subset \mathbf{G}$ an open subset. For $N \in \mathbb{N}$ consider functions $u : \Omega \rightarrow \mathbb{R}^N$. For $0 < \alpha < 1$ we will use the Folland-Stein class

$$\Gamma^\alpha(\Omega) = \left\{ u \in L^\infty(\Omega) \left| \sup_{p_1, p_2 \in \Omega, p_1 \neq p_2} \frac{|u(p_1) - u(p_2)|}{d(p_1, p_2)^\alpha} < \infty \right. \right\},$$

and its local version $\Gamma_{\text{loc}}^\alpha(\Omega) = \{u | \eta u \in \Gamma^\alpha(\Omega) \text{ for every } \eta \in C_0^\infty(\Omega)\}$, see [FS] and [F]. For $k \in \mathbb{N}$ the symbol $\Gamma_{\text{loc}}^{k, \alpha}(\Omega)$ will denote the set of functions having horizontal derivatives up to order k in $\Gamma_{\text{loc}}^\alpha(\Omega)$. Observe that $\Gamma^\alpha(\Omega)$ becomes a Banach space when endowed with the norm

$$\|u\|_{\Gamma^\alpha(\Omega)} = \sup_{p_1, p_2 \in \Omega, p_1 \neq p_2} \frac{|u(p_1) - u(p_2)|}{d(p_1, p_2)^\alpha} + \|u\|_{L^\infty(\Omega)}.$$

We will need the horizontal Sobolev spaces

$$\mathcal{L}_{\text{loc}}^{1, q}(\Omega) = \left\{ f : \Omega \rightarrow \mathbb{R}^N \mid f \in L_{\text{loc}}^q(\Omega), Xf = (X_1 f, \dots, X_m f) \in L_{\text{loc}}^q(\Omega) \right\},$$

$$1 \leq q < \infty.$$

Here, $X_j f$, $j = 1, \dots, m$ denotes the weak derivative of $f \in L_{\text{loc}}^1(\Omega)$ along X_j defined by the formula

$$\int_{\Omega} X_j f(p) \phi(p) dp = - \int_{\Omega} f(p) X_j \phi(p) dp, \quad \phi \in C_0^\infty(\Omega).$$

We recall that in any Carnot group one has $X_j^* = -X_j$ for the formal adjoint of X_j , see [F]. We will also use the Sobolev space of functions with zero trace, $\mathcal{L}_o^{1, q}(\Omega) = \overline{C_o^1(\Omega)}$, where the closure is taken in the $\mathcal{L}^{1, q}$ -norm, $\|u\|_{\mathcal{L}^{1, q}(\Omega)}^q = \int_{\Omega} (|u|^q + |Xu|^q) dx$. Let h be an integer and denote by $I_h = (i_1, \dots, i_h)$ a multi-index with $1 \leq i_j \leq m$. Given a function f we indicate with $X^{I_h} f = X_{i_1} \dots X_{i_h} f$ the corresponding horizontal derivative of order h . The higher order Sobolev spaces are

$$\mathcal{L}_{\text{loc}}^{k, q}(\Omega) = \left\{ f : \Omega \rightarrow \mathbb{R}^N \mid \sum_{I_h} |X^{I_h} f| \in L_{\text{loc}}^q(\Omega), \text{ for any } 0 \leq h \leq k \right\},$$

$$1 \leq q < \infty.$$

We next recall that in a group of step two the gauge pseudo-balls are *NTA* domains, see [CG], hence, in particular, they are Poincaré-Sobolev (*PS*-) domains. Thereby, such sets support the Poincaré-Sobolev inequality, i.e., for every

$1 \leq q < Q$ there exists $S_q > 0$ such that for every $u \in S^{1,q}(B(p_o, r))$ one has

$$\left(\frac{1}{|B(p_o, r)|} \int_{B(p_o, r)} |u - u_{p_o, r}|^{\frac{Qq}{Q-q}} dx \right)^{\frac{Q-q}{Qq}} \leq S_q r \left(\frac{1}{|B(p_o, r)|} \int_{B(p_o, r)} |Xu|^q dx \right)^{\frac{1}{q}},$$

where we have let $u_{p_o, r}$ denote the average of u over the pseudo-ball $B(p_o, r)$. See, e.g., [Lu], [GN] or also [HK]. The fact that gauge pseudo-balls are (PS)-domains allows to use Theorem 1.27 in [GN] and obtain the following compact embedding

$$\mathcal{L}^{1,q}(B(p_o, r)) \hookrightarrow L^q(B(p_o, r)), \quad 1 \leq q < \infty.$$

We will need the analogue of the Morrey-Campanato embedding, proved by Folland (see Theorem 5.15 [F]): *For every $p_o \in \mathbf{G}$, $k \in \mathbb{N}$, $1 \leq q < \infty$, and $r > 0$, one has*

$$\mathcal{L}^{k,q}(B(p_o, r)) \subset \Gamma_{\text{loc}}^\alpha(B(p_o, r)),$$

where $\alpha = k - \frac{Q}{q}$.

The spaces $\mathcal{L}_{\text{loc}}^{1,q}(\Omega)$, with $q < \infty$, can be characterized in terms of difference quotients. Let Z be a smooth vector field in Ω . A function $f \in L_{\text{loc}}^1(\Omega)$ is weakly differentiable in Ω along the direction Z , if there exists a locally integrable function g such that for any $\phi \in C_0^\infty(\Omega)$ the identity

$$\int_{\Omega} f(p) Z^* \phi(p) dp = \int_{\Omega} g(p) \phi(p) dp$$

holds. Here Z^* denotes the formal adjoint of Z . For $0 < s < \varepsilon$, where ε is suitably small, denote with pe^{sZ} the integral curve of Z starting at p . If $K \subset\subset \Omega$ is an open set, then we can think of e^{sZ} as a diffeomorphism $e^{sZ} : K \rightarrow e^{sZ}(K) \subset \Omega$. We also let $e^{sZ} f(p) = f(pe^{sZ})$ when $p \in K$ and f is a real valued function defined on $e^{sZ}(K)$.

Proposition 2.1 *Let $u \in L_{\text{loc}}^q(\Omega)$, $1 \leq q < \infty$, and $K \subset\subset \Omega$. If there exist positive constants ε, C such that*

$$(2.3) \quad \sup_{0 < |s| < \varepsilon} \int_K \frac{|u(xe^{sZ}) - u(x)|^q}{|s|^q} dx \leq C^q,$$

then $Zu \in L^q(K)$ and $\|Zu\|_{L^q(K)} \leq C$. On the other hand, if $Zu \in L_{\text{loc}}^q(\Omega)$, then (2.3) holds with $C = 2\|Zu\|_{L^q(K)}$ for any $K \subset\subset \Omega$.

Proof. The proof of this proposition is standard, see for instance Proposition 2.3 in [C1]. \square

In the following we will need estimates on the Lebesgue norm of fractional derivatives of functions in the direction of commutators T . Let Ω, Z be as in Proposition 2.1. For $w \in L^2(\Omega)$ with compact support in Ω , and $\alpha \in (0, 1)$, we define the seminorm

$$(2.4) \quad |w|_{Z, \alpha}^2 = \sup_{|h| < \varepsilon_0} \int_{\Omega} |h|^{-2\alpha} |w(qe^{hZ}) - w(q)|^2 dq,$$

where ε_0 is chosen sufficiently small. Let $n = m + k = \dim(\mathbf{g})$, and identify \mathbf{g} with the Euclidean space \mathbb{R}^n . Now consider a function $w \in C_0^\infty(\mathbb{R}^n)$. For $\alpha \in (0, 1)$ the fractional derivative of w of order α , along the direction e_j , $j = 1, \dots, n$, is given by

$$\frac{\widehat{\partial^\alpha w}}{\partial e_j} = |h|^\alpha \widehat{w}(x_1, \dots, x_{j-1}, h, x_{j+1}, \dots, x_n),$$

where we have denoted by \widehat{w} the partial Fourier transform in the variable e_j . By the Plancherel formula we have

$$(2.5) \quad \left\| \frac{\partial^\alpha w}{\partial e_j} \right\|_{L^2(\mathbf{g})} = \int_{\mathbb{R}^n} |h|^{2\alpha} |\widehat{w}(x_1, \dots, x_{j-1}, h, x_{j+1}, \dots, x_n)|^2 dx dh,$$

The following theorem of Peetre [P] (see also sec. 2.3.3, Remark 4, p. 180 in [Tr]) will play an important role in the sequel.

Theorem 2.2 *Let $0 < \beta < \alpha < 1$, and $w \in C_0^\infty(\mathbb{R}^n)$. There exists a positive constant C , depending only on α and β , such that for each vector e_j , $j = 1, \dots, n$, of the standard basis of \mathbb{R}^n , one has*

$$C \left\| \frac{\partial^\beta w}{\partial e_j} \right\|_{L^2(\mathbb{R}^n)} \leq |w|_{\frac{\partial}{\partial e_j}, \alpha} \leq C^{-1} \left\| \frac{\partial^\alpha w}{\partial e_j} \right\|_{L^2(\mathbb{R}^n)}.$$

The next result, which is a direct consequence of the formula of Baker-Campbell-Hausdorff, is a special case of Theorem 4.3 in Hörmander's fundamental paper [H].

Theorem 2.3 *Let \mathbf{G} be a Carnot group of step two, let $\Omega \subset \mathbf{G}$ be an open set and $B(p_o, r) \subset \Omega$ a gauge pseudo-ball. If $\eta \in C_0^\infty(B(p_o, r))$, and $w \in C^\infty(\Omega)$, then there exists a positive constant C such that for each $T \in V^2$.*

$$|w\eta|_{T, \frac{1}{2}} \leq C \sum_{l=1}^{2n} |w\eta|_{X_l, 1}.$$

At this point we turn our attention to weak solutions to (1.6). To state our next result we introduce a definition. For a given left-invariant vector field Z and a function $f : \mathbf{G} \rightarrow \mathbb{R}$ we let for $s \neq 0$ fixed

$$(f)_{Z, \gamma}(p) = \frac{f(pe^{sZ}) - f(p)}{|s|^\gamma}.$$

Theorem 2.4 *Let \mathbf{G} be a Carnot group of step two and $\Omega \subset \mathbf{G}$ be an open set. If $u \in \mathcal{L}_{\text{loc}}^{1,2}(\Omega)$ is a weak solution of (1.6) in Ω , then for every left-invariant vector field Z , the function $(u)_{Z,\gamma} = ((u^1)_{Z,\gamma}, \dots, (u^N)_{Z,\gamma})$ satisfies the following system*

$$(2.6) \quad \begin{aligned} & \sum_{\beta=1}^N \sum_{i,j=1}^m \int_{\Omega} \left[A_{ij}^{\alpha\beta} (X_j u^\beta)_{Z,\gamma}(p) + (f_i^\alpha)_{Z,\gamma}(p) \right] X_i \phi^\alpha(p) dp \\ & + \sum_{\beta=1}^N \sum_{i,j=1}^m \int_{\Omega} \left[A_{ij}^{\alpha\beta} X_j u^\beta(pe^{sZ}) + f_i^\alpha(pe^{sZ}) \right] \frac{s|s|^{-\gamma}}{2} [Z, X_i] \phi^\alpha(p) dp \\ & = \int_{\Omega} f^\alpha(p) (\phi^\alpha)_{-Z,\gamma}(p) dp, \quad \alpha = 1, \dots, N, \end{aligned}$$

for every $\phi = (\phi^1, \dots, \phi^N) \in C_o^\infty(\Omega, \mathbb{R}^N)$, $\gamma \in (0, 1)$, and $|s|$ small enough.

Proof. The proof is based on the Baker-Campbell-Hausdorff formula. The first step consists in writing the weak form of (1.6)

$$\sum_{\beta=1}^N \sum_{i,j=1}^m \int_{\Omega} \left(A_{ij}^{\alpha\beta} X_j u^\beta + f_i^\alpha \right) (p) X_i \phi^\alpha(p) dp = \int_{\Omega} f^\alpha(p) \phi^\alpha(p) dp,$$

where $\phi = (\phi^1, \dots, \phi^N) \in C_o^\infty(\Omega, \mathbb{R}^N)$. We also write a similar equation with test function $\tilde{\phi}(p) = \phi(pe^{-sZ})$, subtract one of the other, divide by $|s|^\gamma$ and change variable $p \rightarrow pe^{sZ}$, to obtain

$$\begin{aligned} & \sum_{\beta=1}^N \sum_{i,j=1}^m \int_{\Omega} \left[A_{ij}^{\alpha\beta} (X_j u^\beta)_{Z,\gamma}(p) + (f_i^\alpha)_{Z,\gamma}(p) \right] X_i \phi^\alpha(p) dp \\ & + \sum_{\beta=1}^N \sum_{i,j=1}^m \int_{\Omega} \left[A_{ij}^{\alpha\beta} X_j u^\beta(pe^{sZ}) + f_i^\alpha(pe^{sZ}) \right] \\ & \quad \times \frac{X_i[\phi^\alpha(pe^{-sZ})](pe^{sZ}) - X_i \phi^\alpha(p)}{|s|^\gamma} dp \\ & = \int_{\Omega} f^\alpha(p) (\phi^\alpha)_{-Z,\gamma}(p) dp, \end{aligned}$$

(recall that the Jacobian of the translation $p \rightarrow pe^{sZ}$ is the constant function with value one, see for instance [CoG]). The Baker-Campbell-Hausdorff formula yields

$$\frac{X_i[\phi^\alpha(pe^{-sZ})](pe^{sZ}) - X_i \phi^\alpha(p)}{|s|^\gamma} = \frac{s|s|^{-\gamma}}{2} [Z, X_i] \phi^\alpha(p),$$

concluding the proof. \square

In several occasions we will make use of a basic energy estimate, which in the standard elliptic case is usually referred to as Caccioppoli inequality. In the next lemma we will state such estimate for weak solutions of the system (1.3) (of which

the constant coefficient system (1.6) is a special case), but it should be remarked here that, as in the classical setting (see for instance [Gi]), it holds in much greater generality.

Lemma 2.5 *Let \mathbf{G} be a Carnot group of step two, and let $\Omega \subset \mathbf{G}$ be a bounded open set. If $u \in \mathcal{L}_{\text{loc}}^{1,2}(\Omega)$ is a weak solution of (1.3) in Ω , with the assumptions (1.4), (1.5) and with $f^\alpha = 0$, $\alpha = 1, \dots, N$, then the following Caccioppoli type inequality holds for $B(p_o, 2r) \subset \Omega$*

$$(2.7) \quad \int_{B(p_o, r)} |Xu|^2 dp \leq C \frac{1}{r^2} \int_{B(p_o, 2r)} |u|^2 dp.$$

Proof. The proof follows along the lines of the classical argument, see, e.g., [Gi]. One multiplies both sides of (1.3) by $u\eta^2$, where $\eta \in C_o^\infty(B(p_o, 2r))$, $\eta \geq 0$, $\eta \equiv 1$ in $B(p_o, r)$, and $|X\eta| \leq C/r$. After integrating on $B(p_o, 2r)$, and using Hölder inequality and the coercivity assumptions (1.4), the conclusion follows. The existence of test functions η with the properties described above can be easily established using the group translations and the pseudo-norm $\|\cdot\|$. \square

3. Systems with constant coefficients

In this section we prove the hypoellipticity of the constant coefficient system (1.6)

$$-\sum_{\beta=1}^N \sum_{i,j=1}^m X_i \left(A_{ij}^{\alpha\beta} X_j u^\beta + f_i^\alpha \right) = f^\alpha, \alpha = 1, \dots, N.$$

Our approach is based on appropriate sub-elliptic versions, contained in Theorems 3.5 and 3.6, of the classical $W^{2,2}$ estimates. Following the classical scheme we differentiate the system with respect to a vector field Z , and show that the function Zu solves a new system, similar in structure to (1.6). A major obstacle is constituted by the non-trivial commutation relations involving the horizontal basis $X = \{X_1, \dots, X_m\}$. In fact, when (1.6) is differentiated new terms arise which include derivatives of u along the directions of the commutators $[X_i, Z]$. Since a-priori there is no control on such terms, it is necessary to develop new ideas. In essence, we combine in a delicate way a repeated use of the fractional difference quotients along the integral curves of Z , with the energy inequality in Lemma 2.5, and the interpolation estimates in Theorems 2.2 and 2.3. In the end, the iteration of the sub-elliptic $W^{2,2}$ estimates allows to establish the smoothness of weak solutions.

To conclude the section, in Theorem 3.9 we will prove a priori estimates for quasi-linear systems with quadratic growth in the horizontal gradient. Such estimates, which are derived by the same technique employed in Theorems 3.1 and 3.2, will be used in the proof of Theorem 4.11. Henceforth, we use the notation $f_i = (f_i^1, \dots, f_i^N)$, $i = 1, \dots, m$, $f = (f^1, \dots, f^N)$.

Theorem 3.1 *Let \mathbf{G} be a Carnot group of step two, and $u \in \mathcal{L}_{\text{loc}}^{1,2}(\Omega)$ be a weak solution to (1.6) in an open set $\Omega \subset \mathbf{G}$. If $f^\alpha, T_j f^\alpha, f_i^\alpha, T_j f_i^\alpha \in L_{\text{loc}}^2(\Omega)$, for every $i = 1, \dots, m, j = 1, \dots, k$, and $\alpha = 1, \dots, N$, then for each $T_{j_o} \in V^2$ one has*

$$(i) \quad T_{j_o} u \in \mathcal{L}_{\text{loc}}^{1,2}(\Omega).$$

Furthermore, for every pseudo-ball $B(p_o, 2r) \subset \Omega$ and $\eta \in C_o^\infty(B(p_o, 2r))$, the following estimate holds for $w = \eta T_{j_o} u$

$$(ii) \quad \int_{B(p_o, r)} \left(r^4 |w|^2 + r^6 |Xw|^2 \right) dp \\ \leq C \int_{B(p_o, 2r)} \left(r^2 |Xu|^2 + |u|^2 + r^4 |f|^2 + r^8 |T_{j_o} f|^2 \right. \\ \left. + \sum_{i=1}^m \left[r^2 |f_i|^2 + r^6 |T_{j_o} f_i|^2 \right] \right) dp,$$

for some positive constant C depending only on \mathbf{g} and on the structural parameter in (1.7). Moreover, $T_{j_o} u$ is a weak solution in Ω to

$$(iii) \quad - \sum_{\beta=1}^N \sum_{i,j=1}^m X_i \left(A_{i,j}^{\alpha\beta} X_j (T_{j_o} u)^\beta + T_{j_o} f_i^\alpha \right) = T_{j_o} f^\alpha,$$

for $\alpha = 1, \dots, N$.

Proof. Since the class of systems satisfying (1.7) is invariant under the action of the group left-translations it is enough to show that $w \in \mathcal{L}_{\text{loc}}^{1,2}(B(e, 4r))$, where e is the group identity and $B(e, 4r) \subset \Omega$. By dilation invariance we can assume $r = 1$. Once (ii) has been proved for $r = 1$, the general case follows by rescaling. We use difference quotients in the T_{j_o} direction of order (roughly) one half. For $0 < \gamma < 1$ define

$$u_{(T_{j_o}, \gamma)}^\beta(p) = |\sigma|^{-2\gamma} \left(u^\beta(pe^{\sigma^2 T_{j_o}}) - u^\beta(p) \right).$$

The differentiated system (2.6) takes a simpler form in the case $Z = T_{j_o}$ and $\gamma = 1/2$. For every $\phi \in C_o^\infty(\Omega, \mathbb{R}^N)$ one has

$$(3.1) \quad \sum_{\beta=1}^N \sum_{i,j=1}^m \int_{\Omega} \left[A_{i,j}^{\alpha\beta} (X_j u^\beta)_{(T_{j_o}, 1/2)} + (f_i^\alpha)_{(T_{j_o}, 1/2)} \right] X_i \phi^\alpha dp = \int_{\Omega} f^\alpha \phi_{(-T_{j_o}, 1/2)}^\alpha dp.$$

Now observe that (3.1) continues to hold for test functions in $\mathcal{L}_o^{1,2}(\Omega, \mathbb{R}^N)$. Since T_{j_o} commutes with the X_j 's, then the function $p \rightarrow u(pe^{\sigma^2 T_{j_o}})$ is in $\mathcal{L}^{1,2}(B(e, 4))$. If we set $\phi^\alpha(p) = (u^\alpha)_{(T_{j_o}, 1/2)}(p) \eta^2(p)$, with η as in the statement of the theorem,

then ϕ is a good test function for (3.1). From the coercivity hypothesis, and with this choice of test function, we obtain from (3.1) the following energy inequality

$$(3.2) \quad \int_{\Omega} |(Xu^\alpha)_{(T_{j_0}, 1/2)}|^2 \eta^2 dp \\ \leq C \int_{\Omega} \left[|(u^\alpha)_{(T_{j_0}, 1/2)}|^2 (\eta^2 + |X\eta|^2) \right. \\ \left. + \sum_{i=1}^m |(f_i)_{(T_{j_0}, 1/2)}|^2 \eta^2 + |(f)_{(T_{j_0}, 1/2)}|^2 \eta \right] dp.$$

Noting that $|\eta_{(T_{j_0}, 1/2)}| \leq C$, then (3.2) and Theorems 2.2 and 2.3 yield

$$(3.3) \quad \int_{\Omega} |\partial_{t_{j_0}}^{\frac{1}{2}-\varepsilon} X(\eta u^\alpha)|^2 dp \\ \leq C \int_{B(e, 2)} \left(|u|^2 + |Xu|^2 + \sum_{i=1}^m [|f_i|^2 + |T_{j_0} f_i|^2] + |f|^2 + |T_{j_0} f|^2 \right) dp.$$

To improve in the order of differentiability of u we use Theorem 2.3 once more

$$(3.4) \quad \int_{\Omega} |\partial_{t_{j_0}}^{1-2\varepsilon} \eta u^\alpha|^2 dp \\ \leq C \int_{B(e, 2)} \left(|u|^2 + |Xu|^2 + \sum_{i=1}^m [|f_i|^2 + |T_{j_0} f_i|^2] + |f|^2 + |T_{j_0} f|^2 \right) dp.$$

We now apply (3.4) with $\varepsilon_0 = 1/12$ to gain control over the L^2 norm of $\partial_{t_{j_0}}^{\frac{5}{6}} u^\alpha$. Next, we consider the difference quotients $u_{(T_{j_0}, 5/6)}$, repeat the argument in (3.1)–(3.3), and use (3.4) to obtain

$$\int_{\Omega} \left(|\partial_{t_{j_0}}^{\frac{5}{6}-\varepsilon} (\eta u^\alpha)|^2 + |X \partial_{t_{j_0}}^{\frac{5}{6}-\varepsilon} (\eta u^\alpha)|^2 \right) dp \\ \leq C \int_{B(e, 2)} \left(|u|^2 + |Xu|^2 + \sum_{i=1}^m [|f_i|^2 + |T_{j_0} f_i|^2] + |f|^2 + |T_{j_0} f|^2 \right) dp.$$

At this point we are ready to estimate

$$(3.5) \quad \int_{\Omega} |T(\eta u^\alpha)|^2 dp \leq \int_{\Omega} |\partial_{t_{j_0}}^{\frac{1}{2}-\varepsilon_0} \partial_{t_{j_0}}^{\frac{5}{6}-\varepsilon_0} (\eta u^\alpha)|^2 dp \leq \int_{\Omega} |X \partial_{t_{j_0}}^{\frac{5}{6}-\varepsilon_0} (\eta u^\alpha)|^2 dp \\ \leq \int_{B(e, 2)} \left(|u|^2 + |Xu|^2 + \sum_{i=1}^m [|f_i|^2 + |T_{j_0} f_i|^2] + |f|^2 + |T_{j_0} f|^2 \right) dp.$$

This proves the first part of the theorem. To establish (ii) it now suffices to consider the following difference quotients of order one: $(u^\alpha)_{(T_{j_0}, 1)}$, and apply Lemma 2.5 to (2.6) for this special case. Part (iii) follows from part (i) and from (2.6). \square

In the sequel for a function f defined in \mathbf{G} we indicate with ∇f the standard Riemannian gradient of f , whereas we let $|Tf|^2 = (\sum_{i=1}^k (T_i f)^2)$.

Theorem 3.2 *Let \mathbf{G} , u , and Ω be as in Theorem 3.1. If $f^\alpha, T_j f^\alpha \in L_{\text{loc}}^2(\Omega)$ and $f_i^\alpha, \nabla f_i^\alpha \in L_{\text{loc}}^2(\Omega)$, for $i = 1, \dots, m, j = 1, \dots, k$ and $\alpha = 1, \dots, N$, then for every $i_o = 1, \dots, m$ one has*

$$(i) \quad X_{i_o} u \in \mathcal{L}_{\text{loc}}^{1,2}(\Omega).$$

Moreover, for any pseudo-ball $B(p_o, 2r) \subset \Omega$ and $\eta \in C_0^\infty(B(p_o, 2r))$ one has for $v = \eta X_{i_o} u$

$$(ii) \quad \int_{B(p_o, r)} r^4 |Xv|^2 dp \\ \leq C \int_{B(p_o, 2r)} \left(r^2 |Xu|^2 + |u|^2 \right. \\ \left. + \sum_{i=1}^m \left[r^2 |f_i|^2 + r^4 |Xf_i|^2 + r^6 |Tf_i|^2 \right] + r^4 |f|^2 + r^8 |Tf|^2 \right) dp,$$

for some positive constant C depending only on \mathbf{G} and the structural assumptions in (1.6). Finally, if we also assume that $f^\alpha \in \mathcal{L}_{\text{loc}}^{1,2}(\Omega)$, then $X_{i_o} u$ is a weak solution in Ω to

$$(iii) \quad \sum_{\beta=1}^N \sum_{i,j=1}^m X_i \left(A_{i,j}^{\alpha\beta} X_j X_{i_o} u^\beta + A_{i,j}^{\alpha\beta} [X_{i_o}, X_j] u^\beta + X_{i_o} f_i^\alpha \right) \\ + \sum_{\beta=1}^N \sum_{i,j=1}^m [X_{i_o}, X_i] \left(A_{i,j}^{\alpha\beta} X_j u^\beta + f_i^\alpha \right) = X_{i_o} f^\alpha, \quad \alpha = 1, \dots, N.$$

We mention explicitly that an important consequence of Theorems 3.1 and 3.2 is that

$$X_i u \in W_{\text{loc}}^{1,2}(\Omega),$$

where $W_{\text{loc}}^{1,2}(\Omega)$ denotes the ordinary Sobolev space.

Before proving Theorem 3.2 we need a preliminary lemma.

Lemma 3.3 *Let \mathbf{G} be a Carnot group of step two and $\Omega \subset \mathbf{G}$ be an open set. If $f \in \mathcal{L}_{\text{loc}}^{1,2}(\Omega)$ is such that $T_j f \in L_{\text{loc}}^2(\Omega)$, for $j = 1, \dots, k$, then for each $\theta \in \mathbf{G}$, the function $g(p) = f(p\theta)$ is in $\mathcal{L}_{\text{loc}}^{1,2}(\tilde{\Omega})$, where $\tilde{\Omega} = \Omega\theta^{-1}$. Moreover, one has the identity:*

$$X_i (f(pe^S)) = X_i f(pe^S) - [S, X_i] f(pe^S),$$

where without loss of generality we have set $\theta = e^S$.

Proof. In the Heisenberg group the proof is a simple application of the Baker-Campbell-Hausdorff formula. For the general case the argument is slightly different and is based on the following identity (see Chap. 1 in [CoG], or Sect. 2.5 in [V]): For $X, Y \in \mathfrak{g}$

$$(\text{Ad } \exp X)Y = \sum_{k=0}^{\infty} \frac{1}{k!} (\text{ad } X)^k Y = e^{\text{ad } X} Y,$$

where, if $P \in \mathbf{G}$, then $\text{Ad } P : \mathfrak{g} \rightarrow \mathfrak{g}$ is the differential of the map $Q \rightarrow PQP^{-1}$, and $(\text{ad } X)Y = [X, Y]$. In a nilpotent Lie group the infinite sum will simply be a finite sum, the highest order of non-zero commutators being the step of the group (in our setting hence, a sum of two terms).

At this point we observe that

$$\begin{aligned} X_i(f(pe^S)) &= \frac{d}{dt} f(pe^{tX_i} e^S)|_{t=0} = \frac{d}{dt} f(pe^S(e^{-S})e^{tX_i}e^S)|_{t=0} \\ &= \left(((\text{Ad } e^{-S})X_i)f \right)(pe^S) = \left(\sum_{k=0}^{\text{Step of } \mathbf{G}} (-1)^k \frac{(\text{ad } S)^k X_i}{k!} f \right)(pe^S). \end{aligned}$$

In our case, with a group of step two the formula becomes

$$X_i(f(pe^S)) = X_i f(pe^S) - [S, X_i]f(pe^S).$$

The proof is concluded once we observe that $[S, X_i] \in V^2$ and the vectors T_k are a basis of V^2 . \square

Thanks to Theorem 3.1 (i), Lemma 3.3 allows to use difference quotients in the direction of X_i , $i = 1, \dots, m$.

Remark 3.4 Using the argument in Lemma 3.3 one can easily prove

$$X_i(u_{(X_{i_0}, 1)}) = (X_i u)_{(X_{i_0}, 1)} + [X_i, X_{i_0}]u(pe^{\sigma X_{i_0}}).$$

In particular, this shows that if $u \in W_{\text{loc}}^{1,2}(\Omega)$, then $u_{(X_{i_0}, 1)} \in \mathcal{L}_{\text{loc}}^{1,2}(\Omega)$.

Proof of Theorem 3.2. The proof is based on Theorem 3.1. In particular, what makes things work is the fact that the solution u is not only in the *horizontal* Sobolev space $\mathcal{L}_{\text{loc}}^{1,2}(\Omega)$, but also in the usual Euclidean Sobolev space $W_{\text{loc}}^{1,2}(\Omega)$. This follows from Theorem 3.1, which asserts the L^2 -integrability of the derivatives in the directions of the second layer V^2 , and from the representation of the vector fields in exponential coordinates (see section two).

As in the proof of Theorem 3.1, it suffices to establish (ii) in a fixed pseudo-ball $B(e, 1)$, and then use left-translations and dilations to obtain (ii) for a general pseudo-ball $B(p_\sigma, r)$. Observe that by continuity (2.6) still holds for any ϕ in the ordinary Sobolev space $W_o^{1,2}(\Omega, \mathbb{R}^N)$. Let η be as in the statement of the theorem. By Remark 3.4 we may use $\phi = u_{(X_{i_0}, 1)}\eta^2$ as a test function in (2.6), thus we

proceed to substitute it in the equation. By Remark 3.4 and the coercivity condition (1.7) we derive the energy inequality

$$(3.6) \quad \int_{\Omega} |(X\eta u^\alpha)_{(X_{i_0}, 1)}|^2 dp \leq C \int_{B(e, 2)} \left(|u|^2 + |Xu|^2 + |T(u)_{(X_{i_0}, 1)}|^2 + \sum_{i=1}^m [|Xf_i|^2 + |f_i|^2] + |f^\alpha| |(u^\alpha)_{(X_{i_0}, 1)} \eta^2|_{(X_{i_0}, 1)} \right) dp.$$

The conclusion now follows using Theorem 3.1 and Proposition 2.1. \square

When studying the regularity theory of weak solutions to the quasi-linear system (1.3) one is naturally led to analyze the hypoellipticity of linear sub-elliptic systems with constant coefficients of the type (1.6). In doing this one naturally encounters systems such as (2.6). For $A_{ij}^{\alpha\beta}$ as in (1.7) and f^α a smooth function we thus consider weak solutions of

$$(3.7) \quad \sum_{\beta=1}^N \sum_{i,j=1}^m A_{ij}^{\alpha\beta} X_i X_j u^\beta = f^\alpha,$$

for $\alpha = 1, \dots, N$. By iterating Theorems 3.1 and 3.2 we obtain the following.

Theorem 3.5 *Let \mathbf{G} be a Carnot group of step two and $\Omega \subset \mathbf{G}$ be open. Let $u \in \mathcal{L}_{\text{loc}}^{1,2}(\Omega)$ be a weak solution of (3.7) in Ω . For each $l, h \in \mathbb{N}$ and multi-indices $I_h = (i_1, \dots, i_h) \in \{1, 2, \dots, m\}^h$, and $J_l = (j_1, \dots, j_l) \in \{1, 2, \dots, k\}^l$, define*

$$v_{l,h}^\beta = T^{J_l} X^{I_h} u^\beta = (T_{j_1} \dots T_{j_l})(X_{i_1} \dots X_{i_h}) u^\beta.$$

If $f = (f^1, \dots, f^N) \in C^\infty(\Omega, \mathbb{R}^N)$ then

$$(3.8) \quad v_{l,h} \in \mathcal{L}_{\text{loc}}^{1,2}(\Omega).$$

Moreover, $v_{l,h}$ is a solution to the system

$$(3.9) \quad \sum_{\beta=1}^N \sum_{i,j=1}^m X_i \left(A_{ij}^{\alpha\beta} X_j v_{l,h} + f_{i,(l,h)}^\alpha \right) = f_{(l,h)}^\alpha, \quad \alpha = 1, \dots, N,$$

where $f_{(l,0)}^\alpha = T^{J_l} f^\alpha$, $f_{i,(l,0)}^\alpha = 0$,

$$f_{i,(l,h)}^\alpha = T_{j_l} X_{i_h} f_{i,(l-1,h-1)}^\alpha + A_{ij}^{\alpha\beta} [X_{i_h}, X_j] T_{j_l} v_{l-1,h-1}^\beta,$$

and

$$f_{l,h}^\alpha = T_{j_l} \left\{ [X_{i_h}, X_{i_l}] \left(A_{ij}^{\alpha\beta} X_j v_{l-1,h-1}^\beta + f_{i,(l-1,h-1)}^\alpha \right) + X_{i_h} f_{(l-1,h-1)}^\alpha \right\}.$$

We define the semi-norm $\|g\|_{l,h,R} = \sum_{I_h, J_l} \|T^{J_l} X^{I_h} g\|_{L^2(B(p_o, r))}$, where I_h and J_l respectively range over all the elements of $\{1, 2, \dots, m\}^h$ and of $\{1, 2, \dots, k\}^l$. As a consequence of Theorems 3.1, 3.2 and 3.5 one obtains the following.

Theorem 3.6 *Let G , Ω , $B(p_o, r)$, and u be as in the previous theorem. For each $h, l \in \mathbb{N}$, one has*

$$(3.10) \quad r^{2(2l+2h+2)} \|u\|_{l,2h+2,r}^2 \leq C_{h,l} \left\{ \|u\|_{0,0,2r}^2 + r^{2(2+2(l+3h+1)+2h)} \|f\|_{l+3h+1,2h,2r}^2 \right\}.$$

Proof. By means of left-translations and rescalings we can assume that $p_o = e$, and $r = 1$. Using an induction argument and Theorem 3.5 one deduces that

$$(3.11) \quad \|f_{i,(l,h)}^\alpha\|_{\gamma,v,1} \leq C \|u\|_{l+1+\gamma, h-1+v,2},$$

and

$$(3.12) \quad \|f_{(l,h)}^\alpha\|_{\gamma,v,1} \leq C \left(\|f^\alpha\|_{l+\gamma, h+v,2} + \|u\|_{l+1+\gamma, h+v,2} + \|u\|_{l+2+\gamma, h-2+v,2} \right).$$

Now choose a non-negative cut-off function $\eta_1 \in C_0^\infty(B(e, 2), \mathbb{R}^N)$, with $\eta_1 = 1$ in $B(e, 1)$. Recall that $\eta \in C_0^\infty(B(0, 2))$ is also non-negative. Using Theorem 3.2 we obtain

$$(3.13) \quad \|\eta u\|_{l,h+2,1}^2 \leq C \left(\|u\|_{h,l,2}^2 + \sum_{i=1}^m \left[\|f_{i,(l,h)}\|_{0,1,2}^2 + \|f_{i,(l,h)}\|_{1,0,2}^2 \right] + \|f_{(l,h)}\|_{1,0,2}^2 \right).$$

From (3.11)–(3.12) one has

$$(3.14) \quad \|u\|_{l,h+2,1}^2 \leq C \left(\|u\|_{l+3, h+1,2}^2 + \|f\|_{l+1, h,2}^2 \right).$$

Applying the energy inequality to (3.9) we obtain that for each l and h the following estimate holds

$$(3.15) \quad \|\eta_1 v_{l,h+1}\|_{0,0,1}^2 \leq C \left(\|v_{l,h}\|_{0,0,2}^2 + i \sum_{i=1}^m \|f_{i,(l,h)}\|_{0,0,2}^2 + \|f_{(l,h)}\|_{0,0,2}^2 \right).$$

By (3.11)–(3.15) we infer

$$(3.16) \quad \|\eta u\|_{l,h+2,1}^2 \leq C \left(\|u\|_{l+3, h,2}^2 + \|f\|_{l+1, h,2}^2 \right).$$

If $h = 2\tilde{h}$, then iterating (3.16) \tilde{h} times one finds

$$(3.17) \quad \|u\|_{l,2\tilde{h}+2,1}^2 \leq C_{h,l} \left(\|u\|_{l+3\tilde{h}, 2,2}^2 + \|f\|_{l+3\tilde{h}-2, 2\tilde{h},2}^2 \right).$$

On the other hand, if we take l -th order derivatives of the original system (3.7) in the directions $T_1, \dots, T_k \in V^2$, we obtain the new system

$$(3.18) \quad \sum_{\beta=1}^N \sum_{i,j=1}^m A_{ij}^{\alpha\beta} X_i X_j v_{l,0}^\beta = f_{l,0}^\alpha, \quad \alpha = 1, \dots, N.$$

Theorem 3.2 applied to (3.18) gives

$$\|u\|_{l,2,1}^2 \leq C(\|u\|_{l,1,2} + \|f\|_{l+1,0,2}).$$

Applying the energy inequality to (3.18) one may assume that in the right-hand side of the last inequality there is the norm $\|u\|_{l,0,2}$, and not $\|u\|_{l,1,2}$.

From a repeated application of Theorem 3.1 (l times) one concludes

$$(3.19) \quad \|u\|_{l,2,1}^2 \leq C_l(\|u\|_{0,0,2}^2 + \|f\|_{l+1,0,2}^2).$$

Now (3.10) follows at once from (3.17) and (3.19). \square

Remark. It is important to observe that as a consequence of (3.10) one obtains the hypoellipticity of the system

$$\sum_{\beta=1}^N \sum_{i,j=1}^m A_{ij}^{\alpha\beta} X_i X_j u^\beta = 0, \quad \alpha = 1, \dots, N.$$

Hence, in the special setting of Carnot groups of step two, Theorem 3.6 generalizes to non-diagonal systems the basic results of Kohn [K] and Hörmander [H].

We now pass to the study of some consequences of Theorem 3.2 that will be useful in the study of the regularity of solutions of non-linear systems. The following corollaries are direct consequences of Theorem 3.2 and of the embedding of the Folland-Stein Sobolev spaces $\mathcal{L}^{h,2}$ in the non-isotropic Hölder classes (see section two and thm 5.15 in [F] for details).

Corollary 3.7 *Let \mathbf{G} be a Carnot group of step two and $\Omega \subset \mathbf{G}$ be an open subset. If $u \in \mathcal{L}_{\text{loc}}^{1,2}(\Omega)$ is a solution of the constant-coefficient system*

$$(3.20) \quad \sum_{i,j=1}^m \sum_{\beta=1}^N A_{ij}^{\alpha\beta} X_i X_j u^\beta = 0, \quad \alpha = 1, \dots, N,$$

in the gauge ball $B(p_o, 3R) \subset \Omega$, then u is smooth in $B(p_o, 3R)$. Moreover, there exists a positive constant C such that

$$\sup_{B(p_o, R)} \left(|u|^2 + R^2 |Xu|^2 + R^4 \sum_{i,j=1}^m |X_i X_j u|^2 \right) \leq C \frac{1}{|B(p_o, 2R)|} \int_{B(p_o, 2R)} |u|^2 dp.$$

Corollary 3.8 *Let \mathbf{G} , Ω and u be as in the previous corollary, then there exists a positive constant C such that for each $0 < r < R < 2$ one has*

$$(3.21) \quad \int_{B(p_o, r)} |u - u_{o, r}|^2 dp \leq C \left(\frac{r}{R}\right)^{Q+2} \int_{B(p_o, R)} |u - u_{o, R}|^2 dp,$$

$$(3.22) \quad \int_{B(p_o, r)} |u|^2 dp \leq \left(\frac{r}{R}\right)^Q \int_{B(p_o, R)} |u|^2 dp,$$

$$(3.23) \quad \int_{B(p_o, r)} |Xu - (Xu)_{p_o, r}|^2 dp \leq C \left(\frac{r}{R}\right)^{Q+2} \int_{B(p_o, R)} |Xu - (Xu)_{p_o, R}|^2 dp,$$

$$(3.24) \quad \int_{B(p_o, r)} \sum_{i, j=1}^m |X_i X_j u|^2 dp \leq C \left(\frac{r}{R}\right)^Q \int_{B(p_o, R)} \sum_{i, j=1}^m |X_i X_j u|^2 dp,$$

where $f_{p_o, r}$ denotes the average of the function f on the ball $B(p_o, r)$.

Proof. The estimates (3.21), (3.22) follow from Corollary 3.7 similarly to the classical elliptic case. The key observation is that the function

$$v(p) = u(p) - u_{p_o, R} - \sum_{i=1}^m (x_i(p) - x_i(p_o))(X_i u)_{p_o, R}$$

is still a solution of the system (3.20). This can be easily proved using the Baker-Campbell-Hausdorff formula (2.1) (see also the explicit expression of X'_i 's following (2.1)). Moreover, one has

$$\int_{B(p_o, R)} v dp = 0, \quad \int_{B(p_o, R)} X v dp = 0.$$

The inequality (3.23) requires a new proof, since, because of the non-commutativity of the X_j 's, we cannot simply say that $X_{i_o} u$ satisfies the same system as u . Instead, we have

$$\sum_{\beta=1}^N \sum_{i, j=1}^m A_{ij}^{\alpha\beta} X_i X_j (X_{i_o} u^\beta) = \sum_{j=1}^m \sum_{\beta=1}^N c_{k, j}^\alpha X_j T_k u^\beta, \quad \alpha = 1, \dots, N.$$

In order to deal with this system one has to make use of the results in Theorems 3.1 and 3.2, and Corollary 3.7. The first step consists in proving inequality (3.24). By homogeneity it suffices to do this for $R = 1$, i.e.

$$\int_{B(p_o, r)} \sum_{i, j=1}^m |X_i X_j u|^2 dp \leq r^Q \int_{B(p_o, 1)} \sum_{i, j=1}^m |X_i X_j u|^2 dp,$$

for $r < 1$. In the case $0 < r < 1/2$, the latter inequality is obtained from Corollary 3.7 and the sub-elliptic Poincaré inequality, see [J] or [Va], as follows

$$\begin{aligned} r^{-Q} \int_{B(p_o, r)} \sum_{i, j=1}^m |X_i X_j u|^2 dp &= r^{-Q} \int_{B(p_o, r)} \sum_{i, j=1}^m |X_i X_j v|^2 dp \leq C \sup_{B(p_o, 1/2)} \sum_{i, j=1}^m |X_i X_j v|^2 \\ &\leq C \int_{B(p_o, 1)} |v|^2 dp \leq C \int_{B(p_o, 1)} |Xv|^2 dp \leq C \int_{B(p_o, 1)} \sum_{i, j=1}^m |X_i X_j v|^2 dp. \end{aligned}$$

The remaining case $1/2 < r < 1$ is trivial and requires no proof. This establishes (3.24). Once this is done formula (3.23) follows from the Caccioppoli inequality for $w = X_i u$ (see Theorem 3.2), and from the Poincaré inequality. \square

In the diagonal case proofs of (3.21)–(3.24) have appeared in [XZ].

Remark. In Sect. 5 we will also need the following estimate: Let $u \in \mathcal{L}_{loc}^{1,2}(\Omega, \mathbb{R}^N)$ be a solution of (3.20). There exists $C > 0$ such that

$$(3.25) \quad \int_{B(0, r)} |Xu|^2 dp \leq C \left(\frac{r}{R}\right)^Q \int_{B(0, R)} |Xu|^2 dp.$$

This estimate too follows from Corollary 3.7 and from the Poincaré inequality.

To conclude the section we show how the method of difference quotients allows to establish L^2 estimates for some quasi-linear systems. Let $F : \mathbb{R}^{mN} \rightarrow \mathbb{R}^+$ be a C^2 convex function satisfying the following structural assumption: There exists a positive constant C such that for every $\xi \in \mathbb{R}^{mN}$ one has

$$C^{-1} |\xi|^2 \leq \sum_{i, j=1}^m \sum_{\alpha, \beta=1}^N \partial_{\xi_i^\alpha} \partial_{\xi_j^\beta} F(\xi) \leq C |\xi|^2.$$

As an example, we can consider the functional in (1.2) when the gradient $|Xu|$ is uniformly bounded by two positive constants.

Theorem 3.9 *Let $u \in \mathcal{L}_{loc}^{1,2}(\Omega)$, be a weak solution to the quasi-linear system*

$$(3.26) \quad \sum_{i=1}^m \int_{\Omega} F_{\xi_i^\alpha}(Xu) X_i \phi^\alpha dp = 0,$$

for each $\phi \in C_0^\infty(\Omega, \mathbb{R}^N)$. The solution u is weakly differentiable (in the L^2 sense) along the V^2 directions and moreover for each $i = 1, \dots, m$ and $j = 1, \dots, k$, we have

$$(3.27) \quad X_i u, T_j u \in \mathcal{L}_{loc}^{1,2}(\Omega).$$

Proof. The first step consists in differentiating (3.26) (in terms of difference quotients) along one of the T_{k_0} directions, in the same spirit as (3.1). In this setting we will obtain the system

$$\sum_{\beta=1}^N \sum_{i,j=1}^m \int_{\Omega} \partial_{\xi_i}^{\alpha} \partial_{\xi_j}^{\beta} F(Xu) [X_j u^{\beta}]_{(T_{k_0}, \frac{1}{2})} X_i \phi^{\alpha} dp = 0,$$

for each $\phi \in \mathcal{L}^{1,2}(\Omega, \mathbb{R}^N)$. From the Caccioppoli inequality and Theorems 2.3 and 2.2, as in (3.3) we obtain the existence of fractional derivatives of order $\frac{1}{2} - \varepsilon$ in the T_{k_0} direction

$$\int_{\Omega} |\partial_{T_{k_0}}^{\frac{1}{2}-\varepsilon} X(\eta u^{\alpha})|^2 dp \leq C \int_{B(0,2)} (|u|^2 + |Xu|^2) dp.$$

If we use Theorem 2.3 once more we obtain

$$\int_{\Omega} |\partial_{T_{k_0}}^{1-2\varepsilon} (\eta u_h^{\alpha})|^2 dp \leq C \int_{B(0,2)} (|u|^2 + |Xu|^2) dp.$$

Repeating the argument in (3.5) we can easily prove that $T_{k_0} u$ exists in L^2 and moreover

$$\int_{\Omega} |T_{k_0}(\eta u^{\alpha})|^2 dp \leq C \int_{B(0,2)} (|u|^2 + |Xu|^2) dp.$$

To continue we differentiate the system (3.26) along X_j with difference quotients of order one. As in Remark 3.4, commutators will appear. However, at this point we already have estimates on the commutators and this leads us directly to

$$\int_{\Omega} |(X(\eta u^{\alpha}))_{(X_j, 1)}|^2 dp \leq C \int_{B(0,2)} (|u|^2 + |Xu|^2) dp.$$

The proof is concluded once we apply Proposition 2.1. For more details we refer the reader to the paper [C1], where the scalar case was considered. \square

Since (3.26) is the Euler-Lagrange system of (1.1), the previous theorem has an immediate corollary.

Corollary 3.10 *Let F be as in Theorem 3.9. If $u \in \mathcal{L}_{\text{loc}}^{1,2}(\Omega)$ is a critical point of the functional*

$$\mathcal{F}(u, \Omega) = \int_{\Omega} F(Xu) dp,$$

then u is differentiable along the V^2 directions and moreover for each $i = 1, \dots, m$ and $j = 1, \dots, k$, we have

$$(3.27) \quad X_i u, T_j u \in \mathcal{L}_{\text{loc}}^{1,2}(\Omega).$$

4. Partial regularity of homogeneous systems and of local minimizers of convex functionals. The blow-up technique

In this section using Theorem 3.6 we develop the blow-up analysis and prove partial regularity of weak solutions of non-linear systems of the form

$$(4.1) \quad \sum_{\beta=1}^N \sum_{i,j=1}^m X_i \left(A_{ij}^{\alpha\beta}(x, u) X_j u^\beta \right) = 0,$$

where $A_{ij}^{\alpha\beta}(x, u)$ are bounded continuous or uniformly continuous functions satisfying

$$\sum_{\alpha,\beta=1}^N \sum_{i,j=1}^m A_{ij}^{\alpha\beta}(x, u) \xi_i^\alpha \xi_j^\beta \geq |\xi|^2,$$

for every $x \in \Omega$, and for every $N \times m$ matrix ξ . In section five we will allow in (4.1) a non zero right-hand side satisfying natural growth conditions. However, a method different from the one developed here will be needed, along with the additional hypothesis that u be bounded. Before presenting the proof of the main theorem we need to introduce some preliminary definitions and results.

Definition. Let $\lambda > -\frac{Q-2}{2}$. A function $u \in L^2(\Omega)$ is said to belong to the *Morrey space* $\mathcal{M}^{2,\lambda}(\Omega)$ if

$$r^{-Q} \int_{B(p,r) \cap \Omega} |u|^2 dp \leq Cr^{2(\lambda-1)},$$

for every $p \in \Omega$ and $0 < r < \text{diam}(\Omega)$.

Definition. Let $\lambda \geq -\frac{Q}{2}$. A function $u \in L^2(\Omega)$ is said to belong to the *Campanato space* $\mathcal{C}^{2,\lambda}(\Omega)$ if for every $p \in \overline{\Omega}$ and $0 < r < \text{diam}(\Omega)$ one has $U(p_0, R) \leq CR^{2\lambda}$, where we have let

$$U(p_0, R) = \frac{1}{|B(p_0, R)|} \int_{B(p_0, R) \cap \Omega} |u(p) - u_{p_0, R}|^2 dp.$$

Remark. The classical characterization of Hölder continuous functions via integral means can be generalized to the setting of Carnot groups (see for instance [C2] or [Ga] and references therein). In particular, adapting the arguments in [DP] (see also Chapter 3, Theorem 1.2 in [Gi]), one has the following: *Let $0 < \lambda < 1$. If Ω has the A -property, i.e., there exists $A > 0$ such that for $R > 0$ and $p \in \Omega$*

$$|\Omega \cap B(p, R)| \geq A|B(p, R)|,$$

then $\mathcal{C}^{2,\lambda}(\Omega) \subset \Gamma_{\text{loc}}^\lambda(\Omega)$. In particular, the gauge pseudo-balls have the A -property (for a proof of this fact see for instance [CG]), and therefore for every $p_0 \in \mathbf{G}$, and every $r > 0$ one has

$$\mathcal{C}^{2,\lambda}(B(p_0, r)) \subset \Gamma_{\text{loc}}^\lambda(B(p_0, r)) \text{ for } 0 < \lambda < 1.$$

This continuous inclusion will be henceforth tacitly used.

Our first result in this section generalizes a beautiful theorem of Giusti and Miranda [GM].

Lemma 4.1 *Let u be a weak solution of (4.1), with continuous $A_{ij}^{\alpha\beta}$'s. For each $0 < \tau < 1$, and $M > 0$, there exist ε_o , and $R_o > 0$ such that if for some $p_o \in \Omega$, one has $U(p_o, R) < \varepsilon_o^2$, and $|u_{p_o, R}| \leq M$ for $R \leq \min(R_o, d(p_o, \partial\Omega))$, then*

$$U(p_o, \tau R) \leq C\tau^2 U(p_o, R).$$

Here, C is the same as in (3.21).

Proof. We argue by contradiction, assuming that there exists $0 < \tau < 1$, $M > 0$, sequences $(p_n) \in \Omega$, $\varepsilon_n \rightarrow 0$, $R_n \rightarrow 0$, and a sequence u^n (the upper index is not an exponent) of weak solutions to (4.1), such that

$$(4.2) \quad |u_{p_o, R_n}^n| \leq M,$$

$$(4.3) \quad U^n(p_n, R_n) = \varepsilon_n^2 \rightarrow 0,$$

and

$$(4.4) \quad U^n(p_n, \tau R_n) > 2C\tau^2 \varepsilon_n^2.$$

Set

$$(4.5) \quad v^n(q) = \varepsilon_n^{-1} [u^n(p_n \delta_{R_n}(q)) - u_{p_n, R_n}^n].$$

From the homogeneity of the X_i 's we have

$$(4.6) \quad (X_i u^n)(p_n \delta_{R_n}(q)) = \varepsilon_n R_n^{-1} (X_i v^n)(q).$$

For $R > 0$, set $B_R = B(e, R)$, where e is the group identity. From (4.1), and (4.6) we infer that for every $\phi \in C_0^\infty(B_1)$ the function v^n satisfies the system

$$(4.7) \quad \sum_{\beta=1}^N \sum_{i,j=1}^m \int_{B_1} A_{ij}^{\alpha\beta} (p_n \delta_{R_n}(q), \varepsilon_n v^n + u_{p_n, R_n}^n) X_j (v^n)^\beta(q) X_i \phi^\alpha(q) dq = 0,$$

$$\alpha = 1, \dots, N.$$

Moreover, $v_{B_1}^n = 0$,

$$(4.8) \quad V^n(e, 1) = \frac{1}{|B_1|} \int_{B_1} |v^n|^2 dq = 1,$$

and

$$(4.9) \quad V^n(e, \tau) > 2C\tau^2.$$

From (4.8), taking a sub-sequence that we still denote as v^n , one obtains the existence of $v \in L^2(B_1)$ such that $v_n \rightarrow v$ weakly in $L^2(B_1)$, while $\varepsilon_n v^n \rightarrow 0$ a.e. in B_1 . From the continuity assumptions on $A_{ij}^{\alpha\beta}$ and (4.2) one finds

$$A_{ij}^{\alpha\beta} \left(p_n \delta_{R_n}(q), \varepsilon_n v^n + u_{p_n, R_n}^n \right) \rightarrow B_{ij}^{\alpha\beta}.$$

At this point we recall the compact embedding

$$\mathcal{L}^{1,2}(B_R) \hookrightarrow L^2(B_R),$$

noted in section two. This result, coupled with the Caccioppoli inequality Lemma 2.5, implies $Xv^n \rightarrow Xv$ weakly in $L^2(B_1)$ and $v^n \rightarrow v$ strongly in $L^2(B_1)$. All these facts, together with (4.9), give for every $\phi \in C_0^\infty(B_1)$

$$(4.10) \quad \sum_{\beta=1}^N \sum_{i,j=1}^m \int_{B_1} B_{ij}^{\alpha\beta} X_i^i v^\beta X_i^1 \phi^\alpha dp = 0, \quad \alpha = 1, \dots, N.$$

The contradiction now follows from (4.8), (4.9) and from (3.21). □

Lemma 4.3 *Choose $\tau^2 \leq 2C$ and let u be a weak solution of (4.1) in Ω . If at a point $p_o \in \Omega$ one has for some $R < R_o = R_o(M)$,*

$$(4.11) \quad |u_{p_o, R}| < M,$$

and

$$(4.12) \quad U(p_o, R) \leq \eta_o^2(M),$$

with $\eta_o(M) \leq \min\{\varepsilon_o(M), M\tau^{(Q+2)/2}(1 - \tau\sqrt{2C})\}$, then for every integer k we have

$$U(p_o, \tau^k R) \leq (2C\tau^2)^k U(p_o, R).$$

We omit the proof of Lemma 4.3, since it is completely analogous to that of the elliptic case in [GM] (see also Chap. 4, p.95, Theorem 1.1 in [Gi]). The key step is an iteration based on Lemma 4.1.

Theorem 4.4 *Let $u \in \mathcal{L}_{loc}^{1,2}(\Omega)$ be a weak solution of (4.1) in Ω . Denote by*

$$S_1 = \{p_o \in \Omega \mid \sup_R |u_{p_o, R}| = \infty\},$$

$$S_2 = \{p_o \in \Omega \mid \liminf_{R \rightarrow 0^+} U(p_o, R) > 0\},$$

and we let $S = S_1 \cup S_2$ indicate the singular set. There exists $0 < \alpha < 1$ (depending only on \mathbf{G} , Ω , and on the structural constants in (1.4)) such that $u \in \Gamma_{loc}^{0,\alpha}(\Omega \setminus S)$.

Proof of Theorem 4.4. The proof of the Hölder regularity follows from Lemma 4.3, by adapting the arguments in [GM]. If $p \notin S$, then there exist $M > 0$ and $\eta_o(M)$ such that in a small neighborhood $B(p, r_o)$ of p the inequalities (4.11) and (4.12) hold. From Lemma 4.3 we conclude for each $R > 0$ small enough and for some $0 < a < 1$

$$(4.13) \quad U(p, R) \leq CR^{2a}.$$

This implies $u \in \mathcal{C}^{2,a}(B(p, r_o))$, the sub-elliptic Campanato class introduced above, and thereby the Hölder continuity of u in $B(p, r_o)$ follows. \square

Our next goal is to estimate the Hausdorff dimension (with respect to the group pseudo-metric) of the singular set (compare with [Giu]).

Lemma 4.5 *Let $u \in L^2_{\text{loc}}(\Omega)$. If (4.13) holds for some $a > 0$ then $\lim_{R \rightarrow 0^+} u_{B(p,R)}$ exists and is finite. Consequently, p is regular if and only if (4.13) holds for some $a > 0$.*

The elementary proof of Lemma 4.5 is left to the reader.

Remark. In the case where $A_{ij}^{\alpha\beta}$ are uniformly continuous the singular set coincides with S_2 and a point p is regular if and only if (4.13) holds with $a = 0$.

Lemma 4.6 *Let $f \in L^1_{\text{loc}}(\Omega)$ and set for $0 < \alpha < Q$*

$$F_\alpha = \left\{ p_o \in \Omega \mid \limsup_{r \rightarrow 0^+} r^{-\alpha} \int_{B(p_o, r)} |f(p)| dp > 0 \right\}.$$

One has $H_\alpha(F_\alpha) = 0$.

Proof. It suffices to show that, if for $s > 0$ we let

$$F_s = \left\{ p_o \in \Omega \mid \limsup_{r \rightarrow 0^+} r^{-\alpha} \int_{B(p_o, r)} |f(p)| dp > s^{-1} \right\},$$

then $H_\alpha(F^s) = 0$. Let $K \subset \Omega$ be a compact set and consider an open set O such that $K \subset O \subset \overline{O} \subset \Omega$. Fix $0 < \varepsilon < D(K, \partial O)$. For every $p \in F_s$ there exists $r(p) \in (0, \varepsilon)$ such that

$$r(p)^{-\alpha} \int_{B(p, r(p))} |f(q)| dq > s^{-1}.$$

By the covering lemma of Wiener type in a space of homogeneous type (see [CW]) there exists a constant η , depending only on the homogeneous dimension Q , and a sequence of pairwise disjoint pseudo-balls $B_i = B(p_i, r(p_i))$, such that $F^s \subset \bigcup_{i=1}^{\infty} \eta B_i$. One has

$$(4.14) \quad \sum_{i=1}^{\infty} r_i^\alpha \leq s \int_{\bigcup_{i=1}^{\infty} B_i} |f(q)| dq,$$

and

$$(4.15) \quad \left| \bigcup_{i=1}^{\infty} B_i \right| \leq C\eta\varepsilon^{Q-\alpha_s} \int_O |f(q)|dq.$$

Letting $\varepsilon \rightarrow 0$ in (4.15) we infer $H_\alpha(F^s) = 0$. By the arbitrariness of $s > 0$ the conclusion follows. \square

Theorem 4.7 *Let u be a weak solution of (4.1) in Ω . For every $\varepsilon > 0$ we have*

$$H_{Q-2+\varepsilon}(S) = 0.$$

Moreover, if the coefficients $A_{ij}^{\alpha\beta}$ are uniformly continuous then one may take $\varepsilon = 0$.

Proof. Set

$$E_\varepsilon = \{p_o \in \Omega \mid \limsup_{r \rightarrow 0^+} r^{-\varepsilon} U(p_o, r) > 0\}.$$

Lemma 4.5 implies $S = \bigcap_{\varepsilon > 0} E_\varepsilon$. We want to prove that for each $\varepsilon > 0$

$$H_{Q-2+\varepsilon}(E_\varepsilon) = 0.$$

From the Poincaré inequality [J], [Va] one has

$$r^{-\varepsilon} U(p_o, r) \leq Cr^{-Q+2-\varepsilon} \int_{B(p_o, r)} |Xu(p)|^2 dp.$$

Since $|Xu|^2$ is in $L^1_{\text{loc}}(\Omega)$, the conclusion follows from Lemma 4.6. \square

Remark 4.8 It is easy to see that the solutions of (4.1) satisfy for some $r < 2$ the following reverse Hölder inequality

$$(4.16) \quad \int_{B_1} (1 + |Xu|)^2 dp \leq \left(\int_{B_2} (1 + |Xu|)^r dp \right)^{\frac{2}{r}}.$$

An adaptation of a well known argument due to Gehring (see [Ge], [BI]), based on the Calderón-Zygmund decomposition in spaces of homogeneous type (see [CW], [Ch]) allows to establish the higher integrability of Xu (see also [C3], [FSC], [GMo] and [GG3]). Since one can easily generalize Theorem 4.7 to the L^p setting, with $p \geq 2$, then one infers that for some $p > 2$

$$H_{Q-p+\varepsilon}(S) = 0.$$

We conclude this section with the proof of Theorem 1.1 which we will divide in two steps, Theorems 4.10 and 4.11. We stress that in the proof of Theorem 4.10 the convexity assumption on the energy function F , stated in the introduction, is only used in the study of the constant coefficient system arising as limit of the blow-ups of the Euler-Lagrange equations of (1.1). Such study was conducted in section three. Since the functional (1.1) is convex, then in particular it is *strictly*

quasi-convex [E], i.e., there exists $\gamma > 0$ such that for each $N \times m$ constant matrix A and every $\phi \in C_0^\infty(\Omega, \mathbb{R}^N)$

$$(4.17) \quad \int_{\Omega} \left[F(A) + \gamma(1 + |X\phi|^{q-2})|X\phi|^2 \right] dp \leq \int_{\Omega} F(A + X\phi) dp,$$

We can thus adapt to the present sub-Riemannian setting the ideas in L.C. Evans' paper [E] on quasi-convex elliptic integrands. In doing so, we will avoid any repetition, highlighting only those points where different arguments are needed. These arguments are collected in the following Propositions 4.8 and 4.9 which respectively contain a new Caccioppoli type inequality and the essence of the blow-up technique.

In the following, we will denote points $p \in \mathbf{G}$ by exponential coordinates $p = (x, t)$ where $x \in V^1$ and $t \in V^2$. Observe that if A is an $N \times m$ matrix, then $X(A \cdot x)(p) = A$ for every $p \in \mathbf{G}$. If f is a function, we let $\int_{\Omega} f(p) dp = \frac{1}{|\Omega|} \int_{\Omega} f(p) dp$.

Proposition 4.9 *Let $u \in \mathcal{L}_{\text{loc}}^{1,q}(\Omega)$ be a local minimizer of (1.1). For each $L > 0$ there exists $C_1(L) > 0$ such that*

$$(4.18) \quad \int_{B(p_o, r)} (1 + |Xu - A|^{q-2}) |Xu|^2 dp \leq C(L) \left(\frac{1}{r^2} \int_{B(p_o, r)} |u(p) - a - A \cdot (x - x_o)|^2 dp + \frac{1}{r^q} \int_{B(p_o, r)} |u(p) - a - A \cdot (x - x_o)|^q dp \right),$$

for each $B(p_o, 2r) \subset \Omega$, $a \in \mathbb{R}^N$, and any $N \times m$ constant matrix A with $|A| \leq L$.

Proof. Without loss of generality we can assume $p_o = e$, the group identity. For $r \leq t < s \leq 2r$ choose a test function $\xi \in C_0^\infty(B(p_o, s))$, with $\xi = 1$ in $B(p_o, t)$. Define two auxiliary functions $\phi(p) = \xi(u - a - A \cdot x)$, and $\varphi(p) = (1 - \xi)(u - a - A \cdot x)$. Following Evans [E], and using (4.17) we can easily prove

$$(4.19) \quad \int_{B(e, s)} \left[F(A) + \gamma(1 + |X\phi|^{q-2})|X\phi|^2 \right] dp \leq \int_{B(e, s)} \left[F(Xu) - F'(Xu)X\varphi + C(1 + |Xu|^{q-2} + |X\varphi|^{q-2})|X\varphi|^2 \right] dp.$$

Since u is a minimizer of (1.1), using $|A| \leq L$, and $D^2 F(Xu) \leq C(1 + |Xu|)^{q-2}$, we obtain

$$\int_{B(e, s)} F(Xu) dp \leq \int_{B(e, s)} \left[F(A) + F'(A)X\varphi + C((1 + |Xu|)^{q-2} |X\varphi|^2) \right] dp.$$

The latter inequality and (4.19) give

$$\begin{aligned}
(4.20) \quad & \gamma \int_{B(e,s)} (1 + |X\phi|^{q-2}) |X\phi|^2 dp \\
& \leq C \int_{B(e,s)} \left[(1 + |Xu|)^{q-2} |Xu - A| |X\phi| \right. \\
& \qquad \qquad \qquad \left. + C(1 + |Xu|^{q-2} + |X\phi|^{q-2}) |X\phi|^2 \right] dp.
\end{aligned}$$

Recalling the definition of ξ and using (4.20) we easily arrive at

$$\begin{aligned}
& \int_{B(e,t)} (1 + |Xu - A|^{q-2}) |Xu|^2 dp \\
& \leq C \left\{ \int_{B(e,s) \setminus B(0,t)} (1 + |Xu - A|)^{q-2} |Xu|^2 dp \right. \\
& \qquad \qquad \qquad \left. + \int_{B(e,2r)} \left[\frac{|u - a - A \cdot x|^2}{(s-t)^2} + \frac{|u - a - A \cdot x|^q}{(s-t)^q} \right] dp \right\}.
\end{aligned}$$

The conclusion follows from a standard argument based on the so-called *hole filling* technique (see [Gi] or [E, Lemma 5.1 on p. 239]). \square

Proposition 4.10 *Let u be as in Proposition 4.8. For each $L > 0$, there exists $C_2(L) > 0$ such that for every $0 < \theta < 1/4$ one can find $\varepsilon(\theta, L) > 0$ for which, if $|(Xu)_{p_o,r}| \leq L$, $|(Xu)_{p_o,\theta r}| \leq L$, and $U(p_o, r) \leq \varepsilon$, then*

$$U(p_o, \theta r) \leq C_2(L) \theta^2 U(p_o, r),$$

where

$$U(p_o, r) = \int_{B(p_o,r)} (1 + |Xu - (Xu)_{p_o,r}|)^{q-2} |Xu - (Xu)_{p_o,r}|^2 dp.$$

Proof. Similarly to Lemma 4.3 the proof proceeds by contradiction. If the assertion of the proposition were false, then there would exist a sequence of balls $B(p_k, r_k) \subset \Omega$, such that $|(Xu)_{p_k,r_k}| \leq L$, $|(Xu)_{p_k,\theta r_k}| \leq L$, $U(p_k, r_k) = \lambda_k^2 \rightarrow 0$ as $k \rightarrow \infty$, and

$$(4.21) \quad U(p_k, \theta r_k) > C_2(L) \theta^2 \lambda_k^2.$$

We now consider the blow-up sequence defined by

$$v^k(p) = \frac{u(p_k \delta_{r_k}(p)) - (u)_{p_k,r_k} - (Xu)_{p_k,r_k} \cdot x}{\lambda_k r_k}.$$

By a simple rescaling, using the fact that the vector fields X_i 's are homogeneous of degree one with respect to the non-isotropic group dilations, one obtains from (4.21) the inequality

$$(4.22) \quad \int_{B(e,\theta)} (1 + |\lambda_k^{q-2} |Xv^k - (Xv^k)_{e,\theta}|^{q-2}) |Xv^k - (Xv^k)_{e,\theta}|^2 dp > C_2(L) \theta^2.$$

Using the compact embedding $\mathcal{L}^{1,2}(B(e, 2)) \hookrightarrow L^2(B(e, 2))$, one obtains the existence of $v \in \mathcal{L}^{1,2}(B(e, 2))$, such that $v^k \rightarrow v$ weakly in $\mathcal{L}^{1,2}(B(e, 2))$, and strongly in $L^2(B(e, 2))$. Moreover, $(Xu)_{p_k, r_k} \rightarrow A \in \mathbb{R}^{N \times m}$. Repeating the argument in [E, Lemma 6.2], we infer that v is a weak solution of a constant coefficient sub-elliptic system of the form (4.10), with coefficients depending on A . This system satisfies (1.7) since the original functional is convex. It is at this point that the sub-elliptic estimates from section three enter into play. In particular, by means of Corollary (3.7) and the Poincaré inequality [J], [Va], one has

$$\int_{B(e, 2\theta)} |Xv - (Xv)_{B(e, \theta)}|^2 \leq C(L)\theta^2$$

(note that the average in the integral is on a smaller ball). Using the latter inequality and the Sobolev-Poincaré inequalities recalled in section two, one obtains

$$(4.23) \quad \lim_{k \rightarrow \infty} \lambda_k^{q-2} \int_{B(e, 2\theta)} |v^k - (v^k)_{e, 2\theta} - (Xv^k)_{e, \theta} \cdot x|^q dp = 0.$$

The inequality (4.18) in Proposition 4.8, and (4.23), finally give

$$\limsup_{k \rightarrow \infty} \int_{B(e, \theta)} \left(1 + |\lambda_k^{q-2} |Xv^k - (Xv^k)_{e, \theta}|^{q-2}\right) |Xv^k - (Xv^k)_{e, \theta}|^2 dp \leq C_3(L)\theta^2.$$

Since the latter inequality contradicts (4.22), we have completed the proof. \square

Theorem 4.11 *Let $u \in \mathcal{L}_{\text{loc}}^{1,q}(\Omega)$ be a local minimizer of (1.1). There exists a closed subset $S \subset \Omega$, with $|S| = 0$, such that Xu is locally Hölder continuous in $\Omega \setminus S$.*

Proof. With Propositions 4.8 and 4.9 in hands, the proof of Theorem 4.10 follows along the lines of that of Lemma 4.3 without major modifications. For more details we refer the reader to Evans' paper [E]. \square

Proposition 4.9 characterizes the singular set S . One has $S = S_1 \cup S_2$, where

$$S_1 = \{p_o \in \Omega \mid \lim_{R \rightarrow 0^+} R^{-Q} \int_{B(p_o, R)} |Xu - Xu_{p_o, R}|^q dp > 0\},$$

and

$$S_2 = \{p_o \in \Omega \mid \lim_{R \rightarrow 0^+} Xu_{p_o, R} \text{ does not exist}\}.$$

Clearly, S has measure zero. In the case $q = 2$ we can say more about the singular set.

Theorem 4.12 *Let u be as in Theorem 4.10, with $q = 2$. For every $\varepsilon > 0$ one has*

$$H_{Q-2+\varepsilon}(S) = 0.$$

Proof. By virtue of Theorem 3.9 and Corollary 3.10 we have that for each $i = 1, \dots, m$ and $j = 1, \dots, k$

$$(4.26) \quad X_i u, T_j u \in \mathcal{L}_{\text{loc}}^{1,2}(\Omega).$$

The conclusion now follows from an argument similar to that in Theorem 4.7. One needs to use the results from Lemma 4.5, 4.6, the characterization of the singular set from Proposition 4.10, the Poincaré inequality (see [J], [Va]), and (4.26). \square

Theorem 4.11 generalizes a result of Giusti and Miranda [GM].

5. Partial regularity for non-homogeneous systems via the Morrey-Campanato approach

In this section we establish the partial regularity of weak solutions of (1.3). Our results generalize to the sub-Riemannian setting those of Giaquinta and Giusti in [GG2]. Here, in order to treat systems whose right hand side satisfies the natural growth conditions (1.5), we need to assume a priori the boundedness of the weak solution $|u| \leq M$. In connection with (1.4), (1.5), we also assume

$$(5.1) \quad a(M)M \leq \lambda(M)/2.$$

For the necessity of (5.1) see [Gi] Remark 1.3, Chap. 3. Since for the most part the proofs of the results in this section follow closely those of their Euclidean predecessors, we provide only sketches of the arguments, emphasizing only the parts where the results of section three are used.

We begin by introducing several preliminary results. The first is a maximum principle for constant coefficients systems which extends a result of Cannarsa [Ca].

Lemma 5.1 *Let $\phi \in \mathcal{L}_{\text{loc}}^{1,2}(B(e, 1)) \cap L^\infty(B(e, 1))$, and assume that $X\phi \in \mathcal{M}^{2,0}(B(e, 1))$. Set*

$$L^2 = \sup_{p_0 \in B(e, 1), r > 0} r^{2-Q} \int_{B(p_0, r) \cap B(e, 1)} |X\phi|^2 dp < \infty.$$

Let $v \in \mathcal{L}_{\text{loc}}^{1,2}(B(e, 1))$ be a weak solution in $B(e, 1)$ of the constant coefficient system

$$\sum_{i,j=1}^m \sum_{\beta=1}^N A_{ij}^{\alpha\beta} X_i X_j v^\beta = 0,$$

with $v - \phi \in \mathcal{L}_0^{1,2}(B(e, 1))$, and $A_{ij}^{\alpha\beta}$ satisfying the coercivity condition (1.7). There exists a constant $C > 0$ independent of v , and ϕ such that

$$\sup_{B(e, 1)} |v| \leq C \left(\sup_{B(e, 1)} |\phi| + L \right).$$

Proof. The proof is divided in two parts. First we show that

$$(i) \sup_{B(e,1)} |v|^2 \leq C \left(\sup_{B(e,1)} |\phi|^2 + L^2 + \sup_{p_o \in \overline{B(e,1)}, r > 0} r^{2-Q} \int_{B(p_o, r) \cap B(e,1)} |Xv|^2 dp \right).$$

To conclude the proof we then establish the following Campanato-type estimate

$$(ii) \sup_{p_o \in \overline{B(e,1)}, r > 0} r^{2-Q} \int_{B(p_o, r) \cap B(e,1)} |Xv|^2 dp \leq CL^2.$$

In order to prove (i) let $p_o \in B(e, 1)$, and set $d = d(p_o, \partial B(e, 1)) = d(p_o, z_o)$ for some $z_o \in \partial B(e, 1)$. In virtue of (3.22) for any $0 < r < R = \min(d, 1)$ one has

$$\begin{aligned} r^{-Q} \int_{B(p_o, r)} |v|^2 dp &\leq CR^{-Q} \int_{B(p_o, R)} |v|^2 dp \leq \\ &C \left[\sup_{B(e,1)} |\phi|^2 + R^{-Q} \int_{B(z_o, 2R) \cap B(e,1)} |\phi - v|^2 dp \right]. \end{aligned}$$

From the Poincaré inequality we infer

$$r^{-Q} \int_{B(p_o, r)} |v|^2 dp \leq C \left[\sup_{B(e,1)} |\phi|^2 + L^2 + R^{2-Q} \int_{B(z_o, 2R) \cap B(e,1)} |Xv|^2 dp \right].$$

The latter inequality gives the estimate (i). In order to prove (ii) we observe that the function $w = \phi - v \in \mathcal{L}_o^{1,2}(B(e, 1))$ is a weak solution of the system

$$\sum_{i,j=1}^m \sum_{\beta=1}^N A_{ij}^{\alpha\beta} X_i X_j w = \sum_{i=1}^m X_i F_i^\alpha \alpha = 1, \dots, N$$

where

$$F_i^\alpha := \sum_{j=1}^m \sum_{\beta=1}^N A_{ij}^{\alpha\beta} X_j \phi \alpha = 1, \dots, N$$

in $B(e, 1)$. Let $p_o \in B(e, 1)$, and $r > 0$. Choose a test function $\eta \in C_o^\infty(B(p_o, 2r))$, such that $\eta = 1$ in $B(p_o, r)$ and $|X\eta| \leq Cr^{-1}$. If we multiply both sides of the latter equation by $w\eta^2$ and integrate by parts, then by (1.7) we obtain the following energy inequality

$$\begin{aligned} \int_{B(p_o, r) \cap B(e,1)} |X\phi - Xv|^2 dp &\leq C \int_{B(p_o, 2r) \cap B(e,1)} |F|^2 dp \\ &\leq C \int_{B(p_o, 2r) \cap B(e,1)} |X\phi|^2 dp. \end{aligned}$$

From Schwarz inequality one infers

$$\int_{B(p_o, r) \cap B(e,1)} |Xv|^2 dp \leq C \int_{B(p_o, 2r) \cap B(e,1)} |X\phi|^2 dp.$$

The estimate (ii) follows from the latter inequality, once we multiply both sides by r^{2-Q} and take the supremum over all gauge balls $B(p_o, r)$. \square

One of the main ingredients in the proof of partial regularity for non-homogeneous systems is the higher integrability of the gradient of the solution.

Lemma 5.2 *Let $u \in \mathcal{L}_{\text{loc}}^{1,2}(\Omega)$ be a bounded weak solution of (1.3), and assume that (1.4), (1.5) and (5.1) hold. There exists $q > 2$ such that for every $B(p_0, 2r) \subset \Omega$ one has*

$$(5.2) \quad \|Xu\|_{L^q(B(p_0,r))} \leq r^{Q(\frac{1}{q}-\frac{1}{2})} C \|(1 + |Xu|)\|_{L^2(B(p_0,2r))}.$$

Moreover, if the right hand side of (1.3) is zero (i.e. $a(M) = b = 0$), then one has

$$(5.3) \quad \|Xu\|_{L^q(B(p_0,r))} \leq r^{Q(\frac{1}{q}-\frac{1}{2})} C \|Xu\|_{L^2(B(p_0,2r))}.$$

Proof. The hypothesis (1.4), (1.5) allow to prove a reverse Hölder inequality for the function $(1 + |Xu|)$. The proof is standard and rests on the Caccioppoli and Poincaré-Sobolev inequalities in section two. Once we have the reverse Hölder inequality, then a well-known argument (see [Ge], [BI], [FSC] and Chap. 5 in [Gi]), yields the higher integrability of the horizontal gradient. In the homogeneous case, when $a(M) = b = 0$, the reverse Hölder inequality can be proved directly for the function $|Xu|$, without need of the additional constant. The higher integrability argument mentioned above yields inequality (5.3) (see also Remark 4.8 for more references). \square

The next lemma is the main step in the regularity proof, it substitutes the blow-up argument in section four. We shall suppose that the coefficients $A_{ij}^{\alpha\beta}$ are continuous functions of their arguments, and we shall denote by $\omega(r)$ a bounded continuous function such that

$$|A(x, u) - A(y, v)| \leq \omega(|x - y|^2 + |u - v|^2),$$

whenever $x, y \in \Omega$ and $|u|, |v| \leq M$. Without loss of generality we may assume that ω is increasing, concave and $\omega(0) = 0$ (unless the coefficients A are uniformly continuous one cannot expect ω to be continuous).

Lemma 5.3 *Let \mathbf{G} be a Carnot group of step two and $\Omega \subset \mathbf{G}$ an open set. Let $u \in \mathcal{L}_{\text{loc}}^{1,2}(\Omega)$ be a bounded weak solution of (1.3), and assume that (1.4), (1.5) and (5.1) hold. There exists $R_0 = R_0(M, a, b) > 0$ such that for each point $p_0 \in \Omega$ and for every $0 < r < R < \min(R_0, d(p_0, \partial\Omega))$, one has*

$$(5.4) \quad \int_{B(p_0,r)} (1 + |Xu|^2) dp \leq C \int_{B(p_0,R)} (1 + |Xu|^2) dp \left[\left(\frac{r}{R}\right)^Q + \chi(p_0, R) \right],$$

where

$$\chi(p_0, R) = R + \phi(R)^{\frac{q-2}{2q}} + \omega(C\phi(R))^{\frac{q-2}{q}},$$

and $\phi(R) = R^{2-Q} \int_{B(p_0,R)} (1 + |Xu|^2) dp$, with $q > 2$ is as in (5.2).

Proof. First we prove that there exists $R > 0$ such that

$$(5.5) \quad \sup_{r < R, B(p_0, r) \subset \Omega} r^{2-Q} \int_{B(p_0, r)} |Xu|^2 dp \leq C \sup_{\Omega} |u|.$$

Let $B(p_0, 2r) \subset \Omega$. Set $\phi = \eta^2 u$, with $\eta \in C_0^\infty(B(p_0, 2r))$ $\eta = 1$ in $B(p_0, r)$, $|X\eta| \leq Cr^{-1}$ and use it as a test function in the weak form of (1.3),

$$\begin{aligned} & \sum_{\beta=1}^N \sum_{i,j=1}^m \int_{B(p_0, 2r)} \left(A_{ij}^{\alpha\beta}(x, u) X_j u^\beta \right) X_i \phi^\alpha(x) dx \\ &= \int_{B(p_0, 2r)} f^\alpha(x, u, Xu) \phi^\alpha(x) dx. \alpha = 1, \dots, N \end{aligned}$$

By means of (1.4), and (1.5) we obtain the Caccioppoli inequality

$$\begin{aligned} \int_{B(p_0, 2r)} |Xu|^2 \eta^2(x) dx &\leq Cr^2 \int_{B(p_0, 2r)} |u|^2(x) dx + \int_{B(p_0, 2r)} |f|(x) |u|(x) \eta^2(x) dx. \\ &\leq Cr^2 \int_{B(p_0, 2r)} |u|^2(x) dx + \int_{B(p_0, 2r)} (a(M)|u| |Xu|^2 \eta^2 + b(M)|u| \eta^2) dx. \end{aligned}$$

Now we use (5.1) and obtain for every $p \in \Omega$ and small $r > 0$,

$$(5.6) \quad \int_{B(p, r)} |Xu|^2 \leq Cr^{Q-2} [M^2 + br^2 M].$$

The estimate (5.5) follows from (5.6) if we choose $R > 0$ small enough so that $br^2 \leq M$. Inequality (5.4) is obvious for $\frac{R}{2} < r < R$, hence we will consider only the case $0 < r < \frac{R}{2}$.

Next, we split u as $u = v + w$, where $v \in \mathcal{L}_{\text{loc}}^{1,2}(B(p_0, R/2), \mathbb{R}^n)$ is a solution of

$$\sum_{i,j=1}^{2n} \sum_{\beta=1}^N A_{ij}^{\alpha\beta}(p_0, (u)_{p_0, R/2}) X_i X_j v^\beta = 0,$$

with boundary data $v - u \in \mathcal{L}_0^{1,2}(B(p_0, R/2))$.

Recall that by (3.25) we have that for any $0 < r < R/2$ the following inequality holds

$$\int_{B(0, r)} |Xv|^2 dx \leq C \left(\frac{r}{R} \right)^Q \int_{B(p_0, R/2)} |Xv|^2 dx.$$

From Lemma 5.1 and (5.5) one obtains

$$(5.7) \quad \sup_{B(p_0, R/2)} |v| \leq CM.$$

From (5.7) we deduce that $|w| \leq CM$.

Observe that $w \in \mathcal{L}_0^{1,2}(B(p_0, R/2), \mathbb{R}^n)$, is a weak solution of the system

$$\begin{aligned} & \sum_{\beta=1}^N \sum_{i,j=1}^m \int_{B(p_0, R/2)} \left(A_{ij}^{\alpha\beta}(p_0, u_{p_0, R/2}) X_j w^\beta \right) X_i \phi^\alpha(x) dx \\ &= \int_{B(p_0, R/2)} \left(A_{ij}^{\alpha\beta}(p_0, u_{p_0, R/2}) - A_{ij}^{\alpha\beta}(x, u) \right) X_j w^\beta X_i \phi^\alpha(x) dx \\ & \quad + \int_{B(p_0, R/2)} f^\alpha(x, u, Xu) \phi^\alpha(x) dx. \quad \alpha = 1, \dots, N \end{aligned}$$

Set $\phi = w$ and we obtain the Caccioppoli inequality

$$(5.8) \quad \int_{B(p_0, R/2)} |Xw|^2 dp \leq C \left[\int_{B(p_0, R/2)} (\omega^2 + |w|) |Xu|^2 + \int_{B(p_0, R/2)} |w| dp \right],$$

where $\omega^2 = \omega^2(R^2 + |u - u_{p_0, R/2}|^2)$.

From the Poincaré inequality we have

$$\begin{aligned} \frac{1}{|B(p_0, R/2)|} \int_{B(p_0, R/2)} |w| dx &\leq CR \frac{1}{|B(p_0, R/2)|} \int_{B(p_0, R/2)} |Xw| dx \\ &\leq CR \left(\varepsilon \int_{B(p_0, R/2)} |Xw|^2 dx + \frac{1}{\varepsilon} \right). \end{aligned}$$

From (5.8) and the Poincaré inequality we obtain

$$(5.9) \quad \frac{1}{|B(p_0, R/2)|} \int_{B(p_0, R/2)} |w| dx \leq CR \left(\frac{1}{|B(p_0, R/2)|} \int_{B(p_0, R/2)} (1 + |Xu|^2) dx \right).$$

To complete the estimate of the terms in (5.8) we use the results in Lemma 5.2 and obtain

$$(5.10) \quad \begin{aligned} & \frac{1}{|B(p_0, R/2)|} \int_{B(p_0, R/2)} \omega^2 |Xu|^2 dx \\ & \leq C \frac{1}{|B(p_0, R)|} \int_{B(p_0, R)} (1 + |Xu|^2) dx \left(\frac{1}{|B(p_0, R/2)|} \int_{B(p_0, R/2)} \omega^{2q/(q-2)} dx \right)^{\frac{q-2}{q}}, \end{aligned}$$

where $q > 2$ is as in (5.2). Since ω is bounded we have $\omega^{2q/(q-2)} \leq C\omega$. Being ω concave, then Jensen inequality and (5.10) yield

$$\begin{aligned}
(5.11) \quad & \frac{1}{|B(p_0, R/2)|} \int_{B(p_0, R/2)} \omega^2 |Xu|^2 dx \\
& \leq C \frac{1}{|B(p_0, R)|} \int_{B(p_0, R)} (1 + |Xu|^2) dx \\
& \quad \times \omega \left(CR^2 + \frac{1}{|B(p_0, R)|} \int_{B(p_0, R)} |u - u_{p_0, R}|^2 dx \right) \\
& \leq C \frac{1}{|B(p_0, R)|} \int_{B(p_0, R)} (1 + |Xu|^2) dx \\
& \quad \times \omega \left(CR^2 \frac{1}{|B(p_0, R)|} \int_{B(p_0, R)} (1 + |Xu|^2) dx \right).
\end{aligned}$$

Similarly we obtain

$$\begin{aligned}
(5.12) \quad & \frac{1}{|B(p_0, R/2)|} \int_{B(p_0, R)} w |Xu|^2 dx \\
& \leq C \frac{1}{|B(p_0, R)|} \int_{B(p_0, R)} (1 + |Xu|^2) dx \left(\frac{1}{|B(p_0, R)|} \int_{B(p_0, R)} |w|^2 dx \right)^{\frac{q-2}{q}},
\end{aligned}$$

where the constant C depends on M as well. At this point we recall (5.8), that together with (5.9)–(5.12) and the Poincaré inequality yields

$$(5.13) \quad \int_{B(p_0, R/2)} |Xw|^2 dp \leq C \chi(p_0, R) \int_{B(p_0, R)} (1 + |Xu|^2) dp.$$

To conclude the proof we use (3.25) and obtain for $0 < r < \frac{R}{2}$,

$$\begin{aligned}
& \int_{B(p_0, r)} (1 + |Xu|^2) dx \leq C \left[\int_{B(p_0, r)} (1 + |Xv|^2) dx + \int_{B(p_0, r)} |Xw|^2 dx \right] \\
& \leq C \left[\left(\frac{r}{R} \right)^Q \int_{B(p_0, R/2)} (1 + |Xv|^2) dx + \int_{B(p_0, R/2)} |Xw|^2 dx \right].
\end{aligned}$$

The conclusion follows at once from the latter, from the argument in the proof of (ii) in Lemma 5.1, and from (5.13). \square

We are now ready to prove the main result of the section.

Theorem 5.4 *Let \mathbf{G} be a Carnot group of step two and $\Omega \subset \mathbf{G}$ an open set. Let $u \in \mathcal{L}_{\text{loc}}^{1,2}(\Omega)$ be a bounded weak solution of (1.3), and assume that (1.4), (1.5) and (5.1) hold. There exists a closed set $S \subset \Omega$ such that u is α -Hölder continuous in $\Omega \setminus S$ for every $0 < \alpha < 1$ and $H_{Q-2-\varepsilon}(S) = 0$ for some $\varepsilon \geq 0$.*

Proof. From Lemma 5.3 one has that for $0 < \theta < 1$ the following estimate holds

$$(5.14) \quad \phi(\theta R) \leq C\theta^2\phi(R)(1 + \chi(p_0, R)\theta^{-Q}).$$

Fix $\alpha \in (0, 1)$. Choose θ in such a way that $2C\theta^{2-2\alpha} - 1$. If, for some $R > 0$ small enough

$$(5.15) \quad \chi(p_0, R) < \theta^Q,$$

then we can iterate (5.14) starting from

$$\phi(\theta R) \leq C\tau^{2\alpha}\phi(R),$$

and obtain

$$\phi(\theta^k R) \leq C\tau^{2k\alpha}\phi(R),$$

which implies

$$\int_{B(p_0, r)} |Xu|^2 dx \leq Cr^{Q-2+2\alpha} \left[\int_{\Omega} (1 + |Xu|^2) dx \right].$$

From the latter and the Poincaré inequality it follows that u is α -Hölder continuous in a neighborhood of p_0 . The singular set of u is characterized by inequality (5.15), from which it is easy to infer that

$$S = \{p_0 \in \Omega \mid \liminf_{R \rightarrow 0^+} R^{2-Q} \int_{B(p_0, R)} |Xu|^2 dp > 0\}.$$

The estimate on the Hausdorff dimension of S can be proved arguing as in Theorem 4.7. \square

Theorem 5.1 holds for coefficients which are merely continuous. When $A_{ij}^{\alpha\beta}$ are Hölder continuous, then the horizontal derivatives of u are Hölder continuous too.

Theorem 5.5 *Let u , Ω and S be as in Theorem 5.4. If $A_{ij}^{\alpha\beta}$ are Hölder continuous with exponent σ , then Xu is locally Hölder continuous with the same Hölder exponent in the set $\Omega \setminus S$.*

Proof. We follow the argument in [GG2]. Since $0 < \sigma < 1$ is the Hölder coefficient of the $A_{ij}^{\alpha\beta}$, then we have

$$(5.16) \quad \omega(r) \leq Cr^{\sigma/2}.$$

Let $\Omega_1 \subset\subset \Omega \setminus S$ and let $d = \frac{1}{2}d(\Omega_1, \partial(\Omega \setminus S))$. Let $p_0 \in \Omega_1$ and $R < \min(1, d)$. Following the argument in Lemma 5.3, we split $u = v + w$ in $B(p_0, R)$.

Using (3.23) we obtain for $0 < r < R$, the estimate

$$(5.17) \quad \int_{B(p_0, r)} |Xu - (Xu)_{p_0, r}|^2 dx \\ \leq C \left[\left(\frac{r}{R} \right)^{Q+2} \int_{B(p_0, R)} |Xu - (Xu)_{p_0, R}|^2 dx + \int_{B(p_0, R)} |Xw|^2 dx \right].$$

From (5.13), (5.16) and as a consequence of (5.15) we obtain

$$\int_{B(p_0, R)} |Xw|^2 dx \leq CR^{Q-2+2\alpha+\alpha\sigma\frac{(q-2)}{q}},$$

for every $0 < \alpha < 1$. Consequently (5.17) yields

$$(5.18) \quad \int_{B(p_0, r)} |Xu - (Xu)_{p_0, r}|^2 dx \\ \leq C \left[\left(\frac{r}{R} \right)^{Q+2} \int_{B(p_0, R)} |Xu - (Xu)_{p_0, R}|^2 dx + R^{Q-2+2\alpha+\alpha\sigma\frac{(q-2)}{q}} \right].$$

If α is close enough to 1, so that $2\alpha + \alpha\sigma\frac{(q-2)}{q} > 2$, then using a standard argument (see Lemma 2.1 in [Gi], Chap. 3) we can readily deduce that Xu is Hölder continuous in Ω_1 with some exponent. In particular, Xu is locally bounded in the same set. Now we want to prove that Xu is locally Hölder continuous in Ω_1 with exponent σ . Using the boundedness of Xu in (5.8) along with the Poincaré inequality, we can prove

$$\int_{B(p_0, R)} |Xw|^2 dx \leq CR^{Q+2\sigma}.$$

The proof is concluded once we introduce the latter inequality in (5.17), and use the iteration lemma 2.1 in [Gi], Chap. 3. \square

For the scalar (diagonal) case the reader should consult [XZ], [C1] and [C2].

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