# STOKES' THEOREM, VOLUME GROWTH AND PARABOLICITY 

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(Received June 28, 2010, revised January 14, 2011)


#### Abstract

We present some new Stokes' type theorems on complete non-compact manifolds that extend, in different directions, previous works by Gaffney and Karp and also the so called Kelvin-Nevanlinna-Royden criterion for $p$-parabolicity. Applications to comparison and uniqueness results involving the $p$-Laplacian are deduced.


1. Introduction. In 1954, Gaffney [4], extended the famous Stokes' theorem to complete $m$-dimensional Riemannian manifolds $M$ by proving that, given a $C^{1}$ vector field $X$ on $M$, we have $\int_{M} \operatorname{div} X=0$ provided $X \in L^{1}(M)$ and $\operatorname{div} X \in L^{1}(M)$ (but in fact $(\operatorname{div} X)_{-}=\max \{-\operatorname{div} X, 0\} \in L^{1}(M)$ is enough). This result was later extended by Karp [13], who showed that the assumption $X \in L^{1}(M)$ can be weakened to

$$
\liminf _{R \rightarrow+\infty} \frac{1}{R} \int_{B_{2 R} \backslash B_{R}}|X| d V_{M}=0 .
$$

Here and on, having fixed a reference origin $o \in M$ on a non-compact manifold $M$, we set $r(x)=\operatorname{dist}_{M}(x, o)$ and we denote by $B_{t}$ and $\partial B_{t}$ the geodesic ball and sphere of radius $t>0$ centered at $o$. Moreover $d V_{M}$ is the Riemannian volume measure on $M$.

It turns out that the completeness of a manifold is analogous to the $p=\infty$ case of $p$-parabolicity, i.e., $M$ is complete if and only if it is $\infty$-parabolic. We recall the concept of $p$-parabolicity. The $p$-Laplacian of a real valued function $u: M \rightarrow \boldsymbol{R}$ is defined by $\Delta_{p} u=\operatorname{div}\left(|\nabla u|^{p-2} \nabla u\right)$. A function $u \in W_{\text {loc }}^{1, p}(M)$ is said to be a $p$-subsolution if $\Delta_{p} u \geq 0$ weakly on $M$. When any bounded above $p$-subsolution is necessarily constant we say that the manifold $M$ is $p$-parabolic. A very useful characterization of (non-) $p$-parabolicity goes under the name of the Kelvin-Nevanlinna-Royden criterion. In the linear setting $p=2$ it was proved in a paper by Lyons and Sullivan [16]. See also Pigola, Rigoli and Setti [18, Theorem 7.27]. The non-linear extension, due to Gol'dshtein and Troyanov [7], states that a manifold $M$ is not $p$-parabolic if and only if there exists a vector field $X$ on $M$ such that (a) $|X| \in L^{p /(p-1)}(M)$, (b) div $X \in L_{\mathrm{loc}}^{1}(M)$ and $(\operatorname{div} X)_{-} \in L^{1}(M)$ (in the weak sense) and (c) $0<\int_{M} \operatorname{div} X \leq+\infty$. In particular this result shows that if $M$ is $p$-parabolic and $X$ is a vector field on $M$ satisfying (a) and (b), then $\int_{M} \operatorname{div} X=0$, thus giving a $p$-parabolic analogue of the Gaffney result. Hence, it is natural to ask whether there exists a $p$-parabolic analogue of the Karp theorem, i.e., if it is possible to weaken the assumptions on the vector field $X$ and still conclude that $\int_{M} \operatorname{div} X d V_{M}=0$.

2000 Mathematics Subject Classification. Primary 35B05; Secondary 31C45.
Key words and phrases. p-parabolicity, Stokes' theorem, Kelvin-Nevanlinna-Royden criterion.
This work is partially supported by GNAMPA-INdAM.

In this paper we will present two different ways to get this result. The first one, Theorem 2.2, is presented in Section 2 and relies on the existence of special exhaustion functions. It has a more theoretical taste and gives the desired $p$-parabolic analogue of Karp's theorem, at least in case either $p=2$ or $p>1$ and $M$ is a model manifold. The second one, Theorem 3.5, is more suitable for explicit applications. It is presented in Section 3 and avoids the parabolicity assumption on $M$ by requiring some connections between the $q$-norm of $X, q=p /(p-1)$, and the volume growth of geodesic balls in $M$. In some sense, specified in Remark 3.8, this result is optimal. In Section 4 we use these techniques to generalize some results involving the $p$-Laplacian comparison and uniqueness theorems on the $p$-harmonic representative in a homotopy class. In particular, in Theorem 4.6, we extend a $p$-Laplacian comparison for vector valued maps on $p$-parabolic manifolds recently obtained in [12]. We point out that the new proof admits $C^{0} \cap W_{\text {loc }}^{1, p}$ maps, instead of the smooth maps of [12]. This is a relevant improvement since, as opposed to the linear setting where 2-harmonic maps are necessarily smooth, for $p \neq 2$ one can at most ensure $p$-harmonic maps to be $C^{1, \alpha}$ [24], [9], [27]. Hence, $C^{1, \alpha}$ seems to be the natural class of functions to consider when dealing with the $p$-Laplacian.

Acknowledgments. We specially thank professors Stefano Pigola, Alberto Setti and Ilkka Holopainen for the useful hints and conversations that helped shape this article. Moreover we would like to thank professor Frank Morgan for suggesting some improvements to the first version of this paper.
2. Exhaustion functions and parabolicity. Given a continuous exhaustion function $f: M \rightarrow \boldsymbol{R}^{+}$in $W_{\text {loc }}^{1, p}(M)$, set

$$
C(r)=f^{-1}[0,2 r) \backslash f^{-1}[0, r) .
$$

Definition 2.1. We say that a vector field $X \in L_{\mathrm{loc}}^{q}(M)$ (with $1 / p+1 / q=1$ ) satisfies the condition $\mathcal{E}_{M, p}$ if:

$$
\begin{equation*}
\left.\left.\left.\left.\liminf _{r \rightarrow \infty} \frac{1}{r}\left|\int_{C(r)}\right| \nabla f\right|^{p} d V_{M}\right|^{1 / p}\left|\int_{C(r)}\right| X\right|^{p /(p-1)} d V_{M}\right|^{(p-1) / p}=0 \tag{1}
\end{equation*}
$$

THEOREM 2.2. Let $f: M \rightarrow \boldsymbol{R}^{+}$be a continuous exhaustion function in $W_{\text {loc }}^{1, p}(M)$. If $X$ is a $L_{\mathrm{loc}}^{q}(M)$ vector field with $(\operatorname{div} X)_{-} \in L^{1}(M), \operatorname{div}(X) \in L_{\mathrm{loc}}^{1}(M)$ in the weak sense and $X$ satisfies the condition $\mathcal{E}_{M, p}$, then $\int_{M} \operatorname{div}(X) d V_{M}=0$.

Proof. Note that $(\operatorname{div} X)_{-} \in L^{1}(M)$ and $\operatorname{div}(X) \in L_{\text {loc }}^{1}(M)$ in the weak sense is the most general hypothesis under which $\int_{M} \operatorname{div} X d V$ is well defined (possibly infinite). For $r>0$, consider the $W^{1, p}$ functions defined by

$$
f_{r}(x):=\max \{\min \{2 r-f(x), r\}, 0\},
$$

i.e., $f_{r}$ is a function identically equal to $r$ on $D(r):=f^{-1}[0, r)$, with the support in $D(2 r)$ and such that $\nabla f_{r}=-\chi_{C(r)} \nabla f$, where $\chi_{C(r)}$ is the characteristic function of $C(r)$. By dominated and monotone convergence we can write

$$
\int_{M} \operatorname{div} X d V_{M}=\lim _{r \rightarrow \infty} \frac{1}{r} \int_{D(2 r)} f_{r} \operatorname{div} X d V_{M} .
$$

Since $f$ is an exhaustion function, $f_{r}$ has a compact support, by definition of weak divergence we get

$$
\begin{aligned}
\int_{M} \operatorname{div} X d V_{M} & =\lim _{r \rightarrow \infty} \frac{1}{r} \int_{D(2 r)}\left\langle\nabla f_{r}, X\right\rangle d V_{M} \\
& \leq \liminf _{r \rightarrow \infty} \frac{1}{r}\left(\int_{C(r)}|\nabla f|^{p} d V_{M}\right)^{1 / p}\left(\int_{C(r)}|X|^{q} d V_{M}\right)^{1 / q}=0 .
\end{aligned}
$$

This proves that $(\operatorname{div} X)_{+}:=\operatorname{div} X+(\operatorname{div} X)_{-} \in L^{1}(M)$ and, by exchanging $X$ with $-X$, the claim follows.

Note that setting $p=\infty$ and $f(x)=r(x)$, one gets exactly the statement of Karp [13].
REMARK 2.3. From the proof of Theorem 2.2, it is easy to see that condition $\mathcal{E}_{M, p}$ can be generalized a little. In fact, if there is a function (without regularity assumptions) $g:(0, \infty) \rightarrow(0, \infty)$ such that $g(t)>t$ and

$$
\left.\left.\left.\left.\liminf _{r \rightarrow \infty} \frac{1}{g(r)-r}\left|\int_{G(r)}\right| \nabla f\right|^{p} d V_{M}\right|^{1 / p}\left|\int_{G(r)}\right| X\right|^{q} d V_{M}\right|^{1 / q}=0,
$$

where $G(r) \equiv f^{-1}[0, g(r)) \backslash f^{-1}[0, r]$, the conclusion of Theorem 2.2 is still valid with the same proof, only needlessly complicated by an awkward notation.

The smaller the value of $\int_{C(r)}|\nabla f|^{p} d V_{M}$ is, the more powerful the conclusion of the theorem is, and since $p$-harmonic functions are in some sense minimizers of the $p$-Dirichlet integral, it is natural to look for such functions as candidates for the role of $f$. Of course, if $M$ is $p$-parabolic it does not admit any positive nonconstant $p$-harmonic function defined on all $M$. Anyway, since we are interested only in the behaviour at infinity of functions and vector fields involved (i.e., the behaviour in $C(r)$ for $r$ large enough), it would be enough to have a $p$-harmonic function $f$ which is defined outside a compact set (inside it could be given any value without changing the conclusions of the theorem).

For example, Sario and Nakai proved that for every 2-parabolic surface $M$, and every relatively compact set $\Omega$, there exists an Evans' potential (see [21, Theorems 12 F and 12G]), i.e., a positive harmonic exhaustion function $E: M \backslash \Omega \rightarrow \boldsymbol{R}^{+}$with $\left.E\right|_{\partial \Omega}=0$ and such that for any $c>0$

$$
\int_{\{E(x) \leq c\}}|\nabla E|^{2} d V_{M} \leq c .
$$

If we let $f=E$ in condition $\mathcal{E}_{M, p}$ (with $p=2$ ), the last inequality allow us to conclude that on a 2-parabolic surface a vector field $X \in L_{\mathrm{loc}}^{2}(M)$ with $(\operatorname{div} X)_{-} \in L^{1}(M)$ and $\operatorname{div}(X) \in$ $L_{\text {loc }}^{1}(M)$ has zero divergence over $M$ provided

$$
\lim _{r \rightarrow \infty} \frac{\int_{\{r \leq E(x) \leq 2 r\}}|X|^{2} d V_{M}}{r}=0
$$

This result is very similar (at least formally) to Karp's, except for the different exponents and for the presence of the Evans potential $E(\cdot)$ that plays the role of the geodesic distance $r(\cdot)$.

It can be proved that an Evans potential exists not only on (2-)parabolic surfaces but also on (2-)parabolic manifolds (see [26] for a complete proof). Unfortunately, no similar existence results have been proved yet in the generic non-linear case ( $p \neq 2$ ). Moreover, even for $p=2$, in general there is no explicit characterization of the function $E$ which can help to estimate its level sets, and therefore the quantity $\int_{\{r \leq E(x) \leq 2 r\}}|X|^{q} d V_{M}$. A very special case is given by the model manifolds (in the sense of Greene and Wu [8]). In this setting the radial function $f: M \backslash B_{1} \rightarrow \boldsymbol{R}$ defined as

$$
\begin{equation*}
f(x):=\int_{1}^{r(x)} \frac{1}{A\left(\partial B_{s}\right)^{1 /(p-1)}} d s \tag{2}
\end{equation*}
$$

is $p$-harmonic and it holds

$$
\begin{equation*}
\int_{B_{r_{2}} \backslash B_{r_{1}}}|\nabla f|^{p}=\int_{r_{1}}^{r_{2}} \frac{1}{A\left(\partial B_{s}\right)^{1 /(p-1)}} d s=f\left(r_{2}\right)-f\left(r_{1}\right) \quad \text { for } r_{2}>r_{1}>1 \tag{3}
\end{equation*}
$$

Here and in what follows $A\left(\partial B_{t}\right)$ stands for the ( $m-1$ )-dimensional Hausdorff measure of $\partial B_{t}$ and is a.e. continuous as a function of $t$. Moreover, the model manifold $M$ is $p$-parabolic if and only if $f(\infty)=\infty$ (see [5] and [25]) and hence in this case $f$ is the Evans' potential we looked for.
3. Stokes' theorem under volume growth assumptions. We now return to consider manifolds $M$ which are not necessarily spherically symmetric. In this general situation, the function $f(r)$ defined in (2) is not $p$-harmonic. Moreover, since $f(+\infty)=+\infty$ is not in general a necessary condition for $p$-parabolicity though it is sufficient (see [25], [20], [11], and Remark 3.2 below), we are not ensured that $f$ is an exhaustion function even for $p$ parabolic $M$. However $f$ is still well defined and relation (3) still holds. Hence, we are led to generalize Theorem 2.2 to generic manifolds, i.e., without parabolicity assumptions, and the generalization we obtain is optimal in some sense (see Theorem 3.5 and Remark 3.8 below). Obviously, in this new result the conclusion will depend on the volume growth of geodesics ball of $M$. Key tools are the estimates from above of the capacity of the condenser ( $\bar{B}_{r_{1}}, B_{r_{2}}$ ) with surface and volume comparisons, as shown in the next proposition. Before that, we briefly recall the definition of $p$-capacity.

Given a Riemannian manifold $M$, let $\Omega$ be a connected domain in $M$ and $D \subset \Omega$ a compact set. For $p \geq 1$, the $p$-capacity of $D$ in $\Omega$ is defined by

$$
\operatorname{Cap}_{p}(D, \Omega):=\inf \left\{\int_{\Omega}|d u|^{p} ; u \in W_{0}^{1, p}(\Omega) \cap C_{0}^{0}(\Omega), u \geq 1 \text { on } D\right\}
$$

It is well known that a manifold $M$ is $p$-parabolic if and only if $\operatorname{Cap}_{p}(D, M)=0$ for every compact set $D \subset M$ (or equivalently for a compact set $D$ with nonvoid interior part) [10]. For $p>1$, we define the functions $a_{p}(t), b_{p}^{\left(r_{1}\right)}:(0,+\infty) \rightarrow(0,+\infty)$ as

$$
a_{p}(t):=A\left(\partial B_{t}\right)^{-1 /(p-1)}, \quad b_{p}^{\left(r_{1}\right)}(t):=\left(\frac{s-r_{1}}{V(t)-V\left(r_{1}\right)}\right)^{1 /(p-1)}
$$

Since the volume $V\left(B_{t}\right)$ of a geodesic ball seen as a function of $t$ is continuous and differentiable almost everywhere with $V^{\prime}\left(B_{t}\right)=A\left(\partial B_{t}\right)$ (see for example [3, Proposition III.3.2]), both functions are a.e. continuous in $(0,+\infty)$, and $b_{p}$ is also differentiable almost everywhere.

In [6, p.3], A. Grigor'yan proves the following inequalities for the $p$-capacity of a spherical condenser, inequalities that link the $p$-parabolicity of a manifold to the area and volume growth of its geodesic balls.

Proposition 3.1. Given a complete Riemannian manifold $M$, the capacity of the condenser ( $B_{r_{1}}, B_{r_{2}}$ ) is bounded from above by

$$
\begin{gather*}
\operatorname{Cap}_{p}\left(\bar{B}_{r_{1}}, B_{r_{2}}\right) \leq\left(\int_{r_{1}}^{r_{2}} a_{p}(t) d t\right)^{1-p},  \tag{4}\\
\operatorname{Cap}_{p}\left(\bar{B}_{r_{1}}, B_{r_{2}}\right) \leq 2^{p}\left(\int_{r_{1}}^{r_{2}} b_{p}^{\left(r_{1}\right)}(t) d t\right)^{1-p} .
\end{gather*}
$$

REmARK 3.2. We observed before that $a_{p} \notin L^{1}(0,+\infty)$ implies $M$ is $p$-parabolic. This is easily obtained by letting $r_{2}$ go to infinity in (4).

REMARK 3.3. In [5], the author proves similar inequalities in the case $p=2$, but with a different proof. This proof can be easily adapted to obtain a better constant in inequality (5), in fact $2^{p}$ can be replaced by $p$.

The functions $a_{p}$ and $b_{p}$ can be used to construct special cut-off functions with controlled $p$-Dirichlet integral. Using these cut-offs, with an argument similar to the one we used in the proof of Theorem 2.2, we get a more suitable and manageable condition on a vector field $X$ in order to guarantee that $\int_{M} \operatorname{div}(X) d V_{M}=0$.

Definition 3.4. We say that a real function $f: M \rightarrow \boldsymbol{R}$ satisfies the condition $\mathcal{A}_{M, p}$ on $M$ for some $p>1$ if there exists a function $g:(0,+\infty) \rightarrow(0,+\infty)$ such that

$$
\begin{equation*}
\liminf _{R \rightarrow \infty}\left(\int_{B_{(R+g(R))} \backslash B_{R}} f d V_{M}\right)\left(\int_{R}^{R+g(G)} a_{p}(s) d s\right)^{-1}=0 \tag{6}
\end{equation*}
$$

The next result gives the announced generalization under volume growth assumption of the Kelvin-Nevanlinna-Royden criterion.

Theorem 3.5. Let $(M,\langle\rangle$,$) be a non-compact Riemannian manifold. Let X$ be a vector field on $M$ such that

$$
\begin{equation*}
\operatorname{div} X \in L_{\mathrm{loc}}^{1}(M) \quad \text { and } \quad \max (-\operatorname{div} X, 0)=(\operatorname{div} X)_{-} \in L^{1}(M) \tag{7}
\end{equation*}
$$

If $|X|^{p /(p-1)}$ satisfies the condition $\mathcal{A}_{M, p}$ on $M$, then

$$
\int_{M} \operatorname{div} X d V_{M}=0
$$

Accordingly, if $X$ is a vector field on $M$ such that $|X|^{p /(p-1)}$ satisfies the condition $\mathcal{A}_{M, p}$, $\operatorname{div} X \in L_{\text {loc }}^{1}(M)$, and $\operatorname{div} X \geq 0$ on $M$, then we must necessarily conclude that $\operatorname{div} X=0$ on $M$. As a matter of fact, even if $\operatorname{div} X \notin L_{\mathrm{loc}}^{1}(M)$, we can obtain a similar conclusion as shown in the next proposition, inspired by [12, Proposition 9].

Proposition 3.6. Let $(M,\langle\rangle$,$) be a non-compact Riemannian manifold. Let X$ be a vector field on $M$ such that

$$
\begin{equation*}
\operatorname{div} X \geq f \tag{8}
\end{equation*}
$$

in the sense of distributions for some $f \in L_{\mathrm{loc}}^{1}(M)$ with $f_{-} \in L^{1}(M)$. If $|X|^{p /(p-1)}$ satisfies condition $\mathcal{A}_{M, p}$ on $M$ for some $p>1$, then

$$
\begin{equation*}
\int_{M} f \leq 0 \tag{9}
\end{equation*}
$$

REMARK 3.7. Combining the following proof with the proof of [12, Proposition 9 and Remark 10], one obtains the validity of (9) when $M$ is $p$-parabolic and $|X|^{p /(p-1)} \in L^{1}(M)$ instead of satisfying $\mathcal{A}_{M, p}$.

Proof. Fix $r_{2}>r_{1}>0$ to be chosen later. Define the functions $\hat{\varphi}=\hat{\varphi}_{r_{1}, r_{2}}: B_{r_{2}} \backslash B_{r_{1}}$ $\rightarrow \boldsymbol{R}$ as

$$
\begin{equation*}
\hat{\varphi}(x):=\left(\int_{r_{1}}^{r_{2}} a_{p}(s) d s\right)^{-1} \int_{r(x)}^{r_{2}} a_{p}(s) d s \tag{10}
\end{equation*}
$$

and let $\varphi=\varphi_{r_{1}, r_{2}}: M \rightarrow \boldsymbol{R}$ be defined as

$$
\varphi(x):= \begin{cases}1 & r(x)<r_{1}  \tag{11}\\ \hat{\varphi}(x) & r_{1} \leq r(x) \leq r_{2} \\ 0 & r_{2}<r(x)\end{cases}
$$

A straightforward calculation yields

$$
\int_{M}|\nabla \varphi|^{p} d V_{M}=\int_{B_{r_{2}} \backslash B_{r_{1}}}|\nabla \hat{\varphi}|^{p} d V_{M}=\left(\int_{r_{1}}^{r_{2}} a_{p}(s) d s\right)^{1-p}
$$

By standard density results we can use $\varphi \in W_{0}^{1, p}(M)$ as a test function in the weak relation
(8). Thus we obtain

$$
\begin{align*}
\int_{M} \varphi f d V_{M} & \leq(\operatorname{div} X, \varphi)  \tag{12}\\
& =-\int_{M}\langle X, \nabla \varphi\rangle \\
& \leq\left(\int_{\operatorname{supp}(\nabla \varphi)}|X|^{p /(p-1)}\right)^{(p-1) / p}\left(\int_{M}|\nabla \varphi|^{p}\right)^{1 / p} \\
& \leq\left\{\left(\int_{B_{r_{2}} \backslash B_{r_{1}}}|X|^{p /(p-1)}\right)\left(\int_{r_{1}}^{r_{2}} a_{p}(s) d s\right)^{-1}\right\}^{(p-1) / p},
\end{align*}
$$

where we have applied Hölder in the next-to-last inequality. Now, let $\left\{R_{k}\right\}_{k=1}^{\infty}$ be a sequence such that $R_{k} \rightarrow \infty$, which realizes the lim inf in condition (6). Up to passing to a subsequence, we can suppose $R_{k+1} \geq R_{k}+g\left(R_{k}\right)$. Hence, the sequence of cut-offs $\varphi_{k}:=\varphi_{R_{k}, R_{k}+g\left(R_{k}\right)}$ converges monotonically to 1 and applying monotone and dominated convergence, we have

$$
\lim _{k \rightarrow \infty} \int_{M} \varphi_{k} f=\lim _{k \rightarrow \infty} \int_{M} \varphi_{k} f_{+}-\lim _{k \rightarrow \infty} \int_{M} \varphi_{k} f_{-}=\int_{M} f_{+}-\int_{M} f_{-}=\int_{M} f
$$

Taking limits as $k \rightarrow \infty$ in inequality (12), assumption $\mathcal{A}_{M, p}$ finally gives

$$
\int_{M} f \leq 0 .
$$

Proof of Theorem 3.5. Choosing $f=\operatorname{div} X$ in Proposition 3.6 we get $\int_{M} \operatorname{div} X \leq$ 0 and $(\operatorname{div} X)_{+} \in L^{1}(M)$. Hence, we can repeat the proof replacing $X$ with $-X$.

REMARK 3.8. We point out that one could easily obtain results similar to Theorem 3.5 and Proposition 3.6 replacing $\varphi_{r_{1}, r_{2}}$ in the proofs with standard cut-off functions $0 \leq \xi_{r_{1}, r_{2}} \leq$ 1 defined for any $\varepsilon>0$ in such a way that

$$
\xi_{r_{1}, r_{2}} \equiv 1 \text { on } B_{r_{1}}, \quad \xi_{r_{1}, r_{2}} \equiv 0 \text { on } M \backslash B_{r_{2}}, \quad\left|\nabla \xi_{r_{1}, r_{2}}\right| \leq \frac{1+\varepsilon}{r_{2}-r_{1}} .
$$

Nevertheless $\varphi_{r_{1}, r_{2}}$ gives better results than the standard cut-offs. For example, consider a 2 -dimensional model manifold with the Riemannian metric

$$
d s^{2}=d t^{2}+g^{2}(t) d \theta^{2}
$$

where $g(t)=e^{-t}$ outside a neighborhood of 0 . Then the $p$-energy of $\varphi_{r_{1}, r_{2}}$ and $\xi_{r_{1}, r_{2}}$ are respectively

$$
\begin{gathered}
\int_{M}\left|\nabla \varphi_{r_{1}, r_{2}}\right|^{p} d v_{M}=\left(\int_{r_{1}}^{r_{2}} a_{p}(s) d s\right)^{1-p}=2 \pi(p-1)^{1-p}\left(e^{r_{2} /(p-1)}-e^{r_{1} /(p-1)}\right)^{1-p} \\
\int_{M}\left|\nabla \xi_{r_{1}, r_{2}}\right|^{p} d v_{M} \leq\left(\frac{1+\varepsilon}{r_{2}-r_{1}}\right)^{p} \int_{r_{1}}^{r_{2}} A\left(\partial B_{s}\right) d s=2 \pi\left(\frac{1+\varepsilon}{r_{2}-r_{1}}\right)^{p}\left(e^{-r_{1}}-e^{-r_{2}}\right) .
\end{gathered}
$$

If we choose $r_{2}=2 r_{1} \equiv 2 r$ and let $r \rightarrow \infty$, we get

$$
\begin{gathered}
\int_{M}\left|\nabla \varphi_{r, 2 r}\right|^{p} d v_{M} \sim c e^{-r}\left(e^{r /(p-1)}-1\right)^{1-p} \sim c e^{-2 r} \\
\int_{M}\left|\nabla \xi_{r, 2 r}\right|^{p} d v_{M} \sim \frac{c^{\prime}}{r^{p}} e^{-r}\left(1-e^{-r}\right) \sim \frac{c^{\prime}}{r^{p}} e^{-r}
\end{gathered}
$$

for some positive constants $c, c^{\prime}$. Using $\varphi_{r, 2 r}$ in the proof of Proposition 3.6 (in particular in inequality (12)), we can conclude that $f=0$ provided

$$
\lim _{r \rightarrow \infty}\left(\int_{B_{2 r \backslash B_{r}}}|X|^{p /(p-1)}\right)^{(p-1) / p} e^{-2 r}=0,
$$

while using $\xi_{r, 2 r}$ we get a weaker result, i.e., $f=0$ provided

$$
\lim _{r \rightarrow \infty}\left(\int_{B_{2 r \backslash B_{r}}}|X|^{p /(p-1)}\right)^{(p-1) / p} \frac{1}{r^{p}} e^{-r}=0 .
$$

One could ask if there exist even better cut-offs than the ones we chose. First, note that restricting to model manifolds, the cut-offs $\varphi_{r_{1}, r_{2}}$ are $p$-harmonic on $B_{r_{2}} \backslash B_{r_{1}}$, and so their $p$-energy is minimal. In general this is not true. Anyway, it turns out that the functions $\varphi_{r_{1}, r_{2}}$ are optimal at least in the class of radial functions, in fact they minimize the $p$-energy in this class, and this makes the condition $\mathcal{A}_{M, p}$ sharp. To prove this fact, consider any radial cut-off $\psi:=\psi_{r_{1}, r_{2}}$ satisfying $\psi \equiv 1$ on $B_{r_{1}}, \psi \equiv 0$ on $M \backslash B_{r_{2}}$. By Jensen's inequality we have

$$
\begin{aligned}
\int\left|\nabla \varphi_{r_{1}, r_{2}}\right|^{p} d v_{M} & =\left(\int_{r_{1}}^{r_{2}} a_{p}(s) d s\right)^{1-p}=c_{\psi}^{1-p}\left(\int_{r_{1}}^{r_{2}} \frac{a_{p}(s)}{\left|\psi^{\prime}(s)\right|} \frac{\left|\psi^{\prime}(s)\right| d s}{c_{\psi}}\right)^{1-p} \\
& \leq c_{\psi}^{-p} \int_{r_{1}}^{r_{2}}\left|\psi^{\prime}(s)\right|^{p} A\left(\partial B_{s}\right) d s \leq \int|\nabla \psi|^{p} d v_{M}
\end{aligned}
$$

where $\psi^{\prime}$ is the radial derivative of $\psi$ and $c_{\psi}=\int_{r_{1}}^{r_{2}}\left|\psi^{\prime}(s)\right| d s \geq 1$.
4. Applications. Theorem 3.5, and Proposition 3.6, can be naturally applied to those situations where the standard Kelvin-Nevanlinna-Royden criterion is used to deduce information on $p$-parabolic manifolds. First, we present a global comparison result for the $p$ Laplacian of real valued functions. The original result assuming $p$-parabolicity appears in [12, Theorem 1].

THEOREM 4.1. Let $(M,\langle\rangle$,$) be a connected, non-compact Riemannian manifold. As-$ sume that $u, v \in W_{\text {loc }}^{1, p}(M) \cap C^{0}(M), p>1$, satisfy

$$
\Delta_{p} u \geq \Delta_{p} v \quad \text { weakly on } M
$$

and that $|\nabla u|^{p}$ and $|\nabla v|^{p}$ satisfy the condition $\mathcal{A}_{M, p}$ on $M$. Then, $u=v+A$ on $M$, for some constant $A \in \boldsymbol{R}$.

Choosing a constant function $v$, we immediately deduce the following result, which generalize [17, Corollary 3].

Corollary 4.2. Let $(M,\langle\rangle$,$) be a connected, non-compact Riemannian manifold.$ Assume that $u \in W_{\mathrm{loc}}^{1, p}(M) \cap C^{0}(M), p>1$, is a weak $p$-subharmonic function on $M$ such that $|\nabla u|^{p}$ satisfies the condition $\mathcal{A}_{M, p}$ on $M$. Then $u$ is constant.

In [22], Schoen and Yau considered the problem of uniqueness of the 2-harmonic representative with finite energy in a (free) homotopy class of maps from a complete manifold $M$ of finite volume to a complete manifold $N$ of non-positive sectional curvature. In particular, they obtained that if the sectional curvature of $N$ is negative then a given harmonic map $u$ is unique in its homotopy class, unless $u(M)$ is contained in a geodesic of $N$. Moreover, if $\operatorname{Sect}_{N} \leq 0$ and two homotopic harmonic maps $u$ and $v$ with finite energy are given, then $u$ and $v$ are homotopic through harmonic maps. In [17], Pigola, Setti and Rigoli noticed that the
assumption $V(M)<\infty$ can be replaced by asking $M$ to be 2-parabolic. In Schoen and Yau's result, the finite energy of the maps is used in two fundamental steps of the proof:
(1) to prove that a particular subharmonic map of finite energy is constant;
(2) to construct the homotopy via harmonic maps.

Using the $p=2$ case of Corollary 4.2, we can deal with step (1) and thus obtain the following theorem. If $\operatorname{Sect}_{N} \leq 0$, weakening finite energy assumption in step (2) does not seem trivial to us, but we can still get a result for maps with fast $p$-energy decay without parabolicity assumption.

Theorem 4.3. Suppose $M$ and $N$ are complete manifolds.

1) Suppose $\operatorname{Sect}_{N}<0$. Let $u: M \rightarrow N$ be a harmonic map such that $|\nabla u|^{2}$ satisfies the condition $\mathcal{A}_{M, 2}$ on $M$. Then there is no other harmonic map homotopic to $u$ satisfying the condition $\mathcal{A}_{M, 2}$ unless $u(M)$ is contained in a geodesic of $N$.
2) Suppose $\operatorname{Sect}_{N} \leq 0$. Let $u, v: M \rightarrow N$ be homotopic harmonic maps such that $|\nabla u|^{2},|\nabla v|^{2} \in L^{1}(M)$ satisfy the condition $\mathcal{A}_{M, 2}$ on $M$. Then there is a smooth one parameter family $u_{t}: M \rightarrow N$ for $t \in[0,1]$ of harmonic maps with $u_{0}=u$ and $u_{1}=v$. Moreover, for each $x \in M$, the curve $\left\{u_{t}(x) ; t \in[0,1]\right\}$ is a constant (independent of $x$ ) speed parametrization of a geodesic.

REMARK 4.4. While the existence in a homotopy class of a harmonic representative with finite energy is ensured by a further result by Schoen and Yau [23], in the setting of Theorem 4.3 we are not able to guarantee that there exists at least one harmonic map whose energy satisfies $\mathcal{A}_{M, 2}$.

An interesting task is to extend Schoen and Yau's uniqueness results to the nonlinear $(p \neq 2)$ setting. In [17], the authors take a first step in this direction by proving that a map $u: M \rightarrow N$ with finite $p$-energy and homotopic to a constant is constant provided $M$ is $p$-parabolic and $\operatorname{Sect}_{N} \leq 0$. Using Theorem 3.5 in the proof of their result, we easily obtain the following

ThEOREM 4.5. Let $\left(M,\langle,\rangle_{M}\right)$ and $\left(N,\langle,\rangle_{N}\right)$ be complete Riemannian manifolds. Assume that $M$ is non-compact and that $N$ has non-positive sectional curvatures. If $u: M \rightarrow N$ is a p-harmonic map homotopic to a constant and with energy density $|d u|^{p}$ satisfying the condition $\mathcal{A}_{M, p}$, then $u$ is a constant map.

In [12], the authors apply the Kelvin-Nevanlinna-Royden criterion to obtain a vector valued version of their comparison theorem. In some sense this result is a further step in treating the problem of the uniqueness of $p$-harmonic representative. In particular, if $M$ is $p$-parabolic and $u, v: M \rightarrow \boldsymbol{R}^{n}$ are $C^{\infty}$ maps satisfying $\Delta_{p} u=\Delta_{p} v$ and $|d u|,|d v| \in L^{p}$, then $u=v+A$ for some constant vector $A \in \boldsymbol{R}^{n}$. To prove it, they construct a vector field $X$ depending on $d u$ and $d v$, whose divergence is such that

$$
\begin{equation*}
0 \neq(\operatorname{div} X)_{-} \leq C\left(|d u|^{p}+|d v|^{p}\right) \in L^{1} . \tag{13}
\end{equation*}
$$

As a matter of fact, in their proof, $X$ is defined in such a way that the smoothness of $u$ and $v$ seems to be strictly necessary to do computations. In order to generalize their result, in the direction of Proposition 3.6, apparently assumption (13) can not be dropped. Nevertheless, in case $|d u|,|d v| \in L^{p}$ and their $L^{p}$-norms decay fast with respect to the volume in the $\mathcal{A}_{M, p}$ sense, we obtain a similar result for not $p$-parabolic manifolds and for maps with low regularity.

THEOREM 4.6. Suppose that $(M,\langle\rangle$,$) is a complete non-compact Riemannian mani-$ fold. For $p>1$, let $u, v: M \rightarrow \boldsymbol{R}^{n}$ be $C^{0} \cap W_{\text {loc }}^{1, p}(M)$ maps satisfying

$$
\begin{equation*}
\Delta_{p} u=\Delta_{p} v \quad \text { on } M \tag{14}
\end{equation*}
$$

in the sense of distributions on $M$ and

$$
|d u|,|d v| \in L^{p}(M)
$$

Suppose either $M$ is p-parabolic or $|d u|^{p}$ and $|d v|^{p}$ satisfy the condition $\mathcal{A}_{M, p}$ on $M$. Then $u=v+C$, for some constant $C \in \boldsymbol{R}^{n}$.

REMARK 4.7. In Proposition 3.1, we saw that the capacity of a condenser $\left(B_{r_{1}}, B_{r_{2}}\right)$ can be estimated from above using either the behaviour of $V\left(B_{s}\right)$ or the behaviour of $A\left(\partial B_{s}\right)$. This suggests that the condition $\mathcal{A}_{M, p}$ should have an analogue in which $A\left(\partial B_{s}\right)$ is replaced by $V\left(B_{s}\right)$. This fact is useful since it is usually easier to verify and to handle volume growth assumptions than area growth conditions.

DEFINITION 4.8. We say that a real function $f: M \rightarrow \boldsymbol{R}$ satisfies the condition $\mathcal{V}_{M, p}$ on $M$ for some $p>1$ if there exists a function $g:(0,+\infty) \rightarrow(0,+\infty)$ such that

$$
\begin{equation*}
\liminf _{R \rightarrow \infty}\left(\int_{B_{(R+g(R))} \backslash B_{R}} f d V_{M}\right)\left(\int_{R}^{R+g(R)}\left(\frac{t}{V(t)}\right)^{1 /(p-1)} d s\right)^{-1}=0 \tag{15}
\end{equation*}
$$

Indeed, it turns out that in every proposition stated in Section 4 , the condition $\mathcal{A}_{M, p}$ can be replaced by the condition $\mathcal{V}_{M, p}$, and the proofs remain almost the same.
5. Proof of Theorem 4.6. We can now proceed to prove Theorem 4.6. Recall that, by definition, (14) holds in the sense of distributions on $M$ if

$$
\left.\left.\int_{M}\langle\eta,| d v\right|^{p-2} d v-|d u|^{p-2} d u\right\rangle_{H S}=0
$$

for every compactly supported $\eta \in T^{*} M \otimes \boldsymbol{R}^{n}$. Here $\langle,\rangle_{H S}$ stands for the Hilbert-Schmidt scalar product on $T^{*} M \otimes \boldsymbol{R}^{n}$. Moreover, we mention the following lemma, derived by a basic inequality by Lindqvist [15], which will be useful later.

LEMMA 5.1 ([12, Corollary 5]). Let $(V,\langle\rangle$,$) be a finite dimensional real vector space$ endowed with a positive definite scalar product and let $p>1$. Then, for every $x, y \in V$, it holds

$$
\begin{equation*}
\left.\left.\langle | x\right|^{p-2} x-|y|^{p-2} y, x-y\right\rangle \geq 2 C(p) \Psi(x, y) \tag{16}
\end{equation*}
$$

where

$$
\Psi(x, y):= \begin{cases}|x-y|^{p} & p \geq 2 \\ |x-y|^{2} /(|x|+|y|)^{2-p} & 1<p<2\end{cases}
$$

and $C(p)$ is a positive constant depending only on $p$.
Proof of Theorem 4.6. First, we assume that $M$ is $p$-parabolic. We suppose that at least one of $u$ or $v$ is non-constant, for otherwise there is nothing to prove. Fix $q_{0} \in M$ and set $C:=u\left(q_{0}\right)-v\left(q_{0}\right) \in \boldsymbol{R}^{n}$. Replacing $v$ with $\tilde{v}:=v+C$ if necessary, we can suppose $C=0$. Introduce the radial function $r: \boldsymbol{R}^{n} \rightarrow \boldsymbol{R}$ defined as $r(x)=|x|$. For $A>1$, consider the weakly differentiable vector field $X_{A}$ defined as

$$
X_{A}(x):=\left[\left.d h_{A}\right|_{(u-v)(x)} \circ\left(|d u(x)|^{p-2} d u(x)-|d v(x)|^{p-2} d v(x)\right)\right]^{\sharp}, \quad x \in M,
$$

where $h_{A} \in C^{\infty}\left(\boldsymbol{R}^{n}, \boldsymbol{R}\right)$ is the function

$$
h_{A}(y):=\sqrt{A+r^{2}(y)}
$$

and $\sharp$ denotes the usual musical isomorphism defined by $\left\langle\omega^{\sharp}, V\right\rangle=\omega(V)$ for all differential 1forms $\omega$ and vector fields $V$. We observe that $X_{A}$ is well defined since there exists a canonical identification

$$
T_{(u-v)(q)} \boldsymbol{R}^{n} \cong T_{u(q)} \boldsymbol{R}^{n} \cong T_{v(q)} \boldsymbol{R}^{n} \cong \boldsymbol{R}^{n}
$$

Compute

$$
d h_{A}=\frac{d r^{2}}{2 \sqrt{A+r^{2}}}
$$

and observe that, because of the special structure of $\boldsymbol{R}^{n}$, for each vector field $Y$ on $\boldsymbol{R}^{n}$ it holds

$$
\left.\left(d r^{2}\right)\right|_{(u-v)(x)}(Y)=2\langle(u-v)(x), Y\rangle_{\boldsymbol{R}^{n}}
$$

By definition of weak divergence, for each test function $0 \leq \phi \in C_{c}^{\infty}(M)$, we have

$$
\begin{aligned}
& -\left(\operatorname{div} X_{A}, \phi\right)=\int_{M}\left\langle X_{A}, \nabla \phi\right\rangle_{M} \\
& \quad=\int_{M}\left\langle\left[\left.d h_{A}\right|_{(u-v)(x)} \circ\left(|d u(x)|^{p-2} d u(x)-|d v(x)|^{p-2} d v(x)\right)\right]^{\sharp}, \nabla \phi\right\rangle_{M} \\
& \quad=\left.\int_{M} \frac{d r^{2}}{2 \sqrt{A+r^{2}}}\right|_{(u-v)(x)} \circ\left(\left.|d u(x)|^{p-2} d u\right|_{x}-\left.|d v(x)|^{p-2} d v\right|_{x}\right)(\nabla \phi) \\
& \quad=\int_{M} \frac{1}{\sqrt{A+r^{2}(u-v)(x)}}\left|(u-v)(x),\left(\left.|d u(x)|^{p-2} d u\right|_{x}-\left.|d v(x)|^{p-2} d v\right|_{x}\right)(\nabla \phi(x))\right\rangle_{R^{n}} .
\end{aligned}
$$

Since $u, v \in W_{\text {loc }}^{1, p}(M)$, assumption (14) implies that

$$
\begin{align*}
0 & \left.=\left.\int_{M}\left\langle d\left(\frac{(u-v) \phi}{\sqrt{A+r^{2}(u-v)}}\right),\right| d u\right|^{p-2} d u-|d v|^{p-2} d v\right\rangle_{H S}  \tag{17}\\
& \left.=\left.\int_{M} \frac{1}{\sqrt{A+r^{2}(u-v)}}\langle d \phi \otimes(u-v),| d u\right|^{p-2} d u-|d v|^{p-2} d v\right\rangle_{H S}
\end{align*}
$$

$$
\begin{aligned}
& +\int_{M} \frac{\phi}{\sqrt{A+r^{2}(u-v)}}\left|d u-d v,|d u|^{p-2} d u-|d v|^{p-2} d v\right|_{H S} \\
& -\int_{M} \frac{\phi}{2\left(A+r^{2}(u-v)\right)^{3 / 2}} \\
& \times\left.\left\langle\left. d r^{2}\right|_{(u-v)} \circ(d u-d v) \otimes(u-v),\right| d u\right|^{p-2} d u-\left.|d v|^{p-2} d v\right|_{H S} \\
\geq & -\left(\operatorname{div} X_{A}, \phi\right) \\
& +\int_{M} \frac{2 C(p) \phi}{\sqrt{A+r^{2}(u-v)}} \Psi(d u, d v) \\
& -\int_{M} \frac{\phi r^{2}(u-v)}{\left(A+r^{2}(u-v)\right)^{3 / 2}}(|d u|+|d v|)\left(|d u|^{p-1}+|d v|^{p-1}\right),
\end{aligned}
$$

where we have used Lemma 5.1 for the second term and Cauchy-Schwarz inequality for the third one. Setting

$$
f_{A}:=\frac{2 C(p)}{\sqrt{A+r^{2}(u-v)}} \Psi-\frac{2 r^{2}(u-v)}{\left(A+r^{2}(u-v)\right)^{3 / 2}}\left(|d u|^{p}+|d v|^{p}\right),
$$

by Young's inequality, (17) gives

$$
\begin{equation*}
\operatorname{div} X_{A} \geq f_{A} \tag{18}
\end{equation*}
$$

in the sense of distributions.
Let us now compute the $L^{p /(p-1)}$-norm of $X_{A}$. Since

$$
\begin{aligned}
\left||d u|^{p-2} d u-|d v|^{p-2} d v\right|^{p /(p-1)} & \leq\left(|d u|^{p-1}+|d v|^{p-1}\right)^{p /(p-1)} \\
& \leq 2^{1 /(p-1)}\left(|d u|^{p}+|d v|^{p}\right),
\end{aligned}
$$

we have

$$
\begin{aligned}
\left|X_{A}\right|^{p /(p-1)} & =\left.\left|\sqrt{\frac{r^{2}(u-v)}{A+r^{2}(u-v)}}\right|^{p /(p-1)}| | d u\right|^{p-2} d u-\left.|d v|^{p-2} d v\right|^{p /(p-1)} \\
& \leq 2^{1 /(p-1)}\left(|d u|^{p}+|d v|^{p}\right) \in L^{1}(M) .
\end{aligned}
$$

Hence $X_{A}$ is a weakly differentiable vector field with $\left|X_{A}\right| \in L^{p /(p-1)}(M)$.
To apply Proposition 3.6 (in the $p$-parabolic version pointed out in Remark 3.7), it remains to show that $\left(f_{A}\right)_{-} \in L^{1}(M)$. To this purpose, we note that

$$
\begin{align*}
\left(f_{A}\right)_{-} & \leq \frac{2 r^{2}(u-v)}{\left(A+r^{2}(u-v)\right)^{3 / 2}}\left(|d u|^{p}+|d v|^{p}\right)  \tag{19}\\
& \leq \frac{r^{2}(u-v)}{A+r^{2}(u-v)} \frac{2}{\sqrt{A+r^{2}(u-v)}}\left(|d u|^{p}+|d v|^{p}\right) \\
& \leq \frac{2}{\sqrt{A}}\left(|d u|^{p}+|d v|^{p}\right) \in L^{1}(M) .
\end{align*}
$$

Then, the assumptions of Proposition 3.6 are satisfied and we get, for every $A>1$,

$$
\begin{align*}
0 & \geq \int_{M} f_{A}  \tag{20}\\
& =\int_{M}\left[\frac{2 C(p)}{\sqrt{A+r^{2}(u-v)}} \Psi-\frac{2 r^{2}(u-v)}{\left(A+r^{2}(u-v)\right)^{3 / 2}}\left(|d u|^{p}+|d v|^{p}\right)\right] .
\end{align*}
$$

Fix $T>0$ and define

$$
M_{T}:=\{x \in M ; r(u-v)(x) \leq T\} \quad \text { and } \quad M^{T}:=M \backslash M_{T} .
$$

Then, we can write (20) as
(21) $0 \geq \int_{M^{T}} f_{A}+\int_{M_{T}} \frac{2 C(p)}{\sqrt{A+r^{2}(u-v)}} \Psi-\int_{M_{T}} \frac{2 r^{2}(u-v)}{\left(A+r^{2}(u-v)\right)^{3 / 2}}\left(|d u|^{p}+|d v|^{p}\right)$.

Note that

$$
\begin{equation*}
\int_{M^{T}} f_{A} \geq-\int_{M^{T}} \frac{2}{\sqrt{A+T^{2}}}\left(|d u|^{p}+|d v|^{p}\right)=-\frac{2}{\sqrt{A+T^{2}}} \int_{M^{T}}\left(|d u|^{p}+|d v|^{p}\right) . \tag{22}
\end{equation*}
$$

On the other hand, to deal with $\int_{M_{T}} f_{A}$, observe that

$$
\begin{equation*}
\int_{M_{T}} \frac{2 C(p)}{\sqrt{A+r^{2}(u-v)}} \Psi \geq \frac{2 C(p)}{\sqrt{A+T^{2}}} \int_{M_{T}} \Psi \tag{23}
\end{equation*}
$$

Furthermore, the real function $t \mapsto 2 t /(A+t)^{3 / 2}$ has a global maximum at $t=2 A$ and is increasing in $(0,2 A)$. Hence, up to choosing $A>T^{2} / 2$, we have also

$$
\begin{equation*}
\int_{M_{T}} \frac{2 r^{2}(u-v)}{\left(A+r^{2}(u-v)\right)^{3 / 2}}\left(|d u|^{p}+|d v|^{p}\right) \leq \frac{2 T^{2}}{\left(A+T^{2}\right)^{3 / 2}} \int_{M_{T}}\left(|d u|^{p}+|d v|^{p}\right) . \tag{24}
\end{equation*}
$$

Inserting (22), (23) and (24) in (21), we get

$$
\begin{aligned}
\frac{2 C(p)}{\sqrt{A+T^{2}}} \int_{M_{T}} \Psi \leq & \frac{2 T^{2}}{\left(A+T^{2}\right)^{3 / 2}} \int_{M_{T}}\left(|d u|^{p}+|d v|^{p}\right) \\
& +\frac{2}{\sqrt{A+T^{2}}} \int_{M^{T}}\left(|d u|^{p}+|d v|^{p}\right)
\end{aligned}
$$

which gives

$$
C(p) \int_{M_{T}} \Psi \leq \int_{M^{T}}\left(|d u|^{p}+|d v|^{p}\right)+\frac{T^{2}}{\sqrt{A+T^{2}}} \int_{M_{T}}\left(|d u|^{p}+|d v|^{p}\right),
$$

for all $A>\max \left\{1, T^{2} / 2\right\}$. Letting $A \rightarrow+\infty$, this latter yields

$$
\begin{equation*}
C(p) \int_{M_{T}} \Psi \leq \int_{M^{T}}\left(|d u|^{p}+|d v|^{p}\right) \tag{25}
\end{equation*}
$$

Since $\left(|d u|^{p}+|d v|^{p}\right) \in L^{1}(M)$, for $T \rightarrow \infty$ we can apply respectively dominated convergence on the right-hand side and monotone convergence on the left-hand side of (25), and we get

$$
C(p) \int_{M} \Psi=0
$$

which in turn gives $|d(u-v)| \equiv 0$ on $M$, that is, $u-v \equiv u\left(q_{0}\right)-v\left(q_{0}\right)=C$ on $M$. This conclude the first part of the proof.

If $M$ is not $p$-parabolic, proceeding as above, the condition $\mathcal{A}_{M, p}$ and Proposition 3.6 give (20). From there on, we can repeat the proof of the $p$-parabolic case.
6. Final remarks. So far, we applied Theorem 3.5 and Proposition 3.6 to generalize the results for which originally a $p$-parabolicity assumption on the domain manifold was used to apply the Kelvin-Nevanlinna-Royden criterion and deduce a Stokes' type conclusion. In particular we showed that the $p$-parabolicity can be replaced by suitable volume growth estimates. As a matter of fact, the functions $\varphi_{r_{1}, r_{2}}$ defined in the proof of Proposition 3.1 seem to naturally appear each time when a control on the $L^{p}$-norm of the gradient of a cutoff function is required. For instance, consider a $m$-dimensional manifold $M$ supporting an Euclidean type Sobolev inequality, i.e.,

$$
\begin{equation*}
\|\eta\|_{p}^{q} \leq S_{M}^{q}\|\nabla \eta\|_{q}^{q} \quad \text { for all } \eta \in C_{c}^{\infty}(M) \tag{26}
\end{equation*}
$$

holds for some positve constant $S_{M}$ and for some $1<q<m$ and $p=m q /(m-q)$. For such manifolds, Carron [2], proved that there is an almost Euclidean lower bound for the volume growth of the geodesic ball. Namely, there exists an explicit positive $\gamma>0$ depending on $m$ and $q$ such that

$$
\begin{equation*}
V\left(B_{r}\right) \geq \gamma r^{m} \tag{27}
\end{equation*}
$$

for all $r>0$ (but to the best of our knowledge no lower control on $A\left(\partial B_{r}\right)$ is given). Moreover, Cao, Shen and Zhu [1] and Li and Wang [14] observed that the validity of (26) for $q=2$ implies $M$ is 2-hyperbolic (see also [19] for the $q \neq 2$ case). Observe that, by a standard density argument, inequality (26) holds for all $\eta \in W_{0}^{1, q}(M)$. Hence we can choose $\eta=\varphi_{r_{1}, r_{2}}$ for some $0<r_{1}<r_{2}$ obtaining

$$
\begin{equation*}
V\left(B_{r_{1}}\right)^{q / p} \leq \frac{S_{M}^{q}}{\left(\int_{r_{1}}^{r_{2}} a_{q}(s) d s\right)^{q-1}} \tag{28}
\end{equation*}
$$

In particular, letting $r_{2} \rightarrow \infty$ gives $a_{q} \in L^{1}(0,+\infty)$, since otherwise $V\left(B_{r_{1}}\right) \equiv 0$ for all $r_{1}>0$. Even if this conclusion is immediately implied by the $q$-hyperbolicity of $M$, we can combine inequality (28) with Carron's estimate (27) to obtain a slightly improved result.

PROPOSITION 6.1. Let $M$ be an m-dimensional complete non-compact manifold supporting the euclidean Sobolev inequality (26) for some $q<p$ and $p=m q /(m-q)$. Then there exists a positive constant $0<C_{S}=\gamma^{-(m-q) /(m(q-1))} S_{M}^{q /(q-1)}$ such that

$$
\begin{equation*}
r^{(m-q) /(q-1)} \int_{r}^{\infty} a_{q}(s) d s \leq C_{S} \quad \text { for all } r>0 \tag{29}
\end{equation*}
$$

REMARK 6.2. We underline that Proposition 6.1 is non-trivial, in the sense that, in the absence of the validity of (26), there exist manifolds satisfying (27) and $a_{q} \in L^{1}(0,+\infty)$, for which (29) does not hold. For instance, for fixed $1<q<(m+1) / 2$, an example is given by the $m$-dimensional model manifold $(t, \theta) \in M=(0,+\infty) \times S^{m-1}$ endowed with
the Riemannian metric $\langle,\rangle_{M}=d t^{2}+h^{2}(t) d \theta^{2}$, with warping function $h \in C^{\infty}((0,+\infty))$ chosen such that

$$
\begin{cases}h(0)=0, \quad h^{\prime}\left(0^{+}\right)=1, \\ h(t) \geq t^{\beta}, & \text { for } k=0,1,2, \ldots \\ \left.h(t)\right|_{[4 k+3,4 k+4]} \equiv t^{\beta} & \text { for } k=0,1,2, \ldots \\ \left.h(t)\right|_{[4 k+1,4 k+2]} \equiv H t & \end{cases}
$$

for some constants

$$
\frac{q-1}{m-1}<\beta<\frac{m-q}{m-1}<1 \quad \text { and } \quad H>\frac{1}{4}\left(\frac{m 10^{m} \gamma}{A\left(\partial \boldsymbol{B}_{1}^{R^{m}}\right)}\right)^{1 /(m-1)}
$$

## References

[1] H.-D. CAO, Y. Shen and S. ZhU, The structure of stable minimal hypersurfaces in $\boldsymbol{R}^{n+1}$, Math. Res. Lett. 4 (1997), 637-644.
[2] G. CARRON, Inégalités isopérimétriques et inégalités de Faber-Krahn. Séminaire de Théorie Spectrale et Géométrie, No. 13, Année 1994-1995, 63-66.
[ 3 ] I. CHAVEL, Riemannian geometry: A modern introduction, Cambridge Tracts in Math. 108, Cambridge University Press, Cambridge, 1993.
[4] M. Gaffney, A special Stokes' Theorem for complete Riemannian manifolds, Ann. of Math. 60 (1954), 140-145.
[5] A. GRIGOR'YAN, Analytic and geometric background of recurrence and non-explosion of the Brownian motion on Riemannian manifolds, Bull. Amer. Math. Soc. 36 (1999),135-249.
[6] A. GRIGOR' YAN, Isoperimetric inequalities and capacities on Riemannian manifolds The Mazya anniversary collection, Vol. 1 (Rostock, 1998), 139-153, Oper. Theory Adv. Appl, 109, Birkhäuser, Basel, 1999.
[7] V. Gol'dshtein and M. Troyanov, The Kelvin-Nevanlinna-Royden criterion for p-parabolicity, Math Z. 232 (1999), 607-619.
[8] R. Greene and H. H. Wu, Function theory on manifolds which possess a pole, Lecture Notes in Math. 699, Springer Verlag, Berlin, 1979.
[9] R. Hardt and F.-H. Lin, Mappings minimizing the $L^{p}$ norm of the gradient, Comm. Pure Appl. Math. 40 (1987), 555-588.
[10] I. Holopainen, Nonlinear potential theory and quasiregular mappings on Riemannian manifolds, Ann. Acad. Sci. Fenn. Ser. A I Math. Dissertationes 74 (1990), 45 pp.
[11] I. Holopainen, Quasiregular mappings and the p-Laplace operator, Heat kernels and analysis on manifolds, graphs, and metric spaces (Paris, 2002), 219-239, Contemp. Math. 338, Amer. Math. Soc., Providence, RI, 2003.
[12] I. Holopainen, S. Pigola and G. Veronelli, Global comparison principles for the $p$-Laplace operator on Riemannian manifolds, Potential Analysis 34 (2011), 371-384.
[13] L. Karp, On Stokes' Theorem for noncompact manifolds, Proc. Amer. Math. Soc. 82 (1981), 487-490.
[14] P. Li and J. WANG, Minimal hypersurfaces with finite index, Math. Res. Lett. 9 (2002), 95-103.
[15] P. LindQVISt, On the equation $\operatorname{div}\left(|\nabla u|^{p-2} \nabla u\right)+\lambda|u|^{p-2} u=0$, Proc. Amer Math. Soc. 109 (1990), 157-164.
[16] T. LyOns and D. SULLIVAN, Function theory, random paths and covering spaces, J. Differential Geom. 19 (1984), 229-323.
[17] S. Pigola, M. Rigoli and A. G. Setti, Constancy of $p$-harmonic maps of finite $q$-energy into nonpositively curved manifolds, Math. Z. 258 (2008), 347-362.
[18] S. Pigola, M. Rigoli and A. G. Setti, Vanishing and finiteness results in geometric analysis, A generalization of the Bochner technique, Prog. Math. 266 (2008), Birkhäuser Verlag, Basel 2008.
[19] S. Pigola, A. G. Setti and M. Troyanov, The topology at infinity of a manifold supporting an $L^{p, q_{-}}$ Sobolev inequality, Submitted. Preliminary version: arXiv:1007.1761v1
[20] M. Rigoli and A. G. Setti, Liouville type theorems for $\varphi$-subharmonic functions, Rev. Mat. Iberoamericana 17 (2001), 471-520.
[21] L. Sario and M. Nakai, Classification theory of Riemann surfaces, Die Grundlehren der mathematischen Wissenschaften, Bomd 164, Springer Verlag, New York-Berlin, 1970.
[22] R. Schoen and S. T. Yau, Compact group actions and the topology of manifolds with nonpositive curvature, Topology 18 (1979), 361-380.
[23] R. Schoen and S. T. Yau, Harmonic maps and the topology of stable hypersurfaces and manifolds with non-negative Ricci curvature, Comment. Math. Helv. 51 (1976), 333-341.
[24] P. TolkSDORF, Everywhere-regularity for some quasilinear systems with a lack of ellipticity, Ann. Mat. Pura Appl. (4) 134 (1983), 241-266.
[25] M. Troyanov, Parabolicity of manifolds, Siberian Adv. Math. 9 (1999), 125-150.
[26] D. Valtorta, Potenziali di Evans su varietà paraboliche, Master Thesis, 2009.
[27] S. W. WEI and C.-M. YaU, Regularity of $p$-energy minimizing maps and $p$-superstrongly unstable indices, J. Geom. Anal. 4 (1994), 247-272.

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