# Recognizing the flat torus among $\operatorname{RCD}^{*}(0, N)$ spaces via the study of the first cohomology group 

Nicola Gigli * Chiara Rigoni ${ }^{\dagger}$

May 12, 2017


#### Abstract

We prove that if the dimension of the first cohomology group of a $\operatorname{RCD}^{*}(0, N)$ space is $N$, then the space is a flat torus.

This generalizes a classical result due to Bochner to the non-smooth setting and also provides a first example where the study of the cohomology groups in such synthetic framework leads to geometric consequences.


## Contents

1 Introduction ..... 2
2 Preliminaries ..... 5
$2.1 \quad L^{2}$ - and $L^{0}$-normed modules and basis of differential calculus ..... 5
2.2 Sobolev spaces for locally integrable objects ..... 6
2.3 Regular Lagrangian Flows and continuity equation ..... 11
3 Calculus Tools ..... 15
3.1 Local bounded compression/deformation ..... 15
3.1.1 Pullback module through a map of local bounded compression ..... 15
3.1.2 Localized pullback of 1-forms ..... 18
3.2 Calculus on product spaces ..... 20
3.2.1 Cotangent module and product of spaces ..... 20
3.2.2 Other differential operators in product spaces ..... 25
3.3 Flow of harmonic vector fields on $\operatorname{RCD}(0, \infty)$ spaces ..... 28

[^0]4 Proof of the Main Result ..... 34
4.1 Setting ..... 34
4.2 Preliminary considerations ..... 34
4.3 An explicit formula for Regular Lagrangian Flows on X ..... 36
4.4 Further properties of T and conclusion ..... 40
A Notes on the Hessian on product spaces ..... 44

## 1 Introduction

A classical result due to Bochner concerning manifolds with non-negative Ricci curvature is:

Theorem 1.1 (Bochner). Let $M$ be a compact, smooth and connected Riemannian manifold with non-negative Ricci curvature. Then:
i) The dimension of the first cohomology group is bounded above by the dimension of the manifold,
ii) If these two dimensions are equal, then $M$ is a flat torus.

The key observation that leads to $(i)$ is the fact that under the stated assumptions every harmonic 1 -form must be parallel and thus determined by its value at any given point $x \in M$. Since Hodge's theorem grants that the $k$-th cohomology group is isomorphic to the space of harmonic k-forms, the claim follows.

For (ii), the typical argument starts with the observation that if an $n$-dimensional manifold admits $n$ independent parallel vector fields, then such manifold must be flat. Hence its universal cover, equipped with the pullback of the metric tensor, must be the Euclidean space and the fundamental group $\pi_{1}(M)$ acts on it via isometries. Since $M \sim \mathbb{R}^{n} / \pi_{1}(M)$, all is left to show is that $\pi_{1}(M) \sim \mathbb{Z}^{n}$, which can be obtained by 'soft' considerations about the structure of the isometries of $\mathbb{R}^{n}$ and the fact that $\mathbb{R}^{n} / \pi_{1}(M)$ is, by assumption, compact and smooth (see e.g. [26] for the details).

This paper is about the generalization of the above result to the non-smooth setting of $\operatorname{RCD}^{*}(0, N)$ spaces ([5], [17]). Our starting point is the paper [15] by the first author where a differential calculus on such spaces has been built. Among other things, the vocabulary proposed there allows to speak of vector fields, $k$-forms, covariant derivative, Hodge laplacian and cohomology groups $H_{\mathrm{dR}}^{k}$. In particular, a quite natural version of Hodge's theorem exists in this non-smooth setting, so that we know that cohomology classes are in correspondence with their unique harmonic representative. The basic structure around which the theory is built - and which offers a counterpart for the space of $L^{2}$ section of a normed vector bundle over a smooth manifold - is the one of $L^{2}$-normed $L^{\infty}$-module.

In searching for an analogous of Theorem 1.1 in the nonsmooth setting, one thing to discuss is the notion of dimension. To this aim, let us recall that given a generic $L^{2}$ normed $L^{\infty}$-module $\mathscr{M}$ over a space (X, d, $\mathfrak{m}$ ), there exists a unique, up to negligible sets, Borel partition $\left(E_{i}\right)_{i \in \mathbb{N} \cup\{\infty\}}$ of X such that for every $i \in \mathbb{N}$ the restriction of $M$ to $E_{i}$ has dimension $i$, and for no $F \subset E_{\infty}$ with positive measure the restriction of $\mathscr{M}$ to $F$ has finite dimension. We consider such partition $\left(E_{i}\right)$ for $M$ being the tangent module of the space X and think to the dimension of tangent module on $E_{i}$ as the dimension of the space X on the same set. Hence we shall call $\operatorname{dim}_{\min }(\mathrm{X})\left(\mathrm{resp}\right.$. $\left.\operatorname{dim}_{\max }(\mathrm{X})\right)$ the minimal (resp. supremum) of indexes $i \in \mathbb{N} \cup\{\infty\}$ such that $\mathfrak{m}\left(E_{i}\right)>0$. Then in [15] the following result has been obtained:

Theorem 1.2. Let $(\mathrm{X}, \mathrm{d}, \mathfrak{m})$ be a $\operatorname{RCD}^{*}(0, \infty)$ space. Then $\operatorname{dim}\left(H_{\mathrm{dR}}^{k}\right)(\mathrm{X}) \leq \operatorname{dim}_{\min }(\mathrm{X})$.
The proof of this fact closely follows the argument for point $(i)$ in Theorem 1.1 every harmonic form is proved to be parallel, so that the dimension of the first cohomology group is bounded by how many independent (co)vector fields we can find on any region of our space. Notice that the compactness assumption is not present because, in the terminology of [15, harmonic forms are by definition in $L^{2}$.

We also point out that a priori such result might be empty, in the sense that without any additional assumption it is very possible that $\operatorname{dim}_{\min }(\mathrm{X})=\infty$. In fact, the natural assumption on the space X is not that it is a $\operatorname{RCD}^{*}(0, \infty)$ space, but rather a $\operatorname{RCD}^{*}(0, N)$ one: given that the number $N \in[1, \infty]$ represents, in some sense, an upper bound for the dimension of the space we expect it to bound from above $\operatorname{dim}_{\max }(\mathrm{X})$. This is indeed the case, as it has been proved by Han in [23] (see also [21] for an alternative argument):

Theorem 1.3. Let $(\mathrm{X}, \mathrm{d}, \mathfrak{m})$ be a $\mathrm{RCD}^{*}(K, N)$ space. Then $\operatorname{dim}_{\max }(\mathrm{X}) \leq N$.
Coupling Theorems 1.2 and 1.3 we deduce that:
Proposition 1.4. Let $(\mathrm{X}, \mathrm{d}, \mathfrak{m})$ be a $\mathrm{RCD}^{*}(0, N)$ space. Then $\operatorname{dim}\left(H_{\mathrm{dR}}^{k}\right)(\mathrm{X}) \leq N$.
This statement is a perfect analogue of point $(i)$ in Theorem 1.1. The aim of the present manuscript is to complete this analogy between the smooth and nonsmooth setting by proving:

Theorem 1.5. Let $(\mathrm{X}, \mathrm{d}, \mathfrak{m})$ be a $\mathrm{RCD}^{*}(0, N)$ space such that $\operatorname{dim}\left(H_{\mathrm{dR}}^{k}\right)(\mathrm{X})=N$ (so that in particular $N$ is integer). Then it is isomorphic to the flat $N$-dimensional torus.

Here 'isomorphic' means that there exists a measure preserving isometry from the torus equipped with its Riemannian distance and a constant multiple of the induced volume measure to our space. Unlike the proof of Theorem 1.2 which closely mimics the original one, here we cannot adapt the 'smooth arguments' to the current setting: the problem is that it is not known whether RCD spaces admit a universal cover or not (there are some
results in this direction - see [25] - but it is unclear to us whether they can effectively be used for our current purposes).

We therefore have to pursue a different strategy, the starting steps of our argument being:

- We start studying the flow of an harmonic vector field on our space and, using the fact that in particular such vector field must be parallel and divergence-free, we prove that in accordance with the smooth case such flow is made of measure preserving isometries. Here by 'flow' we intend in fact 'Regular Lagrangian Flow' in the sense of Ambrosio-Trevisan [8] who adapted to the setting of RCD spaces the analogous notion developed by Ambrosio in [1] in connection with the Di Perna-Lions theory [12]. We remark that our appears to be the first application of Ambrosio-Trevisan theory to vector fields which are not gradients.
- We prove that given two such vector fields $X, Y$, for their flows $\mathrm{Fl}_{t}^{X}$ and $\mathrm{Fl}_{s}^{Y}$ we have the formula $\mathrm{Fl}_{t}^{X} \circ \mathrm{Fl}_{s}^{Y}=\mathrm{Fl}_{1}^{t X+s Y}$ for any $t, s \in \mathbb{R}$.
- Our assumption on the space ( $\mathrm{X}, \mathrm{d}, \mathfrak{m}$ ) grants that there are $N$ independent and orthogonal vector fields $X_{1}, \ldots, X_{N}$ which are parallel and divergence-free, hence we can define the map $\mathrm{T}: \mathrm{X} \times \mathbb{R}^{N} \rightarrow \mathrm{X}$ by

$$
\left(x, a_{1}, \ldots, a_{N}\right) \quad \mapsto \quad \mathrm{Fl}_{a_{1}}^{X_{1}}\left(\cdots \mathrm{Fl}_{a_{N}}^{X_{N}}(x)\right) .
$$

What previously proved ensures that this map can be seen as an action of $\mathbb{R}^{N}$ on X by isomorphisms.

Analyzing the properties of the map T will lead to the desired isomorphism with the torus. The hardest part will be the proof of the fact that the action is transitive: to obtain this will require a sharpening of the calculus tools available in the nonsmooth setting and, in particular, we will analyze the structure of the (co)tangent modules on product spaces which we believe to be interesting on its own.

We conclude recalling that Honda proved in [24] that the dimension of the first cohomology group is upper semicontinuous along a non-collapsing sequence of manifolds with same dimension and a uniform lower bound on the Ricci. This result hints at the possibility of obtaining an almost rigidity statement of our theorem in the context of RCD spaces, which would informally read as
'if a RCD space almost fulfils the assumption of Theorem 1.5, then it is mGHclose to a flat torus'.

In this direction it is worth to emphasize that in the smooth category more is known: as Colding proved in [11], $N$-dimensional manifolds with Ricci $\geq-\varepsilon$ and first cohomology group of dimension $N$, not only must be mGH-close to the torus, but also homeomorphich to it if $N \neq 3$ and homotopic if $N=3$. Such topological information is out of reach of simple arguments based on the mGH-compactness of the class of RCD spaces.

## 2 Preliminaries

To keep the presentation at reasonable length we shall assume the reader familiar with the notions of Sobolev functions (see [10], [28, [4), of differential calculus on metric measure spaces (see [15], [13]) and with the notion of Regular Lagrangian Flow on metric measure spaces ( $[8],[9])$. Here we shall only recall a few facts mainly to fix the notation.

For us a metric measure space ( $\mathrm{X}, \mathrm{d}, \mathfrak{m}$ ) will always be a complete and separable metric space equipped with a reference non-negative (and non-zero) Borel measure $\mathfrak{m}$ which is finite on bounded sets and with full support, i.e. $\mathfrak{m}(\Omega)>0$ for every non-empty open set $\Omega \subset \mathrm{X}$. This latter assumption is not really necessary, but simplifies some statements.

## $2.1 \quad L^{2}$ - and $L^{0}$-normed modules and basis of differential calculus

Definition $2.1\left(L^{2}(\mathrm{X})\right.$-normed modules). A $L^{2}(\mathrm{X})$-normed $L^{\infty}(\mathrm{X})$-module, or simply a $L^{2}(\mathrm{X})$-normed module, is a structure $(\mathscr{M},\|\cdot\|, \cdot,|\cdot|)$ where
i) $(\mathscr{M},\|\cdot\|)$ is a Banach space
ii) • is a bilinear map from $L^{\infty}(\mathrm{X}) \times \mathscr{M}$ to $\mathscr{M}$, called multiplication by $L^{\infty}(\mathrm{X})$ functions, such that

$$
\begin{align*}
f \cdot(g \cdot v) & =(f g) \cdot v,  \tag{2.1.1a}\\
1 \cdot v & =v, \tag{2.1.1b}
\end{align*}
$$

for every $v \in \mathscr{M}$ and $f, g \in L^{\infty}(\mathrm{X})$, where $\mathbf{1}$ is the function identically equal to 1 .
iii) $|\cdot|$ is a map from $\mathscr{M}$ to $L^{2}(\mathrm{X})$, called pointwise norm, such that

$$
\begin{align*}
|v| & \geq 0 \quad \mathfrak{m} \text {-a.e. }  \tag{2.1.2a}\\
|f v| & =|f||v| \quad \mathfrak{m} \text {-a.e. }  \tag{2.1.2b}\\
\|v\| & =\sqrt{\int|v|^{2} \mathrm{dm}} \tag{2.1.2c}
\end{align*}
$$

An isomorphism between two $L^{2}(\mathrm{X})$-normed modules is a linear bijection which preserves the norm, the product with $L^{\infty}(\mathrm{X})$ functions and the pointwise norm.

Definition 2.2 ( $L^{0}$-normed module). A $L^{0}$-normed module is a structure $(\mathscr{M}, \tau, \cdot,|\cdot|)$ where:
i) - is a bilinear map, called multiplication with $L^{0}$ functions, from $L^{0}(\mathrm{X}) \times \mathscr{M}$ to $\mathscr{M}$ for which 2.1.1a, 2.1.1b hold for any $f \in L^{0}(\mathrm{X}), v \in \mathscr{M}$,
ii) $|\cdot|: \mathscr{M} \rightarrow L^{0}(\mathrm{X})$, called pointwise norm, satisfies 2.1.2a and 2.1.2b for any $f \in L^{0}(\mathrm{X}), v \in \mathscr{M}$,
iii) for some Borel partition $\left(E_{i}\right)$ of X into sets of finite $\mathfrak{m}$-measure, $\mathscr{M}$ is complete w.r.t. the distance

$$
\begin{equation*}
\mathrm{d}_{0}(v, w):=\sum_{i} \frac{1}{2^{i} \mathfrak{m}\left(E_{i}\right)} \int_{E_{i}} \min \{1,|v-w|\} \mathrm{d} \mathfrak{m} \tag{2.1.3}
\end{equation*}
$$

and $\tau$ is the topology induced by the distance.
An isomorphims of $L^{0}$-normed modules is a linear homeomorphism preserving the pointwise norm and the multiplication with $L^{0}$-functions.

It is readily checked that the choice of the partition $\left(E_{i}\right)$ in (iii) does not affect the completeness of $\mathscr{M}$ nor the topology $\tau$.

Theorem/Definition 2.3 ( $L^{0}$ completion of a module). Let $\mathscr{M}$ be a $L^{2}$-normed module. Then there exists a unique couple $\left(\mathscr{M}^{0}, \iota\right)$, where $\mathscr{M}^{0}$ is a $L^{0}$-normed module and $\iota: \mathscr{M} \rightarrow$ $\mathscr{M}^{0}$ is linear, preserving the pointwise norm and with dense image.

Uniqueness is intended up to unique isomorphism, i.e.: if ( $\left.\tilde{\mathscr{M}}^{0}, \tilde{\iota}\right)$ has the same properties, then there exists a unique isomorphism $\Phi: \mathscr{M}^{0} \rightarrow \tilde{\mathscr{M}}^{0}$ such that $\tilde{\iota}=\Phi \circ \iota$.

### 2.2 Sobolev spaces for locally integrable objects

Given a metric measure space ( $\mathrm{X}, \mathrm{d}, \mathfrak{m}$ ), by $L_{\text {loc }}^{2}(\mathrm{X})$ we mean the space of (equivalence classes w.r.t. $\mathfrak{m}$-a.e. equality of) Borel functions $f: \mathrm{X} \rightarrow \mathbb{R}$ such that $\chi_{B} f \in L^{2}(\mathrm{X})$ for every bounded Borel set $B \subset \mathrm{X}$. A curve $t \mapsto f_{t} \in L_{\text {loc }}^{2}(\mathrm{X})$ will be called continuous (resp. absolutely continuous, Lipschitz, $C^{1}$ ) provided for any bounded Borel set $B \subset \mathrm{X}$ the curve $t \mapsto \chi_{B} f_{t} \in L^{2}(\mathrm{X})$ is continuous (resp. absolutely continuous, Lipschitz, $C^{1}$ ).

Recall that $\boldsymbol{\pi} \in \mathscr{P}(C([0,1], \mathrm{X}))$ is called test plan provided

$$
\begin{aligned}
\left(\mathrm{e}_{t}\right)_{*} \boldsymbol{\pi} & \leq C \mathfrak{m} \quad \forall t \in[0,1], \quad \text { for some } C>0, \\
\iint_{0}^{1}\left|\dot{\gamma}_{t}\right|^{2} \mathrm{~d} t \mathrm{~d} \boldsymbol{\pi} & <\infty .
\end{aligned}
$$

The Sobolev class $\mathcal{S}^{2}(\mathrm{X})\left(\right.$ resp. $\mathcal{S}_{\text {loc }}^{2}(\mathrm{X})$ ) is the space of all Borel functions $f: \mathrm{X} \rightarrow \mathbb{R}$ for which there exists $G \in L^{2}(\mathrm{X})$ (resp. $G \in L_{\text {loc }}^{2}(\mathrm{X})$ ) non-negative, called weak upper gradient, such that

$$
\int\left|f\left(\gamma_{1}\right)-f\left(\gamma_{0}\right)\right| \mathrm{d} \boldsymbol{\pi}(\gamma) \leq \iint_{0}^{1} G\left(\gamma_{t}\right)\left|\dot{\gamma}_{t}\right| \mathrm{d} t \mathrm{~d} \boldsymbol{\pi}(\gamma) \quad \text { for every test plan } \boldsymbol{\pi}
$$

It turns out that $f \in \mathcal{S}_{\text {loc }}^{2}(\mathrm{X})$ and $G$ is a weak upper gradient if and only if for every test plan $\boldsymbol{\pi}$ we have that for $\boldsymbol{\pi}$-a.e. $\gamma$ the map $t \mapsto f\left(\gamma_{t}\right)$ is in $W^{1,1}(0,1)$ and

$$
\begin{equation*}
\left|\frac{\mathrm{d}}{\mathrm{~d} t} f\left(\gamma_{t}\right)\right| \leq G\left(\gamma_{t}\right)\left|\dot{\gamma}_{t}\right| \quad \text { a.e. } t \in[0,1] . \tag{2.2.1}
\end{equation*}
$$

From this characterization it follows that there exists a minimal weak upper gradient in the $\mathfrak{m}$-a.e. sense: it will be called minimal weak upper gradient and denoted by $|\mathrm{D} f|$. With a simple cut-off argument, (2.2.1) also shows that $f \in \mathcal{S}_{\text {loc }}^{2}(\mathrm{X})$ if and only if $\eta f \in \mathcal{S}^{2}(\mathrm{X})$ for every $\eta$ Lipschitz and with bounded support.

The Sobolev space $W^{1,2}(\mathrm{X})$ (resp. $W_{\text {loc }}^{1,2}(\mathrm{X})$ ) is defined as $L^{2} \cap \mathcal{S}^{2}(\mathrm{X})$ (resp. $L_{\text {loc }}^{2} \cap$ $\mathcal{S}_{\text {loc }}^{2}(\mathrm{X})$ ) and again one can check that $f \in W_{\text {loc }}^{1,2}(\mathrm{X})$ if and only if $\eta f \in W^{1,2}(\mathrm{X})$ for every $\eta$ Lipschitz and with bounded support. $W^{1,2}(\mathrm{X})$ is a Banach space when endowed with the norm

$$
\|f\|_{W^{1,2}(\mathrm{X})}^{2}:=\|f\|_{L^{2}(\mathrm{X})}^{2}+\|\mid \mathrm{D} f\|_{L^{2}(\mathrm{X})}^{2}
$$

( $\mathrm{X}, \mathrm{d}, \mathfrak{m}$ ) is said infinitesimally Hilbertian provided $W^{1,2}(\mathrm{X})$ is a Hilbert space.
Among others, minimal weak upper gradients have the following important locality property:

$$
|\mathrm{D} f|=|\mathrm{D} g| \quad \mathfrak{m}-\text { a.e. on }\{f=g\} \quad \forall f, g \in \mathcal{S}_{\text {loc }}^{2}(\mathrm{X}) .
$$

Also, it will be useful to keep in mind that

$$
\begin{equation*}
\forall f \in \mathcal{S}_{\text {loc }}^{2}(\mathrm{X}) \text { there exists }\left(f_{n}\right) \subset \mathcal{S}^{2}(\mathrm{X}) \text { such that } \mathfrak{m}\left(\mathrm{X} \backslash \cup_{n}\left\{f=f_{n}\right\}\right)=0 \tag{2.2.2}
\end{equation*}
$$

and for every $n \in \mathbb{N}$ the function $f_{n}$ is bounded and with bounded support.
Such sequence can be obtained noticing that for $f \in \mathcal{S}_{\text {loc }}^{2}(\mathrm{X})$ the truncated function ( $f \wedge$ $(-c)) \vee c$ also belongs to $\mathcal{S}_{\text {loc }}^{2}(\mathrm{X})$ for every $c>0$ and then proceeding with a cut-off argument.

From the notion of minimal weak upper gradient it is possible to extract the one of differential via the following result:

Theorem/Definition 2.4. There exists a unique couple ( $L^{0}\left(T^{*} \mathrm{X}\right)$, d) with $L^{0}\left(T^{*} \mathrm{X}\right)$ being a $L^{0}(\mathrm{X})$-normed module and $\mathrm{d}: \mathcal{S}_{\text {loc }}^{2}(\mathrm{X}) \rightarrow L^{0}\left(T^{*} \mathrm{X}\right)$ linear and such that
i) $|\mathrm{d} f|=|\mathrm{D} f| \mathfrak{m}$-a.e. for every $f \in \mathcal{S}_{\text {loc }}^{2}(\mathrm{X})$,
ii) $L^{0}\left(T^{*} \mathrm{X}\right)$ is generated by $\left\{\mathrm{d} f: f \in \mathcal{S}_{\text {loc }}^{2}(\mathrm{X})\right\}$, i.e. $L^{0}$-linear combinations of objects of the form $\mathrm{d} f$ are dense in $L^{0}\left(T^{*} \mathrm{X}\right)$.

Uniqueness is intended up to unique isomorphism, i.e. if ( $\left.\mathscr{M}, \mathrm{d}^{\prime}\right)$ is another such couple, then there is a unique isomorphism $\Phi: L^{0}\left(T^{*} \mathrm{X}\right) \rightarrow \mathscr{M}$ such that $\Phi(\mathrm{d} f)=\mathrm{d}^{\prime} f$ for every $f \in \mathcal{S}_{\text {loc }}^{2}(\mathrm{X})$.

Some remarks:
a) Using the approximation property 2.2 .2 one can show that $\mathcal{S}_{\text {loc }}^{2}(\mathrm{X})$ can be replaced with either one of $\mathcal{S}^{2}(\mathrm{X}), W^{1,2}(\mathrm{X})$ in the above statement
b) If one chooses to replace $\mathcal{S}_{\text {loc }}^{2}(\mathrm{X})$ with either $\mathcal{S}^{2}(\mathrm{X})$ of $W^{1,2}(\mathrm{X})$, then it is also possible to replace the $L^{0}$-normed module with a $L^{2}$-normed module in the statement and in this case in (ii) ' $L^{0}$-linear' should be replaced by ' $L^{\infty}$-linear' (notice that the choice of the module also affects the topology considered, whence the possibility of having two different uniqueness results). The proof is unaltered: compare for instance Theorem 3.2 with the construction of pullback module given in [15] and [13].
c) Call $\left(L^{2}\left(T^{*} \mathrm{X}\right), \underline{\mathrm{d}}\right)$ the outcome of Theorem 2.4 written for $L^{2}$-normed modulus and one of the spaces $\mathcal{S}^{2}(\mathrm{X}), W^{1,2}(\mathrm{X})$. Then its $L^{0}$-completion can be fully identified with the couple $\left(L^{0}\left(T^{*} \mathrm{X}\right), \mathrm{d}\right)$ given by Theorem 2.4 in the sense that: there is a unique linear map $\iota: L^{2}\left(T^{*} \mathrm{X}\right) \rightarrow L^{0}\left(T^{*} \mathrm{X}\right)$ sending $\underline{\mathrm{d}} f$ to $\mathrm{d} f$ and preserving the pointwise norm, moreover such map has dense image.
This is trivial to check from the definitions and for this reason we won't use a distinguished notation for the differential coming from the ' $L^{2}$ ' formulation of the statement.

Let us now discuss other differentiation operators defined for objects with $L_{\text {loc }}^{2}$ integrability.
The space of vector fields $L^{0}(T \mathrm{X})$ is defined as the dual of the $L^{0}$-normed module $L^{0}\left(T^{*} \mathrm{X}\right)$. Equivalently, it is the $L^{0}$-completion of the dual $L^{2}(T \mathrm{X})$ of the $L^{2}$-normed module $L^{2}\left(T^{*} \mathrm{X}\right)$ (see [15], [13]). $L_{\text {loc }}^{2}(T \mathrm{X}) \subset L^{0}(T \mathrm{X})$ is the space of $X$ 's such that $|X| \in$ $L_{\text {loc }}^{2}(\mathrm{X})$.

We say that $X \in L_{\text {loc }}^{2}(T \mathrm{X})$ has divergence in $L_{\text {loc }}^{2}$, and write $X \in \mathrm{D}\left(\operatorname{div}_{\text {loc }}\right)$, if there exists $h \in L_{\text {loc }}^{2}(\mathrm{X})$ such that for every $f \in W^{1,2}(\mathrm{X})$ with bounded support it holds

$$
\int f h \mathrm{~d} \mathfrak{m}=-\int \mathrm{d} f(X) \mathrm{d} \mathfrak{m} .
$$

In this case we put $\operatorname{div} X:=h$.
Let us now assume that ( $\mathrm{X}, \mathrm{d}, \mathfrak{m}$ ) is infinitesimally Hilbertian, so that the pointwise norms in $L^{0}\left(T^{*} \mathrm{X}\right), L^{0}(T \mathrm{X})$ induce pointwise scalar products. Recall that in this case the modules $L^{0}\left(T^{*} \mathrm{X}\right)$ and $L^{0}(T \mathrm{X})$ are canonically isomorphic via the Riesz (musical) isomorphism

$$
\begin{equation*}
b: L^{0}(T \mathrm{X}) \rightarrow L^{0}\left(T^{*} \mathrm{X}\right) \quad \text { and } \quad \sharp: L^{0}\left(T^{*} \mathrm{X}\right) \rightarrow L^{0}(T \mathrm{X}) . \tag{2.2.3}
\end{equation*}
$$

The gradient of a function $f \in W_{\text {loc }}^{1,2}(\mathrm{X})$ is defined as $\nabla f:=(\mathrm{d} f)^{\sharp} \in L_{\text {loc }}^{2}(T \mathrm{X})$.
We say that $f \in L_{\text {loc }}^{2}(\mathrm{X})$ has Laplacian in $L_{\text {loc }}^{2}$, and write $f \in \mathrm{D}\left(\Delta_{\text {loc }}\right)$ if there exists $h \in L_{\text {loc }}^{2}(\mathrm{X})$ such that for every $g \in W^{1,2}(\mathrm{X})$ with bounded support it holds

$$
\int g h \mathrm{~d} \mathfrak{m}=-\int\langle\nabla f, \nabla g\rangle \mathrm{d} \mathfrak{m} .
$$

In this case we put $\Delta f:=h$. If $f, h \in L^{2}(\mathrm{X})$ we shall write $f \in D(\Delta)$ instead of $f \in$ $D\left(\Delta_{\text {loc }}\right)$. It is not hard to check that in this case this notion is equivalent to the more
familiar notion of Laplacian as infinitesimal generator of the Dirichlet form

$$
\mathrm{E}(f):= \begin{cases}\frac{1}{2} \int|\mathrm{~d} f|^{2} \mathrm{~d} \mathfrak{m} & \text { if } f \in W^{1,2}(\mathrm{X}) \\ +\infty & \text { otherwise }\end{cases}
$$

(see [4], 5]).
To continue this introduction, we shall now assume that ( $\mathrm{X}, \mathrm{d}, \mathfrak{m}$ ) is a $\mathrm{RCD}(K, \infty)$ space for some $K \in \mathbb{R}$. A relevant property of Sobolev functions in relations to the metric of such spaces is the following result, proved in [5] (see also [14], [16] for the given formulation):

Theorem 2.5. Let $\left(\mathrm{X}_{1}, \mathrm{~d}_{1}, \mathfrak{m}_{1}\right)$ and $\left(\mathrm{X}_{2}, \mathrm{~d}_{2}, \mathfrak{m}_{2}\right)$ be two $\operatorname{RCD}(K, \infty)$ spaces with $\mathfrak{m}_{1}, \mathfrak{m}_{2}$ having full support and $T: \mathrm{X}_{1} \rightarrow \mathrm{X}_{2}$ and $S: \mathrm{X}_{2} \rightarrow \mathrm{X}_{1}$ be Borel maps such that

$$
T \circ S=\operatorname{Id}_{\mathrm{X}_{2}} \quad \mathfrak{m}_{2}-\text { a.e. } \quad S \circ T=\operatorname{Id}_{\mathrm{X}_{1}} \quad \mathfrak{m}_{1}-\text { a.e. }
$$

and

$$
T_{*} \mathfrak{m}_{1}=\mathfrak{m}_{2} \quad \mathrm{E}_{\mathrm{X}_{1}}(f \circ T)=\mathrm{E}_{\mathrm{X}_{2}}(f) \quad \forall f \in L^{2}\left(\mathrm{X}_{2}\right)
$$

Then, up to modifications in a $\mathfrak{m}_{1}$-negligible set, $T$ is an isometry.
Recall that the class of test functions (see [27]) is defined as

$$
\operatorname{Test}(\mathrm{X}):=\left\{f \in L^{\infty} \cap W^{1,2}(\mathrm{X}) \cap D(\Delta):|\mathrm{d} f| \in L^{\infty}(\mathrm{X}), \Delta f \in W^{1,2}(\mathrm{X})\right\}
$$

Crucial properties of test functions are that they form an algebra and that $|\mathrm{d} f|^{2} \in W^{1,2}(\mathrm{X})$ for $f \in \operatorname{Test}(\mathrm{X})$ (see [27]). For our discussion, it is also useful to keep in mind that
for every $K \subset \Omega \subset \mathrm{X}$ with $\mathrm{d}(x, y) \geq c$ for some $c>0$ and every $x \in K, y \in \Omega^{c}$ there exists $f \in \operatorname{Test}(\mathrm{X})$ with $\operatorname{supp}(f) \subset \Omega$ and identically 1 on $K$,
see e.g. [7].
With this said, we can define the space $W_{\text {loc }}^{2,2}(\mathrm{X})$ as the space of $f \in W_{\text {loc }}^{1,2}(\mathrm{X})$ for which there exists $A \in L_{\text {loc }}^{2}\left(\left(T^{*}\right)^{\otimes 2} \mathrm{X}\right)$ such that

$$
\begin{align*}
& 2 \int h A(\nabla g, \nabla \tilde{g}) \mathrm{d} \mathfrak{m} \\
& \quad=-\int\langle\nabla f, \nabla g\rangle \operatorname{div}(h \nabla \tilde{g})+\langle\nabla f, \nabla \tilde{g}\rangle \operatorname{div}(h \nabla g)+h\langle\nabla f, \nabla(\langle\nabla g, \nabla \tilde{g}\rangle)\rangle \mathrm{d} \mathfrak{m} \tag{2.2.5}
\end{align*}
$$

for every $g, \tilde{g}, h \in \operatorname{Test}(\mathrm{X})$ with bounded support. In this case we call $A$ the Hessian of $f$ and denote it by $\operatorname{Hess}(f)$. If $f \in W^{1,2}(\mathrm{X})$ and $\operatorname{Hess}(f) \in L^{2}\left(\left(T^{*}\right)^{\otimes 2} \mathrm{X}\right)$ we say that $f \in W^{2,2}(\mathrm{X})$. Noticing that the two sides of 2.2 .5 are continuous in $h$ w.r.t. the $W^{1,2_{-}}$ norm and using property 2.2 .4 it is easy to check that this definition of $W^{2,2}(\mathrm{X})$ coincides
with the one given in [15]. Also, $W^{2,2}(\mathrm{X})$ is a separable Hilbert space when endowed the norm

$$
\|f\|_{W^{2,2}(\mathrm{X})}^{2}:=\|f\|_{L^{2}(\mathrm{X})}^{2}+\||\mathrm{D} f|\|_{L^{2}(\mathrm{X})}^{2}+\left\||\operatorname{Hess}(f)|_{\mathrm{Hs}}\right\|_{L^{2}(\mathrm{X})}^{2} .
$$

We recall (see [15]) that $D(\Delta) \subset W^{2,2}(\mathrm{X})$ and

$$
\begin{equation*}
\int|\operatorname{Hess}(f)|_{\text {HS }}^{2} \mathrm{~d} \mathfrak{m} \leq \int|\Delta f|^{2}-K|\mathrm{~d} f|^{2} \mathrm{~d} \mathfrak{m}, \tag{2.2.6}
\end{equation*}
$$

so that in particular $\operatorname{Test}(\mathrm{X}) \subset W^{2,2}(\mathrm{X})$. We then define $H^{2,2}(\mathrm{X})$ as the $W^{2,2}$-closure of $\operatorname{Test}(\mathrm{X})$ and similarly $H_{\mathrm{loc}}^{2,2}(\mathrm{X})$ as the $W_{\mathrm{loc}}^{2,2}(\mathrm{X})$-closure of $\operatorname{Test}(\mathrm{X})$, i.e.: $f \in H_{\mathrm{loc}}^{2,2}(\mathrm{X}) \subset$ $W_{\text {loc }}^{2,2}(\mathrm{X})$ provided there exists a sequence $\left(f_{n}\right) \subset \operatorname{Test}(\mathrm{X})$ such that $f_{n}, \mathrm{~d} f_{n}, \operatorname{Hess}\left(f_{n}\right)$ converge to $f, \mathrm{~d} f, \operatorname{Hess}(f)$ in $L_{\text {loc }}^{2}(\mathrm{X}), L_{\text {loc }}^{2}\left(T^{*} \mathrm{X}\right), L_{\text {loc }}^{2}\left(\left(T^{*}\right)^{\otimes 2} \mathrm{X}\right)$ respectively.

The space of Sobolev vector fields $W_{C, \text { loc }}^{1,2}(T \mathrm{X})$ is defined as the space of $X \in L_{\text {loc }}^{2}(T \mathrm{X})$ for which there is $T \in L^{2}\left(T^{\otimes 2} \mathrm{X}\right)$ such that

$$
\int h T(\nabla g, \nabla \tilde{g}) \mathrm{d} \mathfrak{m}=\int-\langle X, \nabla \tilde{g}\rangle \operatorname{div}(h \nabla g)+h \operatorname{Hess}(\tilde{g})(X, \nabla g) \mathrm{d} \mathfrak{m}
$$

for every $h, g, \tilde{g} \in \operatorname{Test}(\mathrm{X})$ with bounded support. In this case we call $T$ the covariant derivative of $X$ and denote it as $\nabla X$. If $|X|,|\nabla X|_{\text {HS }} \in L^{2}(\mathrm{X})$ we shall say that $X \in$ $W_{C}^{1,2}(T X)$; again, it is not hard to check that this definition of $W_{C}^{1,2}(T X)$ coincides with the one given in [15]. Vector fields of the form $g \nabla f$ for $f, g \in \operatorname{Test}(\mathrm{X})$ are in $W_{C}^{1,2}(T \mathrm{X})$ and the $W_{C}^{1,2}$-closure of the linear span of such vector fields will be denoted $H_{C}^{1,2}(T \mathrm{X})$. The space $H_{C, \text { loc }}^{1,2}(T \mathrm{X})$ is then equivalently defined either as the subspace of $L_{\text {loc }}^{2}(T \mathrm{X})$ made of vectors $X$ of such that $f X \in H_{C}^{1,2}(T X)$ for every $f \in \operatorname{Test}(\mathrm{X})$ with bounded support or as the $W_{C, \text { loc }}^{1,2}$-closure of $H_{C}^{1,2}(T \mathrm{X})$, i.e. as the space of vector fields $X \in W_{C, \text { loc }}^{1,2}(T \mathrm{X})$ such that there is $\left(X_{n}\right) \subset H_{C}^{1,2}(T \mathrm{X})$ such that $X_{n} \rightarrow X$ and $\nabla X_{n} \rightarrow \nabla X$ in $L_{\text {loc }}^{2}(T \mathrm{X})$ and $L_{\text {loc }}^{2}\left(T^{\otimes 2} \mathrm{X}\right)$ as $n \rightarrow \infty$.

The 'local versions' of the exterior differential, codifferential and Hodge Laplacian are defined in the same way. For us it will be relevant to know that

$$
\begin{equation*}
f \in D\left(\Delta_{\text {loc }}\right) \quad \Leftrightarrow \quad \mathrm{d} f \in D\left(\Delta_{H, \text { loc }}\right) \quad \text { and in this case } \quad \mathrm{d} \Delta f=-\Delta_{H} \mathrm{~d} f \tag{2.2.7}
\end{equation*}
$$

where the minus sign is due to the usual sign convention $\Delta f=-\Delta_{H} f$; this is a direct consequence of the analogous identity valid in for objects in $D(\Delta), D\left(\Delta_{H}\right)$ and a cut-off argument. Also, we shall use the fact that if $X^{b} \in D\left(\Delta_{H, \text { loc }}\right)$ then $X \in W_{\text {loc }}^{1,2}(T \mathrm{X})$ and the Bochner inequality

$$
\begin{equation*}
\Delta \frac{|X|^{2}}{2} \geq|\nabla X|_{\mathrm{HS}}^{2}-\left\langle X^{b}, \Delta_{H} X^{b}\right\rangle+K|X|^{2} \tag{2.2.8}
\end{equation*}
$$

holds in the weak form, i.e.:

$$
\int \frac{|X|^{2}}{2} \Delta g \mathrm{~d} \mathfrak{m} \geq \int g\left(|\nabla X|_{\mathrm{HS}}^{2}-\left\langle X^{b}, \Delta_{H} X^{b}\right\rangle+K|X|^{2}\right) \mathrm{d} \mathfrak{m}
$$

for every $g \in \operatorname{Test}(\mathrm{X})$ non-negative and with bounded support. As before, this follows with a cut-off argument starting with the analogous inequalities valid for the various objects in $L^{2}$. We remark that 2.2.7) for a vector field $X \in D\left(\Delta_{H}\right)$ implies, in particular, by integration that

$$
\begin{equation*}
\int|\nabla X|_{\mathrm{HS}}^{2} \mathrm{~d} \mathfrak{m} \leq \int\left\langle X^{b}, \Delta_{H} X^{b}\right\rangle-K|X|^{2} \mathrm{~d} \mathfrak{m} \tag{2.2.9}
\end{equation*}
$$

and recall that

$$
X^{b} \in D\left(\Delta_{H}\right) \text { with } \Delta_{H} X=0 \quad \Leftrightarrow \quad\left\{\begin{array}{c}
X^{b} \in D(\mathrm{~d}), X \in D(\text { div })  \tag{2.2.10}\\
\text { with } \mathrm{d} X^{b}=0, \operatorname{div} X=0
\end{array}\right.
$$

indeed $\Leftarrow$ follows from the definition of $\Delta_{H}$ and $\Rightarrow$ from the identity

$$
\int\left\langle X^{b}, \Delta_{H} X^{b}\right\rangle \mathrm{d} \mathfrak{m}=\int\left|\mathrm{d} X^{b}\right|^{2}+|\operatorname{div} X|^{2} \mathrm{~d} \mathfrak{m}
$$

### 2.3 Regular Lagrangian Flows and continuity equation

Here we recall the concept of Regular Lagrangian Flow of Sobolev vector fields on RCD spaces, which provides a non-smooth analogous of the concept of solution of the ODE

$$
\gamma_{t}^{\prime}=X_{t}\left(\gamma_{t}\right)
$$

This notion has been introduced by Ambrosio ([1) in the Euclidean space in the context of Di Perna-Lions theory ([12]). Then Ambrosio-Trevisan ([8], see also [9]) showed that theory could be developed in RCD spaces.

Definition 2.6 (Regular Lagrangian flow). Let $\left(X_{t}\right) \in L^{2}\left([0,1], L_{\text {loc }}^{2}(T X)\right)$. We say that $\mathrm{Fl}^{\left(X_{t}\right)}:[0,1] \times \mathrm{X} \rightarrow \mathrm{X}$ is a Regular Lagrangian Flow for $\left(X_{t}\right)$ provided:
i) There is $C>0$ such that

$$
\begin{equation*}
\left(\mathrm{Fl}_{s}^{\left(X_{t}\right)}\right)_{*} \mathfrak{m} \leq C \mathfrak{m} \quad \forall s \in[0,1] . \tag{2.3.1}
\end{equation*}
$$

ii) For $\mathfrak{m}$-a.e. $x \in \mathrm{X}$ the curve $[0,1] \ni s \mapsto \mathrm{Fl}_{s}^{\left(X_{t}\right)}(x) \in \mathrm{X}$ is continuous and such that $\mathrm{Fl}_{0}^{\left(X_{t}\right)}(x)=x$.
iii) for every $f \in W^{1,2}(\mathrm{X})$ we have: for $\mathfrak{m}$-a.e. $x \in \mathrm{X}$ the function $s \mapsto f\left(\mathrm{Fl}_{s}^{\left(X_{t}\right)}(x)\right)$ belongs to $W^{1,1}(0,1)$ and it holds

$$
\begin{equation*}
\frac{\mathrm{d}}{\mathrm{~d} s} f\left(\mathrm{Fl}_{s}^{\left(X_{t}\right)}(x)\right)=\mathrm{d} f\left(X_{s}\right)\left(\mathrm{Fl}_{s}^{\left(X_{t}\right)}(x)\right) \quad \mathfrak{m} \times\left.\mathcal{L}^{1}\right|_{[0,1]} \text {-a.e. }(x, s) \tag{2.3.2}
\end{equation*}
$$

We shall also deal with Regular Lagrangian Flows for $X_{t} \equiv X$ : in this case, we shall mainly consider flows defined on the whole $[0, \infty)$ rather than only on $[0,1]$ and the various statements below should be read with this implicit assumption in mind.

Notice that it is due to property (i) that property (iii) makes sense. Indeed, for given $X_{s} \in L^{2}(T \mathrm{X})$ and $f \in W^{1,2}(\mathrm{X})$ the function $\mathrm{d} f\left(X_{s}\right) \in L^{1}(\mathrm{X})$ is only defined $\mathfrak{m}$-a.e., so that (part of) the role of 2.3 .1 is to grant that $\mathrm{d} f\left(X_{s}\right) \circ F_{s}$ is well defined $\mathfrak{m}$-a.e.. Also, it is known (see Lemmas 7.4 and 9.2 in [8]) that for $\mathfrak{m}$-a.e. $x$ the curve $s \mapsto \mathrm{Fl}_{s}^{\left(X_{t}\right)}(x)$ is absolutely continuous and for its metric speed $\left|\dot{\mathrm{F}}_{s}^{\left(X_{t}\right)}(x)\right|$ we have

$$
\begin{equation*}
\left|\dot{\mathrm{Fl}}_{s}^{\left(X_{t}\right)}(x)\right|=\left|X_{s}\right|\left(\mathrm{Fl}_{s}^{\left(X_{t}\right)}(x)\right) \quad \text { a.e. } s \in[0,1] . \tag{2.3.3}
\end{equation*}
$$

We shall always work with vector fields $\left(X_{t}\right)$ such that

$$
\begin{equation*}
\left|X_{t}\right| \in L^{\infty}\left([0,1], L^{\infty}(\mathrm{X})\right) \tag{2.3.4}
\end{equation*}
$$

In this case, a simple property, valid for any $p \in[1, \infty)$, of Regular Lagrangian Flows that we shall occasionally use is the following:

$$
\begin{equation*}
f_{s} \rightarrow f \quad \text { in } L^{p}(\mathrm{X}) \text { as } s \rightarrow 0 \quad \Rightarrow \quad f_{s} \circ \mathrm{Fl}_{s}^{\left(X_{t}\right)} \rightarrow f \quad \text { in } L^{p}(\mathrm{X}) \text { as } s \rightarrow 0 . \tag{2.3.5}
\end{equation*}
$$

Indeed, for any Lipschitz function $\tilde{f}$ with bounded support we have

$$
\begin{aligned}
& \left\|f_{s} \circ \mathrm{Fl}_{s}^{\left(X_{t}\right)}-f\right\|_{L^{p}} \leq\left\|f_{s} \circ \mathrm{Fl}_{s}^{\left(X_{t}\right)}-\tilde{f} \circ \mathrm{Fl}_{s}^{\left(X_{t}\right)}\right\|_{L^{p}}+\left\|\tilde{t} \circ \mathrm{Fl}_{s}^{\left(X_{t}\right)}-\tilde{f}\right\|_{L^{p}}+\|\tilde{f}-f\|_{L^{p}} \\
& \text { by (2.3.1) } \leq\left(C^{1 / p}+1\right)\left\|f_{t}-\tilde{f}\right\|_{L^{p}}+\left\|\tilde{f} \circ \mathrm{Fl}_{t}^{\left(X_{t}\right)}-\tilde{f}\right\|_{L^{p} .} .
\end{aligned}
$$

Since for $\mathfrak{m}$-a.e. $x \in \mathrm{X}$ the curve $s \mapsto \mathrm{Fl}_{s}^{\left(X_{t}\right)}(x)$ is Lipschitz and with metric speed bounded above by $\left\|\left|X_{t}\right|\right\|_{L_{t}^{\infty}\left(L_{x}^{\infty}\right)}$, we have

$$
\left|\tilde{f}\left(\mathrm{Fl}_{t}^{X}(x)\right)-\tilde{f}(x)\right| \leq|t| \operatorname{Lip}(\tilde{f})\left\|\left|X_{t}\right|\right\|_{L_{t}^{\infty}\left(L_{x}^{\infty}\right)} \quad \mathfrak{m}-\text { a.e. } x \in \mathrm{X}
$$

hence letting $t \rightarrow 0$ in the above we obtain

$$
\varlimsup_{t \rightarrow 0}\left\|f_{t} \circ \mathrm{Fl}_{t}^{X}-f\right\|_{L^{2}} \leq\left(C^{1 / p}+1\right)\|\tilde{f}-f\|_{L^{p}},
$$

so that (2.3.5) follows from the arbitrariness of $\tilde{f}$ and the density of Lipschitz functions with bounded support in $L^{p}(\mathrm{X})$.

We shall mainly use Regular Lagrangian Flows via the following characterization:
Proposition 2.7. Let $\left(X_{t}\right) \in L^{2}\left([0,1], L_{\text {loc }}^{2}(T \mathrm{X})\right)$ be such that 2.3.4 holds and $F$ : $[0,1] \times \mathrm{X} \rightarrow \mathrm{X}$ be a Borel map satisfying (i), (ii) of Definition 2.6. Then the following are equivalent:
a) (iii) of Definition 2.6 holds, i.e. $F$ is a Regular Lagrangian flow for $\left(X_{t}\right)$.
b) for every $f \in W^{1,2}(\mathrm{X})$ the map $[0,1] \ni t \mapsto f \circ F_{t} \in L^{2}(\mathrm{X})$ is Lipschitz and for a.e. $t \in[0,1]$ it holds

$$
\begin{equation*}
\lim _{h \rightarrow 0} \frac{f \circ F_{t+h}-f \circ F_{t}}{h}=\mathrm{d} f\left(X_{t}\right) \circ F_{t}, \tag{2.3.6}
\end{equation*}
$$

the limit being intended in $L^{2}(\mathrm{X})$.
c) for every $f \in W_{\mathrm{loc}}^{1,2}(\mathrm{X})$ the map $[0,1] \ni t \mapsto f \circ F_{t} \in L_{\mathrm{loc}}^{2}(\mathrm{X})$ is Lipschitz and for a.e. $t \in[0,1]$ 2.3.6 holds with the limit being intended in $L_{\text {loc }}^{2}$.

Moreover, if these holds and $X_{t} \equiv X$, then 'Lipschitz' in (b), (c) can be replaced by ' $C$ ', and (2.3.6) holds for every $t \in[0,1]$.
proof
$(\mathbf{a}) \Rightarrow(\mathbf{b})$ From 2 2.3.2 and Fubini's theorem we have that for $\mathfrak{m}$-a.e. $x$ and $\left(\left.\mathcal{L}^{1}\right|_{[0,1]}\right)^{2}$-a.e. $\left(s_{0}, s_{1}\right)$ it holds

$$
\begin{equation*}
f\left(\mathrm{Fl}_{s_{1}}^{\left(X_{t}\right)}(x)\right)-f\left(\mathrm{Fl}_{s_{0}}^{\left(X_{t}\right)}(x)\right)=\int_{s_{0}}^{s_{1}} \mathrm{~d} f\left(X_{s}\right)\left(\mathrm{Fl}_{s}^{\left(X_{t}\right)}(x) \mathrm{d} s .\right. \tag{2.3.7}
\end{equation*}
$$

By the uniform bound (2.3.4) and 2.3.1) we deduce that $\left(\mathrm{d} f(X.) \circ \mathrm{Fl}^{\left(X_{t}\right)}\right) \in L^{\infty}\left([0,1], L^{2}(\mathrm{X})\right)$, and thus the Bochner integral $\int_{s_{0}}^{s_{1}} \mathrm{~d} f\left(X_{s}\right) \circ \mathrm{Fl}_{s}^{\left(X_{t}\right)} \mathrm{d} s$ is a well defined function in $L^{2}(\mathrm{X})$ which vary continuously in $s_{0}, s_{1}$. By 2.3 .5 we also deduce that $s \mapsto f \circ \mathrm{Fl}_{s}^{\left(X_{t}\right)} \in L^{2}(\mathrm{X})$ is continuous, thus from 2.3.7) we obtain that

$$
f \circ \mathrm{Fl}_{s_{1}}^{\left(X_{t}\right)}-f \circ \mathrm{Fl}_{s_{0}}^{\left(X_{t}\right)}=\int_{s_{0}}^{s_{1}} \mathrm{~d} f\left(X_{s}\right) \circ \mathrm{Fl}_{s}^{\left(X_{t}\right)} \mathrm{d} s, \quad \forall s_{0}, s_{1} \in[0,1], s_{0}<s_{1},
$$

where the identity is intended in $L^{2}(\mathrm{X})$ and the integral in the right hand side is the Bochner one. The Lipschitz continuity of $s \mapsto f \circ \mathrm{Fl}_{s}^{\left(X_{t}\right)} \in L^{2}(\mathrm{X})$ and 2.3.6 follow.
$(\mathbf{b}) \Rightarrow(\mathbf{c})$ The assumption $(2.3 .4)$ together with $(2.3 .3)$ grant finite speed of propagation. Then the claim follows by a simple cut-off argument. (c) $\Rightarrow$ (a) By assumption for every bounded set $B \subset \mathrm{X}$ we have

$$
\chi_{B}\left(f \circ \mathrm{Fl}_{s_{1}}^{\left(X_{t}\right)}-f \circ \mathrm{Fl}_{s_{0}}^{\left(X_{t}\right)}\right)=\int_{s_{0}}^{s_{1}} \chi_{B} \mathrm{~d} f\left(X_{s}\right) \circ \mathrm{Fl}_{s}^{\left(X_{t}\right)} \mathrm{d} s, \quad \forall s_{0}, s_{1} \in[0,1], s_{0}<s_{1}
$$

and thus from the arbitrariness of $B$ and Fubini's theorem we conclude that for $\mathfrak{m}$-a.e. $x$ it holds
$f\left(\mathrm{Fl}_{s_{1}}^{\left(X_{t}\right)}(x)\right)-f\left(\mathrm{Fl}_{s_{0}}^{\left(X_{t}\right)}(x)\right)=\int_{s_{0}}^{s_{1}} \mathrm{~d} f\left(X_{s}\right)\left(\mathrm{Fl}_{s}^{\left(X_{t}\right)}(x)\right) \mathrm{d} s, \quad \mathcal{L}^{2}-$ a.e. $s_{0}, s_{1} \in[0,1], s_{0}<s_{1}$.

Applying Lemma 2.1 of [3] we deduce that for $\mathfrak{m}$-a.e. $x$ the function $t \mapsto f\left(\mathrm{Fl}_{s_{1}}^{\left(X_{t}\right)}(x)\right)$ belongs to $W^{1,1}(0,1)$ and it is now obvious that its distributional derivative is given by $\mathrm{d} f\left(X_{s}\right)\left(\mathrm{Fl}_{s}^{\left(X_{t}\right)}(x)\right)$, thus concluding the proof.
$\mathbf{C}^{1}$ regularity It is sufficient to prove that $s \mapsto \mathrm{~d} f(X) \circ \mathrm{Fl}_{s}^{(X)} \in L^{2}(\mathrm{X})$ (resp. $L_{\text {loc }}^{2}$ ) is continuous for $f \in W^{1,2}(\mathrm{X})$ (resp. $W_{\mathrm{loc}}^{1,2}(\mathrm{X})$ ). This is a direct consequence of (2.3.5) applied to the functions $f_{s}=\mathrm{d} f(X)$ (resp. $\chi_{B} \mathrm{~d} f(X)$ for $B \subset \mathrm{X}$ Borel and bounded).

With this said, we shall now recall the main result of $[8]$ in the form we will use it:
Theorem 2.8. Let $\left(X_{t}\right) \in L^{2}\left([0,1], W_{C, \text { loc }}^{1,2}(T \mathrm{X})\right) \cap L^{\infty}\left([0,1], L^{\infty}(T \mathrm{X})\right)$ be such that $X_{t} \in$ $D\left(\right.$ div $\left._{\text {loc }}\right)$ for a.e. $t \in[0,1]$, with

$$
\begin{equation*}
\int_{0}^{1}\||\nabla X|\|_{L^{2}(\mathrm{X})}+\left\|\operatorname{div}\left(X_{t}\right)\right\|_{L^{2}(\mathrm{X})}+\left\|\left(\operatorname{div}\left(X_{t}\right)\right)^{-}\right\|_{L^{\infty}(\mathrm{X})} \mathrm{d} t<\infty . \tag{2.3.8}
\end{equation*}
$$

Then a Regular Lagrangian Flow $F_{s}^{\left(X_{t}\right)}$ for $\left(X_{t}\right)$ exists and is unique, in the sense that if $\tilde{F}^{\left(X_{t}\right)}$ is another flow, then for $\mathfrak{m}$-a.e. $x \in \mathrm{X}$ it holds $F_{s}(x)=\tilde{F}_{s}(x)$ for every $s \in[0,1]$. Moreover it holds the quantitative bound

$$
\begin{equation*}
\left(F_{s}^{\left(X_{t}\right)}\right)_{*} \mathfrak{m} \leq \exp \left(\int_{0}^{s}\left\|\left(\operatorname{div}\left(X_{t}\right)\right)^{-}\right\|_{L^{\infty}(\mathrm{X})} \mathrm{d} t\right) \mathfrak{m} \quad \forall s \in[0,1] . \tag{2.3.9}
\end{equation*}
$$

Notice that the uniqueness part of the statement grants in particular that for $X_{t}$ independent on $t$, say $X_{t}=X$, the Regular Lagrangian Flow (that in this situation is well defined for any $t \geq 0$ ) satisfies the semigroup property

$$
\begin{equation*}
\mathrm{Fl}_{t}^{(X)} \circ \mathrm{Fl}_{s}^{(X)}=\mathrm{Fl}_{t+s}^{(X)} \quad \mathfrak{m}-\text { a.e. } \quad \forall t, s \geq 0 \tag{2.3.10}
\end{equation*}
$$

The uniqueness part of Theorem 2.8 is tightly linked to the uniqueness of solutions of the continuity equation ([8],[19):

Definition 2.9 (Solutions of the continuity equation). Let $t \mapsto \mu_{t} \in \mathscr{P}(\mathrm{X})$ and $t \mapsto X_{t} \in$ $L^{0}(T \mathrm{X}), t \in[0,1]$, be Borel maps. We say that they solve the continuity equation

$$
\begin{equation*}
\frac{\mathrm{d}}{\mathrm{~d} t} \mu_{t}+\operatorname{div}\left(X_{t} \mu_{t}\right)=0 \tag{2.3.11}
\end{equation*}
$$

provided:
i) $\mu_{t} \leq C \mathfrak{m}$ for every $t \in[0,1]$ and some $C>0$,
ii) we have

$$
\int_{0}^{1} \int\left|X_{t}\right|^{2} \mathrm{~d} \mu_{t} \mathrm{~d} t<\infty
$$

iii) for any $f \in W^{1,2}(\mathrm{X})$ the map $t \mapsto \int f \mathrm{~d} \mu_{t}$ is absolutely continuous and it holds

$$
\frac{\mathrm{d}}{\mathrm{~d} t} \int f \mathrm{~d} \mu_{t}=\int \mathrm{d} f\left(X_{t}\right) \mathrm{d} \mu_{t} \quad \text { a.e. } t .
$$

The following uniqueness result and representation formula has been proved in [8].
Theorem 2.10 (Uniqueness of solutions of the continuity equation). Let $\left(X_{t}\right)$ be as in Theorem 2.8 and $\bar{\mu} \in \mathscr{P}(\mathrm{X})$ be such that $\mu_{0} \leq C \mathfrak{m}$ for some $C>0$.

Then there exists a unique $\left(\mu_{t}\right)$ such that $\left(\mu_{t}, X_{t}\right)$ solves the continuity equation 2.3.11) in the sense of Definition 2.9 and for which $\mu_{0}=\bar{\mu}$. Moreover, such $\left(\mu_{t}\right)$ is given by

$$
\mu_{s}=\left(\mathrm{Fl}_{s}^{\left(X_{t}\right)}\right)_{*} \bar{\mu} \quad \forall s \in[0,1] .
$$

## 3 Calculus Tools

### 3.1 Local bounded compression/deformation

Here we shall extend the constructions of pullback module and pullback of 1-forms provided in [15] (see also [13]) to maps which are locally of bounded compression/deformation. For this, it is technically convenient to work with $L^{0}$-normed modules rather than with $L^{2}$ normed ones.

### 3.1.1 Pullback module through a map of local bounded compression

Definition 3.1 (Maps of local bounded compression). Let ( $\mathrm{X}_{1}, \mathfrak{m}_{1}$ ) and ( $\mathrm{X}_{2}, \mathfrak{m}_{2}$ ) be two $\sigma$-finite measured spaces. We say that $\varphi: \mathrm{X}_{1} \rightarrow \mathrm{X}_{2}$ is a map of local bounded compression provided for every $B \subset \mathscr{B}\left(\mathrm{X}_{1}\right)$ with $\mathfrak{m}_{1}(B)<+\infty$ there exists a constant $C_{B} \geq 0$ such that

$$
\begin{equation*}
\varphi_{*}\left(\mathfrak{m}_{\left.1\right|_{B}}\right) \leq C_{B} \mathfrak{m}_{2} . \tag{3.1.1}
\end{equation*}
$$

Theorem/Definition 3.2. Let $\left(\mathrm{X}_{1}, \mathfrak{m}_{1}\right)$ and $\left(\mathrm{X}_{2}, \mathfrak{m}_{2}\right)$ be two $\sigma$-finite measured spaces, $\varphi: \mathrm{X}_{1} \rightarrow \mathrm{X}_{2}$ be of local bounded compression and $\mathscr{M}$ a $L^{0}\left(\mathrm{X}_{2}\right)$-normed module. Then there exists a unique couple $\left(\left[\varphi^{*}\right] \mathscr{M},\left[\varphi^{*}\right]\right)$ where $\left[\varphi^{*}\right] \mathscr{M}$ is a $L^{0}\left(\mathrm{X}_{1}\right)$-module and $\left[\varphi^{*}\right]: \mathscr{M} \rightarrow$ [ $\left.\varphi^{*}\right] \mathscr{M}$ is a linear map such that:

1) $\left|\left[\varphi^{*}\right] v\right|=|v| \circ \varphi \mathfrak{m}_{1}$-a.e. for every $v \in \mathscr{M}$,
2) $\left[\varphi^{*}\right] \mathscr{M}$ is generated by $\left\{\left[\varphi^{*}\right](v): v \in \mathscr{M}\right\}$.

Uniqueness is intended up to unique isomorphism, i.e.: if $\left(\widetilde{\left[\varphi^{*}\right]} \mathfrak{M},\left[\tilde{\varphi}^{*}\right]\right)$ is another such couple, then there exists a unique isomorphism $\Phi:\left[\varphi^{*}\right] \mathscr{M} \rightarrow\left[\varphi^{*}\right] \mathscr{M}$ such that $\left[\tilde{\varphi}^{*}\right]=$ $\left[\varphi^{*}\right] \circ \Phi$.
proof
Existence We define the 'pre-pullback' set Ppb as
$\mathrm{Ppb}:=\left\{\left(v_{i}, A_{i}\right)_{i=1, \ldots, n}: n \in \mathbb{N},\left(A_{i}\right) \subset \mathscr{B}\left(\mathrm{X}_{1}\right)\right.$ is a partition of $\mathrm{X}_{1}$ and $\left.v_{i} \in \mathscr{M} \forall i=1, \ldots, n\right\}$
and an equivalence relation on it by declaring $\left(v_{i}, A_{i}\right) \sim\left(w_{j}, B_{j}\right)$ provided

$$
\left|v_{i}-w_{j}\right| \circ \varphi=0 \quad \mathfrak{m}_{1} \text {-a.e. on } A_{i} \cap B_{j} \quad \forall i, j .
$$

It s readily verified that it is actually an equivalence relation on Ppb : we shall denote by $\left[\left(v_{i}, A_{i}\right)\right]$ the equivalence class of $\left(v_{i}, A_{i}\right)$.

We endow $\mathrm{Ppb} / \sim$ with a vector space structure by putting

$$
\begin{aligned}
{\left[\left(v_{i}, A_{i}\right)\right]+\left[\left(w_{j}, B_{j}\right)\right]: } & :\left[\left(v_{i}+w_{j}, A_{i} \cap B_{j}\right)\right] \\
\lambda\left[\left(v_{i}, A_{i}\right)\right]: & =\left[\left(\lambda v_{i}, A_{i}\right)\right]
\end{aligned}
$$

for every $\left[\left(v_{i}, A_{i}\right)\right],\left[\left(w_{j}, B_{j}\right)\right] \in \mathrm{Ppb} / \sim$ and $\lambda \in \mathbb{R}$. Notice that these are well defined. Moreover, we define a pointwise norm on $\mathrm{Ppb} / \sim$ and a multiplication with simple functions by putting

$$
\begin{aligned}
& \left|\left[\left(v_{i}, A_{i}\right)\right]\right|:=\sum_{i} \chi_{A_{i}}\left|v_{i}\right| \circ \varphi \\
& g\left[\left(v_{i}, A_{i}\right)\right]:=\left[\left(\alpha_{j} v_{i}, A_{i} \cap E_{j}\right)\right] \quad \text { for } \quad g=\sum_{j} \alpha_{j} \chi_{E_{j}} .
\end{aligned}
$$

Again, these are easily seen to be well defined; then we fix a partition $\left(E_{i}\right) \subset \mathscr{B}\left(\mathrm{X}_{1}\right)$ of $\mathrm{X}_{1}$ made of sets of finite $\mathfrak{m}_{1}$-measure and define the distance $\mathrm{d}_{0}$ on $\mathrm{Ppb} / \sim$ as

$$
\mathrm{d}_{0}\left(\left[\left(v_{i}, A_{i}\right)\right],\left[\left(w_{j}, B_{j}\right)\right]\right):=\sum_{k \in \mathbb{N}} \frac{1}{2^{k} \mathfrak{m}_{1}\left(E_{k}\right)} \int_{E_{k}}\left|\left[\left(v_{i}, A_{i}\right)\right]-\left[\left(w_{j}, B_{j}\right)\right]\right| \mathrm{d} \mathfrak{m}_{1}
$$

We then define the space $\left[\varphi^{*}\right] \mathscr{M}$ as the completion of ( $\mathrm{Ppb} / \sim, \mathrm{d}_{0}$ ), equipped with the induced topology and the pullback map $\left[\varphi^{*}\right]: \mathscr{M} \rightarrow\left[\varphi^{*}\right] \mathscr{M}$ as $\left[\varphi^{*}\right] v:=\left[v, \mathrm{X}_{1}\right]$. The (in)equalities

$$
\begin{aligned}
\left|\left[\left(v_{i}+w_{j}, A_{i} \cap B_{j}\right)\right]\right| & \leq\left|\left[\left(v_{i}, A_{i}\right)\right]\right|+\left|\left[\left(w_{j}, B_{j}\right)\right]\right|, \\
\left|\lambda\left[\left(v_{i}, A_{i}\right)\right]\right| & =|\lambda|\left|\left[\left(v_{i}, A_{i}\right)\right]\right|, \\
\left|g\left[\left(v_{i}, A_{i}\right)\right]\right| & =|g|\left|\left[\left(v_{i}, A_{i}\right)\right]\right|,
\end{aligned}
$$

valid $\mathfrak{m}_{1}$-a.e. for every $\left[\left(v_{i}, A_{i}\right)\right],\left[\left(w_{j}, B_{j}\right)\right] \in \mathrm{Ppb} / \sim, \lambda \in \mathbb{R}$ and simple function $g$ grant that the vector space structure, the pointwise norm and the multiplication by simple functions can all be extended by continuity to the whole $\left[\varphi^{*}\right] \mathscr{M}$ and it is then clear with these operations such space is a $L^{0}$-normed module.

Property (1) then follows by the very definitions of pullback map and pointwise norm, while property (2) from the fact that $\mathrm{Ppb} / \sim$ is dense in $\left[\varphi^{*}\right] \mathscr{M}$ and the typical element $\left[v_{i}, A_{i}\right]$ of $\mathrm{Ppb} / \sim$ is equal to $\sum_{i} \chi_{A_{i}}\left[\varphi^{*}\right] v_{i}$.
Uniqueness The requirement for $\Phi:\left[\varphi^{*}\right] \mathscr{M} \rightarrow \widetilde{\left[\varphi^{*}\right] \mathscr{M}}$ to be $L^{0}\left(\mathrm{X}_{1}\right)$-linear and such that $\left[\tilde{\varphi}^{*}\right]=\left[\varphi^{*}\right] \circ \Phi$ force the definition

$$
\begin{equation*}
\Phi(V):=\sum_{i} \chi_{A_{i}}\left[\tilde{\varphi}^{*}\right] v_{i} \quad \text { for } \quad V=\sum_{i} \chi_{A_{i}}\left[\varphi^{*}\right] v_{i} . \tag{3.1.2}
\end{equation*}
$$

The identity

$$
|\Phi(V)|=\sum_{i} \chi_{A_{i}}\left|\left[\tilde{\varphi}^{*}\right] v_{i}\right| \stackrel{(1) \text { for }\left[\widetilde{\left.\varphi^{*}\right] / \mathscr{M}}\right.}{=} \sum_{i} \chi_{A_{i}}\left|v_{i}\right| \circ \varphi^{(1) \text { for }\left[\varphi^{*}\right] / \mathscr{M}} \sum_{i} \chi_{A_{i}}\left|\left[\varphi^{*}\right] v_{i}\right|=|V|
$$

shows in particular that the definition of $\Phi(V)$ is well-posed, i.e. it depends only on $V$ and not on the particular way to represent it as sum. It also shows that it preserves the pointwise norm and thus it is continuous. Since the space of $V$ 's of the form $\sum_{i} \chi_{A_{i}}\left[\varphi^{*}\right] v_{i}$ is dense in $\left[\varphi^{*}\right] \mathscr{M}$ (property (2) for $\left[\varphi^{*}\right] \mathscr{M}$ ), such $\Phi$ can be uniquely extended to a continuous map on the whole $\left[\varphi^{*}\right] \mathscr{M}$ and such extension is clearly linear, continuous, and preserves the pointwise norm. By definition, it also holds $\Phi(g V)=g \Phi(V)$ for $g$ simple and $V=$ $\sum_{i} \chi_{A_{i}}\left[\varphi^{*}\right] v_{i}$ and thus by approximation we see that the same holds for general $g \in L^{0}\left(\mathrm{X}_{1}\right)$ and $V \in\left[\varphi^{*}\right] \mathscr{M}$.

It remains to show that the image of $\Phi$ is the whole $\widetilde{\left[\varphi^{*}\right] \mathscr{M}}$ : this follows from the fact that elements of the form $\sum_{i} \chi_{A_{i}}\left[\tilde{\varphi}^{*}\right] v_{i}$, which by definition are in the image of $\Phi$, are dense in $\widetilde{\left.\varphi^{*}\right] \cdot \mathscr{M}}$ by property (2) for $\widetilde{\left[\varphi^{*}\right] \mathscr{M}}$.

We shall now provide an explicit representation of such pullback module in the case when $\varphi$ is a projection.

Thus let $\left(\mathrm{X}_{1}, \mathfrak{m}_{1}\right)$ and $\left(\mathrm{X}_{2}, \mathfrak{m}_{2}\right)$ be two $\sigma$-finite measured spaces and let $\mathscr{M}$ be a $L^{0}\left(\mathrm{X}_{1}\right)$ module over $\mathrm{X}_{1}$.

We shall denote by $L^{0}\left(\mathrm{X}_{2}, \mathscr{M}\right)$ the space of (equivalence classes up to $\mathfrak{m}_{2}$-a.e. equality of) strongly measurable ( $=$ Borel and essentially separably valued) functions from $\mathrm{X}_{2}$ to $\mathscr{M}$ and claim that such space canonically carries the structure of $L^{0}\left(\mathrm{X}_{1} \times \mathrm{X}_{2}\right)$-normed module. The multiplication of an element in $L^{0}\left(\mathrm{X}_{2}, \mathscr{M}\right)$ by a function $f \in L^{0}\left(\mathrm{X}_{1} \times \mathrm{X}_{2}\right)$ is defined as the map $\mathrm{X}_{2} \ni x_{2} \mapsto f\left(\cdot, x_{2}\right) \cdot v\left(\cdot, x_{2}\right) \in \mathscr{M}$. By approximating $f$ in $L^{0}\left(\mathrm{X}_{1} \times \mathrm{X}_{2}\right)$ with functions having finite range it is easy to check that $x_{2} \mapsto f\left(\cdot, x_{2}\right) \cdot v\left(\cdot, x_{2}\right)$ has separable range if $x_{2} \mapsto v\left(\cdot, x_{2}\right)$ does, so that this definition is well posed. Similarly, the pointwise norm of $v \in L^{0}\left(\mathrm{X}_{2}, \mathscr{M}\right)$ is obtained composing the map $x_{2} \mapsto v\left(\cdot, x_{2}\right) \in \mathscr{M}$ with the pointwise norm on $\mathscr{M}$, thus providing an element of $L^{0}\left(\mathrm{X}_{2}, L^{0}\left(\mathrm{X}_{1}\right)\right) \sim L^{0}\left(\mathrm{X}_{1} \times \mathrm{X}_{2}\right)$. Finally, we use such pointwise norm to define the topology of $L^{0}\left(\mathrm{X}_{2}, \mathscr{M}\right)$ as in point (iii) of Definition 2.2.

It is not hard to check that sequences converging in this topology are made of maps $v_{n}\left(\cdot, x_{2}\right)$ which are $\mathfrak{m}_{2}$-a.e. converging and that with these definitions $L^{0}\left(\mathrm{X}_{2}, \mathscr{M}\right)$ is indeed a $L^{0}\left(\mathrm{X}_{1} \times \mathrm{X}_{2}\right)$-normed module.

In what will come next, we shall often implicitly use the following identification:
Proposition $3.3\left(\left[\pi_{1}\right]^{*} \mathscr{M}\right.$ is isomorphic to $\left.L^{0}\left(\mathrm{X}_{2}, \mathscr{M}\right)\right)$. Let $\left(\mathrm{X}_{1}, \mathfrak{m}_{1}\right)$ and $\left(\mathrm{X}_{2}, \mathfrak{m}_{2}\right)$ be two $\sigma$-finite measured spaces, $\mathscr{M}$ a $L^{0}\left(\mathrm{X}_{1}\right)$-module over $\mathrm{X}_{1}$ and $\pi_{1}: \mathrm{X}_{1} \times \mathrm{X}_{2} \rightarrow \mathrm{X}_{1}$ the canonical projection.

Then there exists a unique isomorphism from $\left[\pi_{1}{ }^{*}\right](\mathscr{M})$ to $L^{0}\left(\mathrm{X}_{2}, \mathscr{M}\right)$ which for every $v \in \mathscr{M}$ sends $\left[\pi_{1}^{*}\right] v$ to the function in $L^{0}\left(\mathrm{X}_{2}, \mathscr{M}\right)$ constantly equal to $v$.
proof For $v \in \mathscr{M}$ let $\hat{v} \in L^{0}\left(\mathrm{X}_{2}, \mathscr{M}\right)$ be the function constantly equal to $v$. It is clear that $|\hat{v}|=|v| \circ \pi_{1} \mathfrak{m}_{1} \times \mathfrak{m}_{2}$-a.e.. Moreover, from the fact that functions in $L^{0}\left(\mathrm{X}_{2}, \mathscr{M}\right)$ are essentially separably valued it follows by standard means in vector-space integration that $\{\hat{v}: v \in \mathscr{M}\}$ generate the whole $L^{0}\left(\mathrm{X}_{2}, \mathscr{M}\right)$.

The conclusion comes from Theorem 3.2.

### 3.1.2 Localized pullback of 1-forms

Definition 3.4 (Maps of local bounded deformation). Let ( $\mathrm{X}_{1}, \mathrm{~d}_{1}, \mathfrak{m}_{1}$ ) and ( $\mathrm{X}_{2}, \mathrm{~d}_{2}, \mathfrak{m}_{2}$ ) be two metric measure spaces. We say that a map $\varphi: \mathrm{X}_{1} \rightarrow \mathrm{X}_{2}$ is of local bounded deformation if for every bounded set $B \subset \mathrm{X}_{1}$ there are constants $L(B), C(B)>0$ such that:

$$
\begin{aligned}
& \varphi \text { is } L(B) \text {-Lipschitz on } B \\
& \varphi_{*}\left(\mathfrak{m}_{\left.1\right|_{B}}\right) \leq C(B) \mathfrak{m}_{2} .
\end{aligned}
$$

Recalling that the local Lipschitz constant $\operatorname{lip} \varphi: \mathrm{X}_{1} \rightarrow[0, \infty]$ is defined as

$$
\operatorname{lip} \varphi(x):=\varlimsup_{y \rightarrow x} \frac{\mathrm{~d}_{2}(\varphi(x), \varphi(y))}{\mathrm{d}_{1}(x, y)}
$$

if $x$ is not isolated, 0 otherwise, we have the following simple statement:
Proposition 3.5. Let $\left(\mathrm{X}_{1}, \mathrm{~d}_{1}, \mathfrak{m}_{1}\right)$ and $\left(\mathrm{X}_{2}, \mathrm{~d}_{2}, \mathfrak{m}_{2}\right)$ be two metric measure spaces and $\varphi$ : $\mathrm{X}_{1} \rightarrow \mathrm{X}_{2}$ be a map of local bounded deformation. Then for any $f \in \mathcal{S}_{\text {loc }}^{2}\left(\mathrm{X}_{2}\right)$, we have $f \circ \varphi \in \mathcal{S}_{\text {loc }}^{2}\left(\mathrm{X}_{1}\right)$ and

$$
\begin{equation*}
|\mathrm{d}(f \circ \varphi)| \leq \operatorname{lip} \varphi|\mathrm{d} f| \circ \varphi, \quad \mathfrak{m}_{1}-\text { a.e.. } \tag{3.1.3}
\end{equation*}
$$

proof Fix a point $\bar{x} \in \mathrm{X}_{1}$, let $A_{n} \subset C\left([0,1], \mathrm{X}_{1}\right)$ be defined as

$$
A_{n}:=\left\{\gamma \in C\left([0,1], \mathrm{X}_{1}\right): \gamma_{t} \in B_{n}(\bar{x}) \forall t \in[0,1]\right\}
$$

and notice that $\cup_{n} A_{n}=C\left([0,1], \mathrm{X}_{1}\right)$. Now let $\boldsymbol{\pi}$ be a test plan on $\mathrm{X}_{1}$ and notice that for $n$ sufficiently large the measure $\boldsymbol{\pi}_{n}:=\boldsymbol{\pi}\left(A_{n}\right)^{-1} \boldsymbol{\pi}_{A_{n}}$ is well defined and a test plan. By construction we have

$$
\left(\mathrm{e}_{t}\right)_{*} \varphi_{*} \boldsymbol{\pi}_{n} \leq \boldsymbol{\pi}\left(A_{n}\right)^{-1} C\left(B_{n}(\bar{x})\right) C(\boldsymbol{\pi}) \mathfrak{m}_{2} \quad \forall t \in[0,1],
$$

where $C(\boldsymbol{\pi})$ is such that $\left(\mathrm{e}_{t}\right)_{*} \boldsymbol{\pi} \leq C(\boldsymbol{\pi}) \mathfrak{m}_{1}$ for every $t \in[0,1]$, and taking into account the trivial bound

$$
\begin{equation*}
\left\{\text { metric speed of } t \mapsto \varphi\left(\gamma_{t}\right)\right\} \leq \operatorname{lip} \varphi\left(\gamma_{t}\right)\left|\dot{\gamma}_{t}\right| \quad \text { a.e. } t \tag{3.1.4}
\end{equation*}
$$

we also have

$$
\iint_{0}^{1}\left|\dot{\gamma}_{t}\right|^{2} \mathrm{~d} t \mathrm{~d} \varphi_{*} \boldsymbol{\pi}_{n}(\gamma) \leq \boldsymbol{\pi}\left(A_{n}\right)^{-1} L^{2}\left(B_{n}(\bar{x})\right) \iint_{0}^{1}\left|\dot{\gamma}_{t}\right|^{2} \mathrm{~d} t \mathrm{~d} \boldsymbol{\pi}_{n}(\gamma)
$$

Hence $\varphi_{*} \boldsymbol{\pi}_{n}$ is a test plan on $\mathrm{X}_{2}$ and thus for $\varphi_{*} \boldsymbol{\pi}_{n}$-a.e. $\tilde{\gamma}$ we have that $t \mapsto f(\tilde{\gamma})$ is in $W^{1,1}(0,1)$ with

$$
\left|\frac{\mathrm{d}}{\mathrm{~d} t} f\left(\tilde{\gamma}_{t}\right)\right| \leq\left|\dot{\tilde{\gamma}}_{t}\right||\mathrm{d} f|\left(\tilde{\gamma}_{t}\right)
$$

Recalling (3.1.4) this means that for $\boldsymbol{\pi}_{n}$-a.e. $\gamma$ the map $t \mapsto f\left(\varphi\left(\gamma_{t}\right)\right)$ belongs to $W^{1,1}(0,1)$ with

$$
\begin{equation*}
\left|\frac{\mathrm{d}}{\mathrm{~d} t} f\left(\varphi\left(\gamma_{t}\right)\right)\right| \leq \operatorname{lip} \varphi\left(\gamma_{t}\right)\left|\dot{\gamma}_{t}\right||\mathrm{d} f|\left(\varphi\left(\gamma_{t}\right)\right) \tag{3.1.5}
\end{equation*}
$$

Being this true for every $n \in \mathbb{N}$ sufficiently large, (3.1.5) holds also for $\boldsymbol{\pi}$-a.e. $\gamma$ and since $\operatorname{lip} \varphi|\mathrm{d} f| \circ \varphi \in L_{\text {loc }}^{2}\left(\mathrm{X}_{1}\right)$, by the characterization (2.2.1) of Sobolev functions the proof is completed.

Theorem/Definition 3.6 (Pullback of 1-forms). Let $\left(\mathrm{X}_{1}, \mathrm{~d}_{1}, \mathfrak{m}_{1}\right)$ and $\left(\mathrm{X}_{2}, \mathrm{~d}_{2}, \mathfrak{m}_{2}\right)$ be two metric measure spaces and $\varphi: \mathrm{X}_{1} \rightarrow \mathrm{X}_{2}$ be a map of local bounded deformation.

Then there exists a unique linear and continuous map $\varphi^{*}: L^{0}\left(T^{*} \mathrm{X}_{2}\right) \rightarrow L^{0}\left(T^{*} \mathrm{X}_{1}\right)$ such that

$$
\begin{array}{ll}
\varphi^{*}(\mathrm{~d} f)=\mathrm{d}(f \circ \varphi), & \forall f \in \mathcal{S}_{\text {loc }}^{2}\left(\mathrm{X}_{2}\right), \\
\varphi^{*}(g \omega)=g \circ \varphi \varphi^{*} \omega, & \forall g \in L^{0}\left(\mathrm{X}_{2}\right), \omega \in L^{0}\left(T^{*} \mathrm{X}_{2}\right), \tag{3.1.6}
\end{array}
$$

and such map satisfies

$$
\begin{equation*}
\left|\varphi^{*} \omega\right| \leq \operatorname{lip} \varphi|\omega| \circ \varphi \quad \mathfrak{m}_{1} \text {-a.e., } \quad \forall \omega \in L^{0}\left(T^{*} \mathrm{X}_{2}\right) \tag{3.1.7}
\end{equation*}
$$

proof The requirements (3.1.6) force the definition

$$
\begin{equation*}
\varphi^{*} \omega:=\sum_{i} \chi_{\varphi^{-1}\left(E_{i}\right)} \mathrm{d}\left(f_{i} \circ \varphi\right) \quad \text { for } \quad \omega=\sum_{i} \chi_{E_{i}} \mathrm{~d} f_{i} \tag{3.1.8}
\end{equation*}
$$

for $\left(E_{i}\right)$ finite Borel partition of $\mathrm{X}_{2}$ and $\left(f_{i}\right) \subset \mathcal{S}_{\text {loc }}^{2}\left(\mathrm{X}_{2}\right)$. The bound

$$
\begin{equation*}
\left|\varphi^{*} \omega\right|=\sum_{i} \chi_{\varphi^{-1}\left(E_{i}\right)}\left|\mathrm{d}\left(f_{i} \circ \varphi\right)\right| \stackrel{(3.1 .3)}{\leq} \operatorname{lip} \varphi \sum_{i} \chi_{E_{i}} \circ \varphi\left|\mathrm{~d} f_{i}\right| \circ \varphi=\operatorname{lip} \varphi|\omega| \circ \varphi, \tag{3.1.9}
\end{equation*}
$$

grants both that the definition of $\varphi^{*} \omega$ is well-posed (i.e. its value depends only on $\omega$ and not in how it is written as $\sum_{i} \chi_{E_{i}} \mathrm{~d} f_{i}$ ) and that $\varphi^{*}$ is continuous from the space of $\omega$ 's as in (3.1.8) with the $L^{0}\left(T^{*} \mathrm{X}_{2}\right)$-topology to $L^{0}\left(T^{*} \mathrm{X}_{1}\right)$. Since the class of such $\omega$ 's is dense in $L^{0}\left(T^{*} \mathrm{X}_{2}\right)$, we can be uniquely extend $\varphi^{*}$ to a continuous map from $L^{0}\left(T^{*} \mathrm{X}_{2}\right)$ to $L^{0}\left(T^{*} \mathrm{X}_{1}\right)$.

The resulting extension satisfies the first in (3.1.6) by definition, while (3.1.7) comes from (3.1.9). The second in (3.1.6) for simple functions $g$ is a direct consequence of the definition (3.1.8), then the general case follows by approximation.

### 3.2 Calculus on product spaces

### 3.2.1 Cotangent module and product of spaces

Let $\left(\mathrm{X}_{1}, \mathrm{~d}_{1}, \mathfrak{m}_{1}\right)$ and $\left(\mathrm{X}_{2}, \mathrm{~d}_{2}, \mathfrak{m}_{2}\right)$ be two metric measure spaces. Aim of this section is to relate the cotangent modules of $\mathrm{X}_{1}, \mathrm{X}_{2}$ to that of the product space $\mathrm{X}_{1} \times \mathrm{X}_{2}$, which will be always implicitly endowed with the product measure and the distance

$$
\left(\mathrm{d}_{1} \otimes \mathrm{~d}_{2}\right)^{2}\left(\left(x_{1}, x_{2}\right),\left(y_{1}, y_{2}\right)\right):=\mathrm{d}_{1}^{2}\left(x_{1}, y_{1}\right)+\mathrm{d}_{2}^{2}\left(x_{2}, y_{2}\right) .
$$

Let $\pi_{i}: \mathrm{X}_{1} \times \mathrm{X}_{2} \rightarrow \mathrm{X}_{i}, i=1,2$ be the canonical projections, observe that they are of local bounded deformation and recall from Proposition 3.3 and the discussion before it that $L^{0}\left(\mathrm{X}_{2}, L^{0}\left(T^{*} \mathrm{X}_{1}\right)\right) \sim\left[\pi_{2}^{*}\right] L^{0}\left(T^{*} \mathrm{X}_{1}\right)$ canonically carries the structure of $L^{0}\left(\mathrm{X}_{1} \times \mathrm{X}_{2}\right)$ normed module.

We then start with the following simple and general fact:
Proposition 3.7. Let $\left(\mathrm{X}_{1}, \mathrm{~d}_{1}, \mathfrak{m}_{1}\right)$ and $\left(\mathrm{X}_{2}, \mathrm{~d}_{2}, \mathfrak{m}_{2}\right)$ be two metric measure spaces. Then there exists a unique $L^{0}\left(\mathrm{X}_{1} \times \mathrm{X}_{2}\right)$-linear and continuous map $\Phi_{1}: L^{0}\left(\mathrm{X}_{2}, L^{0}\left(T^{*} \mathrm{X}_{1}\right)\right) \rightarrow$ $L^{0}\left(T^{*}\left(\mathrm{X}_{1} \times \mathrm{X}_{2}\right)\right)$ such that

$$
\begin{equation*}
\Phi_{1}(\widehat{\mathrm{~d} g})=\mathrm{d}\left(g \circ \pi_{1}\right) \quad \forall g \in \mathcal{S}_{\text {loc }}^{2}\left(\mathrm{X}_{1}\right), \tag{3.2.1}
\end{equation*}
$$

where $\widehat{\mathrm{d} g}: \mathrm{X}_{2} \rightarrow L^{0}\left(T^{*} \mathrm{X}_{1}\right)$ is the function identically equal to $\mathrm{d} g$. Such map preserves the pointwise norm.

Similarly, there is a unique $L^{0}\left(\mathrm{X}_{1} \times \mathrm{X}_{2}\right)$-linear and continuous map $\Phi_{2}: L^{0}\left(\mathrm{X}_{1}, L^{0}\left(T^{*} \mathrm{X}_{2}\right)\right) \rightarrow$ $L^{0}\left(T^{*}\left(\mathrm{X}_{1} \times \mathrm{X}_{2}\right)\right)$ such that

$$
\Phi_{2}(\widehat{\mathrm{~d} h})=\mathrm{d}\left(h \circ \pi_{1}\right) \quad \forall h \in \mathcal{S}_{\text {loc }}^{2}\left(\mathrm{X}_{2}\right),
$$

where $\widehat{\mathrm{d} h}: \mathrm{X}_{1} \rightarrow L^{0}\left(T^{*} \mathrm{X}_{2}\right)$ is the function identically equal to $\mathrm{d} h$, and such map preserves the pointwise norm.
proof We shall prove the claims for $\Phi_{1}$, as then the ones for $\Phi_{2}$ follow by symmetry. The required $L^{0}$-linearity and (3.2.1) force the definition

$$
\begin{equation*}
\Phi_{1}(W):=\sum_{i, j} \chi_{A_{i}} \circ \pi_{1} \chi_{B_{j}} \circ \pi_{2} \mathrm{~d}\left(g_{i, j} \circ \pi_{1}\right) \quad \text { for } \quad W=\sum_{i, j} \chi_{B_{j}}\left(\chi_{A_{i}} \mathrm{~d} g_{i, j}\right) \tag{3.2.2}
\end{equation*}
$$

where $\left(A_{i}\right),\left(B_{j}\right)$ are finite Borel partitions of $\mathrm{X}_{1}, \mathrm{X}_{2}$ respectively and $\left(g_{i, j}\right) \subset \mathcal{S}_{\text {loc }}^{2}\left(\mathrm{X}_{1}\right)$. Since $\pi_{1}$ is 1-Lipschitz we have

$$
\begin{equation*}
\left|\Phi_{1}(W)\right|=\sum_{i} \chi_{B_{i}} \circ \pi_{2}\left|\mathrm{~d}\left(g_{i} \circ \pi_{1}\right)\right| \stackrel{\sqrt[3.1 .3]{\leq}}{\leq} \sum_{i} \chi_{B_{i}} \circ \pi_{2}\left|\mathrm{~d} g_{i}\right| \circ \pi_{1}=|W| \tag{3.2.3}
\end{equation*}
$$

which shows both that 3.2 provides a good definition for $\Phi_{1}(W)$, in the sense that $\Phi_{1}(W)$ depends only on $W$ and not on the way we write it as $\sum_{i, j} \chi_{B_{j}}\left(\chi_{A_{i}} \mathrm{~d} g_{i, j}\right)$, and that it is continuous. The definition also ensures that $\Phi_{1}(f W)=f \Phi_{1}(W)$ for $f \in L^{0}\left(\mathrm{X}_{1} \times \mathrm{X}_{2}\right)$ of the form $\sum_{i, j} \alpha_{i, j} \chi_{A_{i}} \chi_{B_{j}}$ for $\left(A_{i}\right),\left(B_{j}\right)$ finite Borel partitions of $\mathrm{X}_{1}, \mathrm{X}_{2}$ respectively and $\left(\alpha_{i, j}\right) \subset \mathbb{R}$.

Since these functions are dense in $L^{0}\left(\mathrm{X}_{1} \times \mathrm{X}_{2}\right)$ and $W$ 's as in (3.2.2) are dense in $L^{0}\left(\mathrm{X}_{1}, L^{0}\left(T^{*} \mathrm{X}_{2}\right)\right)$, this is enough to show existence and uniqueness of a $L^{0}$-linear and continuous $\Phi_{1}$ for which (3.2.1) holds and, from (3.2.3), that for such $\Phi_{1}$ we have

$$
\begin{equation*}
\left|\Phi_{1}(W)\right| \leq|W| \quad \mathfrak{m}_{1} \times \mathfrak{m}_{2}-\text { a.e.. } \tag{3.2.4}
\end{equation*}
$$

Thus to conclude it is sufficient to show that equality holds and, by the very same arguments just given, to this aim it is sufficient to show that

$$
\begin{equation*}
|\mathrm{d} g| \circ \pi_{1}=\left|\mathrm{d}\left(g \circ \pi_{1}\right)\right| \quad \mathfrak{m}_{1} \times \mathfrak{m}_{2}-\text { a.e. } \quad \forall g \in \mathcal{S}_{\mathrm{loc}}^{2}\left(\mathrm{X}_{1}\right) \tag{3.2.5}
\end{equation*}
$$

It is now convenient to consider the map

$$
\begin{aligned}
\mathrm{T}: C\left([0,1], \mathrm{X}_{1}\right) \times \mathrm{X}_{2} & \rightarrow C\left([0,1], \mathrm{X}_{1} \times \mathrm{X}_{2}\right) \\
\left(t \rightarrow \gamma_{t}, x_{2}\right) & \mapsto
\end{aligned} t \rightarrow\left(\gamma_{t}, x_{2}\right) .
$$

Notice that
for any $x_{2} \in \mathrm{X}_{2}$ the speed of $T\left(\gamma, x_{2}\right)$ is equal to the speed of $\gamma$ for a.e. $t$
and fix $\mu \in \mathscr{P}\left(\mathrm{X}_{2}\right)$ such that $\mu \leq \tilde{C} \mathfrak{m}_{2}$ for some $\tilde{C}>0$. Then for an arbitrary test plan $\boldsymbol{\pi}$ on $\mathrm{X}_{1}$ define

$$
\tilde{\boldsymbol{\pi}}:=\mathrm{T}_{*}(\boldsymbol{\pi} \times \mu) \in \mathscr{P}\left(C\left([0,1], \mathrm{X}_{1} \times \mathrm{X}_{2}\right)\right)
$$

and observe that $\tilde{\boldsymbol{\pi}}$ is a test plan on $\mathrm{X}_{1} \times \mathrm{X}_{2}$. Hence for any $g \in \mathcal{S}_{\text {loc }}^{2}\left(\mathrm{X}_{1}\right)$ we have

$$
\begin{aligned}
\int\left|g \circ e_{1}-g \circ e_{0}\right| \mathrm{d} \boldsymbol{\pi} & =\int\left|g \circ \pi_{1} \circ e_{1}-g \circ \pi_{1} \circ e_{0}\right| \mathrm{d} \tilde{\boldsymbol{\pi}} \\
& \leq \int_{0}^{1} \int\left|\mathrm{~d}\left(g \circ \pi_{1}\right)\right|\left(\tilde{\gamma}_{t}\right)\left|\dot{\tilde{\gamma}}_{t}\right| \mathrm{d} \tilde{\boldsymbol{\pi}}(\tilde{\gamma}) \mathrm{d} t \\
& =\int_{0}^{1} \int\left(\int\left|\mathrm{~d}\left(g \circ \pi_{1}\right)\right|\left(\mathrm{T}\left(\gamma, x_{2}\right)_{t}\right) \mathrm{d} \mu\left(x_{2}\right)\right)\left|\dot{\gamma}_{t}\right| \mathrm{d} \boldsymbol{\pi}(\gamma) \mathrm{d} t,
\end{aligned}
$$

so that the arbitrariness of $\boldsymbol{\pi}$ gives that the function $\int\left|\mathrm{d}\left(g \circ \pi_{1}\right)\right|\left(\cdot, x_{2}\right) \mathrm{d} \mu\left(x_{2}\right)$ is a weak upper gradient of $g$. Therefore for $\mathfrak{m}_{1}$-a.e. $x_{1}$ we have

$$
|\mathrm{d} g|\left(x_{1}\right) \leq \int\left|\mathrm{d}\left(g \circ \pi_{1}\right)\right|\left(x_{1}, x_{2}\right) \mathrm{d} \mu\left(x_{2}\right) \stackrel{\sqrt{3.1 .3}}{\leq} \int|\mathrm{d} g| \circ \pi_{1}\left(x_{1}, x_{2}\right) \mathrm{d} \mu\left(x_{2}\right)=|\mathrm{d} g|\left(x_{1}\right) .
$$

Hence the inequalities are equalities and the arbitrariness of $\mu$ gives 3.2.5.
It seems hard to obtain any further relation between the cotangent modules in full generality, for this reason from now on we shall make two structural assumptions. In the following, for any function $f\left(x_{1}, x_{2}\right)$ on the product space $\mathrm{X}_{1} \times \mathrm{X}_{2}$, we define $f^{x_{1}}(\cdot):=$ $f\left(x_{1}, \cdot\right)$ and, similarly, $f_{x_{2}}(\cdot):=f\left(\cdot, x_{2}\right)$.

Definition 3.8 (Tensorization of the Cheeger energy). We say that two metric measure spaces $\left(\mathrm{X}_{1}, \mathrm{~d}_{1}, \mathfrak{m}_{1}\right)$ and $\left(\mathrm{X}_{2}, \mathrm{~d}_{2}, \mathfrak{m}_{2}\right)$ have the property of tensorization of the Cheeger energy provided for any $f \in L^{2}\left(\mathrm{X}_{1} \times \mathrm{X}_{2}\right)$ the following holds: $f \in W^{1,2}\left(\mathrm{X}_{1} \times \mathrm{X}_{2}\right)$ if and only if

$$
\begin{aligned}
& -f^{x_{1}} \in W^{1,2}\left(\mathrm{X}_{2}\right) \text { for } \mathfrak{m}_{1} \text {-a.e. } x_{1} \in \mathrm{X}_{1} \text { and } \iint\left|\mathrm{d} f^{x_{1}}\right|^{2} \mathrm{dm}_{2} \mathrm{~d}_{1}\left(x_{1}\right)<\infty \\
& -f_{x_{2}} \in W^{1,2}\left(\mathrm{X}_{1}\right) \text { for } \mathfrak{m}_{2} \text {-a.e. } x_{2} \in \mathrm{X}_{2} \text { with } \iint\left|\mathrm{d} f_{x_{2}}\right|^{2} \mathrm{~d}_{1} \mathrm{~d}_{2}\left(x_{2}\right)<\infty
\end{aligned}
$$

and in this case it holds

$$
\begin{equation*}
|\mathrm{d} f|^{2}\left(x_{1}, x_{2}\right)=\left|\mathrm{d} f^{x_{1}}\right|^{2}\left(x_{2}\right)+\left|\mathrm{d} f_{x_{2}}\right|^{2}\left(x_{1}\right) \quad \mathfrak{m}_{1} \times \mathfrak{m}_{2} \text { - a.e. }\left(x_{1}, x_{2}\right) \tag{3.2.7}
\end{equation*}
$$

Notice that for $g \in L^{\infty} \cap W^{1,2}\left(\mathrm{X}_{1}\right)$ and $h \in L^{\infty} \cap W^{1,2}\left(\mathrm{X}_{2}\right)$ both with bounded support, the function $g \circ \pi_{1} h \circ \pi_{2}$ has bounded support and is in $L^{\infty} \cap W^{1,2}\left(\mathrm{X}_{1} \times \mathrm{X}_{2}\right)$, its differential being given by

$$
\begin{equation*}
\mathrm{d}\left(g \circ \pi_{1} h \circ \pi_{2}\right)=g \circ \pi_{1} \mathrm{~d}\left(h \circ \pi_{2}\right)+h \circ \pi_{2} \mathrm{~d}\left(g \circ \pi_{1}\right) . \tag{3.2.8}
\end{equation*}
$$

Definition 3.9 (Density of the product algebra). We say that two metric measure spaces $\left(\mathrm{X}_{1}, \mathrm{~d}_{1}, \mathfrak{m}_{1}\right)$ and $\left(\mathrm{X}_{2}, \mathrm{~d}_{2}, \mathfrak{m}_{2}\right)$ have the property of density of product algebra if the set
$\mathcal{A}:=\left\{\sum_{j=1}^{n} g_{j} \circ \pi_{1} h_{j} \circ \pi_{2}: n \in \mathbb{N}, \begin{array}{l}g_{j} \in L^{\infty} \cap W^{1,2}\left(\mathrm{X}_{1}\right) \text { has bounded support } \\ h_{j} \in L^{\infty} \cap W^{1,2}\left(\mathrm{X}_{2}\right) \text { has bounded support }\end{array} \forall j=1, \ldots, n\right\}$
is dense in $W^{1,2}\left(\mathrm{X}_{1} \times \mathrm{X}_{2}\right)$ in the strong topology of $W^{1,2}\left(\mathrm{X}_{1} \times \mathrm{X}_{2}\right)$.
From now on we will always assume the following:
Assumption 3.10. $\left(\mathrm{X}_{1}, \mathrm{~d}_{1}, \mathfrak{m}_{1}\right)$ and $\left(\mathrm{X}_{2}, \mathrm{~d}_{2}, \mathfrak{m}_{2}\right)$ are two metric measure spaces for which both the tensorization of Cheeger energy and the density of the product algebra hold.

It is worth to underline that no couple of spaces $\mathrm{X}_{1}, \mathrm{X}_{2}$ are known for which Assumption 3.10 does not hold. On the other hand, it is unclear if that holds in full generality. The first results about the tensorization of Cheeger energies being given in [5] for the cases of two RCD spaces with finite mass, for our purposes the following result covers the cases of interest:

Proposition 3.11. Let $\left(\mathrm{X}_{1}, \mathrm{~d}_{1}, \mathfrak{m}_{1}\right)$ be a $\operatorname{RCD}(K, \infty)$ spaces and let $\mathrm{X}_{2}$ be the Euclidean space $\mathbb{R}^{N}$ equipped with the Euclidean distance and the Lebesgue measure.

Then both the tensorization of the Cheeger energy and the density of the product algebra hold.
proof In [18] it has been proved that for arbitrary $\left(\mathrm{X}_{1}, \mathrm{~d}_{1}, \mathfrak{m}_{1}\right)$ and for $\mathrm{X}_{2}=\mathbb{R}$ the tensorization of the Cheeger energy holds and the algebra $\mathcal{A}$ is dense in energy, i.e.: for any $f \in W^{1,2}\left(\mathrm{X}_{1} \times \mathbb{R}\right)$ there is $\left(f_{n}\right) \subset \mathcal{A}$ such that $f_{n} \rightarrow f$ and $\left|\mathrm{d} f_{n}\right| \rightarrow|\mathrm{d} f|$ in $L^{2}\left(\mathrm{X}_{1} \times \mathbb{R}\right)$.

If $\left(\mathrm{X}_{1}, \mathrm{~d}_{1}, \mathfrak{m}_{1}\right)$ is infinitesimally Hilbertian (which is the case for RCD spaces), then the tensorization of the Cheeger energy ensures that $W^{1,2}\left(\mathrm{X}_{1} \times \mathbb{R}\right)$ is a Hilbert space, so that the uniform convexity of the norm grants that convergence in energy implies strong $W^{1,2}$-convergence.

Thus the thesis is true for $\mathrm{X}_{2}=\mathbb{R}$. The general case follows by a simple induction argument.

With this said, we shall now continue the investigation of the relation between cotangent modules and products of spaces. We start with the following result; notice that the density of the product algebra is used to show that for $f \in \mathcal{S}_{\text {loc }}^{2}\left(\mathrm{X}_{1} \times \mathrm{X}_{2}\right)$ the map $x_{2} \mapsto \mathrm{~d} f_{x_{2}} \in$ $L^{0}\left(T^{*} \mathrm{X}_{1}\right)$ is essentially separably valued.

Lemma 3.12. Let $\left(\mathrm{X}_{1}, \mathrm{~d}_{1}, \mathfrak{m}_{1}\right)$ and $\left(\mathrm{X}_{2}, \mathrm{~d}_{2}, \mathfrak{m}_{2}\right)$ be two metric measure spaces satisfying Assumption 3.10 .

Then for every $f \in \mathcal{S}_{\text {loc }}^{2}\left(\mathrm{X}_{1} \times \mathrm{X}_{2}\right)$ we have that $f_{x_{2}} \in \mathcal{S}_{\text {loc }}^{2}\left(\mathrm{X}_{1}\right)$ for $\mathfrak{m}_{2}$-a.e. $x_{2}$ and the map $x_{2} \mapsto \mathrm{~d} f_{x_{2}}$ belongs to $L^{0}\left(\mathrm{X}_{2}, L^{0}\left(T^{*} \mathrm{X}_{1}\right)\right)$. Moreover, for $\left(f_{n}\right) \subset \mathcal{S}_{\text {loc }}^{2}\left(\mathrm{X}_{1} \times \mathrm{X}_{2}\right)$ we have

$$
\begin{equation*}
\mathrm{d} f_{n} \rightarrow \mathrm{~d} f \text { in } L^{0}\left(T^{*}\left(\mathrm{X}_{1} \times \mathrm{X}_{2}\right)\right) \quad \Rightarrow \quad \mathrm{d}\left(f_{n}\right) . \rightarrow \mathrm{d} f . \text { in } L^{0}\left(\mathrm{X}_{2}, L^{0}\left(T^{*} \mathrm{X}_{1}\right)\right) \tag{3.2.10}
\end{equation*}
$$

Similarly for the roles of $\mathrm{X}_{1}$ and $\mathrm{X}_{2}$ inverted. Finally, the identity (3.2.7) holds for any $f \in \mathcal{S}_{\text {loc }}^{2}\left(\mathrm{X}_{1} \times \mathrm{X}_{2}\right)$.
proof Let $f=g \circ \pi_{1} h \circ \pi_{2}$ for some $g \in L^{\infty} \cap W^{1,2}\left(\mathrm{X}_{1}\right)$ and $h \in L^{\infty} \cap W^{1,2}\left(\mathrm{X}_{2}\right)$ with bounded supports and notice that $\mathrm{d} f_{x_{2}}=h\left(x_{2}\right) \mathrm{d} g$ for every $x_{2} \in \mathrm{X}_{2}$. Hence $x_{2} \mapsto \mathrm{~d} f_{x_{2}} \in$ $L^{2}\left(\mathrm{X}_{2}, L^{2}\left(T^{*} \mathrm{X}_{1}\right)\right)$.

By linearity, the same holds for a generic $f \in \mathcal{A}$. Now notice that for an arbitrary $f \in W^{1,2}\left(\mathrm{X}_{1} \times \mathrm{X}_{2}\right)$, the identity (3.2.7) yields

$$
\begin{equation*}
\left|\mathrm{d} f_{x_{2}}\right|^{2}\left(x_{1}\right) \leq|\mathrm{d} f|^{2}\left(x_{1}, x_{2}\right) \quad \mathfrak{m}_{1} \times \mathfrak{m}_{2}-\text { a.e. }\left(x_{1}, x_{2}\right) \tag{3.2.11}
\end{equation*}
$$

and thus

$$
\begin{equation*}
\|\mid \mathrm{d} f .-\mathrm{d} \tilde{f} \cdot\|_{L^{2}\left(\mathrm{X}_{2}, L^{2}\left(\mathrm{X}_{1}\right)\right)} \leq\|f-\tilde{f}\|_{W^{1,2}\left(\mathrm{X}_{1} \times \mathrm{X}_{2}\right)} . \tag{3.2.12}
\end{equation*}
$$

Hence for $f \in W^{1,2}\left(\mathrm{X}_{1} \times \mathrm{X}_{2}\right)$ arbitrary, using the density of the product algebra we can find $\left(f_{n}\right) \subset \mathcal{A} W^{1,2}$-converging to $f$, so that from (3.2.12) we see that $\mathrm{d} f . \in L^{2}\left(\mathrm{X}_{2}, L^{2}\left(T^{*} \mathrm{X}_{1}\right)\right)$.

For general $f \in \mathcal{S}_{\text {loc }}^{2}\left(\mathrm{X}_{1} \times \mathrm{X}_{2}\right)$, find a sequence $\left(f_{n}\right) \subset W^{1,2}\left(\mathrm{X}_{1} \times \mathrm{X}_{2}\right)$ as in (2.2.2) and use the locality of the differential to get that (3.2.11) holds even for $f \in \mathcal{S}_{\text {loc }}^{2}\left(\mathrm{X}_{1} \times \mathrm{X}_{2}\right)$. Thus, since clearly $\mathrm{d} f_{n} \rightarrow \mathrm{~d} f$ in $L^{0}\left(T^{*}\left(\mathrm{X}_{1} \times \mathrm{X}_{2}\right)\right)$, from (3.2.11) we also get that $\left|\mathrm{d}\left(f_{n}\right) .-\mathrm{d} f.\right| \rightarrow 0$ in $L^{0}\left(\mathrm{X}_{2}, L^{0}\left(\mathrm{X}_{1}\right)\right)$ : this proves both that $\mathrm{d} f . \in L^{0}\left(\mathrm{X}_{2}, L^{0}\left(T^{*} \mathrm{X}_{1}\right)\right)$ and that $\mathrm{d}\left(f_{n}\right) . \rightarrow \mathrm{d} f$. in $L^{0}\left(\mathrm{X}_{2}, L^{0}\left(T^{*} \mathrm{X}_{1}\right)\right)$.

Since this latter convergence does not depend on the particular choice of the sequence $\left(f_{n}\right) \subset \mathcal{S}_{\text {loc }}^{2}\left(\mathrm{X}_{1} \times \mathrm{X}_{2}\right)$ such that $\mathrm{d} f_{n} \rightarrow \mathrm{~d} f$ in $L^{0}\left(T^{*}\left(\mathrm{X}_{1} \times \mathrm{X}_{2}\right)\right)$, we proved also 3.2.10).

The last claim follows along the same approximation argument using the continuity property (3.2.10) (and the analogous one with $\mathrm{X}_{1}$ and $\mathrm{X}_{2}$ inverted).

Let $\mathscr{M}_{1}, \mathscr{M}_{2}$ be two $L^{0}$-normed modules on a space X . Then on the product $\mathscr{M}_{1} \times \mathscr{M}_{2}$ we shall consider the structure of $L^{0}$-normed module given by: the product topology, the multiplication by $L^{0}$-functions given by $f\left(v_{1}, v_{2}\right):=\left(f v_{1}, f v_{2}\right)$ and the pointwise norm defined as

$$
\left|\left(v_{1}, v_{2}\right)\right|^{2}:=\left|v_{1}\right|^{2}+\left|v_{2}\right|^{2} .
$$

It is readily verified that these actually endow $\mathscr{M}_{1} \times \mathscr{M}_{2}$ with the structure of $L^{0}$-normed module.

In particular, $L^{0}\left(\mathrm{X}_{2}, L^{0}\left(T^{*} \mathrm{X}_{1}\right)\right) \times L^{0}\left(\mathrm{X}_{1}, L^{0}\left(T^{*} \mathrm{X}_{2}\right)\right)$ is a $L^{0}\left(\mathrm{X}_{1} \times \mathrm{X}_{2}\right)$-normed module and we can define $\Phi_{1} \oplus \Phi_{2}$ as

$$
\begin{aligned}
\Phi_{1} \oplus \Phi_{2}: \quad L^{0}\left(\mathrm{X}_{2}, L^{0}\left(T^{*} \mathrm{X}_{1}\right)\right) \times L^{0}\left(\mathrm{X}_{1}, L^{0}\left(T^{*} \mathrm{X}_{2}\right)\right) & \rightarrow L^{0}\left(T^{*}\left(\mathrm{X}_{1} \times \mathrm{X}_{2}\right)\right) \\
(\omega, \sigma) & \mapsto \Phi_{1}(\omega)+\Phi_{2}(\sigma)
\end{aligned}
$$

We then have the following result:
Theorem 3.13. Let $\left(\mathrm{X}_{1}, \mathrm{~d}_{1}, \mathfrak{m}_{1}\right)$ and $\left(\mathrm{X}_{2}, \mathrm{~d}_{2}, \mathfrak{m}_{2}\right)$ be two metric measure spaces such that Assumption 3.10 holds. Then $\Phi_{1} \oplus \Phi_{2}$ is an isomorphism of modules, i.e. it is $L^{0}\left(\mathrm{X}_{1} \times \mathrm{X}_{2}\right)$ linear, continuous, surjective and for every $\omega . \in L^{0}\left(\mathrm{X}_{2}, L^{0}\left(T^{*} \mathrm{X}_{1}\right)\right)$ and $\sigma . \in L^{0}\left(\mathrm{X}_{1}, L^{0}\left(T^{*} \mathrm{X}_{2}\right)\right)$ satisfies

$$
\begin{equation*}
\left|\Phi_{1}(\omega .)+\Phi_{2}(\sigma .)\right|^{2}=\left|\omega .\left.\right|^{2}+|\sigma .|^{2} \quad \mathfrak{m}_{1} \times \mathfrak{m}_{2}-\right.\text { a.e.. } \tag{3.2.13}
\end{equation*}
$$

Moreover, for every $f \in \mathcal{S}_{\text {loc }}^{2}\left(\mathrm{X}_{1} \times \mathrm{X}_{2}\right)$ it holds:

$$
\begin{equation*}
\mathrm{d} f=\Phi_{1}(\mathrm{~d} f .)+\Phi_{2}\left(\mathrm{~d} f^{\prime}\right) . \tag{3.2.14}
\end{equation*}
$$

proof From Proposition 3.7 it is clear that $\Phi_{1} \oplus \Phi_{2}$ is $L^{0}\left(\mathrm{X}_{1} \times \mathrm{X}_{2}\right)$-linear and continuous. Taking into account that $L^{0}\left(\mathrm{X}_{2}, L^{0}\left(T^{*} \mathrm{X}_{1}\right)\right)$ is generated by elements of the kind $\widehat{\mathrm{d} g}$ for
$g \in \mathcal{S}^{2}\left(\mathrm{X}_{1}\right)$, where $\widehat{\mathrm{d} g} \in L^{0}\left(\mathrm{X}_{2}, L^{0}\left(T^{*} \mathrm{X}_{1}\right)\right)$ is the function identically equal to $\mathrm{d} g$, and similarly for $L^{0}\left(\mathrm{X}_{1}, L^{0}\left(T^{*} \mathrm{X}_{2}\right)\right)$, to prove (3.2.13) it is sufficient to show that

$$
\begin{equation*}
\left|\Phi_{1}(\widehat{\mathrm{~d} g})+\Phi_{2}(\widehat{\mathrm{~d} h})\right|^{2}=|\mathrm{d} g|^{2} \circ \pi_{1}+|\mathrm{d} h|^{2} \circ \pi_{2} \quad \mathfrak{m}_{1} \times \mathfrak{m}_{2}-\text { a.e. } \tag{3.2.15}
\end{equation*}
$$

for any $g \in \mathcal{S}^{2}\left(\mathrm{X}_{1}\right), h \in \mathcal{S}^{2}\left(\mathrm{X}_{2}\right)$. Fix such $g, h$ and put $f:=g \circ \pi_{1}+h \circ \pi_{2} \in \mathcal{S}_{\text {loc }}^{2}\left(\mathrm{X}_{1} \times \mathrm{X}_{2}\right)$. Notice that trivially $\mathrm{d} f_{x_{2}}=\mathrm{d} g$ and $\mathrm{d} f^{x_{1}}=\mathrm{d} h$ for any $x_{1} \in \mathrm{X}_{1}$ and $x_{2} \in \mathrm{X}_{2}$, hence from the tensorization of Cheeger energy (recall the last claim of Lemma 3.12 above) we have

$$
\left|\Phi_{1}(\widehat{\mathrm{~d} g})+\Phi_{2}(\widehat{\mathrm{~d} h})\right|^{2}=\left|\mathrm{d}\left(g \circ \pi_{1}\right)+\mathrm{d}\left(h \circ \pi_{2}\right)\right|^{2}=|\mathrm{d} f|^{2} \stackrel{\sqrt{3.2 .7}}{=}\left|\mathrm{d} f .\left.\right|^{2}+|\mathrm{d} f|^{2}=|\mathrm{d} g|^{2} \circ \pi_{1}+|\mathrm{d} h|^{2} \circ \pi_{2}\right.
$$

which is 3.2 .15 ). Thus $\Phi_{1} \oplus \Phi_{2}$ preserves the pointwise norm.
Now we prove (3.2.14). Let $g \in L^{\infty} \cap W^{1,2}\left(\mathrm{X}_{1}\right)$ and $h \in L^{\infty} \cap W^{1,2}\left(\mathrm{X}_{2}\right)$ be both with bounded support and consider $f:=g \circ \pi_{1} h \circ \pi_{2}$. Then $f \in W^{1,2}\left(\mathrm{X}_{1} \times \mathrm{X}_{2}\right)$ and the very definition of $\Phi_{1}, \Phi_{2}$ grant that

$$
\mathrm{d} f=h \circ \pi_{2} \mathrm{~d}\left(g \circ \pi_{1}\right)+g \circ \pi_{1} \mathrm{~d}\left(h \circ \pi_{2}\right)=\Phi_{1}(h \mathrm{~d} g)+\Phi_{2}(g \mathrm{~d} h)=\Phi_{1}(\mathrm{~d} f .)+\Phi_{2}(\mathrm{~d} f \cdot),
$$

so that in this case (3.2.14) is proved. By linearity, we get that (3.2.14) holds for general $f \in \mathcal{A}$. Then using first the density of $\mathcal{A}$ in $W^{1,2}\left(\mathrm{X}_{1} \times \mathrm{X}_{2}\right)$ and then property 2.2 .2 , taking into account the convergence property $(3.2 .10)$ we conclude that (3.2.14) holds for general $f \in \mathcal{S}_{\text {loc }}^{2}\left(\mathrm{X}_{1} \times \mathrm{X}_{2}\right)$, as claimed.

It remains to prove that $\Phi_{1} \oplus \Phi_{2}$ is surjective. By (3.2.14) we know that its image contains the space of differential of functions in $\mathcal{S}_{\text {loc }}^{2}\left(\mathrm{X}_{1} \times \mathrm{X}_{2}\right)$, and thus $L^{0}$-linear combinations of them. Since it preserves the pointwise norm, its image must be closed and since $L^{0}\left(T^{*}\left(\mathrm{X}_{1} \times \mathrm{X}_{2}\right)\right)$ is generated by differentials of functions in $\mathcal{S}_{\text {loc }}^{2}\left(\mathrm{X}_{1} \times \mathrm{X}_{2}\right)$, this is sufficient to conclude.

### 3.2.2 Other differential operators in product spaces

In the previous section we have seen how the differential behaves under products of spaces. We shall now investigate other differentiation operators under the assumption that $\mathrm{X}_{1}, \mathrm{X}_{2}$ are infinitesimally Hilbertian.

We start with the following simple orthogonality statement:
Proposition 3.14. Let $\left(\mathrm{X}_{1}, \mathrm{~d}_{1}, \mathfrak{m}_{1}\right)$ and $\left(\mathrm{X}_{2}, \mathrm{~d}_{2}, \mathfrak{m}_{2}\right)$ be infinitesimally Hilbertian spaces such that Assumption 3.10 holds. Then $\mathrm{X}_{1} \times \mathrm{X}_{2}$ is also infinitesimally Hilbertian and for every $\omega_{1}^{1} \in L^{0}\left(\mathrm{X}_{2}, L^{0}\left(T^{*} \mathrm{X}_{1}\right)\right)$ and $\omega^{2} \in L^{0}\left(\mathrm{X}_{1}, L^{0}\left(T^{*} \mathrm{X}_{2}\right)\right)$ we have

$$
\begin{equation*}
\left\langle\Phi_{1}\left(\omega_{.}^{1}\right), \Phi_{2}\left(\omega_{.}^{2}\right)\right\rangle=0 \quad \mathfrak{m}_{1} \times \mathfrak{m}_{2}-\text { a.e.. } \tag{3.2.16}
\end{equation*}
$$

proof The fact that $W^{1,2}\left(\mathrm{X}_{1} \times \mathrm{X}_{2}\right)$ is Hilbert is a direct consequence of the tensorization of the Cheeger energy and the assumption that both $W^{1,2}\left(\mathrm{X}_{1}\right)$ and $W^{1,2}\left(\mathrm{X}_{2}\right)$ are Hilbert. For (3.2.16) notice that

$$
\left|\omega_{.}^{1}\right|^{2}+\left|\omega_{.}^{2}\right|^{2} \stackrel{\sqrt{3.2 .13}}{-}\left|\Phi_{1}\left(\omega_{.}^{1}\right)+\Phi_{2}\left(\omega_{.}^{2}\right)\right|^{2}=\left|\Phi_{1}\left(\omega_{.}^{1}\right)\right|^{2}+\left|\Phi_{2}\left(\omega_{.}^{2}\right)\right|^{2}+2\left\langle\Phi_{1}\left(\omega_{.}^{1}\right), \Phi_{2}\left(\omega_{.}^{2}\right)\right\rangle,
$$

so that the conclusion follows recalling that $\Phi_{1}, \Phi_{2}$ preserve the pointwise norms.
By means of the musical isomorphisms (recall (2.2.3)) the map $\Phi_{1}$ induces a map, still denoted $\Phi_{1}$, from $L^{0}\left(\mathrm{X}_{2}, L^{0}\left(T \mathrm{X}_{1}\right)\right)$ to $L^{0}\left(T\left(\mathrm{X}_{1} \times \mathrm{X}_{2}\right)\right)$ via:

$$
\Phi_{1}(X .):=\Phi_{1}\left(X^{b}\right)^{\sharp} .
$$

Similarly for $\Phi_{2}$. It is clear that these newly defined $\Phi_{1}, \Phi_{2}$ have all the properties we previously proved for the same operators viewed as acting on forms. We also notice that for any $\omega . \in L^{0}\left(\mathrm{X}_{2}, L^{0}\left(T^{*} \mathrm{X}_{1}\right)\right)$ and $X . \in L^{0}\left(\mathrm{X}_{2}, L^{0}\left(T \mathrm{X}_{1}\right)\right)$ we have

$$
\begin{equation*}
\Phi_{1}(\omega .)\left(\Phi_{1}(X .)\right)\left(x_{1}, x_{2}\right)=\omega_{x_{2}}\left(X_{x_{2}}\right)\left(x_{1}\right) \quad \mathfrak{m}_{1} \times \mathfrak{m}_{2}-\text { a.e. }\left(x_{1}, x_{2}\right) . \tag{3.2.17}
\end{equation*}
$$

Indeed, for $\omega . \equiv \mathrm{d} g$ and $X . \equiv \nabla \tilde{g}$ for $g, \tilde{g} \in \mathcal{S}_{\text {loc }}^{2}\left(\mathrm{X}_{1}\right)$ this is a direct consequence of the definition of $\Phi_{1}$ and the fact that $\Phi_{1}$ preserves the pointwise norm (and hence the pointwise scalar product), then the general case follows by $L^{0}\left(\mathrm{X}_{1} \times \mathrm{X}_{2}\right)$-bilinearity and continuity of both sides.

Proposition 3.15. Let $\left(\mathrm{X}_{1}, \mathrm{~d}_{1}, \mathfrak{m}_{1}\right)$ and $\left(\mathrm{X}_{2}, \mathrm{~d}_{2}, \mathfrak{m}_{2}\right)$ be infinitesimally Hilbertian spaces such that Assumption 3.10 holds. Then $X \in D\left(\operatorname{div}_{\text {loc }}, \mathrm{X}_{1}\right)$ if and only if $\Phi_{1}(\hat{X}) \in D\left(\operatorname{div}_{\text {loc }}, \mathrm{X}_{1} \times \mathrm{X}_{2}\right)$, where $\hat{X} \in L^{0}\left(\mathrm{X}_{2}, L^{0}\left(T \mathrm{X}_{1}\right)\right)$ is the function identically equal to $X$, and in this case

$$
\operatorname{div}\left(\Phi_{1}(\hat{X})\right)=\operatorname{div}(X) \circ \pi_{1}
$$

proof From the very definition of divergence it is readily verified that the thesis is equivalent to

$$
\int \mathrm{d} f\left(\Phi_{1}(\hat{X})\right) \mathrm{d}\left(\mathfrak{m}_{1} \times \mathfrak{m}_{2}\right)=\iint \mathrm{d} f .(X) \mathrm{d} \mathfrak{m}_{1} \mathrm{~d} \mathfrak{m}_{2}
$$

for every $f \in W^{1,2}\left(\mathrm{X}_{1} \times \mathrm{X}_{2}\right)$ with bounded support.
For such $f$ we have

$$
\begin{aligned}
\int \mathrm{d} f\left(\Phi_{1}(\hat{X})\right) \mathrm{d}\left(\mathfrak{m}_{1} \times \mathfrak{m}_{2}\right) & \stackrel{\sqrt{3.2 .14}}{=} \int\left(\Phi_{1}(\mathrm{~d} f .)+\Phi_{2}(\mathrm{~d} f \cdot)\right)\left(\Phi_{1}(\hat{X})\right) \mathrm{d}\left(\mathfrak{m}_{1} \times \mathfrak{m}_{2}\right) \\
& \stackrel{(3.2 .16}{=} \int\left(\Phi_{1}(\mathrm{~d} f .)\right)\left(\Phi_{1}(\hat{X})\right) \mathrm{d}\left(\mathfrak{m}_{1} \times \mathfrak{m}_{2}\right) \\
& \stackrel{3.2 .17}{=} \iint \mathrm{d} f .(X) \mathrm{d} \mathfrak{m}_{1} \mathrm{~d} \mathfrak{m}_{2},
\end{aligned}
$$

hence the conclusion.

A related property is the following:
Proposition 3.16. Let $\left(\mathrm{X}_{1}, \mathrm{~d}_{1}, \mathfrak{m}_{1}\right)$ and $\left(\mathrm{X}_{2}, \mathrm{~d}_{2}, \mathfrak{m}_{2}\right)$ be infinitesimally Hilbertian spaces such that Assumption 3.10 holds. Let $X=\Phi_{1}\left(X_{.}^{1}\right)+\Phi_{2}\left(X_{.}^{2}\right) \in L^{2}\left(T\left(\mathrm{X}_{1} \times \mathrm{X}_{2}\right)\right)$ be such that:

$$
\begin{aligned}
& \text { - } X_{x_{2}}^{1} \in D\left(\operatorname{div}, \mathrm{X}_{1}\right) \text { for } \mathfrak{m}_{2} \text {-a.e. } x_{2} \in \mathrm{X}_{2} \text { with } \int\left|\operatorname{div}\left(X_{.}^{1}\right)\right|^{2} \mathrm{~d}\left(\mathfrak{m}_{1} \times \mathfrak{m}_{2}\right)<\infty \\
& \text { - } X_{x_{1}}^{2} \in D\left(\operatorname{div}, \mathrm{X}_{2}\right) \text { for } \mathfrak{m}_{1} \text {-a.e. } x_{1} \in \mathrm{X}_{1} \text { with } \int\left|\operatorname{div}\left(X_{.}^{2}\right)\right|^{2} \mathrm{~d}\left(\mathfrak{m}_{1} \times \mathfrak{m}_{2}\right)<\infty
\end{aligned}
$$

Then $X \in D($ div $)$ and

$$
\begin{equation*}
\operatorname{div}(X)\left(x_{1}, x_{2}\right)=\operatorname{div}\left(X_{x_{2}}^{1}\right)\left(x_{1}\right)+\operatorname{div}\left(X_{x_{1}}^{2}\right)\left(x_{2}\right) \quad \mathfrak{m}_{1} \times \mathfrak{m}_{2}-\text { a.e. }\left(x_{1}, x_{2}\right) \tag{3.2.18}
\end{equation*}
$$

proof For any $f \in W^{1,2}\left(\mathrm{X}_{1} \times \mathrm{X}_{2}\right)$ with bounded support we have

$$
\begin{aligned}
\int \mathrm{d} f(X) \mathrm{d}\left(\mathfrak{m}_{1} \times \mathfrak{m}_{2}\right) & \stackrel{\sqrt[3.2 .14]{=}}{ } \int\left(\Phi_{1}(\mathrm{~d} f .)+\Phi_{2}\left(\mathrm{~d} f^{\prime}\right)\right)\left(\Phi_{1}\left(X_{.}^{1}\right)+\Phi_{2}\left(X_{.}^{2}\right)\right) \mathrm{d}\left(\mathfrak{m}_{1} \times \mathfrak{m}_{2}\right) \\
& \stackrel{(3.2 .16}{=} \int \Phi_{1}(\mathrm{~d} f .) \Phi_{1}\left(X_{.}^{1}\right)+\Phi_{2}\left(\mathrm{~d} f^{\prime}\right) \Phi_{2}\left(X_{.}^{2}\right) \mathrm{d}\left(\mathfrak{m}_{1} \times \mathfrak{m}_{2}\right) \\
& \stackrel{3.2 .17}{=} \int\left(\int \mathrm{d} f .\left(X_{.}^{1}\right) \mathrm{d} \mathfrak{m}_{1}\right) \mathrm{d} \mathfrak{m}_{2}+\int\left(\int \mathrm{d} f \cdot\left(X_{.}^{2}\right) \mathrm{d} \mathfrak{m}_{2}\right) \mathrm{d} \mathfrak{m}_{2}
\end{aligned}
$$

which is the thesis.
These last two statements produce analogous ones for the Laplacian:
Corollary 3.17. Let $\left(\mathrm{X}_{1}, \mathrm{~d}_{1}, \mathfrak{m}_{1}\right)$ and $\left(\mathrm{X}_{2}, \mathrm{~d}_{2}, \mathfrak{m}_{2}\right)$ be infinitesimally Hilbertian spaces such that Assumption 3.10 holds. Then:
i) $f \in D\left(\Delta_{\text {loc }}, \mathrm{X}_{1}\right)$ if and only if $f \circ \pi_{1} \in D\left(\Delta_{\text {loc }}, \mathrm{X}_{1} \times \mathrm{X}_{2}\right)$ and in this case

$$
\Delta\left(f \circ \pi_{1}\right)=(\Delta f) \circ \pi_{1} .
$$

ii) Let $f \in W^{1,2}\left(\mathrm{X}_{1} \times \mathrm{X}_{2}\right)$ be such that

- for $\mathfrak{m}_{1}$-a.e. $x_{1} \in \mathrm{X}_{1}, f^{x_{1}} \in D\left(\Delta, \mathrm{X}_{2}\right)$ with $\int\left\|\Delta f^{x_{1}}\right\|_{L^{2}\left(\mathrm{X}_{2}\right)}^{2} \mathrm{~d} \mathfrak{m}_{1}<\infty$,
- for $\mathfrak{m}_{2}$-a.e. $x_{2} \in \mathrm{X}_{2}, f_{x_{2}} \in D\left(\Delta, \mathrm{X}_{1}\right)$ with $\int\left\|\Delta f_{x_{2}}\right\|_{L^{2}\left(\mathrm{X}_{1}\right)}^{2} \mathrm{~d} \mathfrak{m}_{2}<\infty$.

Then $f \in D\left(\Delta, \mathrm{X}_{1} \times \mathrm{X}_{2}\right)$ and

$$
\begin{equation*}
\Delta f\left(x_{1}, x_{2}\right)=\Delta f_{x_{2}}\left(x_{1}\right)+\Delta f^{x_{1}}\left(x_{2}\right) \quad \mathfrak{m}_{1} \times \mathfrak{m}_{2}-\text { a.e. }\left(x_{1}, x_{2}\right) . \tag{3.2.19}
\end{equation*}
$$

proof For the first claim simply notice that, directly from the definition, we have $f \in$ $D\left(\Delta_{\text {loc }}, \mathrm{X}_{1}\right)$ if and only if $\nabla f \in D\left(\operatorname{div}_{\text {loc }}, \mathrm{X}\right)$ and in this case $\operatorname{div}(\nabla f)=\Delta f$. Similarly for $f \circ \pi_{1}$. Then observe that (3.2.1) grants that $\nabla\left(f \circ \pi_{1}\right)=\Phi_{1}(\widehat{\nabla f})$ and apply Proposition 3.15 above to conclude.

The second claim follows by analogous considerations using Proposition 3.16 and the identity $\nabla f=\Phi_{1}(\nabla f)+.\Phi_{2}\left(\nabla f^{\prime}\right)($ recall 3.2 .14$)$.

### 3.3 Flow of harmonic vector fields on $\operatorname{RCD}(0, \infty)$ spaces

In this section we work on a fixed $\operatorname{RCD}(0, \infty)$ space $(X, d, \mathfrak{m})$ and study the Regular Lagrangian Flow of a fixed non-zero vector field $X \in L^{2}(T X)$ which is harmonic, i.e. $X^{\mathrm{b}} \in D\left(\Delta_{H}\right)$ with $\Delta_{H} X^{\mathrm{b}}=0$. Recalling (2.2.10) we have that $\operatorname{div} X=0$, while (2.2.9) grants that $X$ is parallel, i.e. $X \in H_{C}^{1,2}(T \mathrm{X})$ with $\nabla X=0$. This latter property also implies that $|X|$ is constant (see [15] for the details about this last claim).

We can thus apply Theorem 2.8 to deduce that there exists and is unique the Regular Lagrangian Flow $\left(\mathrm{Fl}_{t}^{(X)}\right)$ of $X$. Aim of this section is to prove that:
i) the $\mathrm{Fl}_{t}^{X}$,s are measure preserving isometries
ii) if $Y$ is another harmonic vector field, then $\mathrm{Fl}_{t}^{X} \circ \mathrm{Fl}_{s}^{Y}=\mathrm{Fl}_{1}^{t X+s Y}$ for any $t, s$.

Notice that by analogy with the smooth case, one would expect to need only the conditions $\operatorname{div} X=0, \nabla X=0$ and that X is a $\operatorname{RCD}(K, \infty)$ space to get the above. Yet, it is unclear to us whether these are really sufficient, (part of) the problem being in the approximation procedure used in Proposition 3.20 which requires our stronger assumptions.

In what comes next we shall occasionally use the following simple fact: for $T, S: \mathrm{X} \rightarrow \mathrm{X}$ Borel we have
$T_{*} \mu=S_{*} \mu \quad \forall \mu \in \mathscr{P}(\mathrm{X})$ with bounded support and density $\quad \Rightarrow \quad T=S \quad \mathfrak{m}$-a.e..
Indeed, if $T \neq S$ on a set of positive measure, for some $r>0$ we would have $\mathrm{d}(T(x), S(x))>$ $2 r$ for a set of $x$ 's of positive measure and thus using the separability of X we would be able to find $\bar{x}$ such that $T_{*} \mathfrak{m}\left(B_{r}(\bar{x})\right)>0$. Thus $\mathfrak{m}\left(T^{-1}\left(B_{r}(\bar{x})\right)\right)>0$ and letting $A \subset T^{-1}\left(B_{r}(\bar{x})\right)$ be any bounded Borel subset of positive $\mathfrak{m}$-measure, for $\mu:=\left.\mathfrak{m}(A)^{-1} \mathfrak{m}\right|_{A}$ we would have that $T_{*} \mu$ and $S_{*} \mu$ are concentrated on disjoint sets, and thus in particular $T_{*} \mu \neq S_{*} \mu$.

With this said, we prove the following result, which shows that the flows of $X$ and $-X$ are one the inverse of the other:

Lemma 3.18. Let $(\mathrm{X}, \mathrm{d}, \mathfrak{m})$ be $a \operatorname{RCD}(0, \infty)$ space and $X$ a harmonic vector field. Then for every $t \geq 0$ the following identities hold $\mathfrak{m}$-a.e.:

$$
\mathrm{Fl}_{t}^{(-X)} \circ \mathrm{Fl}_{t}^{(X)}=\mathrm{Id} \quad \text { and } \quad \mathrm{Fl}_{t}^{(X)} \circ \mathrm{Fl}_{t}^{(-X)}=\mathrm{Id}
$$

proof We shall prove the first identity for $t=1$, as then the rest follows by similar arguments. Let $\mu \in \mathscr{P}(\mathrm{X})$ be with bounded support and density, and consider the curves $[0,1] \ni t \mapsto \mu_{t}, \tilde{\mu}_{t} \in \mathscr{P}(\mathrm{X})$ defined as

$$
\mu_{t}:=\left(\mathrm{Fl}_{1-t}^{(X)}\right)_{*} \mu \quad \text { and } \quad \tilde{\mu}_{t}:=\left(\mathrm{Fl}_{t}^{(-X)}\right)_{*}\left(\mathrm{Fl}_{1}^{(X)}\right)_{*} \mu,
$$

notice that $\mu_{0}=\tilde{\mu}_{0}$ and that they both solve the continuity equation (2.3.11) for $X_{t}=-X$ in the sense of Definition 2.9. By Theorem 2.10 we conclude that $\mu_{1}=\tilde{\mu}_{1}$, i.e.

$$
\mu=\left(\mathrm{Fl}_{1}^{(-X)} \circ \mathrm{Fl}_{1}^{(X)}\right)_{*} \mu
$$

The conclusion follows by the arbitrariness of $\mu$ and (3.3.1).
From this proposition and the semigroup property 2.3.10 of Regular Lagrangian Flows, it follows that defining $\mathrm{Fl}_{-t}^{(X)}:=\mathrm{Fl}_{t}^{(-X)}$ for $t \geq 0$ we have

$$
\begin{equation*}
\mathrm{Fl}_{t}^{(X)} \circ \mathrm{Fl}_{s}^{(X)}=\mathrm{Fl}_{t+s}^{(X)} \quad \text { m-a.e. } \quad \forall t, s \in \mathbb{R} \text {. } \tag{3.3.2}
\end{equation*}
$$

Proposition 3.19 (Preservation of the measure). Let (X, d, m) be a $\operatorname{RCD}(0, \infty)$ space and $X$ a harmonic vector field. Then for every $t \in \mathbb{R}$ we have

$$
\begin{equation*}
\left(\mathrm{Fl}_{t}^{(X)}\right)_{*} \mathfrak{m}=\mathfrak{m} \tag{3.3.3}
\end{equation*}
$$

proof Simply notice that from $\operatorname{div}(X)=\operatorname{div}(-X)=0$ and 2.3.9, for any $t \geq 0$ we have

$$
\mathfrak{m}=\left(\mathrm{Fl}_{t}^{(X)} \circ \mathrm{Fl}_{-t}^{(X)}\right)_{*} \mathfrak{m}=\left(\mathrm{Fl}_{t}^{(X)}\right)_{*}\left(\mathrm{Fl}_{t}^{(-X)}\right)_{*} \mathfrak{m} \leq\left(\mathrm{Fl}_{t}^{(X)}\right)_{*} \mathfrak{m} \leq \mathfrak{m}
$$

forcing the inequalities to be equalities.
Recall that the heat flow $\left(\mathrm{h}_{t}\right)$ on X is firstly defined as the $L^{2}$-gradient flow of the Dirichlet and then extended to a flow on $L^{1}+L^{\infty}$ by monotonicity and continuity. On the other hand, the 'Hodge' heat flow $\left(\mathrm{h}_{H, t}\right)$ is defined on $L^{2}\left(T^{*} \mathrm{X}\right)$ as the gradient flow of the functional

$$
L^{2}\left(T^{*} \mathrm{X}\right) \ni \omega \quad \mapsto \quad \begin{cases}\frac{1}{2} \int|\mathrm{~d} \omega|^{2}+|\delta \omega|^{2} \mathrm{~d} \mathfrak{m} & \text { if } \omega \in H_{H}^{1,2}\left(T^{*} \mathrm{X}\right) \\ +\infty & \text { otherwise }\end{cases}
$$

For us it will be relevant to know the relation

$$
\begin{equation*}
\mathrm{h}_{H, t} \mathrm{~d} f=\mathrm{dh}_{t} f \quad \forall f \in W^{1,2}(\mathrm{X}) \tag{3.3.4}
\end{equation*}
$$

and the improved Bakry-Émery estimate:

$$
\begin{equation*}
\left|\mathrm{h}_{H, t} \omega\right|^{2} \leq e^{-K t} \mathrm{~h}_{t}\left(|\omega|^{2}\right), \quad \mathfrak{m}-a . e . \quad \forall t \geq 0, \omega \in L^{2}\left(T^{*} \mathrm{X}\right) \tag{3.3.5}
\end{equation*}
$$

valid on $\operatorname{RCD}(K, \infty)$ spaces. See [15] for further details about these.
With this said, we can now prove the following lemma, which is key to show that $\mathrm{Fl}_{t}^{(X)}$ is an isometry.

Proposition 3.20 (Euler's equation for $X$ ). Let (X, d, $\mathfrak{m})$ be $a \operatorname{RCD}(0, \infty)$ space and $X$ a harmonic vector field. Then for any $f \in W^{1,2}(\mathrm{X})$ it holds

$$
\begin{equation*}
\mathrm{h}_{t}(\langle\nabla f, X\rangle)=\left\langle\nabla \mathrm{h}_{t} f, X\right\rangle, \quad \mathfrak{m} \text {-a.e. }, \forall t \geq 0 . \tag{3.3.6}
\end{equation*}
$$

Moreover, for every $f \in D(\Delta)$ with $\Delta f \in W^{1,2}(\mathrm{X})$, we have $\langle\nabla f, X\rangle \in D(\Delta)$ and

$$
\begin{equation*}
\Delta\langle\nabla f, X\rangle=\langle\nabla \Delta f, X\rangle, \quad \mathfrak{m}-a . e . . \tag{3.3.7}
\end{equation*}
$$

proof We apply (3.3.5) in our space to the form $X^{b}+\varepsilon \mathrm{d} f$ to obtain

$$
\begin{equation*}
\left|\mathrm{h}_{H, t}\left(X^{b}+\varepsilon \mathrm{d} f\right)\right|^{2} \leq \mathrm{h}_{t}\left(|X+\varepsilon \nabla f|^{2}\right) . \tag{3.3.8}
\end{equation*}
$$

We have already observed that the fact that $X^{b}$ is harmonic grants that $|X|$ is constant, say $|X| \equiv c$. The harmonicity also grants that $\mathrm{h}_{H, t}\left(X^{\mathrm{b}}\right)=X^{\mathrm{b}}$ for every $t \geq 0$, hence we have $\left|\mathrm{h}_{H, t}\left(X^{b}\right)\right|^{2} \equiv c^{2} \equiv \mathrm{~h}_{t}\left(|X|^{2}\right)$ for any $t \geq 0$. Therefore,

$$
c^{2}+2 \varepsilon\left\langle X, \mathrm{~h}_{H, t}(\mathrm{~d} f)\right\rangle+\varepsilon^{2}\left|\mathrm{~h}_{H, t}(\mathrm{~d} f)\right|^{2} \leq c^{2}+2 \varepsilon \mathrm{~h}_{t}\langle X, \nabla f\rangle+\varepsilon^{2} \mathrm{~h}_{t}\left(|\mathrm{~d} f|^{2}\right)
$$

and the arbitrariness of $\varepsilon \in \mathbb{R}$ implies

$$
\left\langle X, \mathrm{~h}_{H, t} \mathrm{~d} f\right\rangle=2 \mathrm{~h}_{t}\langle X, \nabla f\rangle,
$$

which by (3.3.4) is (3.3.6). Then (3.3.7) comes by differentiating (3.3.6) at $t=0$.
Proposition 3.21 (Preservation of the Dirichlet energy). Let (X,d,m) be a $\operatorname{RCD}(0, \infty)$ space and $X$ a harmonic vector field. Then for every $t \in \mathbb{R}$ we have

$$
\begin{equation*}
\mathrm{E}\left(f \circ \mathrm{Fl}_{t}^{X}\right)=\mathrm{E}(f) \quad \forall f \in W^{1,2}(\mathrm{X}) \tag{3.3.9}
\end{equation*}
$$

proof Fix $f \in W^{1,2}(\mathrm{X})$, put $f_{t}:=f \circ \mathrm{Fl}_{t}^{X}$ and notice that since $\mathrm{E}\left(\mathrm{h}_{\varepsilon} g\right) \rightarrow \mathrm{E}(g)$ as $\varepsilon \downarrow 0$ for any $g \in L^{2}(\mathrm{X})$, it is sufficient to prove that for any $\varepsilon>0$ we have

$$
\mathrm{E}\left(\mathrm{~h}_{\varepsilon} f_{t}\right)=\mathrm{E}\left(\mathrm{~h}_{\varepsilon} f\right) \quad \forall t \in \mathbb{R} .
$$

Thus fix $\varepsilon>0$ and notice that Proposition 2.7 grants that $t \mapsto f_{t} \in L^{2}(\mathrm{X})$ is Lipschitz. This in conjunction with the fact that $\mathrm{h}_{\varepsilon}: L^{2}(\mathrm{X}) \rightarrow W^{1,2}(\mathrm{X})$ is continuous ensures that $t \mapsto \mathrm{~h}_{\varepsilon} f_{t} \in W^{1,2}(\mathrm{X})$ is Lipschitz.

We now compute the derivative of the Lipschitz map $t \mapsto \mathrm{E}\left(\mathrm{h}_{\varepsilon} f_{t}\right)$ and start noticing that

$$
\int\left|\nabla \mathrm{h}_{\varepsilon} f_{t+h}\right|^{2}-\left|\nabla \mathrm{h}_{\varepsilon} f_{t}\right|^{2} \mathrm{~d} \mathfrak{m}=\int\left|\nabla \mathrm{h}_{\varepsilon}\left(f_{t+h}-f_{t}\right)\right|^{2}+2\left\langle\nabla \mathrm{~h}_{\varepsilon} f_{t}, \nabla \mathrm{~h}_{\varepsilon}\left(f_{t+h}-f_{t}\right)\right\rangle \mathrm{d} \mathfrak{m}
$$

so that the Lipschitz regularity of $t \mapsto \mathrm{~h}_{\varepsilon} f_{t} \in W^{1,2}(\mathrm{X})$ grants that for any $t \in \mathbb{R}$ it holds

$$
\lim _{h \rightarrow 0} \int \frac{\left|\nabla \mathbf{h}_{\varepsilon} f_{t+h}\right|^{2}-\left|\nabla \mathbf{h}_{\varepsilon} f_{t}\right|^{2}}{2 h} \mathrm{~d} \mathfrak{m}=\lim _{h \rightarrow 0} \int\left\langle\nabla \mathrm{~h}_{\varepsilon} f_{t}, \nabla \frac{\mathrm{~h}_{\varepsilon} f_{t+h}-\mathrm{h}_{\varepsilon} f_{t}}{h}\right\rangle \mathrm{dm} .
$$

Hence

$$
\begin{aligned}
\frac{\mathrm{d}}{\mathrm{~d} t} \mathrm{E}\left(\mathrm{~h}_{\varepsilon} f_{t}\right) & =-\lim _{h \rightarrow 0} \int \Delta \mathrm{~h}_{\varepsilon} f_{t} \frac{\mathrm{~h}_{\varepsilon} f_{t+h}-\mathrm{h}_{\varepsilon} f_{t}}{h} \mathrm{dm} \\
& =-\lim _{h \rightarrow 0} \int \Delta \mathrm{~h}_{2 \varepsilon} f_{t} \frac{f_{t} \circ \mathrm{Fl}_{h}^{(X)}-f_{t}}{h} \mathrm{dm} \\
& =-\lim _{h \rightarrow 0} \int \frac{\left(\Delta \mathrm{~h}_{2 \varepsilon} f_{t}\right) \circ \mathrm{Fl}_{-h}^{(X)}-\Delta \mathrm{h}_{2 \varepsilon}\left(f_{t}\right)}{h} f_{t} \mathrm{dm} \\
7 & =-\int\left\langle\nabla \Delta \mathrm{h}_{2 \varepsilon} f_{t}, X\right\rangle f_{t} \mathrm{dm} .
\end{aligned}
$$

To conclude it is therefore sufficient to prove that for any $g \in L^{2}(\mathrm{X})$ it holds

$$
\begin{equation*}
\int\left\langle\nabla \Delta \mathrm{h}_{2 \varepsilon} g, X\right\rangle g \mathrm{~d} \mathfrak{m}=0 \tag{3.3.10}
\end{equation*}
$$

Hence fix $g \in L^{2}(\mathrm{X})$ and notice that

$$
\begin{equation*}
\int\left\langle\nabla \Delta \mathrm{h}_{2 \varepsilon} g, X\right\rangle g \mathrm{~d} \mathfrak{m} \stackrel{\sqrt{3.3 .6}}{=} \int \mathrm{h}_{\varepsilon}\left\langle\nabla \Delta \mathrm{h}_{\varepsilon} g, X\right\rangle g \mathrm{~d} \mathfrak{m}=\int\left\langle\nabla \Delta \mathrm{h}_{\varepsilon} g, X\right\rangle \mathrm{h}_{\varepsilon} g \mathrm{~d} \mathfrak{m} \tag{3.3.11}
\end{equation*}
$$

and, recalling that $\operatorname{div} X=0$, that

$$
\int\left\langle\nabla \Delta \mathrm{h}_{\varepsilon} g, X\right\rangle \mathrm{h}_{\varepsilon} g \mathrm{~d} \mathfrak{m}=-\int \Delta \mathrm{h}_{\varepsilon} g\left\langle X, \nabla \mathrm{~h}_{\varepsilon} g\right\rangle \mathrm{dm} \stackrel{\sqrt{3.3 .7}}{=}-\int \mathrm{h}_{\varepsilon} g\left\langle X, \nabla \Delta \mathrm{~h}_{\varepsilon} g\right\rangle \mathrm{dm} .
$$

This proves (3.3.10) and the theorem.
We therefore can conclude that:
Theorem 3.22. Let $(\mathrm{X}, \mathrm{d}, \mathfrak{m})$ be a $\operatorname{RCD}(0, \infty)$ space and $X$ a harmonic vector field. Then for every $t \in \mathbb{R}$ the map $\mathrm{Fl}_{t}^{(X)}$ has a continuous representative and this representative is a measure preserving isometry.
proof Use the preservation of measure proved in Proposition 3.19 and the one of Dirichlet energy proved in Proposition 3.21 in conjunction with Theorem 2.5.

From now on we shall identify $\mathrm{Fl}_{t}^{(X)}$ with its continuous representative. It is readily verified from the construction that the group property 3.3 .2 holds everywhere.

One of the consequences of the fact that $\mathrm{Fl}_{t}^{(X)}$ is an automorphism of $(\mathrm{X}, \mathrm{d}, \mathfrak{m})$ is that

$$
\begin{equation*}
f \in \operatorname{Test}(\mathrm{X}) \quad \Rightarrow \quad f \circ \mathrm{Fl}_{t}^{(X)} \in \operatorname{Test}(\mathrm{X}) . \tag{3.3.12}
\end{equation*}
$$

This can be seen by noticing that since $\mathrm{Fl}_{t}^{(X)}$ is a measure preserving isometry, directly from the definition of Sobolev space we have

$$
f \in W^{1,2}(\mathrm{X}) \quad \Leftrightarrow \quad f \circ \mathrm{Fl}_{t}^{(X)} \in W^{1,2}(\mathrm{X}) \quad \text { and in this case }|\mathrm{d} f| \circ \mathrm{Fl}_{t}^{(X)}=\left|\mathrm{d}\left(f \circ \mathrm{Fl}_{t}^{(X)}\right)\right| .
$$

From this fact and the definition of Laplacian we then deduce that

$$
f \in D(\Delta) \Leftrightarrow f \circ \mathrm{Fl}_{t}^{(X)} \in D(\Delta) \quad \text { and in this case }(\Delta f) \circ \mathrm{Fl}_{t}^{(X)}=\Delta\left(f \circ \mathrm{Fl}^{(X)}\right) .
$$

A suitable iteration of these arguments then yields (3.3.12).
Recall also (see [15]) that being $\mathrm{Fl}_{t}^{X}$ invertible and of bounded deformation, its differential $\mathrm{dFl}_{t}^{X}$ is a map from $L^{2}(T \mathrm{X})$ into itself (well) defined by:

$$
\begin{equation*}
\mathrm{d} f\left(\mathrm{~d} \mathrm{Fl}_{s}^{X}(Y)\right)=\mathrm{d}\left(f \circ \mathrm{Fl}_{s}^{X}\right)(Y) \circ \mathrm{Fl}_{-s}^{X} \quad \forall f \in W^{1,2}(\mathrm{X}) . \tag{3.3.13}
\end{equation*}
$$

We now want to prove that if $X, Y$ are both harmonic, their flows commute. The proof is based on the following lemma:

Lemma 3.23. Let $(\mathrm{X}, \mathrm{d}, \mathfrak{m})$ be a $\operatorname{RCD}(0, \infty)$ space and $X, Y$ harmonic vector fields. Then

$$
\mathrm{dFl}_{s}^{(X)}(Y)=Y \quad \forall s \in \mathbb{R}
$$

proof Since differential of test functions generate the whole cotangent module, the claim will follow if we show that for any $f \in \operatorname{Test}(\mathrm{X})$ the map $\mathbb{R} \ni s \mapsto \mathrm{~d} f\left(\mathrm{dFl}_{s}^{(X)}(Y)\right)$ is constant.

Taking account of the equality in (3.3.13) and recalling (3.3.12), in order to conclude it is sufficient to prove that for any $f \in \operatorname{Test}(\mathrm{X})$

$$
\frac{\mathrm{d}\left(f \circ \mathrm{Fl}_{h}^{(X)}\right)(Y) \circ \mathrm{Fl}_{-h}^{(X)}-\mathrm{d} f(Y)}{h} \text { goes to } 0 \text { in the strong } L^{2}(\mathrm{X}) \text {-topology as } h \rightarrow 0 .
$$

To this aim, start observing that
$\frac{\mathrm{d}\left(f \circ \mathrm{Fl}_{h}^{(X)}\right)(Y) \circ \mathrm{Fl}_{-h}^{(X)}-\mathrm{d} f(Y)}{h}=\mathrm{d}\left(\frac{f \circ \mathrm{Fl}_{h}^{(X)}-f}{h}\right)(Y) \circ \mathrm{Fl}_{-h}^{(X)}+\frac{\mathrm{d} f(Y) \circ \mathrm{Fl}_{-h}^{(X)}-\mathrm{d} f(Y)}{h}$.
Since $f \in \operatorname{Test}(\mathrm{X})$ and $X \in H_{C}^{1,2}(T \mathrm{X})$ we have $\mathrm{d} f(Y) \in W^{1,2}(\mathrm{X})$ (see [15 for details about this implication) and thus from the last claim in Proposition (2.7) we have

$$
\lim _{s \rightarrow 0} \frac{\mathrm{~d} f(Y) \circ \mathrm{Fl}_{-s}^{(X)}-\mathrm{d} f(Y)}{s}=-\mathrm{d}(\mathrm{~d} f(Y))(X) \quad \text { in } L^{2}(\mathrm{X}),
$$

hence to conclude it is sufficient to show that

$$
\begin{equation*}
\lim _{s \rightarrow 0} \mathrm{~d}\left(\frac{f \circ \mathrm{Fl}_{s}^{(X)}-f}{s}\right)(Y) \circ \mathrm{Fl}_{-s}^{(X)}=\mathrm{d}(\mathrm{~d} f(X))(Y) . \tag{3.3.14}
\end{equation*}
$$

Let us start proving that

$$
\begin{equation*}
\frac{f \circ \mathrm{Fl}_{s}^{(X)}-f}{s} \rightarrow \mathrm{~d} f(X) \quad \text { as } s \rightarrow 0 \quad \text { in } W^{1,2}(\mathrm{X}) \tag{3.3.15}
\end{equation*}
$$

Notice that 2.3.5) grants convergence in $L^{2}(\mathrm{X})$; moreover the bound

$$
\begin{aligned}
\left|\mathrm{d}\left(\frac{f \circ \mathrm{Fl}_{s}^{(X)}-f}{s}\right)\right|^{2} & =\left|\frac{1}{s} \int_{0}^{s} \mathrm{~d}\left(\mathrm{~d} f(X) \circ \mathrm{Fl}_{r}^{(X)}\right) \mathrm{d} r\right|^{2} \\
& \leq \frac{1}{s} \int_{0}^{s}\left|\mathrm{~d}\left(\mathrm{~d} f(X) \circ \mathrm{Fl}_{r}^{(X)}\right)\right|^{2} \mathrm{~d} r=\frac{1}{s} \int_{0}^{s}|\mathrm{~d}(\mathrm{~d} f(X))|^{2} \circ \mathrm{Fl}_{r}^{(X)} \mathrm{d} r
\end{aligned}
$$

and the fact that $\left(\mathrm{Fl}_{r}^{X}\right)_{*} \mathfrak{m}=\mathfrak{m}$ grant that $\varlimsup_{s \rightarrow 0}\left\|\frac{f \circ \mathrm{Fl}_{s}^{(X)}-f}{s}\right\|_{W^{1,2}} \leq\|\mathrm{d} f(X)\|_{W^{1,2}}$ which is sufficient to get (3.3.15).

From (3.3.15) we deduce that

$$
\mathrm{d}\left(\frac{f \circ \mathrm{Fl}_{s}^{(X)}-f}{s}\right)(Y) \quad \rightarrow \quad \mathrm{d}(\mathrm{~d} f(X))(Y) \quad \text { as } s \rightarrow 0 \quad \text { in } L^{2}(\mathrm{X})
$$

hence (3.3.14) follows from (2.3.5).
Theorem 3.24. Let $(\mathrm{X}, \mathrm{d}, \mathfrak{m})$ be a $\operatorname{RCD}(0, \infty)$ space and $X, Y \in L^{2}(T \mathrm{X})$ be two harmonic vector fields. Then for any $t, s \in \mathbb{R}$ it holds

$$
\mathrm{Fl}_{t}^{X} \circ \mathrm{Fl}_{s}^{Y}=\mathrm{Fl}_{1}^{t X+s Y}
$$

proof For any $r \in \mathbb{R}$ consider the map $G_{r}:=\mathrm{Fl}_{r t}^{X} \circ \mathrm{Fl}_{r s}^{Y}$. Now take $f \in W^{1,2}(\mathrm{X})$ and observe that $f \circ \mathrm{Fl}_{r t}^{X} \in W^{1,2}(\mathrm{X})$, as a consequence of Theorem 3.22 , and that from Proposition 2.7 it easily follows that $r \mapsto f \circ G_{r} \in L^{2}(\mathrm{X})$ is Lipschitz. By direct computation we have:

$$
\begin{equation*}
\frac{\mathrm{d}}{\mathrm{~d} r}\left(f \circ G_{r}\right)=\frac{\mathrm{d}}{\mathrm{~d} r}\left(f \circ \mathrm{Fl}_{r t}^{X} \circ \mathrm{Fl}_{r s}^{Y}\right)=s \mathrm{~d}\left(f \circ \mathrm{Fl}_{r t}^{X}\right)(Y) \circ \mathrm{Fl}_{r s}^{Y}+t \mathrm{~d} f(X) \circ \mathrm{Fl}_{r t}^{X} \circ \mathrm{Fl}_{r s}^{Y} . \tag{3.3.16}
\end{equation*}
$$

Using first identity (3.3.13) and then Lemma 3.23 we have

$$
\mathrm{d}\left(f \circ \mathrm{Fl}_{r t}^{X}\right)(Y) \circ \mathrm{Fl}_{r s}^{Y}=\mathrm{d} f\left(\mathrm{dFl}_{r t}^{X}(Y)\right) \circ \mathrm{Fl}_{r t}^{X} \circ \mathrm{Fl}_{r s}^{Y}=\mathrm{d} f(Y) \circ \mathrm{Fl}_{r t}^{X} \circ \mathrm{Fl}_{r s}^{Y}
$$

and thus from (3.3.16) we obtain

$$
\frac{\mathrm{d}}{\mathrm{~d} r}\left(f \circ G_{r}\right)=s \mathrm{~d} f(Y) \circ \mathrm{Fl}_{r t}^{X} \circ \mathrm{Fl}_{r s}^{Y}+t \mathrm{~d} f(X) \circ \mathrm{Fl}_{r t}^{X} \circ \mathrm{Fl}_{r s}^{Y}=\mathrm{d} f(t X+s Y) \circ G_{r}
$$

and since it is obvious by construction that $\left(G_{r}\right)$ has the properties $(i),(i i)$ in Definition 2.6. by Proposition 2.7 we deduce that $\left(G_{r}\right)$ is a Regular Lagrangian Flow of $t X+s Y$ and thus by the uniqueness part of Theorem 2.8 we deduce that for any $r \geq 0$ we have

$$
\mathrm{Fl}_{r t}^{X} \circ \mathrm{Fl}_{r s}^{Y}=\mathrm{Fl}_{r}^{t X+s Y}
$$

$\mathfrak{m}$-a.e.. In particular this holds for $r=1$ and since both sides are continuous functions, equality holds everywhere.

## 4 Proof of the Main Result

### 4.1 Setting

Here we fix the assumptions and notations that will be used in the rest of the text.
$(\mathrm{X}, \mathrm{d}, \mathfrak{m})$ is a $\operatorname{RCD}^{*}(0, N)$ metric measure space with $\operatorname{supp}(\mathfrak{m})=\mathrm{X}, N \in \mathbb{N}, N>0$, and such that $\operatorname{dim}\left(H_{\mathrm{dR}}^{k}\right)(\mathrm{X})=N$.

The theory developed in [15] grants the existence of $N$ harmonic vector fields $X_{1}, \ldots, X_{N}$ which are orthogonal in $L^{2}(T \mathrm{X})$. As in Section 3.3, since the Ricci curvature is nonnegative, these vector fields belong to $H_{C}^{1,2}(T \mathrm{X})$ and are parallel, i.e. $\nabla X_{i} \equiv 0$ for every $i$. It follows that $\left\langle X_{i}, X_{j}\right\rangle \in W^{1,2}(\mathrm{X})$ with

$$
\mathrm{d}\left\langle X_{i}, X_{j}\right\rangle=\nabla X_{i}\left(\cdot, X_{j}\right)+\nabla X_{j}\left(\cdot, X_{i}\right)=0 \quad \mathfrak{m}-\text { a.e. },
$$

which in turn grants that $\left\langle X_{i}, X_{j}\right\rangle$ is $\mathfrak{m}$-a.e. equal to a constant function. Since $\int\left\langle X_{i}, X_{j}\right\rangle \mathrm{d} \mathfrak{m}=$ 0 for $i \neq j$ we conclude that the $X_{i}$ 's are pointwise orthogonal. The same argument also shows that up to normalization we can, and will, assume that $\left|X_{i}\right| \equiv 1 \mathfrak{m}$-a.e. for every $i$. In particular, since these vector fields are in $L^{2}(T \mathrm{X})$, we have

$$
\begin{equation*}
\mathfrak{m}(\mathrm{X})<\infty \tag{4.1.1}
\end{equation*}
$$

We shall work with the product space $\mathrm{X} \times \mathbb{R}^{N}$ which will be equipped with the measure $\mathfrak{m} \times \mathcal{L}^{N}$ and the distance

$$
\left(\mathrm{d} \otimes \mathrm{~d}_{\text {Eucl }}\right)^{2}((x, a),(y, b)):=\mathrm{d}^{2}(x, y)+|a-b|^{2} .
$$

We shall also define vector fields $Y_{i} \in L^{0}\left(\mathrm{X} \times \mathbb{R}^{N}\right), i=1, \ldots, N$ as

$$
Y_{i}:=\Phi_{2}\left(\widehat{\nabla \pi_{i}}\right) \quad \forall i=1, \ldots, N
$$

where $\pi_{i}: \mathbb{R}^{N} \rightarrow \mathbb{R}$ is the projection on the $i$-th coordinate, $\widehat{\nabla \pi_{i}} \in L^{0}\left(\mathrm{X}, L^{0}\left(T \mathbb{R}^{N}\right)\right)$ is the function identically equal to $\nabla \pi^{i} \in L^{0}\left(T \mathbb{R}^{N}\right)$ and where $\Phi_{2}: L^{0}\left(\mathrm{X}, L^{0}\left(T \mathbb{R}^{N}\right)\right) \rightarrow$ $L^{0}\left(T\left(\mathrm{X} \times \mathbb{R}^{N}\right)\right)$ is defined in Proposition 3.7.

We also define the map $\mathrm{T}: \mathrm{X} \times \mathbb{R}^{N} \rightarrow \mathrm{X}$ by

$$
\begin{align*}
\mathrm{T}: \mathrm{X} \times \mathbb{R}^{N} & \rightarrow \mathrm{X} \\
\left(x, \underline{a}=\left(a_{1}, \ldots a_{N}\right)\right) & \mapsto \mathrm{Fl}_{a_{1}}^{\left(X_{1}\right)} \circ \ldots \circ \mathrm{Fl}_{a_{N}}^{\left(X_{N}\right)}(x) . \tag{4.1.2}
\end{align*}
$$

### 4.2 Preliminary considerations

Let us collect some easy consequences of our assumptions that can be derived from the discussion made in the previous sections. Start recalling from [6] (see also [5]) that the product of two $\operatorname{RCD}(0, \infty)$ spaces is also $\operatorname{RCD}(0, \infty)$, so that $\mathrm{X} \times \mathbb{R}^{N}$ is $\operatorname{RCD}(0, \infty)$.

From the fact that the $\nabla \pi_{i}$ are a pointwise ortonormal base for $L^{0}\left(T \mathbb{R}^{N}\right)$ and the fact that $\Phi_{2}$ preserves the pointwise norm we deduce that

$$
\left\langle Y_{i}, Y_{j}\right\rangle=\delta_{i j} \quad \mathfrak{m} \times \mathcal{L}^{N}-\text { a.e. } \quad \forall i, j
$$

and from the very definition of $\Phi_{2}$ we have that

$$
\begin{equation*}
Y_{i}=\nabla\left(\pi_{i} \circ \pi^{\mathbb{R}^{N}}\right) \tag{4.2.1}
\end{equation*}
$$

Since $\pi_{i}: \mathbb{R}^{N} \rightarrow \mathbb{R}$ is harmonic, we have $\nabla \pi_{i} \in D\left(\operatorname{div}_{\text {loc }}, \mathbb{R}^{N}\right)$ with $\operatorname{div}\left(\nabla \pi_{i}\right)=0$ and thus from Propositions 3.15 and 3.11 we deduce that $Y_{i} \in D\left(\operatorname{div}_{\text {loc }}, \mathrm{X} \times \mathbb{R}^{N}\right)$ with

$$
\begin{equation*}
\operatorname{div} Y_{i} \equiv 0 \tag{4.2.2}
\end{equation*}
$$

Taking into account 2.2.7) we also obtain that $Y_{i}^{b} \in D\left(\Delta_{H, \text { loc }}\right)$ with $\Delta_{H} Y_{i}^{b}=0$ and since $\left|\nabla \pi_{i}\right| \equiv 1$ and $\Phi_{2}$ preserves the pointwise norm we also deduce that $\left|Y_{i}\right| \equiv 1$ : these facts together with (2.2.8) grant that $Y_{i} \in W_{C, \text { loc }}^{1,2}(T \mathrm{X})$ with

$$
\begin{equation*}
\nabla Y_{i} \equiv 0 \tag{4.2.3}
\end{equation*}
$$

Concerning the map T , from Theorem 3.24 we deduce that

$$
\begin{equation*}
\mathrm{T}(\mathrm{~T}(x, \underline{a}), \underline{b})=\mathrm{T}(x, \underline{a}+\underline{b}) \quad \forall x \in \mathrm{X}, \underline{a}, \underline{b} \in \mathbb{R}^{N}, \tag{4.2.4}
\end{equation*}
$$

so that we shall occasionally think at T as an action of $\mathbb{R}^{N}$ on X . Theorem 3.22 grants that this action is made of isometries, i.e.

$$
\begin{equation*}
\mathrm{T}(\cdot, \underline{a}): \mathrm{X} \rightarrow \mathrm{X} \quad \text { is an isometry for any } \underline{a} \in \mathbb{R}^{N} . \tag{4.2.5}
\end{equation*}
$$

From Theorem 3.24 we also have

$$
\mathrm{T}(x, \underline{a})=\mathrm{Fl}_{1}^{\left(X_{\underline{a}}\right)}(x) \quad \text { for } \quad X_{\underline{a}}:=\sum_{i=1}^{N} a_{i} X_{i}
$$

and since the pointwise orthonormality of the $X_{i}$ 's gives $\left|X_{\underline{a}}\right|^{2}=\left|\sum_{i=1}^{N} a_{i} X_{i}\right|^{2}=\sum_{i}\left|a_{i}\right|^{2}=$ $|\underline{a}|^{2} \mathfrak{m}$-a.e., from 2.3 .3 we deduce that for $\mathfrak{m}$-a.e. $x$ it holds

$$
\mathrm{d}(x, \mathrm{~T}(x, \underline{a})) \leq \int_{0}^{1}\left|X_{\underline{a}}\right| \circ \mathrm{Fl}_{t}^{\left(X_{\underline{a}}\right)} \mathrm{d} t \leq\left\|\left|X_{\underline{a}}\right|\right\|_{L^{\infty}}=|\underline{a}| .
$$

Now the continuity of $\mathrm{T}(\cdot, \underline{a})$ ensures that the above holds for every $x \in \mathrm{X}$ and thus taking (4.2.4) into account we conclude that

$$
\begin{equation*}
\mathrm{T}(x, \cdot): \mathbb{R}^{N} \rightarrow \mathrm{X} \quad \text { is } 1 \text {-Lipschitz for any } x \in \mathrm{X} \tag{4.2.6}
\end{equation*}
$$

Finally we remark that Proposition 3.19 together with Fubini theorem guarantees that

$$
\begin{equation*}
\mathrm{T}_{*}\left(\mathfrak{m} \times\left.\mathscr{L}^{N}\right|_{A}\right)=\mathscr{L}^{N}(A) \mathfrak{m} \tag{4.2.7}
\end{equation*}
$$

for every $A \subset \mathbb{R}^{N}$ Borel. This identity, 4.2.5 and (4.2.6) grant in particular that T: $\mathrm{X} \times \mathbb{R}^{N} \rightarrow \mathrm{X}$ is of local bounded deformation.

### 4.3 An explicit formula for Regular Lagrangian Flows on X

Aim of this section is to provide, in Proposition 4.4, an explicit representation formula for Regular Lagrangian Flows on X in terms of the map T. The starting point is the following:

Proposition 4.1 (Conjugation property). With the same notations and assumptions as in Section 4.1 the following holds. For any $i=1, \ldots, N, t \in \mathbb{R}$ and $f \in W_{\text {loc }}^{1,2}(\mathrm{X})$ it holds

$$
\begin{equation*}
\mathrm{d}(f \circ \mathrm{~T})\left(Y_{i}\right)=\mathrm{d} f\left(X_{i}\right) \circ \mathrm{T} \quad \mathfrak{m} \times \mathcal{L}^{N}-\text { a.e.. } \tag{4.3.1}
\end{equation*}
$$

proof Let us put $\overline{\mathrm{Fl}}_{t}^{i}(x, \underline{a}):=\left(x, \underline{a}+t e_{i}\right)$ and notice that by the very definition of T and identity (4.2.4) we have

$$
\begin{equation*}
\mathrm{Fl}_{t}^{X_{i}} \circ \mathrm{~T}=\mathrm{T} \circ \overline{\mathrm{Fl}}_{t}^{i} . \tag{4.3.2}
\end{equation*}
$$

Fix $f \in W_{\text {loc }}^{1,2}(\mathrm{X})$ and recall Proposition 2.7 to get

$$
\mathrm{d} f\left(X_{i}\right) \circ \mathrm{T}=\left(\lim _{t \downarrow 0} \frac{f \circ \mathrm{Fl}_{t}^{\left(X_{i}\right)}-f}{t}\right) \circ \mathrm{T} \stackrel{\sqrt{4.3 .2}}{=} \lim _{t \downarrow 0} \frac{f \circ \mathrm{~T} \circ \overline{\mathrm{Fl}}_{t}^{i}-f \circ \mathrm{~T}}{t},
$$

the first limit being in $L_{\text {loc }}^{2}(\mathrm{X})$ and the second in $L_{\text {loc }}^{2}\left(\mathrm{X} \times \mathbb{R}^{N}\right)$. Hence to conclude it is sufficient to show that for any $\tilde{f} \in W_{\text {loc }}^{1,2}\left(\mathrm{X} \times \mathbb{R}^{N}\right)$ and $\rho \in L^{\infty}\left(\mathrm{X} \times \mathbb{R}^{N}\right)$ with bounded support we have

$$
\begin{equation*}
\lim _{t \downarrow 0} \int \frac{\tilde{f} \circ \overline{\mathrm{~F}} t_{t}^{i}-\tilde{f}}{t} \rho \mathrm{~d}\left(\mathfrak{m} \times \mathcal{L}^{N}\right)=\int \mathrm{d} \tilde{f}\left(Y_{i}\right) \rho \mathrm{d}\left(\mathfrak{m} \times \mathcal{L}^{N}\right) \tag{4.3.3}
\end{equation*}
$$

By the linearity in $\rho$ of this expression we can further assume that $\rho$ is a probability density. Then put $\mu:=\rho \mathfrak{m}$ and $\boldsymbol{\pi}:=\left(\overline{\mathrm{F}} ._{.}^{i}\right)_{*} \mu$, where here $\overline{\mathrm{Fl}}{ }^{i}: \mathrm{X} \times \mathbb{R}^{N} \rightarrow C\left([0,1], \mathrm{X} \times \mathbb{R}^{N}\right)$ is the map sending $(x, \underline{a})$ to the curve $[0,1] \ni t \mapsto \mathrm{Fl}_{t}^{i}(x, \underline{a})$. Notice that $\boldsymbol{\pi}$ is a test plan on $\mathrm{X} \times \mathbb{R}^{N}$ which is concentrated on curves with speed constantly equal to 1 , thus for any $\tilde{f} \in W_{\text {loc }}^{1,2}\left(\mathrm{X} \times \mathbb{R}^{N}\right)$ we have

$$
\begin{aligned}
\frac{\int \tilde{f} \mathrm{~d}\left(\overline{\mathrm{~F}} \mathrm{l}_{t}^{i}\right)_{*} \mu-\int \tilde{f} \mathrm{~d} \mu}{t} & =\frac{1}{t} \int \tilde{f}\left(\gamma_{t}\right)-\tilde{f}\left(\gamma_{0}\right) \mathrm{d} \boldsymbol{\pi}(\gamma) \\
& \leq \frac{1}{t} \iint_{0}^{t}|\mathrm{~d} \tilde{f}|\left(\gamma_{s}\right)\left|\dot{\gamma}_{s}\right| \mathrm{d} s \mathrm{~d} \boldsymbol{\pi}(\gamma) \\
& \leq \frac{1}{2 t} \iint_{0}^{t}|\mathrm{~d} \tilde{f}|^{2}\left(\gamma_{s}\right) \mathrm{d} s \mathrm{~d} \boldsymbol{\pi}(\gamma)+\frac{1}{2 t} \iint_{0}^{t}\left|\dot{\gamma}_{s}^{2}\right| \mathrm{d} s \mathrm{~d} \boldsymbol{\pi}(\gamma) \\
& =\frac{1}{2 t} \int_{0}^{t} \int|\mathrm{~d} \tilde{f}|^{2} \circ \overline{\mathrm{~F}}_{s}^{i} \mathrm{~d} \mu \mathrm{~d} s+\frac{1}{2}
\end{aligned}
$$

Recalling (2.3.5) we thus have

$$
\begin{equation*}
\varlimsup_{t \downarrow 0} \frac{\int \tilde{f} \mathrm{~d}\left(\overline{\mathrm{Fl}}{ }_{t}^{i}\right)_{*} \mu-\int \tilde{f} \mathrm{~d} \mu}{t} \leq \frac{1}{2} \int|\mathrm{~d} \tilde{f}|^{2} \mathrm{~d} \mu+\frac{1}{2} \tag{4.3.4}
\end{equation*}
$$

Now put for brevity $f_{i}:=\pi^{i} \circ \pi^{R^{N}}$, so that $f_{i}$ is 1 -Lipschitz, 4.2.1 reads as

$$
\begin{equation*}
\nabla f_{i}=Y_{i} \tag{4.3.5}
\end{equation*}
$$

and by construction it holds $f_{i} \circ \overline{\mathrm{~F}}_{s}^{i}=f_{i}+s$, so that

$$
\begin{equation*}
\varliminf_{t \downarrow 0} \frac{\int f_{i} \mathrm{~d}\left(\overline{\mathrm{Fl}}_{t}^{i}\right)_{*} \mu-\int f_{i} \mathrm{~d} \mu}{t} \geq 1 \geq \frac{1}{2} \int\left|\mathrm{~d} f_{i}\right|^{2} \mathrm{~d} \mu+\frac{1}{2} \tag{4.3.6}
\end{equation*}
$$

(In the terminology of [17] we just proved that $\boldsymbol{\pi}$ represents the gradient of $f_{i}$ and we are now going to use the link between 'horizontal and vertical' derivatives). Writing (4.3.4) for $f+\varepsilon \tilde{f}$ in place of $\tilde{f}$ and subtracting 4.3.6 we obtain

$$
\varlimsup_{t \downarrow 0} \varepsilon \frac{\int \tilde{f} \mathrm{~d}\left(\overline{\mathrm{~F}}_{t}^{i}\right)_{*} \mu-\int \tilde{f} \mathrm{~d} \mu}{t} \leq \frac{1}{2} \int|\mathrm{~d}(f+\varepsilon \tilde{f})|^{2}-|\mathrm{d} f|^{2} \mathrm{~d} \mu \stackrel{(4.3 .5)}{=} \int \varepsilon \mathrm{d} \tilde{f}\left(Y_{i}\right)+\frac{\varepsilon^{2}}{2}|\mathrm{~d} \tilde{f}|^{2} \mathrm{~d} \mu
$$

Dividing by $\varepsilon>0$ (resp. $\varepsilon<0$ ) and letting $\varepsilon \downarrow 0$ (resp. $\varepsilon \uparrow 0$ ) we obtain 4.3.3) and the conclusion.

We now introduce a map $\Psi: L^{0}(T \mathrm{X}) \rightarrow L^{0}\left(T\left(\mathrm{X} \times \mathbb{R}^{N}\right)\right)$ as

$$
\begin{equation*}
\Psi(v):=\sum_{i=1}^{N}\left\langle v, X_{i}\right\rangle \circ \mathrm{T} Y_{i} \tag{4.3.7}
\end{equation*}
$$

Lemma 4.2. With the same notations and assumptions as in Section 4.1 and with $\Psi$ defined as in 4.3.7), the following holds. Let $v \in L^{\infty} \cap W_{C}^{1,2}(T \mathrm{X})$. Then:
i) $\left\langle v, X_{i}\right\rangle \in W^{1,2}(\mathrm{X})$ for every $i$ and $v \in D($ div $)$ with

$$
\begin{equation*}
\operatorname{div}(v)=\sum_{i=1}^{N} \mathrm{~d}\left(\left\langle v, X_{i}\right\rangle\right)\left(X_{i}\right) \tag{4.3.8}
\end{equation*}
$$

ii) $\Psi(v) \in L^{\infty} \cap W_{C, \text { loc }}^{1,2} \cap D\left(\operatorname{div}_{\text {loc }}\right)\left(T\left(\mathrm{X} \times \mathbb{R}^{n}\right)\right)$ with

$$
\begin{aligned}
\nabla(\Psi(v)) & =\sum_{i=1}^{N} \nabla\left(\left\langle v, X_{i}\right\rangle \circ \mathrm{T}\right) \otimes Y_{i}, \\
\operatorname{div}(\Psi(v)) & =\operatorname{div}(v) \circ \mathrm{T} .
\end{aligned}
$$

proof
(i) From [15] we know that the assumptions on $v$ grant that $\langle v, X\rangle \in W_{\text {loc }}^{1,2}(\mathrm{X})$ for every $X \in L^{\infty} \cap H_{C}^{1,2}(T \mathrm{X})$ with

$$
\mathrm{d}\langle v, X\rangle=\nabla v(\cdot, X)+\nabla X(\cdot, v)
$$

Picking $X:=X_{i}$ and recalling that $D\left(\Delta_{H}\right) \subset\left(H_{C}^{1,2}(T \mathrm{X})\right)^{b}$ by the very definition of $\Delta_{H}$, we conclude that $\left\langle v, X_{i}\right\rangle$ belongs to $W^{1,2}(\mathrm{X})$, as claimed. Now put $a_{i}:=\left\langle v, X_{i}\right\rangle$ for brevity, so that $v=\sum_{i} a_{i} X_{i}$, let $f \in W^{1,2}(\mathrm{X})$ and notice that

$$
\int \mathrm{d} f(v) \mathrm{d} \mathfrak{m}=\sum_{i=1}^{N} \int \mathrm{~d} f\left(a_{i} X_{i}\right) \mathrm{d} \mathfrak{m}=-\sum_{i=1}^{N} \int f \operatorname{div}\left(a_{i} X_{i}\right)=-\int \sum_{i=1}^{N} f \mathrm{~d} a_{i}\left(X_{i}\right) \mathrm{d} \mathfrak{m}
$$

having used the fact that $\operatorname{div}\left(X_{i}\right)=0$. This proves both $v \in D$ (div) and 4.3.8).
(ii) The assumption $v \in L^{\infty}(T \mathrm{X})$ trivially yields $a_{i} \in L^{\infty}(\mathrm{X})$ and since $Y_{i} \in L^{\infty} \cap W_{C}^{1,2}(\mathrm{X} \times$ $\left.\mathbb{R}^{N}\right)$ with $\nabla Y_{i}=0($ recall 4.2.3) $)$, we have $\left(a_{i} \circ \mathrm{~T}\right) Y_{i} \in W_{C, \text { loc }}^{1,2}\left(T\left(\mathrm{X} \times \mathbb{R}^{N}\right)\right)$ with

$$
\nabla\left(\left(a_{i} \circ \mathrm{~T}\right) Y_{i}\right)=\nabla\left(a_{i} \circ \mathrm{~T}\right) \otimes Y_{i}+a_{i} \circ \mathrm{~T} \nabla Y_{i}=\nabla\left(a_{i} \circ \mathrm{~T}\right) \otimes Y_{i} .
$$

The fact that $\Phi(v) \in L^{\infty} \cap W_{C, \text { loc }}^{1,2}\left(T\left(\mathrm{X} \times \mathbb{R}^{n}\right)\right.$ and the formula for $\nabla(\Psi(v))$ follow.
We turn to the divergence: for $g \in W^{1,2}\left(\mathrm{X} \times \mathbb{R}^{N}\right)$ with bounded support we have

$$
\begin{aligned}
\int \mathrm{d} g(\Psi(v)) \mathrm{d}\left(\mathfrak{m} \times \mathcal{L}^{N}\right) & =\sum_{i=1}^{N} \int \mathrm{~d} g\left(a_{i} \circ \mathrm{~T} Y_{i}\right) \mathrm{d}\left(\mathfrak{m} \times \mathcal{L}^{N}\right) \\
\text { because } \operatorname{div}\left(Y_{i}\right)=0 & =-\sum_{i=1}^{N} \int g \mathrm{~d}\left(a_{i} \circ \mathrm{~T}\right)\left(Y_{i}\right) \mathrm{d}\left(\mathfrak{m} \times \mathcal{L}^{N}\right) \\
\text { by 4.3.1 } & =-\sum_{i=1}^{N} \int g \mathrm{~d} a_{i}\left(X_{i}\right) \circ \mathrm{T} \mathrm{~d}\left(\mathfrak{m} \times \mathcal{L}^{N}\right),
\end{aligned}
$$

which by 4.3.8) is the conclusion.
Let $\left(v_{t}\right) \in L^{\infty}\left([0,1], L^{2}(T \mathrm{X})\right) \cap L^{2}\left([0,1], W_{C}^{1,2}(T \mathrm{X})\right)$ be such that $\left(\operatorname{div}\left(v_{t}\right)\right) \in L^{\infty}\left([0,1], L^{\infty}(\mathrm{X})\right)$, so that in particular the Regular Lagrangian Flow $\left(\mathrm{Fl}_{s}^{\left(v_{t}\right)}\right)$ is well defined. The integrability condition of $\left(v_{t}\right)$ ensures that $\left(\left\langle v_{t}, X_{i}\right\rangle\right) \in L^{\infty}\left([0,1], L^{2}(\mathrm{X})\right)$ for every $i=1, \ldots, N$ and thus from (2.3.1) we see that $\left(\left\langle v_{s}, X_{i}\right\rangle \circ \mathrm{Fl}_{s}^{\left(v_{t}\right)}\right) \in L^{\infty}\left([0,1], L^{2}(\mathrm{X})\right)$ as well. Hence the functions

$$
\begin{equation*}
A_{i, t}:=\int_{0}^{t}\left\langle v_{s}, X_{i}\right\rangle \circ \mathrm{Fl}_{s}^{\left(v_{t}\right)} \mathrm{d} s \in L^{2}(\mathrm{X}), \quad t \in[0,1], i=1, \ldots, N, \tag{4.3.9}
\end{equation*}
$$

are well defined. We then have the following result:
Lemma 4.3. With the same assumptions and notation as in Section 4.1, let $\left(v_{t}\right) \in$ $L^{\infty}\left([0,1], L^{2}(T \mathrm{X})\right) \cap L^{2}\left([0,1], W_{C}^{1,2}(T \mathrm{X})\right)$ be such that $\left(\operatorname{div}\left(v_{t}\right)\right) \in L^{\infty}\left([0,1], L^{\infty}(\mathrm{X})\right)$ and define $\Psi$ as in 4.3.7).

Then the vector fields $\Psi\left(v_{t}\right)$ satisfy the assumptions of Theorem 2.8 and for any $s \in \mathbb{R}$ the following identities hold $\mathfrak{m} \times \mathcal{L}^{N}$-a.e.:

$$
\begin{align*}
\mathrm{T} \circ \mathrm{Fl}_{s}^{\left(\Psi\left(v_{t}\right)\right)} & =\mathrm{Fl}_{s}^{\left(v_{t}\right)} \circ \mathrm{T},  \tag{4.3.10}\\
\pi^{\mathrm{X}} \circ \mathrm{Fl}_{s}^{\left(\Psi\left(v_{t}\right)\right)} & =\pi^{\mathrm{X}}  \tag{4.3.11}\\
\pi_{i} \circ \pi^{\mathbb{R}^{N}} \circ \mathrm{Fl}_{s}^{\left(\Psi\left(v_{t}\right)\right)} & =\pi_{i} \circ \pi^{\mathbb{R}^{N}}+A_{i, s} \circ \mathrm{~T} . \tag{4.3.12}
\end{align*}
$$

proof The fact that the $\Psi\left(v_{t}\right)$ 's satisfy the assumptions of Theorem 2.8 is a direct consequence of the assumptions and Lemma 4.2. Also, from 4.3.1) we directly deduce

$$
\begin{equation*}
\mathrm{d}(f \circ \mathbf{T})\left(\Psi\left(v_{t}\right)\right)=\mathrm{d} f\left(v_{t}\right) \circ \mathbf{T} \quad \mathfrak{m} \times \mathcal{L}^{N}-\text { a.e. } \tag{4.3.13}
\end{equation*}
$$

for every $t \in[0,1]$ and $f \in W_{\text {loc }}^{1,2}(\mathrm{X})$. Now pick $\bar{\mu}_{0} \in \mathscr{P}\left(\mathrm{X} \times \mathbb{R}^{N}\right)$ with bounded support and such that $\bar{\mu}_{0} \leq C \mathfrak{m} \times \mathcal{L}^{N}$ for some $C>0$ and for $s \in \mathbb{R}$ define

$$
\bar{\mu}_{s}:=\left(\mathrm{Fl}_{s}^{\left(\Psi\left(v_{t}\right)\right)}\right)_{*} \bar{\mu}_{0} \in \mathscr{P}\left(\mathrm{X} \times \mathbb{R}^{N}\right) \quad \text { and } \quad \mu_{s}:=\mathrm{T}_{*} \bar{\mu}_{s} \in \mathscr{P}(\mathrm{X})
$$

We claim that $\left(\mu_{s}\right)$ solves the continuity equation with vector fields $\left(v_{s}\right)$ in the sense of Definition 2.9 and start observing that, locally in $s$, the measures $\bar{\mu}_{s}, \mu_{s}$ have uniformly bounded density. Now pick $f \in W^{1,2}(\mathrm{X})$, so that $f \circ \mathrm{~T} \in W_{\mathrm{loc}}^{1,2}\left(\mathrm{X} \times \mathbb{R}^{N}\right)$ and, from Proposition 2.7, $s \mapsto \int f \mathrm{~d} \mu_{s}=\int f \circ \mathrm{~T} \circ\left(\mathrm{Fl}_{s}^{\left(\Psi\left(v_{t}\right)\right)}\right) \mathrm{d} \bar{\mu}_{0}$ is Lipschitz with

$$
\begin{aligned}
\frac{\mathrm{d}}{\mathrm{~d} s} \int f \mathrm{~d} \mu_{s} & =\frac{\mathrm{d}}{\mathrm{~d} s} \int f \circ \mathrm{~T} \circ\left(\mathrm{Fl}_{s}^{\left(\Psi\left(v_{t}\right)\right)}\right) \mathrm{d} \bar{\mu}_{0} \\
& =\int \mathrm{d}(f \circ \mathrm{~T})\left(\Psi\left(v_{s}\right)\right) \circ\left(\mathrm{Fl}_{s}^{\left(\Psi\left(v_{t}\right)\right)}\right) \mathrm{d} \bar{\mu}_{0} \\
& =\int \mathrm{d}(f \circ \mathrm{~T})\left(\Psi\left(v_{s}\right)\right) \mathrm{d} \bar{\mu}_{s} \\
& =\int \mathrm{d} f\left(v_{s}\right) \circ \mathrm{T} \mathrm{~d} \bar{\mu}_{s} \\
& =\int \mathrm{d} f\left(v_{s}\right) \mathrm{d} \mu_{s} .
\end{aligned}
$$

This proves our claim. Hence by the representation formula in Theorem 2.10 we deduce that

$$
\mathrm{T}_{*}\left(\mathrm{Fl}_{s}^{\left(\Psi\left(v_{t}\right)\right)}\right)_{*} \bar{\mu}_{0}=\left(\mathrm{Fl}_{s}^{\left(v_{t}\right)}\right)_{*} \mathrm{~T}_{*} \bar{\mu}_{0} \quad \forall s \in \mathbb{R}
$$

and from the arbitrariness of $\bar{\mu}_{0}$ and (3.3.1) identity 4.3.10 follows.
To prove 4.3.11) pick $\bar{\mu}_{0}, f$ and define $\bar{\mu}_{s}$ as above. Then we also put

$$
\nu_{s}:=\pi_{*}^{\mathrm{X}} \bar{\mu}_{s} \in \mathscr{P}(\mathrm{X}) \quad \forall s \in \mathbb{R}
$$

and notice that again the $\nu_{s}$ 's have, locally in $s$, uniformly bounded densities and that it holds

$$
\frac{\mathrm{d}}{\mathrm{~d} s} \int f \mathrm{~d} \nu_{s}=\frac{\mathrm{d}}{\mathrm{~d} s} \int f \circ \pi^{\mathrm{x}} \mathrm{~d} \bar{\mu}_{s}=\int \mathrm{d}\left(f \circ \pi^{\mathrm{x}}\right)\left(\Psi\left(v_{s}\right)\right) \mathrm{d} \bar{\mu}_{s}=\int \Phi_{1}(\widehat{\mathrm{~d} f})\left(\Psi\left(v_{s}\right)\right) \mathrm{d} \bar{\mu}_{s}=0
$$

where as usual $\widehat{\mathrm{dff}} \in L^{0}\left(\mathbb{R}^{N}, L^{0}\left(T^{*} \mathrm{X}\right)\right)$ is the function identically equal to $\mathrm{d} f$ and the last identity follows from 3.2.16 and the very definitions of $\Psi\left(v_{t}\right)$ and $Y_{i}$. This shows that
$\left(\nu_{s}\right)$ solves the continuity equation 2.3 .11 with 0 vector fields, hence by the uniqueness of the solutions we deduce that $\left(\nu_{s}\right)$ is constant, i.e.

$$
\pi_{*}^{\mathrm{X}}\left(\mathrm{Fl}_{s}^{\left(\Psi\left(v_{t}\right)\right)}\right)_{*} \bar{\mu}_{0}=\pi_{*}^{\mathrm{X}} \bar{\mu}_{0},
$$

so that again the arbitrariness of $\bar{\mu}_{0}$ and (3.3.1) give 4.3.11).
For 4.3.12), we notice that the two sides agree for $s=0$ and are absolutely continuous as functions of $s$ with values in $L_{\text {loc }}^{2}\left(\mathrm{X} \times \mathbb{R}^{N}\right)$ (recall Proposition 2.7). The conclusion then follows recalling that it holds $\nabla\left(\pi_{i} \circ \pi^{\mathbb{R}^{N}}\right)=Y_{i}$, so that

$$
\begin{aligned}
\frac{\mathrm{d}}{\mathrm{~d} s} \pi_{i} \circ \pi^{\mathbb{R}^{N}} \circ \mathrm{Fl}_{s}^{\left(\Psi\left(v_{t}\right)\right)} & =\mathrm{d}\left(\pi_{i} \circ \pi^{\mathbb{R}^{N}}\right)\left(\Psi\left(v_{s}\right)\right) \circ \mathrm{Fl}_{s}^{\left(\Psi\left(v_{t}\right)\right)}=\left\langle Y_{i}, \Psi\left(v_{s}\right)\right\rangle \circ \mathrm{Fl}_{s}^{\left(\Psi\left(v_{t}\right)\right)} \\
& =\left\langle v_{s}, X_{i}\right\rangle \circ \mathrm{T} \circ \mathrm{Fl}_{s}^{\left(\Psi\left(v_{t}\right)\right)} \stackrel{4.3 .10}{=}\left\langle v_{s}, X_{i}\right\rangle \circ \mathrm{Fl}_{s}^{\left(v_{t}\right)} \circ \mathrm{T}=\frac{\mathrm{d}}{\mathrm{~d} s} A_{i, s} \circ \mathrm{~T} .
\end{aligned}
$$

This is sufficient to conclude.
We can now state the main result of the section:
Proposition 4.4 (Representation formula for $\mathrm{Fl}^{\left(v_{t}\right)}$ ). With the same assumptions and notation as in Section 4.1. let $\left(v_{t}\right) \in L^{\infty}\left([0,1], L^{2}(T \mathrm{X})\right) \cap L^{2}\left([0,1], W_{C}^{1,2}(T \mathrm{X})\right)$ be such that $\left(\operatorname{div}\left(v_{t}\right)\right) \in L^{\infty}\left([0,1], L^{\infty}(\mathrm{X})\right)$ and define the functions $A_{i, t} \in L^{2}(\mathrm{X})$ as in 4.3.9.).

Then for any $s \in[0,1]$ and $\mathfrak{m}$-a.e. $x \in \mathrm{X}$ it holds

$$
\begin{equation*}
\mathrm{Fl}_{s}^{\left(v_{t}\right)}(x)=\mathrm{T}\left(x, \underline{A}_{s}(x)\right), \tag{4.3.14}
\end{equation*}
$$

where $\underline{A}_{s}:=\left(A_{1, s}, \ldots, A_{N, s}\right)$.
proof The identities (4.3.11), 4.3.12) give

$$
\mathrm{Fl}_{s}^{\left(\Psi\left(v_{t}\right)\right)}(x, \underline{a})=\left(x, \underline{a}+\underline{A}_{s}(\mathrm{~T}(x, \underline{a}))\right) \quad \mathfrak{m} \times \mathcal{L}^{N}-a . e .(x, \underline{a}) .
$$

Applying T on both sides and taking into account 4.3.10 and 4.2.4 we obtain

$$
\begin{equation*}
\mathrm{Fl}_{s}^{\left(v_{t}\right)}(\mathrm{T}(x, \underline{a}))=\mathrm{T}\left(\mathrm{~T}(x, \underline{a}), \underline{A}_{s}(\mathrm{~T}(x, \underline{a}))\right) \quad \mathfrak{m} \times \mathcal{L}^{N}-\text { a.e. }(x, \underline{a}) . \tag{4.3.15}
\end{equation*}
$$

Thus for any $A \subset \mathbb{R}^{N}$ Borel we have that (4.3.14) holds for $\mathrm{T}_{*}\left(\mathfrak{m} \times\left.\mathcal{L}^{N}\right|_{A}\right)$-a.e. $x \in \mathrm{X}$ and the conclusion follows from (4.2.7).

### 4.4 Further properties of $T$ and conclusion

We shall need the following result, proved in [20], about $W_{2}$-geodesics and continuity equation. Recall that a measure $\boldsymbol{\pi} \in \mathscr{P}(C([0,1], \mathrm{X}))$ is called lifting of the geodesic $\left(\mu_{t}\right)$ provided

$$
\begin{aligned}
\left(\mathrm{e}_{t}\right)_{*} \boldsymbol{\pi} & =\mu_{t} \quad \forall t \in[0,1] \\
\iint_{0}^{1}\left|\dot{\gamma}_{t}\right|^{2} \mathrm{~d} t \mathrm{~d} \boldsymbol{\pi}(\gamma) & <+\infty
\end{aligned}
$$

and that, on an $\operatorname{RCD}^{*}(K, N)$ space, as soon as either $\mu_{0}$ or $\mu_{1}$ is absolutely continuous w.r.t. $\mathfrak{m}$, there is a unique geodesic connecting them and a unique lifting of it (see [22]).

Proposition 4.5. Let (X, d, m) be a $\mathrm{RCD}^{*}(K, N)$ space and $\left(\mu_{t}\right) \subset \mathscr{P}(\mathrm{X})$ be a $W_{2}$-geodesic such that $\mu_{0}, \mu_{1}$ have both bounded support and density. Then there are vector fields $\left(v_{t}\right) \subset$ $L^{2}(\mathrm{X})$ such that:
i) the continuity equation 2.3.11) is satisfied for $\left(\mu_{t}, v_{t}\right)$ in the sense of Definition 2.9)
ii) letting $\boldsymbol{\pi} \in \mathscr{P}(C([0,1], \mathrm{X}))$ be the lifting of $\left(\mu_{t}\right)$ it holds

$$
\begin{equation*}
\left|v_{t}\right|\left(\gamma_{t}\right)=\left|\dot{\gamma}_{t}\right| \quad \boldsymbol{\pi} \times \mathcal{L}^{1}-\text { a.e. }(\gamma, t), \tag{4.4.1}
\end{equation*}
$$

iii) $v_{t} \in H^{1,2} \cap D(\operatorname{div})(\mathrm{X})$ for every $t \in(0,1)$ and for every $\varepsilon \in(0,1 / 2)$ it holds

$$
\int_{\varepsilon}^{1-\varepsilon}\left\|\left|\nabla v_{t}\right|_{\mathrm{HS}}\right\|_{L^{2}(\mathrm{X})}^{2} \mathrm{~d} \mathfrak{m}+\sup _{t \in(\varepsilon, 1-\varepsilon)}\left(\left\|v_{t}\right\|_{L^{\infty}(\mathrm{X})}+\left\|\operatorname{div}\left(v_{t}\right)\right\|_{L^{\infty}}\right)<\infty
$$

We briefly comment this statement, given that in [20] it is not presented in this form. The vector fields $v_{t}$ are obtained as gradients of solutions $\eta_{t}$ of a double obstacle problem, the obstacles being given by appropriate 'forward' and 'backward' Kantorovich potentials. This grants that (i) holds. Then (ii) is a general property of 'optimal' lifting of solutions of the continuity equation (see e.g. [15]). The estimates in (iii) are the main gain from [20]: the Laplacian comparison for the squared distance and the Lewy-Stampacchia inequality grant the claimed uniform control on $\operatorname{div}\left(v_{t}\right)=\Delta \eta_{t}$. With a cut-off procedure based on the fact that $\mu_{0}, \mu_{1}$ are assumed to have bounded support, one can show that the $\eta_{t}$ 's can be chosen to also have uniformly bounded support: this and the $L^{\infty}$-bound on the Laplacian implies an $L^{2}$-bound on the Laplacian itself, so that from (2.2.6) we get the $L^{2}$-control on $\left|\nabla v_{t}\right|_{\mathrm{HS}}=\left|\operatorname{Hess}\left(\eta_{t}\right)\right|_{\mathrm{HS}}$. Finally, in [20] it has been proved that the $\eta_{t}$ 's are Lipschitz and, although not explicitly mentioned, keeping track of the various constants involved one can see that it is provided a uniform control on the Lipschitz constant for $t \in(\varepsilon, 1-\varepsilon)$, which in turn implies the desired $L^{\infty}$ control on $\left|v_{t}\right|$.

Thanks to this result we can now prove the following crucial statement:
Proposition 4.6. With the same notations and assumptions as in Section 4.1 the following holds. For every $x, y \in \mathrm{X}$ there exists $\underline{a} \in \mathbb{R}^{N}$ such that

$$
\mathrm{T}(x, \underline{a})=y \quad \text { and } \quad|\underline{a}| \leq \mathrm{d}(x, y) .
$$

proof Fix $y \in \mathrm{X}, R>0$, define

$$
\mu_{0}:=\left.\mathfrak{m}\left(B_{R}(y)\right)^{-1} \mathfrak{m}\right|_{B_{R}(y)} \quad \quad \mu_{1}:=\delta_{y}
$$

and let $\left(\mu_{t}\right)$ be the unique $W_{2}$-geodesic connecting $\mu_{0}$ to $\mu_{1}$ and $\boldsymbol{\pi}$ its lifting. Also, fix $\varepsilon \in(0,1 / 2)$. Then we know from [22] that the $W_{2}$-geodesic $t \mapsto \mu_{t}^{\varepsilon}:=\mu_{\varepsilon+(1-2 \varepsilon) t}$ satisfies the assumptions of Proposition 4.5 and that its lifting $\boldsymbol{\pi}^{\varepsilon}$ is given by

$$
\boldsymbol{\pi}^{\varepsilon}=\left(\operatorname{Restr}_{\varepsilon}^{1-\varepsilon}\right)_{*} \boldsymbol{\pi},
$$

where $\operatorname{Restr}_{t_{0}}^{t_{1}}: C([0,1], \mathrm{X}) \rightarrow C([0,1], \mathrm{X})$ is given by

$$
\operatorname{Restr}_{t_{0}}^{t_{1}}(\gamma)_{t}:=\gamma_{(1-t) t_{0}+t t_{1}} \quad \forall \gamma \in C([0,1], \mathrm{X})
$$

Up to pass to a further restriction, Proposition 4.5 grants the existence of vector fields $\left(v_{t}^{\varepsilon}\right)$ satisfying $(i),(i i),(i i i)$ in the statement. In particular, by (iii) we know that the assumptions of Theorem 2.8 are satisfied so that there exists the Regular Lagrangian Flow $\left(\mathrm{Fl}_{s}^{\left(v_{t}^{\varepsilon}\right)}\right)$ of $\left(v_{t}^{\varepsilon}\right)$.

The representation formula for the solutions of the continuity equation given in Theorem 2.10 gives

$$
\begin{equation*}
\mu_{s}^{\varepsilon}=\left(\mathrm{Fl}_{s}^{\left(v_{t}^{\varepsilon}\right)}\right)_{*} \mu_{0}^{\varepsilon}, \quad \forall s \in[0,1] . \tag{4.4.2}
\end{equation*}
$$

Thus letting $A_{i, t}^{\varepsilon}$ be defined by (4.3.9) for the vector fields $\left(v_{t}^{\varepsilon}\right)$, from (4.3.14) we deduce that

$$
\begin{equation*}
\mathrm{T}\left(x, \underline{A}_{1}^{\varepsilon}(x)\right)=\mathrm{Fl}_{1}^{\left(v_{t}^{\varepsilon}\right)}(x) \in \operatorname{supp}\left(\mu_{1}^{\varepsilon}\right) \subset B_{\varepsilon R}(y) \quad \mu_{0}^{\varepsilon}-\text { a.e. } x . \tag{4.4.3}
\end{equation*}
$$

Now notice that $\boldsymbol{\pi}$ is concentrated on constant speed geodesics of length bounded above by $R$, hence the same holds for $\boldsymbol{\pi}^{\varepsilon}$, so that from (4.4.1) and 4.4.2) we deduce that

$$
\begin{equation*}
\left|v_{s}^{\varepsilon}\right| \circ \mathrm{Fl}_{s}^{\left(v_{t}^{\varepsilon}\right)} \leq R \quad \mu_{0}^{\varepsilon}-\text { a.e.. } \tag{4.4.4}
\end{equation*}
$$

Therefore using the trivial inequality

$$
\left|\underline{A}_{1}^{\varepsilon}\right|^{2}=\sum_{i=1}^{N}\left|\underline{A}_{i, 1}^{\varepsilon}\right|^{2} \leq \sum_{i=1}^{N} \int_{0}^{1}\left|\left\langle v_{s}^{\varepsilon}, X_{i}\right\rangle\right|^{2} \circ \mathrm{Fl}_{s}^{\left(v_{t}^{\varepsilon}\right)} \mathrm{d} s=\int_{0}^{1}\left|v_{s}^{\varepsilon}\right|^{2} \circ \mathrm{Fl}_{s}^{\left(v_{t}^{\varepsilon}\right)} \mathrm{d} s \stackrel{\text { (4.4.4|}}{\leq} R^{2}
$$

valid $\mu_{0}^{\varepsilon}$-a.e. in conjunction with 4.4.3) we deduce that for $\mu_{0}^{\varepsilon}=\mu_{\varepsilon}$-a.e. $x$

$$
\begin{equation*}
\text { there exists } \underline{a} \in \mathbb{R}^{N} \text { with }|\underline{a}| \leq R \text { such that } \mathrm{d}(\mathrm{~T}(x, \underline{a}), y) \leq \varepsilon R \tag{4.4.5}
\end{equation*}
$$

and an argument based on the continuity of T and the compactness of $B_{R}(0) \subset \mathbb{R}^{N}$ yields that the same holds for any $x \in \operatorname{supp}\left(\mu_{\varepsilon}\right)$.

Now notice that simple considerations about the structure of $W_{2}$-geodesics grant that the Hausdorff distance between $\operatorname{supp}\left(\mu_{0}\right)=B_{R}(y)$ and $\operatorname{supp}\left(\mu_{\varepsilon}\right)$ is bounded above by $\varepsilon R$, thus for $x \in B_{R}(y)$ there is a sequence $n \mapsto x_{n} \in \operatorname{supp}\left(\mu_{1 / n}\right)$ converging to $x$. Let $\underline{a}_{n}$ be given by (4.4.5) for $x:=x_{n}$ and $\varepsilon:=\frac{1}{n}$ : by the uniform bound $\left|\underline{a}_{n}\right| \leq R$ and up to pass to a non-relabeled subsequence we can assume that $\underline{a}_{n} \rightarrow \underline{a}$ for some $\underline{a} \in \mathbb{R}^{N}$ with $|\underline{a}| \leq R$. Passing to the limit in

$$
\mathrm{d}\left(\mathrm{~T}\left(x_{n}, \underline{a}_{n}\right), y\right) \leq \frac{R}{n}
$$

using the continuity of T we conclude that $\mathrm{T}(x, \underline{a})=y$. By the arbitrariness of $x \in B_{R}(y)$ and of $R>0$ the proof is completed.

Let us now fix a point $\bar{x} \in \mathrm{X}$ and denote by $\mathbb{G} \subset \mathbb{R}^{N}$ its stabilizer, i.e.

$$
\begin{equation*}
\mathbb{G}:=\left\{\underline{a} \in \mathbb{R}^{N}: \mathrm{T}(\bar{x}, \underline{a})=\bar{x}\right\} . \tag{4.4.6}
\end{equation*}
$$

Notice that the last proposition (and the commutativity of $\mathbb{R}^{N}$ ) grants that the stabilizer does not depend on the choice of the particular point $\bar{x}$; moreover $\mathbb{G}$ is a subgroup of $\mathbb{R}^{N}$ which, by the continuity of T , is closed.

Proposition 4.7. With the same notations and assumptions as in Section 4.1 the following holds. The subgroup $\mathbb{G}$ of $\mathbb{R}^{N}$ defined in 4.4.6 is discrete.
proof We argue by contradiction. If it is not discrete, being closed it must contain a line so that for some $\underline{a}=\left(a_{1}, \ldots, a_{N}\right) \neq 0$ in $\mathbb{R}^{N}$ we have $\underline{\underline{a}} \in \mathbb{G}$ for every $t \in \mathbb{R}$. Put $X:=\sum_{i=1}^{N} a_{i} X_{i}$ and notice that $X$ is not identically 0 and harmonic, so that its Regular Lagrangian Flow $\left(\mathrm{Fl}_{t}^{(X)}\right)$ consists of measure preserving isometries of X such that for $\mathfrak{m}$-a.e. $x$ the curve $t \mapsto \mathrm{Fl}_{t}^{(X)}(x)$ has constant positive speed. In particular, for $\mathfrak{m}$-a.e. $x$ such curve is not constant.

On the other hand, the very definition of T yields

$$
\mathrm{T}(x, t \underline{a})=\mathrm{Fl}_{t}^{(X)}(x) \quad \forall x \in \mathrm{X}, t \in \mathbb{R}
$$

and by assumption the left hand side is equal to $x$ for every $t$ : this gives the desired contradiction and the conclusion.

The quotient space $\mathbb{R}^{N} / \mathbb{G}$ is equipped with the only Riemannian metric letting the quotient map be a Riemannian submersion. The distance induced by this metric is

$$
\begin{equation*}
\left.\left.\mathrm{d}_{\mathbb{R}^{N} / \mathbb{G}}([\underline{a}],[\underline{b}])=\min _{\substack{\left.a^{\prime} ; \mid a^{\prime}\right]=\left[a b \\ \underline{b}^{\prime} ;\right.}} \mid \underline{b}^{\prime}\right]\right][\underline{a}]\left|-\underline{b}^{\prime}\right| . \tag{4.4.7}
\end{equation*}
$$

Also, $\mathbb{R}^{N} / \mathbb{G}$ comes with a canonical, up to multiplication with a positive constant, reference measure $\mathfrak{m}_{\mathbb{R}^{N} / \mathbb{G}}$ : the Haar measure, which also coincides with the volume measure induced by the metric.

Finally, the map $T$ passes to the quotient and induces a map $\tilde{T}: \mathbb{R}^{N} / \mathbb{G} \rightarrow \mathrm{X}$ via the formula:

$$
\tilde{\mathrm{T}}([\underline{a}]):=\mathrm{T}(\bar{x}, \underline{a}) .
$$

With this said, we can now conclude the proof of our main result:
Theorem 4.8. With the same notations and assumptions as in Section 4.1 the following holds.
i) The subgroup $\mathbb{G}$ of $\mathbb{R}^{N}$ defined in (4.4.6) is isomorphic to $\mathbb{Z}^{N}$, so that the quotient space $\mathbb{R}^{N} / \mathbb{G}$ is a flat torus $\mathbb{T}^{N}$.
ii) The induced quotient map $\tilde{\mathrm{T}}: \mathbb{T}^{N} \rightarrow \mathrm{X}$ is an isometry such that $\tilde{\mathrm{T}}_{*} \mathfrak{m}_{\mathbb{T}^{N}}=\mathrm{cm}$ for some $c>0$.
proof
$\tilde{T}$ is an isometry From 4.2 .6 and the definition 4.4.7 we get

$$
\begin{equation*}
\mathrm{d}(\tilde{\mathrm{~T}}([\underline{a}]), \tilde{\mathrm{T}}([\underline{b}])) \leq \mathrm{d}_{\mathbb{R}^{N} / \mathbb{G}}([\underline{a}],[\underline{b}]) \quad \forall[\underline{a}],[\underline{b}] \in \mathbb{R}^{N} / \mathbb{G} \tag{4.4.8}
\end{equation*}
$$

Now let $x, y \in \mathrm{X}$ and apply twice Proposition 4.6 to find $\underline{a} \in \mathbb{R}^{N}$ such that $\mathrm{T}(x, \underline{a})=y$ and $|\underline{a}| \leq \mathrm{d}(x, y)$ and $\underline{b} \in \mathbb{R}^{N}$ such that $\mathrm{T}(\bar{x}, \underline{b})=x$. Then we have

$$
\mathrm{d}_{\mathbb{R}^{N} / \mathbb{G}}([\underline{b}],[\underline{a}+\underline{b}]) \leq|\underline{a}| \leq \mathrm{d}(x, y),
$$

and since by construction and from 4.2 .4 we have $\tilde{T}([\underline{b}])=x$ and $\tilde{T}([\underline{a}+\underline{b}])=y$, this inequality together with 4.4 .8 shows that $\tilde{T}: \mathbb{R}^{N} / \mathbb{G} \rightarrow \mathrm{X}$ is an isometry.
Up to a multiplicative constant, $\tilde{T}$ is measure preserving Being an isometry, $\tilde{T}$ is invertible: denote by $\mathrm{S}: \mathrm{X} \rightarrow \mathbb{R}^{N} / \mathbb{G}$ its inverse and put $\mu:=\mathrm{S}_{*} \mathfrak{m}$. For $\underline{a}=\left(a_{1}, \ldots, a_{N}\right) \in$ $\mathbb{R}^{N}$ let $X_{\underline{a}}:=\sum_{i=1}^{N} a_{i} X_{i}$ and notice that 4.2 .4 reads as $\mathrm{T}(\bar{x}, \underline{b}+\underline{a})=\mathrm{Fl}_{1}^{\left(X_{\underline{a}}\right)}(\mathrm{T}(\bar{x}, \underline{b}))$ for every $\underline{a}, \underline{b} \in \mathbb{R}^{N}$. Passing to the quotient we obtain

$$
\begin{equation*}
\tilde{\mathrm{T}}([\underline{b}]+[\underline{a}])=\mathrm{Fl}_{1}^{\left(X_{\underline{a}}\right)}(\tilde{\mathrm{T}}([\underline{b}])) \quad \forall \underline{a}, \underline{b} \in \mathbb{R}^{N} \tag{4.4.9}
\end{equation*}
$$

hence letting $\tau^{[\underline{a}]}: \mathbb{R}^{N} / \mathbb{G} \rightarrow \mathbb{R}^{N} / \mathbb{G}$ be the translation by $[\underline{a}]$ defined by $\tau^{[\underline{a}]}([\underline{b}]):=[\underline{b}]+[\underline{a}]$ we can rewrite 4.4.9 as

$$
\begin{equation*}
\tau^{[\underline{a}]} \circ \mathrm{S}=\mathrm{S} \circ \mathrm{Fl}_{1}^{\left(X_{\underline{a}}\right)} \quad \forall \underline{a} \in \mathbb{R}^{N} \tag{4.4.10}
\end{equation*}
$$

Therefore we have

$$
\tau_{*}^{[\underline{a}]} \mu=\tau_{*}^{[a]} \mathrm{S}_{*} \mathfrak{m} \stackrel{4.4 .10}{-} \mathrm{S}_{*}\left(\mathrm{Fl}_{1}^{\left(X_{\underline{a}}\right)}\right)_{*} \mathfrak{m} \stackrel{(3.3 .3}{-} \mathrm{S}_{*} \mathfrak{m}=\mu \quad \forall \underline{a} \in \mathbb{R}^{N}
$$

This shows that $\mu$ is translation invariant and thus a multiple of the Haar measure $\mathfrak{m}_{\mathbb{R}^{N}} / \mathbb{G}$. Up to a multiplicative constant, $\tilde{T}$ is measure preserving What we just proved and 4.1.1) ensure that $\mathfrak{m}_{\mathbb{R}^{N} / \mathbb{G}}$ is a finite measure. Now recall that, as it is well known and trivial to prove, discrete subgroups of $\mathbb{R}^{N}$ are isomorphic to $\mathbb{Z}^{n}$ for some $n \leq N$ and that $\mathbb{R}^{N} / \mathbb{Z}^{n}$ has finite volume if and only if $n=N$. Being $\mathbb{G}$ discrete (Proposition 4.7), the thesis follows.

## A Notes on the Hessian on product spaces

In this appendix we continue the investigation of differential operators in product spaces by considering products of RCD spaces and the Hessian of those functions depending only
on one variable. Recall from [6] (see also [5]) that the product of two $\operatorname{RCD}(K, \infty)$ spaces is $\operatorname{RCD}(K, \infty)$ and that the tensorization of the Cheeger energy in the sense of Definition 3.8 holds.

On the other hand, it is not clear whether the density of the product algebra in the sense of 3.9 holds or not, and in any case it seems that the following slightly stronger density property is necessary for the current purposes:

Definition A. 1 (Density of the product algebra - strong form). We say that two metric measure spaces $\left(\mathrm{X}_{1}, \mathrm{~d}_{1}, \mathfrak{m}_{1}\right)$ and $\left(\mathrm{X}_{2}, \mathrm{~d}_{2}, \mathfrak{m}_{2}\right)$ have the property of density of product algebra in the strong form if for $\mathcal{A} \subset W^{1,2}\left(\mathrm{X}_{1} \times \mathrm{X}_{2}\right)$ defined as in (3.2.9) it holds: for $f \in$ $L^{\infty} \cap W^{1,2}\left(\mathrm{X}_{1} \times \mathrm{X}_{2}\right)$ there exists $\left(f_{n}\right) \subset \mathcal{A}$ uniformly bounded and $W^{1,2}$-converging to $f$.

Remark A.2. If $X_{1}$ is infinitesimally Hilbertian and $X_{2}$ the Euclidean space, such strong form of density holds. This is a consequence of the construction done in [18], which grants that for $X_{1}$ arbitrary and $X_{2}=\mathbb{R}$, for any $f \in L^{\infty} \cap W^{1,2}\left(X_{1} \times \mathrm{X}_{2}\right)$ we can find $\left(f_{n}\right) \subset \mathcal{A}$ uniformly bounded and such that $\left(f_{n}\right),\left(\left|\mathrm{d} f_{n}\right|\right)$ converge to $f,|\mathrm{~d} f|$ in $L^{2}$ respectively. The infinitesimal Hilbertianity of $\mathrm{X}_{1}$ and the tensorization of the Cheeger energy (proved in [18]) implies the infinitesimal Hilbertianity of $\mathrm{X}_{1} \times \mathrm{X}_{2}$ and in turn this forces the $W^{1,2_{-}}$ convergence of the functions $\left(f_{n}\right)$ above to $f$.

The case $\mathrm{X}_{2}=\mathbb{R}^{n}$ then comes from an induction argument.
This extra density assumption is needed in the following approximation lemma in order to use the $L^{\infty}$-Lip regularization of the heat flow (see [5]). Such lemma is about approximation of test functions in the product with test functions depending on one variable only and in order to formulate the result it is convenient to introduce the algebra $\tilde{\mathcal{A}}$ as

$$
\tilde{\mathcal{A}}:=\left\{\begin{array}{ll}
\sum_{j=1}^{n} g_{1, j} \circ \pi_{1} g_{2, j} \circ \pi_{2}: n \in \mathbb{N}, & \begin{array}{l}
g_{1, j} \in \operatorname{Test}\left(\mathrm{X}_{1}\right) \text { has bounded support } \\
g_{2, j} \in \operatorname{Test}\left(\mathrm{X}_{2}\right) \text { has bounded support. }
\end{array}
\end{array}\right\}
$$

Notice that the calculus rules obtained in Section 3.2 ensure that $\tilde{\mathcal{A}} \subset \operatorname{Test}\left(\mathrm{X}_{1} \times \mathrm{X}_{2}\right)$.
We then have the following lemma about approximation of test functions with ones in $\tilde{\mathcal{A}}$; notice that a two-steps procedure is needed because the required uniform bound on the differentials prevents arguments by diagonalization.

Lemma A.3. Let $\left(\mathrm{X}_{1}, \mathrm{~d}_{1}, \mathfrak{m}_{1}\right)$ and $\left(\mathrm{X}_{2}, \mathrm{~d}_{2}, \mathfrak{m}_{2}\right)$ be two $\operatorname{RCD}(K, \infty)$ metric measure spaces for which the density of the product algebra holds in the strong form (Definition A.1). Let $f \in \operatorname{Test}\left(\mathrm{X}_{1} \times \mathrm{X}_{2}\right)$ be with bounded support and find $\chi_{1} \in \operatorname{Test}\left(\mathrm{X}_{1}\right), \chi_{2} \in \operatorname{Test}\left(\mathrm{X}_{2}\right)$ with bounded support and such that $\operatorname{supp}(f)$ is contained in the interior of $\left\{\chi_{1}=1\right\} \times\left\{\chi_{2}=1\right\}$ (recall (2.2.4) and for $t>0$ put $\tilde{f}_{t}:=\chi_{1} \circ \pi_{1} \chi_{2} \circ \pi_{2} \mathrm{~h}_{t} f$.

Then:
i) It holds
a) $\tilde{f}_{t} \rightarrow f$ in $W^{2,2}\left(\mathrm{X}_{1} \times \mathrm{X}_{2}\right)$ as $t \downarrow 0$,
b) $\Delta \tilde{f}_{t} \rightarrow \Delta f$ in $L^{2}\left(\mathrm{X}_{1} \times \mathrm{X}_{2}\right)$ as $t \downarrow 0$,
c) $\sup _{t \in(0,1)}\| \| \mathrm{d} \tilde{f}_{t} \mid \|_{L^{\infty}}<\infty$,
d) the sets $\operatorname{supp}\left(\tilde{f}_{t}\right)$ are uniformly bounded for $t \in(0,1)$,
ii) For every $t>0$ there exists a sequence $\left(g_{n}\right) \subset \tilde{\mathcal{A}}$ such that:
a) $g_{n} \rightarrow \tilde{f}_{t}$ in $W^{2,2}\left(\mathrm{X}_{1} \times \mathrm{X}_{2}\right)$ as $n \rightarrow \infty$,
b) $\Delta g_{n} \rightarrow \Delta \tilde{f}_{t}$ in $L^{2}\left(\mathrm{X}_{1} \times \mathrm{X}_{2}\right)$ as $n \rightarrow \infty$,
c) $\sup _{n \in \mathbb{N}}\left\|\left|\mathrm{~d} \tilde{g}_{n}\right|\right\|_{L^{\infty}}<\infty$,
d) the sets $\operatorname{supp}\left(g_{n}\right)$ are uniformly bounded,
proof
(i) It is well known that $\mathrm{h}_{t} f \rightarrow f$ in $W^{1,2}\left(\mathrm{X}_{1} \times \mathrm{X}_{2}\right)$ and $\Delta \mathrm{h}_{t} f \rightarrow \Delta f$ in $L^{2}\left(\mathrm{X}_{1} \times \mathrm{X}_{2}\right)$ as $t \downarrow 0$. From the Leibniz rules for the gradient and the Laplacian and taking into account Proposition 3.7 and Corollary 3.17 we then see that $\tilde{f}_{t} \rightarrow f$ in $W^{1,2}\left(\mathrm{X}_{1} \times \mathrm{X}_{2}\right)$ and $\Delta \tilde{f}_{t} \rightarrow \Delta f$ in $L^{2}\left(\mathrm{X}_{1} \times \mathrm{X}_{2}\right)$ as $t \downarrow 0$. Convergence in $W^{2,2}$ then follows by 2.2 .6 . The uniform bounds on the supports is trivial by construction and the uniform bound on the differential follows by the Bakry-Émery estimate (see Theorem 7.2 in [2]) and the fact that $|\mathrm{d} f| \in L^{\infty}$.
(ii) Fix $t>0$ and use 2.2 .4 to find functions $\tilde{\chi}_{1} \in \operatorname{Test}\left(\mathrm{X}_{1}\right), \tilde{\chi}_{2} \in \operatorname{Test}\left(\mathrm{X}_{2}\right)$ with bounded support such that $\operatorname{supp}\left(f_{t}\right)$ is contained in the interior of $\left\{\tilde{\chi}_{1}=1\right\} \times\left\{\tilde{\chi}_{2}=1\right\}$. Also, let $\left(f_{n}\right) \subset \mathcal{A}$ be uniformly bounded and $W^{1,2}$-converging to $f$ and put

$$
g_{n}:=\left(\chi_{1} \tilde{\chi}_{1}\right) \circ \pi_{1}\left(\chi_{2} \tilde{\chi}_{2}\right) \circ \pi_{2} \mathrm{~h}_{t} f_{n} \quad \forall n \in \mathbb{N}
$$

We claim that the $g_{n}$ 's satisfy the thesis. Indeed, from the regularizing properties of the heat flow we know that $\mathrm{h}_{t} f_{n} \rightarrow \mathrm{~h}_{t} f$ in $W^{1,2}\left(\mathrm{X}_{1} \times \mathrm{X}_{2}\right)$ and $\Delta \mathrm{h}_{t} f_{n} \rightarrow \Delta \mathrm{~h}_{t} f$ in $L^{2}\left(\mathrm{X}_{1} \times \mathrm{X}_{2}\right)$. Since $\left(\chi_{1} \tilde{\chi}_{1}\right) \circ \pi_{1}\left(\chi_{2} \tilde{\chi}_{2}\right) \circ \pi_{2} \mathrm{~h}_{t} f=\tilde{f}_{t}$, the same arguments used in the previous step grant that $g_{n} \rightarrow \tilde{f}_{t}$ in $W^{2,2}\left(\mathrm{X}_{1} \times \mathrm{X}_{2}\right)$ and $\Delta g_{n} \rightarrow \Delta \tilde{f}_{t}$ in $L^{2}\left(\mathrm{X}_{1} \times \mathrm{X}_{2}\right)$. The fact that the supports of the $g_{n}$ 's are uniformly bounded is obvious, and the uniform bound on the differentials follows from the uniform bounds on the $f_{n}$ 's and the $L^{\infty}-$ Lip regularization property (see Theorem 7.3 in [2]).

Thus it remains to show that $g_{n} \in \tilde{\mathcal{A}}$ and since test functions form an algebra, to this aim it is sufficient to show that $\mathrm{h}_{t} f_{n} \in \tilde{\mathcal{A}}$. By the linearity of the heat flow, the fact that $\mathrm{h}_{t} h$ is a test function for $h \in L^{\infty}$ and $t>0$ and performing if necessary a truncation argument on the various addends in $f_{n} \in \mathcal{A}$, to conclude it is sufficient to show that for $h_{1} \in L^{\infty}\left(\mathrm{X}_{1}\right)$ and $h_{2} \in L^{\infty}\left(\mathrm{X}_{2}\right)$ it holds

$$
\begin{equation*}
\mathrm{h}_{t}\left(h_{1} \circ \pi_{1} h_{2} \circ \pi_{2}\right)=h_{1} \circ \pi_{1} \mathrm{~h}_{t}^{\mathrm{X}_{2}}\left(h_{2}\right) \circ \pi_{2}+h_{2} \circ \pi_{2} \mathrm{~h}_{t}^{\mathrm{X}_{1}}\left(h_{1}\right) \circ \pi_{2} \quad \forall t>0 \tag{A.0.11}
\end{equation*}
$$

where $\mathrm{h}_{t}^{\mathrm{X}_{1}}, \mathrm{~h}_{t}^{\mathrm{X}_{2}}$ are the heat flows in $\mathrm{X}_{1}, \mathrm{X}_{2}$ respectively. To this aim notice that Corollary 3.17 grants that for $h_{1}, h_{2}$ in the domain of the Laplacian in the respective spaces it holds

$$
\Delta\left(h_{1} \circ \pi_{1} h_{2} \circ \pi_{2}\right)=h_{1} \circ \pi_{1}\left(\Delta h_{2}\right) \circ \pi_{2}+h_{2} \circ \pi_{2}\left(\Delta h_{1}\right) \circ \pi_{1}
$$

then observe that thanks to this fact the map sending $t \geq 0$ to the right hand side of A.0.11), call it $\tilde{h}_{t}$, is absolutely continuous with values in $L^{2}\left(\mathrm{X}_{1} \times \mathrm{X}_{2}\right)$ and its derivative is given by $\Delta \tilde{h}_{t}$. By the uniqueness of the heat flow we conclude that $\tilde{h}_{t}=\mathrm{h}_{t}\left(\tilde{h}_{0}\right)$, which is our claim.

We then have the following result:
Proposition A.4. Let $\left(\mathrm{X}_{1}, \mathrm{~d}_{1}, \mathfrak{m}_{1}\right)$ and $\left(\mathrm{X}_{2}, \mathrm{~d}_{2}, \mathfrak{m}_{2}\right)$ be two $\operatorname{RCD}(K, \infty)$ spaces for which the density of the product algebra holds in the strong form (Definition A.1) and let $f \in$ $W_{\mathrm{loc}}^{2,2}\left(\mathrm{X}_{1}\right)$. Then $f \circ \pi_{1} \in W_{\mathrm{loc}}^{2,2}\left(\mathrm{X}_{1} \times \mathrm{X}_{2}\right)$ and

$$
\begin{equation*}
\operatorname{Hess}\left(f \circ \pi_{1}\right)(\nabla g, \nabla \tilde{g})\left(x_{1}, x_{2}\right)=\operatorname{Hess}(f)\left(\nabla g_{x_{2}}, \nabla \tilde{g}_{x_{2}}\right)\left(x_{1}\right), \quad \mathfrak{m}_{1} \times \mathfrak{m}_{2}-\text { a.e. }\left(x_{1}, x_{2}\right) \tag{A.0.12}
\end{equation*}
$$

for every $g, \tilde{g} \in \operatorname{Test}\left(\mathrm{X}_{1} \times \mathrm{X}_{2}\right)$.
proof It is readily verified that the map sending $g, \tilde{g} \in \operatorname{Test}\left(\mathrm{X}_{1} \times \mathrm{X}_{2}\right)$ to the right hand side of A.0.12 defines an element of $L_{\text {loc }}^{2}\left(\left(T^{*}\right)^{\otimes 2}\left(\mathrm{X}_{1} \times \mathrm{X}_{2}\right)\right)$, hence to conclude it is sufficient to show that for such element the identity (2.2.5) holds.

Now consider the identity (2.2.5) defining the Hessian for functions in $W_{\text {loc }}^{2,2}$ with $g:=g_{n}$, where $\left(g_{n}\right)$ is a sequence of test functions $W^{2,2}$-converging to some limit $g$, such that $\Delta g_{n} \rightarrow \Delta g$ in $L^{2}$ and with $\operatorname{supp}\left(g_{n}\right)$ and $\left\|\left|\mathrm{d} g_{n}\right|\right\|_{L^{\infty}}$ uniformly bounded: it is readily verified that in this case the two sides of (2.2.5) pass to the limit.

Thus by Lemma A.3 above and the bilinearity and symmetry in $g, \tilde{g}$, to conclude it is sufficient to consider $g=\tilde{g}$ of the form $g=g_{1} \circ \pi_{1} g_{2} \circ \pi_{2}$ for $g_{1} \in \operatorname{Test}\left(\mathrm{X}_{1}\right)$ and $g_{2} \in \operatorname{Test}\left(\mathrm{X}_{2}\right)$ both with bounded support. For such $g$ we have $g_{x_{2}}=g_{2}\left(x_{2}\right) g_{1}$, and thus $\nabla g_{x_{2}}=g\left(x_{2}\right) \nabla g_{1}$, so that our aim is to show that for any $h \in \operatorname{Test}\left(\mathrm{X}_{1} \times \mathrm{X}_{2}\right)$ with bounded support it holds

$$
\begin{align*}
-\int\left\langle\nabla\left(f \circ \pi_{1}\right), \nabla g\right\rangle \operatorname{div} & (h \nabla g)+h\left\langle\nabla(f \circ \pi), \nabla \frac{|\nabla g|^{2}}{2}\right\rangle \mathrm{d}\left(\mathfrak{m}_{1} \times \mathfrak{m}_{2}\right)  \tag{A.0.13}\\
& =\int h g_{2}^{2} \circ \pi_{2} \operatorname{Hess}(f)\left(\nabla g_{1}, \nabla g_{1}\right) \circ \pi_{1} \mathrm{~d}\left(\mathfrak{m}_{1} \times \mathfrak{m}_{2}\right) .
\end{align*}
$$

Denoting for clarity $\operatorname{div}_{1}, \operatorname{div}_{2}$ the divergence operators in $\mathrm{X}_{1}, \mathrm{X}_{2}$ respectively and using formulas (3.2.8), (3.2.16) and (3.2.18), for $\mathfrak{m}_{1} \times \mathfrak{m}_{2}$-a.e. $\left(x_{1}, x_{2}\right)$ we have

$$
\begin{align*}
& \left(\left\langle\nabla\left(f \circ \pi_{1}\right), \nabla g\right\rangle \operatorname{div}(h \nabla g)\right)\left(x_{1}, x_{2}\right) \\
& \quad=g_{2}\left(x_{2}\right)\left\langle\nabla f, \nabla g_{1}\right\rangle\left(x_{1}\right)\left(g_{2}\left(x_{2}\right) \operatorname{div}_{1}\left(h_{x_{2}} \nabla g_{1}\right)\left(x_{1}\right)+g_{1}\left(x_{1}\right) \operatorname{div}_{2}\left(h^{x_{1}} \nabla g_{2}\right)\left(x_{2}\right)\right) . \tag{A.0.14}
\end{align*}
$$

From (3.2.7) we have

$$
|\nabla g|^{2}=g_{2}^{2} \circ \pi_{2}\left|\nabla g_{1}\right|^{2} \circ \pi_{1}+g_{1}^{2} \circ \pi_{1}\left|\nabla g_{2}\right|^{2} \circ \pi_{2}
$$

and thus recalling (3.2.16 we obtain

$$
h\left\langle\nabla\left(f \circ \pi_{1}\right), \nabla \frac{|\nabla g|^{2}}{2}\right\rangle=h g_{2}^{2} \circ \pi_{2}\left\langle\nabla f, \nabla \frac{\left|\nabla g_{1}\right|^{2}}{2}\right\rangle \circ \pi_{1}+h\left|\nabla g_{2}\right|^{2} \circ \pi_{2}\left(g_{1}\left\langle\nabla f, \nabla g_{1}\right\rangle\right) \circ \pi_{1} .
$$

Adding up this identity and A.0.14 and integrating, the conclusion A.0.13 follows by the defining property (2.2.5) of $\operatorname{Hess}(f)$ and the trivial identity

$$
\begin{aligned}
\int g_{1}\left(x_{1}\right) g_{2}\left(x_{2}\right)\left\langle\nabla f, \nabla g_{1}\right\rangle\left(x_{1}\right) & \operatorname{div}_{2}\left(h^{x_{1}} \nabla g_{2}\right)\left(x_{2}\right) \mathrm{d}_{1}\left(x_{1}\right) \mathrm{dm}_{2}\left(x_{2}\right) \\
& =\int g_{1}\left(x_{1}\right)\left\langle\nabla f, \nabla g_{1}\right\rangle\left(x_{1}\right) \int g_{2} \operatorname{div}_{2}\left(h^{x_{1}} \nabla g_{2}\right) \operatorname{d\mathfrak {m}_{2}\operatorname {dm}(x_{1})} \\
& =-\int g_{1}\left(x_{1}\right)\left\langle\nabla f, \nabla g_{1}\right\rangle\left(x_{1}\right) \int h^{x_{1}}\left|\nabla g_{2}\right|^{2} \operatorname{dm_{2}} \operatorname{d\mathfrak {m}(x_{1})} \\
& =-\int h\left|\nabla g_{2}\right|^{2} \circ \pi_{2}\left(g_{1}\left\langle\nabla f, \nabla g_{1}\right\rangle\right) \circ \pi_{1} \mathrm{~d}\left(\mathfrak{m}_{1} \times \mathfrak{m}_{2}\right) .
\end{aligned}
$$

## References

[1] L. Ambrosio, Transport equation and Cauchy problem for BV vector fields, Invent. Math., 158 (2004), pp. 227-260.
[2] L. Ambrosio, N. Gigli, A. Mondino, and T. Rajala, Riemannian Ricci curvature lower bounds in metric measure spaces with $\sigma$-finite measure, Trans. Amer. Math. Soc., 367 (2012), pp. 4661-4701.
[3] L. Ambrosio, N. Gigli, and G. Savaré, Density of Lipschitz functions and equivalence of weak gradients in metric measure spaces, Rev. Mat. Iberoam., 29 (2013), pp. 969-996.
[4] ——, Calculus and heat flow in metric measure spaces and applications to spaces with Ricci bounds from below, Invent. Math., 195 (2014), pp. 289-391.
[5] __, Metric measure spaces with Riemannian Ricci curvature bounded from below, Duke Math. J., 163 (2014), pp. 1405-1490.
[6] __, Bakry-Émery curvature-dimension condition and Riemannian Ricci curvature bounds, The Annals of Probability, 43 (2015), pp. 339-404.
[7] L. Ambrosio, A. Mondino, and G. Savaré, On the Bakry-Émery condition, the gradient estimates and the Local-to-Global property of $R C D^{*}(K, N)$ metric measure spaces, The Journal of Geometric Analysis, 26 (2014), pp. 1-33.
[8] L. Ambrosio and D. Trevisan, Well posedness of Lagrangian flows and continuity equations in metric measure spaces, Anal. PDE, 7 (2014), pp. 1179-1234.
[9] __, Lecture notes on the DiPerna-Lions theory in abstract measure spaces. Accepted at Annales Fac. Sc. de Toulouse, arXiv:1505.05292, 2015.
[10] J. Cheeger, Differentiability of Lipschitz functions on metric measure spaces, Geom. Funct. Anal., 9 (1999), pp. 428-517.
[11] T. H. Colding, Ricci curvature and volume convergence, Ann. of Math. (2), 145 (1997), pp. 477-501.
[12] R. J. DiPerna and P.-L. Lions, Ordinary differential equations, transport theory and Sobolev spaces, Invent. Math., 98 (1989), pp. 511-547.
[13] N. Gigli, Lecture notes on differential calculus on RCD spaces. Preprint, arXiv: 1703.06829 .
[14] __, The splitting theorem in non-smooth context. Preprint, arXiv:1302.5555, 2013.
[15] ——, Nonsmooth differential geometry - an approach tailored for spaces with Ricci curvature bounded from below. Accepted at Mem. Amer. Math. Soc., arXiv:1407.0809, 2014.
[16] _—, An overview of the proof of the splitting theorem in spaces with non-negative Ricci curvature, Analysis and Geometry in Metric Spaces, 2 (2014), pp. 169-213.
[17] _—, On the differential structure of metric measure spaces and applications, Mem. Amer. Math. Soc., 236 (2015), pp. vi+91.
[18] N. Gigli and B. Han, Sobolev spaces on warped products. Preprint, arXiv:1512.03177, 2015.
[19] ——, The continuity equation on metric measure spaces, Calc. Var. Partial Differential Equations, 53 (2013), pp. 149-177.
[20] N. Gigli and S. Mosconi, The abstract Lewy-Stampacchia inequality and applications, J. Math. Pures Appl. (9), 104 (2014), pp. 258-275.
[21] N. Gigli and E. Pasqualetto, Equivalence of two different notions of tangent bundle on rectifiable metric measure spaces. Preprint, arXiv:1611.09645, 2016.
[22] N. Gigli, T. Rajala, and K.-T. Sturm, Optimal Maps and Exponentiation on Finite-Dimensional Spaces with Ricci Curvature Bounded from Below, J. Geom. Anal., 26 (2016), pp. 2914-2929.
[23] B. Han, Ricci tensor on $\operatorname{RCD}^{*}(K, N)$ spaces. Preprint, arXiv: 1412.0441.
[24] S. Honda, Elliptic PDEs on compact Ricci limit spaces and applications. Preprint, arXiv:1410.3296.
[25] A. Mondino and G. Wei, On the universal cover and the fundamental group of an $\operatorname{RCD}(K, N)$-space. Accepted at Crelle's journal, arXiv: 1605.02854.
[26] P. Petersen, Riemannian geometry, vol. 171 of Graduate Texts in Mathematics, Springer, Cham, third ed., 2016.
[27] G. Savaré, Self-improvement of the Bakry-Émery condition and Wasserstein contraction of the heat flow in $\operatorname{RCD}(K, \infty)$ metric measure spaces, Discrete Contin. Dyn. Syst., 34 (2014), pp. 1641-1661.
[28] N. Shanmugalingam, Newtonian spaces: an extension of Sobolev spaces to metric measure spaces, Rev. Mat. Iberoamericana, 16 (2000), pp. 243-279.


[^0]:    *SISSA, Trieste. email: ngigli@sissa.it
    ${ }^{\dagger}$ SISSA, Trieste. email: crigoni@sissa.it

