# Transport equation and Cauchy problem for non-smooth vector fields 

Luigi Ambrosio

Scuola Normale Superiore Piazza dei Cavalieri 7, 56126 Pisa, Italy l.ambrosio@sns.it

Cetraro, June 27-July 2, 2005

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## 1. Introduction

In these lectures we study the well-posedness of the Cauchy problem for the homogeneous conservative continuity equation
(PDE)

$$
\frac{d}{d t} \mu_{t}+D_{x} \cdot\left(\boldsymbol{b} \mu_{t}\right)=0 \quad(t, x) \in I \times \mathbb{R}^{d}
$$

and for the transport equation

$$
\frac{d}{d t} w_{t}+\boldsymbol{b} \cdot \nabla w_{t}=c_{t} .
$$

Here $\boldsymbol{b}(t, x)=\boldsymbol{b}_{t}(x)$ is a given time-dependent vector field in $\mathbb{R}^{d}$ : we are interested to the case when $\boldsymbol{b}_{t}(\cdot)$ is not necessarily Lipschitz and has, for instance, a Sobolev or BV regularity. Vector fields with this "low" regularity show up, for instance, in several PDE's describing the motion of fluids, and in the theory of conservation laws.

We are also particularly interested to the well posedness of the system of ordinary differential equations
(ODE)

$$
\left\{\begin{array}{l}
\dot{\gamma}(t)=\boldsymbol{b}_{t}(\gamma(t)) \\
\gamma(0)=x
\end{array}\right.
$$

In some situations one might hope for a "generic" uniqueness of the solutions of ODE, i.e. for "almost every" initial datum $x$.IAn even weaker requirement is the research of a "selection principle", i.e. a strategy to select for $\mathscr{L}^{d}$-almost every $x$ a solution $\boldsymbol{X}(\cdot, x)$ in such a way that this selection is stable w.r.t. smooth approximations of $\boldsymbol{b}$.ll

In other words, we would like to know that, whenever we approximate $\boldsymbol{b}$ by smooth vector fields $\boldsymbol{b}^{h}$, the classical trajectories $\boldsymbol{X}^{h}$ associated to $\boldsymbol{b}^{h}$ satisfy

$$
\lim _{h \rightarrow \infty} \boldsymbol{X}^{h}(\cdot, x)=\boldsymbol{X}(\cdot, x) \quad \text { in } C\left([0, T] ; \mathbb{R}^{d}\right) \text {, for } \mathscr{L}^{d} \text {-a.e. } x .
$$

We will see that there is a close link between the two problems, first investigated in a nonsmooth setting by Di Perna and Lions in [53].

Let us now make some basic technical remarks on the continuity equation and the transport equation:

Remark 1. [Regularity in space of $\boldsymbol{b}_{t}$ and $\mu_{t}$ ] (1) Since the continuity equation (PDE) is in divergence form, it makes sense without any regularity requirement on $\boldsymbol{b}_{t}$ and/or $\mu_{t}$, provided

$$
\begin{equation*}
\int_{I} \int_{A}\left|\boldsymbol{b}_{t}\right| d\left|\mu_{t}\right| d t<+\infty \quad \forall A \subset \subset \mathbb{R}^{d} \tag{1}
\end{equation*}
$$

However, when we consider possibly singular measures $\mu_{t}$, we must take care of the fact that the product $\boldsymbol{b}_{t} \mu_{t}$ is sensitive to modifications of $\boldsymbol{b}_{t}$ in $\mathscr{L}^{d}$-negligible sets. In the Sobolev or $B V$ case we will consider only measures $\mu_{t}=w_{t} \mathscr{L}^{d}$, so everything is well posed. \|
(2) On the other hand, due to the fact that the distribution $\boldsymbol{b}_{t} \cdot \nabla w$ is defined by

$$
\left\langle\boldsymbol{b}_{t} \cdot \nabla w, \varphi\right\rangle:=-\int_{I} \int w\left\langle\boldsymbol{b}_{t}, \nabla \varphi\right\rangle d x d t-\int_{I}\left\langle D_{x} \cdot \boldsymbol{b}_{t}, w_{t} \varphi_{t}\right\rangle d t \quad \varphi \in C_{c}^{\infty}\left(I \times \mathbb{R}^{d}\right)
$$

(a definition consistent with the case when $w_{t}$ is smooth) the transport equation makes sense only if we assume that $D_{x} \cdot \boldsymbol{b}_{t}=\operatorname{div} \boldsymbol{b}_{t} \mathscr{L}^{d}$ for
$\mathscr{L}^{1}$-a.e. $t \in I$. See also [28], [31] for recent results on the transport equation when $\boldsymbol{b}$ satisfies a one-sided Lipschitz condition.

Next, we consider the problem of the time continuity of $t \mapsto \mu_{t}$ and $t \mapsto w_{t}$.

Remark 2. [Regularity in time of $\mu_{t}$ ] For any test function $\varphi \in$ $C_{c}^{\infty}\left(\mathbb{R}^{d}\right)$, condition (1) gives

$$
\frac{d}{d t} \int_{\mathbb{R}^{d}} \varphi d \mu_{t}=\int_{\mathbb{R}^{d}} \boldsymbol{b}_{t} \cdot \nabla \varphi d \mu_{t} \in L^{1}(I)
$$

and therefore the map $t \mapsto\left\langle\mu_{t}, \varphi\right\rangle$, for given $\varphi$, has a unique uniformly continuous representative in $I . \|$ By a simple density argument we can find a unique representative $\tilde{\mu}_{t}$ independent of $\varphi$, such that $t \mapsto\left\langle\tilde{\mu}_{t}, \varphi\right\rangle$ is uniformly continuous in $I$ for any $\varphi \in C_{c}^{\infty}\left(\mathbb{R}^{d}\right)$. We will always work with this representative, so that $\mu_{t}$ will be well defined for all $t$ and even at the endpoints of $I$. \|

An analogous remark applies for solutions of the transport equation. I

There are some other important links between the two equations:
(1) The transport equation reduces to the continuity equation in the case when $c_{t}=-w_{t} \operatorname{div} \boldsymbol{b}_{t}$;
(2) Formally, one can estabilish a duality between the two equations via the (formal) identity

$$
\begin{aligned}
\frac{d}{d t} \int w_{t} d \mu_{t} & =\int \frac{d}{d t} w_{t} d \mu_{t}+\int \frac{d}{d t} \mu_{t} w_{t} \\
& =\int\left(-\boldsymbol{b}_{t} \cdot \nabla w_{t}+c\right) d \mu_{t}+\int \boldsymbol{b}_{t} \cdot \nabla w_{t} d \mu_{t}=\int c d \mu_{t}
\end{aligned}
$$

This duality method is a classical tool to prove uniqueness in a sufficiently smooth setting (but see also [28], [31]). II
(3) Finally, if we denote by $\boldsymbol{Y}(t, s, x)$ the solution of the ODE at time $t$, starting from $x$ at the initial times $s$, i.e.

$$
\frac{d}{d t} \boldsymbol{Y}(t, s, x)=\boldsymbol{b}_{t}(\boldsymbol{Y}(t, s, x)), \quad \boldsymbol{Y}(s, s, x)=x
$$

then $\boldsymbol{Y}(t, \cdot, \cdot)$ are themselves solutions of the transport equation: to see this, it suffices to differentiate the semigroup identity

$$
\boldsymbol{Y}(t, s, \boldsymbol{Y}(s, l, x))=\boldsymbol{Y}(t, l, x)
$$

w.r.t. $s$ to obtain, after the change of variables $y=\boldsymbol{Y}(s, l, x)$, the equation

$$
\frac{d}{d s} \boldsymbol{Y}(t, s, y)+\boldsymbol{b}_{s}(y) \cdot \nabla \boldsymbol{Y}(t, s, y)=0
$$

This property is used in a essential way in [53] to characterize the flow $\boldsymbol{Y}$ and to prove its stability properties.lThe approach developed here, based on [7], is based on a careful analysis of the measures transported by the flow, and ultimately on the homogeneous continuity equation only.

## 2. Transport equation and continuity equation within the Cauchy-Lipschitz framework

In this section we recall the classical representation formulas for solutions of the continuity or transport equation in the case when

$$
\boldsymbol{b} \in L^{1}\left([0, T] ; W^{1, \infty}\left(\mathbb{R}^{d} ; \mathbb{R}^{d}\right)\right)
$$

Under this assumption it is well known that solutions $\boldsymbol{X}(t, \cdot)$ of the ODE are unique and stable. A quantitative information can be obtained by differentiation:

$$
\begin{aligned}
\frac{d}{d t}|\boldsymbol{X}(t, x)-\boldsymbol{X}(t, y)|^{2} & =2\left\langle\boldsymbol{b}_{t}(\boldsymbol{X}(t, x))-\boldsymbol{b}_{t}(\boldsymbol{X}(t, y)), \boldsymbol{X}(t, x)-\boldsymbol{X}(t, y)\right\rangle \\
& \leq 2 \operatorname{Lip}\left(\boldsymbol{b}_{t}\right)|\boldsymbol{X}(t, x)-\boldsymbol{X}(t, y)|^{2}
\end{aligned}
$$

(here Lip $(f)$ denotes the least Lipschitz constant of $f$ ), so that Gronwall lemma immediately gives

$$
\begin{equation*}
\operatorname{Lip}(\boldsymbol{X}(t, \cdot)) \leq \exp \left(\int_{0}^{t} \operatorname{Lip}\left(\boldsymbol{b}_{s}\right) d s\right) \tag{2}
\end{equation*}
$$

Turning to the continuity equation, uniqueness of measure-valued solutions can be proved by the duality method.I Or, following the techniques developed in these lectures, it can be proved in a more general setting for positive measure-valued solutions (via the superposition principle) and for signed solutions $\mu_{t}=w_{t} \mathscr{L}^{d}$ (via the theory of renormalized solutions). So in this section we focus only on the existence and the representation issues.ll

The representation formula is indeed very simple:
Proposition 3. For any initial datum $\bar{\mu}$ the solution of the continuity equation is given by

$$
\begin{equation*}
\mu_{t}:=\boldsymbol{X}(t, \cdot)_{\#} \bar{\mu}, \quad \text { i.e. } \quad \int_{\mathbb{R}^{d}} \varphi d \mu_{t}=\int_{\mathbb{R}^{d}} \varphi(\boldsymbol{X}(t, x)) d \bar{\mu}(x) . \tag{3}
\end{equation*}
$$

Proof. Notice first that we need only to check the distributional identity $\frac{d}{d t} \mu_{t}+D_{x} \cdot\left(\boldsymbol{b}_{t} \mu_{t}\right)=0$ on test functions of the form $\psi(t) \varphi(x)$, so that

$$
\int_{\mathbb{R}} \psi^{\prime}(t)\left\langle\mu_{t}, \varphi\right\rangle d t+\int_{\mathbb{R}} \psi(t) \int_{\mathbb{R}^{d}}\left\langle\boldsymbol{b}_{t}, \nabla \varphi\right\rangle d \mu_{t} d t=0 .
$$

This means that we have to check that $t \mapsto\left\langle\mu_{t}, \varphi\right\rangle$ belongs to $W^{1,1}(0, T)$ for any $\varphi \in C_{c}^{\infty}\left(\mathbb{R}^{d}\right)$ and that its distributional derivative is $\int_{\mathbb{R}^{d}}\left\langle\boldsymbol{b}_{t}, \nabla \varphi\right\rangle d \mu_{t}$.
We show first that this map is absolutely continuous, and in particular $W^{1,1}(0, T)$; then one needs only to compute the pointwise derivative. For every choice of finitely many, say $n$, pairwise disjoint intervals $\left(a_{i}, b_{i}\right) \subset[0, T]$ we have

$$
\begin{aligned}
\sum_{i=1}^{n}\left|\varphi\left(\boldsymbol{X}\left(b_{i}, x\right)\right)-\varphi\left(\boldsymbol{X}\left(a_{i}, x\right)\right)\right| & \leq\|\nabla \varphi\|_{\infty} \int_{\cup_{i}\left(a_{i}, b_{i}\right)}|\dot{\boldsymbol{X}}(t, x)| d t \\
& \leq\|\nabla \varphi\|_{\infty} \int_{\cup_{i}\left(a_{i}, b_{i}\right)} \sup \left|\boldsymbol{b}_{t}\right| d t
\end{aligned}
$$

and therefore an integration with respect to $\bar{\mu}$ gives

$$
\sum_{i=1}^{n}\left|\left\langle\mu_{b_{i}}-\mu_{a_{i}}, \varphi\right\rangle\right| \leq\|\nabla \varphi\|_{\infty} \int_{\cup_{i}\left(a_{i}, b_{i}\right)} \sup \left|\boldsymbol{b}_{t}\right| d t
$$

The absolute continuity of the integral shows that the right hand side can be made small when $\sum_{i}\left(b_{i}-a_{i}\right)$ is small. This proves the absolute continuity. For any $x$ the identity $\dot{\boldsymbol{X}}(t, x)=\boldsymbol{b}_{t}(\boldsymbol{X}(t, x))$ is fulfilled for $\mathscr{L}^{1}$-a.e. $t \in[0, T]$. Then, by Fubini's theorem, we know also that for $\mathscr{L}^{1}$-a.e. $t \in[0, T]$ the previous identity holds for $\bar{\mu}$-a.e. $x$, and therefore

$$
\begin{aligned}
\frac{d}{d t}\left\langle\mu_{t}, \varphi\right\rangle & =\frac{d}{d t} \int_{\mathbb{R}^{d}} \varphi(\boldsymbol{X}(t, x)) d \bar{\mu}(x) \\
& =\int_{\mathbb{R}^{d}}\left\langle\nabla \varphi(\boldsymbol{X}(t, x)), \boldsymbol{b}_{t}(\boldsymbol{X}(t, x))\right\rangle d \bar{\mu}(x) \\
& =\left\langle\boldsymbol{b}_{t} \mu_{t}, \nabla \varphi\right\rangle
\end{aligned}
$$

for $\mathscr{L}^{1}$-a.e. in $[0, T]$.

In the case when $\bar{\mu}=\rho \mathscr{L}^{d}$ we can say something more, proving that the measures $\mu_{t}=\boldsymbol{X}(t, \cdot)_{\#} \bar{\mu}$ are absolutely continuous w.r.t. $\mathscr{L}^{d}$ and computing explicitely their density. Let us start by recalling the classical area formula: if $f: \mathbb{R}^{d} \rightarrow \mathbb{R}^{d}$ is a (locally) Lipschitz map, then

$$
\int_{A} g|J f| d x=\int_{\mathbb{R}^{d}} \sum_{x \in A \cap f^{-1}(y)} g(x) d y
$$

for any Borel set $A \subset \mathbb{R}^{d}$, where $J f=\operatorname{det} \nabla f$ (recall that, by Rademacher theorem, Lipschitz functions are differentiable $\mathscr{L}^{d}$-a.e.).\|Assuming in addition that $f$ is 1-1 and onto and that $|J f|>0 \mathscr{L}^{d}$-a.e. on $A$ we can set $A=f^{-1}(B)$ and $g=\rho /|J f|$ to obtain

$$
\int_{f^{-1}(B)} \rho d x=\int_{B} \frac{\rho}{|J f|} \circ f^{-1} d y
$$

In other words, we have got a formula for the push-forward:

$$
\begin{equation*}
f_{\#}\left(\rho \mathscr{L}^{d}\right)=\frac{\rho}{|J f|} \circ f^{-1} \mathscr{L}^{d} . \tag{4}
\end{equation*}
$$

In our case $f(x)=\boldsymbol{X}(t, x)$ is surely 1-1, onto and Lipschitz. It remains to show that $|J \boldsymbol{X}(t, \cdot)|$ does not vanish: in fact, one can show that $J \boldsymbol{X}>0$ and

$$
\begin{equation*}
\exp \left[-\int_{0}^{t}\left\|\left[\operatorname{div} \boldsymbol{b}_{s}\right]^{-}\right\|_{\infty} d s\right] \leq J \boldsymbol{X}(t, x) \leq \exp \left[\int_{0}^{t}\left\|\left[\operatorname{div} \boldsymbol{b}_{s}\right]^{+}\right\|_{\infty} d s\right] \tag{5}
\end{equation*}
$$

for $\mathscr{L}^{d}$-a.e. $x$, thanks to the following fact, whose proof is left as an exercise.

Exercise 4. If $\boldsymbol{b}$ is smooth, we have

$$
\frac{d}{d t} J \boldsymbol{X}(t, x)=\operatorname{div} \boldsymbol{b}_{t}(\boldsymbol{X}(t, x)) J \boldsymbol{X}(t, x) .
$$

Hint: use the ODE $\frac{d}{d t} \nabla \boldsymbol{X}=\nabla \boldsymbol{b}_{t}(\boldsymbol{X}) \nabla \boldsymbol{X}$.
I
The previous exercise gives that, in the smooth case, $J \boldsymbol{X}(\cdot, x)$ solves a linear ODE with the initial condition $J \boldsymbol{X}(0, x)=1$, whence the estimates on $J \boldsymbol{X}$ follow.IIn the general case the upper estimate on $J \boldsymbol{X}$ still holds by a smoothing argument, thanks to the lower semicontinuity of

$$
\Phi(v):= \begin{cases}\|J v\|_{\infty} & \text { if } J v \geq 0 \mathscr{L}^{d} \text {-a.e. } \\ +\infty & \text { otherwise }\end{cases}
$$

with respect to the $w^{*}$-topology of $W^{1, \infty}\left(\mathbb{R}^{d} ; \mathbb{R}^{d}\right)$. \|This is indeed the supremum of the family of $\Phi_{p}^{1 / p}$, where $\Phi_{p}$ are the polyconvex (and therefore lower semicontinuous) functionals

$$
\Phi_{p}(v):=\int_{B_{p}}|\chi(J v)|^{p} d x
$$

Here $\chi(t)$, equal to $\infty$ on $(-\infty, 0)$ and equal to $t$ on $[0,+\infty)$, is I.s.c. and convex.IThe lower estimate can be obtained by applying the upper one in a time reversed situation.
I
Now we turn to the representation of solutions of the transport equation:
Proposition 5. If $w \in L_{\text {loc }}^{1}\left([0, T] \times \mathbb{R}^{d}\right)$ solves

$$
\frac{d}{d t} w_{t}+\boldsymbol{b} \cdot \nabla w=c \in L_{\mathrm{loc}}^{1}\left([0, T] \times \mathbb{R}^{d}\right)
$$

then, for $\mathscr{L}^{d}$-a.e. $x$, we have

$$
w_{t}(\boldsymbol{X}(t, x))=w_{0}(x)+\int_{0}^{t} c_{s}(\boldsymbol{X}(s, x)) d s \quad \forall t \in[0, T] .
$$

The (formal) proof is based on the simple observation that

$$
\begin{aligned}
\frac{d}{d t} w_{t} \circ \boldsymbol{X}(t, x) & =\frac{d}{d t} w_{t}(\boldsymbol{X}(t, x))+\frac{d}{d t} \boldsymbol{X}(t, x) \cdot \nabla w_{t}(\boldsymbol{X}(t, x)) \\
& =\frac{d}{d t} w_{t}(\boldsymbol{X}(t, x))+\boldsymbol{b}_{t}(\boldsymbol{X}(t, x)) \cdot \nabla w_{t}(\boldsymbol{X}(t, x)) \\
& =c_{t}(\boldsymbol{X}(t, x))
\end{aligned}
$$

In particular, as $\boldsymbol{X}(t, x)=\boldsymbol{Y}(t, 0, x)=[\boldsymbol{Y}(0, t, \cdot)]^{-1}(x)$, we get

$$
w_{t}(y)=w_{0}(\boldsymbol{Y}(0, t, y))+\int_{0}^{t} c_{s}(\boldsymbol{Y}(s, t, y)) d s
$$

We conclude this presentation of the classical theory pointing out two simple local variants of the assumption $\boldsymbol{b} \in L^{1}\left([0, T] ; W^{1, \infty}\left(\mathbb{R}^{d} ; \mathbb{R}^{d}\right)\right)$ made throughout this section.

Remark 6. [First local variant] The theory outlined above still works under the assumptions

$$
\boldsymbol{b} \in L^{1}\left([0, T] ; W_{\mathrm{loc}}^{1, \infty}\left(\mathbb{R}^{d} ; \mathbb{R}^{d}\right)\right), \quad \frac{|\boldsymbol{b}|}{1+|x|} \in L^{1}\left([0, T] ; L^{\infty}\left(\mathbb{R}^{d}\right)\right)
$$

Indeed, due to the growth condition on $\boldsymbol{b}$, we still have pointwise uniqueness of the ODE and a uniform local control on the growth of $|\boldsymbol{X}(t, x)|$, therefore we need only to consider a local Lipschitz condition w.r.t. $x$, integrable w.r.t. $t$.

The next variant will be used in the proof of the superposition principle. Remark 7. [Second local variant] Still keeping the $L^{1}\left(W_{\text {loc }}^{1, \infty}\right)$ assumption, and assuming $\mu_{t} \geq 0$, the second growth condition on $|\boldsymbol{b}|$ can be replaced by a global, but more intrinsic, condition:

$$
\begin{equation*}
\int_{0}^{T} \int_{\mathbb{R}^{d}} \frac{\left|\boldsymbol{b}_{t}\right|}{1+|x|} d \mu_{t} d t<+\infty \tag{6}
\end{equation*}
$$

Under this assumption one can show that for $\bar{\mu}$-a.e. $x$ the maximal solution $\boldsymbol{X}(\cdot, x)$ of the ODE starting from $x$ is defined up to $t=T$ and still the representation $\mu_{t}=\boldsymbol{X}(t, \cdot)_{\#} \bar{\mu}$ holds for $t \in[0, T]$.

## 3. ODE uniqueness versus PDE uniquenessı

In this section we illustrate some quite general principles, whose application may depend on specific assumptions on $\boldsymbol{b}$, relating the uniqueness of the ODE to the uniqueness of the PDE. \|The viewpoint adopted in this section is very close in spirit to Young's theory [85] of generalized surfaces and controls (a theory with remarkable applications also non-linear PDE's [52, 78] and Calculus of Variations [19]) and has also some connection with Brenier's weak solutions of incompressible Euler equations [24], with Kantorovich's viewpoint in the theory of optimal transportation [57, 76] and with Mather's theory [71, 72, 18]: in order to study existence, uniqueness and stability with respect to perturbations of the data of solutions to the ODE, we consider suitable measures in the space of continuous maps, allowing for superposition of trajectories.IThen, in some special situations we are able to show that this superposition actually does not occur, but still this "probabilistic" interpretation is very useful to understand the underlying techniques and to give an intrinsic characterization of the flow. I

The first very general criterion is the following.
Theorem 8. Let $A \subset \mathbb{R}^{d}$ be a Borel set. The following two properties are equivalent:
(a) Solutions of the $O D E$ are unique for any $x \in A$.
(b) Nonnegative measure-valued solutions of the PDE are unique for any $\bar{\mu}$ concentrated in $A$, i.e. such that $\bar{\mu}\left(\mathbb{R}^{d} \backslash A\right)=0$.

Proof. It is clear that (b) implies (a), just choosing $\bar{\mu}=\delta_{x}$ and noticing that two different solutions $\boldsymbol{X}(t), \tilde{\boldsymbol{X}}(t)$ of the ODE induce two different solutions of the PDE, namely $\delta_{\boldsymbol{X}(t)}$ and $\delta_{\tilde{\boldsymbol{X}}(t)}$. IThe converse implication is less obvious and requires the superposition principle that we are going to describe below, and that provides the representation

$$
\int_{\mathbb{R}^{d}} \varphi d \mu_{t}=\int_{\mathbb{R}^{d}}\left(\int_{\Gamma_{T}} \varphi(\gamma(t)) d \boldsymbol{\eta}_{x}(\gamma)\right) d \mu_{0}(x),
$$

with $\boldsymbol{\eta}_{x}$ probability measures concentrated on the absolutely continuous integral solutions of the ODE starting from $x$.lTherefore, when these are unique, the measures $\boldsymbol{\eta}_{x}$ are unique (and are Dirac masses), so that the solutions of the PDE are unique.
I
We will use the shorter notation $\Gamma_{T}$ for the space $C\left([0, T] ; \mathbb{R}^{d}\right)$ and denote by $e_{t}: \Gamma_{T} \rightarrow \mathbb{R}^{d}$ the evaluation maps $\gamma \mapsto \gamma(t), t \in[0, T]$.

Definition 9. [Superposition solutions] Let $\boldsymbol{\eta} \in \mathscr{M}_{+}\left(\mathbb{R}^{d} \times \Gamma_{T}\right)$ be a measure concentrated on the set of pairs $(x, \gamma)$ such that $\gamma$ is an absolutely continuous integral solution of the ODE with $\gamma(0)=x$. We define

$$
\left\langle\mu_{t}^{\eta}, \varphi\right\rangle:=\int_{\mathbb{R}^{d} \times \Gamma_{T}} \varphi\left(e_{t}(\gamma)\right) d \boldsymbol{\eta}(x, \gamma) \quad \forall \varphi \in C_{b}\left(\mathbb{R}^{d}\right)
$$

By a standard approximation argument the identity defining $\mu_{t}^{\eta}$ holds for any Borel function $\varphi$ such that $\gamma \mapsto \varphi\left(e_{t}(\gamma)\right.$ ) is $\boldsymbol{\eta}$-integrable (or equivalently any $\mu_{t}^{\eta}$-integrable function $\varphi$ ).I

Under the (local) integrability condition

$$
\begin{equation*}
\int_{0}^{T} \int_{\mathbb{R}^{d} \times \Gamma_{T}} \chi_{B_{R}}\left(e_{t}\right)\left|\boldsymbol{b}_{t}\left(e_{t}\right)\right| d \boldsymbol{\eta} d t<+\infty \quad \forall R>0 \tag{7}
\end{equation*}
$$

it is not hard to see that $\mu_{t}^{\eta}$ solves the PDE with the initial condition $\bar{\mu}:=\left(\pi_{\mathbb{R}^{d}}\right)_{\#} \boldsymbol{\eta}$ : indeed, let us check first that $t \mapsto\left\langle\mu_{t}^{\eta}, \varphi\right\rangle$ is absolutely continuous for any $\varphi \in C_{c}^{\infty}\left(\mathbb{R}^{d}\right)$. For every choice of finitely many pairwise disjoint intervals $\left(a_{i}, b_{i}\right) \subset[0, T]$ we have

$$
\sum_{i=1}^{n}\left|\varphi\left(\gamma\left(b_{i}\right)\right)-\varphi\left(\gamma\left(a_{i}\right)\right)\right| \leq \operatorname{Lip}(\varphi) \int_{\cup_{i}\left(a_{i}, b_{i}\right)} \chi_{B_{R}}\left(\left|e_{t}(\gamma)\right|\right) \boldsymbol{b}_{t}\left(e_{t}(\gamma)\right) \mid d t
$$

for $\boldsymbol{\eta}$-a.e. $(x, \gamma)$, with $R$ such that $\operatorname{supp} \varphi \subset \bar{B}_{R}$. ITherefore an integration with respect to $\eta$ gives

$$
\sum_{i=1}^{n}\left|\left\langle\mu_{b_{i}}^{\eta}, \varphi\right\rangle-\left\langle\mu_{a_{i}}^{\eta}, \varphi\right\rangle\right| \leq \operatorname{Lip}(\varphi) \int_{\cup_{i}\left(a_{i}, b_{i}\right)} \int_{\mathbb{R}^{d} \times \Gamma_{T}} \chi_{B_{R}}\left(e_{t}\right)\left|\boldsymbol{b}_{t}\left(e_{t}\right)\right| d \boldsymbol{\eta} d t .
$$

The absolute continuity of the integral shows that the right hand side can be made small when $\sum_{i}\left(b_{i}-a_{i}\right)$ is small. This proves the absolute continuity.I

It remains to evaluate the time derivative of $t \mapsto\left\langle\mu_{t}^{\eta}, \varphi\right\rangle$ : we know that for $\boldsymbol{\eta}$-a.e. $(x, \gamma)$ the identity $\dot{\gamma}(t)=\boldsymbol{b}_{t}(\gamma(t))$ is fulfilled for $\mathscr{L}^{1}$-a.e. $t \in[0, T]$. Then, by Fubini's theorem, we know also that for $\mathscr{L}^{1}$-a.e. $t \in[0, T]$ the previous identity holds for $\boldsymbol{\eta}$-a.e. $(x, \gamma)$, and therefore

$$
\begin{aligned}
& \frac{d}{d t}\left\langle\mu_{t}^{\eta}, \varphi\right\rangle=\frac{d}{d t} \int_{\mathbb{R}^{d}} \varphi(\boldsymbol{X}(t, x)) d \bar{\mu}(x) \\
= & \int_{\mathbb{R}^{d} \times \Gamma_{T}}\left\langle\nabla \varphi\left(e_{t}(\gamma)\right), \boldsymbol{b}_{t}\left(e_{t}(\gamma)\right)\right\rangle d \boldsymbol{\eta}=\left\langle\boldsymbol{b}_{t} \mu_{t}, \nabla \varphi\right\rangle \quad \mathscr{L}^{1} \text {-a.e. in }[0, T] .
\end{aligned}
$$

Remark 10. Actually the formula defining $\mu_{t}^{\eta}$ does not contain $x$, and so it involves only the projection of $\boldsymbol{\eta}$ on $\Gamma_{T}$. Therefore one could also consider measures $\sigma$ in $\Gamma_{T}$, concentrated on the set of solutions of the ODE (for an arbitrary initial point $x$ ). These two viewpoints are basically equivalent: given $\eta$ one can build $\sigma$ just by projection on $\Gamma_{T}$, and given $\sigma$ one can consider the conditional probability measures $\boldsymbol{\eta}_{x}$ concentrated on the solutions of the ODE starting from $x$ induced by the random variable $\gamma \mapsto \gamma(0)$ in $\Gamma_{T}$, the law $\bar{\mu}$ (i.e. the push forward) of the same random variable and recover $\boldsymbol{\eta}$ as follows:

$$
\begin{equation*}
\int_{\mathbb{R}^{d} \times \Gamma_{T}} \varphi(x, \gamma) d \boldsymbol{\eta}(x, \gamma):=\int_{\mathbb{R}^{d}}\left(\int_{\Gamma_{T}} \varphi(x, \gamma) d \boldsymbol{\eta}_{x}(\gamma)\right) d \bar{\mu}(x) . \tag{8}
\end{equation*}
$$

Our viewpoint has been chosen just for technical convenience, to avoid the use, wherever this is possible, of the conditional probability theorem.

By restricting $\boldsymbol{\eta}$ to suitable subsets of $\mathbb{R}^{d} \times \Gamma_{T}$, several manipulations with superposition solutions of the continuity equation are possible and useful, and these are not immediate to see just at the level of general solutions of the continuity equation.IThis is why the following result is interesting.

Theorem 11. [Superposition principle] Let $\mu_{t} \in \mathscr{M}_{+}\left(\mathbb{R}^{d}\right)$ solve PDE and assume that

$$
\int_{0}^{T} \int_{\mathbb{R}^{d}} \frac{|\boldsymbol{b}|_{t}(x)}{1+|x|} d \mu_{t} d t<+\infty
$$

Then $\mu_{t}$ is a superposition solution, i.e. there exists $\boldsymbol{\eta} \in \mathscr{M}_{+}\left(\mathbb{R}^{d} \times \Gamma_{T}\right)$ such that $\mu_{t}=\mu_{t}^{\eta}$ for any $t \in[0, T]$. I

In the proof we use the narrow convergence of positive measures, i.e. the convergence with respect to the duality with continuous and bounded functions, and the easy implication in Prokhorov compactness theorem: any tight and bounded family $\mathscr{F}$ in $\mathscr{M}_{+}(X)$ is (sequentially) relatively compact w.r.t. the narrow convergence.l Remember that tightness means:
for any $\epsilon>0$ there exists $K \subset X$ compact s.t. $\mu(X \backslash K)<\epsilon \forall \mu \in \mathscr{F}$.

A necessary and sufficient condition for tightness is the existence of a coercive functional $\Psi: X \rightarrow[0, \infty]$ such that $\int \Psi d \mu \leq 1$ for any $\mu \in \mathscr{F}$. I

Proof. Step 1 (smoothing). [58] We mollify $\mu_{t}$ w.r.t. the space variable with a kernel $\rho$ having finite first moment $M$ and support equal to the whole of $\mathbb{R}^{d}$ (a Gaussian, for instance), obtaining smooth and strictly positive functions $\mu_{t}^{\epsilon} \cdot \|$ We also choose a function $\psi: \mathbb{R}^{d} \rightarrow[0,+\infty)$ such that $\psi(x) \rightarrow+\infty$ as $|x| \rightarrow+\infty$ and

$$
\int_{\mathbb{R}^{d}} \psi(x) \mu_{0} * \rho_{\epsilon}(x) d x \leq 1 \quad \forall \epsilon \in(0,1)
$$

and a convex nondecreasing function $\Theta: \mathbb{R}^{+} \rightarrow \mathbb{R}$ having a more than linear growth at infinity such that

$$
\int_{0}^{T} \int_{\mathbb{R}^{d}} \frac{\Theta\left(\left|\boldsymbol{b}_{t}\right|(x)\right)}{1+|x|} d \mu_{t} d t<+\infty
$$

(the existence of $\Theta$ is ensured by Dunford-Pettis theorem).IDefining

$$
\mu_{t}^{\epsilon}:=\mu_{t} * \rho_{\epsilon}, \quad \boldsymbol{b}_{t}^{\epsilon}:=\frac{\left(\boldsymbol{b}_{t} \mu_{t}\right) * \rho_{\epsilon}}{\mu_{t}^{\epsilon}}
$$

it is immediate that

$$
\frac{d}{d t} \mu_{t}^{\epsilon}+D_{x} \cdot\left(\boldsymbol{b}_{t}^{\epsilon} \mu_{t}^{\epsilon}\right)=\frac{d}{d t} \mu_{t} * \rho_{\epsilon}+D_{x} \cdot\left(\boldsymbol{b}_{t} \mu_{t}\right) * \rho_{\epsilon}=0
$$

and that $\boldsymbol{b}^{\epsilon} \in L^{1}\left([0, T] ; W_{\text {loc }}^{1, \infty}\left(\mathbb{R}^{d} ; \mathbb{R}^{d}\right)\right)$. Therefore Remark 7 can be applied and the representation $\mu_{t}^{\epsilon}=\boldsymbol{X}^{\epsilon}(t, \cdot)_{\#} \mu_{0}^{\epsilon}$ still holds. IThen, we define

$$
\boldsymbol{\eta}^{\epsilon}:=\left(x, \boldsymbol{X}^{\epsilon}(\cdot, x)\right)_{\#} \mu_{0}^{\epsilon}
$$

so that

$$
\begin{align*}
\int_{\mathbb{R}^{d}} \varphi d \mu_{t}^{\eta_{\epsilon}} & =\int_{\mathbb{R}^{d} \times \Gamma_{T}} \varphi(\gamma(t)) d \boldsymbol{\eta}^{\epsilon}  \tag{9}\\
& =\int_{\mathbb{R}^{d}} \varphi\left(\boldsymbol{X}^{\epsilon}(t, x)\right) d \mu_{0}^{\epsilon}(x)=\int_{\mathbb{R}^{d}} \varphi d \mu_{t}^{\epsilon}
\end{align*}
$$

I
Step 2 (tightness). We will be using the inequality

$$
\begin{equation*}
((1+|x|) c) * \rho_{\epsilon} \leq(1+|x|) * \rho_{\epsilon}+\epsilon c * \tilde{\rho}_{\epsilon} \tag{10}
\end{equation*}
$$

for $c$ nonnegative measure and $\tilde{\rho}(y)=|y| \rho(y)$, and

$$
\begin{equation*}
\Theta\left(\left|\boldsymbol{b}_{t}^{\epsilon}(x)\right|\right) \mu_{t}^{\epsilon}(x) \leq\left(\Theta\left(\left|\boldsymbol{b}_{t}\right|\right) \mu_{t}\right) * \rho_{\epsilon}(x) . \tag{11}
\end{equation*}
$$

The proof of the first one is elementary, while the proof of the second one follows by applying Jensen's inequality with the convex l.s.c. function $(z, t) \mapsto \Theta(|z| / t) t$ (set to $+\infty$ if $t<0$, or $t=0$ and $z \neq 0$, and to 0 if $z=t=0$ ) and with the measure $\rho_{\epsilon}(x-\cdot) \mathscr{L}^{d}$. Let us introduce the functional

$$
\Psi(x, \gamma):=\psi(x)+\int_{0}^{T} \frac{\Theta(|\dot{\gamma}|)}{1+|\gamma|} d t
$$

set to $+\infty$ on $\Gamma_{T} \backslash A C\left([0, T] ; \mathbb{R}^{d}\right)$. ॥Using Ascoli-Arzelá theorem, it is not hard to show that $\Psi$ is coercive (it suffices to show that $\max |\gamma|$ is bounded on the sublevels $\{\Psi \leq t\}$ ). ISince

$$
\begin{aligned}
& \int_{\mathbb{R}^{d} \times \Gamma_{T}} \int_{0}^{T} \frac{\Theta(|\dot{\gamma}|)}{1+|\gamma|} d t d \boldsymbol{\eta}^{\epsilon}(x, \gamma)=\int_{0}^{T} \int_{\mathbb{R}^{d}} \frac{\Theta\left(\left|\boldsymbol{b}_{t}^{\epsilon}\right|\right)}{1+|x|} d \mu_{t}^{\epsilon} d t \\
& \stackrel{(10),(11)}{\leq}(1+\epsilon M) \int_{0}^{T} \int_{\mathbb{R}^{d}} \frac{\Theta\left(\left|\boldsymbol{b}_{t}\right|(x)\right)}{1+|x|} d \mu_{t} d t
\end{aligned}
$$

and

$$
\int_{\mathbb{R}^{d} \times \Gamma_{T}} \psi(x) d \boldsymbol{\eta}^{\epsilon}(x, \gamma)=\int_{\mathbb{R}^{d}} \psi(x) d \mu_{0}^{\epsilon} \leq 1
$$

we obtain that $\int \Psi d \boldsymbol{\eta}^{\epsilon}$ is uniformly bounded for $\epsilon \in(0,1)$, and therefore Prokhorov compactness theorem tells us that the family $\boldsymbol{\eta}^{\epsilon}$ is narrowly sequentially relatively compact as $\epsilon \downarrow 0$.IIf $\boldsymbol{\eta}$ is any limit point we can pass to the limit in (9) to obtain that $\mu_{t}=\mu_{t}^{\eta}$.I

Step 3 ( $\boldsymbol{\eta}$ is concentrated on solutions of the ODE). IIt suffices to show that

$$
\begin{equation*}
\int_{\mathbb{R}^{d} \times \Gamma_{T}} \quad \frac{\left|\gamma(t)-x-\int_{0}^{t} \boldsymbol{b}_{s}(\gamma(s)) d s\right|}{1+\max _{[0, T]}|\gamma|} d \boldsymbol{\eta}=0 \tag{12}
\end{equation*}
$$

for any $t \in[0, T]$. The technical difficulty is that this test function, due to the lack of regularity of $\boldsymbol{b}$, is not continuous.l To this aim, we prove first that

$$
\begin{equation*}
\int_{\mathbb{R}^{d} \times \Gamma_{T}} \frac{\left|\gamma(t)-x-\int_{0}^{t} \boldsymbol{c}_{s}(\gamma(s)) d s\right|}{1+\max _{[0, T]}|\gamma|} d \boldsymbol{\eta} \leq \int_{0}^{T} \int_{\mathbb{R}^{d}} \frac{\left|\boldsymbol{b}_{s}-\boldsymbol{c}_{s}\right|}{1+|x|} d \mu_{s} d s \tag{13}
\end{equation*}
$$

for any continuous function $\boldsymbol{c}$ with compact support.ll Then, choosing a sequence ( $\boldsymbol{c}^{n}$ ) converging to $\boldsymbol{b}$ in $L^{1}\left(\nu ; \mathbb{R}^{d}\right)$, with

$$
\int \varphi(s, x) d \nu(s, x):=\int_{0}^{T} \int_{\mathbb{R}^{d}} \frac{\varphi(s, x)}{1+|x|} d \mu_{s}(x) d s
$$

and noticing that

$$
\int_{\mathbb{R}^{d} \times \Gamma_{T}} \int_{0}^{T} \frac{\left|\boldsymbol{b}_{s}(\gamma(s))-\boldsymbol{c}_{s}^{n}(\gamma(s))\right|}{1+|\gamma(s)|} d s d \boldsymbol{\eta}=\int_{0}^{T} \int_{\mathbb{R}^{d}} \frac{\left|\boldsymbol{b}_{s}-\boldsymbol{c}_{s}^{n}\right|}{1+|x|} d \mu_{s} d s \rightarrow 0,
$$

we can pass to the limit in (13) with $\boldsymbol{c}=\boldsymbol{c}^{n}$ to obtain (12). II
It remains to show (13). This is a limiting argument based on the fact that (12) holds for $\boldsymbol{b}^{\epsilon}, \boldsymbol{\eta}^{\epsilon}:$ II

$$
\begin{aligned}
& \int_{\mathbb{R}^{d} \times \Gamma_{T}} \frac{\left|\gamma(t)-x-\int_{0}^{t} \boldsymbol{c}_{s}(\gamma(s)) d s\right|}{1+\max _{[0, T]}|\gamma|} d \boldsymbol{\eta}^{\epsilon} \\
= & \int_{\mathbb{R}^{d}} \frac{\left|\boldsymbol{X}^{\epsilon}(t, x)-x-\int_{0}^{t} \boldsymbol{c}_{s}\left(\boldsymbol{X}^{\epsilon}(s, x)\right) d s\right|}{1+\max _{[0, T]}\left|\boldsymbol{X}^{\epsilon}(\cdot, x)\right|} d \mu_{0}^{\epsilon}(x) \\
= & \int_{\mathbb{R}^{d}} \frac{\left|\int_{0}^{t} \boldsymbol{b}_{s}^{\epsilon}\left(\boldsymbol{X}^{\epsilon}(s, x)\right)-\boldsymbol{c}_{s}\left(\boldsymbol{X}^{\epsilon}(s, x)\right) d s\right|}{1+\max _{[0, T]}\left|\boldsymbol{X}^{\epsilon}(\cdot, x)\right|} d \mu_{0}^{\epsilon}(x) \leq \int_{0}^{t} \int_{\mathbb{R}^{d}} \frac{\left|\boldsymbol{b}_{s}^{\epsilon}-\boldsymbol{c}_{s}\right|}{1+|x|} d \mu_{s}^{\epsilon} d s \\
\leq & \int_{0}^{t} \int_{\mathbb{R}^{d}} \frac{\left|\boldsymbol{b}_{s}^{\epsilon}-\boldsymbol{c}_{s}^{\epsilon}\right|}{1+|x|} d \mu_{s}^{\epsilon} d s+\leq \int_{0}^{t} \int_{\mathbb{R}^{d}} \frac{\left|\boldsymbol{c}_{s}^{\epsilon}-\boldsymbol{c}_{s}\right|}{1+|x|} d \mu_{s}^{\epsilon} d s \\
\leq & \int_{0}^{t} \int_{\mathbb{R}^{d}} \frac{\left|\boldsymbol{b}_{s}-\boldsymbol{c}_{s}\right|}{1+|x|} d \mu_{s} d s+\int_{0}^{t} \int_{\mathbb{R}^{d}} \frac{\left|\boldsymbol{c}_{s}^{\epsilon}-\boldsymbol{c}_{s}\right|}{1+|x|} d \mu_{s}^{\epsilon} d s .
\end{aligned}
$$

In the last inequalities we added and subtracted $\boldsymbol{c}_{t}^{\epsilon}:=\left(\boldsymbol{c}_{t} \mu_{t}\right) * \rho_{\epsilon} / \mu_{t}^{\epsilon}$.

Since $\boldsymbol{c}_{t}^{\epsilon} \rightarrow \boldsymbol{c}_{t}$ uniformly as $\epsilon \downarrow 0$ thanks to the uniform continuity of $\boldsymbol{c}$, passing to the limit in the chain of inequalities above we obtain (13).

I
The applicability of Theorem 8 is strongly limited by the fact that, on one hand, pointwise uniqueness properties for the ODE are known only in very special situations, for instance when there is a Lipschitz or a one-sided Lipschitz (or log-Lipschitz, Osgood...) condition on b.IOn the other hand, also uniqueness for general measure-valued solutions is known only in special situations.IIt turns out that in many cases uniqueness of the PDE can only be proved in smaller classes $\mathscr{L}$ of solutions, and it is natural to think that this should reflect into a weaker uniqueness condition at the level of the ODE.

We will see indeed that there is uniqueness in the "selection sense". In order to illustrate this concept, in the following we consider a convex class $\mathscr{L}_{b}$ of measure-valued solutions $\mu_{t} \in \mathscr{M}_{+}\left(\mathbb{R}^{d}\right)$ of the continuity equation relative to $\boldsymbol{b}$, satifying the following monotonicity property:

$$
\begin{equation*}
0 \leq \mu_{t}^{\prime} \leq \mu_{t} \in \mathscr{L}_{b} \quad \Longrightarrow \quad \mu_{t}^{\prime} \in \mathscr{L}_{b} \tag{14}
\end{equation*}
$$

whenever $\mu_{t}^{\prime}$ still solves the continuity equation relative to $\boldsymbol{b}$, and the integrability condition

$$
\int_{0}^{T} \int_{\mathbb{R}^{d}} \frac{\left|\boldsymbol{b}_{t}(x)\right|}{1+|x|} d \mu_{t}(x) d t<+\infty
$$

The typical application will be with absolutely continuous measures $\mu_{t}=w_{t} \mathscr{L}^{d}$, whose densities satisfy some quantitative and possibly timedepending bound (e.g. $L^{\infty}\left(L^{1}\right) \cap L^{\infty}\left(L^{\infty}\right)$ )..।

Definition 12. [ $\mathscr{L}_{b}$-lagrangian flows] Given the class $\mathscr{L}_{b}$, we say that $\boldsymbol{X}(t, x)$ is a $\mathscr{L}_{b}$-Lagrangian flow starting from $\bar{\mu} \in \mathscr{M}_{+}\left(\mathbb{R}^{d}\right)$ (at time 0 ) if the following two properties hold:
(a) $\boldsymbol{X}(\cdot, x)$ is absolutely continuous solution in $[0, T]$ and satisfies

$$
\boldsymbol{X}(t, x)=x+\int_{0}^{t} \boldsymbol{b}_{s}(\boldsymbol{X}(s, x)) d s \quad \forall t \in[0, T]
$$

for $\bar{\mu}$-a.e. $x$;
(b) $\mu_{t}:=\boldsymbol{X}(t, \cdot)_{\#} \bar{\mu} \in \mathscr{L}_{b}$.

Heuristically $\mathscr{L}_{b}$-Lagrangian flows can be thought as suitable selections of the solutions of the ODE (possibly non unique), made in such a way to produce a density in $\mathscr{L}_{b}$.I

We will show that the $\mathscr{L}_{b}$-Lagrangian flow starting from $\bar{\mu}$ is unique, modulo $\bar{\mu}$-negligible sets, whenever a comparison principle for the PDE holds, in the class $\mathscr{L}_{b}$ (i.e. the inequality between two solutions at $t=0$ is preserved at later times).

Before stating and proving the uniqueness theorem for $\mathscr{L}_{b}$-Lagrangian flows, we state two elementary but useful results.IThe first one is a simple exercise:

Exercise 13. Let $\sigma \in \mathscr{M}_{+}\left(\Gamma_{T}\right)$ and let $D \subset[0, T]$ be a dense set. Show that $\sigma$ is a Dirac mass in $\Gamma_{T}$ iff its projections $(e(t))_{\#} \sigma, t \in D$, are Dirac masses in $\mathbb{R}^{d}$.

The second one is concerned with a family of measures $\boldsymbol{\eta}_{\boldsymbol{x}}$ :
Lemma 14. Let $\boldsymbol{\eta}_{x}$ be a measurable family of positive finite measures in $\Gamma_{T}$ with the following property: for any $t \in[0, T]$ and any pair of disjoint Borel sets $E, E^{\prime} \subset \mathbb{R}^{d}$ we have

$$
\begin{equation*}
\boldsymbol{\eta}_{x}(\{\gamma: \gamma(t) \in E\}) \boldsymbol{\eta}_{x}\left(\left\{\gamma: \gamma(t) \in E^{\prime}\right\}\right)=0 \quad \bar{\mu} \text {-a.e. in } \mathbb{R}^{d} . \tag{15}
\end{equation*}
$$

Then $\boldsymbol{\eta}_{x}$ is a Dirac mass for $\bar{\mu}$-a.e. $x$.
Proof. Taking into account Exercise 13, for a fixed $t \in(0, T]$ it suffices to check that the measures $\lambda_{x}:=\gamma(t)_{\#} \boldsymbol{\eta}_{x}$ are Dirac masses for $\bar{\mu}$-a.e. $x$. ॥Then (15) gives $\lambda_{x}(E) \lambda_{x}\left(E^{\prime}\right)=0 \bar{\mu}$-a.e. for any pair of disjoint Borel sets $E, E^{\prime} \subset \mathbb{R}^{d} . \|$ Let $\delta>0$ and let us consider a partition of $\mathbb{R}^{d}$ in countably many Borel sets $R_{i}$ having a diameter less then $\delta$.l Then, as $\lambda_{x}\left(R_{i}\right) \lambda_{x}\left(R_{j}\right)=0 \mu$-a.e. whenever $i \neq j$, we have a corresponding decomposition of $\bar{\mu}$-almost all of $\mathbb{R}^{d}$ in Borel sets $A_{i}$ such that supp $\lambda_{x} \subset \bar{R}_{i}$ for any $x \in A_{i}$ (just take $\left\{\lambda_{x}\left(R_{i}\right)>0\right\}$ and subtract from him all other sets $\left.\left\{\lambda_{x}\left(R_{j}\right)>0\right\}, j \neq i\right)$. ISince $\delta$ is arbitrary the statement is proved.

Theorem 15. [Uniqueness of $\mathscr{L}_{b}$-Lagrangian flows] Assume that the PDE fulfils the comparison principle in $\mathscr{L}_{b}$. Then the $\mathscr{L}_{b}$-Lagrangian flow starting from $\bar{\mu}$ is unique, i.e. two different selections $\boldsymbol{X}_{1}(t, x)$ and $\boldsymbol{X}_{2}(t, x)$ of solutions of the ODE inducing solutions of the the continuity equation in $\mathscr{L}$ satisfy

$$
\boldsymbol{X}_{1}(\cdot, x)=\boldsymbol{X}_{2}(\cdot, x) \quad \text { in } \Gamma_{T}, \text { for } \bar{\mu} \text {-a.e. } x .
$$

Proof. If the statement were false we could produce a measure $\boldsymbol{\eta}$ not concentrated on a graph inducing a solution $\mu_{t}^{\eta} \in \mathscr{L}_{b}$ of the PDE. This is not possible, thanks to the next result. The measure $\eta$ can be built as follows:

$$
\boldsymbol{\eta}:=\frac{1}{2}\left(\boldsymbol{\eta}^{1}+\boldsymbol{\eta}^{2}\right)=\frac{1}{2}\left(\left(x, \boldsymbol{X}_{1}(\cdot, x)\right)_{\#} \bar{\mu}+\left(x, \boldsymbol{X}_{2}(\cdot, x)\right)_{\# \bar{\mu}}\right) .
$$

Since $\mathscr{L}_{b}$ is convex we still have $\mu_{t}^{\eta}=\frac{1}{2}\left(\mu_{t}^{\eta^{1}}+\mu_{t}^{\eta^{2}}\right) \in \mathscr{L}_{b}$.

Remark 16. In the same vein, one can also show that

$$
\boldsymbol{X}_{1}(\cdot, x)=\boldsymbol{X}_{2}(\cdot, x) \quad \text { in } \Gamma_{T} \text { for } \bar{\mu}_{1} \wedge \bar{\mu}_{2} \text {-a.e. } x
$$

whenever $\boldsymbol{X}_{1}, \boldsymbol{X}_{2}$ are $\mathscr{L}_{\boldsymbol{b}}$-Lagrangian flows starting respectively from $\bar{\mu}_{1}$ and $\bar{\mu}_{2}$.

We used the following basic result, having some analogy with Kantorovich's and Mather's theories.

Theorem 17. Assume that the PDE fulfils the comparison principle in $\mathscr{L}_{b}$. Let $\boldsymbol{\eta} \in \mathscr{M}_{+}\left(\mathbb{R}^{d} \times \Gamma_{T}\right)$ be concentrated on the pairs $(x, \gamma)$ with $\gamma$ absolutely continuous solution of the ODE, and assume that $\mu_{t}^{\eta} \in \mathscr{L}_{b}$. Then $\boldsymbol{\eta}$ is concentrated on a graph, i.e. there exists a function $x \mapsto X(\cdot, x) \in \Gamma_{T}$ such that

$$
\boldsymbol{\eta}=(x, X(\cdot, x))_{\#} \bar{\mu}, \quad \text { with } \quad \bar{\mu}:=\left(\pi_{\mathbb{R}^{d}}\right)_{\#} \boldsymbol{\eta}=\mu_{0}^{\eta} .
$$

Proof. We use the representation (8) of $\eta$, given by the disintegration theorem, the criterion stated in Lemma 14 and argue by contradiction.If the thesis is false then $\boldsymbol{\eta}_{x}$ is not a Dirac mass in a set of $\bar{\mu}$ positive measure and we can find $t \in(0, T]$, disjoint Borel sets $E, E^{\prime} \subset \mathbb{R}^{d}$ and a Borel set $C$ with $\bar{\mu}(C)>0$ such that

$$
\boldsymbol{\eta}_{x}(\{\gamma: \gamma(t) \in E\}) \boldsymbol{\eta}_{x}\left(\left\{\gamma: \gamma(t) \in E^{\prime}\right\}\right)>0 \quad \forall x \in C .
$$

Possibly passing to a smaller set having still strictly positive $\bar{\mu}$ measure we can assume that

$$
\begin{equation*}
0<\boldsymbol{\eta}_{x}(\{\gamma: \gamma(t) \in E\}) \leq M \boldsymbol{\eta}_{x}\left(\left\{\gamma: \gamma(t) \in E^{\prime}\right\}\right) \quad \forall x \in C \tag{16}
\end{equation*}
$$

for some constant $M . \|$ We define measures $\boldsymbol{\eta}^{1}, \boldsymbol{\eta}^{2}$ whose disintegrations $\boldsymbol{\eta}_{x}^{1}, \boldsymbol{\eta}_{x}^{2}$ are given by

$$
\boldsymbol{\eta}_{x}^{1}:=\chi_{C}(x) \boldsymbol{\eta}_{x}\left\llcorner\{\gamma: \gamma(t) \in E\}, \quad \boldsymbol{\eta}_{x}^{2}:=M \chi_{C}(x) \boldsymbol{\eta}_{x}\left\llcorner\left\{\gamma: \gamma(t) \in E^{\prime}\right\}\right.\right.
$$

and denote by $\mu_{t}^{i}$ the (superposition) solutions of the continuity equation induced by $\boldsymbol{\eta}^{i} . \|$ Then

$$
\mu_{0}^{1}=\boldsymbol{\eta}_{x}(\{\gamma: \gamma(t) \in E\}) \bar{\mu}\left\llcorner C, \quad \mu_{0}^{2}=M \boldsymbol{\eta}_{x}\left(\left\{\gamma: \gamma(t) \in E^{\prime}\right\}\right) \bar{\mu}\llcorner C\right.
$$

so that (16) yields $\mu_{0}^{1} \leq \mu_{0}^{2}$.I On the other hand, $\mu_{t}^{1}$ is orthogonal to $\mu_{t}^{2}$ : precisely, denoting by $\boldsymbol{\eta}_{t x}$ the image of $\eta$ under the map $\gamma \mapsto \gamma(t)$, we have

$$
\mu_{t}^{1}=\int_{C} \boldsymbol{\eta}_{t x}\left\llcornerE d \mu ( x ) \perp M \int _ { C } \boldsymbol { \eta } _ { t x } \left\llcorner E^{\prime} d \mu(x)=\mu_{t}^{2} .\right.\right.
$$

Notice also that $\mu_{t}^{i} \leq \mu_{t}$ and so the monotonicity assumption (14) on $\mathscr{L}_{b}$ gives $\mu_{t}^{i} \in \mathscr{L}_{b}$. This contradicts the assumption on the validity of the comparison principle in $\mathscr{L}_{b}$.

Now we come to the existence of $\mathscr{L}_{b}$-Lagrangian flows.
Theorem 18. [Existence of $\mathscr{L}_{b}$-Lagrangian flows] Assume that the PDE fulfils the comparison principle in $\mathscr{L}_{b}$ and that for some $\bar{\mu} \in \mathscr{M}_{+}\left(\mathbb{R}^{d}\right)$ there exists a solution $\mu_{t} \in \mathscr{L}_{b}$ with $\mu_{0}=\bar{\mu}$. Then there exists a (unique) $\mathscr{L}_{b}$-Lagrangian flow starting from $\bar{\mu}$.
I
Proof. By the superposition principle we can represent $\mu_{t}$ as $\left(e_{t}\right)_{\#} \boldsymbol{\eta}$ for some $\boldsymbol{\eta} \in \mathscr{M}_{+}\left(\mathbb{R}^{d} \times \Gamma_{T}\right)$ concentrated on pairs $(x, \gamma)$ solutions of the ODE. Then, Theorem 17 tells us that $\boldsymbol{\eta}$ is concentrated on a graph, i.e. there exists a function $x \mapsto \boldsymbol{X}(\cdot, x) \in \Gamma_{T}$ such that

$$
(x, \boldsymbol{X}(\cdot, x))_{\#} \bar{\mu}=\boldsymbol{\eta} .
$$

Pushing both sides via $e_{t}$ we obtain

$$
\boldsymbol{X}(t, \cdot)_{\#} \bar{\mu}=\left(e_{t}\right)_{\#} \boldsymbol{\eta}=\mu_{t} \in \mathscr{L}_{\boldsymbol{b}},
$$

and therefore $\boldsymbol{X}$ is a $\mathscr{L}_{\boldsymbol{b}}$-Lagrangian flow.

Finally, let us discuss the stability issue. This is particularly relevant, as we will see, in connection with the applications to PDE's.

Definition 19. [Convergence of velocity fields] We define the convergence of $\boldsymbol{b}^{h}$ to $\boldsymbol{b}$ in a indirect way, defining rather a convergence of $\mathscr{L}_{\boldsymbol{b}^{h}}$ to $\mathscr{L}_{b}$ : we require that

$$
\boldsymbol{b}^{h} \mu_{t}^{h} \rightharpoonup \boldsymbol{b} \mu \text { in }(0, T) \times \mathbb{R}^{d} \quad \text { and } \quad \mu_{t} \in \mathscr{L}_{b}
$$

whenever $\mu_{t}^{h} \in \mathscr{L}_{b^{h}}$ and $\mu_{t}^{h} \rightarrow \mu_{t}$ narrowly for all $t \in[0, T]$. I

For instance, in the typical case when $\mathscr{L}$ is bounded and closed, w.r.t the weak* topology, in $L^{\infty}\left(L^{1}\right) \cap L^{\infty}\left(L^{\infty}\right)$, and

$$
\mathscr{L}_{c}:=\mathscr{L} \cap\left\{w: \frac{d}{d t} w+D_{x} \cdot(\boldsymbol{c} w)=0\right\}
$$

the implication is fulfilled whenever $\boldsymbol{b}^{h} \rightarrow \boldsymbol{b}$ strongly in $L_{\mathrm{loc}}^{1}$. II
The natural convergence for the stability theorem is convergence in measure.ll Let us recall that a $Y$-valued sequence $\left(v_{h}\right)$ is said to converge in $\bar{\mu}$-measure to $v$ if

$$
\lim _{h \rightarrow \infty} \bar{\mu}\left(\left\{d_{Y}\left(v_{h}, v\right)>\delta\right\}\right)=0 \quad \forall \delta>0
$$

This is equivalent to the $L^{1}$ convergence to 0 of the $\mathbb{R}^{+}$-valued maps $1 \wedge d_{Y}\left(v_{h}, v\right)$.

Recall also that convergence $\bar{\mu}$-a.e. implies convergence in measure, and that the converse implication is true passing to a suitable subsequence.

Theorem 20. [Stability of $\mathscr{L}$-Lagrangian flows] Assume that
(i) $\mathscr{L}_{b^{h}}$ converge to $\mathscr{L}_{b}$;
(ii) $\boldsymbol{X}^{h}$ are $\mathscr{L}_{b^{h}}$-flows relative to $\boldsymbol{b}^{h}$ starting from $\bar{\mu} \in \mathscr{M}_{+}\left(\mathbb{R}^{d}\right)$;
(iii) setting $\mu_{t}^{h}:=\boldsymbol{X}^{h}(t, \cdot)_{\#} \bar{\mu}$, we have

$$
\begin{equation*}
\int_{0}^{T} \int_{\mathbb{R}^{d}} \frac{\Theta\left(\left|\boldsymbol{b}_{t}^{h}\right|\right)}{1+|x|} d \mu_{t}^{h}(x) d t \leq 1 \tag{17}
\end{equation*}
$$

for some function $\theta: \mathbb{R}^{+} \rightarrow \mathbb{R}$ having a more than linear growth at infinity, and

$$
\begin{equation*}
\limsup _{h \rightarrow \infty} \int_{0}^{T} \int_{\mathbb{R}^{d}} \frac{\left|\boldsymbol{b}_{t}^{h}\right|}{1+|x|} d \mu_{t}^{h} d t \leq \int_{0}^{T} \int_{\mathbb{R}^{d}} \frac{\left|\boldsymbol{b}_{t}\right|}{1+|x|} d \mu_{t} d t ; \tag{18}
\end{equation*}
$$

(iv) the PDE fulfils the comparison principle in $\mathscr{L}_{b}$.

Then $x \mapsto \boldsymbol{X}^{h}(\cdot, x)$ converge to $x \mapsto \boldsymbol{X}(\cdot, x)$ in $\bar{\mu}$-measure, i.e.

$$
\lim _{h \rightarrow \infty} \int_{\mathbb{R}^{d}} 1 \wedge \sup _{[0, T]}\left|\boldsymbol{X}^{h}(\cdot, x)-\boldsymbol{X}(\cdot, x)\right| d \bar{\mu}(x)=0
$$

Proof. Following the same strategy used in the proof of the superposition principle, we push $\bar{\mu}$ onto the graph of the map $x \mapsto \boldsymbol{X}^{h}(\cdot, x)$, i.e.

$$
\boldsymbol{\eta}^{h}:=\left(x, \boldsymbol{X}^{h}(\cdot, x)\right)_{\#} \bar{\mu}
$$

and we obtain, using (17) in (iii) and the same argument used in Step 2 of the proof of the superposition principle, that $\boldsymbol{\eta}^{h}$ is tight in $\mathscr{M}_{+}\left(\mathbb{R}^{d} \times\right.$ $\Gamma_{T}$ ).

Let now $\boldsymbol{\eta}$ be any limit point of $\boldsymbol{\eta}^{h}$. Using the same argument used in Step 3 of the proof of the superposition principle and (18) we obtain that $\boldsymbol{\eta}$ is concentrated on pairs $(x, \gamma)$ with $\gamma$ absolutely continuous solution of the ODE relative to $\boldsymbol{b}$ starting from $x$. Indeed, this argument was using only the property

$$
\lim _{h \rightarrow \infty} \int_{0}^{T} \int_{\mathbb{R}^{d}} \frac{\left|\boldsymbol{b}_{t}^{h}-\boldsymbol{c}_{t}\right|}{1+|x|} d \mu_{t}^{h} d t=\int_{0}^{T} \int_{\mathbb{R}^{d}} \frac{\left|\boldsymbol{b}_{t}-\boldsymbol{c}_{t}\right|}{1+|x|} d \mu_{t} d t
$$

for any continuous function $\boldsymbol{c}$ with compact support in $(0, T) \times \mathbb{R}^{d}$, and this property is ensured by Lemma 22 below.

Let $\mu_{t}:=\left(\boldsymbol{e}_{t}\right)_{\#} \boldsymbol{\eta}$ and notice that $\mu_{t}^{h}=\left(\boldsymbol{e}_{t}\right)_{\#} \boldsymbol{\eta}^{h}$, hence $\mu_{t}^{h} \rightarrow \mu_{t}$ narrowly for any $t \in[0, T] \|$ As $\mu_{t}^{h} \in \mathscr{L}_{b^{h}}$, assumption (i) gives that $\mu_{t} \in \mathscr{L}_{b}$ and assumption (iv) together with Theorem 17 imply that $\boldsymbol{\eta}$ is concentrated on the graph of the map $x \mapsto \boldsymbol{X}(\cdot, x)$, where $\boldsymbol{X}$ is the unique $\mathscr{L}_{b}{ }^{-}$ Lagrangian flow.IWe have thus obtained that

$$
\left(x, \boldsymbol{X}^{h}(\cdot, x)\right)_{\#} \bar{\mu} \quad \rightharpoonup \quad(x, \boldsymbol{X}(\cdot, x))_{\#} \bar{\mu} .
$$

By applying the following general principle we conclude.

Lemma 21. [Narrow convergence and convergence in measure] Let $v_{h}, v: X \rightarrow Y$ be Borel maps and let $\mu \in \mathscr{M}_{+}(X)$. Then $v_{h} \rightarrow v$ in $\bar{\mu}$-measure iff

$$
\left(x, v_{h}(x)\right)_{\#} \bar{\mu} \text { converges to }(x, v(x))_{\#} \bar{\mu} \text { narrowly in } \mathscr{M}_{+}(X \times Y) .
$$

Proof. If $v_{h} \rightarrow v$ in $\bar{\mu}$-measure then $\varphi\left(x, v_{h}(x)\right)$ converges in $L^{1}(\bar{\mu})$ to $\varphi(x, v(x))$, and we immediately obtain the convergence of the pushforward measures.l Conversely, let $\delta>0$ and, for any $\epsilon>0$, let $w \in$ $C_{b}(X ; Y)$ be such that $\bar{\mu}(\{v \neq w\}) \leq \epsilon$. We define

$$
\varphi(x, y):=1 \wedge \frac{d_{Y}(y, w(x))}{\delta} \in C_{b}(X \times Y)
$$

and notice that

$$
\begin{gathered}
\bar{\mu}(\{v \neq w\})+\int_{X \times Y} \varphi d\left(x, v_{h}(x)\right)_{\#} \bar{\mu} \geq \bar{\mu}\left(\left\{d_{Y}\left(v, v_{h}\right)>\delta\right\}\right), \\
\int_{X \times Y} \varphi d(x, v(x))_{\#} \bar{\mu} \leq \bar{\mu}(\{w \neq v\}) .
\end{gathered}
$$

Taking into account the narrow convergence of the push-forward we obtain that

$$
\limsup _{h \rightarrow \infty} \bar{\mu}\left(\left\{d_{Y}\left(v, v_{h}\right)>\delta\right\}\right) \leq 2 \bar{\mu}(\{w \neq v\}) \leq 2 \epsilon
$$

and since $\epsilon$ is arbitrary the proof is achieved.
I
Lemma 22. Let $A \subset \mathbb{R}^{m}$ be an open set, and let $\sigma^{h} \in \mathscr{M}_{+}(A)$ be narrowly converging to $\sigma \in \mathscr{M}_{+}(A)$. Let $\boldsymbol{f}^{h} \in L^{1}\left(A, \sigma^{h}, \mathbb{R}^{k}\right), \boldsymbol{f} \in$ $L^{1}\left(A, \sigma, \mathbb{R}^{k}\right)$ and assume that
(i) $\boldsymbol{f}^{h} \sigma^{h}$ weakly converge, in the duality with $C_{c}\left(A ; \mathbb{R}^{k}\right)$, to $\boldsymbol{f} \sigma$;
(ii) $\int_{A} \Theta\left(\boldsymbol{f}^{h}\right) d \sigma^{h}$ is uniformly bounded, for some function $\Theta: \mathbb{R}^{+} \rightarrow \mathbb{R}^{+}$ having a more than linear growth at infinity;
(iii) $\limsup _{h \rightarrow \infty} \int_{A}\left|\boldsymbol{f}^{h}\right| d \sigma^{h} \leq \int_{A}|\boldsymbol{f}| d \sigma$.

$$
h \rightarrow \infty
$$

Then $\int_{A}\left|\boldsymbol{f}^{h}-\boldsymbol{c}\right| d \sigma^{h} \rightarrow \int_{A}|\boldsymbol{f}-\boldsymbol{c}| d \sigma$ for any $\boldsymbol{c} \in C_{b}\left(A ; \mathbb{R}^{k}\right)$. I

Proof. We consider the measures $\nu^{h}:=\left(x, \boldsymbol{f}^{h}(x)\right)_{\#} \sigma^{h}$ in $A \times \mathbb{R}^{k}$ and we assume, possibly extracting a subsequence, that $\nu^{h} \rightharpoonup \nu$, with $\nu \in \mathscr{M}_{+}\left(A \times \mathbb{R}^{k}\right)$, in the duality with $C_{c}\left(A \times \mathbb{R}^{k}\right)$.IUsing condition (ii), the narrow convergence of $\sigma^{h}$ and a truncation argument it is easy to see that the convergence actually occurs for any continuous test function $\psi(x, y)$ having a linear growth w.r.t. $y$ (uniform w.r.t. $x$ ).l Then, choosing test functions $\psi=\psi(x) \in C_{b}(A)$, the convergence of $\sigma^{h}$ to $\sigma$ gives

$$
\int_{A \times \mathbb{R}^{k}} \psi d \nu=\int_{A} \psi d \sigma
$$

$l$ and therefore, according to the disintegration theorem, we can represent $\nu$ as

$$
\begin{equation*}
\int_{A \times \mathbb{R}^{k}} \psi(x, y) d \nu(x, y)=\int_{A}\left(\int_{\mathbb{R}^{k}} \psi(x, y) d \nu_{x}(y)\right) d \sigma(x) \tag{19}
\end{equation*}
$$

for a suitable Borel family of probability measures $\nu_{x}$ in $\mathbb{R}^{k} . \|$ Next, we can use $\psi(x) y_{j}$ as test functions and assumption (i), to obtain
$\lim _{h \rightarrow \infty} \int_{A} \boldsymbol{f}_{j}^{h} \psi d \mu^{h}=\lim _{h \rightarrow \infty} \int_{A \times \mathbb{R}^{k}} \psi(x) y_{j} d \nu^{h}=\int_{A} \psi(x)\left(\int_{\mathbb{R}^{k}} y_{j} d \nu_{x}(y)\right) d \sigma(x$
$\|$ As $\psi$ and $j$ are arbitrary, this means that the first moment $\nu_{x}$, i.e.
$\int y d \nu_{x}$, is equal to $\boldsymbol{f}(x)$ for $\sigma$-a.e. $x$.l

On the other hand, choosing $|y|$ as test function, assumption (iii) gives
$\int_{A} \int_{\mathbb{R}^{k}}|y| d \nu_{x}(y) d \sigma(x)=\lim _{h \rightarrow \infty} \int_{A \times \mathbb{R}^{k}}|y| d \nu^{h}=\lim _{h \rightarrow \infty} \int_{A}\left|\boldsymbol{f}^{h}\right| d \sigma^{h}=\int_{A}|\boldsymbol{f}| d \sigma$,
hence $\int|y| d \nu_{x}=\boldsymbol{f}(x)=\left|\int y d \nu_{x}\right|$ for $\sigma$-a.e. $x$. This can happen only if $\nu_{x}=\delta_{f(x)}$ for $\sigma$-a.e. $x$.

Finally, taking into account the representation (19) of $\nu$ with $\nu_{x}=$ $\delta_{f(x)}$, the convergence statement can be achieved just choosing the test function $\psi(x, y)=|y-\boldsymbol{c}(x)|$.

## 4. Vector fields with a Sobolev spatial regularity

Here we discuss the well-posedness of the continuity or transport equations assuming the $\boldsymbol{b}_{t}(\cdot)$ has a Sobolev regularity, following [53]. Then, the general theory previously developed provides existence, uniqueness and stability of the $\mathscr{L}$-Lagrangian flow, with $\mathscr{L}:=L^{\infty}\left(L^{1}\right) \cap L^{\infty}\left(L^{\infty}\right)$. We denote by $I \subset \mathbb{R}$ an open interval.

Definition 23. [Renormalized solutions] Let $\boldsymbol{b} \in L_{\text {loc }}^{1}\left(I ; L_{\text {loc }}^{1}\left(\mathbb{R}^{d} ; \mathbb{R}^{d}\right)\right)$ be such that $D \cdot \boldsymbol{b}_{t}=\operatorname{div} \boldsymbol{b}_{t} \mathscr{L}^{d}$ for $\mathscr{L}^{1}$-a.e. $t \in I$, with

$$
\operatorname{div} \boldsymbol{b}_{t} \in L_{\mathrm{loc}}^{1}\left(I ; L_{\mathrm{loc}}^{1}\left(\mathbb{R}^{d}\right)\right)
$$

Let $w \in L_{\text {loc }}^{\infty}\left(I ; L_{\text {loc }}^{\infty}\left(\mathbb{R}^{d}\right)\right)$ and assume that

$$
\begin{equation*}
c:=\frac{d}{d t} w+\boldsymbol{b} \cdot \nabla w \in L_{\mathrm{loc}}^{1}\left(I \times \mathbb{R}^{d}\right) . \tag{20}
\end{equation*}
$$

IThen, we say that $w$ is a renormalized solution of (20) if

$$
\frac{d}{d t} \beta(w)+\boldsymbol{b} \cdot \nabla \beta(w)=c \beta^{\prime}(w) \quad \forall \beta \in C^{1}(\mathbb{R})
$$

Equivalently, recalling the definition of the distribution $\boldsymbol{b} \cdot \nabla w$, the definition could be given in a conservative form, writing

$$
\frac{d}{d t} \beta(w)+D_{x} \cdot(\boldsymbol{b} \beta(w))=c \beta^{\prime}(w)+\operatorname{div} \boldsymbol{b}_{t} \beta(w) .
$$

Notice also that the concept makes sense, choosing properly the class of "test" functions $\beta$, also for $w$ that do not satisfy (20), or are not even locally integrable. This is particularly relevant in connection with DiPerna-Lions's existence theorem for Boltzmann equation, or with the case when $w$ is the characteristic of an unbounded vector field $\boldsymbol{b}$. This concept is also reminiscent of Kruzkhov's concept of entropy solution for a scalar conservation law

$$
\frac{d}{d t} u+D_{x} \cdot(\boldsymbol{f}(u))=0 \quad u:(0,+\infty) \times \mathbb{R}^{d} \rightarrow \mathbb{R}
$$

In this case a distributional one-sided inequality is required:

$$
\frac{d}{d t} \eta(u)+D_{x} \cdot(\boldsymbol{q}(u)) \leq 0
$$

for any convex entropy-entropy flux pair $(\eta, \boldsymbol{q})$ (i.e. $\eta$ is convex and $\eta^{\prime} \boldsymbol{f}^{\prime}=\boldsymbol{q}^{\prime}$ ).

Remark 24. [Time continuity] Using the fact that both $t \mapsto w_{t}$ and $t \mapsto \beta\left(w_{t}\right)$ have a uniformly continuous representative (w.r.t. the $w^{*}-$ $L_{\text {loc }}^{\infty}$ topology), we obtain that, for any renormalized solution $w, t \mapsto w_{t}$ has a unique representative which is continuous w.r.t. the $L_{\text {loc }}^{1}$ topology. The proof follows by a classical weak-strong convergence argument:

$$
f_{n} \rightharpoonup f, \quad \beta\left(f_{n}\right) \rightharpoonup \beta(f) \quad \Longrightarrow \quad f_{n} \rightarrow f
$$

provided $\beta$ is strictly convex. In the case of scalar conservation laws there are analogous results [82], [73].

Using the concept of renormalized solution we can prove a comparison principle in the following natural class $\mathscr{L}$ :

$$
\begin{align*}
\mathscr{L}:= & \left\{w \in L^{\infty}\left([0, T] ; L^{1}\left(\mathbb{R}^{d}\right)\right) \cap L^{\infty}\left([0, T] ; L^{\infty}\left(\mathbb{R}^{d}\right)\right):\right.  \tag{21}\\
& \left.w \in C\left([0, T] ; w^{*}-L^{\infty}\left(\mathbb{R}^{d}\right)\right)\right\} .
\end{align*}
$$

Theorem 25. [Comparison principle] Assume that

$$
\begin{equation*}
\frac{|\boldsymbol{b}|}{1+|x|} \in L^{1}\left([0, T] ; L^{\infty}\left(\mathbb{R}^{d}\right)\right)+L^{1}\left([0, T] ; L^{1}\left(\mathbb{R}^{d}\right)\right) \tag{22}
\end{equation*}
$$

that $D \cdot \boldsymbol{b}_{t}=\operatorname{div} \boldsymbol{b}_{t} \mathscr{L}^{d}$ for $\mathscr{L}^{1}$-a.e. $t \in[0, T]$, and that

$$
\begin{equation*}
\left[\operatorname{div} \boldsymbol{b}_{t}\right]^{-} \in L_{\mathrm{loc}}^{1}\left([0, T) \times \mathbb{R}^{d}\right) \tag{23}
\end{equation*}
$$

Setting $\boldsymbol{b}_{t} \equiv 0$ for $t<0$, assume in addition that any solution of (20) in $(-\infty, T) \times \mathbb{R}^{d}$ is renormalized. Then the comparison principle for the continuity equation holds in the class $\mathscr{L}$.

Proof. By the linearity of the equation, it suffices to show that $w \in \mathscr{L}$ and $w_{0} \leq 0$ implies $w_{t} \leq 0$ for any $t \in[0, T] . \|$ We extend first the PDE to negative times, setting $w_{t}=w_{0}$.IThen, fix a cut-off function $\varphi \in C_{c}^{\infty}\left(\mathbb{R}^{d}\right)$ with $\operatorname{supp} \varphi \subset \bar{B}_{2}(0)$ and $\varphi \equiv 1$ on $B_{1}(0)$, and the renormalization functions

$$
\beta_{\epsilon}(t):=\sqrt{\epsilon^{2}+\left(t^{+}\right)^{2}}-\epsilon \in C^{1}(\mathbb{R}) .
$$

॥Notice that

$$
\begin{equation*}
\beta_{\epsilon}(t) \uparrow t^{+} \quad \text { as } \epsilon \downarrow 0, \quad t \beta_{\epsilon}^{\prime}(t)-\beta_{\epsilon}(t) \in[0, \epsilon] . \tag{24}
\end{equation*}
$$

We know that

$$
\frac{d}{d t} \beta_{\epsilon}\left(w_{t}\right)+D_{x} \cdot\left(\boldsymbol{b} \beta_{\epsilon}\left(w_{t}\right)\right)=\operatorname{div} \boldsymbol{b}_{t}\left(\beta_{\epsilon}\left(w_{t}\right)-w_{t} \beta_{\epsilon}^{\prime}\left(w_{t}\right)\right)
$$

in the sense of distributions in $(-\infty, T) \times \mathbb{R}^{d} . \|$ Plugging $\varphi_{R}(\cdot):=\varphi(\cdot / R)$, with $R \geq 1$, into the PDE we obtain
$\frac{d}{d t} \int_{\mathbb{R}^{d}} \varphi_{R} \beta_{\epsilon}\left(w_{t}\right) d x=\int_{\mathbb{R}^{d}} \beta_{\epsilon}\left(w_{t}\right)\left\langle\boldsymbol{b}_{t}, \nabla \varphi_{R}\right\rangle d x+\int_{\mathbb{R}^{d}} \varphi_{R} \operatorname{div} \boldsymbol{b}_{t}\left(\beta_{\epsilon}\left(w_{t}\right)-w_{t} \beta_{\epsilon}^{\prime}\left(w_{t}\right.\right.$
Splitting $\boldsymbol{b}$ as $\boldsymbol{b}_{1}+\boldsymbol{b}_{2}$, with

$$
\frac{\boldsymbol{b}_{1}}{1+|x|} \in L^{1}\left([0, T] ; L^{\infty}\left(\mathbb{R}^{d}\right)\right) \quad \text { and } \quad \frac{\boldsymbol{b}_{2}}{1+|x|} \in L^{1}\left([0, T] ; L^{1}\left(\mathbb{R}^{d}\right)\right)
$$

and using the inequality

$$
\frac{1}{R} \chi_{\{R \leq|x| \leq 2 R\}} \leq \frac{3}{1+|x|} \chi_{\{R \leq|x|\}}
$$

we can estimate the first integral in the right hand side with $3\|\nabla \varphi\|_{\infty}\left\|\frac{\boldsymbol{b}_{1 t}}{1+|x|}\right\|_{\infty} \int_{\{|x| \geq R\}}\left|w_{t}\right| d x+3\|\nabla \varphi\|_{\infty}\left\|w_{t}\right\|_{\infty} \int_{\{|x| \geq R\}} \frac{\left|\boldsymbol{b}_{1 t}\right|}{1+|x|} d x$. The second integral can be estimated with

$$
\epsilon \int_{\mathbb{R}^{d}} \varphi_{R}\left[\operatorname{div} \boldsymbol{b}_{t}\right]^{-} d x
$$

Passing to the limit first as $\epsilon \downarrow 0$ and then as $R \rightarrow+\infty$ and using the integrability assumptions on $b$ and $w$ we get

$$
\frac{d}{d t} \int_{\mathbb{R}^{d}} w_{t}^{+} d x \leq 0
$$

in the distribution sense in $\mathbb{R}$. Since the function vanishes for negative times, this suffices to conclude using Gronwall lemma.

I
Remark 26. It would be nice to have a completely non-linear comparison principle between renormalized solutions, as in the Kruzkhov theory. Here, on the other hand, we rather used the fact that the difference of the two solutions is renormalized.

I
In any case, Di Perna and Lions proved that all distributional solutions are renormalized when there is a Sobolev regularity with respect to the spatial variables.

Theorem 27. Let $b \in L_{\mathrm{loc}}^{1}\left(I ; W_{\mathrm{loc}}^{1,1}\left(\mathbb{R}^{d} ; \mathbb{R}^{d}\right)\right)$ and let $w \in L_{\mathrm{loc}}^{\infty}\left(I \times \mathbb{R}^{d}\right)$ be a distributional solution of (20). Then $w$ is a renormalized solution. I

Proof. We mollify with respect to the spatial variables and we set

$$
r^{\epsilon}:=(\boldsymbol{b} \cdot \nabla w) * \rho_{\epsilon}-\boldsymbol{b} \cdot\left(\nabla\left(w * \rho_{\epsilon}\right)\right), \quad w^{\epsilon}:=w * \rho_{\epsilon}
$$

to obtain

$$
\frac{d}{d t} w^{\epsilon}+\boldsymbol{b} \cdot \nabla w^{\epsilon}=c * \rho_{\epsilon}-r^{\epsilon}
$$

By the smoothness of $w^{\epsilon}$ w.r.t. $x$, the PDE above tells us that $\frac{d}{d t} w_{t}^{\epsilon} \in$ $L_{\mathrm{loc}}^{1}$, therefore $w^{\epsilon} \in W_{\mathrm{loc}}^{1,1}\left(I \times \mathbb{R}^{d}\right)$ and we can apply the standard chain rule in Sobolev spaces, getting

$$
\frac{d}{d t} \beta\left(w^{\epsilon}\right)+\boldsymbol{b} \cdot \nabla \beta\left(w^{\epsilon}\right)=\beta^{\prime}\left(w^{\epsilon}\right) c * \rho_{\epsilon}-\beta^{\prime}\left(w^{\epsilon}\right) r^{\epsilon}
$$

When we let $\epsilon \downarrow 0$ the convergence in the distribution sense of all terms in the identity above is trivial, with the exception of the last one. To ensure its convergence to zero, it seems necessary to show that $r^{\epsilon} \rightarrow 0$ strongly in $L_{\text {loc }}^{1}$ (remember that $\beta^{\prime}\left(w^{\epsilon}\right)$ is locally equibounded w.r.t. $\epsilon$ ). This is indeed the case, and it is exactly here that the Sobolev regularity plays a role.

## Proposition 28. [Strong convergence of commutators] If

 $w \in L_{\mathrm{loc}}^{\infty}\left(I \times \mathbb{R}^{d}\right)$ and $\boldsymbol{b} \in L_{\mathrm{loc}}^{1}\left(I ; W_{\mathrm{loc}}^{1,1}\left(\mathbb{R}^{d} ; \mathbb{R}^{d}\right)\right)$ we have$$
L_{\mathrm{loc}}^{1}-\lim _{\epsilon \downarrow 0}(\boldsymbol{b} \cdot \nabla w) * \rho_{\epsilon}-\boldsymbol{b} \cdot\left(\nabla\left(w * \rho_{\epsilon}\right)\right)=0 .
$$

Proof. Playing with the definitions of $\boldsymbol{b} \cdot \nabla w$ and convolution product of a distribution and a smooth function, one proves first the identity

$$
\begin{equation*}
r^{\epsilon}(t, x)=\int_{\mathbb{R}^{d}} w(t, x-\epsilon y) \frac{\left(\boldsymbol{b}_{t}(x-\epsilon y)-\boldsymbol{b}_{t}(x)\right) \cdot \nabla \rho(y)}{\epsilon} d y-\left(w \operatorname{div} \boldsymbol{b}_{t}\right) * \rho_{\epsilon}(x) . \tag{25}
\end{equation*}
$$

Introducing the commutators in the (easier) conservative form

$$
R^{\epsilon}:=\left(D_{x} \cdot(\boldsymbol{b} w)\right) * \rho_{\epsilon}-D_{x} \cdot\left(\boldsymbol{b} w^{\epsilon}\right)
$$

(here we set again $w^{\epsilon}:=w * \rho_{\epsilon}$ ) it suffices to show that $R^{\epsilon}=L^{\epsilon}-$ $w^{\epsilon} \operatorname{div} \boldsymbol{b}_{t}$, where

$$
L^{\epsilon}(t, x):=\int_{\mathbb{R}^{d}} w(t, z)\left(\boldsymbol{b}_{t}(x)-\boldsymbol{b}_{t}(z)\right) \cdot \nabla \rho_{\epsilon}(z-x) d z
$$

Indeed, for any test function $\varphi$, we have that $\left\langle R^{\epsilon}, \varphi\right\rangle$ is given by

$$
\begin{aligned}
& -\int_{I} \int w \boldsymbol{b} \cdot \nabla \rho_{\epsilon} * \varphi d y-\int_{I} \int \varphi \boldsymbol{b} \cdot \nabla \rho_{\epsilon} * w d x-\int_{I} \int w^{\epsilon} \varphi \operatorname{div} \boldsymbol{b}_{t} \\
= & -\int_{I} \iint w_{t}(y) \boldsymbol{b}_{t}(y) \cdot \nabla \rho_{\epsilon}(y-x) d x d y \\
& -\int_{I} \iint \boldsymbol{b}_{t}(x) \nabla \rho_{\epsilon}(x-y) w_{t}(y) \varphi(x) d y d x-\int_{I} \int w^{\epsilon} \varphi \operatorname{div} \boldsymbol{b}_{t} \\
= & \int_{I} \int L^{\epsilon} \varphi-\int_{I} \int w^{\epsilon} \operatorname{div} \boldsymbol{b}_{t}
\end{aligned}
$$

(in the last equality we used the fact that $\nabla \rho$ is odd).

Then, one uses the strong convergence of translations in $L^{p}$ and the strong convergence of the difference quotients (a property that characterizes functions in Sobolev spaces)

$$
\frac{u(x+\epsilon z)-u(x)}{\epsilon} \rightarrow \nabla u(x) z \quad \text { strongly in } L_{\mathrm{loc}}^{1} \text {, for } u \in W_{\mathrm{loc}}^{1,1}
$$

to obtain that $r^{\epsilon}$ strongly converge in $L_{\mathrm{loc}}^{1}\left(I \times \mathbb{R}^{d}\right)$ to

$$
-w(t, x) \int_{\mathbb{R}^{d}}\left\langle\nabla \boldsymbol{b}_{t}(x) y, \nabla \rho(y)\right\rangle d y-w(t, x) \operatorname{div} \boldsymbol{b}_{t}(x)
$$

The elementary identity

$$
\int_{\mathbb{R}^{d}} y_{i} \frac{\partial \rho}{\partial y_{j}} d y=-\delta_{i j}
$$

then shows that the limit is 0 (this can also be derived by the fact that, in any case, the limit of $r^{\epsilon}$ in the distribution sense should be 0 ).

In this context, given $\bar{\mu}=\rho \mathscr{L}^{d}$ with $\rho \in L^{1} \cap L^{\infty}$, the $\mathscr{L}$-Lagrangian flow starting from $\bar{\mu}$ (at time 0 ) is defined by the following two properties:
(a) $\boldsymbol{X}(\cdot, x)$ is absolutely continuous in $[0, T]$ and satisfies

$$
\boldsymbol{X}(t, x)=x+\int_{0}^{t} \boldsymbol{b}_{s}(\boldsymbol{X}(s, x)) d s \quad \forall t \in[0, T]
$$

for $\bar{\mu}$-a.e. $x$;
(b) $\boldsymbol{X}(t, \cdot)_{\#} \bar{\mu} \leq C \mathscr{L}^{d}$ for all $t \in[0, T]$, with $C$ independent of $t$.

Summing up what we obtained so far, the general theory provides us with the following existence and uniqueness result.

Theorem 29. [Existence and uniqueness of $\mathscr{L}$-Lagrangian flows] Let $\boldsymbol{b} \in L^{1}\left([0, T] ; W_{\mathrm{loc}}^{1,1}\left(\mathbb{R}^{d} ; \mathbb{R}^{d}\right)\right)$ be satisfying
(i) $\frac{|\boldsymbol{b}|}{1+|x|} \in L^{1}\left([0, T] ; L^{1}\left(\mathbb{R}^{d}\right)\right)+L^{1}\left([0, T] ; L^{\infty}\left(\mathbb{R}^{d}\right)\right)$;
(ii) $\left[\operatorname{div} \boldsymbol{b}_{t}\right]^{-} \in L^{1}\left([0, T] ; L^{\infty}\left(\mathbb{R}^{d}\right)\right)$.

Then the $\mathscr{L}$-Lagrangian flow relative to $\boldsymbol{b}$ exists and is unique. I

Proof. By the previous results, the comparison principle holds for the continuity equation relative to $\boldsymbol{b}$. Therefore the general theory
previously developed applies, and Theorem 15 provides uniqueness of the $\mathscr{L}$-Lagrangian flow. $\|$ As for the existence, still the general theory (Theorem 18) tells us that it can be achieved provided we are able to solve, within $\mathscr{L}$, the continuity equation

$$
\begin{equation*}
\frac{d}{d t} w+D_{x} \cdot(\boldsymbol{b} w)=0 \tag{26}
\end{equation*}
$$

for any nonnegative initial datum $w_{0} \in L^{1} \cap L^{\infty}$. IThe existence of these solutions can be immediately achieved by a smoothing argument: we approximate $\boldsymbol{b}$ in $L_{\text {loc }}^{1}$ by smooth $\boldsymbol{b}^{h}$ with a uniform bound in $L^{1}\left(L^{\infty}\right)$ for $\left[\operatorname{div} \boldsymbol{b}_{t}^{h}\right]^{-}$. This bound, in turn, provides a uniform lower bound on $J \boldsymbol{X}^{h}$ and finally a uniform upper bound on $w_{t}^{h}=\left(w_{0} / J \boldsymbol{X}_{t}^{h}\right) \circ\left(\boldsymbol{X}_{t}^{h}\right)^{-1}$, solving

$$
\frac{d}{d t} w^{h}+D_{x} \cdot\left(\boldsymbol{b}^{h} w^{h}\right)=0
$$

Therefore, any weak limit of $w^{h}$ solves (26). I

Notice also that, choosing for instance a Gaussian, we obtain that the $\mathscr{L}$ Lagrangian flow is well defined up to $\mathscr{L}^{d}$-negligible sets (and independent of $\bar{\mu} \ll \mathscr{L}^{d}$, thanks to Remark 16).I

It is interesting to compare our characterization of Lagrangian flows with the one given in [53]. Heuristically, while the Di Perna-Lions one is based on the semigroup of transformations $x \mapsto \boldsymbol{X}(t, x)$, our one is based on the properties of the map $x \mapsto \boldsymbol{X}(\cdot, x)$.

Remark 30. The definition of the flow in [53] is based on the following three properties:
(a) $\frac{\partial \boldsymbol{Y}}{\partial t}(t, s, x)=b(t, \boldsymbol{Y}(t, s, x))$ and $\boldsymbol{Y}(s, s, x)=x$ in the distribution sense in $(0, T) \times \mathbb{R}^{d}$;
(b) the image $\lambda_{t}$ of $\mathscr{L}^{d}$ under $\boldsymbol{Y}(t, s, \cdot)$ satisfies

$$
\frac{1}{C} \mathscr{L}^{d} \leq \lambda_{t} \leq C \mathscr{L}^{d} \quad \text { for some constant } C>0
$$

(c) for all $s, s^{\prime}, t \in[0, T]$ we have

$$
\boldsymbol{Y}\left(t, s, \boldsymbol{Y}\left(s, s^{\prime}, x\right)\right)=\boldsymbol{Y}\left(t, s^{\prime}, x\right) \quad \text { for } \mathscr{L}^{d} \text {-a.e. } x .
$$

Then, $\boldsymbol{Y}(t, s, x)$ corresponds, in our notation, to the flow $\boldsymbol{X}^{s}(t, x)$ starting at time $s$ (well defined even for $t<s$ if one has two-sided $L^{\infty}$ bounds on the divergence). IIn our setting condition (c) can be recovered as a consequence with the following argument: assume to fix the ideas that $s^{\prime} \leq s \leq T$ and define

$$
\tilde{\boldsymbol{X}}(t, x):= \begin{cases}\boldsymbol{X}^{s^{\prime}}(t, x) & \text { if } t \in\left[s^{\prime}, s\right] ; \\ \boldsymbol{X}^{s}\left(t, \boldsymbol{X}^{s^{\prime}}(s, x)\right) & \text { if } t \in[s, T]\end{cases}
$$

It is immediate to check that $\tilde{\boldsymbol{X}}(\cdot, x)$ is an integral solution of the ODE in $\left[s^{\prime}, T\right]$ for $\mathscr{L}^{d}$-a.e. $x$ and that $\tilde{\boldsymbol{X}}(t, \cdot)_{\#} \bar{\mu}$ is bounded by $C^{2} \mathscr{L}^{d}$. Then, Theorem 29 (with $s^{\prime}$ as initial time) gives $\tilde{\boldsymbol{X}}(\cdot, x)=\boldsymbol{X}\left(\cdot, s^{\prime}, x\right)$ in $\left[s^{\prime}, T\right]$ for $\mathscr{L}^{d}$-a.e. $x$, whence (c) follows.
I
Moreover, the stability Theorem 20 can be read in this context as follows.
Theorem 31. [Stability] Let $\boldsymbol{b}^{h}, \boldsymbol{b} \in L^{1}\left([0, T] ; W_{\mathrm{loc}}^{1,1}\left(\mathbb{R}^{d} ; \mathbb{R}^{d}\right)\right)$, let $\boldsymbol{X}^{h}, \boldsymbol{X}$ be the $\mathscr{L}$-Lagrangian flows relative to $\boldsymbol{b}^{h}$, $\boldsymbol{b}$, let $\bar{\mu}=\rho \mathscr{L}^{d} \in$ $\mathscr{M}_{+}\left(\mathbb{R}^{d}\right)$ and assume that
(i) $\boldsymbol{b}^{h} \rightarrow \boldsymbol{b}$ in $L_{\mathrm{loc}}^{1}\left((0, T) \times \mathbb{R}^{d}\right)$;
(ii) setting $\mu_{t}^{h}:=\boldsymbol{X}^{h}(t, \cdot)_{\#} \rho \mathscr{L}^{d}$, we have

$$
\sup _{h} \int_{0}^{T} \int_{\mathbb{R}^{d}} \frac{\Theta\left(\left|\boldsymbol{b}^{h}\right|\right)}{1+|x|} d \mu_{t}^{h} d t<+\infty
$$

for some function $\Theta: \mathbb{R}^{+} \rightarrow \mathbb{R}^{+}$having a more than linear growth at infinity;
(iii) $\left[\operatorname{div} \boldsymbol{b}_{t}^{h}\right]^{-}$is bounded in $L^{1}\left([0, T] ; L^{\infty}\left(\mathbb{R}^{d}\right)\right)$.

Then,

$$
\lim _{h \rightarrow \infty} \int_{\mathbb{R}^{d}} \max _{[0, T]}\left|\boldsymbol{X}^{h}(\cdot, x)-\boldsymbol{X}(\cdot, x)\right| \wedge \rho(x) d x=0
$$

Taking into account the uniform upper bounds on $\mu_{t}^{h}$, ensured by (iii), condition (ii) is easily seen to be fulfilled if one has the decomposition $\boldsymbol{b}^{h}=\boldsymbol{c}^{h}+\boldsymbol{d}^{h}$, with

$$
\frac{\boldsymbol{c}^{h}}{1+|x|} \text { converging strongly in } L^{1}\left([0, T] ; L^{1}\left(\mathbb{R}^{d}\right)\right)
$$

and $\boldsymbol{d}^{h} /(1+|x|)$ bounded in $L^{1}\left([0, T] ; L^{\infty}\left(\mathbb{R}^{d}\right)\right)$. $\|$
Finally, we conclude this section with the illustration of some recent results [64], [13], [14] that seem to be more specific of the Sobolev case, concerned with the "differentiability" w.r.t. to $x$ of the flow $\boldsymbol{X}(t, x)$. These results provide a sort of bridge with the standard Cauchy-Lipschitz calculus:

Theorem 32. There exist Borel maps $L_{t}: \mathbb{R}^{d} \rightarrow M^{d \times d}$ satisfying

$$
\lim _{h \rightarrow 0} \frac{\boldsymbol{X}(t, x+h)-\boldsymbol{X}(t, x)-h L_{t}(x)}{|h|}=0 \quad \text { locally in measure }
$$

for any $t \in[0, T] . I I f$, in addition, we assume that

$$
\int_{0}^{T} \int_{B_{R}}\left|\nabla \boldsymbol{b}_{t}\right| \ln \left(2+\left|\nabla \boldsymbol{b}_{t}\right|\right) d x d t<+\infty \quad \forall R>0
$$

then the flow has the following "local" Lipschitz property: for any $\epsilon>0$ there exists a Borel set $A$ with $\bar{\mu}\left(\mathbb{R}^{d} \backslash A\right)<\epsilon$ such that $\left.\boldsymbol{X}(\cdot, t)\right|_{A}$ is Lipschitz for any $t \in[0, T]$.

According to this result, $L$ can be thought as a (very) weak derivative of the flow $\boldsymbol{X}$.IIt is still not clear whether the local Lipschitz property holds in the $W_{\text {loc }}^{1,1}$ case, or in the $B V_{\text {loc }}$ case discussed in the next section.

## 5. Vector fields having a $B V$ spatial regularity

In this section we prove the renormalization Theorem 27 under the weaker assumption of a $B V$ dependence w.r.t. the spatial variables, but still assuming that

$$
\begin{equation*}
D \cdot \boldsymbol{b}_{t} \ll \mathscr{L}^{d} \quad \text { for } \mathscr{L}^{1} \text {-a.e. } t \in(0, T) . \tag{27}
\end{equation*}
$$

I
Theorem 33. Let $\boldsymbol{b} \in L_{\mathrm{loc}}^{1}\left((0, T) ; B V_{\mathrm{loc}}\left(\mathbb{R}^{d} ; \mathbb{R}^{d}\right)\right)$ be satisfying (27). Then any distributional solution $w \in L_{\text {loc }}^{\infty}\left((0, T) \times \mathbb{R}^{d}\right)$ of

$$
\frac{d}{d t} w+D_{x} \cdot(\boldsymbol{b} w)=c \in L_{\mathrm{loc}}^{1}\left((0, T) \times \mathbb{R}^{d}\right)
$$

is a renormalized solution.
I
We try to give reasonably detailed proof of this result, referring to the original paper [7] for minor details. Before doing that we set up some notation, denoting by $D \boldsymbol{b}_{t}=\nabla \boldsymbol{b}_{t} \mathscr{L}^{d}+D^{s} \boldsymbol{b}_{t}$ the Radon-Nikodym decomposition of $D b_{t}$ in absolutely continuous and singular part w.r.t. $\mathscr{L}^{d}$.IWe also introduce the measures $|D \boldsymbol{b}|$ and $\left|D^{s} \boldsymbol{b}\right|$ by integration w.r.t. the time variable, i.e.

$$
\begin{aligned}
\int \varphi(t, x) d|D \boldsymbol{b}| & :=\int_{0}^{T} \int_{\mathbb{R}^{d}} \varphi(t, x) d\left|D \boldsymbol{b}_{t}\right| d t \\
\int \varphi(t, x) d\left|D^{s} \boldsymbol{b}\right| & :=\int_{0}^{T} \int_{\mathbb{R}^{d}} \varphi(t, x) d\left|D^{s} \boldsymbol{b}_{t}\right| d t
\end{aligned}
$$

We shall also assume, by the locality of the arguments involved, that $\|w\|_{\infty} \leq 1$.

We are going to find two estimates on the commutators, quite sensitive to the choice of the convolution kernel, and then combine them in a (pointwise) kernel optimization argument. II

Step 1 (anisotropic estimate). Let us start from the expression
$r^{\epsilon}(t, x)=\int_{\mathbb{R}^{d}} w(t, x-\epsilon y) \frac{\left(\boldsymbol{b}_{t}(x-\epsilon y)-\boldsymbol{b}_{t}(x)\right) \cdot \nabla \rho(y)}{\epsilon} d y-\left(w \operatorname{div} \boldsymbol{b}_{t}\right) * \rho_{\epsilon}(x)$
of the commutators $(\boldsymbol{b} \cdot \nabla w) * \rho_{\epsilon}-\boldsymbol{b} \cdot\left(\nabla\left(w * \rho_{\epsilon}\right)\right)$ : since $\boldsymbol{b}_{t} \notin W^{1,1}$ we cannot use anymore the strong convergence of the difference quotients.ll However, for any function $u \in B V_{\text {loc }}$ and any $z \in \mathbb{R}^{d}$ we have a classical $L^{1}$ estimate on the difference quotients

$$
\int_{K}|u(x+z)-u(x)| d x \leq\left|D_{z} u\right|\left(K_{\epsilon}\right) \quad \text { for any } K \subset \mathbb{R}^{d} \text { compact, }
$$

where $D u=\left(D_{1} u, \ldots, D_{d} u\right)$ stands for the distributional derivative of $u$, $D_{z} u=\langle D u, z\rangle=\sum_{i} z_{i} D_{i} u$ denotes the component along $z$ of $D u$ and $K_{\epsilon}$ is the open $\epsilon$-neighbourhood of $K$. Its proof follows from an elementary smoothing and lower semicontinuity argument.

We notice that

$$
D_{z}\left\langle b_{t}, \nabla \rho(z)\right\rangle=\left\langle M_{t}(\cdot) z, \nabla \rho(z)\right\rangle|D b| \quad \forall z \in \mathbb{R}^{d}
$$

and therefore the $L^{1}$ estimate on difference quotients gives the anisotropic

## estimate

$$
\begin{equation*}
\limsup _{\epsilon \downarrow 0} \int_{K}\left|r^{\epsilon}\right| d x \leq \int_{K} \int_{\mathbb{R}^{d}}\left|\left\langle M_{t}(x) z, \nabla \rho(z)\right\rangle\right| d z d|D b|(t, x)+d\left|D^{a} b\right|(K) \tag{29}
\end{equation*}
$$

for any compact set $K \subset(0, T) \times \mathbb{R}^{d}$.
Step 2 (isotropic estimate). On the other hand, a different estimate of the commutators that reduces to the standard one when $b(t, \cdot) \in W_{\text {loc }}^{1,1}$ can be achieved as follows. Let us start from the case $d=1$ : if $\mu$ is a $\mathbb{R}^{m}$-valued measure in $\mathbb{R}$ with locally finite variation, then by Jensen's inequality the functions

$$
\hat{\mu}_{\epsilon}(t):=\frac{\mu([t, t+\epsilon])}{\epsilon}=\mu * \frac{\chi_{[-\epsilon, 0]}}{\epsilon}(t), \quad t \in \mathbb{R}
$$

satisfy

$$
\begin{equation*}
\int_{K}\left|\hat{\mu}_{\epsilon}\right| d t \leq|\mu|\left(K_{\epsilon}\right) \quad \text { for any compact set } K \subset \mathbb{R} \tag{30}
\end{equation*}
$$

where $K_{\epsilon}$ is again the open $\epsilon$ neighbourhood of $K \| \mathrm{A}$ density argument based on (30) then shows that $\hat{\mu}_{\epsilon}$ converge in $L_{\text {loc }}^{1}(\mathbb{R})$ to the density of $\mu$ with respect to $\mathscr{L}^{1}$ whenever $\mu \ll \mathscr{L}^{1}$. If $u \in B V_{\text {loc }}$ and $\epsilon>0$ we know that

$$
\frac{u(x+\epsilon)-u(x)}{\epsilon}=\frac{D u([x, x+\epsilon])}{\epsilon}=\frac{D^{a} u([x, x+\epsilon])}{\epsilon}+\frac{D^{s} u([x, x+\epsilon])}{\epsilon}
$$

for $\mathscr{L}^{1}$-a.e. $x$ (the exceptional set possibly depends on $\epsilon$ ).In this way we have canonically split the difference quotient of $u$ as the sum of two
functions, one strongly converging to $\nabla u$ in $L_{\text {loc }}^{1}$, and the other one having an $L^{1}$ norm on any compact set $K$ asymptotically smaller than $\left|D^{s} u\right|(K)$. I

If we fix the direction $z$ of the difference quotient, the slicing theory of $B V$ functions gives that this decomposition can be carried on also in $d$ dimensions, showing that the difference quotients

$$
\frac{\boldsymbol{b}_{t}(x+\epsilon z)-\boldsymbol{b}_{t}(x)}{\epsilon}
$$

can be canonically split into two parts, the first one strongly converging in $L_{\text {loc }}^{1}\left(\mathbb{R}^{d}\right)$ to $\nabla \boldsymbol{b}_{t}(x) z$, and the second one having an $L^{1}$ norm on $K$ asymptotically smaller than $\left|\left\langle D^{s} \boldsymbol{b}_{t}, z\right\rangle\right|(K)$.IThen, repeating the DiPernaLions argument and taking into account the error induced by the presence of the second part of the difference quotients, we get the isotropic estimate

$$
\begin{equation*}
\underset{\epsilon \downarrow 0}{\limsup } \int_{K}\left|r^{\epsilon}\right| d x \leq\left(\int_{K} \int_{\mathbb{R}^{d}}|z||\nabla \rho(z)| d z\right) d\left|D^{s} \boldsymbol{b}\right|(t, x) \tag{31}
\end{equation*}
$$

for any compact set $K \subset(0, T) \times \mathbb{R}^{d} . \|$
Step 3 (reduction to a pointwise optimization problem). Roughly speaking, the isotropic estimate is useful in the regions where the absolutely continuous part is the dominant one, so that $\left|D^{s} b\right|(K) \ll 1$, while the anisotropic one turns out to be useful in the regions where the dominant part is the singular one. ILet us see how the two estimates can be combined: coming back to the smoothing scheme, we have

$$
\begin{equation*}
\frac{d}{d t} \beta\left(w^{\epsilon}\right)+\boldsymbol{b} \cdot \nabla \beta\left(w^{\epsilon}\right)-\beta^{\prime}\left(w^{\epsilon}\right) c * \rho_{\epsilon}=\beta^{\prime}\left(w^{\epsilon}\right) r^{\epsilon} \tag{32}
\end{equation*}
$$

Let $L$ be the supremum of $\left|\beta^{\prime}\right|$ on $[-1,1]$. Then, since $K$ is an arbitrary compact set, (31) tells us that any limit measure $\nu$ of $\left|\beta^{\prime}\left(w^{\epsilon}\right) r^{\epsilon}\right| \mathscr{L}^{d}$ as $\epsilon \downarrow 0$ satisfies

$$
\nu \leq L I(\rho)\left|D^{s} \boldsymbol{b}\right| \quad \text { with } \quad I(\rho):=\int_{\mathbb{R}^{d}}|z \| \nabla \rho(z)| d z .
$$

and, in particular, is singular with respect to $\mathscr{L}^{d} . \|$ On the other hand, the estimate (29) tells also us that

$$
\nu \leq L \int_{\mathbb{R}^{d}}|\langle M .(\cdot) z, \nabla \rho(z)\rangle| d z|D \boldsymbol{b}| .
$$

The second estimate and the singularity of $\nu$ with respect to $\mathscr{L}^{d}$ give

$$
\begin{equation*}
\nu \leq L \int_{\mathbb{R}^{d}}|\langle M .(\cdot) z, \nabla \rho(z)\rangle| d z\left|D^{s} \boldsymbol{b}\right| . \tag{33}
\end{equation*}
$$

Notice that in this way we got rid of the potentially dangerous term $I(\rho)$ : in fact, we are going to choose very anisotropic kernels $\rho$ on which $I(\rho)$ can be arbitrarily large. IThe measure $\nu$ can of course depend on the choice of $\rho$, but (32) tells us that the "defect" measure

$$
\sigma:=\frac{d}{d t} \beta\left(w_{t}\right)+\boldsymbol{b} \cdot \nabla \beta\left(w_{t}\right)-c_{t} \beta^{\prime}\left(w_{t}\right),
$$

clearly independent of $\rho$, satisfies $|\sigma| \leq \nu$. IEventually we obtain

$$
\begin{equation*}
|\sigma| \leq L \Lambda(M .(\cdot), \rho)\left|D^{s} \boldsymbol{b}\right| \quad \text { with } \quad \Lambda(N, \rho):=\int_{\mathbb{R}^{d}}|\langle N z, \nabla \rho(z)\rangle| d z . \tag{34}
\end{equation*}
$$

For $(x, t)$ fixed, we are thus led to the minimum problem

$$
\begin{equation*}
G(N):=\inf \left\{\Lambda(N, \rho): \rho \in C_{c}^{\infty}\left(B_{1}\right), \rho \geq 0, \int_{\mathbb{R}^{d}} \rho=1\right\} \tag{35}
\end{equation*}
$$

with $N=M_{t}(x)$. Notice that (34) gives

$$
|\sigma| \leq L \inf _{\rho \in D} \Lambda(M .(\cdot), \rho)\left|D^{s} \boldsymbol{b}\right|
$$

for any countable set $D$ of kernels $\rho$, and the continuity of $\rho \mapsto \Lambda(N, \rho)$ w.r.t. the $W^{1,1}\left(B_{1}\right)$ norm and the separability of $W^{1,1}\left(B_{1}\right)$ give

$$
\begin{equation*}
|\sigma| \leq L G(M .(\cdot))\left|D^{s} \boldsymbol{b}\right| \tag{36}
\end{equation*}
$$

$\|$ Notice now that the assumption that $D \cdot \boldsymbol{b}_{t} \ll \mathscr{L}^{d}$ for $\mathscr{L}^{1}$-a.e. $t \in(0, T)$ gives

$$
\operatorname{trace} M_{t}(x)\left|D^{s} \boldsymbol{b}_{t}\right|=0 \quad \text { for } \mathscr{L}^{1} \text {-a.e. } t \in(0, T)
$$

$\|$ Hence, recalling the definition of $\left|D^{s} \boldsymbol{b}\right|$, the trace of $M_{t}(x)$ vanishes for $\left|D^{s} \boldsymbol{b}\right|$-a.e. $(t, x)$. Applying the following lemma, a courtesy of Alberti, and using (36) we obtain that $\sigma=0$, thus concluding the proof.

Lemma 34. [Alberti] For any $d \times d$ matrix $N$ the infimum in (35) is $\mid$ trace $N \mid$.
I

Proof. Notice first that the lower bound follows immediately by the identity

$$
\int_{\mathbb{R}^{d}}\langle N z, \nabla \rho(z)\rangle d z=\int_{\mathbb{R}^{d}}-\rho(z) \operatorname{div} N z+\operatorname{div}(\rho(z) N z) d z=-\operatorname{trace} N
$$

Hence, we have to show only the upper bound. Again, by the identity

$$
\langle N z, \nabla \rho(z)\rangle=\operatorname{div}(N z \rho(z))-\operatorname{trace} N \rho(z)
$$

it suffices to show that for any $T>0$ there exists $\rho$ such that

$$
\begin{equation*}
\int_{\mathbb{R}^{d}}|\operatorname{div}(N z \rho(z))| d z \leq \frac{2}{T} \tag{37}
\end{equation*}
$$

The heuristic idea is (again...) to build $\rho$ as the superposition of elementary probability measures associated to the curves $e^{t N} x, 0 \leq t \leq T$, on which the divergence operator can be easily estimated. Given a smooth convolution kernel $\theta$ with compact support, it turns out that the function

$$
\begin{equation*}
\rho(z):=\frac{1}{T} \int_{0}^{T} \theta\left(e^{-t N} z\right) e^{-t \operatorname{traceN}} d t \tag{38}
\end{equation*}
$$

has the required properties (here $e^{t N} x=\sum_{i} t^{i} N^{i} x / i$ ! is the solution of the ODE $\dot{\gamma}=N \gamma$ with the initial condition $\gamma(0)=x)$.I Indeed, it is immediate to check that $\rho$ is smooth and compactly supported. To estimate the divergence of $N z \rho(z)$, we notice that $\rho=\int \theta(x) \mu_{x} d x$, where $\mu_{x}$ are the probability 1 -dimensional measures concentrated on the image of the curves $t \mapsto e^{t N} x$ defined by

$$
\mu_{x}:=\left(e^{\cdot N} x\right)_{\#}\left(\frac{1}{T} \mathscr{L}_{\llcorner[0, T]}^{1}\right)
$$

Indeed, for any $\varphi \in C_{c}^{\infty}\left(\mathbb{R}^{d}\right)$ we have

$$
\begin{aligned}
\int_{\mathbb{R}^{d}} \theta(x)\left\langle\mu_{x}, \varphi\right\rangle d x & =\frac{1}{T} \int_{0}^{T} \int_{\mathbb{R}^{d}} \theta(x) \varphi\left(e^{t N} x\right) d x d t \\
& =\frac{1}{T} \int_{0}^{T} \int_{\mathbb{R}^{d}} \theta\left(e^{-t N} y\right) e^{-t \text { traceN }} \varphi(y) d y d t \\
& =\int_{\mathbb{R}^{d}} \rho(y) \varphi(y) d y
\end{aligned}
$$

By the linearity of the divergence operator, it suffices to check that

$$
\left|D_{z} \cdot\left(N z \mu_{x}\right)\right|\left(\mathbb{R}^{d}\right) \leq \frac{2}{T} \quad \forall x \in \mathbb{R}^{d}
$$

But this is elementary, since
$\int_{\mathbb{R}^{d}}\langle N z, \nabla \varphi(z)\rangle d \mu_{x}(z)=\frac{1}{T} \int_{0}^{T}\left\langle N e^{t N} x, \nabla \varphi\left(e^{t N} x\right)\right\rangle d t=\frac{\varphi\left(e^{T N} x\right)-\varphi(x)}{T}$
for any $\varphi \in C_{c}^{\infty}\left(\mathbb{R}^{d}\right)$, so that $T D_{z} \cdot\left(N z \mu_{x}\right)=\delta_{x}-\delta_{e^{T N} x}$.
The original argument in [7] was slightly different and used, instead of Lemma 34, a much deeper result, still due to Alberti, saying that for a $B V_{\text {loc }}$ function $u: \mathbb{R}^{d} \rightarrow \mathbb{R}^{m}$ the matrix $M(x)$ in the polar decomposition $D u=M|D u|$ has rank 1 for $\left|D^{s} u\right|$-a.e. $x$, i.e. there exist unit vectors $\xi(x) \in \mathbb{R}^{d}$ and $\eta(x) \in \mathbb{R}^{m}$ such that $M(x) z=\eta(x)\langle z, \xi(x)\rangle$. IIn this case the asymptotically optimal kernels are much easier to build, by mollifying in the $\xi$ direction much faster than in all other ones.IThis is precisely what Bouchut and Lions did in some particular cases (respectively "Hamiltonian" vector fields and piecewise Sobolev ones).

As in the Sobolev case we can now obtain from the general theory given in Section 3 existence and uniqueness of $\mathscr{L}$-Lagrangian flows, with $\mathscr{L}=L^{\infty}\left(L^{1}\right) \cap L^{\infty}\left(L^{\infty}\right)$ : we just replace in the statement of Theorem 29 the assumption $\boldsymbol{b} \in L^{1}\left([0, T] ; W_{\operatorname{loc}}^{1,1}\left(\mathbb{R}^{d} ; \mathbb{R}^{d}\right)\right)$ with $\boldsymbol{b} \in$ $L^{1}\left([0, T] ; B V_{\text {loc }}\left(\mathbb{R}^{d} ; \mathbb{R}^{d}\right)\right)$, assuming as usual that $D \cdot \boldsymbol{b}_{t} \ll \mathscr{L}^{d}$ for $\mathscr{L}^{1}$ a.e. $t \in[0, T]$. $\|$

Analogously, with the same replacements in Theorem 31 (for $\boldsymbol{b}$ and $\boldsymbol{b}^{h}$ ) we obtain stability of $\mathscr{L}$-Lagrangian flows.

## 6. Applications

6.1. A system of conservation laws.ILet us consider the Cauchy problem (studied in one space dimension by Keyfitz-Kranzer in [63])

$$
\begin{equation*}
\frac{d}{d t} u+\sum_{i=1}^{d} \frac{\partial}{\partial x_{i}}\left(\boldsymbol{f}_{i}(|u|) u\right)=0, \quad u: \mathbb{R}^{d} \times(0,+\infty) \rightarrow \mathbb{R}^{k} \tag{39}
\end{equation*}
$$

with the initial condition $u(\cdot, 0)=\bar{u}$. Here $\boldsymbol{f}: \mathbb{R} \rightarrow \mathbb{R}^{d}$ is a $C^{1}$ function.I In a recent paper [32] Bressan showed that the problem can be ill-posed for $L^{\infty}$ initial data and he conjectured that it could be well posed for $B V$ initial data, suggesting to extend to this case the classical method of characteristics.II In [8] we proved that this procedure can really be implemented, thanks to the results in [7], for initial data $\bar{u}$ such that $\bar{\rho}:=|\bar{u}| \in B V \cap L^{\infty}$, with $1 /|\bar{u}| \in L^{\infty}$.I Later on, in a joint work with Bouchut and De Lellis [10], we proved that the lower bound on $\bar{\rho}$ is not necessary and, moreover, we proved that the solution built in [8] is unique in a suitable class of admissible functions: those whose modulus $\rho$ satisfies the scalar PDE

$$
\begin{equation*}
\frac{d}{d t} \rho+\sum_{i=1}^{d} \frac{\partial}{\partial x_{i}}\left(\boldsymbol{f}_{i}(\rho) \rho\right)=0 \tag{40}
\end{equation*}
$$

in the Kruzhkov sense (i.e. $\eta(\rho)_{t}+D_{x} \cdot(\boldsymbol{q}(\rho)) \leq 0$ for any convex entropy-entropy flux pair $(\eta, \boldsymbol{q})$, here $\left.s \boldsymbol{f}^{\prime}(s) \eta^{\prime}(s)=\boldsymbol{q}^{\prime}(s)\right)$, with the initial condition $\rho(0, \cdot)=\bar{\rho} . \|$

Notice that the regularity theory for this class of solutions gives that $\rho \in$ $L^{\infty} \cap B V_{\text {loc }}\left([0,+\infty) \times \mathbb{R}^{d}\right)$, due to the $B V$ regularity and the boundedness of $|\bar{u}|$. Furthermore the maximum principle gives $0<1 / \rho \leq 1 /|\bar{u}| \in L^{\infty}$.l

In order to obtain the (or, better, a) solution $u$ we can formally decouple the system, writing

$$
u=\theta \rho, \quad \bar{u}=\bar{\theta} \bar{\rho}, \quad|\theta|=|\bar{\theta}|=1,
$$

thus reducing the problem to the system (decoupled, if one neglects the constraint $|\theta|=1$ ) of transport equations

$$
\begin{equation*}
\theta_{t}+\sum_{i=1}^{d} \frac{\partial}{\partial x_{i}}\left(\boldsymbol{f}_{i}(\rho) \theta\right)=0 \tag{41}
\end{equation*}
$$

with the initial condition $\theta(0, \cdot)=\bar{\theta} . \|$
A formal solution of the system, satisfying also the constraint $|\theta|=1$, is given by

$$
\theta(t, x):=\bar{\theta}\left([\boldsymbol{X}(t, \cdot)]^{-1}(x)\right),
$$

where $\boldsymbol{X}(t, \cdot)$ is the flow associated to $\boldsymbol{f}(\rho) . \|$ Notice that the nonautonomous vector field $\boldsymbol{f}(\rho)$ is bounded and of class $B V_{\text {loc }}$, but the theory illustrated in these lectures is not immediately applicable because its divergence is not absolutely continuous with respect to $\mathscr{L}^{d+1}$.In this case, however, a simple argument still allows the use of the theory, representing $\boldsymbol{f}(\rho)$ as a part of the autonomous vector field $\boldsymbol{b}:=(\rho, \rho \boldsymbol{f}(\rho))$ in $\mathbb{R}^{+} \times \mathbb{R}^{d}$. This new vector field is still $B V_{\text {loc }}$ and bounded, and it is divergence-free due to (40).II

At this point, it is not hard to see that the reparameterization of the flow $(t(s), \boldsymbol{x}(s))$ associated to $\boldsymbol{b}$

$$
(\dot{t}(s), \dot{\boldsymbol{x}}(s))=(\rho(t(s), \boldsymbol{x}(s)), \boldsymbol{f}(\rho(t(s), \boldsymbol{x}(s))) \rho(t(s), \boldsymbol{x}(s)))
$$

defined by $\tilde{\boldsymbol{x}}(t)=\boldsymbol{x}\left(t(s)^{-1}(t)\right)$ (and here we use the assumption $\rho>0$ ) defines a flow for the vector field $\boldsymbol{f}(\rho)$ we were originally interested to.ll

In this way we get a kind of formal, or pointwise, solution of the system (40), that could indeed be very far from being a distributional solution.

But here comes into play the stability theorem, showing that all formal computations above can be justified just assuming first ( $\rho, \boldsymbol{f}(\rho)$ ) smooth, and then by approximation (see [8] for details).
6.2. Lagrangian solutions of semi-geostrophic equations.I The semigeostrophic equations are a simple model of the atmosphere/ocean flows [45], described by the system of transport equations
(SGE)

$$
\left\{\begin{array}{l}
\frac{d}{d t} \partial_{2} p+\boldsymbol{u} \cdot \nabla \partial_{2} p=-\boldsymbol{u}_{2}+\partial_{1} p \\
\frac{d}{d t} \partial_{1} p+\boldsymbol{u} \cdot \nabla \partial_{1} p=-\boldsymbol{u}_{1}-\partial_{2} p \\
\frac{d}{d t} \partial_{3} p+\boldsymbol{u} \cdot \nabla \partial_{3} p=0 .
\end{array}\right.
$$

Here $\boldsymbol{u}$, the velocity, is a divergence-free field, $p$ is the pressure and $\rho:=-\partial_{3} p$ represents the density of the fluid. We consider the problem in $[0, T] \times \Omega$, with $\Omega$ bounded and convex. Initial conditions are given on the pressure and a no-flux condition through $\partial \Omega$ is imposed for all times.ll

Introducing the modified pressure $P_{t}(x):=p_{t}(x)+\left(x_{1}^{2}+x_{2}^{2}\right) / 2$, (SGE) can be written in a more compact form as

$$
\frac{d}{d t} \nabla P+\boldsymbol{u} \cdot \nabla P=J(\nabla P-x) \quad \text { with } \quad J:=\left(\begin{array}{ccc}
0 & -1 & 0  \tag{42}\\
1 & 0 & 0 \\
0 & 0 & 0
\end{array}\right) .
$$

## I

Existence (and uniqueness) of solutions are still open for this problem. In [20] and [46], existence results have been obtained in the so-called dual coordinates, where we replace the physical variable $x$ by $X=\nabla P_{t}(x)$. Under this change of variables, and assuming $P_{t}$ to be convex, the system becomes

$$
\begin{equation*}
\frac{d}{d t} \alpha_{t}+D_{x} \cdot\left(\boldsymbol{U}_{t} \alpha_{t}\right)=0 \quad \text { with } \quad \boldsymbol{U}_{t}(X):=J\left(X-\nabla P_{t}^{*}(X)\right) \tag{43}
\end{equation*}
$$

with $\alpha_{t}:=\left(\nabla P_{t}\right)_{\#}\left(\mathscr{L}_{\Omega}\right)$ (here we denote by $\mathscr{L}_{\Omega}$ the restriction of $\mathscr{L}^{d}$ to $\Omega$ ). II Indeed, for any test function $\varphi$ we can use the fact that $\boldsymbol{u}$ is divergence-free to obtain:

$$
\begin{aligned}
\frac{d}{d t} \int_{\mathbb{R}^{d}} \varphi d \alpha_{t} & =\int_{\mathbb{R}^{d}} \nabla \varphi\left(\nabla P_{t}\right) \cdot \frac{d}{d t} \nabla P_{t} d x \\
& =\int_{\mathbb{R}^{d}} \nabla \varphi\left(\nabla P_{t}\right) \cdot J\left(\nabla P_{t}-x\right) d x+\int_{\mathbb{R}^{d}} \nabla \varphi\left(\nabla P_{t}\right) \nabla^{2} P_{t} \cdot \boldsymbol{u} d x \\
& =\int_{\mathbb{R}^{d}} \nabla \varphi \cdot J\left(X-\nabla P_{t}^{*}\right) d \alpha_{t}+\int_{\mathbb{R}^{d}} \nabla\left(\varphi \circ \nabla P_{t}\right) \cdot \boldsymbol{u} d x \\
& =\int_{\mathbb{R}^{d}} \nabla \varphi \cdot \boldsymbol{U}_{t} d \alpha_{t}
\end{aligned}
$$

Existence of a solution to (43) can be obtained by a suitable time discretization scheme.INow the question is: can we go back to the original physical variables ?!An important step forward has been achieved by Cullen and Feldman in [47], with the concept of Lagrangian solution of (SGE).II
Taking into account that the vector field $\boldsymbol{U}_{t}(X)=J\left(X-\nabla P_{t}^{*}(X)\right)$ is $B V$, bounded and divergence-free, there is a well defined, stable and measure preserving flow $\boldsymbol{X}(t, X)=\boldsymbol{X}_{t}(X)$ relative to $\boldsymbol{U}$. IThis flow can be carried back to the physical space with the transformation

$$
F_{t}(x):=\nabla P_{t}^{*} \circ \boldsymbol{X}_{t} \circ \nabla P_{0}(x),
$$

thus defining maps $F_{t}$ preserving $\mathscr{L}_{\Omega}^{d}$.I

Using the stability theorem can also show that $Z_{t}(x):=\nabla P_{t}\left(F_{t}(x)\right)$ solve, in the distributions sense, the Lagrangian form of (42), i.e.

$$
\begin{equation*}
\frac{d}{d t} Z_{t}(x)=J\left(Z_{t}-F_{t}\right) \tag{44}
\end{equation*}
$$

This provides us with a sort of weak solution of (42), and it is still an open problem how the Eulerian form could be recovered (see Section 7).

## 7. Open problems, bibliographical notes, and references

Section 2. The material contained in this section is classical. Good references are [56], Chapter 8 of [12], [29] and [53]. For the proof of the area formula, see for instance [6], [55], [60].

The proof of the second local variant, under the stronger assumption $\int_{0}^{T} \int_{\mathbb{R}^{d}}\left|\boldsymbol{b}_{t}\right| d \mu_{t} d t<+\infty$, is given in Proposition 8.1.8 of [12]. The same proof works under the weaker assumption (6).

Section 3. Many ideas of this section, and in particular the idea of looking at measures in the space of continuous maps to characterize the flow and prove its stability, are borrowed from [7], dealing with $B V$ vector fields. Later on, the arguments have been put in a more general form, independent of the specific class of vector fields under consideration, in [9]. Here we present a more refined version of [9].

The idea of a probabilistic representation is of course classical, and appears in many contexts (particularly for equations of diffusion type); to my knoledge the first reference in the context of conservation laws and fluid mechanics is [24], where a similar approach is proposed for the incompressible Euler equation (see also [25], [26], [27]): in this case the compact (but neither metrizable, nor separable) space $X^{[0, T]}$, with $X \subset \mathbb{R}^{d}$ compact, has been considered.

This approach is by now a familiar one also in optimal transport theory, where transport maps and transference plans can be thought in a natural way as measures in the space of minimizing geodesics [76], and in the so called irrigation problems, a nice variant of the optimal transport problem [22]. See also [18] for a similar approach within Mather's theory. The

Lecture Notes [84] (see also the Appendix of [69]) contain, among several other things, a comprehensive treatment of the topic of measures in the space of action-minimizing curves, including at the same time the optimal transport and the dynamical systems case (this unified treatment was inspired by [21]). Another related reference is [50].

The superposition principle is proved, under the weaker assumption $\int_{0}^{T} \int_{\mathbb{R}^{d}}\left|\boldsymbol{b}_{t}\right|^{p} d \mu_{t} d t<+\infty$ for some $p>1$, in Theorem 8.2.1 of [12], see also [70] for the extension to the case $p=1$ and to the non-homogeneous continuity equation. Very closely related results, relative to the representation of a vector field as the superposition of "elementary" vector fields associated to curves, appear in [77], [18].

In [16] an interesting variant of the stability Theorems 20 and 31 is discussed, peculiar of the case when the limit vector field $\boldsymbol{b}$ is a sufficiently regular gradient. In this case it has been proved in [16] that narrow convergence of $\mu_{t}^{h}$ to $\mu_{t}$ for all $t \in[0, T]$ and the energy estimate

$$
\limsup _{h \rightarrow \infty} \int_{0}^{T} \int_{\mathbb{R}^{d}}\left|\boldsymbol{b}_{t}^{h}\right|^{2} d \mu_{t}^{h} d t \leq \int_{0}^{T} \int_{\mathbb{R}^{d}}\left|\boldsymbol{b}_{t}\right|^{2} d \mu_{t} d t<+\infty
$$

are sufficient to obtain the stability property. This is due to the fact that, given $\mu_{t}$, gradient vector fields minimize $\int_{0}^{T} \int\left|\boldsymbol{c}_{t}\right|^{2} d \mu_{t}$ among all velocity fields $\boldsymbol{c}_{t}$ for which the continuity equation $\frac{d}{d t} \mu_{t}+D_{x} \cdot\left(\boldsymbol{c}_{t} \mu_{t}\right)=0$ holds (see Chapter 8 of [12] for a general proof of this fact, and for references to earlier works of Otto, Benamou-Brenier).

The convergence result in [16] can be used to answer positively a question raised in [59], concerning the convergence of the implicit Euler scheme

$$
\boldsymbol{u}_{k+1} \in \operatorname{Argmin}\left[\frac{1}{2 h} \int_{\Omega}\left|\boldsymbol{u}-\boldsymbol{u}_{k}\right|^{2}+\int_{\Omega} F(\nabla \boldsymbol{u}) d x\right]
$$

(here $\Omega, \Omega^{\prime}$ are bounded open in $\mathbb{R}^{d}$ and $\boldsymbol{u}: \Omega \rightarrow \Omega^{\prime}$ ) in the case when $F(\nabla \boldsymbol{u})$ depends only, in a convex way, only on the determinant of $\nabla \boldsymbol{u}$. It turns out that, representing as in [59] $\boldsymbol{u}_{k}$ as the composition of $k$ optimal transport maps, $\boldsymbol{u}_{[t / h]}$ converge as $h \downarrow 0$ to the solution $\boldsymbol{u}_{t}$ of

$$
\frac{d}{d t} \boldsymbol{u}_{t}=\operatorname{div}\left(\nabla F\left(\nabla \boldsymbol{u}_{t}\right)\right)
$$

built in [59] by purely differential methods (coupling a nonlinear diffusion equation for the measures $\beta_{t}:=\left(\boldsymbol{u}_{t}\right)_{\#}\left(\mathscr{L}_{\Omega}\right)$ in $\Omega^{\prime}$ to a transport equation for $\boldsymbol{u}_{t}^{-1}$ ). Existence of solutions (via differential or variational methods) for wider classes of energy densities $F$ is a largely open problem.

Section 4. The definition of renormalized solution and the strong convergence of commutators are entirely borrowed from [53]. See also [54] for the relevance of this concept in connection with the existence theory for Boltzmann equation. The proof of the comparison principle assuming only an $L^{1}\left(L_{\text {loc }}^{1}\right)$ bound (instead of an $L^{1}\left(L^{\infty}\right)$ one, as in [53], [7]) on the divergence was suggested to me by G.Savaré. The differentiability properties of the flow have been found in [64]: later on, this differentiability property has been characterized and compared with the more classical approximate differentiability [60] in [14], while [13] contains the proof of the stronger "local" Lipschitz properties. Theorem 32 summarizes all these results. The paper [44] contains also more explicit Lipschitz estimates and an independent proof of the compactness of flows. See also [37] for a proof, using radial convolution kernels, of the renormalization property for vector fields satisfying $D_{i} \boldsymbol{b}^{j}+D_{j} \boldsymbol{b}^{i} \in L_{\mathrm{loc}}^{1}$.

Both methods, the one illustrated in these notes and the DiPerna-Lions one, are based on abstract compactness arguments and do not provide a rate of convergence in the stability theorem. It would be interesting to
find an explicit rate of convergence (in mean with respect to $x$ ) of the trajectories. This problem is open even for autonomous, bounded and Sobolev (but not Lipschitz) vector fields.

No general existence result for Sobolev (or even $B V$ ) vector fields seems to be known in the infinite-dimensional case: the only reference we are aware of is [23]. Also the investigation of non-Euclidean geometries, e.g. Carnot groups and horizontal vector fields, could provide interesting results.

Finally, notice that the theory has a natural invariance, namely if $\boldsymbol{X}$ is a flow relative to $\boldsymbol{b}$, then $\boldsymbol{X}$ is a flow relative to $\tilde{\boldsymbol{b}}$ whenever $\{\tilde{\boldsymbol{b}} \neq \boldsymbol{b}\}$ is $\mathscr{L}^{1+d}$-negligible in $(0, T) \times \mathbb{R}^{d}$. So a natural question is whether the uniqueness "in the selection sense" might be enforced by choosing a canonical representative $\tilde{b}$ in the equivalence class of $b$ : in other words we may think that, for a suitable choice of $\tilde{\boldsymbol{b}}$, the ODE $\dot{\gamma}(x)=\tilde{\boldsymbol{b}}_{t}(\gamma(t))$ has a unique absolutely continuous solution starting from $x$ for $\mathscr{L}^{d}$-a.e. $x$.

Section 5. Here we followed closely [7]. The main idea of this section, i.e. the adaptation of the convolution kernel to the local behaviour of the vector field, has been used at various level of generality in [30], [66], [41] (see also [38], [39] for related results independent of this technique), until the general result [7].

The optimal regularity condition on $\boldsymbol{b}$ ensuring the renormalization property, and therefore the validity of the comparison principle in $\mathscr{L}_{b}$, is still not known. New results, both in the Sobolev and in the $B V$ framework, are presented in [11], [64], [65].
In [15] we investigate in particular the possibility to prove the renormalization property for nearly incompressible $B V_{\text {loc }} \cap L^{\infty}$ fields $\boldsymbol{b}$ : they are defined by the property that there exists a positive function $\rho$, with $\ln \rho \in L^{\infty}$, such that the space-time field $(\rho, \rho \boldsymbol{b})$ is divergence free. As in the case of
the Keyfitz-Kranzer system, the existence a function $\rho$ with this property seems to be a natural replacement of the condition $D_{x} \cdot \boldsymbol{b} \in L^{\infty}$ (and is actually implied by it ); as explained in [10], a proof of the renormalization property in this context would lead to a proof of a conjecture, due to Bressan, on the compactness of flows associated to a sequence of vector fields bounded in $B V_{t, x}$.

Section 6. In connection with the Keyfitz-Kranzer system there are several open questions: in particular one would like to obtain uniqueness (and stability) of the solution in more general classes of admissible functions (partial results in this direction are given in [10]). A strictly related problem is the convergence of the vanishing viscosity method to the solution built in [8]. Also, very little about the regularity of solutions is presently known: we know [49] that $B V$ estimates do not hold and, besides, that the contruction in [8] seems not applicable to more general systems of triangular type, see the counterexample in [43].

In connection with the semi-geostrophic problem, the main problem is the existence of solutions in the physical variables, i.e. in the Eulerian form. A formal argument suggests that, given $P_{t}$, the velocity $\boldsymbol{u}$ should be defined by

$$
\partial_{t} \nabla P_{t}^{*}\left(\nabla P_{t}(x)\right)+\nabla^{2} P_{t}^{*}\left(\nabla P_{t}(x)\right) J\left(\nabla P_{t}(x)-x\right) .
$$

On the other hand, the a-priori regularity on $\nabla P_{t}$ (ensured by the convexity of $P_{t}$ ) is a $B V$ regularity, and it is still not clear how this formula could be rigorously justified. In this connection, an important intermediate step could be the proof of the $W^{1,1}$ regularity of the maps $\nabla P_{t}$ (see also [33], [34], [35], [36], [80], [81] for the regularity theory of optimal transport maps under regularity assumptions on the initial and final densities).

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