

COMN	UINICAT	IONS
Comm	IUNICAI	
	IN	
	PARTIAL	
DIE	EEDENTI	AL
	FENENTI	
	UATION	S
at here	test in this online edition	10
	Number 10	

Communications in Partial Differential Equations

ISSN: (Print) (Online) Journal homepage: https://www.tandfonline.com/loi/lpde20

A superposition principle for the inhomogeneous continuity equation with Hellinger–Kantorovich-regular coefficients

Kristian Bredies, Marcello Carioni & Silvio Fanzon

To cite this article: Kristian Bredies, Marcello Carioni & Silvio Fanzon (2022) A superposition principle for the inhomogeneous continuity equation with Hellinger–Kantorovich-regular coefficients, Communications in Partial Differential Equations, 47:10, 2023-2069, DOI: 10.1080/03605302.2022.2109172

To link to this article: https://doi.org/10.1080/03605302.2022.2109172

9

© 2022 The Author(s). Published with license by Taylor and Francis Group, LLC



Published online: 15 Sep 2022.

_	_
ſ	
L	67.
L	<u> </u>

Submit your article to this journal 🕝

Article views: 550

Q

View related articles 🗹



View Crossmark data 🗹



👌 OPEN ACCESS ! 🖲 🤇

Check for updates

A superposition principle for the inhomogeneous continuity equation with Hellinger–Kantorovich-regular coefficients

Kristian Bredies^a, Marcello Carioni^b, and Silvio Fanzon^a

^aInstitute of Mathematics and Scientific Computing, University of Graz, Graz, Austria; ^bDepartment of Applied Mathematics, University of Twente, Enschede, The Netherlands

ABSTRACT

We study measure-valued solutions of the inhomogeneous continuity equation $\partial_t \rho_t + \text{div} (v \rho_t) = q \rho_t$ where the coefficients v and q are of low regularity. A new superposition principle is proven for positive measure solutions and coefficients for which the recently-introduced dynamic Hellinger-Kantorovich energy is finite. This principle gives a decomposition of the solution into curves $t \mapsto h(t) \delta_{v(t)}$ that satisfy the characteristic system $\dot{\gamma}(t) = v(t, \gamma(t)), \ \dot{h}(t) = g(t, \gamma(t))h(t)$ in an appropriate sense. In particular, it provides a generalization of existing superposition principles to the low-regularity case of g where characteristics are not unique with respect to h. Two applications of this principle are presented. First, uniqueness of minimal total-variation solutions for the inhomogeneous continuity equation is obtained if characteristics are unique up to their possible vanishing time. Second, the extremal points of dynamic Hellinger-Kantorovich-type regularizers are characterized. Such regularizers arise, for example, in the context of dynamic inverse problems and dynamic optimal transport.

ARTICLE HISTORY

Received 26 February 2021 Accepted 31 July 2022

KEYWORDS

Continuity equation; dynamic inverse problems; Hellinger–Kantorovich energy; optimal transport regularization; superposition principle; uniqueness

2020 MATHEMATICS SUBJECT CLASSIFICATION 28A50; 35C15; 35F05; 65J20; 35L03

1. Introduction

The main objective of this article is to present a new superposition principle for positive measure solutions to the linear inhomogeneous continuity equation, assuming natural regularity on the velocity field and on the source term. Such assumptions are substantially weaker than what is currently available in the literature, as we will discuss below. To be more precise, given $\Omega \subset \mathbb{R}^d$ the closure of an open bounded domain, we consider narrowly continuous curves of positive measures $t \mapsto \rho_t$ in $\mathcal{M}^+(\Omega)$ solving

$$\partial_t \rho_t + \operatorname{div}(\nu \rho_t) = g \rho_t \qquad \text{in} \qquad (0, 1) \times \Omega$$
 (1)

in the sense of distributions, where $\nu : (0,1) \times \Omega \to \mathbb{R}^d$ is a velocity field satisfying no flux boundary conditions on $\partial\Omega$ and $g : (0,1) \times \Omega \to \mathbb{R}$ is a source term encoding the inhomogeneity of the equation. We assume that the coefficients ν and g are Hellinger-Kantorovich-regular, namely, they are Borel measurable and satisfy the bound

CONTACT Marcello Carioni 🖾 m.c.carioni@utwente.nl 🗈 Department of Applied Mathematics, University of Twente, P.O. Box 217, 7500 AE Enschede, The Netherlands

 $[\]ensuremath{\mathbb{C}}$ 2022 The Author(s). Published with license by Taylor and Francis Group, LLC

This is an Open Access article distributed under the terms of the Creative Commons Attribution License (http://creativecommons.org/licenses/ by/4.0/), which permits unrestricted use, distribution, and reproduction in any medium, provided the original work is properly cited.

$$\int_{0}^{1} \int_{\Omega} |v(t,x)|^{2} + |g(t,x)|^{2} \ d\rho_{t}(x) \ dt < \infty.$$
⁽²⁾

In the following we will clarify the role of (2) in connection to recent advancements in the theory of Unbalanced Optimal Transport. Our task is to provide a superposition principle for (1) that allows to represent any positive solution $t \mapsto \rho_t$ as a superposition of elementary solutions, that is, curves of measures of the form $t \mapsto h(t)\delta_{\gamma(t)}$, where the trajectories $\gamma : [0,1] \to \Omega$ and the weights $h : [0,1] \to [0,\infty)$ solve, in an appropriate sense, the system of characteristics for (1):

(*i*)
$$\dot{\gamma}(t) = \nu(t, \gamma(t))$$
 (*ii*) $\dot{h}(t) = g(t, \gamma(t))h(t)$ in (0, 1). (3)

Notice that (*i*) describes all possible elementary trajectories which follow the flow given by v, while (*ii*) encodes the lack of mass preservation for solutions to (1), due to the inhomogeneity. The precise statement of such superposition principle is given in Theorem 1.1 below. Subsequently we provide two applications of the superposition principle for (1). First we prove uniqueness for minimal norm solutions to (1) under the assumption of uniqueness for solutions to (3) up to their possible vanishing time (see Theorem 1.2); Second, we characterize extremal points of regularizers closely related to the energy at (2), and apply such result to sparsity for dynamic inverse problems regularized via unbalanced optimal transport (see Theorem 1.3).

Concerning relevant literature, we mention that the superposition principle for narrowly continuous curves of probability measures $t \mapsto \rho_t$ solving the homogeneous continuity equation

$$\partial_t \rho_t + \operatorname{div}(\nu \rho_t) = 0 \quad \text{in} \quad (0, 1) \times \Omega$$
(4)

is by now classical. It was first introduced in the Euclidean setting by Ambrosio in [1], where it was employed to investigate uniqueness and stability of Lagrangian flows in the context of DiPerna-Lions Theory [2]. Since then it has been applied to different tasks [3–7] and extended to various settings [8–11]. In [12] the velocity field v is assumed to satisfy

$$\int_0^1 \int_{\Omega} |v(t,x)|^2 \ d\rho_t(x) \ dt < \infty.$$
(5)

An elementary solution to (4) is of the form $t \mapsto \delta_{\gamma(t)}$ where $\gamma : [0, 1] \to \Omega$ is an absolutely continuous curve solving the characteristic equation (*i*) in (3). Due to the lack of regularity of ν , solutions to the initial value problem associated to (*i*) are not unique. Such non-uniqueness is reflected in the superposition formula, which in this case is achieved by constructing a probability measure σ on the set $\Gamma := C([0, 1]; \Omega)$. To be more precise, it can be shown that if $\rho_t \in \mathcal{M}^+(\Omega)$ is a narrowly continuous solution to (4) and ν satisfies (5), then there exists a measure $\sigma \in \mathcal{M}^+(\Gamma)$ concentrated on absolutely continuous curves satisfying (*i*), with the property that ρ_t can be represented by the pushforward of σ via the evaluation map $e_t(\gamma) := \gamma(t)$, that is,

$$\int_{\Omega} \varphi(x) \ d\rho_t(x) = \int_{\Gamma} \varphi(\gamma(t)) \ d\sigma(\gamma) \quad \text{for all} \quad \varphi \in C(\Omega), \quad t \in [0, 1].$$
(6)

We refer the reader to [12, Theorem 8.2.1] for a proof of (6) with $\Omega = \mathbb{R}^d$ and to [13, Theorem 7] for the case of Ω being the closure of a bounded domain.

A generalization of (14) for positive measure solutions to the inhomogeneous continuity Equation (1) in $\Omega = \mathbb{R}^d$ is presented in [10]. Specifically, the following is proven in [10, Theorem 4.1]: suppose that $\rho_t \in \mathcal{M}^+(\Omega)$ is a narrowly continuous solution to (1), that ν satisfies (5) and g is bounded; then there exists a representing measure $\sigma \in$ $\mathcal{M}^+(\Gamma \times \Omega)$, concentrated on pairs (γ, x) with γ absolutely continuous curve solving (*i*) in (3) with the initial condition $\gamma(0) = x$, and such that ρ_t is represented via the implicit formula

$$\int_{\Omega} \varphi(x) \ d\rho_t(x) = \int_{\Gamma \times \Omega} \varphi(\gamma(t)) \ d\sigma(\gamma, x) + \int_0^t \left(\int_{\Omega} \int_{\Gamma} \varphi(\gamma(t)) \ d\sigma_s^x(\gamma) g(s, x) \ d\rho_s(x) \right) \ ds,$$
(7)

for all $\varphi \in C(\Omega)$, where for fixed *t*, the family $\{\sigma_t^x\}_{x\in\Omega}$ is the disintegration of σ with respect to $(\tilde{e}_t)_{\#}\sigma \in \mathcal{M}^+(\Omega)$, with $\tilde{e}_t(\gamma, x) := \gamma(t)$. There are two main drawbacks with the superposition principle from [10]: First, the representation formula (7) is implicit; Second, the source term *g* is required to be bounded. Such assumption on *g* is substantial, as it implies uniqueness of solutions to (*ii*) in (3) along any trajectory. This fact essentially allows the author of [10] to construct the measure σ in (7) in the same way as the one in (6). Another limitation of [10] is that it is not possible to provide a representation via (7) for solutions with mass that is vanishing or generating from zero during the evolution (for an example, see Remark 4.6).

The main focus of this article is to obtain a superposition principle for (1) which overcomes the above mentioned limitations of [10]. Indeed we obtain an explicit representation formula for (1) that resembles (6). In addition, we remove the boundedness assumption on g, and we replace it by the growth condition (2). Removing such assumption on g is far from straightforward, as it requires a new functional analytic framework for constructing a representation measure σ . In fact, the low regularity of gimplies non-uniqueness for the initial value problem associated with (*ii*) in (3). This suggests that a measure σ representing a solution $t \mapsto \rho_t$ to (15) has to account for nonuniqueness both for the trajectories γ and the weights h. Therefore, σ cannot just be a measure on Γ , but rather on a space of pairs (γ , h), as discussed in Theorem 1.1 below.

We now discuss the coupling of the continuity equation at (1) with the energy at (2), which is at the center of recent important developments in the theory of Unbalanced Optimal Transport. The classical theory of Optimal Transport, in its Monge-Kantorovich formulation [16–18], concerns the problem of transporting mass from a probability measure into a target one, while minimizing a given cost. Benamou and Brenier [19] made the crucial observation that the classical formulation of optimal transport has a dynamic counterpart, which links the continuity Equation (4) with the energy at (5). More precisely they observed that it is possible to compute the optimal transport between two probability measures ρ_0 and ρ_1 by minimizing the dissipation at (5) among all the curves of probability measures $t \mapsto \rho_t$ and velocity fields v solving the continuity Equation (4) with initial and final conditions given by ρ_0 and ρ_1 respectively. Such dynamic formulation makes possible to endow the space of probability measures with a differentiable structure [12], bringing to light deep connections between optimal transport and functional analytic issues, such as the characterization of differential equations as gradient flows in spaces of measures [12, 20-25] or the derivation of sharp inequalities [15, 26-30]. Particularly in connection to applications, the assumption of mass preservation during the evolution is quite restrictive. Overcoming this limitation is at the core of the so-called unbalanced optimal transport theory. Among the various formulations, we highlight the one introduced in [31-33]. There, transporting a positive measure ρ_0 into a target one ρ_1 corresponds to minimize a weighted version of (2) among all curves of positive measures $t \mapsto \rho_t$ and fields v, g satisfying the inhomogeneous continuity Equation (1) with initial and final conditions given by ρ_0 and ρ_1 , respectively. The quantity at (2) takes the name of Wasserstein-Fisher-Rao or Hellinger-Kantorovich energy in the literature. Such an approach has been successfully employed in applications where mass preservation is violated [34-37]. In particular in [33] it is shown that the above minimization procedure induces a distance which is compatible with a differentiable structure on the space $\mathcal{M}^+(\Omega)$. This distance can also be derived from the dynamic formulation of the Logarithmic-Entropy Optimal Transport problem [33] or can be regarded as dissipation energy for a certain class of scalar reaction-diffusion equations [36].

We conclude this introduction by discussing in more details the superposition principle we propose for (1), as well as the applications provided in this article. The rest of the manuscript is organized as follows. In Section 2 we introduce basic notations, as well as presenting some results on continuity equations and optimal transport energies. In Section 3 we set the functional analytic framework needed in order to prove our principle. In particular we investigate properties superposition of the Hellinger-Kantorovich energy (2) when restricted to elementary solutions to (1). In Section 4 we provide a proof for the main result of this article, that is, the superposition principle in Theorem 1.1 below. Finally, in Sections 5, 6 we detail applications of the superposition principle to uniqueness for solutions to (1) and to sparsity for dynamic inverse problems with Hellinger-Kantorovich-type regularizers.

1.1. Main result

To obtain a superposition principle for (1) under the energy bound (2) we construct a positive measure σ on the set \mathscr{G}_{Ω} of narrowly continuous curves $t \mapsto \rho_t$ with values in $\mathscr{G}_{\Omega} := \{h\delta_{\gamma} \in \mathcal{M}(\Omega) : h \ge 0, \gamma \in \Omega\}$. We endow \mathscr{G}_{Ω} with the flat distance of measures and \mathscr{G}_{Ω} with the respective supremum distance. In this way \mathscr{G}_{Ω} becomes a separable metric space. Notice that \mathscr{G}_{Ω} plays the role of the set of continuous curves Γ in (6). As we will see, c.f. Remark 3.2, the construction of \mathscr{G}_{Ω} closely resembles the cone space introduced in [33, 36] to study absolutely continuous curves with respect to the Hellinger–Kantorovich distance. It is immediate to check that elements of \mathscr{G}_{Ω} can be represented by $\rho_t = h(t)\delta_{\gamma(t)}$, for some non-negative weight $h \in C[0, 1]$ and curve $\gamma \in C(\{h > 0\}; \Omega)$, where we set $\{h > 0\} := \{t \in [0, 1] : h(t) > 0\}$. Thus, the mass of the elements of \mathscr{G}_{Ω} is varying continuously in time and is allowed to vanish, reflecting the behavior of solutions to (1). The measure σ we construct is concentrated on elements $\rho_t = h(t)\delta_{\gamma(t)} \in \mathscr{G}_{\Omega}$, with h and γ solving the system of ODEs:

(*i*)
$$\dot{\gamma}(t) = \nu(t, \gamma(t))$$
 a.e. in $\{h > 0\}$ (*ii*) $h(t) = g(t, \gamma(t))h(t)$ a.e. in (0,1).
(8)

Notice that, in comparison to the system of characteristics at (3), we are restricting the first ODE to the set $\{h > 0\}$. Indeed, if h(t) = 0, then $\rho_t = 0$ and thus we lose any information on the trajectories for that time instant. The above observations are formalized in the following theorem, which is the main result of our article (c.f. Theorem 4.3).

Theorem 1.1. Let $\Omega \subset \mathbb{R}^d$ be the closure of an open bounded domain. Let $\rho_t : [0,1] \to \mathcal{M}^+(\Omega)$ be a narrowly continuous solution to (1) for some Borel measurable $v : (0,1) \times \Omega \to \mathbb{R}^d$, $g : (0,1) \times \Omega \to \mathbb{R}$ satisfying (2) and such that v has no flux on $\partial\Omega$. Then there exists a measure $\sigma \in \mathcal{M}^+(\mathscr{S}_{\Omega})$ concentrated on curves of measures $\rho_t = h(t)\delta_{\gamma(t)}$ with h, γ solving (8) and such that

$$\int_{\Omega} \varphi(x) \ d\rho_t(x) = \int_{\mathscr{S}_{\Omega}} h(t)\varphi(\gamma(t)) \ d\sigma(\gamma,h) \quad \text{for all} \quad \varphi \in C(\Omega), \quad t \in [0,1].$$
(9)

Conversely, assume that $\sigma \in \mathcal{M}^+(\mathscr{S}_{\Omega})$ is concentrated on solutions to (8) and satisfies

$$\int_{0}^{1} \int_{\mathscr{S}_{\Omega}} h(t) \left(1 + |\nu(t,\gamma(t))| + |g(t,\gamma(t))| \right) \, d\sigma(\gamma,h) \, dt < \infty.$$

$$\tag{10}$$

Then (9) defines a narrowly continuous curve of positive measures solving (1).

Notice that the growth condition (10) is natural, in the sense that if a measure σ represents ρ_t and (2) holds, then automatically σ satisfies (10). We refer the reader to Remark 4.4 below for more details. We also remark that the set Ω in Theorem 1.1 is required to be bounded. Indeed it would be interesting to extend our result to unbounded domains, in the spirit of [10, 12] where $\Omega = \mathbb{R}^d$ is considered. However, it seems that a different proof strategy or stronger assumptions are required, see Remark 4.7 below for details. Moreover, similarly to [10, 12], it should be possible to prove a version of Theorem 1.1 in which (2) is replaced by an L^p bound for $1 \le p \le \infty$. Such analysis falls outside the scope of our article.

The proof of Theorem 1.1 is presented in Section 4. It is based on a similar smoothing strategy as the one employed in [1] to prove (6). However in this case there are two main differences: first one needs to establish compactness properties for a coercive version of the Hellinger-Kantorovich energy when restricted to elements of \mathscr{C}_{Ω} , see Proposition 3.10; second the smoothing needs to take into account the possibility of the measure ρ_t vanishing at some time instance, as detailed in Remark 4.6 below.

1.2. Uniqueness of solutions to the continuity equation

In Section 5 we present the first application of the superposition principle of Theorem 1.1. Our aim is to show that uniqueness of solutions for the system of ODEs at (3), up to their possible vanishing time, implies uniqueness for measure solutions to the inhomogeneous continuity equation (1) satisfying the bound (2) and with minimal total variation. The key ingredient of the proof is formula (9), which allows to decompose any solution of (1) satisfying the bound (2) into a superposition

of elementary curves $t \mapsto h(t) \delta_{\gamma(t)}$ such that (γ, h) are solutions to (8). Such representation tation allows to link uniqueness for (3) with the one for (1). The main difference between our result and the classical one for the homogeneous continuity equation [14, Theorem 9] lies in the fact that elementary solutions $\rho_t = h(t)\delta_{\gamma(t)}$ are allowed to vanish in time. In this case uniqueness for (3) is not enough to ensure uniqueness of solutions to the inhomogeneous continuity equation. Indeed, when the mass of a solution vanishes at a given time instant $\overline{t} \in (0,1)$, the uniqueness assumption for (3) is not providing any information on the behavior of the solutions for $t > \overline{t}$: this is because the measure σ is concentrated on solutions to (8) where i) is only valid in the set $\{h > 0\}$. Therefore, in order to recover uniqueness for (1), we impose an extra constraint on the total variation of its solutions. More precisely, we show that solutions to (1) with minimal mass can be represented, invoking Theorem 1.1, by a measure σ concentrated on curves $t \mapsto h(t) \delta_{\gamma(t)}$ such that (γ, h) solves (8) and h is strictly positive in an interval $[0, \tau) \cap [0, 1]$ for some $\tau \in \mathbb{R}$. Such observation allows to employ uniqueness for the system of characteristics at (3), up to their possible vanishing time, to infer uniqueness for measure solutions to (1) with minimal total variation. We obtain the following theorem, c.f. Theorem 5.1.

Theorem 1.2. Let $v : (0,1) \times \Omega \to \mathbb{R}^d$, $g : (0,1) \times \Omega \to \mathbb{R}$ be Borel measurable functions and $A \subset \Omega$ be a Borel measurable set. Suppose that:

(Hyp) For each $x \in A$ the solution of the system of ODEs (8) with initial value (x, 1) is unique in $[0, \tau)$ for every $\tau \in (0, 1)$ such that $[0, \tau) \subset \{h > 0\}$.

Then, for any initial datum $\rho_0 \in \mathcal{M}^+(\Omega)$ concentrated on A, the inhomogeneous continuity equation (1) admits at most one positive narrowly continuous solution $t \mapsto \rho_t$ satisfying (2), with initial datum ρ_0 , and such that $||\rho||_{\mathcal{M}} \leq ||\tilde{\rho}||_{\mathcal{M}}$ for every $t \mapsto \tilde{\rho}_t$ positive narrowly continuous solution to (1) satisfying (2), and such that $\tilde{\rho}_0 = \rho_0$.

1.3. Extremal points of the Hellinger-Kantorovich energy

In the context of inverse problems, the knowledge of the structure of extremal points of the regularizer allows to numerically reconstruct sparse solutions, that is, solutions given by the superposition of finitely many extremal points [38,39]. It has been recently proposed [34] to regularize dynamic inverse problems via an energy related to the one at (2). To be more specific, the energy at (2) can be recast into a convex functional B_{δ} over the space $\mathcal{M} := \mathcal{M}((0,1) \times \Omega)^{d+2}$ defined by

$$B_{\delta}(\rho, m, \mu) := \frac{1}{2} \int_{0}^{1} \int_{\Omega} \left| \frac{dm}{d\rho} \right|^{2} + \delta^{2} \left| \frac{d\mu}{d\rho} \right|^{2} d\rho$$
(11)

if $\rho \ge 0$, $m, \mu \ll \rho$, and set to ∞ otherwise, where $\delta > 0$ is a parameter (c.f. Section 2.2). The regularizer studied in [34] consists in the energy at (11) to which the total variation of ρ is added, while enforcing the continuity equation constraint $\partial_t \rho + \text{div } m = \mu$. An analysis of the extremal points of such energy is currently missing in the literature: Therefore, in this article, we employ the superposition principle of Theorem 1.1 to characterize the extremal points of the set

$$\mathscr{B} = \{(\rho, m, \mu) : \partial_t \rho + \operatorname{div} m = \mu, \quad \beta B_\delta(\rho, m, \mu) + \alpha ||\rho||_{\mathcal{M}} \le 1\},$$
(12)

where $\alpha, \beta > 0$ are parameters. Notice that we do not impose boundary conditions in the continuity equation at (12). Moreover the total variation of ρ is added to the functional B_{δ} , in order to enforce coercivity, and thus compactness of \mathcal{B} . We prove the following result (c.f. Theorem 6.3).

Theorem 1.3. The extremal points of the set defined in (12) are exactly given by the zero measure (0, 0, 0) and the triples of measures (ρ, m, μ) such that $\rho = h(t) dt \otimes \delta_{\gamma(t)}$, $m = \dot{\gamma}(t)\rho$, $\mu = \dot{h}(t) dt \otimes \delta_{\gamma(t)}$ with the following properties:

- (a) $h, \sqrt{h} \in AC^2[0, 1], \gamma \in C(\{h > 0\}; \Omega) \text{ and } \sqrt{h}\gamma \in AC^2([0, 1]; \mathbb{R}^d),$
- (b) the set $\{h > 0\}$ is connected,
- (c) the energy satisfies $\beta B_{\delta}(\rho, m, \mu) + \alpha ||\rho||_{\mathcal{M}} = 1$.

In the above we denote by AC^2 the set of absolutely continuous functions with a.e. derivative in L^2 (see [12, Section 1.1] for a precise definition).

Theorem 1.3 is a generalization of the results obtained in [13], where the Benamou-Brenier energy with homogeneous continuity equation constraint is considered. In Section 6.2 we apply Theorem 1.3 to understand the structure of sparse solutions for dynamic inverse problems with unbalanced optimal transport regularization. In particular, we consider the inverse problem proposed in [34], where the minimization of the energy at (12) is coupled with a fidelity term penalizing the distance between ρ and some fixed observation. Applying recent results on sparsity [40,41] we show that the minimization problem in [34] admits a solution which is a finite linear combination of extremal points of \mathcal{B} , that is, of curves as described in Theorem 1.3.

2. Preliminaries

For measure theory notations and definitions we follow [42]. Given a metric space Y we denote by $\mathcal{M}(Y)$, $\mathcal{M}(Y; \mathbb{R}^d)$, $\mathcal{M}^+(Y)$ the spaces of bounded Borel measures, bounded vector Borel measures, bounded positive Borel measures on Y, respectively. Throughout the article, whenever we say that a set or a function is measurable, we always intend Borel measurable, that is, measurability with respect to the Borel σ -algebra. For a measure μ we denote its total variation measure by $|\mu|$. We say that a sequence of measures $\{\mu_n\}_n$ on Y converges narrowly to μ if $\int_Y \varphi(y) d\mu_n(y) \to \int_Y \varphi(y) d\mu(y)$ for all $\varphi \in C_b(Y)$, where $C_b(Y)$ denotes the set of real valued continuous and bounded functions on Y.

Let $\Omega \subset \mathbb{R}^d$ be the closure of a bounded domain, with $d \in \mathbb{N}$, $d \ge 1$, and define the time-space domain $X_{\Omega} := (0,1) \times \Omega$. We say that $\rho \in \mathcal{M}(X_{\Omega})$ disintegrates with respect to time if there exists a Borel family of measures $\{\rho_t\}_{t \in [0,1]} \subset \mathcal{M}(\Omega)$ such that $\int_{X_{\Omega}} \varphi \ d\rho = \int_0^1 \int_{\Omega} \varphi(t,x) \ d\rho_t(x) \ dt$ for all $\varphi \in L^1_\rho(X_{\Omega})$. The disintegration is denoted by $\rho = dt \otimes \rho_t$. Further, a curve of measures $t \in [0,1] \mapsto \rho_t \in \mathcal{M}(\Omega)$ is narrowly continuous if the map $t \mapsto \int_{\Omega} \varphi(x) \ d\rho_t(x)$ is continuous for each fixed $\varphi \in C(\Omega)$. The family of narrowly continuous, by the principle of uniform boundedness, it follows that $\rho := dt \otimes \rho_t$

belongs to $\mathcal{M}(X_{\Omega})$. We also introduce $C_{w}([0,1]; \mathcal{M}^{+}(\Omega))$ as the family of narrowly continuous curves with values into $\mathcal{M}^{+}(\Omega)$. The above definitions extend verbatim to the case $\Omega = \mathbb{R}^{d}$.

2.1. Continuity equation

Set $\mathcal{M}_{\Omega} := \mathcal{M}(X_{\Omega}) \times \mathcal{M}(X_{\Omega}; \mathbb{R}^d) \times \mathcal{M}(X_{\Omega})$. We say that the triple $(\rho, m, \mu) \in \mathcal{M}_{\Omega}$ solves the continuity equation

$$\partial_t \rho + \operatorname{div} m = \mu \quad \text{in} \quad X_\Omega,$$
 (13)

whenever (13) holds in the sense of distributions, that is,

$$\int_{X_{\Omega}} \partial_t \varphi \ d\rho + \int_{X_{\Omega}} \nabla \varphi \cdot dm + \int_{X_{\Omega}} \varphi \ d\mu = 0 \quad \text{for all} \quad \varphi \in C^{\infty}_c(X_{\Omega}).$$
(14)

Here, ρ represents a density, *m* a momentum field advecting ρ , while μ is a source term accounting for mass change. The above definition also holds for unbounded spatial domains, for example., $\Omega = \mathbb{R}^d$. Moreover the time interval (0, 1) can be replaced by (0, T) with T > 0. We remark that (14) includes no flux boundary conditions for *m* on $\partial \Omega$, and no initial conditions for ρ are prescribed. Moreover (14) can be equivalently tested with maps in $C_c^1(X_\Omega)$ [12, Remark 8.1.1]. The following lemma provides some properties of solutions to (14) which will be needed in the coming analysis. The statement holds both in bounded domains as well as in \mathbb{R}^d . For a proof in bounded domains see, for example, Propositions 2.2, 2.4 in [34], which can be easily generalized to \mathbb{R}^d .

Lemma 2.1. Assume that $(\rho, m, \mu) \in \mathcal{M}_{\Omega}$ satisfies (14) with $\rho \in \mathcal{M}^+(X_{\Omega})$. Then $\rho = dt \otimes \rho_t$, where $\rho_t \in \mathcal{M}^+(\Omega)$ for a.e. t in (0, 1). Moreover the map $t \mapsto \rho_t(\Omega)$ belongs to BV(0, 1), with distributional derivative given by $\pi_{\#}\mu$, where $\pi : X_{\Omega} \to (0, 1)$ is the projection on the time coordinate. If in addition $m = v\rho$, $\mu = g\rho$ for some measurable $v : X_{\Omega} \to \mathbb{R}^d$, $g : X_{\Omega} \to \mathbb{R}$ with

$$\int_0^1 \int_\Omega |\nu(t,x)| + |g(t,x)| \ d\rho_t(x) \ dt < \infty,$$

then there exists a unique curve $t \mapsto \tilde{\rho}_t$ in $C_w([0,1]; \mathcal{M}^+(\Omega))$ such that $\rho_t = \tilde{\rho}_t$ a.e. in (0, 1).

In the rest of the article we will identify ρ_t with its narrowly continuous representative $\tilde{\rho}_t$, whenever the assumptions of Lemma 2.1 hold.

2.2. Optimal transport energy

We now introduce the Wasserstein–Fisher-Rao energy, also known as the Hellinger–Kantorovich energy, as originally done in [31–33]. To this end, let $\delta > 0$ be a fixed parameter. Define the convex, one-homogeneous and lower semi-continuous map $\Psi_{\delta} : \mathbb{R} \times \mathbb{R}^d \times \mathbb{R} \to [0, \infty]$ by setting

COMMUNICATIONS IN PARTIAL DIFFERENTIAL EQUATIONS 👄 2031

$$\Psi_{\delta}(t, x, y) := \begin{cases} \frac{|x|^2 + \delta^2 y^2}{2t} & \text{if } t > 0, \\ 0 & \text{if } t = |x| = y = 0, \\ \infty & \text{otherwise,} \end{cases}$$
(15)

where $\infty y^2 = \infty$ for $y \neq 0$ and $\infty y^2 = 0$ for y = 0. The Wasserstein-Fisher-Rao energy is given by the map $B_{\delta} : \mathcal{M}_{\Omega} \to [0, \infty]$ defined by

$$B_{\delta}(\rho, m, \mu) := \int_{X_{\Omega}} \Psi_{\delta}\left(\frac{d\rho}{d\lambda}, \frac{dm}{d\lambda}, \frac{d\mu}{d\lambda}\right) d\lambda, \tag{16}$$

where $\lambda \in \mathcal{M}^+(X_{\Omega})$ is an arbitrary measure such that $\rho, m, \mu \ll \lambda$. Definition (16) does not depend on the choice of λ , as Ψ_{δ} is one-homogeneous. Properties of the energy B_{δ} which are relevant in the following analysis are summarized in Lemma A.4 (for a proof see [34, Proposition 2.6]). We now introduce a coercive version of B_{δ} : Set

$$\mathcal{D}_{\Omega} := \{ (\rho, m, \mu) \in \mathcal{M}_{\Omega} : \partial_t \rho + \text{div } m = \mu \text{ in the sense of } (14) \},\$$

and define the functional $J_{\alpha,\,\beta,\,\delta}:\mathcal{M}_\Omega\to[0,\infty]$ as

$$J_{\alpha,\beta,\delta}(\rho,m,\mu) := \begin{cases} \beta B_{\delta}(\rho,m,\mu) + \alpha \|\rho\|_{\mathcal{M}(X_{\Omega})} & \text{if } (\rho,m,\mu) \in \mathcal{D}_{\Omega}, \\ \infty & \text{otherwise,} \end{cases}$$
(17)

where $\alpha > 0$ and $\beta > 0$ are fixed constants. We remark that adding the total variation of ρ to B_{δ} enforces the balls of $J_{\alpha,\beta,\delta}$ to be compact in the weak* topology of \mathcal{M}_{Ω} . Such property, together with others, is the object of Lemma A.5. The content of Lemma A.5 is based on results proven in [34, Lemmas 4.5, 4.6].

2.3. Characteristics theory for the continuity equation

We start by recalling a classical result on the theory of ordinary differential equations in \mathbb{R}^d [12, Lemma 8.1.4].

Proposition 2.2. Let $v : [0,1] \times \mathbb{R}^d \to \mathbb{R}^d$ be measurable and such that

$$\int_{0}^{1} \sup_{x \in \mathbb{R}^{d}} |v(t,x)| + \operatorname{Lip}(v(t,\cdot), \mathbb{R}^{d}) \, dt < \infty.$$
(18)

Then for each $x \in \mathbb{R}^d$ the ODE

$$\dot{X}_x(t) = v(t, X_x(t))$$
 for a.e. $t \in (0, 1)$, $X_x(0) = x$, (19)

admits a unique absolutely continuous solution $t \mapsto X_x(t)$ defined for all $t \in [0, 1]$.

Next we provide a representation formula for measure solutions of the continuity equation (13). This is the analogue of [12, Lemma 8.1.6] for the inhomogeneous continuity equation, and a generalization of [10, Proposition 3.6] to the case of g unbounded.

Proposition 2.3. Let $v : [0,1] \times \mathbb{R}^d \to \mathbb{R}^d$, $g : [0,1] \times \mathbb{R}^d \to \mathbb{R}$ be measurable. Assume that

$$\int_{0}^{1} \sup_{x \in \mathbb{R}^{d}} |g(t,x)| + \operatorname{Lip}(g(t,\cdot), \mathbb{R}^{d}) \, dt < \infty$$
⁽²⁰⁾

and (18) hold. Let $\rho_0 \in \mathcal{M}^+(\mathbb{R}^d)$ and denote by $t \mapsto X_x(t)$ the unique solution to (19) defined for all $t \in [0, 1]$ and $x \in \mathbb{R}^d$. Then, the map

$$t \mapsto \rho_t := (X_{(\cdot)}(t))_{\#} \left(\rho_0 \ e^{\int_0^t g(s, X_{(\cdot)}(s)) ds} \right)$$

$$\tag{21}$$

is a narrowly continuous solution to the continuity equation $\partial_t \rho_t + \text{div}(v\rho_t) = g\rho_t$ in $(0,1) \times \mathbb{R}^d$ in the sense of (14), where the push-forward in (21) is with respect to the space variable.

Proof. Narrow continuity of $t \mapsto \rho_t$ follows immediately from (20), dominated convergence and the continuity of $t \mapsto X_x(t)$ for each x. Let now $\varphi \in C_c^1((0,1) \times \mathbb{R}^d)$. Then for ρ_0 -a.e. x in \mathbb{R}^d , the map $t \mapsto \varphi(t, X_x(t))$ is absolutely continuous in (0, 1), with a.e. derivative given by

$$\frac{d}{dt} \ \varphi(t, X_x(t)) = \partial_t \varphi(t, X_x(t)) + \nabla \varphi(t, X_x(t)) \cdot v(t, X_x(t)),$$
(22)

thanks to Proposition 2.2. By (20) we also have that $t \mapsto \varphi(t, X_x(t))e^{\int_0^t g(s, X_x(s))ds}$ is absolutely continuous in (0, 1), and for a.e. $t \in (0, 1)$ it holds

$$\frac{d}{dt}\left(\varphi(t,X_x(t))e^{\int_0^t g(s,X_x(s))ds}\right) = \left(\frac{d}{dt} \ \varphi(t,X_x(t)) + \varphi(t,X_x(t)) \ g(t,X_x(t))\right)e^{\int_0^t g(s,X_x(s))ds}.$$
(23)

In particular, it is immediate to check that

$$\int_{0}^{1} \int_{\mathbb{R}^{d}} \left| \frac{d}{dt} \left(\varphi(t, X_{x}(t)) e^{\int_{0}^{t} g(s, X_{x}(s)) ds} \right) \right| \ d\rho_{0}(x) \ dt \leq \|\varphi\|_{C^{1}} \ \rho_{0}(\mathbb{R}^{d}) \ e^{M_{g}} (1 + M_{\nu} + M_{g}),$$

where $M_v := \int_0^1 \sup_{x \in \mathbb{R}^d} |v(t,x)| dt$, $M_g := \int_0^1 \sup_{x \in \mathbb{R}^d} |g(t,x)| dt$, which are finite by (18, 20). Therefore, we can apply Fubini's theorem and (21–23), to compute

$$\int_{X_{\mathbb{R}^d}} (\partial_t \varphi + \nabla \varphi \cdot \nu + \varphi \ g) \ d\rho = \int_{\mathbb{R}^d} \int_0^1 \frac{d}{dt} \left(\varphi(t, X_x(t)) e^{\int_0^t g(s, X_x(s)) ds} \right) \ dt \ d\rho_0(x),$$

where $\rho = dt \otimes \rho_t$. Now notice that the above right-hand side vanishes since φ is compactly supported, concluding the proof.

The next proposition states that, under some regularity assumptions, every solution of (14) can be represented as in (21).

Proposition 2.4. Assume that $\rho_t : [0,1] \to \mathcal{M}^+(\mathbb{R}^d)$ is a narrowly continuous solution to the continuity equation $\partial_t \rho_t + \text{div} (v\rho_t) = g\rho_t$ in $(0,1) \times \mathbb{R}^d$ in the sense of (14), for some Borel maps $v : [0,1] \times \mathbb{R}^d \to \mathbb{R}^d$, $g : [0,1] \times \mathbb{R}^d \to \mathbb{R}$ satisfying (18, 20) and

$$\int_{0}^{1} \int_{\mathbb{R}^{d}} |v(t,x)| + |g(t,x)| \ d\rho_{t}(x) \ dt < \infty.$$
(24)

Then for ρ_0 -a.e. $x \in \mathbb{R}^d$ the ODE (19) admits a solution $X_x(t)$ for $t \in [0, 1]$, and

$$\rho_t = (X_{(\cdot)}(t))_{\#} \left(\rho_0 \ e^{\int_0^t g(s, X_{(\cdot)}(s)) ds} \right) \quad \text{for each} \quad t \in [0, 1],$$

where the push-forward is with respect to the space variable.

Proof. Define the map $t \mapsto \mu_t := (X_{(\cdot)}(t))_{\#}(\rho_0 \ e^{\int_0^t g(s, X_{(\cdot)}(s))ds})$. Proposition 2.3 implies that μ_t is a narrowly continuous solution to the continuity equation in $(0, 1) \times \mathbb{R}^d$. Moreover $\mu_0 = \rho_0$ by construction. It is immediate to check that $\mu_t - \rho_t$ and $\rho_t - \mu_t$ satisfy (121). As $\mu_t - \rho_t$ and $\rho_t - \mu_t$ both satisfy the continuity equation, we can apply (twice) the comparison principle in Proposition A.6 to deduce that $\mu_t = \rho_t$ for every $t \in [0, 1]$.

3. Functional analytic setting

In this section we discuss the functional analytic setting that is instrumental in proving the superposition principle in Theorem 1.1. Throughout the section, V will be the closure of a bounded domain of \mathbb{R}^d , with $d \in \mathbb{N}$, $d \ge 1$. We recall the notations $X_V :=$ $(0,1) \times V$ and $\mathcal{M}_V := \mathcal{M}(X_V) \times \mathcal{M}(X_V; \mathbb{R}^d) \times \mathcal{M}(X_V)$.

3.1. Curves in cones of measures

We start by introducing the set

$$\mathscr{C}_{V} := \left\{ h\delta_{\gamma} \in \mathcal{M}(V) : h \ge 0, \ \gamma \in V \right\}$$

$$(25)$$

and the space of narrowly continuous curves with values in \mathscr{C}_V , that is,

$$\mathscr{S}_{V} := \left\{ (t \mapsto \rho_{t}) \in C_{w}([0,1]; \mathcal{M}^{+}(V)) : \rho_{t} \in \mathscr{C}_{V} \text{ for all } t \in [0,1] \right\}.$$

$$(26)$$

Notice that if $t \mapsto \rho_t$ belongs to \mathscr{S}_V , then $\rho := dt \otimes \rho_t$ belongs to $\mathcal{M}(X_V)$. With a little abuse of notation, in what follows, we will denote by ρ both the curve $t \mapsto \rho_t$ and the measure $dt \otimes \rho_t$.

Remark 3.1. If $\rho \in \mathscr{S}_V$, then $\rho_t = h(t)\delta_{\gamma(t)}$ for $h: [0,1] \to [0,\infty)$ and $\gamma: [0,1] \to V$, where γ is uniquely determined in the set $\{h > 0\}$.

We endow the set \mathscr{C}_V with the flat distance on $\mathcal{M}(V)$, that is, for $\rho^i \in \mathscr{C}_V$ we set

$$d_{F}(\rho^{1},\rho^{2}) := \sup \left\{ \int_{V} \varphi \ d(\rho^{1} - \rho^{2}) : \varphi \in C(V), \|\varphi\|_{\infty} \le 1, \quad \operatorname{Lip}(\varphi, V) \le 1 \right\}.$$
(27)

We then define a distance over \mathscr{S}_V , by setting

$$\mathbf{d}(\rho^1, \rho^2) := \sup_{t \in [0, 1]} \mathbf{d}_F(\rho_t^1, \rho_t^2).$$
(28)

Remark 3.2. In [33, 36] the authors introduced the cone space over V given by $C_V := (V \times [0, \infty))/\sim$, where \sim is the equivalence relationship such that the pairs (γ_1, h_1) and (γ_2, h_2) are identified if and only if $\gamma_1 = \gamma_2$ and $h_1 = h_2$, or if $h_1 = h_2 = 0$. Notice that C_V is in one-to-one correspondence with \mathscr{C}_V . However in [33, 36] the cone space

is equipped with the cone distance

$$H^{2}(\rho^{1},\rho^{2}) := \begin{cases} h_{1} + h_{2} - 2\sqrt{h_{1}h_{2}}\cos\left(|\gamma_{1} - \gamma_{2}|\right) & \text{if } |\gamma_{1} - \gamma_{2}| \le \pi, \\ h_{1} + h_{2} + 2\sqrt{h_{1}h_{2}} & \text{otherwise,} \end{cases}$$

for all $\rho^1, \rho^2 \in \mathscr{C}_V$. By elementary calculations, and employing (29) below, it is possible to show that H^2 and d_F induce equivalent topologies on \mathscr{C}_V , for example, there exists a constant C > 0 such that

$$\frac{1}{C}H^{2}(\rho^{1},\rho^{2}) \leq d_{F}(\rho^{1},\rho^{2}) \leq C\sqrt{(h_{1}+h_{2})H^{2}(\rho^{1},\rho^{2})}.$$

The following characterization for d_F holds.

Lemma 3.3. For $\rho^1, \rho^2 \in \mathscr{C}_V$ we have

$$d_F(\rho^1, \rho^2) = \begin{cases} |h_1 - h_2| + \min(h_1, h_2)|\gamma_1 - \gamma_2| & \text{if } |\gamma_1 - \gamma_2| \le 2, \\ h_1 + h_2 & \text{otherwise.} \end{cases}$$
(29)

Proof. By definition it follows that

$$d_F(\rho^1,\rho^2) = \sup_{c_1,c_2 \in \mathbb{R}} \{h_1c_1 - h_2c_2 : |c_1|, |c_2| \le 1, |c_1 - c_2| \le |\gamma_1 - \gamma_2|\}$$

By symmetry we can assume $h_1 \ge h_2$. For all $c_1, c_2 \in \mathbb{R}$ such that $|c_1|, |c_2| \le 1$ and $|c_1 - c_2| \le |\gamma_1 - \gamma_2|$, we estimate

$$h_1c_1 - h_2c_2 \le |h_1c_1 - h_2c_1| + |h_2c_1 - h_2c_2| \le |h_1 - h_2| + \min(h_1, h_2)|\gamma_1 - \gamma_2|.$$

The thesis follows since the supremum is achieved by $(1, 1 - |\gamma_1 - \gamma_2|)$ if $|\gamma_1 - \gamma_2| \le 2$ and by (1, -1) otherwise.

We will now show that the metric space (\mathscr{G}_V, d) can be identified with $C([0, 1]; \mathscr{C}_V)$, where \mathscr{C}_V is equipped with d_F and $C([0, 1]; \mathscr{C}_V)$ inherits the relative topology as a subset of $C([0, 1]; \mathcal{M}_{\text{flat}}(V))$, $\mathcal{M}_{\text{flat}}(V)$ being the space $\mathcal{M}(V)$ equipped with the flat norm. In order to achieve that, we need a preliminary lemma.

Lemma 3.4. Let $\rho_t : [0,1] \to \mathscr{C}_V$. Then the following statements are equivalent:

- (i) ρ_t is narrowly continuous,
- (ii) $\rho_t = h(t)\delta_{\gamma(t)}$ with $h \in C[0,1]$ and $\gamma \in C(\{h > 0\}; \mathbb{R}^d)$.

Proof. Assume (i), so that the map $t \mapsto h(t)\varphi(\gamma(t))$ is continuous for each $\varphi \in C(V)$. By choosing $\varphi \equiv 1$ we conclude that h is continuous. If we pick $\varphi(x) := x_i$ coordinate function, for all i = 1, ..., d, we also infer continuity for $h\gamma$, so that γ is continuous in $\{h > 0\}$. Conversely, assume (ii). Let $\varphi \in C(V)$ and $\hat{t} \in [0, 1]$. If $h(\hat{t}) = 0$, we conclude continuity of $t \mapsto h(t)\varphi(\gamma(t))$ at \hat{t} by boundedness of φ and continuity of h, while if $h(\hat{t}) > 0$, we conclude by (ii).

Proposition 3.5. Assume that $\rho_t : [0,1] \to \mathscr{C}_V$. Then the following statements are equivalent:

- (i) ρ_t is narrowly continuous,
- (ii) ρ_t is continuous with respect to d_F .

In particular, we have that (\mathscr{G}_V, d) is a metric space that can be identified with $C([0, 1]; \mathscr{C}_V)$.

Proof. Assume (i), so that $h \in C[0, 1]$ and $\gamma \in C(\{h > 0\}; \mathbb{R}^d)$ by Lemma 3.4. Fix $t \in [0, 1]$ and $t_n \to t$. If h(t) = 0, by continuity of h and (29) we infer $d_F(\rho_{t_n}, \rho_t) = h(t_n) \to 0$. If instead h(t) > 0, by continuity of γ in t, it holds that $|\gamma(t_n) - \gamma(t)| \leq 2$ for n sufficiently large. By continuity of h we conclude (ii). Conversely, assume (ii). In order to show (i), we prove that $h \in C[0, 1]$ and $\gamma \in C(\{h > 0\}; \mathbb{R}^d)$ (Lemma 3.4). From (29) we have $|h(t_1) - h(t_2)| \leq d_F(\rho_{t_1}, \rho_{t_2})$ for all $t_1, t_2 \in [0, 1]$, so that h is continuous by (ii). Let us now fix $t \in \{h > 0\}$ and $t_n \to t$. Since h(t) > 0, it is immediate to check by contradiction that $|\gamma(t_n) - \gamma(t)| \leq 2$ eventually, and hence

$$|h(t_n) - h(t)| + \min(h(t_n), h(t))|\gamma(t_n) - \gamma(t)| = d_F(\rho_{t_n}, \rho_t),$$
(30)

for sufficiently large *n*. By continuity of *h*, (ii), and the assumption h(t) > 0, we conclude continuity for γ , and hence (i). The final part of the statement follows from the first part and from the definition of d.

For the space $(\mathcal{G}_V, \mathbf{d})$ the following holds.

Proposition 3.6. We have that (\mathscr{G}_V, d) is a complete separable metric space.

The above statement is somewhat classical. However, due to the lack of a reference, we provide a proof in Section A.4. We conclude this section with a useful lemma that provides sufficient conditions for continuity and measurability for scalar maps on (\mathscr{G}_V, d) .

Lemma 3.7. Let $\varphi : V \times [0, \infty) \to \mathbb{R}$ be such that $\varphi(x, 0) = 0$ for all $x \in V$. For $t \in [0, 1]$ define the map $\Psi_t : \mathscr{S}_V \to \mathbb{R}$ by $\Psi_t(\rho) := \varphi(\gamma(t), h(t))$, where $\rho_t = h(t)\delta_{\gamma(t)}$. If φ is measurable (resp. continuous), then Ψ_t is measurable (resp. continuous) with respect to d.

Proof. Notice that the condition $\varphi(x,0) = 0$ for all $x \in V$ implies that Ψ_t is well defined. Suppose first that φ is continuous and assume that $d(\rho^n, \rho) \to 0$ as $n \to \infty$. By (29) we have $|h_n(t) - h(t)| \leq d_F(\rho_t^n, \rho_t)$, so that $h_n(t) \to h(t)$. If h(t) = 0, then $\rho_t = 0$ and $\Psi_t(\rho) = 0$. By continuity of φ and compactness of V we infer that $\Psi_t(\rho^n) \to 0$. If h(t) > 0, the usual argument by contradiction implies that $|\gamma_n(t) - \gamma(t)| \leq 2$ for n sufficiently large. Thus by (29) and the convergences $\min(h_n(t), h(t)) \to h(t) > 0$ and $d_F(\rho_t^n, \rho_t) \to 0$, we have that $\gamma_n(t) \to \gamma(t)$. By continuity of φ we conclude $\Psi_t(\rho^n) \to \Psi_t(\rho)$. Suppose now that φ is measurable. Define the evaluation map $e_t : \mathscr{S}_V \to \mathscr{C}_V$ by $e_t(\rho) := \rho_t$ and the projection $\pi : \mathscr{C}_V \to V \times [0, \infty)$ by

$$\pi(h\delta_{\gamma}) := (\gamma, h) \ \chi_{\mathscr{C}_V \setminus \{(0,0)\}}(\gamma, h) + (p,0) \ \chi_{\{(0,0)\}}(\gamma, h),$$

where $p \in V$ is arbitrary but fixed. Notice that by construction e_t is continuous from $(\mathscr{S}_V, \mathbf{d})$ into $(\mathscr{C}_V, \mathbf{d}_F)$. Additionally the map $h\delta_{\gamma} \mapsto (\gamma, h)$ is continuous in $\mathscr{C}_V \setminus \{(0, 0)\}$ by repeating the above arguments. Hence π is measurable, being sum of measurable functions. Noting that $\Psi_t = \varphi \circ \pi \circ e_t$, we see that Ψ_t is measurable.

3.2. Properties of the Hellinger–Kantorovich energy over \mathscr{C}_V

In this section we investigate some properties of the coercive version of the Hellinger-Kantorovich energy at (17) when restricted to measures belonging to \mathscr{S}_V . To be more precise, we consider the functional $\mathscr{F} : \mathscr{G}_V \to [0, \infty]$ defined by

$$\mathscr{F}(\rho) := \inf\{J_{\alpha,\beta,\delta}(\rho,m,\mu) : (m,\mu) \in \mathcal{M}(X_V;\mathbb{R}^d) \times \mathcal{M}(X_V)\},\tag{31}$$

where $J_{\alpha,\beta,\delta}$ is defined at (17) and $\alpha,\beta,\delta > 0$. We start by introducing the subset of \mathscr{S}_V

$$\mathcal{H}_{V} := \left\{ \rho_{t} = h(t)\delta_{\gamma(t)} \in \mathscr{S}_{V} : h, \ \sqrt{h} \in \mathrm{AC}^{2}[0,1], \ \sqrt{h}\gamma \in \mathrm{AC}^{2}([0,1];\mathbb{R}^{d}) \right\}.$$
(32)

As already mentioned in the introduction, we denote by AC^2 the set of absolutely continuous functions with a.e. derivative in L^2 (see [12, Section 1.1] for a precise definition).

Lemma 3.8. Let
$$\rho_t = h(t)\delta_{\gamma(t)} \in \mathcal{H}_V$$
, $b \in C^1(V)$. Then $h(b \circ \gamma) \in AC^2[0, 1]$ with
 $(h(t)b(\gamma(t)))' = \dot{h}(t)b(\gamma(t)) + h(t)\nabla b(\gamma(t)) \cdot \dot{\gamma}(t)$ a.e. in (0,1). (33)

Proof. By definition of \mathcal{H}_V , it follows that $h\gamma \in \mathrm{AC}^2([0,1];\mathbb{R}^d)$. For every $0 \le t \le s \le 1$

$$\begin{split} h(t)b(\gamma(t)) - h(s)b(\gamma(s))| &\leq \operatorname{Lip}(b) \ h(t)|\gamma(s) - \gamma(t)| + ||b||_{\infty}|h(s) - h(t)| \\ &\leq \operatorname{Lip}(b) \ |h(s)\gamma(s) - h(t)\gamma(t)| + ||\gamma||_{\infty}\operatorname{Lip}(b) \ |h(s) - h(t)| + ||b||_{\infty}|h(s) - h(t)| \end{split}$$

Hence $h(b \circ \gamma) \in AC^2[0, 1]$. From the regularity assumed, we immediately infer the product rule at (33) for a.e. $t \in \{h > 0\}$. Moreover, using that $h(b \circ \gamma) \in AC^2([0, 1]; \mathbb{R}^d)$, we have $(h(t)b(\gamma(t)))' = 0$ almost everywhere in $\{h = 0\}$ ([43, Theorem 4.4]), so that (33) follows.

Proposition 3.9. Let $\rho_t = h(t)\delta_{\gamma(t)} \in \mathscr{G}_V$ and $(m,\mu) \in \mathcal{M}(X_V;\mathbb{R}^d) \times \mathcal{M}(X_V)$ be such that $J_{\alpha,\beta,\delta}(\rho,m,\mu) < \infty$. Then the following properties hold:

- (i) There exist $v: X_V \to \mathbb{R}^d$, $g: X_V \to \mathbb{R}$ measurable maps such that $m = v\rho$, $\mu = g\rho$,
- (ii) $\dot{\gamma}(t) = v(t, \gamma(t))$ for a.e. $t \in \{h > 0\}$ and $\dot{h}(t) = g(t, \gamma(t))h(t)$ for a.e. $t \in (0, 1)$,
- (iii) The curve $t \mapsto \rho_t$ belongs to \mathcal{H}_V .

Moreover the energy $J_{\alpha,\beta,\delta}$ can be computed by

$$J_{\alpha,\beta,\delta}(\rho,m,\mu) = \int_{\{h>0\}} \frac{\beta}{2} |\dot{\gamma}(t)|^2 h(t) + \frac{\beta\delta^2}{2} \frac{|\dot{h}(t)|^2}{h(t)} + \alpha h(t) dt.$$
(34)

Conversely, let $\rho_t = h(t)\delta_{\gamma(t)} \in \mathcal{H}_V$ and set $m := h(t)\dot{\gamma}(t)$ $dt \otimes \delta_{\gamma(t)}$, $\mu := \dot{h}(t)$ $dt \otimes \delta_{\gamma(t)}$. Then (ρ, m, μ) belongs to \mathcal{M}_V and solves the continuity equation (14) in X_V . Moreover $J_{\alpha,\beta,\delta}(\rho, m, \mu) < \infty$ and (34) holds.

Proof. Assume $\rho_t = h(t)\delta_{\gamma(t)} \in \mathscr{S}_V$, $(m,\mu) \in \mathcal{M}(X_V; \mathbb{R}^d) \times \mathcal{M}(X_V)$ and $J_{\alpha,\beta,\delta}(\rho,m,\mu) < \infty$. In particular, by definition of $J_{\alpha,\beta,\delta}$, we have that (ρ,m,μ) solves (14). By Lemma A.4 we deduce (*i*). We now show that the second ODE in (*ii*) holds. By Lemma 2.1 we

have $h \in BV(0,1)$, with distributional derivative given by $\pi_{\#}(g\rho)$, where $\pi: X_V \to (0,1)$ is the projection on the time coordinate. Thus, for all $\varphi \in C_c^{\infty}(0,1)$,

$$\pi_{\#}(g\rho)(\phi) = \int_{0}^{1} \int_{V} \phi(t)g(t,x) \ d\rho_{t}(x) \ dt = \int_{0}^{1} \phi(t)g(t,\gamma(t))h(t) \ dt.$$

Since $J_{\alpha,\beta,\delta}(\rho,m,\mu) < \infty$, by (119) and continuity of h, we conclude that $\dot{h}(t) = g(t,\gamma(t))h(t)$ almost everywhere and $h \in AC^2[0,1]$. We will now show that the first ODE in (*ii*) holds. By testing (14) against $\varphi(t,x) := a(t)b(x)$ with $a \in C_c^1(0,1), b \in C^1(V)$, we obtain

$$\frac{d}{dt}\int_{V}b(x) \ d\rho_t(x) = \int_{V} \left(\nabla b(x) \cdot v(t,x) + b(x)g(t,x)\right) \ d\rho_t(x), \quad \text{for a.e.} \quad t \in (0,1),$$

since the right-hand side belongs to $L^2(0,1)$, thanks to Jensen's inequality, (119) and the assumption $J_{\alpha,\beta,\delta}(\rho,m,\mu) < \infty$. In particular, choosing *b* as the coordinate functions, we deduce that $h\gamma \in AC^2([0,1]; \mathbb{R}^d)$ with $(h\gamma)'(t) = h(t)[\nu(t,\gamma(t)) + \gamma(t)g(t,\gamma(t))]$. In particular $\gamma \in AC^2(\{h \ge c\}; \mathbb{R}^d)$ for every c > 0, given that *V* is bounded. Consider now the test function $\varphi \in C_c^1((\{h > 0\} \cap (0,1)) \times V)$. Using that $\dot{h}(t) = g(t,\gamma(t))h(t)$ almost everywhere, it is easy to check that the equation $\partial_t \rho_t + \text{div} (\nu \rho_t) = g\rho_t$ can be rewritten as

$$\int_0^1 \frac{d}{dt} (h(t)\varphi(t,\gamma(t))) \ dt + \int_0^1 \ \nabla\varphi(t,\gamma(t)) \cdot (\nu(t,\gamma(t)) - \dot{\gamma}(t)) \ h(t) \ dt = 0.$$
(35)

Notice that the first integral in (35) vanishes, as φ is compactly supported. Set $\varphi(t, x) := a(t)x_i$ with $a \in C_c^1(\{h > 0\} \cap (0, 1))$ and x_i coordinate function. Testing (35) against φ yields (*ii*). By (*ii*), Lemma A.4, and the energy bound, we also see that (34) holds. We are left to show (*iii*). First we claim that $\sqrt{h}\gamma \in AC^2([0, 1]; \mathbb{R}^d)$. Indeed, for $\varepsilon > 0$ and $\varphi \in C_c^{\infty}(0, 1)$, an integration by parts yields

$$\int_{0}^{1} h(t)\gamma(t) \frac{1}{\sqrt{h(t)+\varepsilon}} \dot{\varphi}(t) \ dt = -\int_{\{h>0\}} \left[\frac{(h(t)\gamma(t))'}{\sqrt{h(t)+\varepsilon}} - \frac{h(t)\gamma(t)\dot{h}(t)}{2(h(t)+\varepsilon)^{3/2}} \right] \ \varphi(t) \ dt, \tag{36}$$

where we used that $h\gamma \in AC^2([0,1]; \mathbb{R}^d)$, $(h\gamma)' = 0$ a.e. in $\{h = 0\}$ (see, e.g., [43, Theorem 4.4]) and $(h\gamma)' = \dot{h}\gamma + h\dot{\gamma}$ a.e. in $\{h > 0\}$. By (34), continuity of h, boundedness of V, we can invoke dominated convergence and pass to the limit as $\varepsilon \to 0$ in (36), thus concluding that $\sqrt{h}\gamma \in AC^2([0,1]; \mathbb{R}^d)$ with derivative given by $2^{-1}\chi_{\{h>0\}}\dot{h}\gamma/\sqrt{h} + \sqrt{h}\dot{\gamma}$. A similar argument shows that $\sqrt{h} \in AC^2[0,1]$ with derivative given by $2^{-1}\chi_{\{h>0\}}\dot{h}/\sqrt{h}$, concluding the proof of (*iii*) and of the direct implication.

Conversely, assume that $\rho_t = h(t)\delta_{\gamma(t)} \in \mathcal{H}_V$ and set $m := h(t)\dot{\gamma}(t) dt \otimes \delta_{\gamma(t)}, \mu := \dot{h}(t) dt \otimes \delta_{\gamma(t)}$. It is clear that $(\rho, m, \mu) \in \mathcal{M}_V$, as a consequence of the regularity on h and γ . We claim that (ρ, m, μ) solves (14) in X_V . Fix $b \in C^1(V)$. By Lemma 3.8 we have that $h(b \circ \gamma) \in AC^2[0, 1]$ and (33) holds. Thus, for all $a \in C_c^1(0, 1)$,

$$\int_{X_V} a'(t)b(x) \ d\rho_t(x) \ dt = -\int_{X_V} a(t)b(x) \ d\mu - \int_{X_V} a(t)\nabla b(x) \cdot dm.$$

Therefore, by employing a standard density argument, (ρ, m, μ) solves (14) in $X_{V.}$ Finally, by the regularity of h, γ and (119), we conclude that $J_{\alpha,\beta,\delta}(\rho, m, \mu) < \infty$ and (34) holds.

Proposition 3.10. Let $\mathscr{F} : (\mathscr{G}_V, \mathbf{d}) \to [0, \infty]$ be the functional defined at (31). The domain of \mathscr{F} is given by \mathcal{H}_V , where we have

$$\mathscr{F}(\rho) = \mathscr{F}(\gamma, h) = \int_{\{h>0\}} \frac{\beta}{2} |\dot{\gamma}(t)|^2 h(t) + \frac{\beta \delta^2}{2} \frac{|h(t)|^2}{h(t)} + \alpha h(t) dt.$$
(37)

Moreover \mathcal{F} is lower semi-continuous and its sublevel sets are compact.

Proof. We start by showing that the domain of \mathscr{F} is given by \mathcal{H}_V and that (37) holds. Assume first that $\rho^* \in \mathscr{S}_V$ and $\mathscr{F}(\rho^*) < \infty$. We claim that exists a pair $(m^*, \mu^*) \in \mathcal{M}(X_V; \mathbb{R}^d) \times \mathcal{M}(X_V)$ such that

$$\mathscr{F}(\rho^*) = J_{\alpha,\beta,\delta}(\rho^*, m^*, \mu^*).$$
(38)

Indeed the functional $(m,\mu) \mapsto J_{\alpha,\beta,\delta}(\rho,m,\mu)$ is weak^{*} lower semi-continuous by Lemma A.5. Invoking (120) and the direct method, we conclude that the infimum at (31) is achieved, showing (38). Hence we can apply the direct implication of Proposition 3.9 to (ρ^*, m^*, μ^*) to obtain that $\rho^* \in \mathcal{H}_V$ and that (37) holds. Conversely, assume that $\rho^*_t = h(t)\delta_{\gamma(t)} \in \mathcal{H}_V$ and set $m := \dot{\gamma}\rho^*$, $\mu := (\dot{h}/h)\rho^*$. By the converse implication of Proposition 3.9 we know that $(\rho^*, m, \mu) \in \mathcal{M}_V$ and $J_{\alpha,\beta,\delta}(\rho^*, m, \mu) < \infty$, from which we infer $\mathscr{F}(\rho^*) < \infty$. Thus there exists a pair $(m^*, \mu^*) \in \mathcal{M}(X_V; \mathbb{R}^d) \times \mathcal{M}(X_V)$ such that (38) holds. An application of the direct implication of Proposition 3.9 to (ρ^*, m^*, μ^*) yields (37).

We now prove that \mathscr{F} is lower semi-continuous with respect to d. To this end, assume that $d(\rho^n, \rho) \to 0$ as $n \to \infty$. We claim that $dt \otimes \rho_t^n \rightharpoonup^* dt \otimes \rho_t$ weakly^{*} in $\mathcal{M}(X_V)$. By density, it is sufficient to prove convergence for test functions $\varphi(t, x) = a(t)b(x)$ with $a \in C_c(0, 1), b \in C(V)$. Moreover, it is not restrictive to assume that $\|b\|_{\infty} \leq 1$. For a fixed $\varepsilon > 0$ there exists $c \in C^1(V)$ such that $\|c\|_{\infty} \leq 1$ and $\|b - c\|_{\infty} < \varepsilon$. For $t \in [0, 1]$ we have

$$\begin{aligned} \left| \int_{V} b(x) \ d(\rho_{t}^{n} - \rho_{t}) \right| &\leq \|b - c\|_{\infty}(\|\rho_{t}^{n}\|_{\mathcal{M}(V)} + \|\rho_{t}\|_{\mathcal{M}(V)}) + \left| \int_{V} c(x) \ d(\rho_{t}^{n} - \rho_{t}) \right| \\ &\leq \varepsilon(d(\rho^{n}, 0) + d(\rho, 0)) + \operatorname{Lip}(c) \ d(\rho^{n}, \rho) \leq \varepsilon C + \operatorname{Lip}(c) \ d(\rho^{n}, \rho) \end{aligned}$$

where the first term in the first line was estimated by (29), and the second one by (27). Since the estimate does not depend on t, and ε is arbitrary, we conclude that $dt \otimes \rho_t^n \longrightarrow^* dt \otimes \rho_t$. We now claim that \mathscr{F} is weak^{*} lower semi-continuous in \mathscr{S}_V considered as a subset of $\mathcal{M}(X_V)$: Indeed assume that $\rho^n \longrightarrow^* \rho$ in $\mathcal{M}(X_V)$. Without loss of generality we can assume that $\sup_n \mathscr{F}(\rho_n) < \infty$ along a subsequence, so that there exist $(m^n, \mu^n) \in \mathcal{M}(X_V; \mathbb{R}^d) \times \mathcal{M}(X_V)$ such that, up to subsequences, $\mathscr{F}(\rho^n) = J_{\alpha,\beta,\delta}(\rho^n, m^n, \mu^n)$. By (120) we infer the existence of a pair (m, μ) such that, up to subsequences, $m^n \longrightarrow^* m$, $\mu^n \longrightarrow^* \mu$. We can now invoke weak^{*} lower semi-continuity of $J_{\alpha,\beta,\delta}$ (Lemma A.5) to conclude weak^{*} lower semi-continuity of \mathscr{F} . Since $dt \otimes \rho_t^n \longrightarrow^* dt \otimes \rho_t$ in $\mathcal{M}(X_V)$ whenever $d(\rho^n, \rho) \to 0$, we infer lower semi-continuity of \mathscr{F} with respect to d.

Finally, we show that the sublevel sets of \mathscr{F} are compact with respect to d. As $\mathscr{F} \geq 0$ and is positively one-homogeneous, it is enough to show that $S_{\mathscr{F}} := \{\rho \in \mathscr{S}_V : \mathscr{F}(\rho) \leq 1\}$ is compact. Let $\rho_t = h(t)\delta_{\gamma(t)} \in S_{\mathscr{F}}$, so that, in particular, $\rho \in \mathcal{H}_V$. In order to show compactness of $S_{\mathscr{F}}$ we first provide some preliminary estimates for the maps h and $h\gamma$. By (37) we immediately infer that $||h||_1 \leq 1/\alpha$. Let $0 \leq t_1 \leq t_2 \leq 1$. There holds

$$h(t_{2}) - h(t_{1}) \leq \int_{t_{1}}^{t_{2}} |\dot{h}(s)| ds = \int_{(t_{1}, t_{2}) \cap \{h > 0\}} |\dot{h}(s)| ds = \int_{(t_{1}, t_{2}) \cap \{h > 0\}} \frac{|\dot{h}(s)|}{\sqrt{h(s)}} \sqrt{h(s)} ds$$

$$\leq \left(\int_{\{h > 0\}} \frac{|\dot{h}(s)|^{2}}{h(s)} ds \right)^{1/2} \left(\int_{t_{1}}^{t_{2}} h(s) ds \right)^{1/2} \leq \frac{2}{\beta \delta^{2}} \left(\int_{t_{1}}^{t_{2}} h(s) ds \right)^{1/2},$$
(39)

where we used that $\dot{h} = 0$ almost everywhere in $\{h = 0\}$ ([43, Theorem 4.4]), Hölder's inequality, and the fact that $\mathscr{F}(\rho) \leq 1$ in conjunction with (37). Since $h \geq 0$, choosing $t_1 \in \arg \min h$ in the above estimate yields

$$\|h\|_{\infty} \leq \frac{2}{\beta\delta^2} \|h\|_1^{1/2} + \|h\|_1 \leq C, \qquad \|h\gamma\|_{\infty} \leq CR,$$
 (40)

where $R := \max\{|p| : p \in V\}, C := 2/(\beta \delta^2 \sqrt{\alpha}) + 1/\alpha$. Recall that $R < \infty$ as V is bounded. Thus, by (39) and (40),

$$|h(t_2) - h(t_1)| \le \left(\int_{t_1}^{t_2} h(s)ds\right)^{1/2} \le C |t_1 - t_2|^{1/2} \quad \text{for all} \quad 0 \le t_1 \le t_2 \le 1.$$
(41)

Moreover, by (39)-(40) we can estimate

$$\int_{t_1}^{t_2} |\dot{h}(s)\gamma(s)| ds \le R \int_{t_1}^{t_2} |\dot{h}(s)| ds \le R \left(\int_{t_1}^{t_2} h(s) ds \right)^{1/2} \le CR |t_1 - t_2|^{1/2}.$$

and also

$$\int_{t_1}^{t_2} |h(s)\dot{\gamma}(s)| ds \leq \left(\int_{t_1}^{t_2} h(s) ds\right)^{1/2} \left(\int_0^1 |\dot{\gamma}(s)|^2 h(s) ds\right)^{1/2} \leq \frac{2C}{\beta} |t_1 - t_2|^{1/2}$$

where we used Hölder's inequality, (37, 41), and $\mathscr{F}(\rho) \leq 1$. By Lemma 3.8 and the above estimates we thus infer

$$|h(t_1)\gamma(t_1) - h(t_2)\gamma(t_2)| \le \int_{t_1}^{t_2} |\dot{h}(s)\gamma(s)| ds + \int_{t_1}^{t_2} |h(s)\dot{\gamma}(s)| ds \le C(R + 2\beta^{-1})|t_1 - t_2|^{1/2}$$
(42)

for every $0 \le t_1 \le t_2 \le 1$. Hence, considering a sequence $\{\rho^n\}_n$ in $S_{\mathscr{F}}$ with $\rho_t^n = h_n(t)\delta_{\gamma_n(t)}$, by (40)–(42) we have that h_n and $h_n\gamma_n$ are equibounded and equicontinuous. Therefore Ascoli–Arzelà's theorem implies that, up to subsequences, $h_n \to h$ and $\gamma_n h_n \to f$ uniformly, where $h \in C[0, 1], h \ge 0$ and $f \in C([0, 1]; \mathbb{R}^d)$. Define $\gamma(t) := f(t)/h(t)$ if h(t) > 0. By the uniform convergence $h_n \to h$ we have that $\gamma(t) \in V$ for $t \in C[0, 1]$.

 $\{h > 0\}$. Therefore, by setting $\rho_t := h(t)\delta_{\gamma(t)}$, Lemma 3.4 implies that $\rho \in \mathscr{S}_V$. Since $h_n \to h$ pointwise and $\gamma_n \to \gamma$ pointwise in $\{h > 0\}$, and since $\|h_n\|_{\infty} \leq C$, by dominated convergence one immediately concludes that $dt \otimes \rho_t^n \rightharpoonup^* dt \otimes \rho_t$ in $\mathcal{M}(X_V)$. We can then invoke the weak* lower semi-continuity of \mathscr{F} to conclude that $\rho \in S_{\mathscr{F}}$. We are left to prove that $\rho^n \to \rho$ with respect to d. Fix $\varepsilon > 0$. By the uniform convergences $h_n \to h$ and $h_n \gamma_n \to h \gamma$, there exists $N(\varepsilon) \in \mathbb{N}$ such that

$$|h_n(t) - h(t)| < \frac{\varepsilon}{R}, \quad |h_n(t)\gamma_n(t) - h(t)\gamma(t)| < \varepsilon, \quad \text{for all} \quad n \ge N(\varepsilon), \quad t \in [0, 1].$$
(43)

Let $t \in \{h \ge \varepsilon\}$ and $n \ge N(\varepsilon)$. Using the above condition we infer

$$|\gamma_n(t) - \gamma(t)| \le \left|\frac{h_n(t)}{h(t)}\gamma_n(t) - \gamma(t)\right| + |\gamma_n(t)| \left|\frac{h_n(t)}{h(t)} - 1\right| < \frac{\varepsilon}{h(t)} + R \quad \frac{\varepsilon}{Rh(t)} \le 2.$$

Set $m_n(t) := \min(h_n(t), h(t))$. Then, by (29),

$$\begin{aligned} \mathsf{d}_F(\rho_t^n,\rho_t) &< \frac{\varepsilon}{R} + m_n(t) \ |\gamma_n(t) - \gamma(t)| \\ &\leq \frac{\varepsilon}{R} + |\gamma_n(t)| \ |m_n(t) - h_n(t)| + |h_n(t)\gamma_n(t) - h(t)\gamma(t)| + |\gamma(t)| \ |m_n(t) - h(t)| \\ &\leq \frac{\varepsilon}{R} + 2R|h_n(t) - h(t)| + |h_n(t)\gamma_n(t) - h(t)\gamma(t)| < (R^{-1} + 3)\varepsilon. \end{aligned}$$

Let now $t \in \{h \le \varepsilon\}$. By triangle inequality and (29, 43)

$$d_F(\rho_t^n, \rho_t) \le h_n(t) + h(t) \le |h_n(t) - h(t)| + 2h(t) \le \varepsilon(R^{-1} + 2)$$

In total we infer $d(\rho^n, \rho) < C\varepsilon$ for $n \ge N(\varepsilon)$, concluding the proof.

4. The main decomposition theorem

In this section we will prove the decomposition result in Theorem 1.1 anticipated in the introduction. Specifically, the proof is presented in Sections 4.1, 4.3, while Section 4.2 contains auxiliary results which are instrumental to the proof.

For reader's convenience we will recall a few notations and the statement of Theorem 1.1. Let $d \in \mathbb{N}$, $d \ge 1$ and $V \subset \mathbb{R}^d$ be the closure of a bounded domain of \mathbb{R}^d . We denote the time-space cylinder by $X_V := (0,1) \times V$. We also recall the definitions of \mathscr{C}_V and \mathscr{S}_V at (25)–(26). The set \mathscr{C}_V is equipped with the flat metric d_F defined at (27), while \mathscr{S}_V is equipped with the supremum distance d defined at (28). We remind the reader that (\mathscr{S}_V, d) is a complete metric space (Proposition 3.6). Moreover we will also consider the set \mathcal{H}_V introduced at (32). Let $v : X_V \to \mathbb{R}^d$, $g : X_V \to \mathbb{R}$ be given measurable maps and consider the system of ODEs

$$\dot{\gamma}(t) = v(t, \gamma(t))$$
 a.e. in $\{h > 0\},$ (O1)

$$\dot{h}(t) = g(t, \gamma(t))h(t)$$
 a.e. in (0,1). (O2)

For *v* and *g* as above, we define the following subset of H_V :

$$\mathcal{H}_{V}^{\nu,g} := \left\{ \rho_{t} = h(t)\delta_{\gamma(t)} \in \mathcal{H}_{V} : (h,\gamma) \text{ satisfy } (O1) - (O2) \right\}$$

Also define the subset of \mathscr{S}_V

$$\mathcal{H}_{V}^{1} := \{(h, \gamma) \in \mathscr{S}_{V} : ||h||_{1} = 1\}.$$
(45)

Finally, define the subset of $\mathcal{M}^+(\mathscr{G}_V)$:

$$\mathcal{M}_1^+(\mathscr{S}_V) := \bigg\{ \sigma \in \mathcal{M}^+(\mathscr{S}_V) \; : \; \int_{\mathscr{S}_V} ||h||_\infty \; d\sigma(\gamma,h) < \infty \bigg\},$$

where the notation $d\sigma(\gamma, h)$ is a shorthand for expressing that the integral is computed on all curves $\rho_t = h(t)\delta_{\gamma(t)} \in \mathscr{S}_V$.

Definition 4.1. For a measure $\sigma \in \mathcal{M}^+(\mathscr{S}_V)$ we define the set function ρ_t^{σ} as

$$\rho_t^{\sigma}(E) := \int_{\mathscr{S}_V} h(t) \ \chi_E(\gamma(t)) \ d\sigma(\gamma, h)$$
(46)

for all Borel sets $E \subset V$ and $t \in [0, 1]$.

Remark 4.2. The map $(\gamma, h) \mapsto h(t)\chi_E(\gamma(t))$ at (46) is measurable in (\mathscr{S}_V, d) by Lemma 3.7; therefore the integral is well defined, possibly unbounded. Assume in addition that $\sigma \in \mathcal{M}_1^+(\mathscr{S}_V)$. It is easy to check that ρ_t^{σ} at (46) belongs to $\mathcal{M}^+(V)$ for all $t \in [0, 1]$. Moreover, if $\varphi \in L^1_{\rho_t^{\sigma}}(V)$ for some fixed $t \in [0, 1]$, then the map $(\gamma, h) \mapsto h(t)\varphi(\gamma(t))$ belongs to $L^1_{\sigma}(\mathscr{S}_V)$ and

$$\int_{V} \varphi(x) \ d\rho_{t}^{\sigma}(x) = \int_{\mathscr{S}_{V}} h(t) \ \varphi(\gamma(t)) \ d\sigma(\gamma, h).$$
(47)

This fact can be shown by mimicking the proof of [44, Theorem 3.6.1], in conjunction with Lemma 3.7. Similarly, if $\varphi: V \to \mathbb{R} \cup \{\pm \infty\}$ is measurable and the map $(\gamma, h) \mapsto h(t)\varphi(\gamma(t))$ belongs to $L^1_{\sigma}(\mathscr{S}_V)$, then $\varphi \in L^{1}_{\rho_t^{\sigma}}(V)$ and (47) holds.

We are now ready to state the main decomposition result of the article.

Theorem 4.3. Assume that $\Omega \subset \mathbb{R}^d$ is the closure of a bounded domain, with $d \in \mathbb{N}$, $d \ge 1$. Let $\rho_t \in C_w([0,1]; \mathcal{M}^+(\Omega))$ be a measure solution of the continuity equation $\partial_t \rho_t + \operatorname{div}(v\rho_t) = g\rho_t$ in X_Ω in the sense of (14), for some measurable maps $v : X_\Omega \to \mathbb{R}^d$, $g : X_\Omega \to \mathbb{R}$ satisfying

$$\int_{0}^{1} \int_{\Omega} |v(t,x)|^{2} + |g(t,x)|^{2} \ d\rho_{t}(x) \ dt < \infty.$$
(48)

Then there exists a measure $\sigma \in \mathcal{M}_1^+(\mathscr{S}_\Omega)$ concentrated on $\mathcal{H}_\Omega^{\nu,g} \cap \mathcal{H}_\Omega^1$ and such that $\rho_t = \rho_t^{\sigma}$ for all $t \in [0, 1]$, where ρ_t^{σ} is defined at (46), that is,

$$\int_{\Omega} \varphi(x) \ d\rho_t(x) = \int_{\mathscr{S}_{\Omega}} h(t)\varphi(\gamma(t)) \ d\sigma(\gamma,h) \quad \text{for all} \quad \varphi \in C(\Omega).$$
(49)

Conversely, assume that $\sigma \in \mathcal{M}^+(\mathscr{S}_{\Omega})$ is concentrated on $\mathcal{H}^{^{v,g}}_{\Omega}$ and satisfies the bound

$$\int_{0}^{1} \int_{\mathscr{S}_{\Omega}} h(t) \left(1 + |v(t,\gamma(t))| + |g(t,\gamma(t))| \right) \, d\sigma(\gamma,h) \, dt < \infty.$$
(50)

Then σ belongs to $\mathcal{M}_1^+(\mathscr{S}_{\Omega})$ and ρ_t^{σ} defined by (46) belongs to $C_w([0,1]; \mathcal{M}^+(\Omega))$ and satisfies $\partial_t \rho_t^{\sigma} + \operatorname{div}(v \rho_t^{\sigma}) = g \rho_t^{\sigma}$ in X_{Ω} .

Remark 4.4. Condition (50) is natural in the following sense. If ρ_t satisfies the assumptions of Theorem 4.3, then in particular the map $\varphi(t,x) := 1 + |\nu(t,x)| + |g(t,x)|$ belongs to $L^1_{\rho^{\sigma}_t}(\Omega)$ for a.e. $t \in (0, 1)$, thanks to (48, 49) and narrow continuity of ρ_t . Therefore, by applying Remark 4.2, we see that the measure σ representing ρ_t satisfies (50).

4.1. Proof of the converse implication of Theorem 4.3

We now prove the converse statement in Theorem 4.3. To this end, assume that $\sigma \in \mathcal{M}^+(\mathscr{S}_{\Omega})$ is concentrated on $\mathcal{H}^{\nu,g}_{\Omega}$ and (50) holds. Let us first show that $\sigma \in \mathcal{M}^+_1(\mathscr{S}_{\Omega})$. Let $\rho_t = h(t)\delta_{\gamma(t)} \in \mathscr{S}_{\Omega}$ and $t^* \in \arg \min h$, which exists by continuity of h (see Lemma 3.4). Using the definition of $\mathcal{H}^{\nu,g}_{\Omega}$ we can estimate

$$h(t) = h(t^*) + \int_{t^*}^t \dot{h}(\tau) \ d\tau \leq \int_0^1 h(\tau) \ d\tau + \int_{t^*}^t g(\tau, \gamma(\tau))h(\tau) \ d\tau \qquad \sigma - \text{a.e. in } \mathscr{S}_{\Omega},$$

for all $t \in [0, 1]$. In particular,

$$\|h\|_{\infty} \leq \int_0^1 h(t)(1+|g(t,\gamma(t))|) dt \qquad \sigma-\text{a.e. in } \mathscr{S}_{\Omega}, \tag{51}$$

concluding that $\sigma \in \mathcal{M}_1^+(\mathscr{S}_\Omega)$, thanks to (50). We now show that the curve $t \mapsto \rho_t^\sigma$ defined by (46) belongs to $C_w([0,1]; \mathcal{M}^+(\Omega))$. First, Remark 4.2 implies that $\rho_t^\sigma \in \mathcal{M}^+(\Omega)$ for all $t \in [0,1]$. For the narrow continuity, fix $\varphi \in C(\Omega)$ and notice that by definition the map $t \mapsto h(t)\varphi(\gamma(t))$ is continuous for all $\rho_t = h(t)\delta_{\gamma(t)} \in \mathscr{S}_\Omega$. Since $\sigma \in \mathcal{M}_1^+(\mathscr{S}_\Omega)$ we can apply dominated convergence and conclude that also $t \mapsto \int_\Omega \varphi(x) \ d\rho_t^\sigma(x)$ is continuous. We are left to show that ρ^σ solves the continuity equation $\partial_t \rho_t^\sigma + \operatorname{div}(v\rho_t^\sigma) = g\rho_t^\sigma$ in X_Ω . To this end, fix $b \in C^1(\Omega)$. By Lemma 3.8 the map $t \mapsto h(t)b(\gamma(t))$ is differentiable almost everywhere and (33) holds. Therefore, for all $0 \leq s \leq t \leq 1$ the following holds

$$\begin{split} \int_{\Omega} b \ d\rho_t^{\sigma} &- \int_{\Omega} b \ d\rho_s^{\sigma} = \int_{\mathscr{S}_{\Omega}} \int_s^t \frac{d}{d\tau} [h(\tau) b(\gamma(\tau))] \ d\tau \ d\sigma(\gamma, h) \\ &= \int_{\mathscr{S}_{\Omega}} \int_s^t \dot{h}(\tau) b(\gamma(\tau)) + h(\tau) \nabla b(\gamma(\tau)) \cdot \dot{\gamma}(\tau) \ d\tau \ d\sigma(\gamma, h) \\ &= \int_s^t \int_{\mathscr{S}_{\Omega}} h(\tau) [b(\gamma(\tau)) g(\tau, \gamma(\tau)) + \nabla b(\gamma(\tau)) \cdot v(\tau, \gamma(\tau))] \ d\sigma(\gamma, h) \ d\tau, \end{split}$$

where in the last equality we used that σ is concentrated on $\mathcal{H}_{\Omega}^{v,g}$ and applied Fubini's Theorem, which we are allowed to do as the integrand is absolutely integrable by (50), triangle inequality, and the fact that $b \in C^1(\Omega)$. In particular, the map $t \mapsto \int_{\Omega} b(x) \ d\rho_t^{\sigma}(x)$ is absolutely continuous with almost everywhere derivative given by

$$\frac{d}{dt} \int_{\Omega} b(x) \ d\rho_t^{\sigma}(x) = \int_{\mathscr{S}_{\Omega}} h(t) \left[b(\gamma(t))g(t,\gamma(t)) + \nabla b(\gamma(t)) \cdot v(t,\gamma(t)) \right] \ d\sigma(\gamma,h)
= \int_{\Omega} b(x)g(t,x) + \nabla b(x) \cdot v(t,x) \ d\rho_t^{\sigma}(x).$$
(52)

The second equality in (52) follows because v and g are measurable and hence $\Psi(t, x) := b(x)g(t, x) + \nabla b(x) \cdot v(t, x)$ is measurable in Ω for a.e. t fixed. From (50) we have that $(\gamma, h) \mapsto h(t)\Psi(t, \gamma(t))$ belongs to $L^1_{\sigma}(\mathscr{S}_{\Omega})$ for a.e. t, and hence by Remark 4.2 we can apply (47) to $\Psi(t, \cdot)$ and obtain the second equality in (52). Identity (52) implies that ρ^{σ}_t solves the continuity equation in X_{Ω} in the sense of (14), for all $\varphi \in C^1_c(X_{\Omega})$ of the form $\varphi(t, x) = a(t)b(x)$ for $a \in C^1_c(0, 1), b \in C^1(\Omega)$, and hence, by density, for all the elements of $C^1_c(X_{\Omega})$.

4.2. Regularized solutions of the continuity equation

Before starting the proof of the direct statement in Theorem 4.3, we provide some smoothing arguments which will be employed to construct the measure σ . To this end, let $\Omega \subset \mathbb{R}^d$, $d \in \mathbb{N}$, $d \ge 1$ be the closure of a bounded domain. Let $v : X_\Omega \to \mathbb{R}^d$, $g : X_\Omega \to \mathbb{R}$ be given measurable maps, and $\rho_t \in C_w([0,1]; \mathcal{M}^+(\Omega))$ be such that $\partial_t \rho_t + \operatorname{div}(v\rho_t) = g\rho_t$ in X_Ω in the sense of (14). We extend v, g to zero to the space $(0,1) \times \mathbb{R}^d$. Similarly extend ρ_t to zero so that $\rho_t \in \mathcal{M}^+(\mathbb{R}^d)$. Notice that the extensions (ρ, v, g) satisfy the continuity equation in $(0, 1) \times \mathbb{R}^d$, due to the no-flux boundary conditions. For $x \in \mathbb{R}^d$, r > 0 let $B_r(x) := \{x \in \mathbb{R}^d : |x| < r\}$ and let $\xi \in C^\infty(\mathbb{R}^d)$ be such that $\xi \ge 0$, supp $\xi \subset B_1(0)$ and $\int_{\mathbb{R}^d} \xi dx = 1$. For every $0 < \varepsilon < 1$ and $x \in \mathbb{R}^d$ set $\xi_\varepsilon(x) :$ $= \varepsilon^{-d}\xi(x\varepsilon^{-1})$. Note that supp $\xi_\varepsilon \subset B_\varepsilon(0)$. Let R > 0 be such that

$$\{x \in \mathbb{R}^d : \operatorname{dist}(x, \Omega) \le 2\} \subset V$$
, $V := \overline{B_R(0)}$, (53)

and define

$$\rho_t^{\varepsilon} := (\rho_t * \xi_{\varepsilon}) + \eta_{\varepsilon}, \quad \eta_{\varepsilon} := \varepsilon \, \chi_V, \quad v_t^{\varepsilon} := \frac{(v_t \rho_t) * \xi_{\varepsilon}}{\rho_t^{\varepsilon}}, \quad g_t^{\varepsilon} := \frac{(g_t \rho_t) * \xi_{\varepsilon}}{\rho_t^{\varepsilon}}, \tag{54}$$

where v_t^{ε} and g_t^{ε} are set to be zero in the region where $\rho_t^{\varepsilon}(x) = 0$, i.e., in $(0,1) \times (\mathbb{R}^d \setminus V)$. Here, with a little abuse of notation, we denote $v_t = v(t, \cdot), v_t^{\varepsilon} = v^{\varepsilon}(t, \cdot), g_t = g(t, \cdot), g_t^{\varepsilon} = g^{\varepsilon}(t, \cdot)$.

Lemma 4.5. Let $\rho_t \in C_w([0,1]; \mathcal{M}^+(\Omega))$ and $v : X_\Omega \to \mathbb{R}^d$, $g : X_\Omega \to \mathbb{R}^d$ be measurable. Suppose that $\partial_t \rho_t + \operatorname{div}(v\rho_t) = g\rho_t$ in X_Ω in the sense of (14) and that (48) holds. Let $(\rho_t^\varepsilon, v_t^\varepsilon, g_t^\varepsilon)$ be defined as in (54). Then $(\rho_t^\varepsilon dx, v_t^\varepsilon, g_t^\varepsilon)$ is a solution to $\partial_t \rho_t^\varepsilon dx + \operatorname{div}(v_t^\varepsilon \rho_t^\varepsilon dx) = g_t^\varepsilon \rho_t^\varepsilon dx$ in $(0,1) \times \mathbb{R}^d$ and $\rho_t^\varepsilon dx \to \rho_t$ narrowly in $\mathcal{M}(V)$ as $\varepsilon \to 0$, for all $t \in [0,1]$. Moreover v^ε and g^ε satisfy (18) and (20), respectively. Finally, for every $t \in [0,1]$ there holds

$$\int_{\mathbb{R}^d} |v^{\varepsilon}(t,x)|^2 \rho_t^{\varepsilon}(x) dx \leq \int_{\Omega} |v(t,x)|^2 d\rho_t(x),$$

$$\int_{\mathbb{R}^d} |g^{\varepsilon}(t,x)|^2 \rho_t^{\varepsilon}(x) dx \leq \int_{\Omega} |g(t,x)|^2 d\rho_t(x).$$
(55)

Proof. By the interplay between weak differentiation and mollification, it is immediate to check that $(\rho_t^{\varepsilon} dx, v_t^{\varepsilon}, g_t^{\varepsilon})$ solves the continuity equation in $(0, 1) \times \mathbb{R}^d$ for all $0 < \varepsilon < 1$. The fact that $\rho_t^{\varepsilon} dx \to \rho_t$ narrowly is an immediate consequence of the properties of convolutions and of the convergence $\eta_{\varepsilon} \to 0$ as $\varepsilon \to 0$. We now prove that v^{ε} satisfies (18). Notice that by definition $v_t^{\varepsilon}(x) = 0$ in $\mathbb{R}^d \setminus (\overline{\Omega + B_1(0)})$ for every $t \in [0, 1]$. Moreover $\rho_t^{\varepsilon} \ge \varepsilon$ in V for all t. Therefore

$$\int_0^1 \sup_{x \in \mathbb{R}^d} |v^{\varepsilon}(t,x)| \ dt \leq \frac{1}{\varepsilon} \int_0^1 \sup_{x \in \overline{\Omega + B_1(0)}} |(v_t \rho_t) * \xi_{\varepsilon}(x)| \ dt \leq \frac{1}{\varepsilon^{d+1}} \int_0^1 \int_{\mathbb{R}^d} |v_t(y)| d\rho_t(y) \ dt < \infty$$

by (48). By direct calculation $\nabla v_t^{\varepsilon} = [((v_t \rho_t) * \nabla \xi_{\varepsilon})\rho_t^{\varepsilon} - ((v_t \rho_t) * \xi_{\varepsilon})(\rho_t * \nabla \xi_{\varepsilon})]/(\rho_t^{\varepsilon})^2$, so that

$$\begin{split} |\nabla v_t^{\varepsilon}| &\leq \varepsilon^{-1} |(v_t \rho_t) * \nabla \xi_{\varepsilon}| + \varepsilon^{-2} |(v_t \rho_t) * \xi_{\varepsilon}) ||\rho_t * \nabla \xi_{\varepsilon}| \\ &\leq \left[\varepsilon^{-1} \|\nabla \xi_{\varepsilon}\|_{\infty} + \varepsilon^{-2} \|\xi_{\varepsilon}\|_{\infty} \|\nabla \xi_{\varepsilon}\|_{\infty} \rho_t(\Omega) \right] \int_{\Omega} |v_t(x)| \ d\rho_t(x) \end{split}$$

As $t \mapsto \rho_t(\Omega)$ is continuous, the quantity $C(\rho) := \max_t |\rho_t(\Omega)|$ is well defined. Therefore

$$\int_0^1 \operatorname{Lip}(\nu^{\varepsilon}(t,\cdot),\mathbb{R}^d) \ dt = \int_0^1 \sup_{x\in\mathbb{R}^d} |\nabla\nu^{\varepsilon}_t(x)| dt \leq C(\varepsilon)C(\rho) \int_0^1 \int_{\Omega} \nu_t(x)d\rho_t(x) dt < \infty,$$

where the last term is finite by (48). By similar computations and by (48), one can easily show that g^{ε} satisfies (20). We now prove the first estimate in (55). Fix $t \in [0, 1]$. If $\rho_t = 0$, there is nothing to prove. Otherwise we have

$$\begin{split} \int_{\mathbb{R}^d} |v^{\varepsilon}(t,x)|^2 \ \rho^{\varepsilon}_t(x) \ dx &= \int_{\Omega+B_{\varepsilon}(0)} \frac{|(v_t\rho_t) * \xi_{\varepsilon}|^2}{\rho_t * \xi_{\varepsilon} + \varepsilon} \ dx \leq \int_{\Omega+B_{\varepsilon}(0)} \frac{|(v_t\rho_t) * \xi_{\varepsilon}|^2}{\rho_t * \xi_{\varepsilon}} \ dx \\ &= \int_{\mathbb{R}^d} \left| \frac{(v_t\rho_t) * \xi_{\varepsilon}}{\rho_t * \xi_{\varepsilon}} \right|^2 (\rho_t * \xi_{\varepsilon}) \ dx \leq \int_{\mathbb{R}^d} |v(t,x)|^2 \ d\rho_t(x), \end{split}$$

where in the last inequality we used Proposition A.7. Since $v(t, \cdot)$ vanishes in $\mathbb{R}^d \setminus \Omega$, we conclude the first estimate in (55). A similar argument yields the second estimate in (55).

Remark 4.6. Notice that there exist nontrivial $\rho_t \in C_w([0,1]; \mathcal{M}^+(\Omega))$ which solve the continuity equation (14) for v and g that satisfy the bound (48), but such that $\rho_t = 0$ on an open interval in [0,1]. For example, consider $\Omega := [0,1]^2$, v(t,x) := (0,0), $g(t,x) := -(t-1/2)^{-2}\chi_{(0,1/2)}(t)$, $\gamma(t) := (1/2, 1/2)$, $h(t) := \exp(2 - 2(1 - 2t)^{-1})\chi_{(0,1/2)}(t)$. It is easy to check that $\rho_t := h(t)\delta_{\gamma(t)}$ belongs to $C_w([0,1]; \mathcal{M}^+(\Omega))$, solves (14) and (48) holds.

The above is the reason why we add η_{ε} to the definition of ρ_t^{ε} in (54), since otherwise, we could have $\rho_t * \xi = 0$ for some *t*, independently on the chosen mollifier. We remark that the addition of η_{ε} is the main difference to the smoothing results [12, Lemma 8.1.9] and [10, Lemma 3.10], where narrowly continuous measure solutions ρ_t to (14) are smoothed via $\rho_t * \xi$ with ξ being a mollifier.

4.3. Proof of the direct implication of Theorem 4.3

We divide the proof of the direct implication of Theorem 4.3 into two steps: First we construct a measure $\sigma \in \mathcal{M}^+(\mathscr{S}_{\Omega})$ satisfying (49); Then we prove that σ is concentrated on $\mathcal{H}_{\Omega}^{\nu,g}$.

Step 1 - Construction of the measure σ .

Let $V := B_R(0)$, with R > 0 as in (53). For each $0 < \varepsilon < 1$ define $\rho_t^{\varepsilon}, v_t^{\varepsilon}, g_t^{\varepsilon}$ according to (54). By Lemma 4.5 the triple $(\rho_t^{\varepsilon} dx, v_t^{\varepsilon}, g_t^{\varepsilon})$ solves $\partial_t \rho_t^{\varepsilon} dx + \text{div} (v_t^{\varepsilon} \rho_t^{\varepsilon} dx) = g_t^{\varepsilon} \rho_t^{\varepsilon} dx$ in $(0, 1) \times \mathbb{R}^d$ and satisfies the bounds (18, 20, 55). As (48) holds, we can then apply Proposition 2.4 and obtain the representation

$$\rho_t^{\varepsilon} dx = (X_{(\cdot)}^{\varepsilon}(t))_{\#} \Big[R_{(\cdot)}^{\varepsilon}(t) dx \Big], \quad R_x^{\varepsilon}(t) = \rho_0^{\varepsilon}(x) e^{\int_0^t g^{\varepsilon}(s, X_x^{\varepsilon}(s))ds}, \tag{56}$$

where X_x^{ε} and R_x^{ε} are the unique solutions to the ODEs system

$$\begin{cases} \dot{X}_{x}^{\varepsilon}(t) = v^{\varepsilon}(t, X_{x}^{\varepsilon}(t)), \\ X_{x}^{\varepsilon}(0) = x, \end{cases} \begin{cases} \dot{R}_{x}^{\varepsilon}(t) = g^{\varepsilon}(t, X_{x}^{\varepsilon}(t)) \ R_{x}^{\varepsilon}(t), \\ R_{x}^{\varepsilon}(0) = \rho_{0}^{\varepsilon}(x), \end{cases}$$
(57)

for all $t \in [0, 1]$. We define σ^{ε} by duality as

$$\int_{\mathscr{S}_{V}} \varphi(\gamma, h) \ d\sigma^{\varepsilon}(\gamma, h) := \int_{V} \varphi\left(X_{x}^{\varepsilon}, \frac{R_{x}^{\varepsilon}}{\int_{0}^{1} R_{x}^{\varepsilon}(t) \ dt}\right) \left(\int_{0}^{1} R_{x}^{\varepsilon}(t) \ dt\right) \ dx, \tag{58}$$

for all $\varphi \in C_b(\mathscr{S}_V)$. Here we adopted the notation $\varphi(\gamma, h)$ to denote that φ is evaluated on the curve $t \mapsto h(t)\delta_{\gamma(t)}$. We claim that $\sigma^{\varepsilon} \in \mathcal{M}^+(\mathscr{S}_V)$. First, we show that σ_{ε} is welldefined. Indeed, notice that $\rho_t^{\varepsilon} \ge \varepsilon$ in V by construction. Hence by (56) and (20) we estimate

$$\int_{0}^{1} R_{x}^{\varepsilon}(t) \ dt = \rho_{0}^{\varepsilon}(x) \int_{0}^{1} e^{\int_{0}^{t} g^{\varepsilon}(s, X_{x}^{\varepsilon}(s)) ds} \ dt \ge C(\varepsilon) \rho_{0}^{\varepsilon}(x) \ge C(\varepsilon) \varepsilon > 0, \tag{59}$$

for all $x \in V$, where $C(\varepsilon) > 0$ is a constant depending only on ε . Also, by construction, $v^{\varepsilon}(t,x) = 0$ for $x \in \mathbb{R}^d \setminus \overline{\Omega + B_1(0)}$ and $t \in [0,1]$. Therefore from (57) we deduce that $X_x^{\varepsilon}(t) \in V$ for each initial datum $x \in V$ and $0 < \varepsilon < 1$. Thanks to Lemma 3.4, we then obtain that the curve $t \mapsto (\int_0^1 R_x^{\varepsilon}(s) ds)^{-1} R_x^{\varepsilon}(t) \delta_{X_x^{\varepsilon}(t)}$ belongs to \mathscr{S}_V for all $x \in V$. Moreover the map

$$x \mapsto \left(t \mapsto \left(\int_0^1 R_x^{\varepsilon}(s) \ ds \right)^{-1} R_x^{\varepsilon}(t) \delta_{X_x^{\varepsilon}(t)} \right)$$

is continuous from \mathbb{R}^d to (\mathscr{S}_V, d) , which is a consequence of the stability of solutions to (57) with respect to the initial datum $x \in \mathbb{R}^d$, and of the fact that uniform convergence of weights and curves implies d-convergence of measures in \mathscr{S}_V , thanks to (29). This proves that the definition at (58) is well posed. We now estimate the total variation of σ^{ε} . By (54) and standard properties of convolutions we have that $\|\rho_t^{\varepsilon} dx\|_{\mathcal{M}(V)} \leq$ $\|\rho_t\|_{\mathcal{M}(\Omega)} + \varepsilon |V|$ for all $t \in [0, 1], \varepsilon \in (0, 1)$. Hence, by testing σ^{ε} against $\varphi \equiv 1$ and using (56) we infer

$$\|\sigma^{\varepsilon}\|_{\mathcal{M}(\mathscr{S}_{V})} \leq \int_{V} \int_{0}^{1} R_{x}^{\varepsilon}(t) dt dx = \int_{0}^{1} \|\rho_{t}^{\varepsilon} dx\|_{\mathcal{M}(V)} dt \leq \|\rho\|_{\mathcal{M}(X_{\Omega})} + \varepsilon |V|.$$
(60)

Moreover $\sigma^{\varepsilon} \geq 0$ by (59), showing that $\sigma_{\varepsilon} \in \mathcal{M}^+(\mathscr{G}_V)$. We also remark that σ^{ε} is concentrated on \mathcal{H}_V , given that the curve $t \mapsto (\int_0^1 R_x^{\varepsilon}(s) ds)^{-1} R_x^{\varepsilon}(t) \delta_{X_x^{\varepsilon}(t)}$ belongs to \mathcal{H}_V for each $x \in V$, thanks to the regularity of solutions to (57).

We now show that the family σ^{ε} is tight as $0 < \varepsilon < 1$, by proving that

$$\sup_{0<\varepsilon<1} \int_{\mathscr{S}_{V}} \mathscr{F}(\gamma,h) \ d\sigma^{\varepsilon}(\gamma,h) < \infty, \tag{61}$$

where $\mathscr{F}: (\mathscr{S}_V, d) \to [0, \infty]$ is the functional defined at (31): Indeed assume that (61) holds; by Proposition 3.10 we know that \mathscr{F} is d-measurable and its sublevels are compact. Moreover (\mathscr{S}_V, d) is a complete separable metric space (see Proposition 3.6). Thus we can apply Proposition A.1 to conclude tightness for σ^{ε} . Let us proceed with the proof of (61). First notice that (58) can be tested against \mathscr{F} , as $\sigma^{\varepsilon} \ge 0$ and \mathscr{F} is lower semi-continuous with respect to the metric d (Proposition 3.10). Since σ^{ε} is concentrated on \mathcal{H}_V , by formula (37) and one-homogeneity of \mathscr{F} with respect to h we have

$$\int_{\mathscr{S}_{V}} \mathscr{F}(\gamma, h) \ d\sigma^{\varepsilon}(\gamma, h) = \int_{\mathscr{S}_{V}} \int_{\{h>0\}} \frac{\beta}{2} \ |\dot{\gamma}(t)|^{2} h(t) + \frac{\beta \delta^{2}}{2} \ \frac{|\dot{h}(t)|^{2}}{h(t)} + \alpha h(t) \ dt \ d\sigma^{\varepsilon}(\gamma, h).$$
(62)

By (56, 57) and (55) we estimate

$$\begin{split} \int_{\mathscr{S}_{V}} \int_{0}^{1} |\dot{\gamma}(t)|^{2} h(t) \ dt \ d\sigma^{\varepsilon} &= \int_{V} \int_{0}^{1} |\dot{X}_{x}^{\varepsilon}(t)|^{2} \ R_{x}^{\varepsilon}(t) \ dt \ dx = \int_{0}^{1} \int_{V} |v^{\varepsilon}(t, X_{x}^{\varepsilon}(t))|^{2} \ R_{x}^{\varepsilon}(t) \ dx \ dt \\ &= \int_{0}^{1} \int_{\mathbb{R}^{d}} |v^{\varepsilon}(t, x)|^{2} \ \rho_{t}^{\varepsilon}(x) \ dx \ dt \leq \int_{0}^{1} \int_{\Omega} |v(t, x)|^{2} \ d\rho_{t}(x) \ dt, \end{split}$$

and, in a similar fashion,

$$\begin{split} \int_{\mathscr{S}_{V}} \int_{\{h>0\}} \frac{|\dot{h}(t)|^{2}}{h(t)} \ dt \ d\sigma^{\varepsilon} &= \int_{V} \int_{\{h>0\}} \frac{|\dot{R}_{x}^{\varepsilon}(t)|^{2}}{R_{x}^{\varepsilon}(t)} \ dt \ dx = \int_{\{h>0\}} \int_{V} |g^{\varepsilon}(t, X_{x}^{\varepsilon}(t))|^{2} \ R_{x}^{\varepsilon}(t) \ dx \ dt \\ &= \int_{\{h>0\}} \int_{\mathbb{R}^{d}} |g^{\varepsilon}(t, x)|^{2} \ \rho_{t}^{\varepsilon}(x) \ dx \ dt \leq \int_{0}^{1} \int_{\Omega} |g(t, x)|^{2} \ d\rho_{t}(x) \ dt. \end{split}$$

Finally, by (60),

$$\int_{\mathscr{S}_{V}} \int_{0}^{1} h(t) dt d\sigma^{\varepsilon} = \int_{V} \int_{0}^{1} R_{x}^{\varepsilon}(t) dt dx \leq \|\rho\|_{\mathcal{M}(X_{\Omega})} + \varepsilon |V|.$$
(63)

From the above estimates, and (48, 62), we conclude (61), proving that $\{\sigma^{\varepsilon}\}_{\varepsilon}$ is tight. Since $\{\sigma_{\varepsilon}\}_{\varepsilon}$ is uniformly bounded by (60), we can apply the compactness result [44, Theorem 8.6.2] to infer the existence of $\sigma \in \mathcal{M}^+(\mathscr{S}_V)$ such that $\sigma_{\varepsilon} \to \sigma$ narrowly as $\varepsilon \to 0$. In particular, as \mathscr{F} is d-lower semi-continuous, $\mathscr{F} \ge 0$ and (61) holds, we can apply (115) in Proposition A.2 to infer $\int_{\mathscr{S}_V} \mathscr{F}(\gamma, h) d\sigma(\gamma, h) < \infty$. From the latter, we see that σ is concentrated on the domain of \mathscr{F} , that is, on the set \mathcal{H}_V (Proposition 3.10).

We now prove that σ satisfies the representation formula (49). To this end, let $\varphi \in C_c(X_V)$ and define the map $\Psi(\gamma, h) := \int_0^1 h(t)\varphi(t, \gamma(t)) dt$ for $\rho = h\delta_{\gamma} \in \mathscr{S}_V$. We claim that σ^{ε} according to (58) can be tested against Ψ : indeed, first notice that Ψ is d-continuous. This is because the map $(\gamma, h) \mapsto h(t)\varphi(t, \gamma(t))$ is continuous for t fixed, by Lemma 3.7; if $d(\rho^n, \rho) \to 0$, then $\|h_n - h\|_{\infty} \leq d(\rho^n, \rho)$ by (29), so that $\{h_n\}_n$ is uniformly bounded; thus by dominated convergence we conclude continuity for Ψ , since φ is bounded, and since $\varphi(t, \gamma_n(t)) \to \varphi(t, \gamma(t))$ when h(t) > 0. Moreover, thanks to (63), we can estimate

$$\int_{V} \left| \Psi \left(X_{x}^{\varepsilon}, \frac{R_{x}^{\varepsilon}}{\int_{0}^{1} R_{x}^{\varepsilon}(t) \ dt} \right) \right| \left(\int_{0}^{1} R_{x}^{\varepsilon}(t) \ dt \right) \ dx \le \|\varphi\|_{\infty} (\|\rho\|_{\mathcal{M}(X_{\Omega})} + \varepsilon |V|), \tag{64}$$

showing that the right-hand side of (58) tested against $|\Psi|$ is finite. The fact that σ^{ε} can be tested against Ψ follows immediately. By (58), the latter yields

$$\int_{\mathscr{S}_{V}} \Psi(\gamma, h) \ d\sigma^{\varepsilon}(\gamma, h) = \int_{V} \int_{0}^{1} \varphi(t, X_{x}^{\varepsilon}(t)) R_{x}^{\varepsilon}(t) \ dt \ dx = \int_{0}^{1} \int_{V} \varphi(t, x) \ \rho_{t}^{\varepsilon}(x) \ dx \ dt, \quad (65)$$

where in the last equality we used (56). We want to pass to the limit as $\varepsilon \to 0$ in (65). Notice that the right-hand side passes to the limit since $dt \otimes \rho_t^\varepsilon dx \longrightarrow^* dt \otimes \rho_t$ in $\mathcal{M}(X_V)$: Indeed $\rho_t^\varepsilon dx \longrightarrow \rho_t$ narrowly in $\mathcal{M}(V)$ for all t (Lemma 4.5) and $\rho_t^\varepsilon dx$ is uniformly bounded in $\mathcal{M}(V)$, as previously shown. Concerning the left-hand side of (65), we first claim that the map $|\Psi|$ is uniformly integrable with respect to σ^ε according to definition (116). To this end, for k > 0 define $A_k := \{(\gamma, h) \in \mathscr{S}_V : |\Psi(\gamma, h)| \ge k\}$. By the definition of σ^ε and by (63) we get

$$\begin{split} \int_{A_k} |\Psi(\gamma,h)| \ d\sigma^{\varepsilon}(\gamma,h) &\leq \frac{1}{k} \int_{\mathscr{S}_V} |\Psi(\gamma,h)|^2 \ d\sigma^{\varepsilon}(\gamma,h) &\leq \frac{\|\varphi\|_{\infty}^2}{k} \int_V \int_0^1 R_x^{\varepsilon}(t) \ dt \ dx \\ &\leq \frac{\|\varphi\|_{\infty}^2}{k} (\|\rho\|_{\mathcal{M}(X_{\Omega})} + |V|), \end{split}$$

concluding uniform integrability for $|\Psi|$. Therefore we can invoke (117) and pass to the limit as $\varepsilon \to 0$ in the left-hand side of (65). After one application of Fubini's Theorem we obtain

$$\int_{0}^{1} \int_{\mathscr{S}_{V}} h(t)\varphi(t,\gamma(t)) \ d\sigma(\gamma,h) \ dt = \int_{0}^{1} \int_{V} \varphi(t,x) \ d\rho_{t}(x) \ dt, \qquad \text{for all} \quad \varphi \in C_{c}(X_{V}).$$
(66)

We claim that (49) descends from (66). In order to show it, we first derive a pointwise in time version of (66). We start by showing that $\Theta(t) := \int_{\mathscr{S}_V} h(t)\varphi(t,\gamma(t)) \, d\sigma(\gamma,h)$ is continuous for all $\varphi \in C_c(X_V)$ fixed. Indeed, the map $t \mapsto h(t)\varphi(t,\gamma(t))$ is continuous for each fixed $(\gamma,h) \in \mathscr{S}_V$, by Lemma 3.4. Moreover, by recalling that σ^{ε} is concentrated on solutions of (57), and by arguing as in the proof of (51), we can show that for all ε it holds that

$$\int_{\mathscr{S}_{V}} \|h\|_{\infty} \ d\sigma^{\varepsilon}(\gamma,h) \leq \int_{\mathscr{S}_{V}} \int_{0}^{1} h(t)(1+|g^{\varepsilon}(t,\gamma(t))|) \ dt \ d\sigma^{\varepsilon}(\gamma,h).$$

Therefore, by employing (55, 56, 58, 63), and setting $C := \|\rho\|_{\mathcal{M}(X_{\Omega})} + |V|$, we obtain

$$\begin{split} \int_{\mathscr{S}_{V}} \|h\|_{\infty} \ d\sigma^{\varepsilon}(\gamma,h) &\leq C + \int_{V} \int_{0}^{1} R_{x}^{\varepsilon}(t) |g^{\varepsilon}(t,X_{x}^{\varepsilon}(t))| \ dt \ dx = C + \int_{0}^{1} \int_{V} |g^{\varepsilon}(t,x)| \ \rho_{t}^{\varepsilon}(x) \ dx \ dt \\ &\leq C + \left(\int_{0}^{1} \|\rho_{t}^{\varepsilon} \ dx\|_{\mathcal{M}(V)} \ dt \right)^{1/2} \left(\int_{0}^{1} \int_{V} |g^{\varepsilon}(t,x)|^{2} \ \rho_{t}^{\varepsilon}(x) \ dx \ dt \right)^{1/2} \\ &\leq C + C^{1/2} \left(\int_{0}^{1} \int_{V} |g(t,x)|^{2} \ d\rho_{t}(x) \ dt \right)^{1/2}, \end{split}$$

and the last term is bounded by assumption (48). Finally, the map $(\gamma, h) \in \mathscr{S}_V \mapsto ||h||_{\infty}$ is d-continuous and non-negative, therefore by the narrow convergence $\sigma^{\varepsilon} \to \sigma$ and (115) we infer $\int_{\mathscr{G}_V} ||h||_{\infty} d\sigma(\gamma, h) < \infty$. By dominated convergence we then conclude continuity of Θ . As a byproduct of this argument, we have additionally shown that $\sigma \in \mathcal{M}_1^+(\mathscr{S}_V)$. Notice that also the map $t \mapsto \int_V \varphi(t, x) d\rho_t(x)$ is continuous, as a consequence of the narrow continuity of $t \mapsto \rho_t$. Testing (66) against $\varphi(t, x) := a(t)b(x)$ for $a \in C_c(0, 1), b \in C(V)$, yields

$$\int_{\mathscr{S}_{V}} h(t)b(\gamma(t)) \ d\sigma(\gamma,h) = \int_{V} b(x) \ d\rho_{t}(x), \qquad \text{for all} \qquad b \in C(V), \quad t \in [0,1].$$
(67)

Fix $t \in [0, 1]$ and $b \in C(V)$ such that b = 0 in Ω and b > 0 in $V \setminus \Omega$. Recalling that ρ_t is concentrated on Ω , from (67) we obtain a set $E_t \subset \mathscr{S}_V$ such that $\sigma(\mathscr{S}_V \setminus E_t) = 0$ and $h(t)b(\gamma(t)) = 0$ for all $(\gamma, h) \in E_t$. In particular, by definition of b,

$$\gamma(t) \in \Omega \quad \text{if} \quad h(t) > 0,$$
 (68)

for all $(\gamma, h) \in E_t$. Let $Q \subset [0, 1]$ be a dense countable subset and define $E := \bigcap_{t \in Q} E_t$, so that $\sigma(\mathscr{G}_V \setminus E) = 0$ and (68) holds for all $(\gamma, h) \in E$, $t \in Q$, that is, $\gamma(\{h > 0\} \cap Q) \subset \Omega$ for σ -a.e. $(\gamma, h) \in \mathscr{G}_V$. By density of Q and continuity of h, γ we deduce $\gamma(\{h > 0\}) \subset \Omega$ Ω for σ -a.e. $(\gamma, h) \in \mathscr{G}_V$, from which we conclude concentration of σ on \mathscr{G}_Ω . Since we already showed that σ is concentrated on \mathcal{H}_V we also infer that σ is concentrated on \mathcal{H}_Ω . It is immediate to check that \mathscr{G}_Ω is d-closed in \mathscr{G}_V , and hence d-measurable. Therefore we can restrict σ to \mathscr{S}_{Ω} to obtain a measure in $\mathcal{M}_{1}^{+}(\mathscr{S}_{\Omega})$ satisfying (49), as claimed.

Step 2 – σ is concentrated on $\mathcal{H}^{\nu,g}_{\Omega}$.

So far we have constructed a measure $\sigma \in \mathcal{M}_1^+(\mathscr{S}_\Omega)$ concentrated on \mathcal{H}_Ω and satisfying (49). We now prove that σ is concentrated on $\mathcal{H}_\Omega^{v,g}$, i.e., that the ODEs (O1)-(O2) hold for σ -a.e. (γ, h) in \mathcal{H}_Ω . This claim follows from two preliminary estimates, whose proof we postpone for a moment: For $\bar{\nu} \in C_c(X_\Omega; \mathbb{R}^d), \bar{g} \in C_c(X_\Omega)$ and any $\varphi \in C_c^1(0, 1)$, there exists a constant C > 0 depending only on φ and on the radius of V, such that

$$\int_{\mathscr{S}_{\Omega}} \left| \int_{0}^{1} h(t)\varphi'(t) + h(t) \ \bar{g}(t,\gamma(t))\varphi(t) \ dt \right| \ d\sigma \leq C \int_{X_{\Omega}} |\bar{g} - g| \ d\rho_t \ dt, \tag{69}$$

$$\int_{\mathscr{S}_{\Omega}} \left| \int_{0}^{1} h(t)\gamma(t) \cdot \varphi'(t) + \Psi_{\bar{\nu},\bar{g}}(t,\gamma,h) \cdot \varphi(t) dt \right| d\sigma \leq C \int_{X_{\Omega}} |\bar{\nu}-\nu| + |\bar{g}-g| d\rho_{t} dt,$$
(70)

where $\Psi_{\bar{\nu},\bar{g}}(t,\gamma,h) := h(t)\gamma(t) \ \bar{g}(t,\gamma(t)) + h(t) \ \bar{\nu}(t,\gamma(t))$. We start by showing (O2). By the energy bound (48) and Hölder's inequality, we can find two sequences $\{\nu_n\}_n$ in $C_c(X_{\Omega}; \mathbb{R}^d)$ and $\{g_n\}_n$ in $C_c(X_{\Omega})$ converging to ν and g in $L^1_{\rho}(X_{\Omega})$, respectively. By (49) and Remark 4.2 we get

$$\int_{\mathscr{S}_{\Omega}} \left| \int_{0}^{1} h(t)(g(t,\gamma(t)) - g_{n}(t,\gamma(t)))\varphi(t) dt \right| d\sigma \leq \|\varphi\|_{\infty} \int_{X_{\Omega}} |g_{n} - g| d\rho_{t} dt.$$
(71)

Hence by (69) with $\bar{g} := g_n$, (71), and triangle inequality we get

$$\int_{\mathscr{S}_{\Omega}} \left| \int_{0}^{1} h(t) \varphi'(t) + h(t)g(t,\gamma(t))\varphi(t) dt \right| d\sigma \leq C \int_{X_{\Omega}} |g_{n} - g| d\rho_{t} dt \to 0,$$

as $n \to \infty$, for every test function $\varphi \in C_c^1(X_\Omega)$. Therefore

$$\int_0^1 h(t)\varphi'(t) + h(t)g(t,\gamma(t))\varphi(t) \, dt = 0 \quad \text{for all} \quad (\gamma,h) \in E_{\varphi}, \tag{72}$$

where $\sigma(\mathscr{S}_{\Omega} \setminus E_{\varphi}) = 0$ and E_{φ} depends on φ . Let $D \subset C_{c}^{1}(0, 1)$ be a dense countable set and $E := \bigcap_{\varphi \in D} E_{\varphi}$, so that $\sigma(\mathscr{S}_{\Omega} \setminus E) = 0$ and (72) holds for all $\varphi \in D$ and $(\gamma, h) \in E$. Consider $\varphi(t, x) := 1 + |v(t, x)| + |g(t, x)|$ and notice that $\varphi(t, \cdot) \in L^{1}_{\rho_{t}}(\Omega)$ for a.e. $t \in (0, 1)$, thanks to (48) and narrow continuity of ρ_{t} . Hence, by Remark 4.2 applied to φ , we conclude that σ satisfies (50). Therefore there exists a set F with $\sigma(\mathscr{S}_{\Omega} \setminus F) = 0$ and such that

$$\int_{0}^{1} h(t)(1+|\nu(t,\gamma(t))|+|g(t,\gamma(t))|) \, dt < \infty$$
(73)

for all $(\gamma, h) \in F$. Consider now $\varphi \in C_c^1(0, 1)$ and $\varphi_n \in D$ such that $\varphi_n \to \varphi$ in $C_c^1(0, 1)$. As a consequence of (72), for any (γ, h) in $E \cap F$ we have

$$\left| \int_{0}^{1} h(t)\varphi'(t) + h(t)g(t,\gamma(t))\varphi(t) \, dt \right| \leq \|\varphi_{n} - \varphi\|_{C^{1}} \int_{0}^{1} h(t)(1 + |g(t,\gamma(t))|) \, dt \to 0,$$

as $n \to \infty$, so that (72) holds for all $\varphi \in C_c^1(0,1)$ and $(\gamma, h) \in E \cap F$. Therefore

$$\dot{h} = g(t, \gamma(t))h(t) \tag{74}$$

in the sense of distributions for σ -a.e. $(\gamma, h) \in \mathscr{S}_{\Omega}$. Since σ is concentrated on \mathcal{H}_{Ω} , the distributional formulation of (74) coincides with the a.e. one, so that (O2) holds. We now prove (O1), which follows by similar arguments. First, by (49) and Remark 4.2, we estimate

$$\int_{\mathscr{S}_{\Omega}} \left| \int_{0}^{1} (\Psi_{\nu,g} \quad (t,\gamma,h) - \Psi_{\nu_{n},g_{n}}(t,\gamma,h)) \cdot \varphi(t) \, dt \right| \, d\sigma$$

$$\leq \|\varphi\|_{\infty} \max\{1,R\} \int_{X_{\Omega}} |\nu_{n} - \nu| + |g_{n} - g| \, d\rho_{t} \, dt,$$

where *R* is as in (53). By applying (70) to $\bar{v} := v_n$, $\bar{g} := g_n$ and by triangle inequality we infer

$$\int_{\mathscr{S}_{\Omega}} \left| \int_{0}^{1} h(t) \gamma(t) \cdot \varphi'(t) + \Psi_{\nu,g}(t,\gamma,h) \cdot \varphi(t) dt \right| d\sigma \leq C \int_{X_{\Omega}} |\nu_n - \nu| + |g_n - g| d\rho_t dt \to 0,$$

as $n \to \infty$, for all $\varphi \in C_c^1((0,1); \mathbb{R}^d)$. By reasoning as above, we can find a countable dense subset \tilde{D} of $C_c^1((0,1); \mathbb{R}^d)$ and a set \tilde{E} with $\sigma(\mathscr{S}_{\Omega} \setminus \tilde{E}) = 0$ such that

$$\int_{0}^{1} h(t)\gamma(t) \cdot \varphi'(t) + \Psi_{\nu,g}(t,\gamma,h) \cdot \varphi(t) \ dt = 0$$
(75)

for all $\varphi \in \tilde{D}$, $(\gamma, h) \in \tilde{E}$. By (75) and (73), we conclude that

$$(h(t)\gamma(t))' = h(t)\gamma(t)g(t,\gamma(t)) + h(t)\nu(t,\gamma(t))$$

in the sense of distributions for all $(\gamma, h) \in \tilde{E} \cap F$. Recall that (74) holds in the sense of distributions in $E \cap F$. Moreover σ is concentrated on \mathcal{H}_{Ω} , whose elements satisfy $h\gamma \in AC^2([0, 1]; \mathbb{R}^d)$ and the product rule holds (see Lemma 3.8). Hence we can find a set \tilde{F} such that $\sigma(\mathscr{S}_{\Omega} \setminus \tilde{F}) = 0$, and that (O2) and

$$h(t)\dot{\gamma}(t) = h(t)v(t,\gamma(t))$$
 for a.e. $t \in (0,1)$,

hold for σ -a.e. $(\gamma, h) \in \mathcal{H}_{\Omega}$. This establishes (O1).

We are left to show (69)-(70). We start by proving (69). First notice that the map

$$\phi(\gamma, h) := \left| \int_0^1 h(t) \varphi'(t) + h(t) \overline{g}(t, \gamma(t)) \varphi(t) \, dt \right| \tag{76}$$

in the left-hand side of (69) is d-continuous: indeed, if $d(\rho^n, \rho) \to 0$, from (29) we have $\|h_n - h\|_{\infty} \leq d(\rho^n, \rho)$, so that h_n is uniformly bounded; by Lemma 3.7 it follows that $(\gamma, h) \mapsto h(t)\varphi'(t) + h(t)\bar{g}(t,\gamma(t))\varphi(t)$ is d-continuous for every *t*; thus continuity of ϕ follows by dominated convergence. Extend \bar{g} to zero outside of Ω and set $\bar{g}^{\varepsilon} := [(\bar{g}\rho_t) * \xi_{\varepsilon}]/\rho_t^{\varepsilon}$, with ξ_{ε} as in (54). As ϕ is continuous, we can test (58) against ϕ , integrate by parts, and use (56)–(57) to get

$$\begin{aligned} \int_{\mathscr{S}_{V}} \left| \int_{0}^{1} 1h(t)\varphi'(t) + h(t)\bar{g}(t,\gamma(t))\varphi(t) dt \right| d\sigma^{\varepsilon}(\gamma,h) \\ &\leq \|\varphi\|_{\infty} \int_{X_{V}} |\bar{g} - g^{\varepsilon}| \ \rho_{t}^{\varepsilon} \ dx \ dt \leq \|\varphi\|_{\infty} \int_{X_{V}} |\bar{g}^{\varepsilon} - g^{\varepsilon}| + |\bar{g}^{\varepsilon} - \bar{g}| \ \rho_{t}^{\varepsilon} \ dx \ dt. \end{aligned}$$
(77)

We then estimate each term separately. First, recalling (54),

$$\int_{X_{V}} |\bar{g}^{\varepsilon} - g^{\varepsilon}| \rho_{t}^{\varepsilon} dx dt = \int_{0}^{1} \int_{\Omega + B_{\varepsilon}(0)} |((g - \bar{g})\rho_{t}) * \xi_{\varepsilon}| dx dt \leq \int_{X_{\Omega}} |g - \bar{g}| d\rho_{t} dt, \quad (78)$$

by standard estimates on convolutions of measures. Second,

$$\begin{aligned} \int_{X_{V}} |\bar{g}^{\varepsilon} - \bar{g}| & \rho_{t}^{\varepsilon} dx dt = \int_{X_{V}} |(\bar{g}\rho_{t}) * \xi_{\varepsilon} - \bar{g}(\rho_{t} * \xi_{\varepsilon} + \varepsilon)| dx dt \\ & \leq \int_{X_{V}} |(\bar{g}\rho_{t}) * \xi_{\varepsilon} - \bar{g}(\rho_{t} * \xi_{\varepsilon})| dx dt + \varepsilon \int_{X_{V}} |\bar{g}| dx dt. \end{aligned}$$

$$\tag{79}$$

The second term in (79) converges to zero as $\varepsilon \to 0$. Moreover, by the uniform continuity of \bar{g} , for each $\zeta > 0$ there exists $\tilde{\zeta} > 0$ such that $|\bar{g}(t,y) - \bar{g}(t,x)| < \zeta$ whenever $|x - y| \leq \tilde{\zeta}, t \in [0, 1]$. Therefore, for all $\varepsilon < \tilde{\zeta}$ and $x \in V$ we have

$$|(\bar{g}\rho_t)*\xi_{\varepsilon}-\bar{g}(\rho_t*\xi_{\varepsilon})|(x)\leq \int_{B_{\varepsilon}(x)}|\bar{g}(t,y)-\bar{g}(t,x)|\xi_{\varepsilon}(x-y)\ d\rho_t(y)\leq \zeta\ (\rho_t*\xi_{\varepsilon})(x),$$

from which we infer

$$\int_{X_V} |(\bar{g}\rho_t) * \xi_\varepsilon - \bar{g}(\rho_t * \xi_\varepsilon)| \ dx \ dt \leq \zeta \int_{X_V} \rho_t * \xi_\varepsilon \ dx \ dt \leq \zeta \|\rho\|_{\mathcal{M}(X_\Omega)}$$

As ζ is arbitrary, by (79) we conclude that $\int_{X_V} |\bar{g}^{\varepsilon} - \bar{g}| \rho_t^{\varepsilon} dx dt \to 0$. Thus, from (77) to (78)

$$\limsup_{\varepsilon\to 0} \int_{\mathscr{S}_{V}} \left| \int_{0}^{1} h(t) \varphi'(t) + h(t) \bar{g}(t, \gamma(t)) \varphi(t) dt \right| d\sigma^{\varepsilon}(\gamma, h) \leq \|\varphi\|_{\infty} \int_{X_{\Omega}} |g - \bar{g}| d\rho_{t} dt.$$

As $\sigma^{\varepsilon} \to \sigma$ narrowly, σ is concentrated on \mathscr{S}_{Ω} , ϕ at (76) is continuous and nonnegative, by (115) we conclude (69). We now show (70). First notice that the function in the left integral of (70) is d-continuous, a fact that can be proven similarly to (76). Set $C := (R+1) \|\phi\|_{\infty}$, with *R* as in (53). Similarly to the above proof of (69), we can integrate by parts and make use of (56)–(58), and estimate

$$\begin{split} \int_{\mathscr{S}_{V}} \left| \int_{0}^{1} h(t)\gamma(t) \cdot \varphi'(t) + \Psi_{\bar{\nu},\bar{g}}(t,\gamma,h) \cdot \varphi(t) \, dt \, \left| d\sigma^{\varepsilon}(\gamma,h) \right| \\ &\leq C \int_{X_{V}} R_{x}^{\varepsilon}(t) \left[\left| g^{\varepsilon}(t,X_{x}^{\varepsilon}(t)) - \bar{g}(t,X_{x}^{\varepsilon}(t)) \right| + \left| v^{\varepsilon}(t,X_{x}^{\varepsilon}(t)) - \bar{\nu}(t,X_{x}^{\varepsilon}(t)) \right| \right] \, dt \, dx \qquad (80) \\ &= C \int_{X_{V}} \left| g^{\varepsilon} - \bar{g} \right| + \left| v^{\varepsilon} - \bar{\nu} \right| \, \rho_{t}^{\varepsilon} \, dx \, dt \leq C \int_{X_{V}} \left| g - \bar{g} \right| + \left| v - \bar{\nu} \right| \, \rho_{t}^{\varepsilon} \, dx \, dt + o(1), \end{split}$$

where in the last inequality we employed (78) and the convergence $\int_{X_V} |\bar{g}^{\varepsilon} - \bar{g}| \rho_t^{\varepsilon} dx dt \to 0$ to estimate the first term, and similar estimates involving v for the second, and where $o(1) \to 0$ as $\varepsilon \to 0$. By passing to the limes superior in (80) and by recalling that $\sigma^{\varepsilon} \to \sigma$ narrowly, we can invoke (115) and obtain (70).

Step 3 – σ is concentrated on \mathcal{H}^1_{Ω} .

We are left to prove that σ is concentrated on $\mathcal{H}_{\Omega}^{1} = \{(h, \gamma) \in \mathscr{S}_{\Omega} : ||h||_{1} = 1\}$. By definition, the measure $\sigma^{\varepsilon} \in \mathcal{M}^{+}(\mathscr{S}_{V})$ introduced at (58) is concentrated on $\mathcal{H}_{V}^{1} = \{(h, \gamma) \in \mathscr{S}_{V} : ||h||_{1} = 1\}$, where *V* is as in (53). Also recall that we have proven $\sigma^{\varepsilon} \to \sigma$ narrowly. Moreover note that the set \mathcal{H}_{V}^{1} is closed in \mathscr{S}_{V} , as an immediate consequence of (29). Hence, by (115), we get $\sigma(\mathscr{S}_{V} \setminus \mathcal{H}_{V}^{1}) \leq \liminf_{\varepsilon \to 0} \sigma^{\varepsilon}(\mathscr{S}_{V} \setminus \mathcal{H}_{V}^{1}) = 0$, showing that σ is concentrated on \mathcal{H}_{V}^{1} . As σ is concentrated on \mathscr{S}_{Ω} , we conclude. This ends the proof of Theorem 4.3.

Remark 4.7. As mentioned in the introduction, it would be interesting to extend Theorem 4.3 to the case of $\Omega = \mathbb{R}^d$. Notice however that our construction of the measure σ is heavily reliant on the boundedness of Ω : first, such assumption is needed in proving compactness of the sublevels of the functional \mathscr{F} (see (40) and estimates after), which in turn allows to show tightness for the family σ^{ε} (see (61) and argument immediately after); second, boundedness of Ω is employed to provide the uniform bound (60) on the norm of σ^{ε} . These arguments are crucial to obtain compactness for σ^{ε} and, consequently, the representing measure σ as their limit.

5. Uniqueness of characteristics and uniqueness for the PDE

The aim of this section is to apply Theorem 4.3 to relate uniqueness of the characteristics with uniqueness of solutions for the continuity equation with given initial data and minimal total variation. Throughout the section $\Omega \subset \mathbb{R}^d$ with $d \ge 1$ is the closure of a bounded domain. We denote $X_{\Omega} := (0, 1) \times \Omega$. Moreover \mathscr{S}_{Ω} denotes the set defined at (26), equipped with the distance d at (28). We remind the reader that (\mathscr{S}_V, d) is a complete metric space (Proposition 3.6). Let $v : X_{\Omega} \to \mathbb{R}^d$ and $g : X_{\Omega} \to \mathbb{R}$ be measurable maps and recall the definition of $\mathcal{H}_{\Omega}^{v,g}$ at (44), that is, the set of regular characteristics of the ODEs system (O1)-(O2). Also recall the definition of \mathcal{H}_{Ω}^1 at (45) Finally, we define the following set

$$\mathcal{D}_{\nu,g} := \left\{ (t \mapsto \rho_t) \in C_w([0,1]; \mathcal{M}^+(\Omega)) : (\rho_t, \nu_t, g_t) \text{ satisfy } \partial_t \rho_t + \text{div } (\nu \rho_t) = g \rho_t \text{ and } (48) \right\}.$$

We will prove the following result:

Theorem 5.1. Let $A \subset \Omega$ be a measurable set. Suppose that:

(Hyp) For each $x \in A$ the solution $(\gamma, h) \in \mathcal{H}_{\Omega}^{\nu,g}$ to (O1)-(O2) with initial value (x, 1) is unique in $[0, \tau)$ for every $\tau \in (0, 1)$ such that $[0, \tau) \subset \{h > 0\}$, i.e., if $(\gamma_1, h_1), (\gamma_2, h_2) \in \mathcal{H}_{\Omega}^{\nu,g}$ solve (O1)-(O2) in $[0, \tau)$ with initial data (x, 1) and $h_1 > 0$, $h_2 > 0$ in $[0, \tau)$, then $h_1 = h_2$ and $\gamma_1 = \gamma_2$ in $[0, \tau)$.

Then, for any initial data $\rho_0 \in \mathcal{M}^+(\Omega)$ concentrated on A, the continuity equation $\partial_t \rho_t + \operatorname{div}(v\rho_t) = g\rho_t$ admits at most one solution $\rho \in \mathcal{D}_{v,g}$ with initial data ρ_0 and such that

$$\|\rho\|_{\mathcal{M}(X_{\Omega})} \le \|\tilde{\rho}\|_{\mathcal{M}(X_{\Omega})} \quad \text{for all} \quad \tilde{\rho} \in \mathcal{D}_{\nu,g} \quad \text{such that} \quad \tilde{\rho}_0 = \rho_0.$$
(81)

In the next section we provide several auxiliary lemmas and definitions, which will be instrumental in proving Theorem 5.1. The proof of Theorem 5.1 will be carried out in Section 5.2.

5.1. Auxiliary results

Define the following subset of \mathscr{S}_{Ω} :

$$\mathscr{S}_{\Omega}^* := \{ (\gamma, h) \in \mathscr{S}_{\Omega} : \{ h > 0 \} = [0, 1] \cap (-\infty, \tau) \text{ for some } \tau \in \mathbb{R} \}.$$
(82)

The first step is to prove that condition (81) implies that the measure σ obtained by Theorem 4.3 is concentrated on \mathscr{G}^*_{Ω} . To this aim, we define a cutoff operator on the space \mathscr{G}_{Ω} .

Definition 5.2 (Cutoff operator). Define the vanishing time map $\tau : \mathscr{S}_{\Omega} \to [0, \infty]$ as

$$\tau(\gamma, h) := \begin{cases} \arg \min \{t \in [0, 1] : h(t) = 0\} & \text{if } \{t \in [0, 1] : h(t) = 0\} \neq \emptyset, \\ \infty & \text{otherwise,} \end{cases}$$

and the cutoff operator $G: \mathscr{S}_{\Omega} \to \mathscr{S}_{\Omega}$ as $G(\gamma, h) := (\gamma, h \chi_{[0, \tau(h, \gamma))}).$

Lemma 5.3. Let τ and G be as in in Definition 5.2. Then τ is lower semi-continuous and G is measurable. Moreover, $G(\mathscr{S}_{\Omega}) \subset \mathscr{S}_{\Omega}^*$, $G(\mathcal{H}_{\Omega}^{\nu,g}) \subset \mathcal{H}_{\Omega}^{\nu,g}$, and the set \mathscr{S}_{Ω}^* is measurable.

Proof. We start by proving that τ is lower semi-continuous. Assume that $d(\rho^n, \rho) \to 0$ as $n \to \infty$. Set $\tau_n := \tau(\rho^n), \tau := \tau(\rho)$. Without loss of generality we can suppose that $\tau^* := \lim_{n\to\infty} \tau_n = \lim_{n\to\infty} \inf_{n\to\infty} \tau_n < \infty$. Thus, by definition, we have that $h_n(\tau_n) = 0$ for n sufficiently large. Therefore

$$h(\tau^*) \le |h(\tau^*) - h(\tau_n)| + |h(\tau_n) - h_n(\tau_n)|.$$
(83)

Notice that $||h_n - h||_{\infty} \to 0$ by (29) and $d(\rho^n, \rho) \to 0$. Hence the second term in (83) converges to zero as $n \to \infty$. Thanks to the continuity of h and the convergence $\tau_n \to \tau^*$ also the first term in (83) is infinitesimal, concluding that $h(\tau^*) = 0$. Thus $\tau \leq \tau^*$ by minimality, from which lower semi-continuity follows.

We now show that G is measurable by constructing measurable maps $G_n : \mathscr{S}_{\Omega} \to \mathscr{S}_{\Omega}$ such that $G_n(\rho) \to G(\rho)$ for all $\rho \in \mathscr{S}_{\Omega}$. Indeed this immediately implies measurability of G (see, e.g., [44, Corollary 6.2.6]). To this end, define the continuous maps $\varphi_n : \mathbb{R} \to [0, 1]$ by setting

$$\varphi_n(t) := \chi_{(-\infty, -1/n]}(t) - nt \ \chi_{(-1/n, 0)}(t).$$

Introduce $T_n : [0, \infty] \to C([0, 1]; [0, 1])$ by $T_n(s)(t) := \varphi_n(t - s)$ for all $s \in [0, \infty], t \in [0, 1]$. It is straightforward to check that T_n is continuous. Thus the map

$$\rho \mapsto ((T_n \circ \tau)(\rho), \rho) \tag{84}$$

from \mathscr{S}_{Ω} into $C([0,1];[0,1]) \times \mathscr{S}_{\Omega}$ is measurable, given that τ lower semi-continuous, and hence measurable. Moreover

$$(\varphi, \rho) \mapsto \varphi \rho \tag{85}$$

from $C([0,1];[0,1]) \times \mathscr{S}_{\Omega}$ into \mathscr{S}_{Ω} is continuous, since by triangle inequality and (29) we can readily check that for all $\rho^{i} = h_{i}\delta_{\gamma_{i}} \in \mathscr{S}_{\Omega}$, $\varphi_{i} \in C([0,1];[0,1])$, i = 1, 2, it holds $d(\varphi_{1}\rho^{1},\varphi_{2}\rho^{2}) \leq d(\varphi_{1}\rho^{1},\varphi_{2}\rho^{1}) + d(\varphi_{2}\rho^{1},\varphi_{2}\rho^{2}) = \|\varphi_{1} - \varphi_{2}\|_{\infty}\|h_{1}\|_{\infty} + \|\varphi_{2}\|_{\infty}d(\rho^{1},\rho^{2}).$

We now define $G_n: \mathscr{G}_{\Omega} \to \mathscr{G}_{\Omega}$ by composing the maps at (84)–(85), that is,

$$G_n(\rho) = (T_n \circ \tau)(\rho) \ \rho = (\gamma, h \ T_n(\tau(\rho))), \text{ for } \rho \in \mathscr{S}_{\Omega}$$

In view of the above, G_n is measurable for all $n \in \mathbb{N}$. We now claim that $G_n \to G$ pointwise in \mathscr{S}_{Ω} . Indeed, by (29), we see that

$$d(G_n(\rho), G(\rho)) = \sup_{t \in [0, 1]} h(t) |\chi_{[0, \tau(\rho))}(t) - \varphi_n(t - \tau(\rho))|, \text{ for all } \rho = h\delta_{\gamma} \in \mathscr{S}_{\Omega}.$$
(86)

Fix $\rho = h\delta_{\gamma} \in \mathscr{S}_{\Omega}$. If $\tau(\rho) = \infty$ it is immediate to check that $T_n(\tau(\rho)) \equiv 1$ in [0, 1], so that $d(G_n(\rho), G(\rho)) = 0$ for all $n \in \mathbb{N}$ by (86). Similarly, if $\tau(\rho) = 0$, then $T_n(\tau(\rho)) \equiv 0$ in [0, 1] and again $d(G_n(\rho), G(\rho)) = 0$ for all $n \in \mathbb{N}$. Finally, assume that $0 < \tau(\rho) < \infty$ and fix $\varepsilon > 0$. By continuity of h there exists $n_0 \in \mathbb{N}$ such that $\tau(\rho) - 1/n_0 > 0$ and

$$h(t) < \varepsilon$$
 for all $t \in [\tau(\rho) - 1/n_0, \tau(\rho)].$ (87)

For all $n \ge n_0$ we can compute

$$\begin{split} h(t)|\chi_{[0,\tau(\rho))}(t) &- \varphi_n(t-\tau(\rho))| = 0, \quad \text{if} \quad t \in [0,\tau(\rho)-1/n] \cup [\tau(\rho),1], \\ h(t)|\chi_{[0,\tau(\rho))}(t) &- \varphi_n(t-\tau(\rho))| = h(t)|1 + n(t-\tau(\rho))| \le h(t), \quad \text{if} \quad t \in [\tau(\rho)-1/n,\tau(\rho)). \end{split}$$

Recalling (86)–(87) we then obtain $d(G_n(\rho), G(\rho)) < \varepsilon$ for $n \ge n_0$, and the proof of the measurability of G is concluded.

The inclusion $G(\mathscr{S}_{\Omega}) \subset \mathscr{S}_{\Omega}^*$ is immediate from the definition of \mathscr{S}_{Ω}^* . For the inclusion $G(\mathcal{H}_{\Omega}^{\nu,g}) \subset \mathcal{H}_{\Omega}^{\nu,g}$, consider $(\gamma,h) \in \mathcal{H}_{\Omega}^{\nu,g}$ and notice that if $\tau(\gamma,h) = \infty$, then $G(\gamma,h) = (\gamma,h)$ and the thesis is immediate. On the other hand, if $\tau(\gamma,h) \in [0,1]$, then for every $\varphi \in C_c(0,1)$ there holds

$$\begin{split} \int_0^1 h(t)\chi_{[0,\tau(\gamma,h))}(t)\dot{\varphi}(t) \ dt &= \int_0^{\tau(\gamma,h)} h(t)\dot{\varphi}(t) \ dt \\ &= -\int_0^{\tau(\gamma,h)} \dot{h}(t)\varphi(t) \ dt + h(\tau(\gamma,h))\varphi(\tau(\gamma,h)) \\ &= -\int_0^{\tau(\gamma,h)} h(t)g(\gamma(t),t)\varphi(t) \ dt \\ &= -\int_0^1 h(t)\chi_{[0,\tau(\gamma,h))}(t)g(\gamma(t),t)\varphi(t) \ dt, \end{split}$$

implying that $h\chi_{[0,\tau(\gamma,h))} \in AC^2[0,1]$ and $(h\chi_{[0,\tau(\gamma,h))})'(t) = \chi_{[0,\tau(\gamma,h))}(t)g(\gamma(t),t)h(t)$ for a.e. $t \in (0,1)$. Noticing that $\sqrt{h\chi_{[0,\tau(\gamma,h))}} = \sqrt{h\chi_{[0,\tau(\gamma,h))}}$, by a similar argument we obtain that $\sqrt{h\chi_{[0,\tau(\gamma,h))}} \in AC^2[0,1]$ and $\sqrt{h\chi_{[0,\tau(\gamma,h))}}\gamma \in AC^2[0,1]$. This shows that $(\gamma,h\chi_{[0,\tau(\gamma,h))}) \in \mathcal{H}_{\Omega}^{\nu,g}$, concluding the claimed inclusion. We finally show that \mathscr{S}_{Ω}^* is measurable. Consider the map $\mathcal{G}: \mathscr{S}_{\Omega} \to \mathscr{S}_{\Omega} \times \mathscr{S}_{\Omega}$ defined as $\mathcal{G}(\gamma,h) :=$ $((\gamma,h), G(\gamma,h))$, where $\mathscr{S}_{\Omega} \times \mathscr{S}_{\Omega}$ is equipped with the Borel σ -algebra of the product space. As \mathscr{S}_{Ω} is a separable metric space (Proposition 3.6) and G is measurable, we deduce that also \mathcal{G} is measurable. Define the diagonal $D := \{((\gamma,h), (\gamma,h)) : (\gamma,h) \in$ $\mathscr{S}_{\Omega}\}$, which is clearly closed. As the set $\operatorname{Fix}(G) := \{(\gamma,h) \in \mathscr{S}_{\Omega} : G(\gamma,h) = (\gamma,h)\}$ coincides with \mathscr{S}_{Ω}^* , we obtain that $\mathscr{S}_{\Omega}^* = \mathcal{G}^{-1}(D)$, implying that \mathscr{S}_{Ω}^* is measurable. **Lemma 5.4.** Let $(t \mapsto \rho_t) \in C_w([0, 1]; \mathcal{M}^+(\Omega))$ be in $\mathcal{D}_{\nu,g}$ and $\sigma \in \mathcal{M}_1^+(\mathscr{S}_\Omega)$ be concentrated on $\mathcal{H}_{\Omega}^{\nu,g}$. Suppose that (49) and (81) hold. Then σ is concentrated on \mathscr{S}_{Ω}^* .

Proof. Suppose by contradiction that $\sigma(\mathscr{S}_{\Omega} \setminus \mathscr{S}_{\Omega}^*) > 0$. Let τ and G be as in Definition 5.2. As G is measurable (Lemma 5.3), we can consider the measure $\hat{\sigma} := G_{\#}\sigma \in \mathcal{M}^+(\mathscr{S}_{\Omega})$. By the inclusion $G(\mathscr{S}_{\Omega}) \subset \mathscr{S}_{\Omega}^*$ in Lemma 5.3, we have that $\hat{\sigma}$ is concentrated on \mathscr{S}_{Ω}^* . Using that σ is concentrated on $\mathcal{H}_{\Omega}^{v,g}$ and the inclusion $G(\mathcal{H}_{\Omega}^{v,g}) \subset \mathcal{H}_{\Omega}^{v,g}$ (Lemma 5.3), we also deduce that $\hat{\sigma}$ is concentrated on $\mathcal{H}_{\Omega}^{v,g}$. By Remark 4.4 and by definition of $\hat{\sigma}$, we get that $\hat{\sigma}$ satisfies (50) with respect to v and g. Therefore we can apply Theorem 4.3 to $\hat{\sigma}$ and obtain a curve $t \mapsto \hat{\rho}_t$ in $C_w([0, 1]; \mathcal{M}^+(\Omega))$ such that (49) hold and $\partial_t \hat{\rho}_t + \text{div} (v \hat{\rho}_t) = g \hat{\rho}_t$ in X_{Ω} . Additionally, using Remark 4.2, (49), and the definition of $\hat{\sigma}$, we obtain that

$$\int_{0}^{1} \int_{\Omega} |v(t,x)|^{2} + |g(t,x)|^{2} \ d\hat{\rho}_{t}(x) \ dt \leq \int_{0}^{1} \int_{\Omega} |v(t,x)|^{2} + |g(t,x)|^{2} \ d\rho_{t}(x) \ dt < \infty,$$

implying that $(t \mapsto \hat{\rho}_t) \in \mathcal{D}_{\nu,g}$. Moreover $\hat{\rho}_0 = \rho_0$, by (49) at time t = 0 and definition of $\hat{\sigma}$. Finally, using again (49), we estimate

$$\begin{aligned} ||\hat{\rho}||_{\mathcal{M}(X_{\Omega})} &= \int_{0}^{1} \int_{\mathscr{S}_{\Omega}} h(t) \ d\hat{\sigma}(\gamma, h) \ dt = \int_{\mathscr{S}_{\Omega}} \int_{0}^{\tau(h, \gamma)} h(t) \ dt \ d\sigma(\gamma, h) \\ &\leq ||\rho||_{\mathcal{M}(X_{\Omega})} - \int_{\mathscr{S}_{\Omega} \setminus \mathscr{S}_{\Omega}^{*}} \int_{\tau(h, \gamma)}^{1} h(t) \ dt \ d\sigma(\gamma, h). \end{aligned}$$

$$(88)$$

Thanks to the continuity of h for every $(\gamma, h) \in \mathscr{S}_{\Omega}$ and the definition of \mathscr{S}_{Ω}^* , we know that $\int_{\tau(h,\gamma)}^{1} h(t) dt > 0$ for all $(\gamma, h) \in \mathscr{S}_{\Omega} \setminus \mathscr{S}_{\Omega}^*$. Hence, as $\sigma(\mathscr{S}_{\Omega} \setminus \mathscr{S}_{\Omega}^*) > 0$, from (88) we conclude that $||\hat{\rho}||_{\mathcal{M}(X_{\Omega})} < ||\rho||_{\mathcal{M}(X_{\Omega})}$ contradicting (81).

Next we show that we can disintegrate any measure obtained by the application of Theorem 4.3 into a family of Borel measures parametrized by $x \in \Omega$ and concentrated on the set

$$E_{x} := \{ (\gamma, h) \in \mathcal{H}_{\Omega}^{\nu, g} \cap \mathscr{S}_{\Omega}^{*} \cap \mathcal{H}_{\Omega}^{1} : \gamma(0) = x \}.$$
(89)

Notice that E_x is measurable for every $x \in \Omega$. Indeed, by employing similar arguments to the ones in Lemma 3.7, one can show that the map $\pi : \mathscr{S}^*_{\Omega} \setminus \{0\} \to \Omega$ defined as $\pi(\gamma, h) := \gamma(0)$ is continuous. Therefore, as $\mathscr{S}^*_{\Omega} \cap \mathcal{H}^1_{\Omega} \subset \mathscr{S}^*_{\Omega} \setminus \{0\}$, we can write $E_x = \pi^{-1}(x) \cap \mathcal{H}^{v,g}_{\Omega} \cap \mathcal{H}^1_{\Omega}$. Thus E_x is measurable, given that \mathcal{H}^1_{Ω} is closed and $\mathcal{H}^{v,g}_{\Omega}$ is measurable by Lemma 5.3.

Lemma 5.5. Let $v: X_{\Omega} \to \mathbb{R}^d$, $g: X_{\Omega} \to \mathbb{R}$ be measurable. Let $\rho \in \mathcal{D}_{v,g}$ and $\sigma \in \mathcal{M}_1^+(\mathscr{S}_{\Omega})$ be such that (49) holds. Then there exists a Borel family of measures $\{\sigma^x\}_{x\in\Omega} \subset \mathcal{M}^+(\mathscr{S}_{\Omega})$ such that for every $f \in L^1_{\sigma}(\mathscr{S}_{\Omega})$ we have

$$\int_{\mathcal{H}_{\Omega}^{\nu,g} \cap \mathcal{H}_{\Omega}^{1} \cap \mathscr{S}_{\Omega}^{*}} f(\gamma,h) \ d\sigma(h,\gamma) = \int_{\Omega} \int_{E_{x}} f(\gamma,h) \ d\sigma^{x}(\gamma,h) \ d\rho_{0}(x), \tag{90}$$

where E_x is defined as in (89). Moreover σ^x is concentrated on E_x for ρ_0 -a.e. $x \in \Omega$.

Proof. Set $\mathcal{H} := \mathcal{H}_{\Omega}^{\nu,g} \cap \mathcal{H}_{\Omega}^{1} \cap \mathscr{S}_{\Omega}^{*}$ and notice that it is measurable thanks to Lemma 5.3. Consider the map $\pi : \mathscr{S}_{\Omega} \to \Omega$ defined by $\pi(\gamma, h) := \gamma(0) \ \chi_{\mathcal{H}}(\gamma, h) + z \ \chi_{\mathscr{S}_{\Omega} \setminus \mathcal{H}}(\gamma, h)$, where $z \in \Omega$ is arbitrary, but fixed. Notice that, as h(0) > 0 for every $(\gamma, h) \in \mathcal{H}$, the map π is well-defined and measurable in \mathscr{S}_{Ω} using similar arguments as in Lemma 3.7. Define then $\tilde{\sigma} := \sigma \sqcup \mathcal{H}$. We aim to show that $\pi_{\#} \tilde{\sigma} \ll \rho_{0}$. To this end, consider a Borel set $B \subset \Omega$ such that $\rho_{0}(B) = 0$. Define $\tilde{B} = \{(\gamma, h) \in \mathcal{H} : \gamma(0) \in B\}$ and notice that \tilde{B} is measurable as $\tilde{B} = \pi^{-1}(B) \cap \mathcal{H}$. Then

$$0 = \rho_0(B) = \int_{\mathscr{S}_\Omega} h(0)\chi_B(\gamma(0)) \ d\sigma \ge \int_{\mathcal{H}} h(0)\chi_B(\gamma(0)) \ d\sigma = \int_{\tilde{B}} h(0) \ d\sigma,$$

implying that $\sigma(\tilde{B}) = 0$, since h(0) > 0 for all $(\gamma, h) \in \mathcal{H}$. By direct calculation we have $(\pi_{\#}\tilde{\sigma})(B) = \sigma(\tilde{B})$, and thus $(\pi_{\#}\tilde{\sigma})(B) = 0$, concluding that $\pi_{\#}\tilde{\sigma} \ll \rho_0$. Hence, as \mathscr{S}_{Ω} is a complete separable metric space by Proposition 3.6, we can apply Theorem A.3 to $\tilde{\sigma} \in \mathcal{M}^+(\mathscr{S}_{\Omega})$, and obtain a Borel family of measures $\{\sigma^x\}_{x\in\Omega} \subset \mathcal{M}^+(\mathscr{S}_{\Omega})$ satisfying the thesis.

5.2. Proof of Theorem 5.1

Assume that $t \mapsto \rho_t$ belongs to $\mathcal{D}_{v,g}$. Moreover suppose that ρ_0 is concentrated on $A \subset \Omega$ and that (81) holds. By Theorem 4.3, there exists $\sigma \in \mathcal{M}_1^+(\mathscr{S}_{\Omega})$ concentrated on $\mathcal{H}^{\nu,g}_{\Omega} \cap \mathcal{H}^{1}_{\Omega}$ that represents ρ_{t} , that is, (49) holds. Using Lemma 5.4, we infer that σ is concentrated on $\mathcal{H} := \mathcal{H}_{\Omega}^{\nu,g} \cap \mathcal{H}_{\Omega}^{1} \cap \mathscr{S}_{\Omega}^{*}$. Thanks to Lemma 5.5 we can disintegrate σ into a Borel family $\{\sigma^x\}_{x\in\Omega} \subset \mathcal{M}^+(\mathscr{G}_\Omega)$ such that (90) holds, with σ^x concentrated on E_x for ρ_0 -a.e. $x \in \Omega$. We claim that assumption (*Hyp*) implies that E_x contains at most one point for all $x \in A$. Indeed, suppose that $(\gamma_1^x, h_1^x), (\gamma_2^x, h_2^x) \in E_x$. As $(\gamma_i^x, h_i^x) \in \mathscr{S}_{\Omega}^* \cap \mathcal{H}_{\Omega}^1$, there exist $\tau_i \in \mathbb{R}$ such that $\{h_i > 0\} = [0, 1] \cap (-\infty, \tau_i)$ and $||h_i^x||_1 = 1$. Assume $\tau_1 \leq \tau_2$. As $(\gamma_i^x, h_i^x) \in \mathcal{H}_{\Omega}^{\gamma,g}$, we have that (γ_i^x, h_i^x) solves (O1)-(O2) in $[0, \tau_1)$. Now notice that by linearity of (O2) and assumption (*Hyp*), we have that $\gamma_1^x(t) = \gamma_2^x(t)$ and $h_1^x(t) = h_2^x(t)h_1^x(0)/h_2^x(0)$ for all $t \in [0, \tau_1)$. As $||h_i^x||_1 = 1$, we then infer $(\gamma_1^x, h_1^x) = (\gamma_2^x, h_2^x)$ in $[0, \tau_1)$ and by the continuity of h_i we also obtain that $h_1^x(\tau_1) = h_2^x(\tau_2) = 0$. By definition of τ_i we conclude that $\tau_1 = \tau_2 = \tau$, so that $(\gamma_1^x, h_1^x) = (\gamma_2^x, h_2^x)$ in $[0, \tau)$. Since $h_1^x(t) = h_2^x(t) = 0$ for all $t \ge \tau$, we conclude that E_x contains at most one point. Thus, for ρ_0 -a.e. $x \in E$, $E := \{x \in \Omega : E_x \neq \emptyset\}$, we have $\sigma^x = c_x \ \delta_{(\gamma^x, h^x)}$, with $c_x := ||\sigma^x||_{\mathcal{M}(\mathscr{G}_\Omega)}, \ (\gamma^x, h^x) \in E_x$. We claim that $c_x = 1/h^x(0)$. Indeed, by definition of E_x , we have $\gamma^x(0) = x$. Using (49, 90), and $\sigma(\mathscr{S}_\Omega \setminus \mathcal{H}) = 0$, we then obtain

$$\int_{\Omega} \varphi(x) \ d\rho_0(x) = \int_E c_x h^x(0) \varphi(x) \ d\rho_0(x),$$

for all $\varphi \in C(\Omega)$, showing that $c_x h^x(0) = 1$ for ρ_0 -a.e. $x \in E$. Again by (49) and (90) we get

$$\int_{\Omega} \varphi(x) \ d\rho_t(x) = \int_E \frac{1}{h^x(0)} h^x(t) \varphi(\gamma^x(t)) \ d\rho_0(x),$$

for every $\varphi \in C(\Omega)$, where we also used that $\sigma^x = \frac{1}{h^x(0)} \delta_{(\gamma^x, h^x)}$ for ρ_0 -a.e. $x \in E$. Thus ρ_t depends only on the initial data ρ_0 , ending the proof.

6. Extremal points of the Wasserstein-Fisher-Rao energy

Let $\Omega \subset \mathbb{R}^d$ with $d \ge 1$ be the closure of a bounded domain of \mathbb{R}^d . Let $\alpha, \beta > 0, \delta \in (0, \infty]$ and define \mathscr{B} to be the unit ball of the functional $J_{\alpha, \beta, \delta}$ defined at (17), that is,

$$\mathscr{B} := \{ (\rho, m, \mu) \in \mathcal{M}_{\Omega} : J_{\alpha, \beta, \delta}(\rho, m, \mu) \leq 1 \}.$$

The aim of this section is to characterize the extremal points Ext \mathscr{B} . Notice that $J_{\alpha,\beta,\infty}$ corresponds to the coercive version of the Benamou-Brenier energy, whose extremal points were characterized in [13]. Hence here we focus on the case $\delta < \infty$. After the characterization of Ext \mathscr{B} is obtained, we will show how this information can be applied to the analysis of dynamic inverse problems which are regularized via the optimal transport energy $J_{\alpha,\beta,\delta}$ [34]. In particular we will obtain a sparse representation formula for regularized solutions to the dynamic problem.

Before stating the characterization theorem we remind the reader the notations $\mathscr{C}_{\Omega}, \mathscr{S}_{\Omega}, \mathcal{H}_{\Omega}$ introduced at (25, 26, 32). In the following \mathscr{S}_{Ω} is equipped with the distance d at (28), making it a complete metric space (Proposition 3.6). We now define the set of characteristics of (14) with energy $J_{\alpha,\beta,\delta} = 1$, which will play a role in the characterization of Ext \mathscr{B} .

Definition 6.1 (Characteristics). Define the set C of all the triples $(\rho, m, \mu) \in \mathcal{M}_{\Omega}$ of the form $\rho = h(t) dt \otimes \delta_{\gamma(t)}$, $m = \dot{\gamma}(t)\rho$, $\mu = \dot{h}(t) dt \otimes \delta_{\gamma(t)}$ that satisfy the following properties:

- (i) $t \mapsto h(t)\delta_{\gamma(t)}$ belongs to \mathcal{H}_{Ω} ,
- (ii) the set $\{h > 0\} := \{t \in [0, 1] : h(t) > 0\}$ is connected,
- (iii) the energy satisfies $J_{\alpha,\beta,\delta}(\rho, m, \mu) = 1$.

The above definition is well-posed since (ρ, m, μ) belongs to \mathcal{M}_{Ω} and solves the continuity equation (14) in X_{Ω} (by the converse of Proposition 3.9 with $V = \Omega$). Hence (iii) is compatible with the definition of $J_{\alpha,\beta,\delta}$.

Remark 6.2. If $(\rho, m, \mu) \in \mathcal{M}_{\Omega}$ with $\rho \in \mathcal{H}_{\Omega}$, then an application of Proposition 3.9 (with $V = \Omega$) yields the representation

$$J_{\alpha,\beta,\delta}(\rho,m,\mu) = J_{\alpha,\beta,\delta}(\gamma,h) = \int_{\{h>0\}} \frac{\beta}{2} h(t) |\dot{\gamma}(t)|^2 + \frac{\beta\delta^2}{2} \frac{h(t)^2}{h(t)} + \alpha h(t) dt.$$
(91)

In particular $J_{\alpha,\beta,\delta}$ is d-measurable, as a consequence of Proposition 3.10. For a measurable set $E \subset [0,1]$ we define the localized energy

$$J_{\alpha,\beta,\delta,E}(\rho,m,\mu):=\int_{E\cap\{h>0\}}\frac{\beta}{2} h(t)|\dot{\gamma}(t)|^2+\frac{\beta\delta^2}{2} \frac{\dot{h}(t)^2}{h(t)}+\alpha h(t) dt.$$

We are now ready to state the characterization theorem.

Theorem 6.3. For parameters $\alpha, \beta, \delta > 0$ we have $\text{Ext } \mathcal{B} = \mathcal{C} \cup \{0\}$, where 0 denotes the null triple in \mathcal{M}_{Ω} .

The proof of Theorem 6.3 will be carried out in the next section, while in Section 6.2 we will detail the application of Theorem 6.3 to dynamic inverse problems.

6.1. Proof of Theorem 6.3

In order to simplify notations, we will denote $J := J_{\alpha,\beta,\delta}$ and $J_E := J_{\alpha,\beta,\delta,E}$ for any $E \subset [0,1]$ measurable.

Step 1. $C \cup \{0\} \subset \text{Ext} \mathscr{B}$: Assume first that $(\rho, m, \mu) = 0$, and that there exists a decomposition

$$(0,0,0) = \lambda(\rho^1, m^1, \mu^1) + (1-\lambda)(\rho^2, m^2, \mu^2)$$
(92)

with $(\rho^j, m^j, \mu^j) \in \mathscr{B}$ and $\lambda \in (0, 1)$. In particular by Lemma A.4 point (i) we have $\rho^j \ge 0$ and $m^j, \mu^j \ll \rho^j$. Therefore (92) immediately implies that $(\rho^j, m^j, \mu^j) = 0$, showing that $0 \in \text{Ext } \mathscr{B}$.

Assume now that $(\rho, m, \mu) \in C$, according to Definition 6.1. In particular the set $\{h > 0\}$ is non-empty, since $J(\rho, m, \mu) = 1$. Assume that $(\rho^1, m^1, \mu^1), (\rho^2, m^2, \mu^2) \in \mathscr{B}$ are such that

$$(\rho, m, \mu) = \lambda(\rho^1, m^1, \mu^1) + (1 - \lambda)(\rho^2, m^2, \mu^2),$$
(93)

for some $\lambda \in (0, 1)$. We need to show that $(\rho, m, \mu) = (\rho^j, m^j, \mu^j)$. By (93), convexity of J (see Lemma A.5), and the fact that $J(\rho^j, m^j, \mu^j) \leq 1$, $J(\rho, m, \mu) = 1$, we have that $J(\rho^j, m^j, \mu^j) = 1$. Thus, by Lemmas 2.1, A.4 we infer $\rho^j = dt \otimes \rho_t^j$ with $t \mapsto \rho_t^j$ in $C_w([0, 1]; \mathcal{M}^+(\Omega))$. Set $h_j(t) := \rho_t^j(\Omega)$ and notice that h_j is continuous by narrow continuity of ρ^j . From the decomposition (93) and the uniqueness of the disintegration, we thus obtain $\rho_t^j = h_j(t)\delta_{\gamma(t)} \in \mathscr{S}_{\Omega}$, and in particular

$$h(t) = \lambda h_1(t) + (1 - \lambda)h_2(t)$$
 for every $t \in [0, 1]$. (94)

We will now show that there exists c > 0 such that

$$h_2(t) = c \ h_1(t) \quad \text{for all} \quad t \in \{h > 0\}.$$
 (95)

We start by defining the sets $E := \{h_1 > 0\} \cap \{h_2 > 0\}, Z_1 := \{h_1 > 0\} \cap \{h_2 = 0\}$ and $Z_2 := \{h_1 = 0\} \cap \{h_2 > 0\}$. These sets are pairwise disjoint, and by (94) we have $\{h > 0\} = E \cup Z_1 \cup Z_2$, where we recall that $\{h > 0\} \neq \emptyset$ is connected by assumption. We claim that $E \neq \emptyset$. Indeed, assume by contradiction that $E = \emptyset$, so that in particular $Z_1 \cup Z_2 = \{h > 0\}$. Notice that Z_1, Z_2 are relatively closed in $\{h > 0\}$ since they can be written as $Z_1 = \{h > 0\} \cap \{h_2 = 0\}, Z_2 = \{h > 0\} \cap \{h_1 = 0\}$, due to (94). As $\{h > 0\}$ is connected, we deduce that either $Z_1 = \emptyset$ or $Z_2 = \emptyset$. If $Z_1 = \emptyset$, then we would have $h_1 \equiv 0$, which in turn would imply $\rho^1 = 0$. Hence by Lemma A.4 point (i) we would obtain $J(\rho^1, m^1, \mu^1) = 0$, contradicting $J(\rho^1, m^1, \mu^1) = 1$. Similarly $Z_2 = \emptyset$ leads to the contradiction $J(\rho^2, m^2, \mu^2) = 0$. We therefore conclude $E \neq \emptyset$.

Claim: h_1/h_2 is constant in each connected component of E.

Proof of Claim: Since $J(\rho^j, m^j, \mu^j) < \infty$, by Proposition 3.9, we have that $\rho^j \in \mathcal{H}_{\Omega}$ and there exist $\nu^j : X_{\Omega} \to \mathbb{R}^d$, $g^j : X_{\Omega} \to \mathbb{R}$ measurable such that $m^j = \nu^j \rho^j$, $\mu^j = g^j \rho^j$ and

$$h_j(t) = g^j(t, \gamma(t))h_j(t)$$
 for a.e. $t \in (0, 1),$ (96)

$$\dot{\gamma}(t) = v^{j}(t, \gamma(t)) \text{ a.e. in } \{h_{j} > 0\}.$$
 (97)

Moreover $J(\rho^{j}, m^{j}, \mu^{j}) = J(h_{j}, \gamma)$ can be computed via (91). By direct calculation, and using (93) and (91) we have

$$\begin{split} J_E(\rho, m, \mu) &= J_E(h, \gamma) = J_E(\lambda h_1 + (1 - \lambda)h_2, \gamma) \\ &= \int_E (\lambda h_1 + (1 - \lambda)h_2) \left(\frac{\beta}{2} |\dot{\gamma}|^2 + \alpha\right) dt + \frac{\beta \delta^2}{2} \int_E \frac{(\lambda \dot{h}_1 + (1 - \lambda)\dot{h}_2)^2}{\lambda h_1 + (1 - \lambda)h_2} dt \\ &= \lambda J_E(h_1, \gamma) + (1 - \lambda)J_E(h_2, \gamma) + \frac{\beta \delta^2}{2} \int_E \frac{(\lambda \dot{h}_1 + (1 - \lambda)\dot{h}_2)^2}{\lambda h_1 + (1 - \lambda)h_2} - \lambda \frac{\dot{h}_1^2}{h_1} - (1 - \lambda)\frac{\dot{h}_2^2}{h_2} dt, \end{split}$$

so that

$$J_{E}(\rho, m, \mu) = \lambda J_{E}(\rho^{1}, m^{1}, \mu^{1}) + (1 - \lambda) J_{E}(\rho^{2}, m^{2}, \mu^{2}) \\ - \frac{\beta \delta^{2}}{2} \lambda (1 - \lambda) \int_{E} \frac{(\dot{h}_{1}h_{2} - h_{1}\dot{h}_{2})^{2}}{(\lambda h_{1} + (1 - \lambda)h_{2})h_{1}h_{2}} dt.$$
⁽⁹⁸⁾

By proceeding as above, one can check that

$$J_{Z_1}(\rho, m, \mu) = \lambda J_{Z_1}(\rho^1, m^1, \mu^1) , \quad J_{Z_2}(\rho, m, \mu) = (1 - \lambda) J_{Z_2}(\rho^2, m^2, \mu^2),$$
(99)

where we used (93, 94), definition of Z_i and [43, Theorem 4.4]. Moreover by definition

$$J(\rho^{1}, m^{1}, \mu^{1}) = J_{Z_{1}}(\rho^{1}, m^{1}, \mu^{1}) + J_{E}(\rho^{1}, m^{1}, \mu^{1}),$$
(100)

$$J(\rho^2, m^2, \mu^2) = J_{Z_2}(\rho^2, m^2, \mu^2) + J_E(\rho^2, m^2, \mu^2).$$
(101)

By combining (98)–(101), we obtain

$$J(\rho, m, \mu) = \lambda J(\rho^{1}, m^{1}, \mu^{1}) + (1 - \lambda) J(\rho^{2}, m^{2}, \mu^{2}) \\ - \frac{\beta \delta^{2}}{2} \lambda (1 - \lambda) \int_{E} \frac{(\dot{h}_{1}h_{2} - h_{1}\dot{h}_{2})^{2}}{(\lambda h_{1} + (1 - \lambda)h_{2})h_{1}h_{2}} dt.$$

Now we can make use of the fact that $J(\rho, m, \mu) = J(\rho^j, m^j, \mu^j) = 1$ to infer $\dot{h}_1 h_2 = h_1 \dot{h}_2$ a.e. in *E*. In particular $(h_1/h_2)' = 0$ a.e. in *E*, and hence the claim follows.

We are now ready to show (95). For an arbitrary C > 0 and $t \in \{h > 0\}$ define the map

$$f(t) := \min igg(rac{h_1(t)}{h_2(t)}, C igg) \chi_{\{h_2 > 0\}} + C \,\, \chi_{\{h_2 = 0\}}.$$

Notice that f is continuous and, since $E \neq \emptyset$, f is not identically zero. Moreover, as $(h_1/h_2)' = 0$ a.e. in E, the image $f(\{h > 0\})$ is at most countable. Assume by contradiction that $Z_2 \neq \emptyset$, and notice that f vanishes on Z_2 . Therefore f assumes at least two

different values on $\{h > 0\}$, which is a contradiction, as $f(\{h > 0\})$ is connected and, consequently, uncountable. Hence $Z_2 = \emptyset$ and $\{h > 0\} = E \cup Z_1$. By interchanging the roles of h_1 and h_2 , we can repeat the same argument and conclude that $Z_1 = \emptyset$, so that $E = \{h > 0\}$. As $\{h > 0\}$ is connected, we thus deduce (95) directly from the fact that $(h_1/h_2)' = 0$ a.e. in *E*.

We are now ready to conclude. Indeed, note that, as $\rho_t^j = h_j(t)\delta_{\gamma(t)}$, condition (95) implies that $\rho^2 = c\rho^1$ and $\{h_j > 0\} = \{h > 0\}$. In particular (97) yields $\nu^j(t, \gamma(t)) = \dot{\gamma}(t)$ a.e. in $\{h > 0\}$, showing that $m^2 = c m^1$. Finally from (96) we infer $g^1(t, \gamma(t)) = g^2(t, \gamma(t))$ a.e. in (0, 1), from which we conclude $\mu^2 = c \mu^1$. In total we have $(\rho^2, m^2, \mu^2) = c (\rho^1, m^1, \mu^1)$, and by $J(\rho^j, m^j, \mu^j) = 1$ and one-homogeneity of J we conclude that c = 1. Therefore (93) yields extremality of (ρ, m, μ) .

Step 2. Ext $\mathscr{B} \subset \mathcal{C} \cup \{0\}$: Let $(\rho, m, \mu) \in \text{Ext } \mathscr{B}$. We can assume that $(\rho, m, \mu) \neq 0$, so that $J(\rho, m, \mu) > 0$. By extremality of (ρ, m, μ) , convexity and 1-homogeneity of J, we conclude that $J(\rho, m, \mu) = 1$. In particular by Lemma A.4 we obtain $\rho \geq 0$ and $m = v\rho$, $\mu = g\rho$ for some measurable maps $v : X_{\Omega} \to \mathbb{R}^d$, $g : X_{\Omega} \to \mathbb{R}$ satisfying

$$J(\rho, m, \mu) = \int_0^1 \int_\Omega \left(\frac{\beta}{2} |v(t, x)|^2 + \frac{\beta \delta^2}{2} |g(t, x)|^2 + \alpha \right) \, d\rho_t(x) \, dt = 1.$$
(102)

By definition of *J*, we then have that $\partial_t \rho_t + \text{div} (v \rho_t) = g \rho_t$ in X_{Ω} . Thanks to Lemma 2.1 we also have $\rho = dt \otimes \rho_t$ with $t \mapsto \rho_t$ in $C_w([0, 1]; \mathcal{M}^+(\Omega))$. Set $h := \rho_t(\Omega)$, and recall that *h* is continuous. We first prove the following claim.

Claim: supp ρ_t is a singleton for every $t \in \{h > 0\}$.

Proof of Claim: Assume by contradiction that there exists $\hat{t} \in \{h > 0\}$ such that supp $\rho_{\hat{t}}$ is not a singleton. Then there exist disjoint Borel sets $E_1, E_2 \subset \Omega$ such that $E_1 \cup E_2 = \Omega$ and $\rho_{\hat{t}}(E_i) > 0$ for i = 1, 2. Invoking Theorem 4.3, there exists a measure $\sigma \in \mathcal{M}_1^+(\mathscr{S}_{\Omega})$ concentrated on $\mathcal{H}_{\Omega}^{v,g}$ which represents ρ_t , that is, (49) holds. Define the sets

$$A_i := \{(\gamma, h) \in \mathscr{S}_{\Omega} : \gamma(\hat{t}) \in E_i, \ h(\hat{t}) > 0\}, \qquad Z := \{(\gamma, h) \in \mathscr{S}_{\Omega} : h(\hat{t}) = 0\},$$

and notice that A_1, A_2, Z are pairwise disjoint and $\mathscr{S}_{\Omega} = A_1 \cup A_2 \cup Z$. Also Z is d-measurable, being d-closed, as it is readily seen by (29). We claim that also A_i is d-measurable. To this end define the maps $e_t : \mathscr{S}_{\Omega} \to \mathscr{C}_{\Omega}$ with $e_t(\rho) := \rho_t$ and $\pi : \mathscr{C}_{\Omega} \to \mathbb{R}^d$ where $\pi(\gamma, h) := \gamma \chi_{\mathscr{C}_{\Omega} \setminus \{0\}}(\gamma, h) + p\chi_{\{0\}}(\gamma, h)$, with $p \in \mathbb{R}^d \setminus \Omega$ arbitrary but fixed. Notice that by construction e_t is continuous from $(\mathscr{S}_{\Omega}, d)$ into $(\mathscr{C}_{\Omega}, d_F)$. Moreover π is measurable since the map $(\gamma, h) \mapsto \gamma$ is d_F -continuous in $\mathscr{C}_{\Omega} \setminus \{0\}$. Since $A_i = (\pi \circ e_i)^{-1}(E_i)$, we have that A_i is measurable. By applying (49), we get

$$0 < \rho_{\hat{t}}(E_i) = \int_{\mathscr{S}_{\Omega}} h(\hat{t}) \chi_{E_i}(\gamma(\hat{t})) \ d\sigma(\gamma, h) = \int_{A_i} h(\hat{t}) \ d\sigma(\gamma, h), \tag{103}$$

which implies $\sigma(A_i) > 0$. Hence setting $\Sigma_1 := A_1$, $\Sigma_2 := A_2 \cup Z$ we obtain a measurable partition of \mathscr{S}_{Ω} with $\sigma(\Sigma_i) > 0$. Notice now that the map $\Psi(t,x) := \beta |\nu(t,x)|^2/2 + \beta \delta^2 |g(t,x)|^2/2 + \alpha$ belongs to $L^1_{\rho_t}(\Omega)$ for a.e. $t \in (0,1)$, thanks to (102). Moreover *J* is non-negative and d-measurable by Remark 6.2. Since σ is concentrated on $\mathcal{H}^{\nu,g}_{\Omega}$, we can apply Remark 4.2 to Ψ and obtain COMMUNICATIONS IN PARTIAL DIFFERENTIAL EQUATIONS 🛞 2061

$$\int_{\mathscr{S}_{\Omega}} J(\gamma, h) \ d\sigma(\gamma, h) = \int_{\mathscr{S}_{\Omega}} \int_{0}^{1} h(t) \Psi(t, \gamma(t)) \ dt \ d\sigma(\gamma, h) = \int_{0}^{1} \int_{\Omega} \Psi(t, x) \ d\rho_{t}(x) \ dt = 1,$$
(104)

where in the last equality we again used (102). Define the coefficients $\lambda_i := \int_{\Sigma_i} J(\gamma, h) \, d\sigma(\gamma, h)$. From (104) we infer $0 \le \lambda_1, \lambda_2 \le 1$ and $\lambda_1 + \lambda_2 = 1$. We claim that $\lambda_i > 0$. Indeed, the map $f_i(t) := \int_{\Sigma_i} h(t) \, d\sigma(\gamma, h)$ for $t \in [0, 1]$ is continuous by dominated convergence and the fact that $\int_{\mathscr{S}_\Omega} ||h||_{\infty} \, d\sigma(\gamma, h) < \infty$, as $\sigma \in \mathcal{M}_1^+(\mathscr{S}_\Omega)$. Notice that by construction $f_i(\hat{t}) > 0$. Therefore by definition of J and continuity of f_i we have

$$\lambda_i = \int_{\Sigma_i} J(\gamma, h) \ d\sigma(\gamma, h) \ge \int_{\Sigma_i} \int_0^1 h(t) \ dt \ d\sigma(\gamma, h) = \int_0^1 f_i(t) \ dt > 0,$$

as claimed. The measure $\sigma|_{\Sigma_i}$ satisfies the hypothesis of the converse in Theorem 4.3, given that (104) holds and σ is concentrated on $\mathcal{H}_{\Omega}^{\nu,g}$. Hence, the curve $t \mapsto \rho_t^i$ defined by

$$\int_{\Omega} \varphi(x) \ d\rho_t^i(x) := \int_{\Sigma_i} h(t) \varphi(\gamma(t)) \ d\sigma(\gamma, h), \quad \text{for all} \quad \varphi \in C(\Omega)$$
(105)

belongs to $C_w([0, 1]; \mathcal{M}^+(\Omega))$ and solves the continuity equation with ν and g. We can now define $(\rho^i, m^i, \mu^i) \in \mathcal{M}$ by setting $\rho^i := dt \otimes \rho^i_t, m^i := \nu \rho^i, \mu^i := g\rho^i$. Note that by (49) and (105) we have that $\rho^i_t \leq \rho_t$ for every $t \in [0, 1]$. Hence

$$\int_0^1 \int_\Omega \left(\frac{\beta}{2} |v(t,x)|^2 + \frac{\beta \delta^2}{2} |g(t,x)|^2 + \alpha \right) \ d\rho_t^i(x) \ dt \leq 1,$$

by (102). Given that the above holds, by repeating the same arguments used to prove (104), but applied to ρ_t^i and $\sigma_{\perp \Sigma_i}$, we have that $J(\rho^i, m^i, \mu^i) = \lambda_i$. Consider the decomposition

$$(\rho, m, \mu) = \lambda_1 \frac{1}{\lambda_1} (\rho^1, m^1, \mu^1) + \lambda_2 \frac{1}{\lambda_2} (\rho^2, m^2, \mu^2),$$
(106)

and notice that $\lambda_i^{-1}(\rho^i, m^i, \mu^i) \in \mathscr{B}$ thanks to the condition $J(\rho^i, m^i, \mu^i) = \lambda_i$ and to the one-homogeneity of *J*. We assert that

$$\frac{1}{\lambda_1}(\rho^1, m^1, \mu^1) \neq \lambda_2 \ \frac{1}{\lambda_2}(\rho^2, m^2, \mu^2).$$
(107)

Indeed we have that $\lambda_1^{-1}\rho^1 \neq \lambda_2^{-1}\rho^2$: If they were equal then by narrow continuity we would have $\lambda_1^{-1}\rho t^1 = \lambda_2^{-1}\rho t^2$. However by (103) it is immediate to check that $\rho t^1(E_1) = \rho_t(E_1) > 0$ and $\rho t^2(E_1) = 0$, yielding a contradiction. Thus (107) holds and (106) gives a non-trivial convex decomposition of (ρ, m, μ) , contradicting extremality. This proves the claim.

In particular, we have shown that $\rho_t = h(t)\delta_{\gamma(t)}$ for some $\gamma : [0,1] \to \Omega$, $h \ge 0$. Thus $t \mapsto \rho_t$ belongs to \mathscr{S}_{Ω} , being narrowly continuous. Hence $\gamma \in C(\{h > 0\}; \mathbb{R}^d)$ thanks to Lemma 3.4. Moreover, as a consequence of (102) and Proposition 3.9, we have that $t \mapsto \rho_t$ belongs to \mathcal{H}_{Ω} , $m = \dot{\gamma}\rho$, $\mu = \dot{h}(t) dt \otimes \delta_{\gamma(t)}$ and

$$J(\rho, m, \mu) = \int_{\{h>0\}} \frac{\beta}{2} h(t) |\dot{\gamma}(t)|^2 + \frac{\beta \delta^2}{2} \frac{\dot{h}(t)^2}{h(t)} + \alpha h(t) \, dt = 1.$$
(108)

In order to prove that $(\rho, m, \mu) \in C$ we are left to show that the set $\{h > 0\}$ is connected. To this end, assume by contradiction that $\{h > 0\} = E_1 \cup E_2$ with E_1 , E_2 relatively open, non-empty and disjoint. For $t \in [0, 1]$ set $\rho_t^i := h(t)\chi_{E_i}(t)\delta_{\gamma(t)}$. Note that as $\{h > 0\}$ is relatively open we have that $\partial^{\{h>0\}}E_i = \partial^{[0,1]}E_i \cap \{h > 0\}$ where we denote by ∂^A the relative boundary with respect to the set A. Hence as $\partial^{\{h>0\}}E_i = \emptyset$ we deduce that h(t) = 0 for every $t \in \partial^{[0,1]}E_i$. In particular the map $t \mapsto h(t)\chi_{E_i}(t)$ is continuous in [0,1]. Moreover $\gamma \in C(\{h\chi_{E_i} > 0\}; \mathbb{R}^d)$, hence Lemma 3.4 ensures that the curve $t \mapsto \rho_t^i$ belongs to \mathcal{S}_{Ω} . We claim that $t \mapsto \rho_t^i$ belongs to \mathcal{H}_{Ω} . In order to show this, we make use of the information $(t \mapsto \rho_t) \in \mathcal{H}_{\Omega}$. Notice that the set E_i is relatively open in [0,1], given that $\{h > 0\}$ is open. Thus $E_i = \bigcup_{n=1}^{\infty} I_n$, where $\{I_n\}_n$ are pairwise disjoint intervals in [0,1]. By dominated convergence

$$\int_{0}^{1} h(t)\chi_{E_{i}}(t)\dot{\varphi}(t) dt = \sum_{n=1}^{\infty} \int_{I_{n}} h(t)\dot{\varphi}(t) dt = -\sum_{n=1}^{\infty} \int_{I_{n}} \dot{h}(t)\varphi(t) dt = \int_{0}^{1} \dot{h}(t)\chi_{E_{i}}(t)\varphi(t) dt$$

for every $\varphi \in C_c^1(0,1)$, where we used that h=0 on $\partial^{[0,1]}I_n$, given that $\partial^{[0,1]}I_n \subset \partial^{[0,1]}E_i$. Since $h \in AC^2[0,1]$, we infer that $h\chi_{E_i} \in AC^2[0,1]$, with derivative $\dot{h}\chi_{E_i}$. Noticing that $\sqrt{h\chi_{E_i}} = \sqrt{h}\chi_{E_i}$ by similar arguments we also deduce that $\sqrt{h} \in AC^2[0,1]$ and $\sqrt{h}\chi_{E_i}\gamma \in AC^2([0,1];\mathbb{R}^d)$, thus concluding $(t \mapsto \rho_t^i) \in \mathcal{H}_\Omega$. Set

$$ho^i:=\chi_{E_i}(t)h(t) \,\,\,dt\otimes\delta_{\gamma(t)}, \quad m^i:=\dot{\gamma}(t)
ho^i, \quad \mu^i=\chi_{E_i}(t)\dot{h}(t)dt\otimes\delta_{\gamma(t)}.$$

Thanks to Proposition 3.9 we have that (ρ^i, m^i, μ^i) belongs to \mathcal{M}_{Ω} and

$$J(\rho^i,m^i,\mu^i)=\int_{E_i}rac{eta}{2}h(t)|\dot{\gamma}(t)|^2+rac{eta\delta^2}{2}\;rac{\dot{h}(t)^2}{h(t)}+lpha h(t)\;dt<\infty.$$

Set $\lambda_i := J(\rho^i, m^i, \mu^i)$ and notice that $0 < \lambda_i < 1$, $\lambda_1 + \lambda_2 = 1$ thanks to (108) and definition of E_i . By construction we have $\chi_{E_1} + \chi_{E_2} = 1$ in $\{h > 0\}$. By recalling that $\dot{h} = 0$ a.e. in $\{h > 0\}$, we have that a decomposition of the form (106) holds. As $\lambda_1^{-1}(\rho^1, m^1, \mu^1) \neq \lambda_2^{-1}(\rho^2, m^2, \mu^2)$ and $\lambda_i^{-1}(\rho^i, m^i, \mu^i) \in \mathcal{B}$, this contradicts the extremality of (ρ, m, μ) . Thus we conclude that the set $\{h > 0\}$ must be connected, ending the proof.

6.2. Sparsity for dynamic inverse problems with optimal transport regularization

In this section we analyze the problem of reconstructing a family of time-dependent Radon measures given a finite number of observations. More precisely, let H be a finite dimensional Hilbert space and $K : C_w([0, 1]; \mathcal{M}(\Omega)) \to H$ be a linear operator which is continuous in the following sense: given a sequence $\{(t \mapsto \rho_t^n)\}_n$ in $C_w([0, 1]; \mathcal{M}(\Omega))$, we require that

$$\rho_t^n \to \rho_t \quad \text{narrowly in} \quad \mathcal{M}(\Omega) \quad \text{for all} \quad t \in [0, 1] \quad \text{implies} \quad K\rho^n \to K\rho \quad \text{in } H.$$
(109)

For a given datum $y \in H$, we aim at finding a solution $\rho \in C_w([0, 1]; \mathcal{M}(\Omega))$ to the ill-posed inverse problem

$$K\rho = y. \tag{110}$$

We regularize (110) via the Hellinger–Kantorovich-type energy $J_{\alpha,\beta,\delta}$ defined at (17), following the approach in [34]. To this end, introduce the space

$$ilde{\mathcal{M}}_{\Omega} := C_{\mathrm{w}}([0,1];\mathcal{M}(\Omega)) imes \mathcal{M}(X_{\Omega};\mathbb{R}^d) imes \mathcal{M}(X_{\Omega};\mathbb{R}),$$

and define the Tikhonov functional $G: \tilde{\mathcal{M}}_{\Omega} \to \mathbb{R} \cup \{\infty\}$ by

$$G(\rho, m, \mu) := F(K\rho) + J_{\alpha, \beta, \delta}(\rho, m, \mu), \qquad (111)$$

where $F: H \to \mathbb{R}$ is a fidelity functional assumed to be convex, lower semi-continuous and bounded from below. We then replace (110) by

$$\min_{(\rho,m,\mu)\in\tilde{\mathcal{M}}_{\Omega}} G(\rho,m,\mu).$$
(112)

Note that G is proper, since $J_{\alpha,\beta,\delta}(0,0,0) = 0$. Moreover under the assumptions on K and F, problem (112) admits a solution: This is indeed an immediate consequence of the direct method and of Lemma A.5.

It is well-known that the finite-dimensionality of the data space *H* promotes sparsity in the reconstruction of solutions to (110), in the sense that there exists a minimizer to (112) which is finite linear combination of extremal points of the ball of the regularizer. This observation was recently made rigorous in the works [40,41] (see also [45,46]). Since in Theorem 6.3 we characterized the extremal points of the ball of $J_{\alpha,\beta,\delta}$, we can specialize the representation results in [40,41] to our setting, and obtain the following statement for sparse minimizers to (112).

Theorem 6.4. There exists a solution $(\hat{\rho}, \hat{m}, \hat{\mu}) \in \tilde{\mathcal{M}}_{\Omega}$ to (112) which is of the form

$$(\hat{\rho}, \hat{m}, \hat{\mu}) = \sum_{i=1}^{p} c_i \ (\rho^i, m^i, \mu^i), \tag{113}$$

where $p \leq \dim(H)$, $c_i > 0$, $\sum_{i=1}^{p} c_i = J_{\alpha,\beta,\delta}(\hat{\rho}, \hat{m}, \hat{\mu})$ and $(\rho^i, m^i, \mu^i) \in C$, with C is as in Definition 6.1.

In order to prove the above theorem, it is sufficient to apply Theorem 6.3 and check validity for the assumptions of Corollary 2 in [40]. The proof is a straightforward adaptation of the one of Theorem 10 in [13] (which deals with the case $\delta = \infty$) and is hence omitted.

We now present an application of Theorem 6.4 to dynamic inverse problems, in a simplified case of the framework introduced in [34]. To be more specific, let $t_1 < ... < t_N$ be a finite discretization of the time interval [0, 1]. The aim is to reconstruct an element of $C_w([0, 1]; \mathcal{M}(\Omega))$ by only making observations at the time instants t_i . Hence let H_i be a family of finite-dimensional Hilbert spaces and set $H = \times_{i=1}^N H_i$, normed by $||y||_H^2 := \sum_{i=1}^N ||y_i||_{H_i}^2$. Let $K_i : \mathcal{M}(\Omega) \to H_i$ be linear and weak* continuous operators.

For a given observation $y \in H$, consider the problem of finding $\rho \in C_w([0,1]; \mathcal{M}(\Omega))$ such that

$$K_i \rho_{t_i} = y_i$$
 for each $i = 1, ..., N_i$

Following [34], we regularize the above problem by

$$\min_{(\rho,m,\mu)\in \tilde{\mathcal{M}}_{\Omega}} \frac{1}{2} \sum_{i=1}^{N} ||K_i \rho_{t_i} - y_i||_{H_i}^2 + J_{\alpha,\beta,\delta}(\rho,m,\mu).$$
(114)

To recast (114) into the form (112), define the linear operator K: $C_w([0,1]; \mathcal{M}(\Omega)) \to H$ as $K\rho := (K_1\rho_{t_1}, ..., K_N\rho_{t_N})$ and note that K is continuous in the sense of (109). Moreover define the fidelity term $F : H \to \mathbb{R}$ by $F(x) := \frac{1}{2}||x - y||_{H^*}^2$ which is convex, lower semi-continuous and bounded from below. In this way (114) is a particular case of (112) and Theorem 6.4 applies, thus showing the existence and characterizing the structure of sparse solutions to the discrete reconstruction problem regularized via the Hellinger-Kantorovich energy.

Acknowledgements

KB and SF are supported by the Christian Doppler Research Association (CDG) and Austrian Science Fund (FWF) through project PIR-27 "Mathematical methods for motion-aware medical imaging" and project P 29192 "Regularization graphs for variational imaging". MC is supported by the Royal Society (Newton International Fellowship NIF\R1\192048). The Institute of Mathematics and Scientific Computing, to which KB and SF are affiliated, is a member of NAWI Graz (http://www.nawigraz.at/en/). The authors KB and SF are members of/associated with BioTechMed Graz (https://biotechmedgraz.at/en/).

Funding

Christian Doppler Forschungsgesellschaft, http://dx.doi.org/10.13039/501100006012; Austrian Science Fund, http://dx.doi.org/10.13039/501100002428; and Royal Society, http://dx.doi.org/ 10.13039/50110000288.

References

- Ambrosio, L. (2004). Transport equation and Cauchy problem for BV vector fields. *Invent. Math.* 158(2):227–260. DOI: 10.1007/s00222-004-0367-2.
- [2] DiPerna, R. J., Lions, P. L. (1989). Ordinary differential equations, transport theory and Sobolev spaces. *Invent Math.* 98(3):511–547. DOI: 10.1007/BF01393835.
- [3] Ambrosio, L., Colombo, M., Figalli, A. (2015). Existence and uniqueness of maximal regular flows for non-smooth vector fields. *Arch Rational Mech. Anal.* 218(2):1043–1081. DOI: 10.1007/s00205-015-0875-9.
- Bernard, P. (2008). Young measures, superposition and transport. *Indiana Univ. Math. J.* 57(1):247–275. DOI: 10.1512/iumj.2008.57.3163.
- [5] Bernard, P., Ambrosio, L. (2008). Uniqueness of signed measures solving the continuity equation for Osgood vector fields. *Atti Della Accad. Naz. Dei Lincei* 19(3):237–245.
- [6] Bianchini, S., Bonicatto, P. (2020). A uniqueness result for the decomposition of vector fields in \mathbb{R}^d . *Invent. Math.* 220(1):255–393.

- [7] Bianchini, S., Bonicatto, P., Gusev, N. A. (2016). Renormalization for autonomous nearly incompressible BV vector fields in two dimensions. SIAM J. Math. Anal. 48(1):1–33. DOI: 10.1137/15M1007380.
- [8] Bonicatto, P., Gusev, N. A. (2018). Superposition principle for the continuity equation in a bounded domain. J. Phys. Conf. Ser. 990:012002. DOI: 10.1088/1742-6596/990/1/012002.
- [9] Lisini, S. (2006). Characterization of absolutely continuous curves in Wasserstein spaces. Calc. Var. 28(1):85–120. DOI: 10.1007/s00526-006-0032-2.
- [10] Maniglia, S. (2007). Probabilistic representation and uniqueness results for measure-valued solutions of transport equations. J. Math. Pures Appli.. 87(6):601–626. DOI: 10.1016/j.matpur.2007.04.001.
- [11] Stepanov, E., Trevisan, D. (2017). Three superposition principles: Currents, continuity equations and curves of measures. J. Funct. Anal. 272(3):1044–1103. DOI: 10.1016/j.jfa. 2016.10.025.
- [12] Ambrosio, L., Gigli, N., Savaré, G. (2005). Gradient Flows: In Metric Spaces and in the Space of Probability Measures. Basel: Birkhäuser.
- [13] Bredies, K., Carioni, M., Fanzon, S., Romero, F. (2021). On the extremal points of the ball of the Benamou-Brenier energy. *Bull. London Math. Soc.* 53(5):1436–1452. DOI: 10.1112/ blms.12509.
- [14] Ambrosio, L., Crippa, G. (2008). Existence, Uniqueness, Stability and Differentiability Properties of the Flow Associated to Weakly Differentiable Vector Fields. Springer Berlin Heidelberg, pp. 3–57.
- [15] Agueh, M., Ghoussoub, N., Kang, X. (2004). Geometric inequalities via a general comparison principle for interacting gases. *Geometric Funct. Anal.* 14(1):215–244. DOI: 10.1007/ s00039-004-0455-x.
- [16] Kantorovitch, L. (1942). On the translocation of masses. Comptes Rendus (Doklady) L'Acad. Des Sci. l'URSS. 37:199–201.
- [17] Santambrogio, F. (2015). Optimal Transport for Applied Mathematicians. Basel: Birkhäuser.
- [18] Villani, C. (2008). Optimal Transport: Old and New. Berlin Heidelberg: Springer.
- [19] Benamou, J.-D., Brenier, Y. (2000). A computational fluid mechanics solution to the Monge-Kantorovich mass transfer problem. *Numer Math.* 84(3):375–393. DOI: 10.1007/ s002110050002.
- [20] Ambrosio, L., Gigli, N., Savaré, G. (2014). Metric measure spaces with Riemannian Ricci curvature bounded from below. *Duke Math. J.* 163(7):1405–1490. DOI: 10.1215/00127094-2681605.
- [21] Ambrosio, L., Gigli, N., Savaré, G. (2014). Calculus and heat flow in metric measure spaces and applications to spaces with Ricci bounds from below. *Invent. Math.* 195(2): 289–391. DOI: 10.1007/s00222-013-0456-1.
- [22] Jordan, R., Kinderlehrer, D., Otto, F. (1997). Free energy and the Fokker–Planck equation. *Physica D*. 107(2-4):265–271. DOI: 10.1016/S0167-2789(97)00093-6.
- [23] Jordan, R., Kinderlehrer, D., Otto, F. (1998). The variational formulation of the Fokker-Planck equation. SIAM J. Math. Anal. 29(1):1–17. DOI: 10.1137/ S0036141096303359.
- [24] Otto, F. (1998). Dynamics of labyrinthine pattern formation in magnetic fluids: A mean-field theory. *Arch. Ration. Mech. Anal.* 141(1):63–103. DOI: 10.1007/s002050050073.
- [25] Otto, F. (2001). The geometry of dissipative evolution equations: the porus medium equation. *Commun. Partial Diff. Eq.* 26(1-2):101–174. DOI: 10.1081/PDE-100002243.
- [26] Cordero-Erausquin, D., Nazaret, B., Villani, C. (2004). A mass-transportation approach to sharp Sobolev and Gagliardo-Nirenberg inequalities. *Adv. Math.* 182(2):307–332. DOI: 10. 1016/S0001-8708(03)00080-X.
- [27] Maggi, F., Villani, C. (2007). Balls have the worst best Sobolev inequalities. Part II: variants and extensions. *Calc. Var.* 31(1):47–74. DOI: 10.1007/s00526-007-0105-x.
- [28] Maggi, F., Villani, C. (2005). Balls have the worst best Sobolev inequalities. J Geom Anal. 15(1):83–121. DOI: 10.1007/BF02921860.

- [29] Nazaret, B. (2006). Best constant in Sobolev trace inequalities on the half-space. *Nonlinear Anal. Theory, Methods Appl.* 65(10):1977–1985. DOI: 10.1016/j.na.2005.05.069.
- [30] Otto, F., Villani, C. (2000). Generalization of an inequality by Talagrand and links with the logarithmic Sobolev inequality. *J. Funct. Anal.* 173(2):361–400. DOI: 10.1006/jfan.1999. 3557.
- [31] Chizat, L., Peyré, G., Schmitzer, B., Vialard, F.-X. (2018). An interpolating distance between optimal transport and Fisher-Rao metrics. *Found Comput. Math.* 18(1):1-44. DOI: 10.1007/s10208-016-9331-y.
- [32] Kondratyev, S., Monsaingeon, L., Vorotnikov, D. (2016). A new optimal transport distance on the space of finite Radon measures. *Adv. Diff. Eq.* 21(11/12):1117–1164.
- [33] Liero, M., Mielke, A., Savaré, G. (2018). Optimal Entropy-Transport problems and a new Hellinger-Kantorovich distance between positive measures. *Invent. Math.* 211(3):969–1117. DOI: 10.1007/s00222-017-0759-8.
- [34] Bredies, K., Fanzon, S. (2020). An optimal transport approach for solving dynamic inverse problems in spaces of measures. *Esaim: M2an.* 54(6):2351–2382. DOI: 10.1051/m2an/ 2020056.
- [35] Chizat, L., Peyré, G., Schmitzer, B., Vialard, F.-X. (2018). Scaling algorithms for unbalanced transport problems. *Math. Comp.* 87(314):2563–2609. DOI: 10.1090/mcom/3303.
- [36] Liero, M., Mielke, A., Savaré, G. (2016). Optimal transport in competition with reaction: The Hellinger-Kantorovich distance and geodesic curves. SIAM J. Math. Anal. 48(4): 2869–2911. DOI: 10.1137/15M1041420.
- [37] Schmitzer, B., Wirth, B. (2019). Dynamic models of Wasserstein-1-type unbalanced transport. *Esaim Cocv.* 25:23. DOI: 10.1051/cocv/2018017.
- [38] Bredies, K., Carioni, M., Fanzon, S., Romero, F. (2022). A generalized conditional gradient method for dynamic inverse problems with optimal transport regularization. *Found Comput. Math.* DOI: 10.1007/s10208-022-09561-z.
- [39] Bredies, K., Carioni, M., Fanzon, S., Walter, D. (2021). Linear convergence of accelerated generalized conditional gradient methods. arXiv e-Prints, arXiv:2110.06756.
- [40] Boyer, C., Chambolle, A., Castro, Y. D., Duval, V., de Gournay, F., Weiss, P. (2019). On representer theorems and convex regularization. SIAM J. Optim. 29(2):1260–1281. DOI: 10.1137/18M1200750.
- [41] Bredies, K., Carioni, M. (2020). Sparsity of solutions for variational inverse problems with finite-dimensional data. *Calc. Var.* 59(1):14. DOI: 10.1007/s00526-019-1658-1.
- [42] Ambrosio, L., Fusco, N., Pallara, D. (2000). Functions of Bounded Variation and Free Discontinuity Problems. Clarendon Press, Oxford: Oxford Science Publications.
- [43] Evans, L. C., Gariepy, R. F. (2015). *Measure Theory and Fine Properties of Functions*. Boca Raton, Florida: CRC Press.
- [44] Bogachev, V. I. (2007). *Measure Theory*. Berlin Heidelberg: Springer-Verlag.
- [45] Unser, M. (2021). A unifying representer theorem for inverse problems and machine learning. Found Comput. Math. 21(4):941–960. DOI: 10.1007/s10208-020-09472-x.
- [46] Unser, M., Fageot, J., Ward, J. P. (2017). Splines are universal solutions of linear inverse problems with generalized TV regularization. SIAM Rev. 59(4):769–793. DOI: 10.1137/ 16M1061199.
- [47] Kechris, A. S. (1995). Classical Descriptive Set Theory. New York: Springer-Verlag.

Appendix A

[A.1.] Properties of narrow convergence

We give some results about narrow convergence of measures. These results are classical and are stated for probability measures in the literature: here we adapt them to positive measures. For a complete separable metric space Y, we say that a family of measures $\mathcal{A} \subset \mathcal{M}(Y)$ is *tight* if for every $\varepsilon > 0$ there exists a compact set $K_{\varepsilon} \subset Y$ such that $|\mu|(Y \setminus K_{\varepsilon}) < \varepsilon$ for all $\mu \in \mathcal{A}$. The next proposition provides a tightness criterion for positive measures. The proof follows as in [12, Remark 5.1.5], and is hence omitted.

Proposition A.1. Let Y be a complete separable metric space and $\mathcal{A} \subset \mathcal{M}^+(Y)$. Suppose that there exists a measurable function $\mathscr{F} : Y \to [0, \infty]$ such that $\{y \in Y : \mathscr{F}(y) \leq c\}$ is compact for each $c \geq 0$ and $\sup_{\mu \in \mathcal{A}} \int_Y \mathscr{F}(y) \ d\mu(y) < \infty$. Then \mathcal{A} is tight.

Finally, we provide a result which clarifies the behavior of narrowly convergent sequences of positive measures when tested against lower semi-continuous, or continuous unbounded integrands. The proof easily follows by combining [12, Lemma 5.1.7] with a scaling argument.

Proposition A.2. Let Y be a complete separable metric space. Assume that $\{\mu_n\}_n$, μ belong to $\mathcal{M}^+(Y)$ and $\mu_n \to \mu$ narrowly as $n \to \infty$. If $g: Y \to [0, \infty]$ is lower semi-continuous then

$$\int_{Y} g(y) \ d\mu(y) \le \liminf_{n \to \infty} \int_{Y} g(y) \ d\mu_n(y).$$
(115)

If $f: Y \to \mathbb{R}$ is continuous with |f| uniformly integrable with respect to $\{\mu_n\}_n$, that is,

$$\lim_{k \to \infty} \sup_{n \in \mathbb{N}} \int_{\{y \in Y: |f(y)| \ge k\}} |f(y)| \ d\mu_n(y) = 0, \tag{116}$$

then it holds

$$\lim_{n \to \infty} \int_{Y} f(y) \ d\mu_n(y) = \int_{Y} f(y) \ d\mu(y).$$
(117)

[A.2] Disintegration of measures

In this section we state and prove the disintegration theorem employed in Section 5. This result is a straightforward consequence of [12, Theorem 5.3.1].

Theorem A.3. Let Z, X be Radon separable metric spaces and let $\mu \in \mathcal{M}^+(Z)$, $\nu \in \mathcal{M}^+(X)$ be given. Let $\pi : Z \to X$ a measurable map such that $\pi_{\#} \mu \ll \nu$. Then there exists a Borel family of measures $\{\mu^x\}_{x\in X} \subset \mathcal{M}^+(Z)$ such that

(i) $\mu^x(Z \setminus \pi^{-1}(x)) = 0$ for ν -a.e. $x \in X$,

(ii) for every function $f \in L^1_\mu(Z)$ there holds

$$\int_{Z} f(z) \ d\mu(z) = \int_{X} \int_{Z} f(z) \ d\mu^{x}(z) \ d\nu(x), \tag{118}$$

(iii) if μ is concentrated on $E \subset Z$, then μ^x is concentrated on $\pi^{-1}(x) \cap E$ for ν -a.e. $x \in X$.

Proof. Without loss of generality we can suppose that $\mu \neq 0$. By a rescaling argument we can assume that $||\mu||_{\mathcal{M}(Z)} = 1$ as well. Thanks to [12, Theorem 5.3.1] there exists a Borel family of measures $\{\tilde{\mu}^x\}_{x\in X} \subset \mathcal{M}^+(Z)$ such that $\tilde{\mu}^x(Z \setminus \pi^{-1}(x)) = 0$ for $(\pi_{\#}\mu)$ -a.e. $x \in X$, and that (118) holds with μ^x and ν replaced by $\tilde{\mu}^x$ and $\pi_{\#}\mu$, respectively, for every Borel function $f: Z \to [0, \infty]$. For all $x \in X$ set $\mu^x := \frac{\partial(\pi_{\#}\mu)}{\partial \nu}(x) \tilde{\mu}^x$. We immediately obtain that $\mu^x \in \mathcal{M}^+(Z)$ is a family

of Borel measures satisfying (i). Moreover, it is easy to check that μ^x satisfies (118) for every Borel function $f: Z \to [0, \infty]$. If $f \in L^1_{\mu}(Z)$, by (118) we get $f \in L^1_{\mu^x}(Z)$ for ν -a.e. $x \in X$, yielding (*ii*). Finally, (*iii*) is implied by (*ii*).

[A.3.] Properties of B_{δ} and J_{α} , β , δ

In this section we gather some of the properties of the functionals B_{δ} and $J_{\alpha,\beta,\delta}$ introduced in Section 2.2. The interested reader can find the proofs of such results in Proposition 2.6 and Lemmas 4.5, 4.6 in [34].

Lemma A.4 (Properties of B_{δ}). The functional B_{δ} defined in (16) is non-negative, convex, onehomogeneous and sequentially lower semi-continuous with respect to the weak^{*} topology on \mathcal{M}_{Ω} . Moreover it satisfies the following properties:

- (i) if $B_{\delta}(\rho, m, \mu) < \infty$, then $\rho \ge 0$ and $m, \mu \ll \rho$, that is, there exist measurable maps $v : X_{\Omega} \to \mathbb{R}^d$, $g : X_{\Omega} \to \mathbb{R}$ such that $m = v\rho$, $\mu = g\rho$,
- (ii) if $\rho \ge 0$ and $m = v\rho$, $\mu = g\rho$ for some measurable $v : X_{\Omega} \to \mathbb{R}^d$, $g : X_{\Omega} \to \mathbb{R}$, then

$$B_{\delta}(\rho, m, \mu) = \int_{X_{\Omega}} \Psi_{\delta}(1, \nu, g) \ d\rho = \frac{1}{2} \int_{X_{\Omega}} \left(|\nu|^2 + \delta^2 g^2 \right) \ d\rho.$$
(119)

Lemma A.5 (Properties of $J_{\alpha,\beta,\delta}$). Let $\alpha, \beta, \delta > 0$. The functional $J_{\alpha,\beta,\delta}$ is non-negative, convex, one-homogeneous and sequentially lower semi-continuous with respect to the weak* topology on \mathcal{M}_{Ω} . For $(\rho, m, \mu) \in \mathcal{M}_{\Omega}$ such that $J_{\alpha,\beta,\delta}(\rho, m, \mu) < \infty$ we have that

$$\max\{\alpha \|\rho\|_{\mathcal{M}(X_{\Omega})}, \ C\|m\|_{\mathcal{M}(X_{\Omega};\mathbb{R}^{d})}, \ C\|\mu\|_{\mathcal{M}(X_{\Omega})}\} \le J_{\alpha,\beta,\delta}(\rho,m,\mu)$$
(120)

where $C := \min\{2\alpha, \beta\min\{1, \delta^2\}\}$. If in addition the sequence $\{(\rho^n, m^n, \mu^n)\}_n$ in \mathcal{M}_Ω is such that $\sup_n J_{\alpha,\beta,\delta}(\rho^n, m^n, \mu^n) < \infty$, then $\rho^n = dt \otimes \rho_t^n$ for some $(t \mapsto \rho_t^n) \in C_w([0,1]; \mathcal{M}^+(\Omega))$ and there exists (ρ, m, μ) in \mathcal{D}_Ω with $\rho = dt \otimes \rho_t$, $(t \mapsto \rho_t) \in C_w([0,1]; \mathcal{M}^+(\Omega))$, such that, up to subsequences, $(\rho^n, m^n, \mu^n) \longrightarrow (\rho, m, \mu)$ weakly* in \mathcal{M}_Ω and $\rho_t^n \to \rho_t$ narrowly in $\mathcal{M}(\Omega)$ for every $t \in [0, 1]$.

[A.4.] Proof of Proposition 3.6

Remember that $\mathscr{G}_V = C([0,1];\mathscr{G}_V)$ by Proposition 3.5. Therefore, in order to prove that (\mathscr{G}_V, d) is complete and separable, it is sufficient to show that (\mathscr{G}_V, d_F) is complete and separable (see Theorem 4.19 in [47]). Let us first prove that (\mathscr{G}_V, d_F) is complete. Hence, let $\rho^n = h_n \delta_{\gamma_n} \in \mathscr{G}_V$ be a Cauchy sequence. By (29) we have $|h_n - h_m| \leq d_F(\rho^n, \rho^m)$ for all $m, n \in \mathbb{N}$. Therefore $h_n \to h$ for some $h \geq 0$. If h = 0, by (29) we have $d_F(\rho^n, 0) = h_n \to 0$, showing that ρ^n converges to $0 \in \mathscr{G}_V$. Assume now that h > 0. Notice that $|\gamma_n - \gamma_m| \leq 2$ for sufficiently large m, n, otherwise we could extract a subsequence (not relabelled) such that $d_F(\rho^n, \rho^m) = h_n + h_m \to 2h > 0$ as $m, n \to \infty$, which contradicts ρ^n being Cauchy. By (29) and the facts that $h_n \to h > 0$ and that ρ^n is Cauchy, we get that γ_n is Cauchy, so that $\gamma_n \to \gamma \in V$. An application of (29) shows that $\rho^n \to \rho := h \delta_\gamma$ with respect to d_F , concluding completeness. The fact that (\mathscr{G}_V, d) is separable is immediate: indeed the countable set $\mathscr{C}'_V := \left\{h\delta_\gamma : h \in [0,\infty) \cap \mathbb{Q}, \ \gamma \in V \cap \mathbb{Q}^d\right\} \subset \mathscr{C}_V$ is d_F -dense in \mathscr{C}_V , since V is the closure of a domain.

[A.5.] Comparison principle

In this section we recall a comparison principle for signed measure solutions of the continuity equation. **Proposition A.6** (Comparison principle). Let $\rho_t : [0,1] \to \mathcal{M}(\mathbb{R}^d)$ be narrowly continuous and $v : (0,1) \times \mathbb{R}^d \to \mathbb{R}^d$, $g : (0,1) \times \mathbb{R}^d \to \mathbb{R}$ be measurable. Suppose that $\partial_t \rho_t + \text{div} (v\rho_t) = g\rho_t$ holds in $(0,1) \times \mathbb{R}^d$ in the sense of (14). Assume that $\rho_0 \leq 0$, as well as (18, 20) and

$$\int_{0}^{1} \int_{\mathbb{R}^{d}} \left(|\nu(t,x)| + |g(t,x)| \right) \, d|\rho_{t}|(x) \, dt < \infty.$$
(121)

Then $\rho_t \leq 0$ for all $t \in [0, 1]$.

A proof of the above proposition can be found in [10, Lemma 3.5]. We just point out that in [10, Lemma 3.5] it is assumed that the narrowly continuous curve $t \mapsto \rho_t \in \mathcal{M}(\mathbb{R}^d)$ satisfies $\int_0^1 |\rho_t|(B) \, dt < \infty$ for all $B \subset \mathbb{R}^d$ compact. However this condition is always fulfilled, since ρ_t automatically satisfies $\sup_{t \in [0,1]} \|\rho_t\|_{\mathcal{M}(\mathbb{R}^d)} < \infty$, as shown in [34, Proposition A.3]. Moreover the statement of [10, Lemma 3.5] also requires that g is bounded: after carefully inspecting the proof, we noticed that such assumption is not needed.

[A.6.] Property of convolutions

Here we recall a result on convolution of measures, which can be found in [10, Lemma 3.9].

Proposition A.7. Let $p \ge 1$, $\rho \in \mathcal{M}^+(\mathbb{R}^d)$, $E \in \mathcal{M}(\mathbb{R}^d, \mathbb{R}^m)$ and ξ be a convolution kernel on \mathbb{R}^d . Suppose that E is absolutely continuous with respect to ρ . Then,

$$\int_{\mathbb{R}^d} \left| \frac{E * \xi}{\rho * \xi} \right|^p \ \left(\rho * \xi \right) \ dx \leq \int_{\mathbb{R}^d} \left| \frac{dE}{d\rho} \right|^p \ d\rho,$$

where $dE/d\rho$ is the Radon–Nikodym derivative of E with respect to ρ .