# A new transportation distance between non-negative measures, with applications to gradients flows with Dirichlet boundary conditions 

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#### Abstract

In this paper we introduce a new transportation distance between non-negative measures inside a domain $\Omega$. This distance enjoys many nice properties, for instance it makes the space of non-negative measures inside $\Omega$ a geodesic space, without any convexity assumption on $\Omega$. Moreover, we will show that the gradient flow of the entropy functional $\int_{\Omega}[\rho \log (\rho)-\rho] d x$ w.r.t. this distance coincides with the heat equation, subject to the Dirichlet boundary condition equal to 1 .

\section*{Résumé}

Dans ce papier, nous introduisons une nouvelle distance sur l'espace des mesures positive dans un domaine $\Omega$. Cette distance satisfait plusieurs propriétés intéressantes : par exemple, elle fait de l'espace des mesures positives dans $\Omega$ un espace géodésique, sans aucune hypothèse de convexité sur le domaine. De plus, on montre que le flot gradient de la fonctionnelle d'entropie $\int_{\Omega}[\rho \log (\rho)-\rho] d x$ par rapport à cette distance donne lieu à l'équation de la chaleur, avec condition de Dirichlet égale à 1 sur le bord.


## 1 Introduction

Nowadays, it is well-know that transportation distances between probability measures can be successfully used to study evolutionary equations. More precisely, one of the most surprisingly achievement of $[8,10,11]$ has been that many evolution equations of the form

$$
\frac{d}{d t} \rho(t)=\operatorname{div}(\nabla \rho(t)-\rho(t) \nabla V-\rho(t)(\nabla W * \rho(t)))
$$

can be seen as gradient flows of some entropy functionals on the space of probability measures with respect to the Wasserstein distance

$$
W_{2}(\mu, v)=\inf \left\{\sqrt{\int|x-y|^{2} d \gamma(x, y)}: \pi_{\#}^{1} \gamma=\mu, \pi_{\#}^{2} \gamma=v\right\} .
$$

[^0]Besides the fact that this interpretation allows to prove entropy estimates and functional inequalities (see [12, 13] for more details on this area, which is still very active and in continuous evolution), this point of view provides a powerful variational method to prove esistence of solutions to the above equations: given a time step $\tau>0$, construct an approximate solution by iteratively minimizing

$$
\rho \quad \mapsto \quad \frac{W_{2}\left(\rho, \rho_{0}\right)}{2 \tau}+\int\left[\rho \log (\rho)+\rho V+\frac{1}{2} \rho(W * \rho)\right] d x
$$

We refer to [2] for a general description of this approach.
Let us observe that the choice of the distance on the space of probability measures plays a key role, and by changing it one can construct solutions to more general classes of evolution equations, see for instance [1, 5, 7]. However, all the distances considered up to now need the two measures to have the same mass (which up to a scaling can always be assumed equal to 1 ). In particular, since the mass remains constant along the evolution, if one restricts to measures concentrated on a bounded domain, then the approach described above will always produce solutions to parabolic equations with Neumann boundary conditions.

Motivated by the intent to find an analogous approach to construct solutions of evolution equations subject to Dirichlet boundary condition, in this paper we introduce a new transportation distance $W b_{2}$ between measures. As we will see, the main features of the distance $W b_{2}$ are:

- It metrizes the weak convergence of positive measures $\mathcal{M}_{+}(\Omega)$ in $\Omega$, see Proposition 2.2. (This is similarly to what happens for the common Wasserstein distances, but without any mass constraint.)
- The resulting metric space $\left(\mathcal{M}_{+}(\Omega), W b_{2}\right)$ is always geodesic, see Proposition 2.8. This is a particularly interesting property compared to what happens in the classical Wasserstein space: indeed the space $\left(\mathscr{P}(\Omega), W_{2}\right)$ is geodesic if and only if $\Omega$ is convex. In our case, the convexity of the open set is not required. (Actually, not even connectedness is needed!)
- The natural approach via minimizing movements to the study of the gradient flow of the entropy leads to weak solution of the heat equation with Dirichlet boundary condition, see Theorem 3.5. Interesting enough, with this approach the regularity of the boundary of $\Omega$ does not play any role.

As a drawback, the entropy functional do not have the same nice properties it has in the classical Wasserstein space. In particular:

- It is not geodesically convex. Still, it has some sort of convexity properties along geodesics, see Remark 3.4.
- Due to the lack of geodesic convexity, we were not able to prove any kind of contractivity result for the flow.
- Actually, we are not even able to prove uniqueness of the limit of the minimizing movements scheme. (Of course one knows by standard PDE techniques that weak solutions of the heat equation with Dirichlet boundary conditions are unique, therefore a posteriori it is clear that the limit has to be unique - what we are saying here is that we do not know whether such uniqueness may be deduced a priori via techniques similar, e.g., to those appeared in [2].)

The distance $W b_{2}$ is defined in the following way (the ' $b$ ' stands to recall that we have some room to play with the boundary of $\Omega$ ). Let $\Omega \subset \mathbb{R}^{d}$ be a bounded open set, and let $\mathcal{M}_{+}(\Omega)$ denote the space of non-negative finite measures on $\Omega$. We define the distance $W b_{2}$ on $\mathcal{M}_{+}(\Omega)$ as a result of the following problem:

Problem 1.1 (A variant of the transportation problem) Let $\mu, v \in \mathcal{M}_{+}(\Omega)$. The set of admissible couplings $\operatorname{ADM}(\mu, v)$ is defined as the set of positive measures $\gamma$ on $\bar{\Omega} \times \bar{\Omega}$ satisfying

$$
\begin{equation*}
\pi_{\#}^{1} \gamma_{\left.\right|_{\Omega}}=\mu, \quad \pi_{\#}^{2} \gamma_{\left.\right|_{\Omega}}=v \tag{1}
\end{equation*}
$$

For any non-negative measure $\gamma$ on $\bar{\Omega} \times \bar{\Omega}$, we define its cost $C(\gamma)$ as

$$
C(\gamma):=\int_{\bar{\Omega} \times \bar{\Omega}}|x-y|^{2} d \gamma(x, y)
$$

The distance $W b_{2}(\mu, v)$ is then defined as:

$$
W b_{2}^{2}(\mu, v):=\inf _{\boldsymbol{\gamma} \in \operatorname{ADM}(\mu, v)} C(\gamma)
$$

The difference between $W b_{2}$ and $W_{2}$ relies on the fact that an admissible coupling is a measure on the closure of $\Omega \times \Omega$, rather than just on $\Omega \times \Omega$, and that the marginals are required to coincide with the given measures only inside $\Omega$. This means that we can use $\partial \Omega$ as an infinite reserve: we can 'take' as mass as we wish from the boundary, or 'give' it back some of the mass, provided we pay the transportation cost. This is why this distance is well defined for measures which do not have the same mass.


Figure 1: Example of admissible transport plan
Although this approach could be applied for more general costs than just $|x-y|^{2}$ and for a wider class of entropy functionals, we preferred to provide a complete result only in the particular case of the heat equation, in order to avoid technicalities and generalizations which would just obscure the main ideas. We refer to Section 4 for some possible generalizations, a comparison between our and the classical $L^{2}$-approach, and some open problems.

## 2 General properties of the distance $W b_{2}$

The aim of this section is to describe the main properties of the distance $W b_{2}$.
Let $\gamma$ be a non-negative on $\bar{\Omega} \times \bar{\Omega}$. We will write $\gamma_{A}^{B}$ for the restriction of $\gamma$ to the rectangle $A \times B \subset \bar{\Omega} \times \bar{\Omega}$. Observe that there is a natural splitting of $\gamma$ into 4 parts:

$$
\gamma=\gamma_{\Omega}^{\Omega}+\gamma_{\Omega}^{\partial \Omega}+\gamma_{\partial \Omega}^{\Omega}+\gamma_{\partial \Omega}^{\partial \Omega} .
$$

We now remark that, if $\gamma \in \operatorname{Adm}(\mu, v)$, then

$$
\gamma-\gamma_{\partial \Omega}^{\partial \Omega} \in \operatorname{ADM}(\mu, v) \quad \text { and } \quad C\left(\gamma-\gamma_{\partial \Omega}^{\partial \Omega}\right) \leq C(\gamma) .
$$

Hence, when looking for optimal plans, it is not restrictive to assume that

$$
\begin{equation*}
\gamma_{\partial \Omega}^{\partial \Omega}=0 \tag{2}
\end{equation*}
$$

This gives the bound

$$
\begin{align*}
\gamma(\bar{\Omega} \times \bar{\Omega}) & =\gamma(\bar{\Omega} \times \bar{\Omega} \backslash \partial \Omega \times \partial \Omega)  \tag{3}\\
& \leq \gamma(\Omega \times \bar{\Omega})+\gamma(\bar{\Omega} \times \Omega)=\mu(\Omega)+v(\Omega)
\end{align*}
$$

from which it follows the compactness of the set of admissible plans satisfying (2) w.r.t. the topology of weak convergence of measures [2]. Thus optimal plans always exist. We will denote the set of optimal plans by $\operatorname{Opt}(\mu, v)$, and we will always assume that an optimal plan satisfies (2).

To show that $W b_{2}$ satisfies the triangle inequality, we first prove a variant of the classical gluing lemma (see [2, Lemma 5.3.2]):

Lemma 2.1 (A variant of the gluing lemma) Fix $\mu_{1}, \mu_{2}, \mu_{3} \in \mathcal{M}_{+}(\Omega)$, and let $\boldsymbol{\gamma}^{12} \in$ $\operatorname{Adm}\left(\mu_{1}, \mu_{2}\right), \gamma^{23} \in \operatorname{Adm}\left(\mu_{2}, \mu_{3}\right)$ such that $\left(\gamma^{12}\right)_{\partial \Omega}^{\partial \Omega}=\left(\gamma^{23}\right)_{\partial \Omega}^{\partial \Omega}=0$. Then there exist $\gamma^{123} \in \mathcal{M}(\bar{\Omega} \times \bar{\Omega} \times \bar{\Omega})$ such that

$$
\begin{aligned}
& \pi_{\#}^{12} \gamma^{123}=\gamma^{12}+\sigma^{12}, \\
& \pi_{\#}^{23} \gamma^{123}=\gamma^{23}+\sigma^{23},
\end{aligned}
$$

where $\boldsymbol{\sigma}^{12}$ and $\boldsymbol{\sigma}^{23}$ are both concentrated on the diagonal of $\partial \Omega \times \partial \Omega$, i.e. on the set of pairs of points $\{(x, x): x \in \partial \Omega\}$.

Let us point out that, in contrast with the classical result, in our case the second marginal of $\gamma^{12}$ on $\bar{\Omega}$ does not necessarily coincides with the first marginal of $\gamma^{23}$, and so the two measures cannot be 'glued' together in a trivial way.
Proof. In order to clarify the structure of the proof, it is convenient to see $\mu_{1}, \mu_{2}, \mu_{3}$ as measures on $\mathcal{M}_{+}\left(\Omega_{1}\right), \mathcal{M}_{+}\left(\Omega_{2}\right), \mathcal{M}_{+}\left(\Omega_{3}\right)$ respectively, where $\Omega_{1}, \Omega_{2}, \Omega_{3}$ are three distinct copies of $\Omega$. In this way we have $\gamma^{12} \in\left(\bar{\Omega}_{1} \times \bar{\Omega}_{2}\right), \gamma^{23} \in\left(\bar{\Omega}_{2} \times \bar{\Omega}_{3}\right)$, and $\gamma^{123} \in \mathcal{M}\left(\bar{\Omega}_{1} \times \bar{\Omega}_{2} \times \bar{\Omega}_{3}\right)$. However, since in fact $\Omega_{1}=\Omega_{2}=\Omega_{3}$, sometimes we will identify $\Omega_{2}$ with $\Omega, \Omega_{1}$, or $\Omega_{3}$. Furthermore, we will use $\pi^{2}$ to denote both the canonical projection from $\bar{\Omega}_{1} \times \bar{\Omega}_{2}$ onto $\bar{\Omega}_{2}$, and the one from $\bar{\Omega}_{2} \times \bar{\Omega}_{3}$ onto $\bar{\Omega}_{2}$.

Let us define

$$
\begin{aligned}
& \sigma^{12}:=\left(\pi^{2}, \pi^{2}\right)_{\#}\left(\left(\gamma^{23}\right)_{\partial \Omega}^{\Omega}\right) \in \mathcal{M}\left(\partial \Omega_{2} \times \partial \Omega_{2}\right)=\mathcal{M}\left(\partial \Omega_{1} \times \partial \Omega_{2}\right), \\
& \sigma^{23}:=\left(\pi^{2}, \pi^{2}\right)_{\#}\left(\left(\gamma^{12}\right)_{\Omega}^{\partial \Omega}\right) \in \mathcal{M}\left(\partial \Omega_{2} \times \partial \Omega_{2}\right)=\mathcal{M}\left(\partial \Omega_{2} \times \partial \Omega_{3}\right),
\end{aligned}
$$

and set

$$
\begin{aligned}
\tilde{\gamma}^{12} & :=\gamma^{12}+\sigma^{12} \in \mathcal{M}\left(\partial \Omega_{1} \times \partial \Omega_{2}\right), \\
\tilde{\gamma}^{23} & :=\gamma^{23}+\sigma^{23} \in \mathcal{M}\left(\partial \Omega_{2} \times \partial \Omega_{3}\right) .
\end{aligned}
$$

Observe that

$$
\begin{equation*}
\left(\pi_{\#}^{2} \tilde{\gamma}^{12}\right)_{\left.\right|_{\Omega}}=\left(\pi_{\#}^{2} \gamma^{12}\right)_{\left.\right|_{\Omega}}+\left(\pi_{\#}^{2} \sigma^{12}\right)_{\left.\right|_{\Omega}}=\mu_{2}, \tag{4}
\end{equation*}
$$

and similarly $\left(\pi_{\#}^{2} \tilde{\gamma}^{23}\right)_{\left.\right|_{\Omega}}=\mu_{2}$. Moreover, since by assumption $\left(\gamma^{12}\right)_{\partial \Omega}^{\partial \Omega}=0$,

$$
\left(\pi_{\#}^{2} \gamma^{12}\right)_{\left.\right|_{\partial \Omega}}=\pi_{\#}^{2}\left(\left(\gamma^{12}\right)_{\Omega}^{\partial \Omega}\right)=\pi_{\#}^{2} \sigma^{23} .
$$

and similarly $\left(\pi_{\#}^{2} \gamma^{23}\right)_{\partial \Omega}=\pi_{\#}^{2} \sigma^{12}$. Therefore we obtain

$$
\begin{equation*}
\left(\pi_{\#}^{2} \tilde{\gamma}^{12}\right)_{\left.\right|_{\partial \Omega}}=\left(\pi_{\#}^{2} \gamma^{12}\right)_{\left.\right|_{\partial \Omega}}+\pi_{\#}^{2} \sigma^{12}=\pi_{\#}^{2} \sigma^{23}+\left(\pi_{\#}^{2} \gamma^{23}\right)_{\left.\right|_{\partial \Omega}}=\left(\pi_{\#}^{2} \tilde{\gamma}^{23}\right)_{\left.\right|_{\partial \Omega}} . \tag{5}
\end{equation*}
$$

Thanks to (4) and (5) we finally obtain $\pi_{\#}^{2} \gamma^{12}=\pi_{\#}^{2} \gamma^{23}$, and so the conclusion follows from the classical gluing lemma, see for instance [2, Lemma 5.3.2].

Proposition 2.2 The function $W b_{2}$ is a distance on $\mathcal{M}_{+}(\Omega)$ which metrizes the weak convergence, i.e., the topology given by duality with $C_{c}(\Omega)$.
Proof. The facts that $W b_{2}(\mu, v)=0$ if and only if $\mu=v$ and the symmetry are obvious. For the triangle inequality we need to use the version of gluing lemma we just proved. Fix $\mu_{1}, \mu_{2}, \mu_{3} \in \mathcal{M}_{+}(\Omega)$ and let $\gamma^{12}, \gamma^{23}$ be two optimal plans from $\mu_{1}$ to $\mu_{2}$ and from $\mu_{2}$ to $\mu_{3}$ respectively. Use lemma 2.1 to find a 3-plan $\gamma^{123}$ such that $\pi_{\#}^{1,2} \gamma^{123}=\gamma^{12}+\sigma^{12}$ and $\pi_{\#}^{2,3} \boldsymbol{\gamma}^{123}=\boldsymbol{\gamma}^{23}+\boldsymbol{\sigma}^{23}$, with $\boldsymbol{\sigma}^{12}$ and $\boldsymbol{\sigma}^{12}$ concentrated on the diagonals of $\partial \Omega \times \partial \Omega$. Then we have $\left(\pi_{\#}^{1} \gamma^{123}\right)_{\left.\right|_{\Omega}}=\left(\pi_{\#}^{1} \gamma^{12}+\sigma^{12}\right)_{\left.\right|_{\Omega}}=\mu_{1}$. Similarly, we have $\left(\pi_{\#}^{3} \gamma^{123}\right)_{\left.\right|_{\Omega}}=\mu_{3}$, therefore $\pi^{1,2} \boldsymbol{\gamma}^{123} \in \operatorname{ADM}\left(\mu_{1}, \mu_{3}\right)$ and it holds

$$
\begin{aligned}
W b_{2}\left(\mu_{1}, \mu_{3}\right) & \leq \sqrt{\int\left|x_{1}-x_{3}\right|^{2} d \gamma^{123}} \\
& \leq \sqrt{\int\left|x_{1}-x_{2}\right|^{2} d \gamma^{123}}+\sqrt{\int\left|x_{2}-x_{3}\right|^{2} d \gamma^{123}} \\
& =\sqrt{\int\left|x_{1}-x_{2}\right|^{2} d\left(\gamma^{12}+\sigma^{12}\right)}+\sqrt{\int\left|x_{2}-x_{3}\right|^{2} d\left(\gamma^{23}+\sigma^{23}\right)} \\
& =\sqrt{\int\left|x_{1}-x_{2}\right|^{2} d \gamma^{12}}+\sqrt{\int\left|x_{2}-x_{3}\right|^{2} d \gamma^{23}} \\
& =W b_{2}\left(\mu_{1}, \mu_{2}\right)+W b_{2}\left(\mu_{2}, \mu_{3}\right),
\end{aligned}
$$

where in the fourth step we used the fact that $\sigma^{12}$ and $\sigma^{23}$ are concentrated on a diagonal. Finally, the fact that $W b_{2}$ metrizes the weak topology can be proved as in [2, Proposition 7.1.5] - we omit the details.

Remark 2.3 Note carefully that we are speaking of weak convergence in duality with functions with compact support in $\Omega$, and not, e.g., with continuous and bounded functions in $\bar{\Omega}$. Indeed, the mass can 'disappear' inside the boundary, so that in general we only have

$$
\liminf _{n \rightarrow \infty} \mu_{n}(\Omega) \geq \mu(\Omega),
$$

for any sequence $\left\{\mu_{n}\right\}_{n \in \mathbb{N}} \subset \mathcal{M}_{+}(\Omega)$ such that $W b_{2}\left(\mu_{n}, \mu\right) \rightarrow 0$.

Proposition 2.4 (Cyclical monotonicity) Let $\mu, v \in \mathcal{M}_{+}(\Omega)$. Then there exists a cyclically monotone set $\Gamma \subset \bar{\Omega} \times \bar{\Omega}$ such that any optimal plan is concentrated on $\Gamma$.

Proof. Recall that we always assume that condition (2) holds for optimal plans. From the uniform bound (3), and the fact that the set of optimal plans is a closed under weak convergence (which can been easily proven by adapting the proof of [13, Theorem 5.20]), we deduce that $\operatorname{OPT}(\mu, v)$ is compact w.r.t. the weak convergence of measures. Let $\left(\gamma_{i}\right)_{i \geq 0} \subset \operatorname{Opt}(\mu, v)$ be a countable dense subset, and define

$$
\gamma:=\sum_{i \geq 0} \frac{1}{2^{i}} \gamma_{i} .
$$

Then it is easy to check by the convexity of the constraints (1), and the linearity of the cost, that $\gamma \in \operatorname{Opt}(\mu, v)$. Furthermore, since its support (i.e. the smallest closed set on which it is concentrated) contains the supports of all the $\gamma_{i}$ 's, and since they are dense inside $\operatorname{Opt}(\mu, v)$, the support of $\gamma$ contains that of any optimal plan.

We now observe that, since $\gamma$ is optimal for the Problem 1.1, it is also optimal for the classical optimal transport problem with cost $|x-y|^{2}$ for the measures $\pi_{\#}^{1} \gamma$ and $\pi_{\#}^{2} \gamma$. (This follows from the fact that any measure $\tilde{\gamma}$ with the same marginals as $\gamma$ belongs to $\operatorname{Adm}(\mu, v)$.) Hence the conclusion follows by the general theory of optimal transport (see for instance [2, Chapter 6]).

Remark 2.5 The idea on which is based the proof of the above proposition is wellknown for the classical transport problem. Recently, the first author used the same tool to prove a similar result for the optimal partial transport problem (see [6]). The result proven here is not covered by previous theorems on the topic, as the marginals of the admissible plans are neither fixed, nor dominated.

The fact that the same idea works in so many different situations, shows the power of the idea itself: observe also that here the fact that the cost function is the squared distance does not play any role. Therefore a similar statement holds for much more general cost functions (we will not stress this point any further).

Remark 2.6 The cyclically monotone subset of $\bar{\Omega} \times \bar{\Omega}$ can always be chosen so that it contains the diagonal of $\partial \Omega \times \partial \Omega$. Indeed, just add to any optimal plan a measure of a given amount of mass, say 1 , concentrated on the diagonal of $\partial \Omega \times \partial \Omega$ : this does not affect neither the cost nor the compactness of the set of optimal plans. Thus the conclusion follows as above.

From now on, $P: \Omega \rightarrow \partial \Omega$ will be a measurable map such that

$$
|x-P(x)|=d(x, \partial \Omega) \quad \forall x \in \Omega
$$

It is well-known that such a map is uniquely defined on $\mathcal{L}^{d}$-a.e. $x \in \Omega$. (Indeed, $P(x)$ is uniquely defined whenever the Lipschitz function $d(\cdot, \partial \Omega)$ is differentiable, and is given by $P(x)=x-\nabla d(x, \partial \Omega)^{2} / 2$.) Here we are just defining it on the whole $\Omega$ via a measurable selection argument (we omit the details).

We will use the notation Id : $\bar{\Omega} \rightarrow \bar{\Omega}$ to denote the identity map on $\bar{\Omega}$.
Proposition 2.7 (Behavior of optimal plans) Let $\mu, v \in \mathcal{M}_{+}(\Omega)$, and fix $\boldsymbol{\gamma} \in \operatorname{OPT}(\mu, v)$. Then:
(i) For $\gamma_{\Omega}^{\partial \Omega}$-a.e. $(x, y)$, we have $|y-x|=d(x, \partial \Omega)$. Similarly for $\gamma_{\partial \Omega}^{\Omega}$.
(ii) If $\mu \ll \mathcal{L}^{d}$, then $\gamma_{\Omega}^{\bar{\Omega}}$ is unique, and it is given by $(\mathrm{Id}, T)_{\#} \mu$, where $T: \Omega \rightarrow \bar{\Omega}$ is the gradient of a convex function. (However, $\gamma$ as a whole may be not uniquely defined as there may be multiple ways of bringing the mass from the boundary to $v$ if no hypothesis on $v$ are made).
(iii) If $\mu, \nu \ll \mathcal{L}^{d}$, then $\gamma$ is unique.

Proof. We start with (i). Let

$$
A:=\{(x, y) \in \Omega \times \partial \Omega:|x-y|>d(x, \partial \Omega)=|x-P(x)|\}
$$

and assume by contradiction that $\gamma_{\Omega}^{\partial \Omega}(A)>0$. Then, we define

$$
\tilde{\gamma}_{\Omega}^{\partial \Omega}:=(\mathrm{Id}, P)_{\#} \pi_{\#}^{1} \gamma_{\Omega}^{\partial \Omega},
$$

and set

$$
\tilde{\gamma}:=\gamma_{\Omega}^{\Omega}+\tilde{\gamma}_{\Omega}^{\partial \Omega}+\gamma_{\partial \Omega}^{\Omega} .
$$

Since $\pi_{\#}^{1} \tilde{\gamma}_{\Omega}^{\partial \Omega}=\pi_{\#}^{1} \gamma_{\Omega}^{\partial \Omega}$ we have $\pi_{\#}^{1} \tilde{\gamma}=\pi_{\#}^{1} \gamma$. Moreover $\pi_{\#}^{2} \tilde{\gamma}_{\Omega_{\Omega}}=\pi_{\#}^{2} \gamma_{\left.\right|_{\Omega}}$ by construction, so that $\tilde{\gamma} \in \operatorname{Adm}\left(\mu_{0}, \mu_{1}\right)$. Since

$$
\begin{aligned}
\int_{\Omega \times \partial \Omega}|x-y|^{2} d \tilde{\gamma}(x, y) & <\int_{\Omega \times \partial \Omega}|x-y|^{2} d \gamma(x, y), \\
\int_{\bar{\Omega} \times \bar{\Omega} \backslash \Omega \times \partial \Omega}|x-y|^{2} d \tilde{\gamma}(x, y) & =\int_{\bar{\Omega} \times \bar{\Omega} \backslash \Omega \times \partial \Omega}|x-y|^{2} d \gamma(x, y),
\end{aligned}
$$

we have $C(\tilde{\gamma})<C(\gamma)$, which gives the desired contradiction. A symmetric argument holds for $\gamma_{\partial \Omega}^{\Omega}$.

The validity of (ii) is a direct consequence of Proposition 2.4: the support of the plan $\gamma$ is contained in a cyclically monotone set which depends only on $\mu$ and $v$. Hence, by Rockefeller's theorem is contained in the subdifferential of a convex function. Thus, as in the classical optimal transport problem with quadratic cost [3, 4, 12], for $\mu$-a.e. $x$ there is a unique $y$ such that $(x, y)$ belongs to the support of $\gamma$. This proves the uniqueness of $\gamma_{\Omega}^{\bar{\Omega}}$ and the fact that it is induced by a map, which is the gradient of a convex function.

Finally, (iii) follows from (ii).
We define the Hit time function $\mathrm{Ht}: \bar{\Omega} \times \bar{\Omega} \rightarrow[0,1]$ as

$$
\operatorname{Ht}(\mathrm{x}, \mathrm{y}):=\inf \{\mathrm{t} \in[0,1]:(1-\mathrm{t}) \mathrm{x}+\mathrm{ty} \notin \Omega\},
$$

where $\operatorname{Ht}(\mathrm{x}, \mathrm{y}):=1$ if $(1-t) x+t y \in \Omega$ for any $t \in[0,1]$. The function Ht is lower semicontinuous, and hence measurable.

Proposition 2.8 (Geodesics) The space $\left(\mathcal{M}_{+}(\Omega), W b_{2}\right)$ is a geodesic space. A curve $[0,1] \ni t \mapsto \mu_{t}$ is a minimizing geodesic with constant speed if and only if there exists $\gamma \in \operatorname{Opt}\left(\mu_{0}, \mu_{1}\right)$ such that

$$
\begin{equation*}
\mu_{t}=\left((1-t) \pi^{1}+t \pi^{2}\right)_{\#} \gamma, \quad \forall t \in(0,1) . \tag{6}
\end{equation*}
$$

Also, given a geodesic $\left(\mu_{t}\right)$, for any $t \in(0,1)$ and $s \in[0,1]$ there is a unique optimal plan $\gamma_{t}^{s}$ from $\mu_{t}$ to $\mu_{s}$, which is given by

$$
\boldsymbol{\gamma}_{t}^{s}:=\left((1-t) \pi^{1}+t \pi^{2},(1-s) \pi^{1}+s \pi^{2}\right)_{\#} \boldsymbol{\gamma}
$$



Figure 2: Geodesic interpolation is always possible in the space $\left(\mathcal{M}_{+}(\Omega), W b_{2}\right)$. Indeed, the mass can 'appear' only at $t=0$, can 'vanish' only at $t=1$, and for $t \in(0,1)$ it moves along straight segments inside $\Omega$. In particular, in the open interval $(0,1)$, a geodesic w.r.t. $W b_{2}$ is also a geodesic w.r.t. $W_{2}$.
where $\gamma \in \operatorname{Opt}\left(\mu_{0}, \mu_{1}\right)$ is the plan which induces the geodesic via Equation (6). Furthermore, the plan $\gamma_{t}^{s}$ is the unique optimal transport plan from $\mu_{t}$ to $\mu_{s}$ for the classical transport problem.

In particular, the space $\left(\mathcal{M}_{+}(\Omega), W b_{2}\right)$ is non-branching, and the mass of $\mu_{t}$ inside $\Omega$ is constant for $t \in(0,1)$.

Observe that Equation (6) does not hold for $t=0,1$, as the marginals of $\gamma$ generally charge also $\partial \Omega$. We further remark that such a statement would be false for the classical Wasserstein distance $W_{2}$. Indeed, if $\gamma$ is an optimal plan for $W_{2}$, then the curve $\mu_{t}$ defined by (6) will not in general belong to $\Omega$, unless $\Omega$ is convex.

Proof. The only new part with respect to the classical case is that, if $\gamma$ is an optimal plan from $\mu_{0}$ to $\mu_{1}$, then the measures $\mu_{t}$ defined by (6) belong to $\Omega$ (and not just to its convex hull). Once this result is proved, the rest of the proof becomes exactly the same as in the standard case of the Wasserstein distance, see [2, Paragraph 7.2]. Hence, we are going to prove only this new part.

The fact that the measures $\mu_{t}$ defined by (6) belong to $\Omega$ is equivalent to say that $\mathrm{Ht}=1 \gamma$-a.e. We argue by contradiction: assume that there exists $E \subset \bar{\Omega} \times \bar{\Omega}$ such that $\gamma(E)>0$ and $\mathrm{Ht}(\mathrm{x}, \mathrm{y})<1$ for any $(x, y) \in E$. Roughly speaking, if this was the case, for any $(x, y) \in E$, rather than moving the mass from $x$ to $y$, we could move the mass from $x$ to $P(x)$ and take the mass from $P(y)$ to $y$, reducing the transportation cost.

More rigorously, define $v_{0}:=\pi_{\#}^{1}\left(\gamma_{\left.\right|_{E}}\right), v_{1}:=\pi_{\#}^{2}\left(\gamma_{\left.\right|_{E}}\right)$, and set

$$
\begin{aligned}
\sigma_{0} & :=(\mathrm{Id}, P)_{\#} v_{0}, \\
\sigma_{1} & :=(P, \mathrm{Id})_{\#} v_{1}, \\
\tilde{\gamma} & :=\gamma_{\left.\right|_{E^{c}}}+\sigma_{0}+\sigma_{1} .
\end{aligned}
$$

Since $\pi_{\#}^{1} \sigma_{0}=v_{0}$ and $\pi_{\#}^{1} \sigma_{1}(\Omega)=0$, we have $\pi_{\#}^{1} \tilde{\gamma}_{\left.\right|_{\Omega}}=\pi_{\#}^{1} \gamma_{\left.\right|_{\Omega}}$. Similarly $\pi_{\#}^{2} \tilde{\gamma}_{\left.\right|_{\Omega}}=\pi_{\#}^{2} \gamma_{\left.\right|_{\Omega}}$,
so that $\tilde{\gamma} \in \operatorname{ADM}\left(\mu_{0}, \mu_{1}\right)$. Observe that, since $(1-\operatorname{Ht}(\mathrm{x}, \mathrm{y})) \mathrm{x}+\mathrm{Ht}(\mathrm{x}, \mathrm{y}) \mathrm{y} \in \partial \Omega$, we have

$$
\begin{aligned}
& |x-P(x)| \leq|x-(1-\operatorname{Ht}(\mathrm{x}, \mathrm{y})) \mathrm{x}+\mathrm{Ht}(\mathrm{x}, \mathrm{y}) \mathrm{y}|=\mathrm{Ht}(\mathrm{x}, \mathrm{y})|\mathrm{x}-\mathrm{y}|, \\
& |y-P(y)| \leq|y-(1-\operatorname{Ht}(\mathrm{x}, \mathrm{y})) \mathrm{x}+\operatorname{Ht}(\mathrm{x}, \mathrm{y}) \mathrm{y}|=(1-\mathrm{Ht}(\mathrm{x}, \mathrm{y}))|\mathrm{x}-\mathrm{y}|,
\end{aligned}
$$

and therefore

$$
\begin{aligned}
C(\tilde{\gamma}) & =\int_{E^{c}}|x-y|^{2} d \tilde{\gamma}+\int_{E}|x-P(x)|^{2}+|y-P(y)|^{2} d \gamma \\
& \leq \int_{E^{c}}|x-y|^{2} d \gamma+\int_{E}|x-y|^{2}\left(\operatorname{Ht}(\mathrm{x}, \mathrm{y})^{2}+(1-\operatorname{Ht}(\mathrm{x}, \mathrm{y}))^{2}\right) d \gamma \\
& <\int_{E^{c}}|x-y|^{2} d \gamma+\int_{E}|x-y|^{2} d \gamma=C(\gamma)
\end{aligned}
$$

which contradicts the optimality of $\gamma$. $\operatorname{Thus} \operatorname{Ht}(\mathrm{x}, \mathrm{y})=1$ for $\gamma$-a.e. $(x, y)$, which implies that the measures $\mu_{t}$ are concentrated in $\Omega$, as desired.

Let $\mu, v \in \mathcal{M}_{+}(\Omega)$ and assume that $\mu(\Omega)=v(\Omega)>0$. Then any plan $\gamma$ which is optimal for the classical transportation cost is admissible for the new one. Therefore we have the inequality:

$$
\begin{equation*}
W b_{2}(\mu, v) \leq W_{2}(\mu, v), \quad \forall \mu, v \in \mathcal{M}_{+}(\Omega) \text { s.t. } \mu(\Omega)=v(\Omega)>0 . \tag{7}
\end{equation*}
$$



Figure 3: For measures with the same amount of mass, the distance $W b_{2}$ is smaller than the classical $W_{2}$ : as the picture shows, it may be much better to exchange the mass with the boundary rather than internally.

Proposition 2.9 (An estimate on the directional derivative) Let $\mu, v \in \mathcal{M}_{+}(\Omega)$ and $\boldsymbol{w}: \Omega \rightarrow \mathbb{R}^{d}$ a bounded vector field with compact support. Also, let $\gamma \in \operatorname{Opt}(\mu, v)$, and define $\mu_{t}:=(\mathrm{Id}+t \boldsymbol{w})_{\#} \mu$. Then

$$
\limsup _{t \rightarrow 0} \frac{W b_{2}^{2}\left(\mu_{t}, v\right)-W b_{2}^{2}(\mu, v)}{t} \leq-2 \int\langle\boldsymbol{w}(x), y-x\rangle d \gamma(x, y) .
$$

Proof. Observe that since $\boldsymbol{w}$ is compactly supported in $\Omega$, for $t>0$ sufficiently small $\mu_{t} \in \mathcal{M}_{+}(\Omega)$, so that the statement makes sense. Now it is simple to check that the plan $\gamma_{t}$ defined by

$$
\gamma_{t}:=\left((\operatorname{Id}+t w) \circ \pi^{1}, \pi^{2}\right)_{\#} \gamma,
$$

belongs to $\operatorname{Adm}\left(\mu_{t}, v\right)$. Hence

$$
\begin{aligned}
W b_{2}^{2}\left(\mu_{t}, v\right) & \leq \int|x-y|^{2} d \gamma_{t}(x, y)=\int|x+t \boldsymbol{w}(x)-y|^{2} d \boldsymbol{\gamma}(x, y) \\
& =W b_{2}^{2}(\mu, v)-2 t \int\langle\boldsymbol{w}(x), y-x\rangle d \boldsymbol{\gamma}(x, y)+t^{2} \int|\boldsymbol{w}(x)|^{2} d \boldsymbol{\gamma}(x, y)
\end{aligned}
$$

and the conclusion follows.

## 3 The heat equation with Dirichlet boundary condition as a 'gradient flow'

This section contains an application of our new transportation distance: we are going to show that the gradient flow of the entropy functional $\int_{\Omega}[\rho \log (\rho)-\rho] d x$ coincides with the heat equation, with Dirichlet boundary condition equal to 1 . To prove such a result, we will first study some of the properties of the entropy, showing in particular a lower bound on its slope, see Proposition 3.2. Then, following the strategy introduced in [8], we will apply the minimizing movement scheme to prove our result. Finally we will show that the discrete solutions constructed by minimizing movements enjoy a comparison principle: if $\left(\rho_{k}^{\tau}\right)_{k \in \mathbb{N}}$ and $\left(\tilde{\rho}_{k}^{\tau}\right)_{k \in \mathbb{N}}$ are two discrete solution for a time step $\tau>0$, and $\rho_{0}^{\tau} \leq \tilde{\rho}_{0}^{\tau}$, then $\rho_{k}^{\tau} \leq \tilde{\rho}_{k}^{\tau}$ for all $k \in \mathbb{N}$. Letting $\tau \rightarrow 0$, this monotonicity result allows to recover the classical maximum principle for the heat equation.

To be clear: we will not state any result concerning existence of the gradient flow of the entropy (we will not identify the slope of the entropy, nor the infinitesimal description of the distance $W b_{2}$ ). What we will do is a work 'by hands': we will show that we have compactness in the minimizing movements scheme and prove that any limit is a weak solution of the heat equation with Dirichlet boundary conditions.

### 3.1 The entropy

The entropy functional $E: \mathcal{M}_{+}(\Omega) \rightarrow \mathbb{R} \cup\{+\infty\}$ is defined as

$$
E(\mu):= \begin{cases}\int_{\Omega} e(\rho(x)) d x & \text { if } \mu=\left.\rho \mathcal{L}^{d}\right|_{\Omega} \\ +\infty & \text { otherwise }\end{cases}
$$

where $e:[0,+\infty) \rightarrow[0,+\infty)$ is given by

$$
e(z):=z \log (z)-z+1 .
$$

From now on, since we will often deal with absolutely continuous measures, and by abuse of notation we will sometimes use $\rho$ to denote the measure $\rho \mathcal{L}^{d}{ }_{\Omega \Omega}$. In particular, we will write $\operatorname{Adm}\left(\rho, \rho^{\prime}\right)$ in place of $\operatorname{Adm}\left(\rho \mathcal{L}^{d}{ }_{{ }_{\Omega}}, \rho^{\prime} \mathcal{L}^{d}{ }_{\left.\right|_{\Omega}}\right)$.

Proposition 3.1 (Semicontinuity and compactness of sublevels) The functional $E$ : $\mathcal{M}_{+}(\Omega) \rightarrow \mathbb{R} \cup\{+\infty\}$ takes value in $[0,+\infty]$, it is lower semicontinuous with respect to $W b_{2}$, and its sublevels are compact.

Proof. If $\mu=\left.\rho \mathcal{L}^{d}\right|_{\Omega}$, thanks to Jensen inequality we have

$$
\begin{equation*}
e\left(\frac{\mu(\Omega)}{|\Omega|}\right)=e\left(\frac{1}{|\Omega|} \int_{\Omega} \rho d x\right) \leq \frac{1}{|\Omega|} \int_{\Omega} e(\rho) d x=\frac{E(\mu)}{|\Omega|} . \tag{8}
\end{equation*}
$$

This inequality bounds the mass of $\rho$ in terms of the entropy, which gives the relative compactness of the sublevels of $E$. The bound $E(\mu) \geq 0$ is immediate as $e \geq 0$. Finally, the lower semicontinuity follows from the convexity and superlinearity of $e$ and from the equivalence between weak convergence and convergence w.r.t. $W b_{2}$.

We recall that the slope of the functional $E$ defined on the metric space $\left(\mathcal{M}_{+}(\Omega), W b_{2}\right)$ is defined as:

$$
|\nabla E|(\mu):=\limsup _{v \rightarrow \mu} \frac{(E(\mu)-E(v))^{+}}{W b_{2}(\mu, v)}
$$

Proposition 3.2 (Bound of the slope in terms of Fisher's information) The slope of $E$ is bounded from below by the square root of the Fisher information $F: \mathcal{M}_{+}(\Omega) \rightarrow$ $[0,+\infty]$ :

$$
F(\mu):= \begin{cases}4 \int_{\Omega}|\nabla \sqrt{\rho}|^{2} d x & \text { if } \mu=\left.\rho \mathcal{L}^{d}\right|_{\Omega} \text { and } \sqrt{\rho} \in H^{1}(\Omega), \\ +\infty & \text { otherwise. }\end{cases}
$$

Proof. Take $\mu \in \mathcal{M}_{+}(\Omega)$, define $m:=\mu(\Omega)$, and let $\mathcal{M}_{m}(\Omega)$ be the set of non-negative measures on $\Omega$ with mass $m$. On $\mathcal{M}_{m}(\Omega)$, we can consider the Wasserstein distance $W_{2}$. Consider the functional $E:\left(\mathcal{M}_{m}(\Omega), W_{2}\right) \rightarrow \mathbb{R} \cup\{+\infty\}$. It is well-known that $|\nabla E|(\mu)=\sqrt{F(\mu)}$ for all $\mu \in \mathcal{M}_{1}$, see [2, Chapter 10]. Then, it is easily checked by a scaling argument that the formula remains true for arbitrary $m \geq 0$. Hence, taking into account inequality (7), we obtain

$$
|\nabla E|(\mu) \geq \limsup _{\mathcal{M}_{m}(\Omega) \nexists v \rightarrow \mu} \frac{(E(\mu)-E(v))^{+}}{W b_{2}(\mu, v)} \geq \lim _{\mathcal{M}_{m}(\Omega) \ngtr v \rightarrow \mu} \frac{(E(\mu)-E(v))^{+}}{W_{2}(\mu, v)}=\sqrt{F(\mu)},
$$

as desired.

Proposition 3.3 (A directional derivative of $E$ ) Let $\mu=\rho \mathcal{L}^{d} \in \mathcal{M}_{+}(\Omega)$ be such that $E(\mu)<+\infty$, and let $\boldsymbol{w}: \Omega \rightarrow \mathbb{R}^{d}$ be a $C^{\infty}$ vector field with compact support. Define $\mu_{t}:=(\operatorname{Id}+t \boldsymbol{w})_{\#} \mu$. Then

$$
\lim _{t \rightarrow 0} \frac{E\left(\mu_{t}\right)-E(\mu)}{t}=\int_{\Omega} \rho \operatorname{div} \boldsymbol{w} d x
$$

Proof. Since $\boldsymbol{w}$ is compactly supported, $\mu_{t} \in \mathcal{M}_{+}(\Omega)$ for sufficiently small $t$, and the proof is exactly the same as the one in the Wasserstein case.

Remark 3.4 [A source of difficulties] It is important to underline that the entropy $E$ is not geodesically convex on the space $\left(\mathcal{M}_{+}(\Omega), W b_{2}\right)$. Indeed, since for instance the mass can disappear at the boundary for $t=1$, it is possible that an high concentration of mass near $\partial \Omega$ gives $\lim _{t \uparrow 1} E\left(\mu_{t}\right)=+\infty$, while $E\left(\mu_{1}\right)<+\infty$. (Observe that, once the mass has reached $\partial \Omega$, it does not contribute any more to the energy!) Still, since for $t, s \in(0,1)$ the optimal transport plan for $W b_{2}$ coincides with the optimal transport plan for $W_{2}$ (Proposition 2.8), $t \mapsto E\left(\mu_{t}\right)$ is convex in the open interval $(0,1)$ (see [2, Chapter 9]).

### 3.2 Minimizing movements for the entropy

In this paragraph we apply the minimizing movements to construct a weak solution to the heat equation with Dirichlet boundary condition.

We briefly review the minimizing movement scheme, referring to [2] for a detailed description and general results. Fix $\rho_{0} \in \mathcal{M}_{+}(\Omega)$ such that $E\left(\rho_{0}\right)<+\infty$ (given the


Figure 4: For typical $\mu_{0}, \mu_{1}$, a geodesic connecting them takes mass from the boundary at $t=0$ and leaves mass at $t=1$. In this case the graph of $t \mapsto E\left(\mu_{t}\right)$ looks like in the picture: in the interval $(0,1)$ the function is convex and converges to $+\infty$ as $t \rightarrow 0,1$. The value of $E\left(\mu_{0}\right)$ and $E\left(\mu_{1}\right)$ has basically no connection with the values in intermediate times.
lack of convexity of $E$, we need to assume that the entropy at the initial point is finite, thus in particular the measure is absolutely continuous), and fix a time step $\tau>0$. Set $\rho_{0}^{\tau}:=\rho_{0}$, and define recursively $\rho_{n+1}^{\tau}$ as the unique minimizer of

$$
\mu \quad \mapsto \quad E(\mu)+\frac{W b_{2}^{2}\left(\mu, \rho_{n}^{\tau}\right)}{2 \tau}
$$

(see Proposition 3.6 below). Then, we define the discrete solution $t \mapsto \rho^{\tau}(t) \in \mathcal{M}_{+}(\Omega)$ by:

$$
\rho^{\tau}(t):=\rho_{n}^{\tau} \quad \text { for } t \in[n \tau,(n+1) \tau) .
$$

We recall that the space $W_{0}^{1,1}(\Omega)$ is defined as the closure of $C_{0}^{\infty}(\Omega)$ w.r.t. the $W^{1,1}-$ norm. (Observe that this definition requires no smoothness assumptions on $\partial \Omega$.) Then we say that $f \in W^{1,1}(\Omega)$ has trace 1 if $f-1 \in W_{0}^{1,1}(\Omega)$. (More in general, given a smooth function $\phi: \bar{\Omega} \rightarrow \mathbb{R}$, one may say that $f \in W^{1,1}(\Omega)$ has trace $\phi$ if $f-\phi \in W_{0}^{1,1}(\Omega)$.)

Our main theorem is the following:
Theorem 3.5 With the above notation, for any sequence $\tau_{k} \downarrow 0$ there exists a subsequence, not relabelled, such that, for any $t \geq 0, \rho^{\tau_{k}}(t)$ converges to some limit measure $\rho(t)$ in $\left(\mathcal{M}_{+}(\Omega), W b_{2}\right)$ as $k \rightarrow \infty$. The map $t \mapsto(\rho(t)-1)$ belongs to $L_{l o c}^{2}\left([0,+\infty), W_{0}^{1,1}(\Omega)\right)$, and $t \mapsto \rho(t)$ is a weak solution of the heat equation

$$
\left\{\begin{align*}
\frac{d}{d t} \rho(t) & =\Delta \rho(t),  \tag{9}\\
\rho(0) & =\rho_{0} .
\end{align*}\right.
$$

We recall that a weakly continuous curve of measure $t \mapsto \mu_{t} \in \mathcal{M}_{+}(\Omega)$ is said to be a weak solution of (9) if

$$
\int_{\Omega} \varphi d \mu_{s}(x)-\int_{\Omega} \varphi d \mu_{t}(x)=\int_{t}^{s}\left(\int_{\Omega} \Delta \varphi d \mu_{r}(x)\right) d r, \quad \forall 0 \leq t<s, \forall \varphi \in C_{c}^{\infty}(\Omega)
$$

In order to prove this theorem, we need the following lemma, which describes the behavior of a single step of the minimizing movements scheme.

Proposition 3.6 (A step of the minimizing movement) Let $\mu \in \mathcal{M}_{+}(\Omega)$ and $\tau>0$. Then there exists a unique minimum $\mu_{\tau} \in \mathcal{M}_{+}(\Omega)$ of

$$
\begin{equation*}
\sigma \mapsto E(\sigma)+\frac{W b_{2}^{2}(\mu, \sigma)}{2 \tau} \tag{10}
\end{equation*}
$$

Such a minimum satisfies:
(i) $\mu_{\tau}=\left.\rho_{\tau} \mathcal{L}^{d}\right|_{\Omega}$, with $\rho_{\tau}-1 \in W_{0}^{1,1}(\Omega)$.
(ii) The restriction to $\Omega \times \bar{\Omega}$ of any optimal transport plan from $\mu_{\tau}$ to $\mu$ is induced by a map $T$, which satisfies

$$
\begin{equation*}
\frac{T(x)-x}{\tau} \rho_{\tau}(x)=-\nabla \rho_{\tau}(x), \quad \mathcal{L}^{d}-\text { a.e. } x \tag{11}
\end{equation*}
$$

Proof. The existence of a minimum follows by a standard compactness-semicontinuity argument, while the uniqueness is a direct consequence of the convexity of $W b_{2}^{2}(\cdot, \mu)$ w.r.t. usual linear interpolation of measures and the strict convexity of $E(\cdot)$.

It is well known that at minimum of (10) the slope is finite (see [2, Lemma 3.1.3]). Hence $\sqrt{\rho_{\tau}} \in H^{1}(\Omega)$ by Proposition 3.2. Hence, thanks to Hölder inequality, $\rho_{\tau} \in$ $W^{1,1}(\Omega)$. Moreover, thanks to (21) below, we have

$$
e^{-d(x, \partial \Omega)^{2} /(2 \tau)} \leq \rho_{\tau}(x) \leq e^{d(x, \partial \Omega)^{2} /(2 \tau)} \quad \forall x \in \Omega,
$$

which easily implies that $\rho_{\tau}$ has trace 1 on $\partial \Omega$. This shows (i).
To prove (ii), we start by observing that Proposition 2.7 and the absolute continuity of $\mu_{\tau}$ guarantees the existence of $T$. Now, choose a $C^{\infty}$ vector field $\boldsymbol{w}$ with compact support in $\Omega$, and define $\rho_{\tau}^{t}:=(\operatorname{Id}+t \boldsymbol{w})_{\#} \rho_{\tau}$. Using the minimality of $\rho$, we get

$$
E\left(\rho_{\tau}^{t}\right)-E\left(\rho_{\tau}\right)+\frac{W b_{2}^{2}\left(\rho_{\tau}^{t}, \mu\right)-W b_{2}^{2}\left(\rho_{\tau}, \mu\right)}{2 \tau} \geq 0 .
$$

Dividing by $t$, and letting $t \downarrow 0$, thanks to Propositions 3.3 and 2.9 we get

$$
\int_{\Omega} \rho \operatorname{div} \boldsymbol{w} d x-\int\left\langle\boldsymbol{w}, \frac{T-\mathrm{Id}}{\tau}\right\rangle \rho d x \geq 0 .
$$

Exchanging $\boldsymbol{w}$ with $-\boldsymbol{w}$ and exploiting the arbitrariness of $\boldsymbol{w}$, the thesis follows.
To prove theorem 3.5 we will use the following a priori bound for the discrete solution, see [2, Lemma 3.2.2 and Equation (3.2.3)]:

$$
\begin{equation*}
\frac{1}{2} \sum_{i=n}^{m-1} \frac{W b_{2}^{2}\left(\rho_{\rho}^{\tau}, \rho_{i+1}^{\tau}\right)}{\tau}+\frac{\tau}{2} \sum_{i=n}^{m-1}|\nabla E|^{2}\left(\rho_{i}^{\tau}\right) \leq E\left(\rho_{m}^{\tau}\right)-E\left(\rho_{n}^{\tau}\right) \quad \forall n \leq m \in \mathbb{N} . \tag{12}
\end{equation*}
$$

Proof of Theorem 3.5. - Compactness argument. Let $\left\{\tau_{k}\right\}_{k \in \mathbb{N}}$ be a sequence converging to 0 . First of all we observe that, thanks to (8) and the inequality $E\left(\rho^{\tau_{k}}(t)\right) \leq E\left(\rho_{0}\right)$, the mass of the measures $\rho^{\tau_{k}}(t)$ is uniformly bounded for all $k \in \mathbb{N}, t \geq 0$. Then a standard diagonal argument shows that there exists a subsequence, not relabelled, such
that $\rho^{\tau_{k}}(t)$ converges to some $\rho(t)$ in $\left(\mathcal{M}_{+}(\Omega), W b_{2}\right)$ for any $t \in \mathbb{Q}_{+}$. Now, thanks to the uniform bound on the discrete speed

$$
\frac{1}{2} \sum_{i=n}^{m-1} \frac{W b_{2}^{2}\left(\rho_{i}^{\tau_{k}}, \rho_{i+1}^{\tau_{k}}\right)}{\tau} \leq E\left(\rho_{m}^{\tau}\right)-E\left(\rho_{n}^{\tau}\right) \leq E\left(\rho_{0}\right),
$$

(which is a direct consequence of (12)), we easily get

$$
\begin{equation*}
W b_{2}\left(\rho^{\tau_{k}}(t), \rho^{\tau_{k}}(s)\right) \leq \sqrt{2 E\left(\rho_{0}\right)\left[t-s+\tau_{k}\right]} \quad \forall 0 \leq s \leq t, \tag{13}
\end{equation*}
$$

which implies the convergence of $\rho^{\tau_{k}}(t)$ for every $t \geq 0$.

- Any limit point is a weak solution of the heat equation. Let $\tau_{k} \downarrow 0$ be a sequence such that $\rho^{\tau_{k}}(t)$ converges to some $\rho(t)$ in $\left(\mathcal{M}_{+}(\Omega), W b_{2}\right)$ for any $t \geq 0$. We want to prove that $t \mapsto \rho(t)$ is a weak solution of the heat equation. For any $\tau>0, n \in \mathbb{N}$, let $T_{n}^{\tau}$ be the map which induces $\left(\gamma_{n}^{\tau}\right)_{\Omega}^{\bar{\Omega}}$, where $\boldsymbol{\gamma}_{n}^{\tau} \in \operatorname{Opt}\left(\rho_{n+1}^{\tau}, \rho_{n}^{\tau}\right)$ (Proposition 2.7(ii)). Fix $\varphi \in C_{c}^{\infty}(\Omega)$ and observe that

$$
\begin{align*}
\int_{\Omega} \varphi \rho_{n+1}^{\tau} d x-\int_{\Omega}\left(\varphi \circ T_{n}^{\tau}\right) \rho_{n+1}^{\tau} d x & =\int_{\Omega}\left(\int_{0}^{1}\left\langle\nabla \varphi \circ\left((1-\lambda) T_{n}^{\tau}+\lambda \mathrm{Id}\right), \mathrm{Id}-T_{n}^{\tau}\right\rangle d \lambda\right) \rho_{n+1}^{\tau} d x \\
& =-\int_{\Omega}\left\langle\nabla \varphi, T_{n}^{\tau}-\mathrm{Id}\right\rangle \rho_{n+1}^{\tau} d x+R(\tau, n) \\
& =\tau \int_{\Omega}\left\langle\nabla \varphi, \nabla \rho_{n+1}^{\tau}\right\rangle d x+R(\tau, n) \\
& =-\tau \int_{\Omega} \Delta \varphi \rho_{n+1}^{\tau} d x+R(\tau, n) \tag{14}
\end{align*}
$$

where at the third step we used (11), and the reminder term $R(\tau, n)$ is bounded by

$$
\begin{equation*}
|R(\tau, n)| \leq(\operatorname{Lip} \nabla \varphi) \int_{\Omega}\left|T_{n}^{\tau}-\operatorname{Id}\right|^{2} \rho_{n+1}^{\tau} d x=\operatorname{Lip}(\nabla \varphi) W b_{2}^{2}\left(\rho_{n}^{\tau}, \rho_{n+1}^{\tau}\right) \tag{15}
\end{equation*}
$$

Now, since the support of $\varphi$ is included in $\Omega$ and $\left.\left(\left(T_{n}^{\tau}\right) \# \rho_{n+1}^{\tau}\right)\right|_{\Omega}=\pi_{\#}^{2}\left(\left(\gamma_{n}^{\tau}\right)_{\Omega}^{\Omega}\right)$, we have

$$
\int_{\Omega} \varphi \rho_{n}^{\tau} d x-\int_{\Omega}\left(\varphi \circ T_{n}^{\tau}\right) \rho_{n+1}^{\tau} d x=\int_{\Omega \times \Omega} \varphi(y) d\left(\gamma_{n}^{\tau}\right)_{\partial \Omega}^{\Omega}(x, y) .
$$

By Proposition 2.7 we have $|x-y|=d(y, \partial \Omega)$ for $\left(\gamma_{n}^{\tau}\right)_{\partial \Omega}^{\Omega}$-a.e. $(x, y)$, which implies

$$
\begin{aligned}
W b_{2}^{2}\left(\rho_{n+1}^{\tau}, \rho_{n}^{\tau}\right) \geq & \int_{\Omega \times \operatorname{supp}(\varphi)}|x-y|^{2} d\left(\boldsymbol{\gamma}_{n}^{\tau}\right)_{\partial \Omega}^{\Omega}(x, y) \\
& =\int_{\Omega \times \operatorname{supp}(\varphi)} d(y, \partial \Omega)^{2} d\left(\boldsymbol{\gamma}_{n}^{\tau}\right)_{\partial \Omega}^{\Omega}(x, y) \geq c_{\varphi} \int_{\Omega \times \operatorname{supp}(\varphi)} d\left(\boldsymbol{\gamma}_{n}^{\tau}\right)_{\partial \Omega}^{\Omega}(x, y),
\end{aligned}
$$

where $c_{\varphi}:=\min _{y \in \operatorname{supp}(\varphi)} d(y, \partial \Omega)^{2}>0$. Hence

$$
\left|\int_{\Omega} \varphi \rho_{n}^{\tau} d x-\int_{\Omega}\left(\varphi \circ T_{n}^{\tau}\right) \rho_{n+1}^{\tau} d x\right| \leq \frac{\|\varphi\|_{\infty}}{c_{\varphi}} W b_{2}^{2}\left(\rho_{n+1}^{\tau}, \rho_{n}^{\tau}\right)
$$

Combining the above estimate with (14) and (15), we obtain

$$
\begin{equation*}
\int_{\Omega} \varphi \rho_{n+1}^{\tau} d x-\int_{\Omega} \varphi \rho_{n}^{\tau} d x=-\tau \int_{\Omega} \Delta \varphi \rho_{n+1}^{\tau} d x+\tilde{R}(\tau, n), \tag{16}
\end{equation*}
$$

where

$$
\begin{equation*}
|\tilde{R}(\tau, n)| \leq\left(\operatorname{Lip}(\nabla \varphi)+\frac{\|\varphi\|_{\infty}}{c_{\varphi}}\right) W b_{2}^{2}\left(\rho_{n}^{\tau}, \rho_{n+1}^{\tau}\right) . \tag{17}
\end{equation*}
$$

Now, choose $0 \leq t<s$, let $\tau=\tau_{k}$, and add up Equation (16) from $n=\left[t / \tau_{k}\right]$ to $m=\left[s / \tau_{k}\right]-1$ to get
$\int_{\Omega} \varphi \rho^{\tau_{k}}(s) d x-\int_{\Omega} \varphi \rho^{\tau_{k}}(t) d x=\int_{\tau_{k}\left[t / \tau_{k}\right]}^{\tau_{k}\left[s / \tau_{k}\right]}\left(\int_{\Omega} \Delta \varphi \rho^{\tau_{k}}(r) d x\right) d x d r+\sum_{n=\left[t / \tau_{k}\right]}^{\left[s / \tau_{k}\right]-1} \tilde{R}\left(\tau_{k},\left[r / \tau_{k}\right]\right)$.
We want to take the limit in the above equation as $\tau_{k} \downarrow 0$. The $W b_{2}$-convergence of $\rho^{\tau_{k}}(r)$ to $\rho(r)$, combined with Proposition 2.2, gives that the left hand side converges to $\int_{\Omega} \varphi \rho(s) d x-\int_{\Omega} \varphi \rho(t) d x$. For the same reason, $\int_{\Omega} \Delta \varphi \rho^{\tau_{k}}(r) d x \rightarrow \int_{\Omega} \Delta \varphi \rho(r) d x$ for any $r \geq 0$. Thus, since the mass of the measures $\rho^{\tau_{k}}(t)$ is uniformly bounded, we get

$$
\int_{\Omega}\left|\Delta \varphi \rho^{\tau_{k}}(r)\right| d x \leq\|\Delta \varphi\|_{\infty} \int_{\Omega} \rho^{\tau_{k}}(r) d x \leq C_{0}
$$

for some positive constant $C_{0}$, so that by the dominated convergence theorem we get

$$
\int_{\tau_{k}\left[t / \tau_{k}\right]}^{\tau_{k}\left[s / \tau_{k}\right]}\left(\int_{\Omega} \Delta \varphi \rho^{\tau_{k}}(r) d x\right) d r \rightarrow \int_{t}^{s}\left(\int_{\Omega} \Delta \varphi \rho(r) d x\right) d r
$$

as $\tau_{k} \downarrow 0$. Finally, thanks to (12) and (17), the reminder term is bounded by

$$
\begin{aligned}
\left|\sum_{n=\left[t / \tau_{k}\right]}^{\left[s / \tau_{k}\right]-1} \tilde{R}\left(\tau_{k},\left[r / \tau_{k}\right]\right)\right| & \leq\left(\operatorname{Lip}(\nabla \varphi)+\frac{\|\varphi\|_{\infty}}{c_{\varphi}}\right)_{n=\left[t / \tau_{k}\right]}^{\left[s / \tau_{k}\right]-1} W b_{2}^{2}\left(\rho_{n}^{\tau_{k}}, \rho_{n+1}^{\tau_{k}}\right) \\
& \leq 2 \tau_{k}\left(\operatorname{Lip}(\nabla \varphi)+\frac{\|\varphi\|_{\infty}}{c_{\varphi}}\right) E\left(\rho_{0}\right)
\end{aligned}
$$

and thus it goes to 0 as $\tau_{k} \downarrow 0$. In conclusion, we proved that

$$
\int_{\Omega} \varphi \rho(s) d x-\int_{\Omega} \varphi \rho(t) d x=\int_{t}^{s}\left(\int_{\Omega} \Delta \varphi \rho(r) d x\right) d r, \quad \forall 0 \leq t<s, \forall \varphi \in C_{c}^{\infty}(\Omega)
$$

Thanks to Equation (13) it is immediate to check that the curve $\left.t \mapsto \rho(t) \mathcal{L}^{d}\right|_{\Omega} \in \mathcal{M}_{+}(\Omega)$ is continuous w.r.t. $W b_{2}$, and therefore weakly continuous. Finally, since $\rho^{\tau}(0)=\rho_{0}$ for any $\tau>0, \rho(0)=\rho_{0}$ and the initial condition is satisfied.

- The curve $t \mapsto(\rho(t)-1)$ belongs to $L_{l o c}^{2}\left([0,+\infty), W_{0}^{1,1}(\Omega)\right)$. From inequality (12), Proposition 3.2, and the fact that the mass of the measures $\rho^{\tau}(t)$ is uniformly bounded, we know that

$$
\int_{0}^{\infty}\left(\int_{\Omega}\left|\nabla \sqrt{\rho^{\tau_{k}}(t)}\right|^{2} d x\right) d t \leq \frac{1}{4} \int_{0}^{\infty}|\nabla E|^{2}\left(\rho^{\tau_{k}}(t)\right) d t \leq C_{0} E\left(\rho_{0}\right)
$$

which means that the functions $t \mapsto \sqrt{\rho^{\tau_{k}}(t)}$ are equibounded in $L_{l o c}^{2}\left([0,+\infty), H_{0}^{1}(\Omega)\right)$, which implies that $t \mapsto \sqrt{\rho(t)}$ belongs to $L_{l o c}^{2}\left([0,+\infty), H^{1}(\Omega)\right)$, so that by Hölder $t \mapsto$ $\rho(t) \in L_{l o c}^{2}\left([0,+\infty), W^{1,1}(\Omega)\right)$. Moreover, thanks to Fatou lemma,

$$
\int_{0}^{\infty} \liminf _{k \rightarrow+\infty}\left(\int_{\Omega}\left|\nabla \sqrt{\rho^{\tau_{k}}(t)}\right|^{2} d x\right) d t<+\infty
$$

which gives

$$
\liminf _{k \rightarrow+\infty} \int_{\Omega}\left|\nabla \sqrt{\rho^{\tau_{k}}(t)}\right|^{2} d x<+\infty \quad \text { for a.e. } t \geq 0,
$$

so that by Hölder inequality

$$
\liminf _{k \rightarrow+\infty} \int_{\Omega}\left|\nabla \rho^{\tau_{k}}(t)\right| d x<+\infty \quad \text { for a.e. } t \geq 0
$$

Now, for any $t$ such that the above liminf is finite, consider a subsequence $k_{n}$ (depending on $t$ ) such that

$$
\sup _{n \in \mathbb{N}} \int_{\Omega}\left|\nabla \rho^{\tau_{k_{n}}}(t)\right| d x<+\infty
$$

Then, recalling that $\rho^{\tau_{k}}(t) \rightarrow \rho(t)$ in $\left(\mathcal{M}_{+}(\Omega), W b_{2}\right)$, since $\rho^{\tau_{k_{n}}}(t)$ is uniformly bounded in $W^{1,1}(\Omega)$ and belong to $W_{0}^{1,1}(\Omega)$ by Proposition 3.6(i) we easily get that $\rho^{\tau_{k}}(t) \rightarrow \rho(t)$ weakly in $W^{1,1}(\Omega)$, and $\rho(t)-1 \in W_{0}^{1,1}(\Omega)$ as desired.

### 3.3 A comparison principle

In this section we prove the following monotonicity result for the minimizing movement scheme of $E$ w.r.t. $W b_{2}$ : if we have two measures $\mu, \tilde{\mu}$ satisfying $\mu \geq \tilde{\mu}$, then $\mu_{\tau} \geq \tilde{\mu}_{\tau}$ for every $\tau \geq 0$, where $\mu_{\tau}, \tilde{\mu}_{\tau}$ are the unique minimizers of (10) for $\mu$ and $\tilde{\mu}$ respectively. It is interesting to underline that:

- Once monotonicity for the single time step is proven, a maximum principle for weak solutions of heat equation can be proved as a direct consequence, see Corollary 3.9.
- Although our strategy is not new (for instance, it has been used in the context of the classical transportation problem in $[9,1]$ to prove a maximum principle), the fact of having no mass constraints makes it more efficient, and the properties of minimizers that we are able to deduce are in some sense stronger.
- The argument that we are going to use holds in much more general situations, see Remark 3.10. (This in not the case when one deals with the classical transportation problem, where the fact that the cost function satisfies $c(x, x) \leq c(x, y)$ for all $x, y \in \Omega$ plays an important role, see [1, 7].)

The proof of the monotonicity relies on a set of inequalities valid for each minimizer of (10). In the next proposition we are going to assume that $\mu=\left.\rho \mathcal{L}^{d}\right|_{\Omega} \in \mathcal{M}_{+}(\Omega)$ is an absolutely continuous measure and that $\tau>0$ is a fixed time step. Also, we will denote by $\mu_{\tau}=\left.\rho_{\tau} \mathcal{L}^{d}\right|_{\Omega}$ the unique minimizer of (10) (which is absolutely continuous by Proposition 3.6), by $\gamma$ the unique optimal plan for $\left(\rho, \rho_{\tau}\right)$, by $T$ the map which induces $\gamma_{\Omega}^{\bar{\Omega}}$, and by $S$ the map which induces $\gamma_{\bar{\Omega}}^{\Omega}$ seen from $\rho_{\tau}$ (see Proposition 2.7).

Proposition 3.7 With the notation above, the following inequalities hold:

- Let $y_{1}, y_{2} \in \Omega$ be Lebesgue points for $\rho_{\tau}$, and assume that $y_{1}$ is also a Lebesgue point for $S$. Then

$$
\begin{equation*}
\log \left(\rho_{\tau}\left(y_{1}\right)\right)+\frac{\left|y_{1}-S\left(y_{1}\right)\right|^{2}}{2 \tau} \leq \log \left(\rho_{\tau}\left(y_{2}\right)\right)+\frac{\left|y_{2}-S\left(y_{1}\right)\right|^{2}}{2 \tau} . \tag{18}
\end{equation*}
$$

- Let $x \in \Omega$ be a Lebesgue point for both $\rho$ and $T$, and assume that $T(x) \in \partial \Omega$. Assume further that $y \in \Omega$ is a Lebesgue point for $\rho_{\tau}$. Then

$$
\begin{equation*}
\frac{|x-T(x)|^{2}}{2 \tau} \leq \log \left(\rho_{\tau}(y)\right)+\frac{|x-y|^{2}}{2 \tau} \tag{19}
\end{equation*}
$$

- Let $y_{1} \in \Omega$ be a Lebesgue point for $\rho_{\tau}$. Then, for any $y_{2} \in \partial \Omega$, we have

$$
\begin{equation*}
\log \left(\rho_{\tau}\left(y_{1}\right)\right)+\frac{\left|y_{1}-S\left(y_{1}\right)\right|^{2}}{2 \tau} \leq \frac{\left|y_{2}-S\left(y_{1}\right)\right|^{2}}{2 \tau} \tag{20}
\end{equation*}
$$

- Let $y \in \Omega$ be a Lebesgue point for $\rho_{\tau}$. Then

$$
\begin{equation*}
\frac{d^{2}(y, \partial \Omega)}{2 \tau} \geq\left|\log \left(\rho_{\tau}(y)\right)\right| \tag{21}
\end{equation*}
$$

- Let $y \in \Omega$ be a Lebesgue point for both $\rho_{\tau}$ and $S$, and assume that $S(y) \in \partial \Omega$. Then

$$
\begin{equation*}
\log \left(\rho_{\tau}(y)\right)+\frac{d^{2}(y, \partial \Omega)}{2 \tau}=0 \tag{22}
\end{equation*}
$$

Proof. - Heuristic arguments. We start with (18). Consider a point $y_{1} \in \Omega$, and observe that the mass $\rho_{\tau}\left(y_{1}\right)$ comes from $S\left(y_{1}\right)$. (It does not matter whether $S\left(y_{1}\right) \in \Omega$ or $S\left(y_{1}\right) \in \partial \Omega$ ) We now make a small perturbation of $\rho_{1}$ in the following way: we pick a small amount of mass from $S\left(y_{1}\right)$ and, instead than moving it to $y_{1}$, we move it to $y_{2}$. In terms of entropy, we are earning $\log \left(\rho_{1}\left(S\left(y_{1}\right)\right)\right)$ because of the less mass in $S\left(y_{1}\right)$ and paying $\log \left(\rho_{1}\left(y_{2}\right)\right)$ because of the greater amount of mass at $y_{2}$. In terms of the transportation cost, we are earning $\frac{\left|y_{1}-S\left(y_{1}\right)\right|^{2}}{2 \tau}$ and paying $\frac{\left|y_{2}-S\left(y_{1}\right)\right|^{2}}{2 \tau}$. But since $\rho_{1}$ is a minimizer of (10), what we are earning must be less or equal to what we are paying, and we get (18).

Inequality (19) is analogous: here we are just considering those points $x$ which are sent to the boundary by $T$. In this case, if we decide to send some small mass at $x$ onto a point $y \in \Omega$, we are not earning in terms of entropy but just paying $\log \left(\rho_{\tau}(y)\right)$, while in terms of cost we are earning $\frac{|x-T(x)|^{2}}{2 \tau}$ and paying $\frac{|x-y|^{2}}{2 \tau}$.

To prove inequality (21) we argue as follows. Consider first a point $y \in \Omega$, and perturb $\rho_{\tau}$ by picking some small mass from one of the nearest point to $y$ on $\partial \Omega$, and putting it onto $y$. In this way we pay $\log \left(\rho_{\tau}(y)\right)$ in terms of entropy, and $\frac{d^{2}(y, \partial \Omega)}{2 \tau}$ in terms of cost, so that by minimality we get

$$
\begin{equation*}
\frac{d^{2}(y, \partial \Omega)}{2 \tau} \geq-\log \left(\rho_{\tau}(y)\right) . \tag{23}
\end{equation*}
$$

The other part of the inequality comes by taking some small mass at $y$ and putting it on one of the nearest point to $y$ on $\partial \Omega$.

The proof of (22) is a sort of converse of (23). Indeed, since $S(y) \in \partial \Omega$, we know that the mass of $y$ is coming from the boundary. Hence we can perturb $\rho_{\tau}$ by taking a bit less of mass from the boundary, so that there is a bit less of mass in $y$. In this way we obtain the opposite of (21), and equality holds.

- Rigorous proof. We will prove rigorously only (18), the proof of the other inequalities being analogous.

Fix $y_{1}, y_{2} \in \Omega$, and two real numbers $r, a$, with $r>0$ small enough so that $B_{r}\left(y_{1}\right) \cup$ $B_{r}\left(y_{2}\right) \subset \Omega$, and $a \in(0,1)$. Let $\operatorname{Tr}: \mathbb{R}^{d} \rightarrow \mathbb{R}^{d}$ be the map defined by $\operatorname{Tr}(y):=y-y_{1}+y_{2}$, and let $\gamma \in \operatorname{Opt}\left(\rho, \rho_{\tau}\right)$ be the unique optimal plan. Define the plan $\gamma^{r, a}$ as

$$
\gamma^{r, a}:=\gamma_{\bar{\Omega}}^{B_{r}\left(y_{1}\right)^{c}}+a \gamma_{\bar{\Omega}}^{B_{r}\left(y_{1}\right)}+(1-a)\left(\left(\pi^{1}, \operatorname{Tr}\right)_{\#} \gamma_{\bar{\Omega}}^{B_{r}\left(y_{1}\right)}\right),
$$

and set

$$
\mu_{\tau}^{r, a}:=\pi_{\#}^{2} \boldsymbol{\gamma}^{r, a} .
$$

Observe that $\pi_{\#}^{1} \boldsymbol{\gamma}^{r, a}=\pi_{\#}^{1} \boldsymbol{\gamma}, \boldsymbol{\gamma}^{r, a} \in \operatorname{ADM}\left(\rho, \mu_{1}^{r, a}\right)$, and $\mu_{\tau}^{r, a}=\rho_{\tau}^{r, a} \mathcal{L}^{d}$, with

$$
\rho_{\tau}^{r, a}(y)= \begin{cases}\rho_{\tau}(y) & \text { if } y \in B_{r}\left(y_{1}\right)^{c} \cap B_{r}\left(y_{2}\right)^{c} \\ a \rho_{\tau}(y) & \text { if } y \in B_{r}\left(y_{1}\right) \\ \rho_{\tau}(y)+(1-a) \rho_{\tau}\left(y-y_{2}+y_{1}\right) & \text { if } y \in B_{r}\left(y_{2}\right)\end{cases}
$$

From the minimality of $\rho_{\tau}$ we get

$$
\int_{\Omega} e\left(\rho_{\tau}\right) d x+\frac{1}{2 \tau} C(\gamma) \leq \int_{\Omega} e\left(\rho_{\tau}^{r, a}\right) d x+\frac{1}{2 \tau} C\left(\gamma^{r, a}\right)
$$

Hence

$$
\begin{aligned}
& \int_{B_{r}\left(y_{1}\right) \cup B_{r}\left(y_{2}\right)} e\left(\rho_{\tau}(y)\right) d y+ \frac{1}{2 \tau} \int_{B_{r}\left(y_{1}\right) \cup B_{r}\left(y_{2}\right)}|y-S(y)|^{2} \rho_{\tau}(y) d y \\
& \leq \int_{B_{r}\left(y_{1}\right)} e\left(a \rho_{\tau}(y)\right) d y+\frac{a}{2 \tau} \int_{B_{r}\left(y_{1}\right)}|y-S(y)|^{2} \rho_{1}(y) d y \\
&+\int_{B_{r}\left(y_{2}\right)} e\left(\rho_{\tau}(y)+(1-a) \rho_{\tau}\left(y-y_{1}+y_{2}\right)\right) d y \\
&+\frac{1}{2 \tau} \int_{B_{r}\left(y_{2}\right)}|y-S(y)|^{2}\left(\rho_{\tau}(y)+(1-a) \rho_{\tau}\left(y-y_{1}+y_{2}\right)\right) d y
\end{aligned}
$$

which we write as

$$
\begin{aligned}
\int_{B_{r}\left(y_{1}\right)}\left(e\left(\rho_{\tau}(y)\right)-\right. & \left.e\left(a \rho_{\tau}(y)\right)+\frac{1-a}{2 \tau}|y-S(y)|^{2} \rho_{\tau}(y)\right) d y \\
\leq & \int_{B_{r}\left(y_{2}\right)}\left(e\left(\rho_{\tau}(y)+(1-a) \rho_{\tau}\left(y-y_{2}+y_{1}\right)\right)-e\left(\rho_{\tau}(y)\right)\right. \\
& \left.\quad+\frac{1-a}{2 \tau}|y-S(y)|^{2} \rho_{\tau}\left(y-y_{2}+y_{1}\right)\right) d y
\end{aligned}
$$

Dividing by $1-a$ and letting $a \uparrow 1$ we obtain

$$
\begin{aligned}
\int_{B_{r}\left(y_{1}\right)} & \left(e^{\prime}\left(\rho_{\tau}(y)\right)+\frac{1}{2 \tau}|y-S(y)|^{2}\right) \rho_{\tau}(y) d y \\
& \leq \int_{B_{r}\left(y_{2}\right)}\left(e^{\prime}\left(\rho_{1}(y)\right)+\frac{1}{2 \tau}|y-S(y)|^{2}\right) \rho_{\tau}\left(y-y_{2}+y_{1}\right) d y
\end{aligned}
$$

Now, since $y_{1}, y_{2}$ are both Lebesgue points of $\rho_{\tau}$, and $y_{1}$ is also a Lebesgue point of $S$, dividing both sides by $\mathcal{L}^{d}\left(B_{r}(0)\right)$, and letting $r \downarrow 0$ we obtain (18).

Proposition 3.8 (Monotonicity) Let $\mu \geq \tilde{\mu} \in \mathcal{M}_{+}(\Omega), \tau>0$, and $\mu_{\tau}, \tilde{\mu}_{\tau}$ the minima of the minimizing problem (10). Then $\mu_{\tau} \geq \tilde{\mu}_{\tau}$.

Proof. From the uniqueness part of Proposition 3.6, it follows easily that the map $\mu \mapsto$ $\mu_{\tau}$ is continuous w.r.t. the weak topology. Therefore, we can assume by approximation that both $\mu$ and $v$ are absolutely continuous, say $\mu=\rho \mathcal{L}^{d}$ and $\tilde{\mu}=\tilde{\rho} \mathcal{L}^{d}$. Also, recall that by Proposition 3.6(i) both $\mu_{\tau}$ and $v_{\tau}$ are absolutely continuous, say $\mu_{\tau}=\rho_{\tau} \mathcal{L}^{d}$ and $\tilde{\mu}_{\tau}=\tilde{\rho}_{\tau} \mathcal{L}^{d}$. Let $\underline{\gamma} \in \operatorname{Opt}\left(\rho, \rho_{\tau}\right)$ and $\tilde{\gamma} \in \operatorname{Opt}\left(\tilde{\rho}, \tilde{\rho}_{\tau}\right)$, and let $T, \tilde{T}$ be the maps which induce $\gamma_{\Omega}^{\bar{\Omega}}$ and $\tilde{\gamma}_{\Omega}^{\bar{\Omega}}$ respectively.

Argue by contradiction, and assume that $A:=\left\{\tilde{\rho}_{\tau}>\rho_{\tau}\right\} \subset \Omega$ satisfies $\tilde{\rho}_{\tau}(A)>0$. Two cases arise: either $\tilde{\gamma}_{\bar{\Omega}}^{A}$ is concentrated on $\Omega \times A$ or it is not, i.e. either the mass of $\tilde{\rho}_{\tau}$ in $A$ comes entirely from $\Omega$ or it is partly taken from the boundary.

Case 1: the mass of $\tilde{\rho}_{\tau}$ in $A$ comes entirely from $\Omega$. Let $B:=\tilde{T}^{-1}(A)$, and observe that $\tilde{\mu}(B)=\tilde{\mu}_{\tau}(A)$. Let $C \subset B$ be the set of points $x \in B$ such that $T(x) \notin A$. We remark that $\mu(C)>0$, as otherwise we would have

$$
\mu_{\tau}(A) \geq \mu_{\tau}(T(B))=\mu\left(T^{-1}(T(B))\right) \geq \mu(B) \geq \tilde{\mu}(B)=\tilde{\mu}_{\tau}(A),
$$

which contradicts the definition of $A$. Define

$$
C_{1}:=\{x \in C: T(x) \in \Omega\}, \quad C_{2}:=\{x \in C: T(x) \in \partial \Omega\} .
$$

Since $C=C_{1} \cup C_{2}$, either $\mu\left(C_{1}\right)>0$ or $\mu\left(C_{2}\right)>0$. Suppose we are in the first case. Then, as both $T_{C_{C_{1}}}$ and $\tilde{T}_{C_{1}}$ map subsets of the support of $\tilde{\rho}$ of positive Lebesgue measure into sets of positive Lebesgue measure, we can find $x \in C_{1}$ a Lebesgue for both $T$ and $\tilde{T}$, such that $T(x)$ and $\tilde{T}(x)$ are Lebesgue points for both $\rho_{\tau}$ and $\tilde{\rho}_{\tau}$. With this choice of $x$, we apply (18) with $y_{1}=T(x)$ and $y_{2}=\tilde{T}(x)$ to get

$$
\log \left(\rho_{\tau}(T(x))\right)+\frac{|x-T(x)|^{2}}{2 \tau} \leq \log \left(\rho_{\tau}(\tilde{T}(x))\right)+\frac{|x-\tilde{T}(x)|^{2}}{2 \tau}
$$

Similarly, using (18) for $\tilde{\rho}_{\tau}$ with $y_{1}=\tilde{T}(x)$ and $y_{2}=T(x)$ we obtain

$$
\log \left(\tilde{\rho}_{\tau}(\tilde{T}(x))\right)+\frac{|x-\tilde{T}(x)|^{2}}{2 \tau} \leq \log \left(\tilde{\rho}_{\tau}(T(x))\right)+\frac{|x-T(x)|^{2}}{2 \tau}
$$

Adding up the last two inequalities, we get

$$
\log \left(\rho_{\tau}(T(x))\right)+\log \left(\tilde{\rho}_{\tau}(\tilde{T}(x))\right) \leq \log \left(\rho_{\tau}(\tilde{T}(x))\right)+\log \left(\tilde{\rho}_{\tau}(T(x))\right)
$$

which contradicts definition of $C_{1}$ and the choice of $x$, as we have:

$$
\begin{aligned}
& T(x) \notin A \quad \Rightarrow \quad \rho_{\tau}(T(x)) \geq \tilde{\rho}_{\tau}(T(x)) \quad \Rightarrow \quad \log \left(\rho_{\tau}(T(x))\right) \geq \log \left(\tilde{\rho}_{\tau}(T(x))\right), \\
& \tilde{T}(x) \in A \quad \Rightarrow \quad \tilde{\rho}_{\tau}(\tilde{T}(x))>\rho_{\tau}(\tilde{T}(x)) \quad \Rightarrow \quad \log \left(\tilde{\rho}_{\tau}(\tilde{T}(x))\right)>\log \left(\rho_{\tau}(\tilde{T}(x))\right) .
\end{aligned}
$$

It remains to exclude the possibility $\mu\left(C_{2}\right)>0$. Fix $x \in C_{2}$ a Lebesgue point for both $T$ and $\tilde{T}$, such that $\tilde{T}(x)$ is a Lebesgue point for both $\rho_{\tau}$ and $\tilde{\rho}_{\tau}$. We apply (19) with $y=\tilde{T}(x)$ to obtain

$$
\frac{|x-T(x)|^{2}}{2 \tau} \leq \log \left(\rho_{\tau}(\tilde{T}(x))\right)+\frac{|x-\tilde{T}(x)|^{2}}{2 \tau}
$$

Now, we use (20) for $\tilde{\rho}_{\tau}$ with $y_{1}=\tilde{T}(x), S\left(y_{1}\right)=x$, and $y_{2}=T(x)$, to get

$$
\log \left(\tilde{\rho}_{\tau}(\tilde{T}(x))\right)+\frac{|x-\tilde{T}(x)|^{2}}{2 \tau} \leq \frac{|x-T(x)|^{2}}{2 \tau}
$$

Since $\tilde{T}(x) \in A$, we have $\rho_{\tau}(\tilde{T}(x))<\tilde{\rho}_{\tau}(\tilde{T}(x))$, which together with the above inequalities implies

$$
\begin{aligned}
\frac{|x-T(x)|^{2}}{2 \tau} & \leq \log \left(\rho_{\tau}(\tilde{T}(x))\right)+\frac{|x-\tilde{T}(x)|^{2}}{2 \tau} \\
& <\log \left(\tilde{\rho}_{\tau}(\tilde{T}(x))\right)+\frac{|x-\tilde{T}(x)|^{2}}{2 \tau} \leq \frac{|x-T(x)|^{2}}{2 \tau}
\end{aligned}
$$

again a contradiction.
Case 2: the mass of $\tilde{\rho}_{\tau}$ in $A$ comes partly from $\partial \Omega$. Let $\tilde{S}$ be the map which induces $\tilde{\gamma}_{\bar{\Omega}}^{\Omega}$ seen from $\tilde{\rho}_{\tau}$, and let $D \subset A$ be the set of points $y$ such that the mass $\tilde{\rho}_{\tau}(y)$ comes from the boundary, i.e. $D:=\{y \in A: \tilde{S}(y) \in \partial \Omega\}$. Fix $y \in D$ a Lebesgue point for $\rho_{\tau}, \tilde{\rho}_{\tau}$, and $\tilde{S}$. Thanks to (21) we have

$$
\log \left(\rho_{\tau}(y)\right)+\frac{d^{2}(y, \partial \Omega)}{2 \tau} \geq 0
$$

while applying (22) with $\tilde{\rho}_{\tau}$ (recall that $\tilde{S}(y) \in \partial \Omega$ ) we obtain

$$
\log \left(\tilde{\rho}_{\tau}(y)\right)+\frac{d^{2}(y, \partial \Omega)}{2 \tau}=0
$$

But this is absurd as $y \in D \subset A$.
Thanks to Proposition 3.8, we immediately obtain the following:
Corollary 3.9 (Comparison principle) Let $\mu_{0}, v_{0} \in \mathcal{M}_{+}(\Omega)$, assume that $\mu_{0} \geq \tilde{\mu}_{0}$, and let $\tau_{k} \downarrow 0$ be a sequence of time steps such that the corresponding discrete solutions $\mu^{\tau_{k}}(t), \tilde{\mu}^{\tau_{k}}(t)$ associated to $\mu_{0}, \tilde{\mu}_{0}$ respectively converge to two solutions $\mu_{t}, \tilde{\mu}_{t}$ of the heat equation, as described in Theorem 3.5. Then $\mu_{t} \geq \tilde{\mu}_{t}$ for all $t \in[0,+\infty)$.

Remark 3.10 [Different energies and costs] The proof of the above theorem relies entirely on the set of inequalities proved in Proposition (3.7). Here we want to point out that a corresponding version of such inequalities is true in more general cases.

Indeed, let $c: \bar{\Omega} \times \bar{\Omega} \rightarrow \mathbb{R} \cup\{+\infty\}$ be a continuous cost function, and define the Cost of transport as the infimum of

$$
\int_{\bar{\Omega} \times \bar{\Omega}} c(x, y) d \gamma(x, y),
$$

among all $\gamma \in \operatorname{ADM}\left(\mu_{0}, \mu\right)$. Let $e:[0,+\infty) \rightarrow \mathbb{R}$ be a superlinear convex function. Then, a minimizer $\rho_{1}$ for

$$
\rho \quad \mapsto \quad \int_{\Omega} e(\rho(x)) d x+\text { Cost of transport }\left(\rho, \mu_{0}\right)
$$

always exists, and arguing as in the proof of Proposition 3.7 it is possible to check that for $\rho_{1}$-a.e. $y_{1}, y_{2}$, and any $x$ such that $\left(x, y_{1}\right)$ belongs to the support of an optimal plan from $\mu$ to $\rho_{1}$, we have

$$
\frac{d}{d t^{-}} e\left(\rho_{1}\left(y_{1}\right)\right)+c\left(x, y_{1}\right) \leq \frac{d}{d t^{+}} e\left(\rho_{1}\left(y_{2}\right)\right)+c\left(x, y_{2}\right)
$$

and similarly for the other inequalities. Then the convexity of $e$ implies that

$$
\frac{d}{d t^{-}} e\left(z_{1}\right) \leq \frac{d}{d t^{+}} e\left(z_{1}\right) \leq \frac{d}{d t^{-}} e\left(z_{2}\right) \quad \forall 0 \leq z_{1}<z_{2}
$$

and the proof of the monotonicity goes on like in the case we analyzed. In particular, it is interesting to observe that the choice $c(x, y)=|x-y|^{2}$ in this setting does not play any role.

## 4 Comments and open problems

- The boundedness assumption on $\Omega$ was done just to make less technical the proofs, in order to clarify the new ideas in this approach. All our results can be generalized to unbounded domains, provided one works on the set of nonnegative measures on $\Omega$ such that

$$
\int d^{2}(x, \partial \Omega) d \mu(x)<+\infty
$$

- All our results could be extended to more general cost function and more general entropies. For instance, by considering $c(x, y)=|x-y|^{p}$ with $p>1$, and $e(z)=$ $z \log (z)-\alpha z$ with $\alpha \in \mathbb{R}$, one can construct a weak solution of

$$
\left\{\begin{aligned}
\frac{d}{d t} \rho(t) & =\Delta_{p} \rho(t) \\
\rho(0) & =\rho_{0}
\end{aligned}\right.
$$

(where $\Delta_{p} \rho$ denotes the $p$-Laplacian of $\rho$ ), subject to the Dirichlet boundary condition

$$
\rho(t)_{\partial \Omega}=e^{\alpha-1}, \quad \text { for a.e. } t \geq 0 .
$$

- It is interesting to observe that our approach allows to introduce a drift term in the diffusion: by considering the entropy $\int_{\Omega}[\rho \log \rho-V \rho] d x$ for some smooth function $V: \bar{\Omega} \rightarrow \mathbb{R}$ we obtain a weak solution of

$$
\left\{\begin{aligned}
\frac{d}{d t} \rho(t) & =\Delta \rho(t)-\operatorname{div}(\rho \nabla V) \\
\rho(0) & =\rho_{0}
\end{aligned}\right.
$$

subject to the Dirichlet boundary condition

$$
\rho(t)_{\mid \partial \Omega}=e^{V}, \quad \text { for a.e. } t \geq 0 .
$$

- A standard approach for constructing weak solutions to the heat equation with Dirichlet boundary condition equal to a function $\phi$ consists viewing the equation as the gradient flow of $\int_{\Omega}|\nabla \rho|^{2}$ on the set of functions $\rho \in H_{\phi}^{1}(\Omega):=\{\rho \in$ $\left.H^{1}(\Omega): \operatorname{trace}(\rho)=\phi\right\}$, with respect to the $L^{2}$-norm. However, although this approach allows to treat general boundary conditions, it cannot be used to add a drift term: given $F=F(x, u, p): \Omega \times \mathbb{R} \times \mathbb{R}^{d} \rightarrow \mathbb{R}$, the gradient flow of a functional of the form $\int_{\Omega} F(x, \rho, \nabla \rho) d x$ is given by

$$
\frac{d}{d t} \rho(t)=\operatorname{div}_{x}\left(F_{p}(x, \rho(t), \nabla \rho(t))\right)-F_{u}(x, \rho(t), \nabla \rho(t))
$$

and it is easy to check by a direct computation that there is no choice of $F$ which allows to obtain $\Delta \rho(t)-\operatorname{div}(\rho \nabla V)$ as the right-hand side.

- Although it is possible to prove uniqueness of solution by purely PDE methods, it is not clear to us if one can use a transportation approach to prove this result. In particular it is not clear if, as in the classical Wasserstein case, $t \mapsto W b_{2}\left(\rho_{t}, \tilde{\rho}_{t}\right)$ is decreasing along gradient flows of the entropy $\int_{\Omega} \rho \log (\rho) d x$.
- In Proposition 2.9 we only proved an upper bound for the derivative of $W b_{2}$. We conjecture that the following formula should be true: let $t \mapsto \mu_{t}$ an absolutely continuous curve with values in $\left(\mathcal{M}_{+}(\Omega), W b_{2}\right)$. Then:
(a) There exists a velocity field $\boldsymbol{w}_{t} \in L_{l o c}^{1}\left([0,+\infty), L^{2}\left(\Omega, \mu_{t}\right)\right)$ such that

$$
\frac{d}{d t} \mu_{t}+\operatorname{div}\left(\boldsymbol{w}_{t} \mu_{t}\right)=0
$$

in $[0,+\infty) \times \Omega$. (Observe that, since by definition the continuity equation can be tested only against smooth functions with support inside $[0,+\infty) \times \Omega$, the mass of $\mu_{t}$ is not necessarily constant.)
(b) Given $\mu \in \mathcal{M}_{+}(\Omega)$, for a.e. $t \geq 0$ we have

$$
\frac{d}{d t} W b_{2}^{2}\left(\mu_{t}, \mu\right)=-2 \int_{\Omega \times \Omega}\left\langle\boldsymbol{w}_{t}, y-x\right\rangle d \boldsymbol{\gamma}(x, y),
$$

where $\gamma$ is any optimal plan between $\mu_{t}$ and $\mu$.

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