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STOCHASTIC OPTIMAL TRANSPORT AND HAMILTON-JACOBI-BELLMAN EQUATIONS ON THE SET OF PROBABILITY MEASURES

CHARLES BERTUCCI ¹

ABSTRACT. We introduce a stochastic version of the optimal transport problem. We provide an analysis by means of the study of the associated Hamilton-Jacobi-Bellman equation, which is set on the set of probability measures. We introduce a new definition of viscosity solutions of this equation, which yields general comparison principles, in particular for cases involving terms modeling stochasticity in the optimal control problem. We are then able to establish results of existence and uniqueness of viscosity solutions of the Hamilton-Jacobi-Bellman equation. These results rely on controllability results for stochastic optimal transport that we also establish.

CONTENTS

Introduction	2
Optimal transport	2
Hamilton-Jacobi-Bellman equation on the set of probability measures	3
Organisation of the paper	3
Notation	4
1. From deterministic to stochastic optimal transport	4
1.1. Optimal transport	4
1.2. Optimal transport through dynamic programming	6
1.3. Warning on the formulation of the HJB equation	8
1.4. Stochastic optimal transport	9
1.5. HJB equations of stochastic optimal transport	9
2. A comparison principle for HJB equations on the set of probability measures	12
2.1. Super-differentials of functions on $\mathcal{P}(\mathbb{T}^d)$	13
2.2. Statement of the problem and definition of viscosity solutions	16
2.3. Comparison principle and uniqueness of viscosity solutions	17
2.4. Extensions to other HJB equations	20
2.5. Comments on our notion of viscosity solution	24
3. The value function is a viscosity solution of the HJB equation	25
3.1. Definition of the value function	25
3.2. The value function is a viscosity solution of the HJB equation	27
3.3. Comparison with the more usual value functions	29

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4.	Bounds on the value of the stochastic optimal transport problem near the singularity	32
4.1.	On the choice of the cost functional	32
4.2.	Controllability of the stochastic problems and L^∞ bounds of the value functions	33
5.	Conclusion and perspectives	41
5.1.	Summary of the results	41
5.2.	Potential approximating schemes for the HJB equation	41
5.3.	More general optimal control problem	41
	Acknowledgments	42
	References	42

INTRODUCTION

This paper introduces a stochastic version of the famous problem of optimal transport. We consider a dynamic formulation of the classical problem as in Benamou and Brenier [5] and are interested in the case in which the target measure is described by a stochastic process. This problem is a state-constrained stochastic optimal control problem, which is set on the set of probability measures. We adopt the dynamic programming approach and study the associated Hamilton-Jacobi-Bellman (HJB for short) equation. In particular, we prove a general comparison principle for viscosity solutions of HJB equations on the set of probability measures. Moreover, the HJB equation associated to the stochastic optimal transport problem is associated to a singularity at initial time which models the state constraint which is to reach the target once the problem is over.

Optimal transport. The optimal transport problem is one of the most famous problems in applied mathematics. It consists in finding the best way to transport a repartition of mass into another one, given a certain cost functional for the transport. Formulated first by Monge in [29], it has proven to be a mathematical problem of tremendous difficulty. The theoretical comprehension of this problem is now quite complete and the interest has now shifted onto more practical and numerical problems. More details on optimal transport can be found in the book of Villani [34] or in the one of Santambrogio [30].

A point of view which has proven to be particularly helpful to attack optimal transport is looking at a dynamic formulation of the problem. In the setting introduced by Benamou and Brenier [5], a time interval $[0, T]$ is given. The problem consists in transporting a measure m_0 into another measure m_T in this time interval $[0, T]$. This formulation is somehow closer to applications as its solutions describe precisely how the mass is going to be transported. Moreover, it naturally leads to the notion of geodesics in certain sets of measures.

In this paper we consider an extension of the aforementioned dynamic reformulation of optimal transport, in which the final repartition of mass, or target as we shall call it, is stochastic. More precisely, we shall assume that there is Markovian stochastic process

$(\nu_t)_{t \geq 0}$ such that the target repartition of mass is ν_T . This problem is introduced in more details in Section 1. Remark that such an optimal transport problem falls in the category of stochastic optimal control problem, in a space of measures, with a terminal state constraint.

Let us insist on the fact that, from the point of view of applications, this stochastic version of the optimal transport problem is natural and should prove to be of interest. Indeed, in our economy, the transportation of goods usually starts before the exact location of the addresses is known. This is for example the case for the delivery of oil in most ports. In a finite state space case, an analogue of this problem was studied in a mean field game framework in Bertucci et al. [7].

Hamilton-Jacobi-Bellman equation on the set of probability measures. As already mentioned, the stochastic optimal transport problem is a stochastic optimal control problem. Hence, naturally, the following study relies at some point on the study of the associated Hamilton-Jacobi-Bellman (HJB) equation. This partial differential equation (PDE) is set on the space of probability measures. We shall prove a comparison principle for viscosity sub and super solutions of this HJB equation. A particularity of the HJB equation associated to optimal transport problems, stochastic or not, is that it is associated with a singular boundary condition in time, namely because of the constraint that the target has to be reached. We provide an analysis of this singularity in Section 4.

Usually, the study of HJB equations relies mostly on the notion of viscosity solutions, introduced in finite dimensional space in Crandall and Lions [13]. We refer to Crandall et al. [16] for a complete presentation of viscosity solutions in finite dimensional space. The study of this notion in cases modeling state constraints, and the singular behaviour they can produce, is largely due to Lasry and Lions [24] and Soner [31, 32]. Infinite dimensional cases have also been investigated. HJB equations on Hilbert spaces have namely been studied in Crandall and Lions [14, 15] for first order equations and in Lions [25] for second order problems. Quite recently, the study of HJB equations set on the space of probability measures has gained a lot of interest, namely because of its link with the study of potential Mean Field Games. A major step in this study is the so-called Hilbertian approach, or lifting, introduced by Lions in [26]. An approach, which we may call more intrinsic, was developed in Gangbo and Swiech [18], Marigonda and Quincampoix [27] and in Jimenez et al. [21]. The links between Lions' Hilbertian approach and this more intrinsic approach is presented (among other things) in Gangbo and Tudorascu [19] and in Jimenez et al. [22]. Several authors have also considered methods relying on finite dimensional approximations of the PDE such as Gangbo et al. [20], Mayorga and Swiech [28] and Cosso et al. [12]. In Cecchin and Delarue [10], the authors used approximation of measures by Fourier sums to prove results of well-posedness for such equations. Finally let us also mention Conforti et al. [11] in which the authors considered a Schrödinger problem whose study leads to the same kind of PDE.

Organisation of the paper. The rest of the paper is organized as follows. In Section 1, we introduce the main problem at interest and derive the associated HJB equation. In Section 2, we provide the definition of viscosity solutions we are going to use as well as general comparison principles. We then show in Section 3 why the value functions of the stochastic optimal transport problem are indeed viscosity solutions of the associated HJB equations. We then proceed to establish some estimates on the behaviour of the value of the problem near the time singularity in Section 4, that we use in Section 5 to summarize our analysis.

Notation. Let us now introduce some notation.

- The d -dimensional torus is denoted by \mathbb{T}^d . The set of probability measures on \mathbb{T}^d is $\mathcal{P}(\mathbb{T}^d)$.
- Consider a function $\phi : \mathcal{P}(\mathbb{T}^d) \rightarrow \mathbb{R}$. When it is defined, we note for $\mu \in \mathcal{P}(\mathbb{T}^d), x \in \mathbb{T}^d$

$$(0.1) \quad \nabla_\mu \phi(\mu, x) = \lim_{\theta \rightarrow 0} \frac{\phi((1-\theta)\mu + \theta\delta_x) - \phi(\mu)}{\theta}.$$

- We note, if it is defined, $D_\mu \phi(\mu, x) = \nabla_x \nabla_\mu \phi(\mu, x) \in \mathbb{R}^d$. The second order derivatives are defined similarly.
- The image measure of a measure μ by a map T is denoted by $T_\# \mu$.
- The set of couplings between μ and ν in $\mathcal{P}(\mathbb{T}^d)$ is $\Pi(\mu, \nu)$.
- The notations usc and lsc stand for respectively upper semi continuous and lower semi continuous. The inf (resp. sup) envelope U_* (resp. U^*) of a locally bounded function U is defined by $U_*(x) = \liminf_{y \rightarrow x} U(y)$ (resp. $\limsup_{y \rightarrow x} U(y)$).
- The law of a random variable X is denoted by $\mathcal{L}(X)$.
- Given a n -uple $x = (x_1, x_2, \dots, x_n)$, we denote by $\pi_k(x) = x_k$.
- The set of $d \times d$ symmetric real matrices is denoted by $S_d(\mathbb{R})$.

1. FROM DETERMINISTIC TO STOCHASTIC OPTIMAL TRANSPORT

We introduce here the main mathematical problems at interest in this paper, starting with the well-known case of optimal transport.

1.1. Optimal transport. The problem of optimal transport consists in finding the best way (for a particular criteria) to transport a certain repartition of mass to another repartition of mass. We give a short presentation of this problem and refer to Villani [34] and Santamborgio [30] for more details on this topic. Given μ and ν two probability measures on measurable sets E_1 and E_2 , the main question of optimal transport is to find optimal measurable maps $T : E_1 \rightarrow E_2$ in the problem

$$(1.1) \quad \inf \{ \tilde{c}(T) \mid T_\# \mu = \nu \},$$

for a given real valued cost function \tilde{c} . Quite often, this cost is taken of the form

$$(1.2) \quad c(T) = \int_{E_1} c(x, T(x)) \mu(dx),$$

where $c : E_1 \times E_2 \rightarrow \mathbb{R}$. This problem lead to numerous mathematical developments since the seminal work of Monge. In the previous form, the problem has no minimizer in general. To observe this, it suffices to consider μ a Dirac mass and ν the Lebesgue

measure on some real interval. Indeed in this case the infimum is taken over an empty set. To address this issue, one usually introduces the relaxation of Kantorovich [23]. In this relaxed version, the typical form of the optimal transport problem becomes

$$(1.3) \quad \inf \int_{E_1 \times E_2} c(x, y) \pi(dx, dy),$$

where the infimum is taken over all couplings π between μ and ν , that is on probability measures on $E_1 \times E_2$ such that for any measurable sets $A \subset E_1, B \subset E_2$, $\pi(A \times E_2) = \mu(A)$ and $\pi(E_1 \times B) = \nu(B)$.

Let us also mention the natural probabilistic interpretation of such a problem. Consider a probabilistic space $(\Omega, \mathcal{A}, \mathbb{P})$. The previous relaxation of the optimal transport problem can be expressed as

$$(1.4) \quad \inf_{(X, Y)} \mathbb{E}[c(X, Y)],$$

where the infimum is taken over all the couples (X, Y) of random variables on $(\Omega, \mathcal{A}, \mathbb{P})$ such that $\mathcal{L}(X) = \mu$ and $\mathcal{L}(Y) = \nu$. Questions of existence, uniqueness and stability of optimal transport maps and optimal couplings have extensively studied since.

In this paper, we are mostly interested in the case $E_1 = E_2 = \mathbb{T}^d$. In this case we note $\Pi(\mu, \nu)$ the set of couplings μ and ν in $\mathcal{P}(\mathbb{T}^d)$. When the cost c is chosen as $c(x, y) = |x - y|^k$, the value of the optimal transport problem defines the Wasserstein distances through

$$(1.5) \quad W_k(\mu, \nu) = \left(\inf_{\pi \in \Pi(\mu, \nu)} \int_{\mathbb{T}^{2d}} |x - y|^k \pi(dx, dy) \right)^{\frac{1}{k}}.$$

The set of optimal couplings for the case $k = 2$ is denoted by $\Pi^{opt}(\mu, \nu)$.

One of the most useful approach for optimal transport problems has been the reformulation of (1.1) into

$$(1.6) \quad \inf_{(\alpha, m)} \int_0^1 \int_{\mathbb{T}^d} L(x, \alpha_t(x)) m_t(dx) dt,$$

where $L : \mathbb{T}^d \times \mathbb{R}^d \rightarrow \mathbb{R}$ is a certain cost function which is assumed to be bounded from below, and the infimum is taken over all pairs (α, m) such that $m : [0, 1] \rightarrow \mathcal{P}(\mathbb{T}^d)$ is a continuous map, $\alpha : [0, 1] \times \mathbb{T}^d \rightarrow \mathbb{R}^d$ is measurable and (α, m) satisfies in the weak sense

$$(1.7) \quad \begin{aligned} \partial_t m + \operatorname{div}(\alpha m) &= 0 \text{ in } (0, 1) \times \mathbb{T}^d, \\ m(0) &= \mu, m(1) = \nu. \end{aligned}$$

Let us insist on the fact that, in general, such a product αm is not well defined as a distribution, and thus the precise sense in which the previous equality holds has to be defined with care, which we postpone for the moment.

This approach is due to Benamou and Brenier [5] and it allows to interpret the optimal transport problem as a dynamic optimal control problem, with a terminal state

constraint, where the controlled state is a measure. As shown in [5], the optimality conditions of the problem (1.6) can be expressed through the following system of PDE

$$(1.8) \quad \begin{aligned} -\partial_t u + H(x, \nabla_x u) &= 0, \text{ in } (0, 1) \times \mathbb{T}^d \\ \partial_t m - \operatorname{div}(D_p H(x, \nabla_x u)m) &= 0 \text{ in } (0, 1) \times \mathbb{T}^d, \\ m(0) = \mu, m(1) &= \nu, \end{aligned}$$

where $H(x, \cdot)$ is the Legendre transform of $L(x, \cdot)$, given by

$$(1.9) \quad H(x, p) := \sup_{\alpha \in \mathbb{R}^d} \{-L(x, p) - \alpha \cdot p\}.$$

Let us remark that in this setting, the fact that the duration of the problem is 1 does not play any sort of role except for fixing some constants. This last approach is similar to the use of Pontryagin's maximum principle in dynamic optimal control.

1.2. Optimal transport through dynamic programming. We give a more dynamical approach, à la Bellman, of the optimal transport problem. The first thing to be said is that in this approach, the time parameter is crucial. This is of course obvious since we are doing dynamic programming. We adopt the convention that the terminal time, i.e. the time at which the target measure has to be reached is $T > 0$.

Let us introduce, formally, the value function U of the optimal transport problem, defined on $(0, \infty) \times \mathcal{P}(\mathbb{T}^d)^2$ by

$$(1.10) \quad U(t, \mu, \nu) = \inf_{\alpha, m} \int_t^T \int_{\mathbb{T}^d} L(x, \alpha(s, x)) m_s(dx) ds,$$

where the infimum is taken over all (α, m) satisfying the same measurability condition as in (1.6) and such that, in the weak sense,

$$(1.11) \quad \begin{aligned} \partial_s m + \operatorname{div}(\alpha m) &= 0 \text{ in } (t, T) \times \mathbb{T}^d, \\ m(t) = \mu, m(T) &= \nu. \end{aligned}$$

It is very tempting to analyze such a value function by a dynamic programming approach and the associated HJB equation. The study of HJB equations is now an extensively studied topic and we refer to the introduction for related works. The expression of the dynamic programming principle usually takes the form, for $0 < \delta < T - t$

$$(1.12) \quad U(t, \mu, \nu) = \inf_{\alpha, m} \left\{ \int_t^{t+\delta} \int_{\mathbb{T}^d} L(x, \alpha(s, x)) m_s(dx) ds + U(t + \delta, m_{t+\delta}, \nu) \right\},$$

where the infimum is taken over the same set as in (1.10) and $m_{t+\delta}$ is the value of m at time $t + \delta$. To obtain the associated HJB equation, the usual method is to divide by δ and to let $\delta \rightarrow 0$ in (1.12), under the assumption that U is smooth. Doing so yields

$$(1.13) \quad \begin{aligned} 0 &= -\partial_t U(t, \mu, \nu) - \\ &\quad - \liminf_{\delta \rightarrow 0} \inf_{\alpha, m} \left\{ \frac{1}{\delta} \int_t^{t+\delta} \int_{\mathbb{T}^d} L(s, x, \alpha(s, x)) m_s(dx) + \int_{\mathbb{T}^d} D_\mu U(t, \mu, \nu, y) \cdot \alpha(s, y) m_s(dy) ds \right\}. \end{aligned}$$

where we have used

(1.14)

$$\begin{aligned}
\delta^{-1}(U(t + \delta, m_{t+\delta}, \nu) - U(t, \mu, \nu)) &= \delta^{-1} \left(U(t + \delta, m_{t+\delta}, \nu) - U(t, m_{t+\delta}, \nu) \right. \\
&\quad \left. + \int_0^1 \int_{\mathbb{T}^d} \nabla_{\mu} U(t, \mu + \theta(m_{t+\delta} - \mu), \nu, y)(m_{t+\delta} - \mu)(dy) d\theta \right) \\
&= -\partial_t U(t, m_{t+\delta}, \nu) + o(1) + \\
&\quad + \delta^{-1} \int_0^1 \int_t^{t+\delta} \int_{\mathbb{T}^d} D_{\mu} U(t, \mu + \theta(m_{t+\delta} - \mu), \nu, y) \cdot \alpha(s, y) m_s(dy) ds d\theta \\
&= -\partial_t U(t, \mu, \nu) + \delta^{-1} \int_t^{t+\delta} \int_{\mathbb{T}^d} D_{\mu} U(t, \mu, \nu, y) \cdot \alpha(s, y) m_s(dy) ds + o(1),
\end{aligned}$$

and the fact that the $o(1)$ is uniform in (α, m) along minimizing sequences of the infimum. We do not insist too much on this assumption which is, in a lot of situations, immediate to verify given that L grows sufficiently fast with the size of α . Moreover, our aim is to derive the HJB equation, not particularly to consider smooth solutions of this PDE.

Recalling that U is assumed to be smooth here, we finally arrive at the HJB equation

$$(1.15) \quad -\partial_t U(t, \mu, \nu) + \mathcal{H}(\mu, D_{\mu} U) = 0 \text{ in } (0, \infty) \times \mathcal{P}(\mathbb{T}^d)^2,$$

where the Hamiltonian $\mathcal{H} : \mathcal{P}(\mathbb{T}^d) \times (\mathbb{T}^d \rightarrow \mathbb{R}^d) \rightarrow \mathbb{R}$ is given by

$$(1.16) \quad \mathcal{H}(\mu, \xi) = \int_{\mathbb{T}^d} H(x, \xi(x)) \mu(dx).$$

Note that in order for this Hamiltonian to be well defined, an integrability assumption has to be made on $x \rightarrow H(x, \xi(x))$ with respect to the measure μ .

Remark 1.1. *To be precise, we emphasize the fact that, a priori, the Hamiltonian \mathcal{H} also depends on ν since the infimum is taken over all admissible controls. Indeed we have not yet proven that, given any pair (α, m) defined on the time interval $[0, \delta]$ we can construct an admissible pair on $[0, t]$ which coincides with (α, m) on $[0, \delta]$. This will be the case under a controllability assumption, namely that from any starting measure μ , we can always transport μ toward ν in time t in finite cost. This will be the case for most of the problem we are interested in but we shall give an example in which this assumption is not verified.*

Moreover, because there is the state constraint at the terminal time T that the state measure μ has to be transported toward ν , we expect that U satisfies

$$(1.17) \quad U(T, \mu, \nu) = \begin{cases} 0 & \text{if } \mu = \nu, \\ +\infty & \text{otherwise.} \end{cases}$$

This is always satisfied by the value function since, if $\mu \neq \nu$, then the set of admissible controls is empty and thus the value infinite. However, as we shall see in Section 4, the behaviour of U as $t \rightarrow T$ might be of a different nature, depending on the nature of the

cost L .

Clearly, in this standard framework, ν is fixed and U is simply a function of t and μ . In the next section, we shall see why the addition of what is only a parameter here, is helpful to model more general problems.

The approach of studying (1.15) seems equivalent to (1.1). However the situation is the same as in standard finite dimensional optimal control. For several deterministic problems, the use of Pontryagin's maximum principle is efficient to provide a complete mathematical analysis. But for a larger class of problems, it is more convenient to use the dynamic programming principle and the associated HJB equation, this is in particular true for the stochastic problems that we are going to introduce later on.

Moreover, as usual in dynamic programming, if one is given a smooth solution U of (1.15), then a (smooth) closed-loop optimal control α^* in (1.10) can be computed using the derivatives of U by using the formula

$$(1.18) \quad \alpha^*(t, \mu, x) = -D_p H(x, D_\mu U(t, \mu, x)) \text{ in } (0, T) \times \mathcal{P}(\mathbb{T}^d) \times \mathbb{T}^d.$$

1.3. Warning on the formulation of the HJB equation. The formal computation which allowed us to derive (1.15) holds under a smoothness assumption on the value function which does not hold in general.

Indeed, if it was the case, then consider the problem of optimal transport which starts at $\mu = \delta_x$ for some $x \in \mathbb{T}^d$ when the time to reach ν is $t > 0$. If U is smooth, then an optimal control α is given as a smooth function of time and space. In particular, the induced trajectory, i.e. the unique solution of (1.11) will stay a Dirac mass at all time. Hence, as soon as the target measure ν is not a Dirac mass, we have a contradiction.

The PDE theory is used to derive the equations for smooth functions, and then provide weaker notion of solutions. However, we emphasize that the previous derivation might lead to a dangerous interpretation of the problem as it could lead to restrict the set of admissible controls. In our opinion, the analogy is very much in the spirit of Kantorovich's relaxation. If we restrict too much the set of admissible controls, we might be missing the only admissible controls. We shall come back later on this fact, as it bears some importance in the choice of the definition of viscosity solutions we are going to take.

Example 1.2. *In the case of a quadratic cost of optimal transport, i.e. when $L(x, p) = \frac{1}{2}|p|^2$, the associated HJB equation is given by*

$$(1.19) \quad -\partial_t U + \frac{1}{2} \int_{\mathbb{T}^d} |D_\mu U(t, \mu, \nu, x)|^2 \mu(dx) = 0 \text{ in } (0, \infty) \times \mathcal{P}(\mathbb{T}^d).$$

In this case, the value function U is simply given by

$$(1.20) \quad U(t, \mu, \nu) = \frac{1}{2(T-t)} W_2^2(\mu, \nu).$$

In particular, U is not smooth, see for instance Alfonsi and Jourdain [2]. We shall explain in Section 2 in which sense it is a viscosity solution of (1.19).

1.4. Stochastic optimal transport. This section introduces the main problem at interest in this paper, namely a stochastic version of (1.1). In this general formulation of the optimal transport problem, it may seem unclear what to do if either the cost function or any of the measure is random. Hence, we focus on the formulation (1.6). We work on a fixed filtered probability space $(\Omega, \mathcal{A}, \mathbb{P}, (\mathcal{F}_t)_{t \geq 0})$ which is assumed rich enough to contain independent Brownian motions.

We want to model the optimal transport of a given measure toward a stochastic target, in the time horizon $T > 0$. We assume here that the target measure is represented by an adapted Markovian stochastic process $(\nu_s)_{s \geq 0}$, valued in $\mathcal{P}(\mathbb{T}^d)$ and the (stochastic) target is given by ν_T . The problem we want to model is the following: the controlled state is a measure μ whose value at time t is μ_t . At time t , the trajectory $(\nu_s)_{s \in [0, t]}$ is known (obviously we do not know the future values of the target process, as this would put us in the usual framework). Then, we want to minimize a certain cost while transporting μ toward ν_T .

As usual in stochastic optimal control, some assumptions have to be made on how the optimization problem takes into account the randomness. To simplify the following discussion, we assume that the problem is risk neutral, hence the problem at interest is given by

$$(1.21) \quad \inf_{(\alpha, m)} \mathbb{E} \left[\int_0^T \int_{\mathbb{T}^d} L(x, \alpha_s(x)) m_s(dx) ds \right],$$

where $\alpha : \Omega \times [0, T] \times \mathbb{T}^d \rightarrow \mathbb{R}^d$ and $m : \Omega \times [0, T] \rightarrow \mathcal{P}(\mathbb{T}^d)$ have to be measurable maps which, almost surely in $\omega \in \Omega$, satisfy in the weak sense, the continuity equation

$$(1.22) \quad \begin{cases} \partial_t m + \operatorname{div}(\alpha m) = 0 & \text{in } (0, T) \times \mathbb{T}^d \\ m(0) = \mu, m(T) = \nu. \end{cases}$$

together with the condition that they have to be adapted processes to the filtration $(\mathcal{F}_t)_{t \geq 0}$.

We now derive the associated HJB equation for different target processes $(\nu_t)_{t \geq 0}$.

1.5. HJB equations of stochastic optimal transport. We define, formally for the moment, the value function U by

$$(1.23) \quad U(t, \mu, \nu) = \inf_{(\alpha, m)} \mathbb{E} \left[\int_t^T \int_{\mathbb{T}^d} L(x, \alpha_s(x)) m_s(dx) ds \middle| \nu_t = \nu \right],$$

where the state process (α, m) has to satisfy the same requirement as previously except that the condition $m_0 = \mu$ is now replaced by $m_t = \mu$. Note that the expectation is conditioned on $\{\nu_t = \nu\}$. Another point of view consists in looking at ν as an uncontrolled state variable of the optimal control problem, that we try to attain at the final time with the controlled state variable.

Depending on the nature of the process $(\nu_s)_{s \geq 0}$, different HJB equations arise for the value U in (1.23). We now give a few examples of such equations.

1.5.1. A constant target process.

Recall briefly that in the simplest case in which the target process $(\nu_s)_{s \geq 0}$ is constant, we recover the usual optimal transport problem and the associated HJB equation is then (1.15).

1.5.2. A Bernoulli like target process.

Consider a case in which at time $T/2$, the process ν is going to take the value ν_1 with probability $p \in (0, 1)$ and ν_2 with probability $1 - p$. It will then remain constant up to the final time T . In this context, after the time $T/2$, the problem is a standard (deterministic) optimal transport problem, whose value function we denote $U_{det}(t, \mu, \nu)$. Because the problem is assumed to be risk neutral, we can compute the value of the problem at time $T/2$, just before the value of ν is revealed. It is simply given by

$$(1.24) \quad U\left(\frac{T}{2}, \mu, \nu\right) = pU_{det}\left(\frac{T}{2}, \mu, \nu_1\right) + (1-p)U_{det}\left(\frac{T}{2}, \mu, \nu_2\right),$$

where U_{det} is the value of the associated deterministic optimal transport problem. We can then compute the value U for time $t \leq T/2$ by using the HJB equation

$$(1.25) \quad \partial_t U(t, \mu, \nu) + \mathcal{H}(\mu, D_\mu U) = 0 \text{ in } \left(0, \frac{T}{2}\right) \times \mathcal{P}(\mathbb{T}^d)^2,$$

together with the condition (1.24). Let us remark that the value of the process before $T/2$ does not matter.

1.5.3. A target process with jumps.

Consider now a case in which the target process $(\nu_s)_{s \geq 0}$ jumps, at times $(s_n)_{n \geq 0}$ given by a Poisson process of intensity $\lambda : [0, T] \rightarrow \mathbb{R}_+$, from ν_{s_n} into $\mathcal{T}\nu_{s_n}$, where $\mathcal{T} : \mathcal{P}(\mathbb{T}^d) \rightarrow \mathcal{P}(\mathbb{T}^d)$ is a given operator. In such a situation, the associated HJB equation is given by

$$(1.26) \quad -\partial_t U(t, \mu, \nu) + \mathcal{H}(\mu, \nabla_\mu U) + \lambda(t)(U(t, \mu, \nu) - U(t, \mu, \mathcal{T}\nu)) = 0 \text{ in } (0, T) \times \mathcal{P}(\mathbb{T}^d)^2.$$

Let us insist on the fact that such type of target process can cover a wide range of models. For instance, if \mathcal{T} is a constant operator, then the framework is quite close to the previous one and as at most two possible values of the target are possible, the initial one and its image by \mathcal{T} . Moreover, we could also consider cases involving more general jumps. For instance, assume that there is an independent sequence of times $(\tilde{s}_n)_{n \geq 0}$ given by a Poisson process of intensity λ_2 , associated to an operator \mathcal{T}_2 such that the previous rule also applies but also the process jumps according to \mathcal{T}_2 . In this case, the associated HJB equation would be

$$(1.27) \quad \begin{aligned} -\partial_t U(t, \mu, \nu) + \mathcal{H}(\mu, \nabla_\mu U) + \lambda(t)(U(t, \mu, \nu) - U(t, \mu, \mathcal{T}\nu)) \\ + \lambda_2(t)(U(t, \mu, \nu) - U(t, \mu, \mathcal{T}_2\nu)) = 0 \text{ in } (0, T) \times \mathcal{P}(\mathbb{T}^d)^2. \end{aligned}$$

1.5.4. *The target process is pushed by a diffusion.*

Consider now a case in which the target process $(\nu_s)_{s \geq 0}$ is given by $\nu_s = (\tau_{W_s})_{\#} \nu$ for a given $\nu \in \mathcal{P}(\mathbb{T}^d)$ and a process $(W_s)_{s \geq 0}$ (where $\tau_h(x) = x + h$ is the translation of h). We assume that $(W_s)_{s \geq 0}$ is given as the solution of the stochastic differential equation (SDE)

$$(1.28) \quad dW_t = \sigma(t)dB_t,$$

where $\sigma : \mathbb{R}_+ \rightarrow \mathbb{R}$ is a given function and $(B_t)_{t \geq 0}$ is a standard Brownian motion on $(\Omega, \mathcal{A}, \mathbb{P}, (\mathcal{F}_t)_{t \geq 0})$. In other words, the actual shape of the target measure is fixed by ν , but it is constantly being translated by the process $(W_t)_{t \geq 0}$. Using an infinite dimensional analogue of Itô's Lemma, such as in Cardaliaguet et al. [9] for instance, we deduce that the HJB equation characterizing the associated value U is given by

$$(1.29) \quad -\partial_t U(t, \mu, \nu) + \mathcal{H}(\mu, D_\mu U) - \frac{\sigma^2(t)}{2} \int_{\mathbb{T}^{2d}} \text{Tr} [D_{\nu\nu}^2 U] dm \otimes dm = 0 \text{ in } (0, \infty) \times \mathcal{P}(\mathbb{T}^d)^2.$$

Contrary to the previous case, the present situation leads to terms involving derivatives of the value function with respect to the variable ν which represents the target measure. This is a general feature of such problems. In some particular situations, including this one as we shall see, the problem can be reduced in such a way that those terms do not appear, however in a general situation we cannot avoid to work directly with them.

Let us remark that, formally, following the computations of Bertucci [6], this HJB equation can be obtained as the limit of the case with jumps for well chosen operators \mathcal{T} and jump rates λ .

1.5.5. *The case of a stochastic cost functional.*

Consider now a slightly different setting. We now assume that the target process $(\nu_s)_{s \geq 0}$ is constant and that the randomness is carried in the cost function L . We assume that this randomness appears through a dependence on the value w of a d dimensional process $(W_s)_{s \geq 0}$ given as the solution of (1.28). In this situation, it is natural to consider a value function U which also depends on the value w of this process. To be more precise, we are considering the value U defined by

$$(1.30) \quad U(t, \mu, \nu, w) = \inf_{(\alpha, m)} \mathbb{E} \left[\int_t^T \int_{\mathbb{T}^d} L(x, \alpha_s(x), W_s) m_s(dx) ds \middle| W_t = w \right],$$

where the state process $(m_s)_{s \in [0, T]}$ evolves as in (1.23), the infimum is carried over the same set.

In this situation, the natural HJB equation satisfied by U is

$$(1.31) \quad -\partial_t U(t, \mu, \nu, w) + \mathcal{H}(\mu, D_\mu U, w) - \frac{\sigma^2(t)}{2} \Delta_w U = 0 \text{ in } (0, T) \times \mathcal{P}(\mathbb{T}^d)^2 \times \mathbb{T}^d.$$

1.5.6. *Reduction of the case in which the target measure is pushed by a Brownian motion.* The equation (1.31) leads us to the following simplification of the case in which the target measure is pushed by the process $(W_s)_{s \geq 0}$. Indeed, as we mentioned the shape of the

final target is fixed at ν and thus only a finite dimensional parameter is sufficient to characterize it. More precisely, we want to make the formal change of variable

$$(1.32) \quad U(t, \mu, (\tau_{W_s})_{\#}\nu) = U(t, \mu, W_s).$$

This leads to the following HJB equation

$$(1.33) \quad -\partial_t U(t, \mu, w) + \mathcal{H}(\mu, D_\mu U) - \frac{\sigma^2(t)}{2} \Delta_w U = 0 \text{ in } (0, T) \times \mathcal{P}(\mathbb{T}^d) \times \mathbb{T}^d,$$

which is thus associated to the slightly more involved terminal condition

$$(1.34) \quad U(T, \mu, w) = \begin{cases} 0 & \text{if } \mu = (\tau_w)_{\#}\nu, \\ +\infty & \text{else.} \end{cases}$$

1.5.7. A comment on modelling.

Let us briefly comment on the choice we make to consider value functions as functions of both μ and ν . There seemed to be a wide range of models for which keeping this distinction is not necessary: for instance if the cost function L does not depend on the variable $x \in \mathbb{T}^d$. Indeed in this case, consider the equation (1.26) and assume that \mathcal{T} is a translation. Then studying (1.26) is equivalent to studying

$$(1.35) \quad -\partial_t U(t, \mu) + \mathcal{H}(\mu, \nabla_\mu U) + \lambda(t)(U(t, \mu) - U(t, \mathcal{T}^{-1}\mu)) = 0 \text{ in } (0, T) \times \mathcal{P}(\mathbb{T}^d).$$

This could have also been observed on the case of Section 1.5.4, by considering the equation

$$(1.36) \quad -\partial_t U(t, \mu) + \mathcal{H}(\mu, \nabla_\mu U) - \frac{\sigma^2(t)}{2} \int_{\mathbb{T}^{2d}} \text{Tr} [D_{\mu\mu}^2 U] dm \otimes dm = 0 \text{ in } (0, T) \times \mathcal{P}(\mathbb{T}^d).$$

We believe that this type of simplification can be helpful in several cases. Although it seems that the intrinsic nature of the HJB equation associated to this stochastic optimal transport is by nature involving these two variables: a controlled one which yields a Hamiltonian, and an uncontrolled one which yields the term associated to the generator of the stochastic evolution of the target process.

2. A COMPARISON PRINCIPLE FOR HJB EQUATIONS ON THE SET OF PROBABILITY MEASURES

As we mentioned in the introduction, the aim of this paper is to study the stochastic optimal transport problem by means of the associated HJB equation. Two main mathematical difficulties arise in this approach. The first one consists in studying the HJB equation in itself while the second one lies in the characterization of the singular terminal condition. In this section, we focus on the first question and postpone the question of the initial condition to the Section 4.

Here, we establish a general comparison principle for HJB equations on $\mathcal{P}(\mathbb{T}^d)$. We analyze first a pure HJB equation and we then explain how to extend it to HJB equations associated to stochastic optimal transport problems as the ones we mentioned before. This study is set on more general Hamiltonians than the one we have introduced before.

The notion of viscosity solution we introduce is different from the one usually used the literature. We believe that the present approach is better suited to study a wide

range of problems. Furthermore, we justify this notion in the next section when defining rigorously the value functions.

2.1. Super-differentials of functions on $\mathcal{P}(\mathbb{T}^d)$. Before presenting our notion of viscosity solution, we have to define a notion of super/sub-differential of functions on $\mathcal{P}(\mathbb{T}^d)$. Even though it is not the notion we are going to use, we start by recalling a common definition of super-differential.

In the literature, it is said that a function $\xi \in L^1((\mathbb{T}^d, \mu), \mathbb{R}^d)$ belongs to the super-differential of $U : \mathcal{P}(\mathbb{T}^d) \rightarrow \mathbb{R}$ at $\mu \in \mathcal{P}(\mathbb{T}^d)$ if for any $\mu' \in \mathcal{P}(\mathbb{T}^d)$, $\pi \in \Pi^{opt}(\mu, \mu')$,¹

$$(2.1) \quad U(\mu') \leq U(\mu) + \int_{\mathbb{T}^d \times \mathbb{T}^d} \xi(x) \cdot (y - x) \pi(dx, dy) + o(W_2(\mu, \mu')).$$

In such a situation we note $\xi \in \partial_{clas}^+ U(\mu)$. The sub-differential $\partial_{clas}^- U(\mu)$ is defined as $\partial_{clas}^- U(\mu) = -\partial_{clas}^+(-U)(\mu)$. When U is a smooth function, one recovers easily that $\partial_{clas}^+ U(\mu) = \partial_{clas}^- U(\mu) = \{D_\mu U(\mu)\}$.

Remark 2.1. *Let us insist on the fact that this notion of smoothness views $(\mathcal{P}(\mathbb{T}^d), W_2)$ as a geometric space whose geodesics are the ones of the optimal transport with cost $c(x, y) = |x - y|^2$. Indeed, in (2.1), we are considering optimal couplings between μ and μ' . Looking at $\mathcal{P}(\mathbb{T}^d)$ as a flat space, would lead to consider super-differentials $\partial_{flat}^+ U(\mu)$ as the set of $\phi \in C^0(\mathbb{T}^d, \mathbb{R})$ such that for all $\mu' \in \mathcal{P}(\mathbb{T}^d)$.*

$$(2.2) \quad U(\mu') \leq U(\mu) + \int_{\mathbb{T}^d} \phi(x)(\mu' - \mu)(dx) + o(W_2(\mu, \mu')).$$

In this case, we would have $\partial_{flat}^+ U(\mu) = \{\nabla_\mu U(\mu)\}$ for smooth functions U .

In this article we generalize the previous notion of super-differential in the spirit of Kantorovich's relaxation of (1.3). The simplest way to proceed is to replace $\xi : \mathbb{T}^d \rightarrow \mathbb{R}^d$ by a map $\psi : \mathbb{T}^d \rightarrow \mathcal{P}(\mathbb{R}^d)$. Equivalently, we can consider a measure $\gamma \in \mathcal{P}(\mathbb{T}^d \times \mathbb{R}^d)$ whose first marginal is $(\pi_1)_\# \gamma = \mu$ where μ is the measure at which we are looking for an element of the super-differential. As in (2.1), when considering the variations of a function U between μ' and μ we have to consider a coupling between the two measures, and not only the difference as in (2.2) for instance. We are not particularly interested with geodesics here so we shall not ask for the coupling to be optimal. Hence, we are lead to consider the following Definition.

Definition 2.2. *Given an upper semi continuous function $U : \mathcal{P}(\mathbb{T}^d) \rightarrow \mathbb{R}$, we say that a measurable map $\psi : \mathbb{T}^d \rightarrow \mathcal{P}(\mathbb{R}^d)$ is in the super-differential of U at the point μ if*

- *There exists $C > 0$ such that, for all $x \in \mathbb{T}^d$, the support of $\psi(x, \cdot)$ is contained in $B(0_{\mathbb{R}^d}, C)$.*
- *For any $\mu' \in \mathcal{P}(\mathbb{T}^d)$, for any $\gamma \in \Pi(\mu, \mu')$, the following holds*

$$(2.3) \quad U(\mu') - U(\mu) \leq \int_{\mathbb{T}^{2d}} \int_{\mathbb{R}^d} z \cdot (y - x) \psi(x, dz) \gamma(dx, dy) + o\left(\left(\int_{\mathbb{T}^{2d}} |x - y|^2 \gamma(dx, dy)\right)^{\frac{1}{2}}\right).$$

¹Recall that $\Pi^{opt}(\mu, \mu')$ is the set of optimal couplings between μ and μ' for the quadratic cost.

In this case we note $\psi \in \partial^+U(\mu)$.

- Remark 2.3.**
- The condition on the boundedness of the support of ψ is too strong at the level of this definition and could have been replaced by $(x \rightarrow \int_{\mathbb{R}^d} z\psi(x, dz)) \in L^1(\mu)$ so that (2.3) still makes sense. However, since we are working on the bounded set \mathbb{T}^d , this condition will not be too restrictive for the rest of the analysis, and will greatly help with the definition of the HJB equation on the elements of the super-differential, hence we leave it here for convenience.
 - The inequality (2.3) is a priori only carrying information when the term in the $o(\cdot)$ vanishes as $W_2(\mu, \mu') \rightarrow 0$.
 - If $\psi \in \partial^+U(\mu)$, then $x \rightarrow \int_{\mathbb{R}^d} z\psi(x, dz) \in \partial_{clas}^+U(\mu)$.
 - If $\xi \in \partial_{clas}^+U(\mu)$ then $x \rightarrow \delta_{\xi(x)} \in \partial^+U(\mu)$.

For a lower semi continuous function U , we define $\partial^-U(\mu) = \{x \rightarrow (-Id)_{\#}\psi(x, \cdot) | \psi \in \partial^+U(\mu)\}$.

We now provide what we believe to be an instructive computation around this notion of super-differentiability. This computation is based on Lions' Hilbertian approach. Consider a smooth function $\Phi : \mathcal{P}(\mathbb{T}^d) \rightarrow \mathbb{R}$, X and Y two \mathbb{T}^d valued random variables such that $\mathcal{L}(X) = \mu$ and $\mathcal{L}(Y) = \mu'$. We then want to evaluate the variations of Φ along the path $(m_s)_{s \in [0,1]}$ defined by $m_s = \mathcal{L}(X + s(Y - X))$.

$$(2.4) \quad \begin{aligned} \Phi(m_t) - \Phi(\mu) &= \int_0^t \mathbb{E}[D_\mu \Phi(m_s, X + s(Y - X)) \cdot (Y - X)] ds \\ &= \int_0^t \mathbb{E}[D_\mu \Phi(\mu, X) \cdot (Y - X)] ds + o(\mathbb{E}[|Y - X|]). \end{aligned}$$

Because Φ is smooth, it defines a smooth mapping on the \mathbb{T}^d valued random variables $\tilde{\Phi} : X \rightarrow \Phi(\mathcal{L}(X))$. On this last computation we see that the derivative $D_\mu \Phi$ is linked to the gradient of $\tilde{\Phi}$. Because Φ is smooth, the gradient of $\tilde{\Phi}$ at X is in fact of the form $\xi(X)$ for some map $\xi : \mathbb{T}^d \rightarrow \mathbb{R}^d$. The notion of super-differential we provided consists in looking for random variables in the super-differential of $\tilde{\Phi}$ without any particular restriction whereas in the super-differential ∂_{clas}^+ , the random variable in the super-differential has to be a function of X .

The following Proposition states the super-differentiability of the squared Wasserstein distance. For the usual notion, such a result was already proved in Ambrosio and Gangbo [3], Proposition 4.3.

Proposition 2.4. *For any $\mu, \nu \in \mathcal{P}(\mathbb{T}^d)$, $\gamma^o \in \Pi^{opt}(\mu, \nu)$, consider the measurable map ψ defined almost everywhere by the disintegration $(\pi_1, \pi_1 - \pi_2)_{\#}\gamma^o = \mu(dx)\psi(x, dz)$. The function $\Phi : \mu' \rightarrow \frac{1}{2}W_2^2(\mu', \nu)$ is such that $\psi \in \partial^+\Phi(\mu)$.*

Proof. Consider $\gamma \in \Pi(\mu, \mu')$ and its disintegration $\gamma(dx, dy) = \mu(dx)k(x, dy)$. By definition of Φ

(2.5)

$$\begin{aligned} 2\Phi(\mu') - 2\Phi(\mu) &\leq \int_{\mathbb{T}^{3d}} |y - z|^2 \gamma^o(dx, dy) k(x, dz) - \int_{\mathbb{T}^{2d}} |x - z|^2 \gamma^o(dx, dy) \\ &= \int_{\mathbb{T}^{3d}} |y - x + x - z|^2 - |x - z|^2 \gamma^o(dx, dy) k(x, dz) \\ &= 2 \int_{\mathbb{T}^{3d}} (y - x) \cdot (x - z) \gamma^o(dx, dy) k(x, dz) + \int_{\mathbb{T}^{2d}} |y - x|^2 \gamma(dx, dz). \end{aligned}$$

From which the result follows. Remark in particular that since \mathbb{T}^d is bounded, the bound on the support is indeed verified. \square

This result of everywhere super-differentiability justifies the use of the 2-Wasserstein distance in the argument of doubling of variables that we are going to make afterwards to obtain a comparison principle.

Remark 2.5. *The previous result can be interpreted in the probabilistic or Hilbertian approach. It states that, given an optimal coupling (X, Y) for the quadratic cost between μ and ν , we can consider an element of the super-differential of Φ which is constructed on the random variable $X - Y$ and not just on $X - \mathbb{E}[Y|X]$.*

Another advantage of this definition of super-differential is that it makes more transparent the link with the so-called Hilbertian approach. Let $(\Omega, \mathcal{A}, \mathbb{P})$ be an atomless standard probabilistic space and K be the set of \mathbb{T}^d valued random variables from this space. We have the following result.

Proposition 2.6. *Consider an usc function $U : \mathcal{P}(\mathbb{T}^d) \rightarrow \mathbb{R}$. Define $\tilde{U} : K \rightarrow \mathbb{R}$ by $\tilde{U}(X) = U(\mathcal{L}(X))$. Take $X \in K$ and assume that $Z \in \mathbb{L}^2(\Omega, \mathbb{R}^d)$ is such that for all $Y \in K$,*

$$(2.6) \quad \tilde{U}(Y) \leq \tilde{U}(X) + \mathbb{E}[Z \cdot (Y - X)] + o\left(\sqrt{\mathbb{E}[|X - Y|^2]}\right).$$

Consider now a function $\psi : \mathbb{T}^d \rightarrow \mathcal{P}(\mathbb{R}^d)$ such that

$$(2.7) \quad \mathcal{L}(X, Z)(dx, dz) = \mathcal{L}(X)(dx) \psi(x, dz).$$

Then, $\psi \in \partial^+ U(\mathcal{L}(X))$.

Proof. Set $\mu = \mathcal{L}(X)$. Take $\gamma \in \Pi(\mu, \mu')$ and disintegrate $\gamma(dx, dy)$ into $\mu(dx)k(x, dy)$.

Consider $\psi : \mathbb{T}^d \times \mathbb{R}^d \rightarrow \mathcal{P}(\mathbb{R}^d)$ such that $\mathcal{L}(X, Z)(dx, dz) = \mu(dx) \psi(x, dz)$. Consider now

(X', Y', Z') such that $\mathcal{L}((X', Y', Z'))(dx, dy, dz) = \mu(dx) \psi(x, dz) k(x, dy)$. Thanks to

classical results, we can in fact consider a sequence $(X_n, Y_n, Z_n)_{n \geq 0}$ such that $\mathcal{L}((X_n, Y_n, Z_n))(dx, dy, dz) =$

$\mu(dx) \psi(x, dz) k(x, dy)$ and $\|(X_n, Z_n) - (X, Z)\|_\infty \leq n^{-1}$. It holds that

$$(2.8) \quad U(\mu') = \tilde{U}(Y_n) \leq U(\mu) + \mathbb{E}[Z \cdot (Y_n - X)] + o\left(\sqrt{\mathbb{E}[|X - Y_n|^2]}\right).$$

On the other hand,

$$(2.9) \quad \begin{cases} |\mathbb{E}[Z \cdot (Y_n - X)] - \mathbb{E}[Z_n \cdot (Y_n - X_n)]| \xrightarrow{n \rightarrow \infty} 0, \\ \mathbb{E}[|X - Y_n|^2] \xrightarrow{n \rightarrow \infty} \int_{\mathbb{T}^{2d}} |x - y|^2 \gamma(dx, dy). \end{cases}$$

and

$$(2.10) \quad \mathbb{E}[Z_n \cdot (Y_n - X_n)] = \int_{\mathbb{T}^{2d} \times \mathbb{R}^d} z \cdot (y - x) \psi(x, dz) \gamma(dx, dy).$$

Hence the result follows from passing to the limit $n \rightarrow \infty$ in (2.8). \square

Remark 2.7. *In other words, we have elements of the super-differential of U which describes all the elements of the super-differential of \tilde{U} .*

We end this Section by explaining on a simple example how to consider super-differential of functions of more variables than a measure $\mu \in \mathcal{P}(\mathbb{T}^d)$. For instance, consider the case of an additional time variable. Given $T > 0$ and an usc function $U : [0, T] \times \mathcal{P}(\mathbb{T}^d) \rightarrow \mathbb{R}$, we note $(\theta, \psi) \in \partial^+ U(t, \mu)$ if

$$(2.11) \quad \begin{aligned} U(s, \mu') \leq & U(t, \mu) + \theta(s - t) + \int_{\mathbb{T}^{2d}} \int_{\mathbb{R}^d} z \cdot (y - x) \psi(x, dz) \gamma(dx, dy) + \\ & + o(|s - t|) + o\left(\left(\int_{\mathbb{T}^{2d}} |x - y|^2 \gamma(dx, dy)\right)^{\frac{1}{2}}\right), \end{aligned}$$

for any $s \leq t, \mu' \in \mathcal{P}(\mathbb{T}^d), \gamma \in \Pi(\mu, \mu')$. Similarly we introduce $\partial^- U(t, \mu) = \{(-\theta, x \rightarrow (-Id)_\# \psi(x, \cdot)) | (\theta, \psi) \in \partial^+(-U)(t, \mu)\}$.

We are now able to define viscosity solutions of (2.12).

2.2. Statement of the problem and definition of viscosity solutions. In this section we want to prove a comparison principle for HJB equations of the form

$$(2.12) \quad \partial_t U(t, \mu, \nu) + \mathcal{H}(t, \mu, D_\mu U) = 0 \text{ in } (0, \infty) \times \mathcal{P}(\mathbb{T}^d),$$

where $\mathcal{H} : \mathbb{R}_+ \times \mathcal{P}(\mathbb{T}^d) \times (\mathbb{T}^d \rightarrow \mathbb{R}^d) \rightarrow \mathbb{R}$ is given by

$$(2.13) \quad \mathcal{H}(t, \mu, \xi) = \int_{\mathbb{T}^d} H(t, x, \mu, \xi(x)) \mu(dx),$$

where $H : \mathbb{R}_+ \times \mathbb{T}^d \times \mathcal{P}(\mathbb{T}^d) \times \mathbb{R}^d$ satisfies

- H is globally continuous.
- There exists $C > 0$ such that for all $p \in \mathbb{R}^d, s, t \in \mathbb{R}_+, \mu, \nu \in \mathcal{P}(\mathbb{T}^d), x, y \in \mathbb{R}^d$ we have

$$(2.14) \quad |H(t, x, \mu, p) - H(t, y, \nu, p)| \leq C(1 + |p|)(|t - s| + W_2(\mu, \nu) + |x - y|).$$

Let us insist on the fact that, to simplify this Section, we have reversed the sense of time compared to the HJB equations derived in Section 1.

As mentioned earlier, the expected candidate to be a solution of (2.12) is in general not smooth and we have to define a notion of weak solution. The natural techniques to study HJB equations such as (2.12) comes from the theory of viscosity solutions. The approach we provide here is somehow close to the one of Marigonda and Quincampoix [27] in the sense that we provide an intrinsic proof of a comparison principle, and it is also close to the point of view of Lions' Hilbertian approach and Gangbo and Tudorascu [19] in the sense that the notion of viscosity solution we are going to provide relies on ideas from this Hilbertian lifting. However, our result is more general than the ones of

Marigonda and Quincampoix [27] and Gangbo and Tudorascu [19] and we believe the proof we provide to be simpler.

In our definition of super-differential, an element of the super differential (with respect to the measure variable) is a map $\mathbb{T}^d \rightarrow \mathcal{P}(\mathbb{R}^d)$ and not simply a function $\mathbb{T}^d \rightarrow \mathbb{R}^d$. Thus, we have to precise how we want to evaluate \mathcal{H} on such elements. We introduce here $\bar{\mathcal{H}} : \mathbb{R}_+ \times \mathcal{P}(\mathbb{T}^d) \times (\mathbb{T}^d \rightarrow \mathcal{P}(\mathbb{R}^d)) \rightarrow \mathbb{R}$ defined by

$$(2.15) \quad \bar{\mathcal{H}}(t, \mu, \psi) = \int_{\mathbb{T}^d \times \mathbb{R}^d} H(t, x, \mu, y) \mu(dx) \psi(x, dy),$$

and we shall work with the following Definition.

Definition 2.8. *An usc (resp. lsc) function $U : \mathbb{R}_+ \times \mathcal{P}(\mathbb{T}^d) \rightarrow \mathbb{R}$ is a viscosity super (resp. sub)-solution of (2.12) if, for any $t > 0, \mu \in \mathcal{P}(\mathbb{T}^d)$ and $(\theta, \psi) \in \partial^+(U)(t, \mu)$ (resp. $\partial^-(U)(t, \mu)$) the following holds*

$$(2.16) \quad \theta + \bar{\mathcal{H}}(t, \mu, \psi) \leq 0 \quad (\text{resp. } \geq 0).$$

A viscosity solution of (2.12) is a locally bounded function such that U^ is a viscosity sub-solution and U_* is a viscosity super-solution.*

Remark 2.9. *The term $\bar{\mathcal{H}}(t, \mu, \psi)$ is well defined because by definition of the super/sub-differential, $\psi(x, dz)$ has bounded support in $z \in \mathbb{R}^d$, uniformly in $x \in \mathbb{T}^d$.*

Remark 2.10. *The choice we made to consider $\bar{\mathcal{H}}$ is not trivial. It shall be justified in the next Section when proving that the value function of the stochastic optimal transport is indeed a viscosity solution of the HJB equation.*

The main advantage of this notion of viscosity solutions, by comparison with other ones in the literature, is that it provides, relatively easily, a comparison principle, as we shall now see.

2.3. Comparison principle and uniqueness of viscosity solutions. As usual in the theory of viscosity solutions, uniqueness of solutions, L^∞ estimates and other stability properties come from a comparison principle. We now establish such a result.

Theorem 2.11. *Assume that U and V are respectively viscosity sub and super-solution of (2.12) such that for all $\mu \in \mathcal{P}(\mathbb{T}^d)$, $U(0, \mu) \leq V(0, \mu)$. Then for all $t \geq 0, \mu \in \mathcal{P}(\mathbb{T}^d)$, $U(t, \mu) \leq V(t, \mu)$.*

Following the standard ways to establish comparison principle, we are going to use the so-called technique of doubling of variable. We now present formally this technique. The proof of the Theorem is postponed to the end of this section. In this setting on $\mathcal{P}(\mathbb{T}^d)$, we introduce, for $\epsilon > 0$, the function

$$(2.17) \quad (t, s, \mu, \mu') \rightarrow V(s, \mu') - U(t, \mu) + \frac{1}{2\epsilon}((t - s)^2 + W_2^2(\mu, \mu')),$$

Arguing by contradiction, we will consider a point (t^*, s^*, μ^*, ν^*) of minimum of this function. Then, using Proposition 2.4, we will prove that $\partial^+U(t^*, \mu^*)$ and $\partial^-V(s^*, \nu^*)$ are non empty. More precisely, we will be able to consider two elements, one in each of

those sets, with some relation between them.

The next, final and main step of the proof consists in arriving at a contradiction by taking the difference of the viscosity relations, i.e. the relations given by the fact that U is a sub-solution and V a super-solution. Before presenting the proof of the final step, we prove the Lemma that we are going to use in order to consider elements of $\partial^+U(s^*, \nu^*)$ and $\partial^-V(t^*, \mu^*)$.

Lemma 2.12. *Consider an usc function U and a continuous function Φ on $\mathbb{R}_+ \times \mathcal{P}(\mathbb{T}^d)$ and $(t, \mu) \in (0, \infty) \times \mathcal{P}(\mathbb{T}^d)$, a point of maximum of $U - \Phi$. Then $(\theta, \psi) \in \partial^+\Phi(t, \mu) \Rightarrow (\theta, \psi) \in \partial^+U(t, \mu)$.*

Proof. Take $t > 0$, $\mu \in \mathcal{P}(\mathbb{T}^d)$ such that $(t, \mu) \in \operatorname{argmax}\{U - \Phi\}$ and also $(\theta, \psi) \in \partial^+\Phi(t, \mu)$. For any $s \leq t$, $\mu' \in \mathcal{P}(\mathbb{T}^d)$ and $\gamma \in \Pi(\mu, \mu')$,

$$(2.18) \quad \begin{aligned} U(s, \mu') - U(t, \mu) &\leq \Phi(s, \mu') - \Phi(t, \mu) \\ &\leq \theta(s - t) + \int_{\mathbb{T}^{2d}} \int_{\mathbb{R}^d} z \cdot (y - x) \psi(x, dz) \gamma(dx, dy) \\ &\quad + o\left(|s - t| + \left(\int_{\mathbb{T}^{2d}} |x - y|^2 \gamma(dx, dy)\right)^{\frac{1}{2}}\right) \end{aligned}$$

□

We are now ready to prove the main result of this section.

Proof. (Of Theorem 2.11.) Assume that the result does not hold. Hence, there exists $\kappa > 0$ such that

$$(2.19) \quad \inf_{t \leq T, \mu \in \mathcal{P}(\mathbb{T}^d)} V(t, \mu) - U(t, \mu) \leq -\kappa.$$

Thus, there exists $\rho > 0$ such that for any $\epsilon > 0$,

$$(2.20) \quad \inf_{t, s \leq T, \mu, \mu'} V(s, \mu') - U(t, \mu) + \frac{1}{2\epsilon} \left((t - s)^2 + W_2^2(\mu, \mu') \right) + \rho(t + s) \leq -\frac{\kappa}{2}.$$

Since U is usc and V lsc, the previous infimum is reached at some point $(t_\epsilon, s_\epsilon, \mu_\epsilon, \mu'_\epsilon)$, since we are minimizing a lsc function on the compact $[0, T]^2 \times \mathcal{P}(\mathbb{T}^d)^2$.

Step 1: using the viscosity properties. We treat first the case $t_\epsilon, s_\epsilon > 0$. Take $\gamma_\epsilon^o \in \Pi^{opt}(\mu_\epsilon, \mu'_\epsilon)$ and denote $\psi_\epsilon : \mathbb{T}^d \rightarrow \mathcal{P}(\mathbb{R}^d)$ such that

$$(2.21) \quad (\pi_1, \epsilon^{-1}(\pi_1 - \pi_2)) \# \gamma_\epsilon^o(dx, dz) = \mu_\epsilon(dx) \psi_\epsilon(x, dz).$$

From Proposition 2.4 and Lemma 2.12, we obtain that

$$(2.22) \quad \begin{cases} (\rho + \epsilon^{-1}(t_\epsilon - s_\epsilon), \psi_\epsilon) \in \partial^+U(t_\epsilon, \mu_\epsilon), \\ (-\rho - \epsilon^{-1}(s_\epsilon - t_\epsilon), \psi_\epsilon) \in \partial^-V(s_\epsilon, \mu'_\epsilon). \end{cases}$$

Since U and V are respectively sub and super-viscosity solution of (2.12), we deduce that

$$(2.23) \quad \rho + \epsilon^{-1}(t_\epsilon - s_\epsilon) + \int_{\mathbb{T}^{2d}} H\left(t_\epsilon, x, \mu_\epsilon, \frac{x - y}{\epsilon}\right) \gamma_\epsilon^o(dx, dy) \leq 0,$$

and that

$$(2.24) \quad -\rho - \epsilon^{-1}(s_\epsilon - t_\epsilon) + \int_{\mathbb{T}^{2d}} H\left(s_\epsilon, y, \mu'_\epsilon, \frac{x-y}{\epsilon}\right) \gamma_\epsilon^o(dx, dy) \geq 0.$$

Step 2: Standard estimates.

Taking the differences of the two previous inequalities leads to

$$(2.25) \quad 2\rho \leq \int_{\mathbb{T}^{2d}} H(s_\epsilon, y, \mu'_\epsilon, \epsilon^{-1}(x-y)) \gamma_\epsilon^o(dx, dy) - \int_{\mathbb{T}^{2d}} H(t_\epsilon, x, \mu_\epsilon, \epsilon^{-1}(x-y)) \gamma_\epsilon^o(dx, dy).$$

Using the regularity assumptions we made on H , we deduce that

$$(2.26) \quad \begin{aligned} 2\rho &\leq C\epsilon^{-1} \int_{\mathbb{T}^{2d}} (|t_\epsilon - s_\epsilon| + |x-y| + W_2(\mu_\epsilon, \mu'_\epsilon)) |x-y| \gamma_\epsilon^o(dx, dy) \\ &\leq C\epsilon^{-1} \left(W_2^2(\mu_\epsilon, \mu'_\epsilon) + (|t_\epsilon - s_\epsilon| + W_2(\mu_\epsilon, \mu'_\epsilon)) \int_{\mathbb{T}^d} |x-y| \gamma_\epsilon^o(dx, dy) \right). \end{aligned}$$

From standard estimates techniques of the method of doubling of variables, see for instance Crandall et al. [16], we know that

$$(2.27) \quad \frac{(t_\epsilon - s_\epsilon)^2}{\epsilon} \longrightarrow 0, \frac{W_2^2(\mu_\epsilon, \mu'_\epsilon)}{\epsilon} \longrightarrow 0, \text{ as } \epsilon \rightarrow 0.$$

Hence if the previous holds for all $\epsilon > 0$, we arrive at a contradiction by taking the limit $\epsilon \rightarrow 0$ since $\rho > 0$ was fixed independently of ϵ .

Step 3: The minimum is at the boundary $t = 0$.

From the previous Step, we deduce that, for ϵ small enough, the minimum is in fact reached at a point such that either $s_\epsilon = 0$ or $t_\epsilon = 0$. Since, we obtain from (2.20) that $\lim_{\epsilon \rightarrow 0} \epsilon^{-1}(t_\epsilon - s_\epsilon)^2 = \lim_{\epsilon \rightarrow 0} \epsilon^{-1}W_2^2(\mu_\epsilon, \mu'_\epsilon) = 0$, we deduce that $\lim_{\epsilon \rightarrow 0} t_\epsilon = \lim_{\epsilon \rightarrow 0} s_\epsilon = 0$. Moreover, extracting a subsequence if necessary, there exists μ_0 , limit of both $(\mu_\epsilon)_{\epsilon > 0}$ and $(\mu'_\epsilon)_{\epsilon > 0}$. Hence, using the lower semi continuity of V and the upper semi continuity of U , we deduce that

$$(2.28) \quad \begin{aligned} V(0, \mu_0) - U(0, \mu_0) &\leq \liminf_{\epsilon \rightarrow 0} V(t_\epsilon, \mu_\epsilon) - U(s_\epsilon, \mu'_\epsilon) - \frac{1}{2\epsilon}((t_\epsilon - s_\epsilon)^2 + W_2^2(\mu_\epsilon, \mu'_\epsilon)) - \alpha(t_\epsilon + s_\epsilon) \\ &\leq -\frac{\kappa}{2}. \end{aligned}$$

The previous being clearly a contradiction, we finally deduce that the Theorem is true. \square

From this comparison principle easily follows the following result of uniqueness.

Theorem 2.13. *Given a continuous initial condition U_0 , there exists at most one viscosity solution U of (2.12) such that for all $\mu \in \mathcal{P}(\mathbb{T}^d)$, $U^*(0, \mu) = U_0(\mu) = U_*(0, \mu)$.*

Proof. By considering two such solutions U and V , using the comparison principle, we immediately arrive at the fact that $V \leq U \leq V$ which proves the claim. \square

More generally, we can use the comparison principle to establish stability results or L^∞ estimates. For instance, the following is an immediate corollary of Theorem 2.11.

Corollary 2.14. *Consider a viscosity subsolution U and a viscosity supersolution V of (2.12), then $\frac{d}{dt} \max(U(t, \cdot) - V(t, \cdot)) \leq 0$.*

2.4. Extensions to other HJB equations. We now explain how to make use of the previous results, or more precisely of their proofs, to study the equations (1.26) and (1.33).

2.4.1. The case of jumps. We focus here on (1.26). Formally, it suffices to remark that the terms in λ in (1.26) do not involve derivatives of the solution and thus are quite easy to treat. Moreover, the fact that the functions depend here on two measures instead of one does not perturb the previous argument as we shall now see. The definition of viscosity solutions of (1.26) takes the form

Definition 2.15. *An usc (resp. lsc) function $U : [0, T] \times \mathcal{P}(\mathbb{T}^d) \rightarrow \mathbb{R}$ is said to be a viscosity sub-solution (resp. super-solution) of (1.26) if for any $t \in [0, T)$, $\mu, \nu \in \mathcal{P}(\mathbb{T}^d)$ and $(\theta, \psi, \psi') \in \partial^+ U(t, \mu, \nu)$ (resp. $\in \partial^- U(t, \mu)$)*

$$(2.29) \quad -\theta + \bar{\mathcal{H}}(t, \mu, \psi) + \lambda(t)(U(t, \mu, \nu) - U(t, \mu, \mathcal{T}\nu)) \leq 0 \text{ (resp. } \geq 0).$$

A viscosity solution of (1.26) is a locally bounded function U such that U_ is a viscosity super-solution and U^* is a viscosity sub-solution.*

Remark 2.16. *Of course the existence of an element in the super-differential in the ν variable is useless here, and could be removed.*

As in the previous case, a comparison principle can be stated.

Proposition 2.17. *Assume that U and V are respectively viscosity sub and super-solutions of (1.26) such that for all μ, ν , $U(T, \mu, \nu) \leq V(T, \mu, \nu)$ and such that they are both bounded functions. Assume also that λ is a continuous non negative function. Then for all time $t \in [0, T]$ and measures $\mu, \nu \in \mathcal{P}(\mathbb{T}^d)$, $U(t, \mu, \nu) \leq V(t, \mu, \nu)$.*

Proof. We only explain how the addition of the term in λ perturbs the proof of Theorem 2.11. As in the previous proof, we consider the function

$$(2.30) \quad Z(t, s, \mu, \nu, \mu', \nu') := V(t, \mu, \nu) - U(s, \mu', \nu') + \frac{1}{2\epsilon} (W_2^2(\mu, \mu') + W_2^2(\nu, \nu') + (t-s)^2) + \rho(2T - t - s).$$

Considering a point of minimum $(t_\epsilon, s_\epsilon, \mu_\epsilon, \nu_\epsilon, \mu'_\epsilon, \nu'_\epsilon)$ of Z , if $t_\epsilon, s_\epsilon < T$, and arguing exactly as we did before we easily arrive at

$$(2.31) \quad 2\rho + \lambda(s_\epsilon)(U(s_\epsilon, \mu'_\epsilon, \nu'_\epsilon) - U(s_\epsilon, \mu'_\epsilon, \mathcal{T}\nu'_\epsilon)) - \lambda(t_\epsilon)(V(t_\epsilon, \mu_\epsilon, \nu_\epsilon) - V(t_\epsilon, \mu_\epsilon, \mathcal{T}\nu_\epsilon)) \leq o(1),$$

where the right side term vanishes as $\epsilon \rightarrow 0$. Let us compute

$$(2.32) \quad \begin{aligned} & \lambda(s_\epsilon)(U(s_\epsilon, \mu'_\epsilon, \nu'_\epsilon) - U(s_\epsilon, \mu'_\epsilon, \mathcal{T}\nu'_\epsilon)) - \lambda(t_\epsilon)(V(t_\epsilon, \mu_\epsilon, \nu_\epsilon) - V(t_\epsilon, \mu_\epsilon, \mathcal{T}\nu_\epsilon)) \\ & \geq -C|\lambda(t_\epsilon) - \lambda(s_\epsilon)| + \lambda(s_\epsilon)(U(s_\epsilon, \mu'_\epsilon, \nu'_\epsilon) - U(s_\epsilon, \mu'_\epsilon, \mathcal{T}\nu'_\epsilon) - (V(t_\epsilon, \mu_\epsilon, \nu_\epsilon) - V(t_\epsilon, \mu_\epsilon, \mathcal{T}\nu_\epsilon))) \\ & \geq -C|\lambda(t_\epsilon) - \lambda(s_\epsilon)| + \lambda(s_\epsilon) \frac{1}{2\epsilon} (W_2^2(\mathcal{T}\nu_\epsilon, \mathcal{T}\nu'_\epsilon) - W_2^2(\nu_\epsilon, \nu'_\epsilon)) \\ & \geq -C|\lambda(t_\epsilon) - \lambda(s_\epsilon)| - \lambda(s_\epsilon) \frac{1}{2\epsilon} W_2^2(\nu_\epsilon, \nu'_\epsilon). \end{aligned}$$

Remark that in the previous, C only depends on the bounds on U and V . The limit of the last lower bound in the previous chain of inequalities is 0 as $\epsilon \rightarrow 0$. Indeed, as in the previous proof of the comparison principle, we also recover that $\lim_{\epsilon \rightarrow 0} s_\epsilon - t_\epsilon = \lim_{\epsilon \rightarrow 0} \epsilon^{-1} W_2^2(\nu_\epsilon, \nu'_\epsilon) = 0$. Hence, using the continuity of λ we obtain the required result by following the same argument as in the proof of Theorem 2.11. \square

Remark 2.18. *No assumption on \mathcal{T} is needed here.*

2.4.2. *The case of a target measure being pushed by a diffusion.* We now turn to the case of (1.33). This equation being of second order in w , the definition of viscosity solution is more involved. Indeed, because we are interested in viscosity solutions of a second order HJB equation, we need to introduce super-jets. We only consider super-jets which are of interest for our problem here, namely we are only interested in second order derivatives in the w variable.

For a function $U : [0, T] \times \mathcal{P}(\mathbb{T}^d) \times \mathbb{T}^d \rightarrow \mathbb{R}$, and $(t, \mu, w) \in [0, T] \times \mathcal{P}(\mathbb{T}^d) \times \mathbb{T}^d$, the super-jet of U at (t, μ, w) $J^+(U)(t, \mu, w)$ is defined as the set of $(\theta, \psi, p, X) \in \mathbb{R} \times (\mathbb{T}^d \rightarrow \mathcal{P}(\mathbb{R}^d)) \times \mathbb{R}^d \times S_d(\mathbb{R})$ such that for any $t' \geq t, \mu' \in \mathcal{P}(\mathbb{T}^d), w' \in \mathbb{T}^d$ and $\gamma \in \Pi(\mu, \mu')$,

$$(2.33) \quad \begin{aligned} U(t', \mu', w') \leq & U(t, \mu, w) + \theta(t' - t) + \int_{\mathbb{T}^{2d}} \int_{\mathbb{R}^d} z \cdot (y - x) \gamma(dx, dy) \psi(x, dz) \\ & + p \cdot (w' - w) + (w' - w) \cdot X \cdot (w' - w) \\ & + o(|t' - t|) + o\left(\left(\int_{\mathbb{T}^{2d}} |x - y|^2 d\gamma(dx, dy)\right)^{\frac{1}{2}}\right) + o(|w' - w|^2). \end{aligned}$$

We can now introduce the notion of viscosity solutions of (1.33).

Definition 2.19. *An usc (resp. lsc) function U is a viscosity sub-solution (resp. super-solution) of (1.33) if for any $(t, \mu, w) \in [0, T] \times \mathcal{P}(\mathbb{T}^d) \times \mathbb{T}^d$, $(\theta, \psi, p, X) \in J^+(U)(t, \mu, w)$ (resp. $J^-(U)(t, \mu, w)$)*

$$(2.34) \quad -\theta + \bar{\mathcal{H}}(t, \mu, \psi) - \frac{\sigma^2(t)}{2} \text{Tr}(X) \leq 0 \text{ (resp. } \geq 0\text{)}.$$

A viscosity solution of (1.33) is a bounded function U such that U_ is a super-solution and U^* is a sub-solution.*

Once again, a comparison principle result holds for this type of equation.

Proposition 2.20. *Assume that in addition to the other requirements on H , there exists $C > 0$ such that for all $t \in [0, T], x \in \mathbb{T}^d, \mu \in \mathcal{P}(\mathbb{T}^d), p \in \mathbb{R}^d$*

$$(2.35) \quad \begin{aligned} |H(t, x, \mu, p)| &\leq C(1 + |p|^2), \\ |D_p H(t, x, \mu, p)| &\leq C(1 + |p|). \end{aligned}$$

Let U and V be respectively a bounded viscosity sub-solution and a bounded viscosity super-solution of (1.33). If $U(0, \mu, w) \leq V(0, \mu, w)$, then $U \leq V$.

Proof. As usual in viscosity solution theory, we argue by contradiction and we assume that there exists $\kappa, \rho > 0$, such that for any $\epsilon > 0$

$$(2.36) \quad \inf \left\{ V(s, \mu', w') - U(t, \mu, w) + \frac{1}{2\epsilon} \left((t - s)^2 + W_2^2(\mu, \mu') + |w - w'|^2 \right) + \rho(2T - t - s) \right\} \leq -\kappa,$$

where the infimum is taken over all $s, t \leq T, w, w' \in \mathbb{T}^d, \mu, \mu' \in \mathcal{P}(\mathbb{T}^d)$.

Step 1: Reformulation of the problem in the Hilbert space.

Let us define $\tilde{V}(t, X, w) = V(t, \mathcal{L}_{\mathbb{T}^d}(X), w)$ and $\tilde{U}'(t, X, w) = U(t, \mathcal{L}_{\mathbb{T}^d}(X), w)$ for $X \in \mathbb{L}^2(\Omega, \mathbb{R}^d)$, and where, for the rest of this proof, for $X \in \mathbb{L}^2(\Omega, \mathbb{R}^d)$, $\mathcal{L}_{\mathbb{T}^d}(X) = \mathcal{L}(\text{proj}(X))$ where $\text{proj} : \mathbb{R}^d \rightarrow \mathbb{T}^d$ is the natural projection.

We can now consider

$$(2.37) \quad \Phi(t, s, X, Y, w, w') := \tilde{V}(s, Y, w') - \tilde{U}(t, X, w) + \frac{1}{2\epsilon} \left((t-s)^2 + \mathbb{E}[|X - Y|^2] + |w - w'|^2 \right) + \rho(2T - t - s) + \alpha(\eta(X) + \eta(Y)),$$

where $\eta(X) := \sqrt{1 + \|X\|^2}$. We obtain that, for $\alpha > 0$ sufficiently small,

$$(2.38) \quad \inf \left\{ \Phi(t, s, X, Y, w, w') \mid t, s \leq T, X, Y \in \mathbb{L}^2(\Omega, \mathbb{R}^d), w, w' \in \mathbb{T}^d \right\} \leq -\frac{\kappa}{2}.$$

Thanks to Stegall's Lemma [33], we know that for any $\delta > 0$, there exists $Z_1, Z_2 \in \mathbb{L}^2(\Omega, \mathbb{R}^d)$, $\|Z_1\|, \|Z_2\| \leq \delta$ and

$$\Phi(t, s, X, Y, w, w') + \mathbb{E}[Z_1 \cdot X + Z_2 \cdot Y]$$

has a unique strict minimum at point $(\bar{t}, \bar{s}, \bar{X}, \bar{Y}, \bar{w}, \bar{w}')$. Assume first that $\bar{t}, \bar{s} < T$. From Lemma 4 in Lions [25], we obtain that there exist sequences $(\omega_n, \omega'_n, \xi_n, \xi'_n, A_n, B_n)_{n \geq 0}$ converging toward 0 and $(t_n, s_n, X_n, Y_n, w_n, w'_n)_{n \geq 0}$ converging toward $(\bar{t}, \bar{s}, \bar{X}, \bar{Y}, \bar{w}, \bar{w}')$ such that ²

$$(2.39) \quad \begin{aligned} &(-\rho + \epsilon^{-1}(\bar{t} - \bar{s}) + \omega_n, \epsilon^{-1}(\bar{X} - \bar{Y}) + \xi_n + Z_1 + \alpha \nabla \eta(\bar{X}), p_n, S + A_n) \in J^+(\tilde{U})(t_n, X_n, w_n), \\ &(\rho + \epsilon^{-1}(\bar{t} - \bar{s}) + \omega'_n, \epsilon^{-1}(\bar{X} - \bar{Y}) + \xi'_n - Z_2 - \alpha \nabla \eta(\bar{Y}), p'_n, S' + B_n) \in J^-(\tilde{V})(s_n, Y_n, w'_n), \\ &S \leq S'. \end{aligned}$$

Step 2: Coming back to the original formulation.

Thanks to Proposition 2.6, (2.39) implies in particular that

$$(2.40) \quad \begin{aligned} &(-\rho + \epsilon^{-1}(\bar{t} - \bar{s}) + \omega_n, \psi_n, p_n, S + A_n) \in J^+(U)(t_n, \mathcal{L}(X_n), w_n), \\ &(\rho + \epsilon^{-1}(\bar{t} - \bar{s}) + \omega'_n, \varphi_n, p'_n, S' + B_n) \in J^-(V)(s_n, \mathcal{L}(Y_n), w'_n), \\ &S \leq S', \end{aligned}$$

where ψ_n and φ_n satisfy

$$(2.41) \quad \begin{aligned} &\mathcal{L}(X_n, \epsilon^{-1}(\bar{X} - \bar{Y}) + \xi_n + Z_1 + \alpha \nabla \eta(\bar{X}))(dx, dz) = \mathcal{L}(X_n)(dx) \psi_n(x, dz), \\ &\mathcal{L}(Y_n, \epsilon^{-1}(\bar{X} - \bar{Y}) + \xi'_n - Z_2 - \alpha \nabla \eta(\bar{Y}))(dx, dz) = \mathcal{L}(Y_n)(dy) \varphi_n(x, dz). \end{aligned}$$

Using the fact that U is a subsolution of (1.33), we obtain that

$$(2.42) \quad \rho - \epsilon^{-1}(\bar{t} - \bar{s}) + \omega_n + \bar{\mathcal{H}}(t_n, \mathcal{L}(X_n), \psi_n) - \frac{1}{2} \text{Tr}(S + A_n) \leq 0.$$

²We use the same notation for the super jets of the functions \tilde{U} and \tilde{V} to avoid introducing a new one.

Using the fact that V is a super solution, we obtain that

$$(2.43) \quad -\rho - \epsilon^{-1}(\bar{t} - \bar{s}) + \omega'_n + \bar{\mathcal{H}}(s_n, \mathcal{L}(Y_n), \varphi_n) - \frac{1}{2}Tr(S' + A_n) \geq 0.$$

Let us compute

$$\begin{aligned} & |\bar{\mathcal{H}}(t_n, \mathcal{L}(X_n), \psi_n) - \mathbb{E}[H(\bar{t}, \bar{X}, \mathcal{L}(\bar{X}), \epsilon^{-1}(\bar{X} - \bar{Y}) + Z_1 + \alpha \nabla \eta(\bar{X}))]| \\ &= |\mathbb{E}[H(t_n, X_n, \mathcal{L}(X_n), \epsilon^{-1}(\bar{X} - \bar{Y}) + \xi_n + Z_1 + \alpha \nabla \eta(\bar{X}))] \\ &\quad - \mathbb{E}[H(\bar{t}, \bar{X}, \mathcal{L}(\bar{X}), \epsilon^{-1}(\bar{X} - \bar{Y}) + Z_1 + \alpha \nabla \eta(\bar{X}))]|. \end{aligned}$$

From the growth assumption on H and the dominated convergence Theorem, we deduce that the previous difference vanishes as $n \rightarrow \infty$. Hence, we can pass to the limit $n \rightarrow \infty$ in (2.42) and (2.43) and we obtain

$$(2.44) \quad \begin{aligned} & \rho - \epsilon^{-1}(\bar{t} - \bar{s}) + \mathbb{E}[H(\bar{t}, \bar{X}, \mathcal{L}(\bar{X}), \epsilon^{-1}(\bar{X} - \bar{Y}) + Z_1 + \alpha \nabla \eta(\bar{X}))] - \frac{1}{2}TrS \leq 0, \\ & -\rho - \epsilon^{-1}(\bar{t} - \bar{s}) + \mathbb{E}[H(\bar{s}, \bar{Y}, \mathcal{L}(\bar{Y}), \epsilon^{-1}(\bar{X} - \bar{Y}) - Z_2 - \alpha \nabla \eta(\bar{Y}))] - \frac{1}{2}TrS' \geq 0. \end{aligned}$$

Step 3: Standard viscosity solutions estimates.

Taking the difference of the two previous inequalities yields

$$(2.45) \quad 2\rho \leq \mathbb{E}[H(\bar{t}, \bar{X}, \mathcal{L}(\bar{X}), \epsilon^{-1}(\bar{X} - \bar{Y}) + Z_1 + \alpha \nabla \eta(\bar{X})) - H(\bar{s}, \bar{Y}, \mathcal{L}(\bar{Y}), \epsilon^{-1}(\bar{X} - \bar{Y}) - Z_2 - \alpha \nabla \eta(\bar{Y}))].$$

From the assumptions on the Hamiltonian, we deduce that

$$\begin{aligned} 2\rho \leq & C\mathbb{E}[\epsilon^{-1}|\bar{X} - \bar{Y}|(|\bar{X} - \bar{Y}| + |\bar{t} - \bar{s}| + \sqrt{\mathbb{E}[|X - Y|^2]}) \\ & + C\mathbb{E}[(1 + \epsilon^{-1}|\bar{X} - \bar{Y}|)(|Z_1 + \alpha \eta(\bar{X})| + |Z_2 + \alpha \eta(\bar{Y})|)]. \end{aligned}$$

From the same argument as in the proof of Theorem 2.11, we obtain that the first term of the right side vanishes as $\epsilon \rightarrow 0$, uniformly in α and δ , recall that δ measures the size of Z_1 and Z_2 . From Cauchy Schwarz inequality we finally deduce that the second term of the right side is bounded by

$$C(1 + \epsilon^{-1}\|\bar{X} - \bar{Y}\|)(\delta + \alpha).$$

Hence, taking the limits $\delta \rightarrow 0$ then $\alpha \rightarrow 0$ then $\epsilon \rightarrow 0$, we conclude that $2\rho \leq 0$, which is a contradiction.

The minimum is reached at T .

The case in which the minimum is reached for $\bar{t} = T$ or $\bar{s} = T$ can be treated as in the proof of Theorem 2.11 to arrive at a similar contradiction and we do not reproduce it. In definitive we have indeed proven that $U \leq V$. \square

Remark 2.21. *The assumptions on H specified in the statement of the Theorem can be removed by the use of techniques which are not particularly new to viscosity solutions theory. However, since such questions are not the core ones of our paper, we leave them for other works.*

These comparison principles are essential tools to characterize functions as viscosity solutions of equations of the form of (1.15). Would the terminal conditions in our problems be continuous functions, the previous results would be enough to develop a a

proper theory of (1.15). However because of the singularity that we expect at the origin, we shall have to characterize the behaviour of the solution near $t = T$ to have proper comparison principle. Namely, we shall use the following result.

Theorem 2.22. *Consider U and V , two viscosity solutions of (1.26), locally bounded in $[0, T] \times \mathcal{P}(\mathbb{T}^d)^2$ such that*

$$(2.46) \quad \lim_{t \rightarrow T} \sup_{\mu, \nu \in \mathcal{P}(\mathbb{T}^d)} |U(t, \mu, \nu) - V(t, \mu, \nu)| = 0.$$

Then $U = V$.

Proof. By a symmetry argument, it is sufficient to prove $U \leq V$. Assume that it is not the case, hence that there exists $\kappa > 0$ and $t \in [0, T)$, $\mu, \nu \in \mathcal{P}(\mathbb{T}^d)$ such that

$$(2.47) \quad U(t, \mu, \nu) - V(t, \mu, \nu) > \kappa.$$

Take $\delta > 0$. By exactly the same argument as in the proof of Proposition 2.17, we deduce that there exists $\rho > 0$, such that for $\epsilon > 0$ small enough the minimum of the function

$$(2.48) \quad (t, s, \mu, \nu, \mu', \nu') \rightarrow V(t, \mu, \nu) - U(s, \mu', \nu') + \frac{1}{2\epsilon} ((t-s)^2 + W_2^2(\mu, \mu') + W_2^2(\nu, \nu')) + \rho(2T - s - t)$$

on the set $[0, T - \delta]^2 \times \mathcal{P}(\mathbb{T}^d)^4$, is reached for either $t = T - \delta$ or $s = T - \delta$ (recall that U and V are both bounded on $[0, T - \delta] \times \mathcal{P}(\mathbb{T}^d)^2$). Taking the limit $\epsilon \rightarrow 0$, we deduce that

$$(2.49) \quad \sup_{\mu, \nu \in \mathcal{P}(\mathbb{T}^d)} U(T - \delta, \mu, \nu) - V(T - \delta, \mu, \nu) > \kappa.$$

Taking the limit $\delta \rightarrow 0$, we obtain a contradiction and thus the result is proved. \square

The same type of result obviously holds true for (1.33).

2.5. Comments on our notion of viscosity solution. In recent years, the study of HJB equations on the set of probability measure been the subject to a lot of works which have failed to establish general comparison principles for HJB equations associated to stochastic problems. On the other hand, the study of HJB equations set on an Hilbert space is a problem which is for the most part solved at the moment, except maybe for some singular problems.

We believe that our approach provides a link between the two problems, namely through the notion of super-differential which we have chosen. In our opinion, this is a strong justification of the so-called Hilbertian approach developed by P.-L. Lions, originally to study MFG master equation.

Recall that, in this approach, a probabilistic space $(\Omega, \mathcal{A}, \mathbb{P})$ is fixed, and the study of (2.12) is replaced by the study of

$$\partial_t \tilde{U} + \mathbb{E}[H(t, X, \mathcal{L}(X), \nabla \tilde{U})] = 0 \text{ in } (0, \infty) \times \mathbb{L}^2(\Omega, \mathbb{R}^d),$$

where formally we have made the change of variable $\tilde{U}(X) = U(\mathcal{L}_{\mathbb{T}^d}(X))$.

This approach hints strongly the notion of super differential that we took. But maybe more importantly, it provides an interpretation for this HJB equation in the Hilbert space. Indeed, the Hilbertian approach can be interpreted as a process of labelling all the elements of mass of the measures, namely by labels $\omega \in \Omega$. This procedure allows

to split mass, by assigning to the elements ω and ω' different velocities $Z(\omega)$ and $Z(\omega')$ even if they are in the same location, that is $X(\omega) = X(\omega')$, which is very reminiscent of Kantorovich's relaxation of the optimal transport problem. We come back on this kind of interpretation in the next Section.

Remark 2.23. *In the choice of super-differential we made, everything could have also been true by using not only a coupling between μ' and μ but a coupling $\Gamma(dx, dy, dz)$ between $\mu'(dy)$ and $\mu(dx)\psi(x, dz)$. Such that we could have said that $\psi \in \partial^+ U(\mu)$ if for all $\Gamma \in \mathcal{P}(\mathbb{T}^d \times \mathbb{T}^d \times \mathbb{R}^d)$ such that $(\pi_1, \pi_3)_\# \Gamma(dx, dz) = \mu(dx)\psi(x, dz)$ and $(\pi_2)_\# \Gamma(dy) = \mu'(dy)$, it holds that*

(2.50)

$$U(\mu') \leq U(\mu) + \int_{\mathbb{T}^{2d} \times \mathbb{R}^d} z \cdot (y - x) \Gamma(dx, dy, dz) + o \left(\left(\int_{\mathbb{T}^{2d} \times \mathbb{R}^d} |x - y|^2 \Gamma(dx, dy, dz) \right)^{\frac{1}{2}} \right).$$

3. THE VALUE FUNCTION IS A VISCOSITY SOLUTION OF THE HJB EQUATION

We start by defining properly the value function, formally introduced in Section 1. In view of the previous Section, we do so by reformulating lightly the way we evaluate the cost of a given trajectory. We then proceed to show that these value functions are indeed viscosity solutions of the associated HJB equations. We end this Section by showing that the more common way to define value functions of such problems indeed coincides with our relaxed version.

3.1. Definition of the value function. In this Section we mainly focus on the value function of the deterministic problem. The main difficulty in defining the value function lies in the definition of the set on which the infimum is taken in (1.10). Indeed, without regularity constraints on α and m , it is not clear how to evaluate the derivative of the product. Furthermore, α and m have to be such that the integral which yields the cost is indeed well defined. These difficulties make that it is difficult to talk about α as the control and about m as the state, as given a control, it is not clear how to define the state, as multiple solutions to the continuity equation can exist.

In order to address this issue, in [5], Benamou and Brenier introduced a reformulation of the problem (1.6) into

$$(3.1) \quad \inf_{m, E} \int_0^1 \int_{\mathbb{T}^d} L \left(x, \frac{E_t}{m_t} \right) m_t(dx) dt,$$

under the constraint that (m, E) solves in the weak sense

$$(3.2) \quad \begin{aligned} \partial_t m + \operatorname{div}(E) &= 0 \text{ in } (0, 1) \times \mathbb{T}^d, \\ m(0) &= \mu \text{ and } m(1) = \nu, \end{aligned}$$

and where $L \left(x, \frac{E}{m} \right)$ is set to $+\infty$ as soon as $E \ll m$ is not satisfied. This (fruitful) approach allows to solve the problem of the singularity of the product αm . However, we claim that we can introduce another way to evaluate the cost of the trajectory given by the solution of (3.2), which we believe turns out to be simpler to interpret.

Our main idea consists in saying that different "controls" can give the same evolution of the state but should yield different costs. To illustrate this, consider the following

situation. The cost L is simply given by $L(x, p) = |p|^2$. The initial state and the terminal constraint are both equal to m_0 , the uniform probability measure on \mathbb{T}^d . Consider now the optimal control which consists in choosing $\alpha \equiv 0$. The associated cost is clearly 0. Consider now the inefficient and formal control which consists in assigning to each particles, or element of mass, a constant speed choosen uniformly in the ball $B(0_{\mathbb{R}^d}, 1)$. Clearly, by a symmetry argument, such a control is also admissible and also induces a constant state. However, it is very tempting to say that its cost is positive.

To make this heuristic more precise, we introduce a problem in which the "control" is now a measurable function $\psi : [0, 1] \rightarrow (\mathbb{T}^d \rightarrow \mathcal{P}(\mathbb{R}^d))$. The measure $\psi_s(x, dz)$ is then interpreted as the repartition of speed we provide to the elements of mass located at x at time s . We want to consider a "state" which is given as a solution of

$$(3.3) \quad \partial_t m_t + \operatorname{div}_x \left(\int_{\mathbb{R}^d} z m_t(dx) \psi_t(x, dz) \right) = 0 \text{ in } (0, 1) \times \mathbb{T}^d.$$

The previous equation is the natural PDE to characterize the density of particles evolving with a repartition of speed ψ . We then want to evaluate the cost of such a pair state/control by

$$(3.4) \quad \int_0^1 \int_{\mathbb{T}^d \times \mathbb{R}^d} L(x, z) \psi_t(x, dz) m_t(dx) dt.$$

Note that we can set $E_t(dx) := \int_{\mathbb{R}^d} z m_t(dx) \psi_t(x, dz)$, in which case (3.3) is of the form of (3.2). We can also set $\alpha_t(x) = \int_{\mathbb{R}^d} z \psi_t(x, dz)$ to realize that (3.3) has the exact form of the usual continuity equation. We are now ready to define properly the value functions.

The value function of the deterministic problem $U_{det} : [0, T] \times \mathcal{P}(\mathbb{T}^d)^2 \rightarrow \mathbb{R}$ is defined, for $t < T, \mu, \nu \in \mathcal{P}(\mathbb{T}^d)$, by

$$(3.5) \quad U_{det}(t, \mu, \nu) = \inf_{(\psi, m)} \int_t^T \int_{\mathbb{T}^d \times \mathbb{R}^d} L(x, z) \psi_t(x, dz) m_t(dx) dt,$$

where the infimum is taken over all pairs (ψ, m) such that

- $m \in \mathcal{C}([t, T], \mathcal{P}(\mathbb{T}^d))$, and $m_t = \mu, m_T = \nu$
- $\psi : [t, T] \times \mathbb{T}^d \rightarrow \mathcal{P}(\mathbb{R}^d)$ is a measurable map.
- The pair (ψ, m) satisfies (3.3) in the weak sense, i.e. for all $\varphi \in \mathcal{C}^1([t, T] \times \mathbb{T}^d, \mathbb{R})$

$$(3.6) \quad \int_{\mathbb{T}^d} \varphi(T, x) \nu(dx) - \int_{\mathbb{T}^d} \varphi(t, x) \mu(dx) = \int_t^T \int_{\mathbb{T}^d} (\partial_t \varphi(t, x) + \int_{\mathbb{R}^d} z \psi(x, dz) \cdot \nabla_x \varphi(t, x)) m_s(dx) ds.$$

We denote by $\mathcal{Adm}(t, \mu, \nu)$ the set of such pairs.

Concerning the value function of the stochastic optimal transport problem, recall that we have fixed a filtered probabilistic space $(\Omega, \mathcal{A}, (\mathcal{F}_t)_{t \geq 0}, \mathbb{P})$ and a Markovian, $\mathcal{P}(\mathbb{T}^d)$ valued process $(\nu_t)_{t \geq 0}$. The value function $U : [0, T] \times \mathcal{P}(\mathbb{T}^d)^2$ is defined for

$t < T, \mu, \nu \in \mathcal{P}(\mathbb{T}^d)$, by

$$(3.7) \quad U(t, \mu, \nu) = \inf_{(\psi, m)} \mathbb{E} \left[\int_t^T \int_{\mathbb{T}^d \times \mathbb{R}^d} L(x, z) \psi_t(x, dz) m_t(dx) dt \middle| \nu_t = \nu \right],$$

where the infimum is taken over all pairs of random variables $(\psi, m) : \Omega \rightarrow \mathcal{Adm}(t, \mu, \nu_T)$ which are adapted to the filtration $(\mathcal{F}_s)_{s \geq t}$.

Remark 3.1. *The definition is exactly similar in the case in which we can make a reduction of variable by replacing ν by w .*

3.2. The value function is a viscosity solution of the HJB equation. We now prove that the value functions defined just above are indeed viscosity solutions of the corresponding HJB equations. Since they are three equations at interest, we do not provide three complete proofs, but rather establish the property of viscosity super-solution in one case and the property of viscosity sub-solution in another one, leaving to the interested reader the adaptation of the proofs we provide.

Proposition 3.2. *Assume that the target process is given by $((\tau_{W_i})_{\#} \nu = \nu_t)_{t \geq 0}$ for W a \mathbb{T}^d valued solution of (1.28). Assume that the cost L grows at least like $|p|^k$ for some $k > 1$. Assume finally that U is locally bounded from below on $[0, T) \times \mathbb{T}^d \times \mathbb{T}^d$. Then, the function U_* is a viscosity super-solution of (1.33) in the sense of Definition 2.19.*

Proof. From the assumptions, U_* is well defined. Consider $(\theta, \psi^*, p, X) \in J^- U_*(t, \mu, w)$ and $(t_n, \mu_n, w_n)_{n \geq 0}$ converging toward (t, μ, w) such that $\lim_{n \rightarrow \infty} U(t_n, \mu_n, w_n) = U_*(t, \mu, w)$. We consider only the case $t < T$. We can write down

$$U(t_n, \mu_n, w_n) = \inf_{(\psi, m)} \mathbb{E} \left[\int_{t_n}^{t_n + \tau} \int_{\mathbb{T}^d \times \mathbb{R}^d} L(x, z) \psi_s(x, dz) m_s(dx) ds + U(t_n + \tau, m_{t_n + \tau}, w_{t_n + \tau}) \right],$$

where the infimum is taken over all $(\psi, m) \in \mathcal{Adm}(t_n, \mu_n, w_T)$. From this we deduce

$$(3.8) \quad U(t_n, \mu_n, w_n) \geq \inf_{(\psi, m)} \mathbb{E} \left[\int_{t_n}^{t_n + \tau} \int_{\mathbb{T}^d \times \mathbb{R}^d} L(x, z) \psi_s(x, dz) m_s(dx) ds + U_*(t_n + \tau, m_{t_n + \tau}, w_{t_n + \tau}) \right].$$

From the growth assumption we made on L , we deduce that we can restrict the infimum to pairs (ψ, m) such that, almost surely, for some $k > 1$,

$$(3.9) \quad \int_{t_n}^T \int_{\mathbb{T}^d} \int_{\mathbb{R}^d} |z|^k \psi_s(x, dz) m_s(dx) ds < \infty,$$

From Jensen's inequality, the previous yields that almost surely

$$(3.10) \quad \int_{t_n}^T \int_{\mathbb{T}^d} \left| \int_{\mathbb{R}^d} z \psi_s(x, dz) \right|^k m_s(dx) ds < \infty,$$

Thanks to Theorem 8.3.1 in Ambrosio et al. [4], this implies in particular that $(m_s)_{s \geq t_n}$ is an absolutely continuous curve, almost surely, and thanks to representation theorems such as Theorem 2.1 in Jimenez et al. [22], we obtain that, almost surely, there exists a \mathbb{T}^d valued random process $(X_s^n)_{s \in [t_n, T]}$ such that, for all $s \in [t_n, T]$ $\mathcal{L}(X_s) = m_s$ and for almost every $s \in [t_n, T]$,

$$(3.11) \quad dX_s^n = \int_{\mathbb{R}^d} z \psi_s(X_s^n, dz) ds.$$

Let us precise that, to avoid any difficulty one may have with the fact a filtered probability space is already used, the best construction, that we do not detail and which is standard, is to extend the already existing probability space so that $(X_s^n)_{s \in [t_n, T]}$ is a random variable defined for almost every $\omega \in \Omega$ (in the original probability space).

We are now going to use the process $(X_s^n)_{s \in [t_n, T]}$ to create a coupling between $m_{t_n+\tau}$ and $m_t = \mu$. If $t_n \leq t$ this does not raise an issue, however, in the opposite case, we have to be more precise. Consider first a couple of random variable (X, Z) such that $\mathcal{L}(X, Z)(dx, dz) = \mu(dx)\psi^*(x, dz)$. Recall that, from Definition 2.2, Z is bounded almost surely. It now suffices to remark that we can always choose the process $(X_s^n)_{s \in [t_n, T]}$ such that $\|X_{t_n}^n - X\|_\infty$ goes to 0 as $n \rightarrow \infty$.

Coming back to (3.8), and using the definition of a sub-differential, we can write

$$U(t_n, \mu_n, w_n) - U_*(t, \mu, w) \geq \tau \sigma^2(t) \text{Tr}(X) + \theta(t_n + \tau - t) + o(\tau) \\ + \inf_{(\psi, m)} \mathbb{E} \left[\int_{t_n}^{t_n+\tau} \int_{\mathbb{T}^d \times \mathbb{R}^d} L(x, z) \psi_s(x, dz) m_s(dx) ds + \mathbb{E}[Z \cdot (X_{t_n+\tau}^n - X)] + o(\sqrt{\mathbb{E}[|X_{t_n+\tau}^n - X|^2]}) \right].$$

This leads to

$$U(t_n, \mu_n, w_n) - U_*(t, \mu, w) \geq \tau \sigma^2(t) \text{Tr}(X) + o(t_n + \tau - t) \\ + \inf_{\psi, m} \mathbb{E} \left\{ \int_{t_n}^{t_n+\tau} \mathbb{E} \left[\int_{\mathbb{R}^d} L(X_s^n, z) \psi_s(X_s^n, dz) + Z \cdot \left(\int_{\mathbb{R}^d} z \psi_s(X_s^n, dz) \right) \right] ds \right. \\ \left. + \theta(t_n + \tau - t) + \mathbb{E}[Z \cdot (X_{t_n}^n - X)] + o(\sqrt{\mathbb{E}[|X_{t_n+\tau}^n - X|^2]}) \right\},$$

From which we obtain

$$U(t_n, \mu_n, w_n) - U_*(t, \mu, w) \geq \tau \sigma^2(t) \text{Tr}(X) + o(t_n + \tau - t) \\ + \inf_{\psi, m} \mathbb{E} \left\{ \int_{t_n}^{t_n+\tau} \mathbb{E}[-H(X_s^n, Z)] ds + \theta(t_n + \tau - t) + \mathbb{E}[Z \cdot (X_{t_n}^n - X)] + o(\sqrt{\mathbb{E}[|X_{t_n+\tau}^n - X|^2]}) \right\}.$$

Hence, sending $n \rightarrow \infty$, then dividing by τ and then letting $\tau \rightarrow 0$, we deduce the required result. \square

We now turn to the viscosity sub-solution property.

Proposition 3.3. *Assume that the target process $(\nu_t)_{t \geq 0}$ is constant and that U is locally bounded from above. Then, U is a viscosity sub-solution of (2.12) in the sense of Definition 2.8.*

Proof. By definition of U , it is usc. Consider $(\theta, \psi^*) \in \partial^+ U(t, \mu, \nu)$ with $t < T$. Consider a couple $(X_s, Z_s)_{s \in [t, T]}$ such that $\mathcal{L}(X_s, Z_s)(dx, dz) = m_s(dx)\psi^*(x, dz)$ and

$$(3.12) \quad X_s = X_t - \int_t^s \int_{\mathbb{R}^d} D_p H(X_u, Z_u) du.$$

Note that we can always choose the process $(X_s, Z_s)_{s \in [t, t+\tau]}$ so that almost surely, it is continuous in time. Set $m_s = \mathcal{L}(X_s)$ and remark that m is a solution of

$$(3.13) \quad \partial_t m - \text{div}_x \left(\int_{\mathbb{R}^d} D_p H(x, z) \psi^*(x, dz) m \right) = 0 \text{ in } (t, t + \tau) \times \mathbb{T}^d,$$

with initial condition μ .

Using this trajectory as an admissible one, we can now write

$$U(t, \mu, \nu) \leq \int_t^{t+\tau} \int_{\mathbb{T}^d \times \mathbb{R}^d} L(x, -D_p H(x, z)) \psi^*(x, dz) m_s(dx) ds + U(t + \tau, m_{t+\tau}, \nu).$$

Using that $(\theta, \psi) \in \partial^+ U(t, \mu, \nu)$, we obtain that

$$(3.14) \quad 0 \leq \int_t^{t+\tau} \int_{\mathbb{T}^d \times \mathbb{R}^d} L(x, -D_p H(x, z)) \psi^*(x, dz) m_s(dx) ds + \theta \tau + \mathbb{E}[Z_t \cdot (X_{t+\tau} - X_t)] \\ + o\left(\tau + \sqrt{\mathbb{E}[|X_{t+\tau} - X_t|^2]}\right).$$

By construction of H , we get that

$$(3.15) \quad 0 \leq \int_t^{t+\tau} \int_{\mathbb{T}^d \times \mathbb{R}^d} -H(x, z) + z \cdot D_p H(x, z) \psi^*(x, dz) m_s(dx) ds + \theta \tau \\ - \mathbb{E}[Z_t \cdot (\int_t^{t+\tau} D_p H(X_s, Z_s) ds)] + o\left(\tau + \sqrt{\mathbb{E}[|X_{t+\tau} - X_t|^2]}\right)$$

Dividing by τ and letting $\tau \rightarrow 0$, we obtain that

$$0 \leq - \int_{\mathbb{T}^d \times \mathbb{R}^d} H(x, z) \psi^*(x, dz) \mu(dx) + \theta + \mathbb{E}[Z \cdot D_p H(X_t, Z_t)] - \mathbb{E}[Z \cdot D_p H(X_t, Z_t)].$$

Hence the result follows. \square

3.3. Comparison with the more usual value functions. In this section, we want to compare the two following functions. The first one is

$$(3.16) \quad U(t, \mu) = \inf_{\psi, m} \int_t^T \int_{\mathbb{T}^d \times \mathbb{R}^d} L(x, z) \psi_s(x, dz) m_s(dx) ds,$$

where $t \in [t, T]$, $\mu \in \mathcal{P}(\mathbb{T}^d)$, and the infimum is taken over $\cup_{\nu \in \mathcal{P}(\mathbb{T}^d)} \mathcal{Adm}(t, \mu, \nu)$. In other words it is the value function defined in the way we just presented, by omitting the terminal condition. The second function we are interested in is

$$(3.17) \quad V(t, \mu) = \inf_{\alpha, m} \int_t^T \int_{\mathbb{T}^d} L(x, \alpha_s(dx)) m_s(dx) ds,$$

where $t \in [t, T]$, $\mu \in \mathcal{P}(\mathbb{T}^d)$, and the infimum is taken over all pairs (α, m) such that $m \in \mathcal{C}([t, T], \mathcal{P}(\mathbb{T}^d))$, $m_t = \mu$ and $\alpha : [t, T] \times \mathbb{T}^d \rightarrow \mathbb{R}^d$ is measurable and for all $\varphi \in \mathcal{C}^1([t, T] \times \mathbb{T}^d, \mathbb{R})$,

$$(3.18) \quad \int_{\mathbb{T}^d} \varphi(T, x) m_T(dx) - \int_{\mathbb{T}^d} \varphi(t, x) \mu(dx) = \int_t^T \int_{\mathbb{T}^d} (\partial_t \varphi(t, x) + \alpha(t, x) \cdot \nabla_x \varphi(t, x)) m_s(dx) ds.$$

In other words, V is the value function of optimal control problem on the set of measures as it is more usually defined.

We want to show that those two functions are equal. In order to do so, we are going to use the comparison principle Theorem 2.11 on the equation

$$(3.19) \quad -\partial_t U + \int_{\mathbb{T}^d} H(x, D_\mu U(t, x)) \mu(dx) = 0 \text{ in } (0, T) \times \mathcal{P}(\mathbb{T}^d).$$

We start by showing the following.

Proposition 3.4. *The function V is a viscosity super-solution of (3.19).*

Proof. From Gangbo and Tudorascu [19], we know that for any $t < T, \mu \in \mathcal{P}(\mathbb{T}^d), (\theta, \xi) \in \partial_{clas}^- V(t, \mu)$

$$(3.20) \quad -\theta + \int_{\mathbb{T}^d} H(x, \xi(x)) \mu(dx) \geq 0.$$

Hence taking $(\theta, \psi) \in \partial^- V(t, \mu)$, we deduce that

$$(3.21) \quad -\theta + \int_{\mathbb{T}^d} H(x, \int_{\mathbb{R}^d} z \psi(x, dz)) \mu(dx) \geq 0,$$

and the result follows from Jensen's inequality. \square

We now turn to the more subtle property of viscosity sub-solution. We are not able to establish it in all its generality, but nonetheless, the following result is sufficient to apply Theorem 2.11.

Proposition 3.5. *Assume that for all $x \in \mathbb{T}^d, \xi, p \in \mathbb{R}^d$ with $\xi \neq 0$*

$$(3.22) \quad \xi D_{pp} H(x, p) \xi > 0.$$

Take $\lambda > 0$ and consider $(\theta, \psi) \in \partial^+ V(t, \mu)$ such that ψ satisfies $\mu(dx) \psi(x, dz) = (\pi_1, \lambda(\pi_1 - \pi_2))_{\#} \gamma^o(dx, dz)$ for $\gamma^o \in \Pi^{opt}(\mu, \nu)$ for some $\nu \in \mathcal{P}(\mathbb{T}^d)$. Then

$$(3.23) \quad -\theta + \int_{\mathbb{T}^d \times \mathbb{R}^d} H(x, z) \psi(x, dz) \mu(dx) \leq 0,$$

Proof. Consider (X, Y) such that $\mathcal{L}(X, Y) = \gamma^o$ and define

$$(3.24) \quad F : \mathbb{T}^d \times \mathbb{R}^d \rightarrow \mathbb{R}^d, (x, p) \rightarrow -D_p H(x, \lambda p).$$

As it is done in Proposition 5.30 in Santambrogio [30] for instance, we want to show that $m_s := \mathcal{L}(X + sF(X, X - Y))$ satisfies a continuity equation for a certain drift, and that evaluating the cost of this drift will yields the required inequality.

Take $s > 0$ and $z \in \text{Supp}(m_s)$. By construction, there exists $(x, y) \in \text{Supp}(\gamma^o)$ such that $z = x + sF(x, x - y)$. We want to show that there exists at most one such couple $(x, y) \in \gamma^o$. Hence, take $(x, y), (x', y') \in \text{Supp}(\gamma^o)$ such that

$$(3.25) \quad x + sF(x, x - y) = x' + sF(x', x' - y').$$

Recall that H is convex and x, x', y, y' all belong to \mathbb{T}^d . Hence we can always assume that F is $C > 0$ Lipschitz continuous, we deduce that

$$(3.26) \quad |x - x'| = s|F(x', x' - y') - F(x, x - y)| \leq Cs(|x - x'| + |(x - y) - (x' - y')|).$$

Let us also write

$$(3.27) \quad x - x' = s(F(x', x' - y') - F(x', x - y)) + s(F(x', x - y) - F(x, x - y)).$$

Moreover, from (3.22), we obtain that $-F$ is $\beta > 0$ monotone in p for some $\beta > 0$. Taking the scalar product against $(x' - y') - (x - y)$, we deduce that

$$(3.28) \quad \langle x - x' + s(F(x, x - y) - F(x', x - y)), (x' - y') - (x - y) \rangle \leq -\beta |(x' - y') - (x - y)|^2.$$

Rearranging, we obtain that

$$(3.29) \quad \beta|(x' - y') - (x - y)|^2 \leq |x - x'|^2 + Cs|x - x'| \cdot |(x' - y') - (x - y)| - \langle x - x', y - y' \rangle.$$

Since $(x, y), (x', y') \in \text{Supp}(\gamma^o)$, we obtain that the last scalar product is non-negative and thus that

$$(3.30) \quad \beta|(x' - y') - (x - y)|^2 \leq |x - x'|^2 + Cs|x - x'| \cdot |(x' - y') - (x - y)|.$$

Plugging this estimate into (3.26), we finally deduce that

$$(3.31) \quad |x - x'| \leq Cs(|x - x'| + Cs \left(\frac{Cs + \sqrt{C^2s^2 + 4\beta}}{2\beta} \right) |x - x'|).$$

Hence for s sufficiently small (compared to a constant which depends only on C and β), we deduce that $2|x - x'| \leq |x - x'|$ and thus that $x = x'$ and hence that $y = y'$. This result of uniqueness has strong consequences. In particular, it allows to define two measurable maps $X_s(z)$ and $Y_s(z)$ such that for any $z \in \text{Supp}(m_s)$, $z = X_s(z) + sF(X_s(z), X_s(z) - Y_s(z))$. This implies in particular that (α, m) solve the continuity equation (at least in short time) with

$$\alpha_s(z) = -D_p H(X_s(z), \lambda(X_s(z) - Y_s(z))).$$

From this we deduce that

$$\begin{aligned} V(t, \mu) &\leq \int_t^{t+\delta} \int_{\mathbb{T}^d} L(z, \alpha_s(z)) m_s(dz) ds + V(t + \delta, m_{t+\delta}) \\ &\leq \int_t^{t+\delta} \int_{\mathbb{T}^d} H(X_s(z), \lambda(X_s(z) - Y_s(z))) - \lambda \alpha_s(z) \cdot (X_s(z) - Y_s(z)) m_s(dz) ds \\ &\quad + \int_t^{t+\delta} \int_{\mathbb{T}^d} L(z, \alpha_s(z)) - L(X_s(z), \alpha_s(z)) m_s(dz) ds + V(t + \delta, m_{t+\delta}). \\ &\leq \int_t^{t+\delta} \mathbb{E}[H(X, \lambda(X - Y)) + \lambda D_p H(X, \lambda(X - Y)) \cdot (X - Y)] ds \\ &\quad + C \int_t^{t+\delta} \mathbb{E}[s D_p H(X, \lambda(X - Y))] ds + V(t + \delta, m_{t+\delta}). \end{aligned}$$

Hence, we deduce that

$$(3.32) \quad \begin{aligned} 0 &\leq \int_t^{t+\delta} \mathbb{E}[H(X, \lambda(X - Y)) + \lambda D_p H(X, \lambda(X - Y)) \cdot (X - Y)] ds \\ &\quad + C \int_t^{t+\delta} \mathbb{E}[s D_p H(X, \lambda(X - Y))] ds + \theta \delta - \delta \mathbb{E}[\lambda(X - Y) \cdot D_p H(X, \lambda(X - Y))] + o(\delta). \end{aligned}$$

We then deduce the result by dividing by δ and taking the limit $\delta \rightarrow 0$. \square

As a consequence of the two previous Propositions as well as of Theorem 2.11, we immediately obtain the following.

Theorem 3.6. *Under the assumptions of Proposition 3.5, the two functions U and V are equal.*

4. BOUNDS ON THE VALUE OF THE STOCHASTIC OPTIMAL TRANSPORT PROBLEM NEAR THE SINGULARITY

In this section, we are interested in proving precise estimates on the behaviour of the value function U introduced in (3.7). We start by a comment on the modelling, we then proceed to study successively the cases in which the target process is of the jump process type and when it is driven by a diffusion.

4.1. On the choice of the cost functional. We explain here, on three simple examples, the effect of the cost functional on the type of behavior we may expect near $t = 0$, in the case of a deterministic problem. Such type of behaviors are well known in the optimal control theory and we shall pass through those examples quite rapidly.

4.1.1. *Cost functional with linear growth.* Assume that the cost function L is given by

$$(4.1) \quad L(x, \alpha) := |\alpha|.$$

In this context, if we are concerned with (1.15), a quick change of variable yields that for any $t \in [0, T]$

$$(4.2) \quad U(t, \mu, \nu) = U(T - t, \mu, \nu).$$

In other words, the cost is sufficiently low for high controls to allow the value to be bounded uniformly in time. In this situation, the state constraint is very easily achieved and there is no singularity at the terminal time.

If such situations may present interests in themselves, we believe that from a modeling perspective, they are not the most interesting ones as the problem of the controller does not get harder as the remaining time shortens. It is not even clear it depends on time. We do not detail it too much but in such situations the randomness of the final target somehow disappear as the we can just wait for the final time to reach instantly the final target.

4.1.2. *The case of bounded controls.* Somehow opposite to the previous situation is the case in which L is given by

$$(4.3) \quad L(x, \alpha) := \begin{cases} 0 & \text{if } |\alpha| \leq 1, \\ +\infty & \text{else.} \end{cases}$$

In this situation, the constrained optimization problem is not necessary controllable and the associated value can be infinite for $t \in [0, T)$. Indeed consider for instance $U(T - \epsilon, \delta_x, \delta_y)$ for $|x - y| > \epsilon$. If such cases present a lot of interest in themselves, they do not in this case in which the final density is constrained. Furthermore, if we were to replace this constraint with a bounded terminal cost, then the study of the associated HJB equation would be rather classical and fall in the scope of the previous section.

4.1.3. *Cost functionals which are powers of distances.* We consider here the cases in which L is given, for $k \geq 1$, by

$$(4.4) \quad L(x, \alpha) := |\alpha|^k.$$

A simple change of variable yields that in this situation, if we are still concerned with (1.15),

$$(4.5) \quad U(t, \mu, \nu) = \frac{U(T-1, \mu, \nu)}{(T-t)^{k-1}} = \frac{W_k(\mu, \nu)^k}{(T-t)^{k-1}}.$$

This type of behaviour is the one we are interested in, hence we shall make assumptions to control the cost function L with powers of α . Furthermore, in view of Alfonsi and Jourdain [2] (which focuses on the case $k = 2$), such a function U is not smooth in neither μ nor ν . This justifies in particular the use of the notion of viscosity solutions introduced in the Section 2.

4.2. Controllability of the stochastic problems and L^∞ bounds of the value functions. We now provide, by means of controllability bounds, estimates on the value functions for stochastic optimal transport problems, near the final time $t = T$. In the previous Section, we recalled that, as soon as the cost L satisfies for some $k \geq 1, C > 0$ and for all $x \in \mathbb{T}^d, p \in \mathbb{R}^d$

$$(4.6) \quad 0 \leq L(x, \alpha) \leq C(1 + |\alpha|^k),$$

then the value of the deterministic problem is bounded. In this Section, we explain how we can compare the value function U defined in (3.7) with the value function U_{det} of the deterministic problem, defined in (3.5). Define the function $\omega : [0, T] \rightarrow \mathbb{R}$ by

$$(4.7) \quad \omega(t) = \sup_{s \leq t, \mu, \nu} U_{det}(s, \mu, \nu).$$

In the two cases that follow, we are going to assume that for all $t \in [0, T], \omega(t) < \infty$ and that the infimum in (3.5) is always reached for $t < T$.

4.2.1. The case of jumps. Assume that the target process $(\nu_t)_{t \geq 0}$ is driven by jumps, which happen at Poisson times associated with the intensity $\lambda : [0, T] \rightarrow \mathbb{R}_+$, and which are described by the operator $\mathcal{T} : \mathcal{P}(\mathbb{T}^d) \rightarrow \mathcal{P}(\mathbb{T}^d)$. We can prove the following.

Proposition 4.1. *Assume that there exists $C > 0$ and $\gamma > -1$ such that for $T-t \leq C^{-1}$*

$$(4.8) \quad \lambda(t)\omega(t) \leq C(T-t)^\gamma,$$

where $\omega(t) = \sup_{s \geq t, \mu, \nu} U_{det}(s, \mu, \nu)$. Assume also that

$$(4.9) \quad \lambda(t) \leq C(T-t)^{-1} \int_t^T \lambda.$$

Then there exists a continuous function $\beta : [0, T] \rightarrow \mathbb{R}_+$, such that

$$(4.10) \quad \sup_{\mu, \nu \in \mathcal{P}(\mathbb{T}^d)} |U(t, \mu, \nu) - U_{det}(t, \mu, \nu)| \leq \beta(t) \xrightarrow{t \rightarrow T} 0.$$

Remark 4.2. *We comment the hypotheses of the result.*

- *The assumption (4.9) is purely technical, it is verified by any function such that $\lambda(t) \sim C(T-t)^\alpha$ for any $C > 0, \alpha > -1$. However it is not automatically verified, as for instance $\lambda(t) = \frac{d}{dt}(e^{-(T-t)^{-2}})$ does not verify it. We do not know whether this can be removed.*

- The requirement (4.8) is quite important in our proof. This assumption yields an integrability condition on the product $\lambda\omega$. Such an integrability condition is crucial. Note for example that if λ is constant, then we require (among other things), that $\omega \in L^1_{loc}$, which is not the case for a quadratic cost. We show an example of a situation where bounds on U does not exist if this integrability fails.

Remark 4.3. Note that no assumption is made on \mathcal{T} in this result, in particular, the result still holds true if \mathcal{T} depends also on t and μ . This is due to the fact that $\mathcal{P}(\mathbb{T}^d)$ is compact. If the problem was to transport elements of $\mathcal{P}(\mathbb{R}^d)$, then some assumptions should be made on \mathcal{T} , namely on its growth.

Proof. Notice first that if $\lambda = 0$ in $L^1((0, \kappa), \mathbb{R}_+)$ for some $\kappa > 0$, then the results holds true trivially. Hence, we focus here on the case

$$(4.11) \quad \forall t > 0, \int_0^t \lambda(s) ds > 0.$$

Consider the problem starting in $\mu \in \mathcal{P}(\mathbb{T}^d)$ at time t and where the target process is equal to ν at t . Let n be the (random) number of jumps in $[t, T]$ and consider the sequence $\tau_0 = t < \tau_1 < \tau_2 < \dots < \tau_n \leq T$ of random times at which the target process jumps. Note that this sequence is finite almost surely, possibly empty and that the event $\{\tau_n = T\}$ shall be ignored since it happens with probability 0. Consider the random pairs $(\psi_s, m_s)_{s \in [t, T]}$, given by

- For $1 \leq i \leq n$, $s \in (t_{i-1}, t_i)$, (ψ_s, m_s) is equal (up to a change of time) to a minimum in the problem $U_{det}(t_{i-1}, m_{t_{i-1}}, \mathcal{T}^{i-1}\nu)$.
- For $s \in (t_n, T)$, (ψ_s, m_s) is given through a minimum in $U_{det}(t_n, m_{t_n}, \mathcal{T}^n\nu)$.

Such a pair is clearly admissible. We now estimate its cost.

$$(4.12) \quad \begin{aligned} & \mathbb{E} \left[\int_t^T \int_{\mathbb{T}^d \times \mathbb{R}^d} L(x, z) \psi_s(x, dz) m_s(dx) ds \right] \\ &= \mathbb{P}(n > 0) \mathbb{E} \left[\int_{t_n}^T \int_{\mathbb{T}^d \times \mathbb{R}^d} L(x, z) \psi_s(x, dz) \mu_s(dx) ds \right. \\ & \quad \left. + \sum_{i=0}^{n-1} \int_{t_i}^{t_{i+1}} \int_{\mathbb{T}^d \times \mathbb{R}^d} L(x, z) \psi_s(x, dz) \mu_s(dx) ds \middle| n > 0 \right] + \mathbb{P}(n = 0) U_{det}(t, \mu, \nu) \\ &\leq \mathbb{P}(n > 0) \mathbb{E} \left[U_{det}(t_n, m_{t_n}, \mathcal{T}^n\nu) + \sum_{i=0}^{n-1} U_{det}(t_i, m_{t_i}, \mathcal{T}^i\nu) \middle| n > 0 \right] + \mathbb{P}(n = 0) U_{det}(t, \mu, \nu) \\ &\leq \mathbb{P}(n > 0) \mathbb{E} \left[\omega(t_n) + \sum_{i=0}^{n-1} \omega(t_i) \middle| n > 0 \right] + \mathbb{P}(n = 0) U_{det}(t, \mu, \nu) \\ &\leq \mathbb{P}(n > 0) \mathbb{E} [(1 + n)\omega(t_n) | n > 0] + \mathbb{P}(n = 0) U_{det}(t, \mu, \nu). \end{aligned}$$

We can now compute

$$(4.13) \quad \mathbb{E} [(1 + n)\omega(t_n) | n > 0] = \sum_{k=1}^{\infty} \mathbb{E} [(1 + n)\omega(t_n) | n = k] \mathbb{P}(n = k | n > 0).$$

Since the $(\tau_n)_{n \geq 0}$ are given by a Poisson process, we have that

$$(4.14) \quad \mathbb{P}(n > 0)\mathbb{P}(n = k | n > 0) = \mathbb{P}(n = k) = \frac{\left(\int_t^T \lambda(s) ds\right)^k}{k!} e^{-\int_t^T \lambda(s) ds},$$

and also that there exists $C > 0$ such that for any $1 \leq k < n$, the law of τ_k conditioned on τ_{k-1} has a density which is bounded by

$$(4.15) \quad s \rightarrow \mathbb{1}_{s \geq \tau_{k-1}} C \frac{\lambda(s)}{\int_{\tau_{k-1}}^T \lambda}.$$

Hence, we can estimate

$$(4.16) \quad \begin{aligned} \mathbb{E}[(1+n)\omega(t_n) | n = k] &\leq C^k (k+1) \int_t^T \int_{t_1}^T \cdots \int_{t_{k-1}}^T \omega(t_k) \lambda(t_k) \frac{dt_k}{\int_{t_{k-1}}^T \lambda} \cdots \frac{\lambda(t_2) dt_2}{\int_{t_1}^T \lambda} \frac{\lambda(t_1) dt_1}{\int_t^T \lambda} \\ &\leq C^k (k+1) \int_t^T \int_{t_1}^T \cdots \int_{t_{k-2}}^T (T - t_{k-1})^{\gamma+1} \frac{\lambda(t_{k-1}) dt_{k-1}}{\int_{t_{k-1}}^T \lambda} \cdots \frac{\lambda(t_2) dt_2}{\int_{t_1}^T \lambda} \frac{\lambda(t_1) dt_1}{\int_t^T \lambda} \\ &\leq \frac{C^k (k+1) (T-t)^{\gamma+1}}{\int_t^T \lambda(s) ds} \end{aligned}$$

From the previous estimate, we deduce that

$$(4.17) \quad \begin{aligned} \mathbb{P}(n > 0) \mathbb{E}[(1+n)\omega(t_n) | n > 0] &\leq C \sum_{k=1}^{\infty} \frac{C^k (k+1)}{k!} e^{-\int_t^T \lambda} \left(\int_t^T \lambda\right)^{k-1} (T-t)^{\gamma+1} \\ &\leq C (T-t)^{\gamma+1}. \end{aligned}$$

We can compute

$$(4.18) \quad (1 - \mathbb{P}(n = 0)) U_{\det}(t, \mu, \nu) \leq C \int_t^T \omega(s) \lambda(s) ds \leq C (T-t)^{\gamma+1}.$$

Hence, setting $\beta(t) = C(T-t)t^{\gamma+1}$, we deduce that

$$(4.19) \quad U(t, \mu, \nu) \leq \mathbb{E} \left[\int_t^T \int_{\mathbb{T}^d \times \mathbb{R}^d} L(x, z) \psi_s(x, dz) m_s(dx) ds \right] \leq U_{\det}(t, \mu, \nu) + \beta(t).$$

Obtaining the lower bound is easier. Indeed, for $\epsilon > 0$, consider an ϵ optimal pair (ψ, m) (which exists since the value is bounded from above). It then follows that

$$\begin{aligned}
(4.20) \quad U(t, \mu, \nu) &\geq \mathbb{E} \left[\int_t^T \int_{\mathbb{T}^d \times \mathbb{R}^d} L(x, z) \psi_s(x, dz) m_s(dx) ds \right] - \epsilon \\
&= \mathbb{P}(n = 0) \mathbb{E} \left[\int_t^T \int_{\mathbb{T}^d \times \mathbb{R}^d} L(x, z) \psi_s(x, dz) m_s(dx) ds \middle| n = 0 \right] - \epsilon \\
&\quad + \mathbb{P}(n > 0) \mathbb{E} \left[\int_t^T \int_{\mathbb{T}^d \times \mathbb{R}^d} L(x, z) \psi_s(x, dz) m_s(dx) ds \middle| n > 0 \right] \\
&\geq \mathbb{P}(n = 0) \mathbb{E} \left[\int_t^T \int_{\mathbb{T}^d \times \mathbb{R}^d} L(x, z) \psi_s(x, dz) m_s(dx) ds \middle| n = 0 \right] - \epsilon \\
&\geq \mathbb{P}(n = 0) U_{det}(t, \mu, \nu) - \epsilon.
\end{aligned}$$

Since ϵ is arbitrary, we deduce that the inequality also holds for $\epsilon = 0$. Hence, we deduce the lower bound

$$(4.21) \quad U(t, \mu, \nu) \geq U_{det}(t, \mu, \nu) - \beta(t)$$

following the same computation as in the part concerning the upper bound. \square

The previous result yields in fact more than just bounds on the value function U . It gives a precise behaviour of the value function near $t = T$. It states that, under the standing assumptions, it behaves as U_{det} near $t = T$.

If the assumptions of the previous Theorem are not satisfied, then we can be in an entirely different situation. Indeed, consider the following example.

Example 4.4. Assume $L(x, \alpha) = |\alpha|^2$, λ is a constant and $\mathcal{T}\nu \neq \nu$ for some $\nu \in \mathcal{P}(\mathbb{T}^d)$. In this context, $\omega(t) = C(T - t)^{-1}$. Consider a time $t > 0$ and assume that at this time, the target process is equal to ν . By conditioning on the number of jumps occurring in the remaining time, we obtain that for any admissible pair (ψ, m) , which is progressively measurable by the filtration generated by $(\nu_s)_{s \in [t, T]}$

$$\begin{aligned}
(4.22) \quad U(t, \mu, \nu) &\geq \mathbb{P}(n = 1) \mathbb{E} \left[\int_t^T \int_{\mathbb{T}^d \times \mathbb{R}^d} |z|^2 \psi_s(x, dz) m_s(dx) ds \middle| n = 1 \right] \\
&\quad + \mathbb{P}(n = 0) \mathbb{E} \left[\int_t^T \int_{\mathbb{T}^d \times \mathbb{R}^d} |z|^2 \psi_s(x, dz) m_s(dx) ds \middle| n = 0 \right].
\end{aligned}$$

Let us denote by ρ the density of the law of the jump τ_1 , conditioned on $\{n = 1\}$. Consider $(\mu_s)_{s \in [t, T]}$, the trajectory in the event $\{n = 0\}$. By definition of the 2-Wasserstein distance, we obtain that for any $t' \in [t, T)$

$$\begin{aligned}
(4.23) \quad \mathbb{E} \left[\int_t^T \int_{\mathbb{T}^d \times \mathbb{R}^d} |z|^2 \psi_s(x, dz) m_s(dx) ds \middle| n = 0 \right] &\geq \mathbb{E} \left[\int_{t'}^T \int_{\mathbb{T}^d \times \mathbb{R}^d} |z|^2 \psi_s(x, dz) m_s(dx) ds \middle| n = 0 \right] \\
&\geq \frac{W_2^2(\mu_{t'}, \nu)}{T - t'}
\end{aligned}$$

We then compute

$$(4.24) \quad \mathbb{E} \left[\int_t^T \int_{\mathbb{T}^d \times \mathbb{R}^d} |z|^2 \psi_s(x, dz) m_s(dx) ds \middle| n = 1 \right] \geq \int_t^T \frac{W_2^2(\mu_s, \mathcal{T}\nu)}{T-s} \rho(s) ds$$

Integrating (4.23) with respect to ρ , we deduce that

$$(4.25) \quad \begin{aligned} U(t, \mu, \nu) &\geq \mathbb{P}(n = 0) \int_t^T \frac{W_2^2(\mu_s, \nu)}{T-s} \rho(s) ds \\ &\quad + \mathbb{P}(n = 1) \int_t^T \frac{W_2^2(\mu_s, \mathcal{T}\nu)}{T-s} \rho(s) ds \\ &\geq \mathbb{P}(n = 1) \int_t^T \frac{W_2^2(\mu_s, \nu) + W_2^2(\mu_s, \mathcal{T}\nu)}{T-s} \rho(s) ds, \end{aligned}$$

if t is sufficiently small so that $\mathbb{P}(n = 1) \leq \mathbb{P}(n = 0)$. The right hand side of the previous inequality is equal to $+\infty$ since $\nu \neq \mathcal{T}\nu$. Hence for any $\mu \in \mathcal{P}(\mathbb{T}^d)$, $U(t, \mu, \nu) = +\infty$.

This last example hints that there is a strong dichotomy : either U is infinite in all the points ν such that $\mathcal{T}\nu \neq \nu$, or either it behaves quite similarly as U_{det} .

4.2.2. *The case of the target pushed by a diffusion.* Consider now that $(\nu_t)_{t \geq 0}$ is given by $\nu_t = (\tau_{W_t})_{\#} \nu$ for $\nu \in \mathcal{P}(\mathbb{T}^d)$ and $(W_t)_{t \geq 0}$ the strong solution of

$$(4.26) \quad dW_t = \sigma(t) dB_t \text{ for } t > 0,$$

with initial condition $W_0 = 0$, where $\sigma : [0, T) \rightarrow \mathbb{R}$ is a smooth bounded function and $(B_t)_{t \geq 0}$ is a standard Brownian motion. We also assume that there exists $C > 0$ such that for all $x \in \mathbb{T}^d, \alpha \in \mathbb{R}^d$

$$(4.27) \quad L(x, \alpha) \leq C(1 + |\alpha|^2).$$

Also recall that we are here interested in the value function U as a function of t, μ and $w \in \mathbb{T}^d$.

We start with the following Lemma.

Lemma 4.5. *Almost surely, there exists a unique \mathbb{T}^d valued solution $(X_t)_{t \leq T}$ of*

$$(4.28) \quad dX_t = \frac{W_t - X_t}{T-t} dt,$$

given an initial condition $X_0 \in \mathbb{T}^d$. Almost surely, it satisfies $X_t \rightarrow W_T$ as $t \rightarrow T$.

Moreover,

$$(4.29) \quad \mathbb{E} \left[\int_0^T \left| \frac{dX_t}{dt} \right|^2 ds \right] = \frac{\mathbb{E}[|W_0 - X_0|^2]}{T} + \int_0^T \frac{\sigma(s)^2}{T-s} ds,$$

which possibly reads $+\infty = +\infty$.

Proof. Let us first remark that the existence and uniqueness of the solution on $[0, T)$ is trivial. Hence we only need to show that the limit holds as $t \rightarrow T$. Remark now that (4.28) can be written

$$(4.30) \quad dX_t = \frac{W_T - X_t}{T-t} dt + \frac{W_t - W_T}{T-t} dt.$$

The previous leads to

$$\frac{d|X_t - W_T|^2}{dt} = -2\frac{|X_t - W_T|^2}{T-t} + 2\frac{(W_t - W_T) \cdot (X_t - W_T)}{T-t}.$$

Integrating this relation yields

$$|X_t - W_T|^2 + 2 \int_0^t \frac{|X_s - W_T|^2}{T-s} ds = 2 \int_0^t \frac{W_t - W_T}{\sqrt{T-s}} \cdot \frac{X_s - W_T}{\sqrt{T-s}} ds + |X_0 - W_T|^2.$$

From the regularity property of the Brownian motion, more precisely that, almost surely, for t and s sufficiently close,

$$(4.31) \quad \forall c > 1, |B_t - B_s| \leq c\sqrt{2|t-s|\log(|t-s|^{-1})},$$

we deduce that there exists $C > 0$ independent of t such that, almost surely,

$$(4.32) \quad |X_t - W_T|^2 + 2 \int_0^t \frac{|X_s - W_T|^2}{T-s} ds \leq C.$$

Hence the first part of the results follows.

Let us now remark that

$$d\frac{W_t - X_t}{T-t} = \frac{dW_t}{T-t} = \frac{\sigma(t)dB_t}{T-t}.$$

Hence, we deduce that

$$\frac{dX_t}{dt} = \frac{W_t - X_t}{T-t} = \frac{W_0 - X_0}{T} + \int_0^t \frac{\sigma(s)}{T-s} dB_s.$$

From which follows

$$\mathbb{E} \left[\int_0^T \left| \frac{dX_t}{dt} \right|^2 ds \right] = \frac{\mathbb{E}[|W_0 - X_0|^2]}{T} + \int_0^T \int_0^t \frac{\sigma(s)^2}{(T-s)^2} ds dt,$$

which yields the result. \square

We can now prove a controllability estimate.

Proposition 4.6. *Assume that there exists $C \geq 0$ such that L satisfies for any x, α*

$$(4.33) \quad L(x, \alpha) \leq C(1 + |\alpha|^2),$$

and that

$$(4.34) \quad \sigma(t) \sim K(T-t)^\gamma \text{ as } t \rightarrow T,$$

for some $K, \gamma > 0$. Then for all $t < T, \mu \in \mathcal{P}(\mathbb{T}^d), w \in \mathbb{T}^d, U(t, \mu, w) < \infty$.

Proof. Consider $(t, \mu, w) \in (0, \infty) \times \mathcal{P}(\mathbb{T}^d) \times \mathbb{T}^d$. Take $\delta > 0$ such that $t + \delta < T$ and define $T_0 = T - (t + \delta)$. Denote by $(W_s)_{s \geq t}$ the strong solution of (1.28) such that $W_t = w$. Consider the control α which transports optimally, according to the quadratic cost, μ into ν in a time δ . Consider the process $(X_s)_{s \geq t+\delta}$ defined by

$$dX_s = \frac{W_s - X_s}{T-s} ds,$$

with initial condition $X_{t+\delta} = 0$. We can now build an admissible control by setting $\psi_s(x, dz) = \delta_{X_s}$. Thanks to Lemma 4.5, this control is admissible. Moreover, thanks to the same result, we have the trivial estimate

$$U(t, \mu, w) \leq C(T-t) + C \frac{W_2^2(\mu, \nu)}{\delta} + C \frac{\mathbb{E}[|w_{t+\delta}|^2]}{T_0} + \int_{t+\delta}^T \frac{\sigma(s)^2}{T-s} ds,$$

which is finite thanks to (4.34). \square

Remark 4.7. *As in the case of jumps, the quadratic case seems to be a limit case for which controllability requires $\sigma(t) \rightarrow 0$ as $t \rightarrow T$. Moreover, similar result can be obtained for outside of this quadratic assumption, namely by using the Burkholder-Davis-Gundy inequalities.*

We now prove a refinement of the previous estimate which yields a more precise result for the behaviour of U near $t = T$.

Proposition 4.8. *Under the assumptions and notations of Proposition 4.6, assume furthermore that $\gamma > \frac{1}{2}$, and that for any $t < s < T$*

$$(4.35) \quad U_{det}(s, \mu, \nu) - U_{det}(t, \mu, \nu) \leq (T-s)^{-2}(s-t).$$

Then, it holds that

$$(4.36) \quad U(t, \mu, w) \leq U_{det}(t, \mu, (\tau_w)_{\#}\nu) + \beta(t),$$

for a continuous function β such that $\beta(t) \rightarrow 0$ as $t \rightarrow T$. Moreover, if L is convex in α we always have

$$(4.37) \quad U_{det}(t, \mu, (\tau_w)_{\#}\nu) \leq U(t, \mu, w).$$

Proof. The proof is similar to the previous one. Consider $(t, \mu, w) \in (0, \infty) \times \mathcal{P}(\mathbb{T}^d) \times \mathbb{T}^d$. Take $\delta(t)$ such that $t + \delta(t) < T$ and define $T_0(t) = T - (t + \delta(t))$. The function δ is to be chosen later on.

Denote by $(W_s)_{s \geq t}$ the strong solution of (1.28) such that $W_t = w$. Consider an optimal trajectory, according to the deterministic problem, which transports μ into $(\tau_w)_{\#}\nu$ in a time $\delta(t)$. Consider the process $(X_s)_{s \geq t+\delta(t)}$ defined by

$$dX_s = \frac{W_s - X_s}{T-s} ds,$$

with initial condition $X_{t+\delta(t)} = w$. Consider now the control which consists in playing the first trajectory in the time $\delta(t)$ and then to translate the state with speed $\frac{dX_s}{ds}$ afterwards. Thanks to Lemma 4.5, this control is admissible. Moreover, thanks to the same result, we can estimate

$$(4.38) \quad U(t, \mu, w) \leq U_{det}(T - \delta(t), \mu, (\tau_w)_{\#}\nu) + C \frac{\mathbb{E}[|w_{t+\delta(t)} - w|^2]}{T - (t + \delta(t))} + \int_{t+\delta(t)}^T \frac{\sigma(s)^2}{T-s} ds.$$

Remark that

$$\mathbb{E}[|w_{t+\delta(t)} - w|^2] = \int_t^{t+\delta(t)} \sigma^2(s) ds \leq K(T-t)^{2\gamma} \delta(t),$$

where the inequality holds for t close to T . Let us now set $\delta(t) = (T - t) - (T - t)^\theta$ for some $\theta > 1$. Coming back to (4.38), we obtain

$$(4.39) \quad U(t, \mu, w) \leq U_{det}(t + (T - t)^\theta, \mu, (\tau_w)_{\#}\nu) + C \frac{K(T - t)^{2\gamma} \delta(t)}{(T - t)^\theta} + \int_{t+\delta(t)}^T \frac{\sigma(s)^2}{T - s} ds$$

The last term vanishes as soon as $\gamma > 0$, the second to last term vanishes as soon as $1 + 2\gamma > \theta$. It then remains to estimate the first term of the right hand side. Namely, we are interested in the difference

$$U_{det}(t + (T - t)^\theta, \mu, (\tau_w)_{\#}\nu) - U_{det}(t, \mu, (\tau_w)_{\#}\nu).$$

From the assumption we made on $\partial_t U_{det}$, we can bound this difference by $C(T - t)^{\theta-2}$ and we deduce finally that

$$U(t, \mu, w) \leq U_{det}(t, \mu, (\tau_w)_{\#}\nu) + C(T - t)^{\theta-2} + C \frac{K(T - t)^{2\gamma} \delta(t)}{(T - t)^\theta} + \int_{t+\delta(t)}^T \frac{\sigma(s)^2}{T - s} ds,$$

which yields the required estimate when $\theta \in (2, 1 + 2\gamma)$.

Finally, let us remark that the estimate $U(t, \mu, w) \geq U_{det}(t, \mu, (\tau_w)_{\#}\nu)$ is simply a consequence of the convexity of L in α . Indeed, fix $t < T, \mu \in \mathcal{P}(\mathbb{T}^d), w \in \mathbb{T}^d$ and consider any (stochastic) admissible pair (ψ, m) . As we already mentioned at the beginning of Section 3, we can always consider an associated pair of measures (m, E) solution of

$$\partial_t m + \operatorname{div}(E) = 0 \text{ in } (t, T) \times \mathbb{T}^d,$$

which satisfies $m_t = \mu, m_T = (\tau_{W_T})_{\#}\nu$. Defining $M_t = \mathbb{E}[m_t]$ and $K_t = \mathbb{E}[E_t]$, we obtain that

$$\partial_t M + \operatorname{div}(K) = 0 \text{ in } (t, T) \times \mathbb{T}^d,$$

together with $M_t = \mu, M_T = (\tau_w)_{\#}\nu$. From this we deduce that

$$(4.40) \quad \begin{aligned} \mathbb{E} \left[\int_t^T \int_{\mathbb{T}^d \times \mathbb{R}^d} L(x, z) \psi_s(x, dz) m_s(dx) ds \right] &\geq \mathbb{E} \left[\int_t^T \int_{\mathbb{T}^d} L \left(x, \frac{E_s}{m_s} \right) m_s(dx) ds \right] \\ &\geq \int_t^T \int_{\mathbb{T}^d} L \left(x, \frac{K_s}{M_s} \right) m_s(dx) ds \\ &\geq U_{det}(t, \mu, (\tau_w)_{\#}\nu), \end{aligned}$$

from which the result follows. □

Remark 4.9. *Although it is not the main objective of the present work, it would be interesting to lower the assumption $\gamma > \frac{1}{2}$ into $\gamma > 0$ so that this result is similar to the one we present for the case of jumps, namely that the controllability directly yields a precise behaviour of U near $t = T$.*

Remark 4.10. *The assumption (4.35) is a bound on the time derivative of U_{det} . It is verified in the case $L(x, p) = |p|^2$ for instance. Since L does not depend explicitly in t , it can be verified by a change of variable in time as soon $D_\alpha L$ has linear growth in $|\alpha|$.*

5. CONCLUSION AND PERSPECTIVES

5.1. Summary of the results. In conclusion, we summarize the results that we have brought on the value function U defined in (3.7).

Assume that the cost L satisfies

$$(5.1) \quad L(x, \alpha) \geq C(|\alpha|^k - 1),$$

for some $k > 1$ and that it is such that the Hamiltonian H is globally continuous and such that there exists $C > 0$ such that

$$(5.2) \quad \forall x, y \in \mathbb{T}^d, p \in \mathbb{R}^d, |H(x, p) - H(y, p)| \leq C(1 + |p|)|x - y|.$$

Assume that the value U_{det} of the deterministic problem satisfies

$$(5.3) \quad \forall t \in [0, T], \mu, \nu \in \mathcal{P}(\mathbb{T}^d), U_{det}(t, \mu, \nu) < \infty.$$

Theorem 5.1. *Assume that the target process $(\nu_t)_{t \geq 0}$ is a jump process described by the operator \mathcal{T} and the intensity λ . Assume that λ is continuous, bounded and satisfies the assumptions of Proposition 4.1. Then U is the unique viscosity solution of (1.26) such that*

$$(5.4) \quad \limsup_{t \rightarrow T} \sup_{\mu, \nu} |U(t, \mu, \nu) - U_{det}(t, \mu, \nu)| = 0.$$

Theorem 5.2. *Assume that the target process $(\nu_t)_{t \geq 0}$ is given by $\nu_t = (\tau_{W_t})_{\#} \nu$ where $(W_t)_{t \geq 0}$ is the strong solution of (1.28). Assume that σ satisfies the assumptions of Proposition 4.8. Assume also that L satisfies (4.33) and that H satisfies (2.35). Then U is the unique viscosity solution of (1.33) which satisfies*

$$(5.5) \quad \limsup_{t \rightarrow T} \sup_{\mu, w} |U(t, \mu, w) - U_{det}(t, \mu, (\tau_w)_{\#} \nu)| = 0.$$

Remark 5.3. *In particular L has exactly a quadratic growth here. As we mentioned above, the case in which it grows faster can be treated similarly by using the Burkholder-Davis-Gundy inequalities in the study of the singularity.*

5.2. Potential approximating schemes for the HJB equation. Discretizing both time and the space of measures, one can arrive at usual discrete scheme for dynamic optimal control problem, namely using the notion of Wasserstein barycenters [1]. It seems that the stability properties of viscosity solutions should be helpful to prove some convergence properties. On the other hand, fast methods to compute Wasserstein barycenters now exist [17] and could lead to a tractable numerical treatment of the problem.

5.3. More general optimal control problem. The techniques developed in Section 2 seem to be well-suited to study more general optimal control problems on the space of probability measures. With Pierre-Louis Lions (Collège de France), we are currently generalizing them to treat the case of the control of the parabolic continuity equation

$$(5.6) \quad \partial_t m - \nu \Delta m + \operatorname{div}(\alpha m) = 0 \text{ in } (0, T) \times \mathbb{T}^d,$$

where the control is still α , but the presence of the term in ν makes the analysis more complex.

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