# FINITE DIMENSIONAL APPROXIMATIONS OF HAMILTON-JACOBI-BELLMAN EQUATIONS IN SPACES OF PROBABILITY MEASURES 

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#### Abstract

We prove that viscosity solutions of Hamilton-Jacobi-Bellman (HJB) equations corresponding either to deterministic optimal control problems for systems of $n$ particles or to stochastic optimal control problems for systems of $n$ particles with a common noise converge locally uniformly to the viscosity solution of a limiting HJB equation in the space of probability measures. We prove uniform continuity estimates for viscosity solutions of the approximating problems which may be of independent interest. We pay special attention to the case when the Hamiltonian is convex in the gradient variable and equations are of first order and provide a representation formula for the solution of the limiting first order HJB equation. We also propose an intrinsic definition of viscosity solution on the Wasserstein space.


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## 1. Introduction

We consider the problem of approximation of Hamilton-Jacobi-Bellman (HJB) equations in spaces of probability measures by equations in finite dimensional spaces. More precisely, we study if appropriately interpreted viscosity solution of such HJB equations can be approximated locally uniformly by viscosity solutions of finite dimensional problems. This is related to the problem of whether value functions of variational or optimal control problems in spaces of probability measures can be approximated by value functions corresponding to problems for finite particle systems. Similar convergence problems have been recently studied in the context of mean field games [15, 16, 18, 20, 21, 28, 37, 41, 50, 53]. In particular it was proved in [18] that classical solutions of finite dimensional second order Nash systems converge, in a suitable sense, to classical solutions of the corresponding master equations. Also convergence of functionals of empirical measures of the marginal laws of particle systems for McKean-Vlasov stochastic differential equations was studied recently in [20, 23] using calculus in the space of measures, stochastic analysis and partial differential equations the space of measures. Explicit convergence estimates were obtained in [20, 23]. The problems investigated there are different from the one here. They studied the case of independent noises and no controls so they dealt with partial differential equations which are linear, have slightly different form and have smooth solutions. We refer to the references in [20, 23] for the discussion of other earlier results in this direction. Regularity and convergence problems for finite dimensional approximations of first order HJB equations in spaces of probability measures were studied in [37, 39, 41, 52], either when solutions were regular or when the Hamiltonian was quadratic. Some results are also mentioned in [15] without proofs, while some results may be considered to be part of a folklore of the theory. In this paper we want to investigate the problem rigorously from the point of view of viscosity solutions.

We will be concerned with first and second order degenerate HJB equations of the form

$$
\left\{\begin{array}{l}
\partial_{t} \mathcal{U}-\kappa \Delta_{w} \mathcal{U}+\mathcal{H}\left(\mu, \mu, \nabla_{\mu} \mathcal{U}\right)+\mathcal{F}(\mu)=0 \quad \text { in }(0, T) \times \mathcal{P}_{2}\left(\mathbb{R}^{d}\right)  \tag{1.1}\\
\mathcal{U}(0, \mu)=\mathcal{U}_{0}(\mu) \quad \text { on } \mathcal{P}_{2}\left(\mathbb{R}^{d}\right),
\end{array}\right.
$$

where $\Delta_{w} \mathcal{U}$ is the partial Laplacian of $\mathcal{U}$, see [24], $T>0, \kappa \geq 0, \mathcal{P}_{2}\left(\mathbb{R}^{d}\right)$ is the Wasserstein space of probability measures on $\mathbb{R}^{d}$ with bounded second moments, $\mathcal{F}: \mathcal{P}_{2}\left(\mathbb{R}^{d}\right) \rightarrow \mathbb{R}^{d}$ and, for $\mu_{1}, \mu_{2} \in \mathcal{P}_{2}\left(\mathbb{R}^{d}\right), \xi \in L_{\mu_{1}}^{2}\left(\mathbb{R}^{d} ; \mathbb{R}^{d}\right)$, the Hamiltonian $\mathcal{H}$ is defined by

$$
\mathcal{H}\left(\mu_{1}, \mu_{2}, \xi\right)=\int_{\mathbb{R}^{d}} H\left(x, \mu_{2}, \xi(x)\right) \mu_{1}(d x)
$$

for some function $H: \mathbb{R}^{d} \times \mathcal{P}_{2}\left(\mathbb{R}^{d}\right) \times \mathbb{R}^{d} \rightarrow \mathbb{R}$ which satisfies Hypothesis 2.1 below. We bring to the reader's attention the fact that the dependence of $\mathcal{H}$ in (1.1) on the first $\mu$ is linear, whereas the dependence on the second $\mu$ comes from the integrand function $H$. Moreover, the function $\mathcal{F}$ cannot be simply absorbed into $\mathcal{H}$ : otherwise, the restriction of $\mu$ to averages of Dirac masses would not yield $\sqrt{1.2}$ below. Indeed, part of the goal of this paper is to justify that the approximating finite dimensional problems should have the form

$$
\left\{\begin{align*}
& \partial_{t} u_{n}-\kappa \operatorname{Tr}\left(A_{n} D^{2} u_{n}\right)+\frac{1}{n} \sum_{i=1}^{n} H\left(x_{i}, \frac{1}{n-1} \sum_{j \neq i}^{n} \delta_{x_{j}}, n D_{x_{i}} u_{n}\right)  \tag{1.2}\\
&+\mathcal{F}\left(\frac{1}{n} \sum_{i=1}^{n} \delta_{x_{i}}\right)=0 \\
& \text { in }(0, T) \times\left(\mathbb{R}^{d}\right)^{n}, \\
& u_{n}\left(0, x_{1}, \cdots, x_{n}\right)= \mathcal{U}_{0}\left(\frac{1}{n} \sum_{i=1}^{n} \delta_{x_{i}}\right) \quad \text { on }\left(\mathbb{R}^{d}\right)^{n},
\end{align*}\right.
$$

where for $n \in \mathbb{N}, A_{n}$ is the $n d \times n d$ matrix composed of $n^{2}$ block matrices $I_{d}$.
Equation (1.1) will be interpreted in the $L^{2}$ sense, that is, we will look at the "lifted" version of (1.1) in the space

$$
E:=L^{2}\left(\Omega ; \mathbb{R}^{d}\right),
$$

where $\Omega$ is an atomless probability space. This technique was introduced in [15, 50] and its detailed exposition and recent developments can be found in [20, 21, 43]. Without loss of generality we can assume that $\Omega=(0,1)$ with the standard Lebesgue measure $\mathcal{L}_{1}$. We denote by $\langle\cdot, \cdot\rangle_{1}$ the inner product in $L^{2}(\Omega ; \mathbb{R})$ and, for $X, Y \in L^{2}\left(\Omega ; \mathbb{R}^{d}\right)$, we set

$$
\begin{gathered}
\langle X, Y\rangle_{d}:=\left(\left\langle X_{1}, Y_{1}\right\rangle_{1}, \cdots,\left\langle X_{d}, Y_{d}\right\rangle_{1}\right), \\
\langle X, Y\rangle:=\left\langle X_{1}, Y_{1}\right\rangle_{1}+\cdots+\left\langle X_{d}, Y_{d}\right\rangle_{1} \quad \text { and } \quad|X|=\sqrt{\langle X, X\rangle}
\end{gathered}
$$

where $X_{1}, \ldots, X_{d}, Y_{1}, \ldots, Y_{d}$ are the components of $X$ and $Y$ respectively. We denote the canonical basis in $\mathbb{R}^{d}$ by $\left\{e_{1}, \cdots, e_{k}\right\}$ and consider its elements as constant functions in $E$. We define the functions $U_{0}, F: E \rightarrow \mathbb{R}$ by

$$
U_{0}(X)=\mathcal{U}_{0}\left(X_{\sharp} \mathcal{L}_{1}\right), \quad F(X)=\mathcal{F}\left(X_{\sharp} \mathcal{L}_{1}\right),
$$

where $\sharp$ denotes pushforward. Thus $X_{\sharp} \mathcal{L}_{1}$ is the law of the random vector $X$, and is an element of $\mathcal{P}_{2}\left(\mathbb{R}^{d}\right)$. If $U: E \rightarrow \mathbb{R}$ is twice differentiable and such that $\mathcal{U}(\mu)=U(X)$ if $\mu$ is the law of $X$ then we have the crucial formula

$$
\Delta_{w} \mathcal{U}(\mu) \circ X=\sum_{k=1}^{d}\left\langle D^{2} U(X) e_{k}, e_{k}\right\rangle ;
$$

see Section 5 below.
For $X, P \in E, \mu \in \mathcal{P}_{2}\left(\mathbb{R}^{d}\right)$, we define

$$
\tilde{H}(X, \mu, P):=\int_{\Omega} H(X(\omega), \mu, P(\omega)) d \omega \text {. }
$$

We rely on the theory of viscosity solutions in Hilbert spaces and consider functions $U$ : $[0, T) \times E \rightarrow \mathbb{R}$ which are viscosity solutions of

$$
\left\{\begin{array}{l}
\partial_{t} U-\kappa \sum_{k=1}^{d}\left\langle D^{2} U e_{k}, e_{k}\right\rangle+\tilde{H}\left(X, X_{\sharp} \mathcal{L}_{1}, D U\right)+F(X)=0 \quad \text { in }(0, T) \times E  \tag{1.3}\\
U(0, X)=U_{0}(X) \quad \text { on } E .
\end{array}\right.
$$

Here, $D U, D^{2} U$ stand for the Fréchet derivatives of $U$ with respect to the $X$ variable. We refer the readers to [29] for the theory of viscosity solutions in Hilbert spaces and extensive references.

We have the following definition.
Definition 1.1. Let $\mathcal{U}:[0, T) \times \mathcal{P}_{2}\left(\mathbb{R}^{d}\right) \rightarrow \mathbb{R}$ and define $U:[0, T) \times E \rightarrow \mathbb{R}$ by $U(t, X)=$ $\mathcal{U}\left(t, X_{\sharp} \mathcal{L}_{1}\right)$.
(i) We say that $\mathcal{U}$ is an L-viscosity subsolution of (1.1) on the Wasserstein space if $U$ is a viscosity subsolution of $(1.3)$ on $[0, T) \times E$.
(ii) We say that $\mathcal{U}$ is an L-viscosity supersolution of (1.1) on the Wasserstein space if $U$ is a viscosity supersolution of (1.3) on $[0, T) \times E$.
(ii) When $\mathcal{U}$ is both an L-viscosity subsolution and an L-viscosity supersolution of (1.1) on the Wasserstein space, we say that it is an L-viscosity solution of (1.1) on the Wasserstein space.

We remark that another definition of viscosity solution to HJB master equations in the Wasserstein space, also called an $L$-viscosity solution, was introduced in 59. The definition in [59] uses the framework of path dependent PDE and is not related to Definition 1.1 here, even though both notions have the same name.

We also propose an intrinsic definition of viscosity solution on the Wasserstein space and show in Section 5 that the notion of $L$-viscosity solution provides a way to select particular intrinsic viscosity solutions. Only when the Hamiltonian is convex in the momentum variables and $\kappa=0$ it is known that the notions of $L$-viscosity solution and intrinsic viscosity solution are equivalent 43].

The main result of the manuscript is the following convergence theorem.
Theorem 1.2. Let Hypothesis 2.1 be satisfied and let $\kappa \geq 0$. Suppose that for $n \geq 1$ the functions $u_{n}:[0, T) \times\left(\mathbb{R}^{d}\right)^{n} \rightarrow \mathbb{R}$ are the viscosity solutions of (1.2) Then, for every bounded set $B$ in $\mathcal{P}_{2}\left(\mathbb{R}^{d}\right)$,

$$
\begin{aligned}
\lim _{n \rightarrow \infty} \sup & \left\{\left|u_{n}\left(t, x_{1}, \cdots, x_{n}\right)-\mathcal{U}\left(t, \frac{1}{n} \sum_{i=1}^{n} \delta_{x_{i}}\right)\right|:\right. \\
& \left.\left(t, x_{1}, \cdots, x_{n}\right) \in[0, T) \times\left(\mathbb{R}^{d}\right)^{n}, \frac{1}{n} \sum_{i=1}^{n} \delta_{x_{i}} \in B\right\}=0,
\end{aligned}
$$

where $\mathcal{U}$ is the unique L-viscosity solution of (1.1) on the Wasserstein space.
The assumptions of Hypothesis 2.1 will be introduced in Section 2. To prove Theorem 1.2 we first obtain appropriate uniform continuity estimates for the solutions $u_{n}$ of 1.2 . This is done in Theorem 3.3 for a more general case when the second order coefficients $A_{n}$ may depend on $x$. Theorem 3.3 is the main technical result of the paper and is of independent interest. We then convert the functions $u_{n}$ into functions of empirical measures by defining new functions

$$
\mathcal{V}_{n}\left(t, \mu_{\mathbf{x}}\right):=u_{n}(t, \mathbf{x}), \quad \text { where } \mu_{\mathbf{x}}:=\frac{1}{n} \sum_{i=1}^{n} \delta_{x_{i}}, \mathbf{x}=\left(x_{1}, \cdots, x_{n}\right)
$$

which are well defined since the functions $u_{n}$ are invariant with respect to permutations of the variables of $\mathbf{x}$. The estimates of Theorem 3.3 guarantee that $\mathcal{V}_{n}$ are uniformly continuous in the topology of $[0, T] \times \mathcal{P}_{r}\left(\mathbb{R}^{d}\right)$, where $1<r<2$. We then extend $\mathcal{V}_{n}$ to $[0, T] \times \mathcal{P}_{2}\left(\mathbb{R}^{d}\right)$ preserving its modulus of continuity and then use the Arzelà-Ascoli theorem to pass to the limit, along a subsequence, to a function $\mathcal{V}$ defined on $[0, T] \times \mathcal{P}_{2}\left(\mathbb{R}^{d}\right)$. We then prove directly that its "lifted" version $V:[0, T] \times E \rightarrow \mathbb{R}$ is a viscosity solution of (1.3). Uniqueness of viscosity solutions of (1.3) then guarantees that the whole sequence $\mathcal{V}_{n}$ converges to $\mathcal{V}$. Thus we completely avoid
dealing with equation (1.1) in the space of probability measures which may not have a unique viscosity solution in the sense of [38] (see [43]). In Section 6] we show that if $\kappa=0$ and $H$ is convex in the gradient variable then the functions $u_{n}$, which are value functions of optimal control problems for $n$-particle systems, converge to the value function of a variational problem in $\mathcal{P}_{2}\left(\mathbb{R}^{d}\right)$, thus giving a representation formula for the solution of (1.1). Finally we prove a few technical results in the Appendix.

Equations (1.2) correspond either to deterministic optimal control problems for systems of $n$ particles or to stochastic optimal control problems for systems of $n$ particles with a common noise. Theorem 1.2 solves the problem of convergence for a large class of general first order HJB equations, even though the identification of the limit as a value function is only obtained for the convex case and $\kappa=0$. However, using the methods of this paper we were not able to obtain a result similar to Theorem 1.2 for other stochastic particle systems, for instance for systems of $n$ particles with non-constant diffusion coefficients, in which case the matrices $A_{n}$ are functions like in Hypothesis 3.1, or for systems of $n$ particles with independent noises, in which case $A_{n}=I_{n d}$. We also remark that some assumptions of Hypothesis 2.1 could be changed or relaxed while some may pose a bigger problem. This is worth investigating. The main challenge is in proving uniform continuity estimates of Theorem 3.3. We do not consider other cases here as Hypothesis 2.1 is sufficiently general and we do not want to overburden the presentation with too many technicalities. Our main goal is to convey the basic ideas. The readers can explore various generalizations.

Hamilton-Jacobi-Bellman equations and master equations for mean field games or mean field control problems in spaces of probability measures have been studied a lot in recent years using various approaches. We refer the readers to [6, 7, 8, 9, 10, 11, 14, 15, 16, 17, 18, 19, 20, 21, 22, 24, [28, 37, 38, 39, 41, 43, 45, 46, 52, 53, 54, 57, 58, 59]. Equations related to differential games were studied in [25, 47], equations related to control problems with partial observation were studied in [4, 5] and HJB equations mostly related to large deviations and fluid dynamics problems were investigated with slightly different techniques in [30, 31, 32, 33, 34]. In particular an abstract method of relaxed-limits for viscosity solutions was introduced in [31] for applications in large deviations and this technique was recently generalized in [48, 49]. HJB equations in metric spaces were studied by various techniques in [1, 12, 13, 39, 40, 44, 51, [55, 56]. Finally we refer the readers to [26, 35] for an introduction to the theory of viscosity solutions of partial differential equations in $\mathbb{R}^{d}$ and to [29] and the references there for the overview of the theory of viscosity solutions of second order HJB equations in Hilbert spaces.

## 2. Notation, assumptions and definitions

We denote by $\mathcal{P}_{r}\left(\mathbb{R}^{d}\right), r \geq 1$ the space of Borel probability measures on $\mathbb{R}^{d}$ with finite $r$-th moments, equipped with the Wasserstein $r$-metric

$$
d_{r}(\mu, \nu):=\inf _{\gamma \in \Gamma(\mu, \nu)}\left(\int_{\mathbb{R}^{d} \times \mathbb{R}^{d}}|x-y|^{r} \gamma(d x, d y)\right)^{\frac{1}{r}}
$$

where $\Gamma(\mu, \nu)$ is the set of all Borel probability measures $\gamma$ on $\mathbb{R}^{d} \times \mathbb{R}^{d}$ with marginals $\mu, \nu$. The set of optimal measures in $\Gamma(\mu, \nu)$ will be denoted by $\Gamma_{0}(\mu, \nu)$. When $\mathbf{x}=\left(x_{1}, \ldots, x_{n}\right) \in\left(\mathbb{R}^{d}\right)^{n}$ we set

$$
\mu_{\mathbf{x}}:=\frac{1}{n} \sum_{i=1}^{n} \delta_{x_{i}}, \quad \text { and } \quad|\mathbf{x}|_{r}=\frac{1}{n^{1 / r}}\left(\sum_{i=1}^{n}\left|x_{i}\right|^{r}\right)^{1 / r}
$$

We have

$$
d_{r}\left(\mu_{\mathbf{x}}, \mu_{\mathbf{y}}\right)=\inf _{\sigma}\left|\mathbf{x}-\mathbf{y}_{\sigma}\right|_{r},
$$

where the infimum is taken over all permutations $\sigma$ of $\{1, \ldots, n\}$ and $\mathbf{y}_{\sigma}=\left(y_{\sigma(1)}, \ldots, y_{\sigma(n)}\right)$.

We use $\mathcal{L}_{1}$ to denote the Lebesgue measure on $\mathbb{R}$. If $X \in L^{r}\left(\Omega ; \mathbb{R}^{d}\right)$, then its $L^{r}$ norm will be denoted by $|X|_{r}$. If $r=2$ we will just write $|X|$ as in the introduction. For $X \in L^{r}\left(\Omega ; \mathbb{R}^{d}\right)$, $\operatorname{law}(X):=X_{\sharp} \mathcal{L}_{1}$ denotes the measure in $\mathcal{P}_{r}\left(\mathbb{R}^{d}\right)$ which is the push forward of $\mathcal{L}_{1}$ by $X$.

When $\mu \in \mathcal{P}_{2}\left(\mathbb{R}^{d}\right)$, we denote as $L_{\mu}^{2}\left(\mathbb{R}^{d} ; \mathbb{R}^{d}\right)$ the set of Borel vector fields $\xi: \mathbb{R}^{d} \rightarrow \mathbb{R}^{d}$ such that $\int_{\mathbb{R}^{d}}|\xi|^{2} \mu(d x)<\infty$. The tangent space at $\mu$, denoted by $T_{\mu} \mathcal{P}_{2}\left(\mathbb{R}^{d}\right)$, is the closure of $\nabla C_{c}^{\infty}\left(\mathbb{R}^{d}\right)$ in $L_{\mu}^{2}\left(\mathbb{R}^{d} ; \mathbb{R}^{d}\right)$.

For $x \in \mathbb{R}^{m}$ we will be also using the notation $|x|$ to denote the standard Euclidean norm in $\mathbb{R}^{m}$ and we will write $x \cdot y$ for $x, y \in \mathbb{R}^{m}$ to denote the dot product in $\mathbb{R}^{m}$.

If $A$ is a matrix or a bounded operator in a Hilbert space, we will write $\|A\|$ to denote the operator norm of $A$. We denote by $S(m)$ the set of $m \times m$ symmetric matrices. If $A \in S(m)$, $\operatorname{Tr}(A)$ means the trace of $A$.

For an open set $\mathscr{O} \subset \mathbb{R}^{m}$, we will write $C^{1}(\mathscr{O}), C^{2}(\mathscr{O})$ for the standard spaces of once and twice continuously differentiable functions on $\mathscr{O}$.

If $W$ is a Hilbert space, we denote by $C^{1,2}((0, T) \times W)$ the space of functions $\varphi:(0, T) \times W \rightarrow$ $\mathbb{R}$, such that $\partial_{t} \varphi, D \varphi, D^{2} \varphi$ are continuous on $(0, T) \times W$, where $D \varphi, D^{2} \varphi$ stand for the Fréchet derivatives of $\varphi$ with respect to the Hilbert space variable.

Throughout the paper we will always identify a Hilbert space with its dual. Thus, with this identification, $D \varphi:(0, T) \times W \rightarrow W$ and $D^{2} \varphi:(0, T) \times W \rightarrow S(W)$, where $S(W)$ is the space of bounded self-adjoint operators in $W$.

We make the following assumptions about the Hamiltonian function $H$.
Hypothesis 2.1. Let $1<r<2$.
(i) The function $H: \mathbb{R}^{d} \times \mathcal{P}_{r}\left(\mathbb{R}^{d}\right) \times \mathbb{R}^{d} \rightarrow \mathbb{R}$ is such that

$$
\begin{equation*}
|H(x, \nu, p)-H(x, \nu, q)| \leq C(1+|p|+|q|)|p-q| \quad \forall p, q, x \in \mathbb{R}^{d}, \nu \in \mathcal{P}_{r}\left(\mathbb{R}^{d}\right) \tag{2.1}
\end{equation*}
$$

$$
\begin{equation*}
|H(x, \mu, p)-H(y, \nu, p)| \leq \sigma\left(\left(|x-y|+d_{r}(\mu, \nu)\right)(1+|p|)\right) \quad \forall p, x, y \in \mathbb{R}^{d}, \mu, \nu \in \mathcal{P}_{r}\left(\mathbb{R}^{d}\right) \tag{2.2}
\end{equation*}
$$

for some concave modulus of continuity $\sigma$ and

$$
\begin{equation*}
|H(x, \mu, p)| \leq C\left(1+|p|^{2}\right) \quad \forall p, x \in \mathbb{R}^{d}, \mu \in \mathcal{P}_{r}\left(\mathbb{R}^{d}\right) \tag{2.3}
\end{equation*}
$$

(ii) The functions $\mathcal{U}_{0}, \mathcal{F} \in U C_{b}\left(\mathcal{P}_{r}\left(\mathbb{R}^{d}\right)\right.$ ) (the space of bounded and uniformly continuous functions).
We notice that it easily follows from (2.1) and (2.3) that for all $X, P, Q \in E$,

$$
\begin{equation*}
|\tilde{H}(X, \mu, P)-\tilde{H}(X, \mu, Q)| \leq C(1+|P|+|Q|)|P-Q| \tag{2.4}
\end{equation*}
$$

and

$$
\begin{equation*}
|\tilde{H}(X, \mu, P)| \leq C\left(1+|P|^{2}\right) \tag{2.5}
\end{equation*}
$$

Moreover, by the concavity of $\sigma$ and Jensen's inequality, we obtain that for all $X, Y, P \in E$, $\mu, \nu \in \mathcal{P}_{r}\left(\mathbb{R}^{d}\right)$,

$$
\begin{align*}
|\tilde{H}(X, \mu, P)-\tilde{H}(Y, \nu, P)| & \leq \int_{\Omega} \sigma\left(\left(|X(\omega)-Y(\omega)|+d_{r}(\mu, \nu)\right)(1+|P(\omega)|)\right) d \omega \\
& \leq \sigma\left(\int_{\Omega}\left(|X(\omega)-Y(\omega)|+d_{r}(\mu, \nu)\right)(1+|P(\omega)|) d \omega\right) \\
& \leq \sigma\left(\left(|X-Y|+d_{r}(\mu, \nu)\right)(1+|P|)\right) \tag{2.6}
\end{align*}
$$

Let $m_{1}$ be a modulus of continuity for $\mathcal{U}_{0}$ and $\mathcal{F}$. Since

$$
\left|U_{0}(X)-U_{0}(Y)\right|=\left|\mathcal{U}_{0}\left(X_{\sharp} \mathcal{L}_{1}\right)-\mathcal{U}_{0}\left(Y_{\sharp} \mathcal{L}_{1}\right)\right| \leq m_{1}\left(d_{r}\left(X_{\sharp} \mathcal{L}_{1}, Y_{\sharp} \mathcal{L}_{1}\right)\right) \leq m_{1}\left(|X-Y|_{r}\right)
$$

we conclude

$$
\begin{equation*}
\left|U_{0}(X)-U_{0}(Y)\right| \leq m_{1}(|X-Y|) \tag{2.7}
\end{equation*}
$$

Similarly,

$$
\begin{equation*}
|F(X)-F(Y)| \leq m_{1}(|X-Y|) . \tag{2.8}
\end{equation*}
$$

Moreover, for $\mathbf{x}, \mathbf{y} \in\left(\mathbb{R}^{d}\right)^{n}$,

$$
\begin{array}{r}
\left|\mathcal{U}_{0}\left(\mu_{\mathbf{x}}\right)-\mathcal{U}_{0}\left(\mu_{\mathbf{y}}\right)\right| \leq m_{1}\left(|\mathbf{x}-\mathbf{y}|_{r}\right), \\
\left|\mathcal{F}\left(\mu_{\mathbf{x}}\right)-\mathcal{F}\left(\mu_{\mathbf{y}}\right)\right| \leq m_{1}\left(|\mathbf{x}-\mathbf{y}|_{r}\right) . \tag{2.10}
\end{array}
$$

Let $W$ be a real Hilbert space with an inner product $\langle\cdot, \cdot\rangle$ and the norm $|\cdot|$. We recall the definitions of parabolic second-order jets (see [26], Section 8 and [27], Section 3).

Given $u:(0, T) \times W \rightarrow \mathbb{R}$ and $(\bar{t}, \bar{x}) \in(0, T) \times W$, the parabolic second-order superjet of $u$ at $(\bar{t}, \bar{x}), \mathcal{P}^{2,+} u(\bar{t}, \bar{x})$ is defined by

$$
\begin{align*}
\mathcal{P}^{2,+} u(\bar{t}, \bar{x}) & :=\{(a, p, S) \in \mathbb{R} \times W \times S(W): \\
& \left.\limsup _{(s, y) \rightarrow(, \bar{x})} \frac{u(s, y)-u(\bar{t}, \bar{x})-a(s-\bar{t})-\langle p, y-\bar{x}\rangle-\frac{1}{2}\langle S(y-\bar{x}), y-\bar{x}\rangle}{|s-\bar{t}|+|y-\bar{x}|^{2}} \leq 0\right\} . \tag{2.11}
\end{align*}
$$

The parabolic second-order subjet of $u$ at $(\bar{t}, \bar{x}), \mathcal{P}^{2,-} u(\bar{t}, \bar{x})$, is defined by reversing the inequality and replacing limsup by liminf in (2.11). The closure $\overline{\mathcal{P}}^{2,+} u(\bar{t}, \bar{x})$ of $\mathcal{P}^{2,+} u(\bar{t}, \bar{x})$, is defined as follows.

$$
\begin{aligned}
& \overline{\mathcal{P}}^{2,+} u(\bar{t}, \bar{x}):=\left\{(a, p, S) \in \mathbb{R} \times W \times S(W): \text { there exist }\left(t_{n}, x_{n}\right)\right. \text { and } \\
& \left.\left(a_{n}, p_{n}, S_{n}\right) \in \mathcal{P}^{2,+} u\left(t_{n}, x_{n}\right) \text { s.t. }\left(t_{n}, x_{n}, u\left(t_{n}, x_{n}\right), p_{n}, S_{n}\right) \rightarrow(\bar{t}, \bar{x}, u(\bar{t}, \bar{x}), a, p, S)\right\},
\end{aligned}
$$

The closure $\overline{\mathcal{P}}^{2,-} u(\bar{t}, \bar{x})$ of $\mathcal{P}^{2,-} u(\bar{t}, \bar{x})$ is defined similarly. We recall the definition of viscosity solution of an equation

$$
\left\{\begin{array}{l}
\partial_{t} u+G\left(t, x, u, D u, D^{2} u\right)=0 \quad \text { in }(0, T) \times W  \tag{2.12}\\
u(0, x)=g(x) \text { on } W,
\end{array}\right.
$$

where $G:(0, T) \times W \times \mathbb{R} \times W \times S(W) \rightarrow \mathbb{R}$ is continuous.
Definition 2.2. An upper semicontinuous function $u:[0, T) \times W \rightarrow \mathbb{R}$ is a viscosity subsolution of (2.12) if $u(0, x) \leq g(x)$ on $W$ and whenever $\varphi \in C^{1,2}((0, T) \times W)$ and $u-\varphi$ has a local maximum at $(t, x) \in(0, T) \times W$, then

$$
\partial_{t} \varphi(t, x)+G\left(t, x, u(t, x), D \varphi(t, x), D^{2} \varphi(t, x)\right) \leq 0 .
$$

A lower semicontinuous function $u:[0, T) \times W \rightarrow \mathbb{R}$ is a viscosity supersolution of (2.12) if $u(0, x) \geq g(x)$ on $W$ and whenever $\varphi \in C^{1,2}((0, T) \times W)$ and $u-\varphi$ has a local minimum at $(t, x) \in(0, T) \times W$, then

$$
\partial_{t} \varphi(t, x)+G\left(t, x, u(t, x), D \varphi(t, x), D^{2} \varphi(t, x)\right) \geq 0 .
$$

A function $u$ is a viscosity solution of (2.12) if it is a viscosity subsolution of (2.12) and a viscosity supersolution of (2.12).

It is easy to see that $(a, p, S) \in \mathcal{P}^{2,+} u(\bar{t}, \bar{x})$ (respectively, $\left.(a, p, S) \in \mathcal{P}^{2,-} u(\bar{t}, \bar{x})\right)$ if and only if there exists $\varphi \in C^{1,2}((0, T) \times W)$ such that $u-\varphi$ has a local maximum at $(\bar{t}, \bar{x})$ (respectively, $u-\varphi$ has a local minimum at $(\bar{t}, \bar{x}))$ and

$$
a=\partial_{t} \varphi(\bar{t}, \bar{x}), \quad p=D \varphi(\bar{t}, \bar{x}), \quad S=D^{2} \varphi(\bar{t}, \bar{x})
$$

The proof when $W=\mathbb{R}^{m}$ is in [35], Lemma 4.1, and it easily generalizes to the case of an infinite dimensional Hilbert space. Thus, since $G$ is continuous, Definition 2.2 is equivalent to the definition using the closures of parabolic jets.
Proposition 2.3. An upper semicontinuous function $u:[0, T) \times W \rightarrow \mathbb{R}$ is a viscosity subsolution of (2.12) if $u(0, x) \leq g(x)$ on $W$ and

$$
a+G(t, x, u(t, x), p, S) \leq 0 \quad \text { for all }(t, x) \in(0, T) \times W \text { and }(a, p, S) \in \overline{\mathcal{P}}^{2,+} u(t, x) .
$$

A lower semicontinuous function $u:[0, T) \times W \rightarrow \mathbb{R}$ is a viscosity supersolution of (2.12) if $u(0, x) \geq g(x)$ on $W$ and

$$
a+G(t, x, u(t, x), p, S) \geq 0 \quad \text { for all }(t, x) \in(0, T) \times W \text { and }(a, p, S) \in \overline{\mathcal{P}}^{2,-} u(t, x) .
$$

Remark 2.4. If equation 2.12) is of first order, that is if $G:(0, T) \times W \times \mathbb{R} \times W \rightarrow \mathbb{R}$ then the test functions $\varphi$ in Definition 2.2 are replaced by functions $C^{1}((0, T) \times W)$ and the parabolic second-order superjet and subjets of $u, \mathcal{P}^{2,+} u$ and $\mathcal{P}^{2,-} u$ are replaced by the first order superdifferentials and subdifferentials $D^{+} u$ and $D^{-} u$ respectively. Proposition 2.3 is then still true if $\overline{\mathcal{P}}^{2,+} u(t, x)$ and $\overline{\mathcal{P}}^{2,-} u(t, x)$ are replaced by the closures $\bar{D}^{+} u(t, x)$ and $\bar{D}^{-} u(t, x)$ which are defined similarly as the closures of the parabolic jets.

## 3. Estimates for finite dimensional equations

In this section we will consider a more general version of equations 1.2, with second order coefficients $A_{n}$ depending on $\mathbf{x}$ or being more general constant matrices.
Hypothesis 3.1. For $\mathbf{x}=\left(x_{1}, \ldots, x_{n}\right)$, let $A_{n}(\mathbf{x}), n=1,2, \ldots$, the $n d \times n d$ matrix composed of $n^{2}$ block matrices $a\left(x_{i}\right) a^{*}\left(x_{j}\right), i, j=1,2, \ldots, n$ such that the function $a: \mathbb{R}^{d} \rightarrow S(d)$ is bounded and there exists $L \geq 0$ such that

$$
\begin{equation*}
\|a(x)-a(y)\| \leq L|x-y| \quad \text { for all } x, y \in \mathbb{R}^{d} . \tag{3.1}
\end{equation*}
$$

The proof of the main theorem of this section uses the following simplified version of a well known lemma (cf. Lemma 3.80 of [29]).

Lemma 3.2. Let $\delta>0$, and let $\sigma_{1}$ be a modulus of continuity. Then there exist a nondecreasing, concave, $C^{2}$ function $\varphi_{\delta}$ on $[0,+\infty)$ such that $\varphi_{\delta}(0)<\delta$ and

$$
\begin{equation*}
\sigma_{1}\left(\varphi_{\delta}^{\prime}(s) s+s\right) \leq \varphi_{\delta}(s) \quad \text { for } 0 \leq s \leq 2 \tag{3.2}
\end{equation*}
$$

Theorem 3.3. Let Hypothesis 2.1 be satisfied. Let $A_{n}(\mathbf{x})$ satisfy Hypothesis 3.1, or, let $A_{n}(\mathbf{x})=$ $A_{n}$, where $A_{n}$ is any sequence of $n d \times n d$ symmetric matrices with constant coefficients such that $A_{n} \geq 0$ and $\operatorname{Tr}\left(A_{n}\right) \leq \tilde{L} n$ for some $\tilde{L} \geq 0$. Suppose that for $n \geq 1$ the functions $u_{n}$ : $[0, T] \times\left(\mathbb{R}^{d}\right)^{n} \rightarrow \mathbb{R}$ are the viscosity solutions of

$$
\left\{\begin{align*}
& \partial_{t} u_{n}-\kappa \operatorname{Tr}\left(A_{n} D^{2} u_{n}\right)+\frac{1}{n} \sum_{i=1}^{n} H\left(x_{i}, \frac{1}{n-1} \sum_{j \neq i}^{n} \delta_{x_{j}}, n D_{x_{i}} u_{n}\right)  \tag{3.3}\\
&+\mathcal{F}\left(\frac{1}{n} \sum_{i=1}^{n} \delta_{x_{i}}\right)=0 \text { in }(0, T) \times\left(\mathbb{R}^{d}\right)^{n}, \\
& u_{n}\left(0, x_{1}, \cdots, x_{n}\right)=\mathcal{U}_{0}\left(\frac{1}{n} \sum_{i=1}^{n} \delta_{x_{i}}\right) \text { on }\left(\mathbb{R}^{d}\right)^{n},
\end{align*}\right.
$$

Then there exists a modulus of continuity $\rho$ such that for every $n$,

$$
\begin{equation*}
\left|u_{n}(t, \mathbf{x})-u_{n}(s, \mathbf{y})\right| \leq \rho\left(|t-s|+|\mathbf{x}-\mathbf{y}|_{r}\right) \quad \forall t, s \in[0, T], \mathbf{x}, \mathbf{y} \in\left(\mathbb{R}^{d}\right)^{n} . \tag{3.4}
\end{equation*}
$$

Proof. We first note that if $u$ is a bounded viscosity subsolution of (3.3) and $v$ is a bounded viscosity supersolution of (3.3) and if we replace $u_{n}(t, \mathbf{x})-u_{n}(t, \mathbf{y})$ by $u(t, \mathbf{x})-v(t, \mathbf{y})$ in the proof of continuity estimates below, the same arguments work and we obtain that $u(t, \mathbf{x})-v(t, \mathbf{x}) \leq 0$ for all $(t, \mathbf{x})$. Thus the comparison theorem holds for bounded viscosity subsolutions and bounded viscosity supersolutions of equation (3.3). Moreover, since the function

$$
\mathbf{x} \longmapsto g(\mathbf{x})=\mathcal{U}_{0}\left(\frac{1}{n} \sum_{i=1}^{n} \delta_{x_{i}}\right)
$$

is bounded and uniformly continuous on $\left(\mathbb{R}^{d}\right)^{n}$, if $\left(g_{m}\right)_{m=1}^{\infty} \subset C^{2}\left(\left(\mathbb{R}^{d}\right)^{n}\right)$ have bounded first and second derivatives and $0 \leq g_{m}-g \leq a_{m} \rightarrow 0$ as $m \rightarrow \infty$, then for sufficiently large $C_{m}>0$ the functions $g_{m}(\mathbf{x})+C_{m} t$ and $g_{m}(\mathbf{x})-a_{m}-C_{m} t$ are respectively viscosity supersolutions and viscosity subsolutions of (3.3) for all $m=1,2, \cdots$ Thus the functions

$$
\bar{u}(t, \mathbf{x}):=\inf _{m \geq 1}\left(g_{m}(\mathbf{x})+C_{m} t\right), \quad \underline{u}(t, \mathbf{x}):=\sup _{m \geq 1}\left(g_{m}(\mathbf{x})-a_{m}-C_{m} t\right)
$$

are respectively a viscosity supersolution and viscosity subsolution of (3.3) such that $\bar{u}(0, \mathbf{x})=$ $\underline{u}(0, \mathbf{x})=g(\mathbf{x})$ on $\left(\mathbb{R}^{d}\right)^{n}$. Therefore the unique bounded continuous viscosity solution $u_{n}$ of (3.3) can be constructed, for instance, by means of Perron's method (see [26]).

It is also easy to see that there exists $M>0$ such that $\left\|u_{n}\right\|_{\infty} \leq M, n=1,2, \cdots$. To show this, let $K>0$ be such that $\|\mathcal{F}\|_{\infty} \leq K,\left\|\mathcal{U}_{0}\right\|_{\infty} \leq K$. Recall that we have $\|H(\cdot, \cdot, 0)\|_{\infty} \leq C$, where $C$ is from 2.3). Then $w_{1}(t, \mathbf{x}):=-K-(C+K) t$ is a bounded viscosity subsolution of (3.3) and $w_{2}(t, \mathbf{x}):=K+(C+K) t$ is a bounded viscosity supersolution of (3.3). Therefore, by comparison we obtain that for every $n$

$$
w_{1} \leq u_{n} \leq w_{2} \quad \text { on }[0, T] \times\left(\mathbb{R}^{d}\right)^{n} .
$$

This gives the required bound with $M=K+(C+K) T$.
For $\delta>0$, let $\varphi_{\delta}$ be the function from Lemma 3.2 applied to the modulus

$$
\sigma_{1}(s)=(1+T) \sigma(3 s)+m_{1}(s)+\left(2 \kappa d L^{2}(1+T)+2 M+1\right) s
$$

In particular we have

$$
\begin{equation*}
\varphi_{\delta}(1) \geq 2 M+1, \quad \varphi_{\delta}(s) \geq m_{1}(s) . \tag{3.5}
\end{equation*}
$$

First, we are going to show that for every $\delta>0$

$$
\begin{equation*}
u_{n}(t, \mathbf{x})-u_{n}(t, \mathbf{y}) \leq \varphi_{\delta}\left(|\mathbf{x}-\mathbf{y}|_{r}\right)(1+t), \tag{3.6}
\end{equation*}
$$

if $t \in[0, T]$ and $(\mathbf{x}, \mathbf{y}) \in\left(\mathbb{R}^{d}\right)^{n}$.
We define smooth approximations of $|\mathbf{z}|_{r}$. For $\gamma>0$, let

$$
\psi_{\gamma}(\mathbf{z})=\frac{1}{n^{1 / r}}\left(\sum_{i=1}^{n}\left(\gamma+\left|z_{i}\right|^{2}\right)^{\frac{r}{2}}\right)^{\frac{1}{r}}
$$

We now set

$$
\varphi(t, \mathbf{x}, \mathbf{y}):=\varphi_{\delta}\left(\psi_{\gamma}(\mathbf{x}-\mathbf{y})\right)(1+t) .
$$

Suppose that there exist $\gamma, \mu>0$ such that

$$
\begin{equation*}
\sup _{\mathbf{x}, \mathbf{y} \in\left(\mathbb{R}^{d}\right)^{n}, t \in[0, T]}\left(u_{n}(t, \mathbf{x})-u_{n}(t, \mathbf{y})-\frac{\mu}{T-t}-\varphi(t, \mathbf{x}, \mathbf{y})\right)>0 . \tag{3.7}
\end{equation*}
$$

Then, for every $\alpha>0$ small enough and $h(\mathbf{x}):=\left(1+|\mathbf{x}|^{2}\right)^{1 / 2}$,

$$
\begin{equation*}
\sup _{\mathbf{x}, \mathbf{y} \in\left(\mathbb{R}^{d}\right)^{n}, t \in[0, T]}\left(u_{n}(t, \mathbf{x})-u_{n}(t, \mathbf{y})-\frac{\mu}{T-t}-\varphi(t, \mathbf{x}, \mathbf{y})-\alpha(h(\mathbf{x})+h(\mathbf{y}))\right)>0 . \tag{3.8}
\end{equation*}
$$

Moreover, since $u_{n}$ is bounded and $\alpha(h(\mathbf{x})+h(\mathbf{y})) \rightarrow+\infty$ as $x, y \rightarrow \infty$, the supremum of the above expression over $[0, T] \times\left(\mathbb{R}^{d}\right)^{n} \times\left(\mathbb{R}^{d}\right)^{n}$ is attained at some point $(\bar{t}, \overline{\mathbf{x}}, \overline{\mathbf{y}})$. Obviously $\bar{t}<T$ and it follows from the definition of the function $\varphi(t, \mathbf{x}, \mathbf{y})$ and (3.5) that we must have $0<\bar{t}$. It also follows from (3.5) that $\bar{s}=\psi_{\gamma}(\overline{\mathbf{x}}-\overline{\mathbf{y}})<1$. Also $\overline{\mathbf{x}} \neq \overline{\mathbf{y}}$ since if $\overline{\mathbf{x}}=\overline{\mathbf{y}}$, the expression in (3.8) is negative. We compute

$$
\left.D_{x_{i}} \varphi(\bar{t}, \overline{\mathbf{x}}, \overline{\mathbf{y}})\right)=\varphi_{\delta}^{\prime}(\bar{s}) \frac{\left(\bar{x}_{i}-\bar{y}_{i}\right)\left(\gamma+\left|\bar{x}_{i}-\bar{y}_{i}\right|^{2}\right)^{\frac{r}{2}-1}}{n \bar{s}^{\frac{r}{r^{\prime}}}}(1+\bar{t}),
$$

where $1 / r+1 / r^{\prime}=1$. Also, using the concavity of $\varphi_{\delta}$, we have

$$
\left.D_{\mathbf{x}}^{2} \varphi(\bar{t}, \overline{\mathbf{x}}, \overline{\mathbf{y}})\right)=B_{1}-B_{2},
$$

where $B_{1}$ is a diagonal matrix composed of $n$ diagonal $d \times d$ blocks

$$
B_{1 i}=\varphi_{\delta}^{\prime}(\bar{s}) \frac{\left(\gamma+\left|\bar{x}_{i}-\bar{y}_{i}\right|^{2}\right)^{\frac{r}{2}-1}}{n \bar{s}^{\frac{r}{r}}}(1+\bar{t}) I_{d},
$$

and $B_{2} \geq 0$ is a symmetric matrix. Therefore

$$
D=D^{2} \varphi(\bar{t}, \overline{\mathbf{x}}, \overline{\mathbf{y}})=\left(\begin{array}{cc}
B_{1} & -B_{1} \\
-B_{1} & B_{1}
\end{array}\right)-\left(\begin{array}{cc}
B_{2} & -B_{2} \\
-B_{2} & B_{2}
\end{array}\right)=: D_{1}-D_{2},
$$

where $D^{2} \varphi$ above is the second derivative of $\varphi$ with respect to the variables $(\mathbf{x}, \mathbf{y})$. We now use Theorem 8.3 of [26] applied to the functions

$$
u^{1}(t, x):=u_{n}(t, x)-\alpha h(x), \quad u^{2}(t, y):=-u_{n}(t, y)-\alpha h(y) .
$$

We notice that, since $u_{n}$ is a viscosity solution of (3.3), condition (8.5) of [26] is satisfied. Therefore, it follows from Theorem 8.3 of [26], applied with $\varepsilon=1 /\left(\left\|D_{1}\right\|+\left\|D_{2}\right\|\right)$, that there exist $b_{1}, b_{2} \in \mathbb{R}$ and $S_{1}, S_{2} \in S(n d)$ such that

$$
\begin{gathered}
\left(b_{1}, D_{\mathbf{x}} \varphi(\bar{t}, \overline{\mathbf{x}}, \overline{\mathbf{y}})+\alpha D h(\overline{\mathbf{x}}), S_{1}+\alpha D^{2} h(\overline{\mathbf{x}})\right) \in \overline{\mathcal{P}}^{2,+} u_{n}(\bar{t}, \overline{\mathbf{x}}), \\
\left(b_{2},-D_{\mathbf{y}} \varphi(\bar{t}, \overline{\mathbf{x}}, \overline{\mathbf{y}})-\alpha D h(\overline{\mathbf{y}}), S_{2}-\alpha D^{2} h(\overline{\mathbf{y}})\right) \in \overline{\mathcal{P}}^{2,-} u_{n}(\bar{t}, \overline{\mathbf{y}}), \\
b_{1}-b_{2}=\varphi_{\delta}(\bar{s})+\frac{\mu}{(T-\bar{t})^{2}},
\end{gathered}
$$

where

$$
\left(\begin{array}{cc}
S_{1} & 0  \tag{3.9}\\
0 & -S_{2}
\end{array}\right) \leq D+\frac{1}{\left\|D_{1}\right\|+\left\|D_{2}\right\|} D^{2} \leq 2 D_{1}
$$

where we used that

$$
D^{2} \leq\left(\left\|D_{1}\right\|+\left\|D_{2}\right\|\right)\left(D_{1}+D_{2}\right)
$$

Inequality (3.9) in particular implies that $S_{1} \leq S_{2}$. Using the definition of viscosity solution and setting

$$
Z(\overline{\mathbf{x}}, \overline{\mathbf{y}}):=\frac{1}{n} \sum_{i=1}^{n} H\left(\bar{x}_{i}, \frac{1}{n-1} \sum_{j \neq i} \delta_{\bar{x}_{j}}, \varphi_{\delta}^{\prime}(\bar{s}) \frac{\left(\bar{x}_{i}-\bar{y}_{i}\right)\left(\gamma+\left|\bar{x}_{i}-\bar{y}_{i}\right|^{2}\right)^{\frac{r}{2}-1}}{\bar{s}^{\frac{r}{r}}}(1+\bar{t})+n \alpha D h(\overline{\mathbf{x}})\right)
$$

and

$$
\tilde{Z}(\overline{\mathbf{x}}, \overline{\mathbf{y}}):=\frac{1}{n} \sum_{i=1}^{n} H\left(\bar{y}_{i}, \frac{1}{n-1} \sum_{j \neq i} \delta_{\bar{y}_{j}}, \varphi_{\delta}^{\prime}(\bar{s}) \frac{\left(\bar{x}_{i}-\bar{y}_{i}\right)\left(\gamma+\left|\bar{x}_{i}-\bar{y}_{i}\right|^{2}\right)^{\frac{r}{2}-1}}{\bar{s}^{\frac{r}{r^{\prime}}}}(1+\bar{t})-n \alpha D h(\overline{\mathbf{y}})\right)
$$

we now have

$$
\begin{equation*}
b_{1}-\kappa \operatorname{Tr}\left(A_{n}(\overline{\mathbf{x}})\left(S_{1}+\alpha D^{2} h(\overline{\mathbf{x}})\right)\right)+Z(\overline{\mathbf{x}}, \overline{\mathbf{y}})+\mathcal{F}\left(\mu_{\overline{\mathbf{x}}}\right) \leq 0 \tag{3.10}
\end{equation*}
$$

and

$$
\begin{equation*}
b_{2}-\kappa \operatorname{Tr}\left(A_{n}(\overline{\mathbf{y}})\left(S_{2}-\alpha D^{2} h(\overline{\mathbf{y}})\right)\right)+\tilde{Z}(\overline{\mathbf{x}}, \overline{\mathbf{y}})+\mathcal{F}\left(\mu_{\overline{\mathbf{y}}}\right) \geq 0 \tag{3.11}
\end{equation*}
$$

We notice that (3.1) and (3.9) imply

$$
\begin{aligned}
& \operatorname{Tr}\left(A_{n}(\overline{\mathbf{x}}) S_{1}\right)-\operatorname{Tr}\left(A_{n}(\overline{\mathbf{y}}) S_{2}\right) \\
& \leq 2 \sum_{i=1}^{n} \operatorname{Tr}\left(\left(a\left(x_{i}\right)-a\left(y_{i}\right)\right)\left(a\left(x_{i}\right)-a\left(y_{i}\right)\right)^{*} \varphi_{\delta}^{\prime}(\bar{s}) \frac{\left(\gamma+\left|\bar{x}_{i}-\bar{y}_{i}\right|^{2}\right)^{\frac{r}{2}-1}}{n \bar{s}^{\frac{r}{r^{\prime}}}}(1+\bar{t}) I_{d}\right) \\
& \leq \frac{1}{n} \sum_{i=1}^{n} \varphi_{\delta}^{\prime}(\bar{s}) \frac{2 d L^{2}\left|\bar{x}_{i}-\bar{y}_{i}\right|^{2}\left(\gamma+\left|\bar{x}_{i}-\bar{y}_{i}\right|^{2}\right)^{\frac{r}{2}-1}}{\bar{s}^{\frac{r}{r^{\prime}}}}(1+\bar{t}) .
\end{aligned}
$$

Note that if $A_{n}$ is a constant matrix then obviously $\operatorname{Tr}\left(A_{n} S_{1}\right)-\operatorname{Tr}\left(A_{n} S_{2}\right) \leq 0$. Let us use the notation

$$
\mu_{\mathrm{x}}^{i}:=\frac{1}{n-1} \sum_{j \neq i} \delta_{x_{j}}, \quad \mu_{\mathrm{y}}^{i}:=\frac{1}{n-1} \sum_{j \neq i} \delta_{y_{j}}
$$

A simple calculation shows that $d_{r}\left(\mu_{\overline{\mathrm{x}}}^{i}, \mu_{\overline{\mathrm{y}}}^{i}\right) \leq 2 \bar{s}$. Subtracting (3.11) from 3.10 and using Hypothesis 2.1, the concavity of $\sigma$ and (2.10), we obtain

$$
\begin{aligned}
\varphi_{\delta}(\bar{s}) & +\frac{\mu}{(T-\bar{t})^{2}} \leq \frac{\kappa}{n} \sum_{i=1}^{n} \varphi_{\delta}^{\prime}(\bar{s}) \frac{2 d L^{2}\left|\bar{x}_{i}-\bar{y}_{i}\right|^{2}\left(\gamma+\left|\bar{x}_{i}-\bar{y}_{i}\right|^{2}\right)^{\frac{r}{2}-1}}{\bar{s}^{\frac{r}{r^{\prime}}}}(1+\bar{t}) \\
& +\frac{1}{n} \sum_{i=1}^{n} \sigma\left(\left(\left|\bar{x}_{i}-\bar{y}_{i}\right|+d_{r}\left(\mu_{\overline{\mathbf{x}}}^{i}, \mu_{\overline{\mathbf{y}}}^{i}\right)\right)\left(1+\varphi_{\delta}^{\prime}(\bar{s}) \frac{\left|\bar{x}_{i}-\bar{y}_{i}\right|\left(\gamma+\left|\bar{x}_{i}-\bar{y}_{i}\right|^{2}\right)^{\frac{r}{2}-1}}{\bar{s}^{\frac{r}{r^{\prime}}}}(1+\bar{t})\right)\right) \\
& +m_{1}(\bar{s})+\sigma_{2}(\alpha),
\end{aligned}
$$

where $\lim _{\alpha \rightarrow 0} \sigma_{2}(\alpha)=0$. Thus,

$$
\begin{aligned}
\varphi_{\delta}(\bar{s})+\frac{\mu}{(T-\bar{t})^{2}} \leq & \frac{\kappa}{n} \sum_{i=1}^{n} \varphi_{\delta}^{\prime}(\bar{s}) \frac{2 d L^{2}\left|\bar{x}_{i}-\bar{y}_{i}\right|^{2}\left(\gamma+\left|\bar{x}_{i}-\bar{y}_{i}\right|^{2}\right)^{\frac{r}{2}-1}}{\bar{s}^{\frac{r}{r^{\prime}}}}(1+\bar{t}) \\
& +\sigma\left(\frac{1}{n} \sum_{i=1}^{n}\left(\left|\bar{x}_{i}-\bar{y}_{i}\right|+2 \bar{s}\right)\left(1+\varphi_{\delta}^{\prime}(\bar{s}) \frac{\left|\bar{x}_{i}-\bar{y}_{i}\right|\left(\gamma+\left|\bar{x}_{i}-\bar{y}_{i}\right|^{2}\right)^{\frac{r}{2}-1}}{\bar{s}^{\frac{r}{r^{\prime}}}}(1+\bar{t})\right)\right) \\
& +m_{1}(\bar{s})+\sigma_{2}(\alpha)
\end{aligned}
$$

By Jensen's inequality

$$
|\overline{\mathbf{x}}-\overline{\mathbf{y}}|_{1} \leq|\overline{\mathbf{x}}-\overline{\mathbf{y}}|_{r} \leq \bar{s}
$$

and also,

$$
\frac{1}{n} \sum_{i=1}^{n} \varphi_{\delta}^{\prime}(\bar{s}) \frac{\left|\bar{x}_{i}-\bar{y}_{i}\right|^{2}\left(\gamma+\left|\bar{x}_{i}-\bar{y}_{i}\right|^{2}\right)^{\frac{r}{2}-1}}{\bar{s}^{\frac{r}{r}}} \leq \varphi_{\delta}^{\prime}(\bar{s}) \bar{s} ;
$$

Using Schwarz's inequality, one verifies that

$$
\begin{aligned}
\frac{1}{n} \sum_{i=1}^{n} 2 \bar{s} \varphi_{\delta}^{\prime}(\bar{s}) \frac{\left|\bar{x}_{i}-\bar{y}_{i}\right|\left(\gamma+\left|\bar{x}_{i}-\bar{y}_{i}\right|^{2}\right)^{\frac{r}{2}-1}}{\bar{s}^{\frac{r}{r}}} & \leq \frac{1}{n} \sum_{i=1}^{n} 2 \varphi_{\delta}^{\prime}(\bar{s}) \bar{s}^{1-r / r^{\prime}}\left(\gamma+\left|\bar{x}_{i}-\bar{y}_{i}\right|^{2}\right)^{\frac{r-1}{2}} \\
& \leq 2 \varphi_{\delta}^{\prime}(\bar{s}) s^{1-\frac{r}{r^{\prime}}}\left(\frac{1}{n} \sum_{i=1}^{n}\left(\gamma+\left|\bar{x}_{i}-\bar{y}_{i}\right|^{2}\right)^{\frac{r}{2}}\right)^{\frac{r-1}{r}} \\
& =2 \varphi_{\delta}^{\prime}(\bar{s}) s^{1-r / r^{\prime}} \bar{s}^{r-1}=2 \varphi_{\delta}^{\prime}(\bar{s}) \bar{s}
\end{aligned}
$$

Collecting these bounds into the inequality obtained for $\varphi_{\delta}(\bar{s})+\frac{\mu}{(T-t)^{2}}$ above, and using the sub-additivity of $\sigma$ and the definition of $\sigma_{1}$, we get

$$
\begin{aligned}
\varphi_{\delta}(\bar{s})+\frac{\mu}{T^{2}} & \leq 2 \kappa d L^{2}(1+T) \varphi_{\delta}^{\prime}(\bar{s}) \bar{s}+(1+T) \sigma\left(3 \varphi_{\delta}^{\prime}(\bar{s}) \bar{s}+3 \bar{s}\right)+m_{1}(\bar{s})+\sigma_{2}(\alpha) \\
& \leq \sigma_{1}\left(\varphi_{\delta}^{\prime}(\bar{s}) \bar{s}+\bar{s}\right)+\sigma_{2}(\alpha) .
\end{aligned}
$$

This gives a contradiction when we let $\alpha \rightarrow 0$, due to $\sqrt{3.2}$. Consequently, for all positive $\gamma$ and $\mu$, (3.7) is false. Letting $\gamma, \mu \rightarrow 0$, we obtain (3.6). Thus we have proved that for all $t \in[0, T],(\mathbf{x}, \mathbf{y}) \in\left(\mathbb{R}^{d}\right)^{n}$,

$$
\begin{equation*}
\left|u_{n}(t, \mathbf{x})-u_{n}(t, \mathbf{y})\right| \leq \inf _{\delta>0} \varphi_{\delta}\left(|\mathbf{x}-\mathbf{y}|_{r}\right)(1+T) . \tag{3.12}
\end{equation*}
$$

We will now obtain the continuity estimate with respect to $t$. We know by (2.9) that

$$
\left|u_{n}(0, \mathbf{x})-u_{n}(0, \mathbf{y})\right| \leq m_{1}\left(|\mathbf{x}-\mathbf{y}|_{2}\right)
$$

Setting $v_{n}(\mathbf{x})=u_{n}(0, \sqrt{n} \mathbf{x})$ we thus have

$$
\left|v_{n}(\mathbf{x})-v_{n}(\mathbf{y})\right| \leq m_{1}(|\mathbf{x}-\mathbf{y}|) .
$$

Approximating the functions $v_{n}(\mathbf{x})$ by supinf-convolutions and then mollifying them, there exist constants $L_{m}, m=1,2, \cdots$ (independent of $n$ ) and $C^{2}$ functions $\tilde{\varphi}_{m}^{n}$ such that $0 \leq \tilde{\varphi}_{m}^{n}(\mathbf{x})-$ $v_{n}(\mathbf{x}) \leq \frac{1}{m}$ on $\left(\mathbb{R}^{d}\right)^{n}$ and

$$
\left|D \tilde{\varphi}_{m}^{n}\right| \leq L_{m}, \quad\left\|D^{2} \tilde{\varphi}_{m}^{n}\right\| \leq L_{m}
$$

Then if $\varphi_{m}^{n}(\mathbf{x})=\tilde{\varphi}_{m}^{n}\left(\frac{1}{\sqrt{n}} \mathbf{x}\right)$, we have

$$
\begin{equation*}
0 \leq \varphi_{m}^{n}(\mathbf{x})-u_{n}(0, \mathbf{x}) \leq \frac{1}{m} \quad \text { on }\left(\mathbb{R}^{d}\right)^{n} \tag{3.13}
\end{equation*}
$$

and

$$
\left|D \varphi_{m}^{n}\right| \leq \frac{L_{m}}{\sqrt{n}}, \quad\left\|D^{2} \varphi_{m}^{n}\right\| \leq \frac{L_{m}}{n}
$$

Recall that $K$ is such that $\|\mathcal{F}\|_{\infty} \leq K$. We set $C_{m}=\kappa \tilde{L} L_{m}+C+C L_{m}^{2}+K$, where $\tilde{L}$ is such that $\operatorname{Tr}\left(A_{n}(\mathbf{x})\right) \leq \tilde{L} n$ and $C$ is from Hypothesis 2.1. We define the functions

$$
\psi_{m}^{n}(t, \mathbf{x}):=\varphi_{m}^{n}(\mathbf{x})+C_{m} t
$$

Then

$$
\begin{aligned}
& \partial_{t} \psi_{m}^{n}(t, \mathbf{x})-\kappa \operatorname{Tr}\left(A_{n}(\mathbf{x}) D^{2} \psi_{m}^{n}(t, \mathbf{x})\right)+\frac{1}{n} \sum_{i=1}^{n} H\left(x_{i}, \mu_{\mathbf{x}}^{i}, n D_{x_{i}} \psi_{m}^{n}(t, \mathbf{x})\right)+\mathcal{F}\left(\mu_{\mathbf{x}}\right) \\
& \geq C_{m}-\frac{\kappa}{n} \operatorname{Tr}\left(A_{n}(\mathbf{x}) D^{2} \tilde{\varphi}_{m}^{n}\left(\frac{1}{\sqrt{n}} \mathbf{x}\right)\right)+\frac{1}{n} \sum_{i=1}^{n} H\left(x_{i}, \mu_{\mathbf{x}}^{i}, \sqrt{n} D_{x_{i}} \tilde{\varphi}_{m}^{n}\left(\frac{1}{\sqrt{n}} \mathbf{x}\right)\right)-K \\
& \geq C_{m}-\frac{\kappa}{n} \operatorname{Tr}\left(A_{n}(\mathbf{x})\right)\left\|D^{2} \tilde{\varphi}_{m}^{n}\left(\frac{1}{\sqrt{n}} \mathbf{x}\right)\right\|-\frac{1}{n} \sum_{i=1}^{n} C\left(1+n\left|D_{x_{i}} \tilde{\varphi}_{m}^{n}\left(\frac{1}{\sqrt{n}} \mathbf{x}\right)\right|^{2}\right)-K \\
& \geq C_{m}-\kappa \tilde{L} L_{m}-C-C \sum_{i=1}^{n}\left|D_{x_{i}} \tilde{\varphi}_{m}^{n}\left(\frac{1}{\sqrt{n}} \mathbf{x}\right)\right|^{2}-K \\
& \geq C_{m}-\kappa \tilde{L} L_{m}-C-C L_{m}^{2}-K=0 .
\end{aligned}
$$

Therefore the functions $\psi_{m}^{n}$ are viscosity supersolutions of equations (3.3). Similarly, the functions

$$
\tilde{\psi}_{m}^{n}(t, \mathbf{x}):=\varphi_{m}^{n}(\mathbf{x})-\frac{1}{m}-C_{m} t
$$

are viscosity subsolutions of equations (3.3). Therefore, by comparison, we have for every $n, m \geq$ 1,

$$
\varphi_{m}^{n}(\mathbf{x})-\frac{1}{m}-C_{m} t \leq u_{n}(t, \mathbf{x}) \leq \varphi_{m}^{n}(\mathbf{x})+C_{m} t
$$

Using (3.13), this implies that

$$
\begin{aligned}
-\frac{1}{m}-C_{m} t \leq \varphi_{m}^{n}(\mathbf{x})-\frac{1}{m}-C_{m} t-u_{n}(0, \mathbf{x}) & \leq u_{n}(t, \mathbf{x})-u_{n}(0, \mathbf{x}) \\
& \leq \varphi_{m}^{n}(\mathbf{x})+C_{m} t-u_{n}(0, \mathbf{x}) \\
& \leq C_{m} t+\frac{1}{m}
\end{aligned}
$$

Therefore we obtain

$$
\left|u_{n}(t, \mathbf{x})-u_{n}(0, \mathbf{x})\right| \leq \tilde{\rho}(t):=\inf \left\{\frac{1}{m}+C_{m} t: m=1,2, \cdots\right\}
$$

for all $(t, \mathbf{x}) \in[0, T] \times\left(\mathbb{R}^{d}\right)^{n}, n=1,2, \cdots$. The function $\tilde{\rho}$ is independent of $n$. We then define for every $h \in(0, T)$ the functions

$$
v_{n}^{h}(t, \mathbf{x})=u_{n}(t+h, \mathbf{x}), \quad(t, \mathbf{x}) \in[0, T-h] \times\left(\mathbb{R}^{d}\right)^{n}
$$

The functions $v_{n}^{h}$ are viscosity solutions of 3.3 on $(0, T-h) \times\left(\mathbb{R}^{d}\right)^{n}$ and

$$
\left|v_{n}^{h}(0, \mathbf{x})-u_{n}(0, \mathbf{x})\right| \leq \tilde{\rho}(h)
$$

By comparison we thus obtain

$$
\begin{equation*}
\left|u_{n}(t+h, \mathbf{x})-u_{n}(t, \mathbf{x})\right|=\left|v_{n}^{h}(t, \mathbf{x})-u_{n}(t, \mathbf{x})\right| \leq \tilde{\rho}(h) \tag{3.14}
\end{equation*}
$$

for $(t, \mathbf{x}) \in(0, T-h) \times\left(\mathbb{R}^{d}\right)^{n}$. We now let

$$
\rho(s)=\tilde{\rho}(s)+\inf _{\delta>0} \varphi_{\delta}(s)(1+T)
$$

Combining (3.14 with (3.12), we obtain (3.4) for this modulus $\rho$, uniformly with respect to $n \in \mathbb{N}$.

## 4. Proof of Theorem 1.2

Proof of Theorem 1.2. Step 1. We define the functions

$$
\mathcal{V}_{n}\left(t, \mu_{\mathbf{x}}\right):=u_{n}(t, \mathbf{x}), \quad \text { where } \quad \mu_{\mathbf{x}}=\frac{1}{n} \sum_{i=1}^{n} \delta_{x_{i}}
$$

This function is well defined since the functions $u_{n}$ are invariant with respect to permutations of the variables in $\mathbf{x}$. This follows from uniqueness of viscosity solutions of $(1.2)$ as these equations are invariant with respect to permutations of the variables in $\mathbf{x}$. The function $\mathcal{V}_{n}(t, \cdot)$ is now defined on the subset of $\mathcal{P}_{2}\left(\mathbb{R}^{d}\right)$ that consists of averages of $n$ Dirac point masses. It follows from (3.4) that

$$
\left|\mathcal{V}_{n}\left(t, \mu_{\mathbf{x}}\right)-\mathcal{V}_{n}\left(s, \mu_{\mathbf{y}}\right)\right| \leq \rho\left(|t-s|+d_{r}\left(\mu_{\mathbf{x}}, \mu_{\mathbf{y}}\right)\right) \text { for all } t, s \in[0, T], \mathbf{x}, \mathbf{y} \in\left(\mathbb{R}^{d}\right)^{n}
$$

For each fixed $n$, we can extend the function $\mathcal{V}_{n}$ to a function on $[0, T] \times \mathcal{P}_{r}\left(\mathbb{R}^{d}\right)$, still denoted by $\mathcal{V}_{n}$, satisfying

$$
\begin{equation*}
\left|\mathcal{V}_{n}(t, \mu)-\mathcal{V}_{n}(s, \nu)\right| \leq \rho\left(|t-s|+d_{r}(\mu, \nu)\right) \text { for all } t, s \in[0, T], \mu, \nu \in \mathcal{P}_{r}\left(\mathbb{R}^{d}\right) \tag{4.1}
\end{equation*}
$$

Since sets

$$
M_{R}^{2}=\left\{\mu \in \mathcal{P}_{r}\left(\mathbb{R}^{d}\right): \int_{\mathbb{R}^{d}}|x|^{2} \mu(d x) \leq R\right\}
$$

are relatively compact in $\mathcal{P}_{r}\left(\mathbb{R}^{d}\right)$, up to a subsequence (still denoted by $\mathcal{V}_{n}$ ), $\mathcal{V}_{n}$ converges uniformly on every set $[0, T] \times M_{R}^{2}$ to a function $\mathcal{V}:[0, T] \times \mathcal{P}_{2}\left(\mathbb{R}^{d}\right) \rightarrow \mathbb{R}$ which satisfies the same estimate (4.1). Define now

$$
\begin{aligned}
V:[0, T] \times E & \longrightarrow \mathbb{R} \\
(t, X) & \longmapsto V(t, X):=\mathcal{V}(t, \operatorname{law}(X)) .
\end{aligned}
$$

We will show that $V$ is a viscosity solution of (1.3). Since equation (1.3) has a unique bounded viscosity solution $U$, we can then conclude that $V=U$. This will prove the theorem since the argument can be done for any subsequence of $\mathcal{V}_{n}$. The proof that $\sqrt{1.3}$ has a unique bounded viscosity solution is included in the Appendix, Theorem 7.4 .

Step 2. Let then $\varphi \in C^{1,2}((0, T) \times E)$ and suppose that $V-\varphi$ has a local maximum at $\left(t_{0}, X_{0}\right) \in(0, T) \times E$. By considering $\varphi(t, X)+\left(t-t_{0}\right)^{2}+\left|X-X_{0}\right|^{2}$ and modifying it outside of a neighborhood of ( $t_{0}, X_{0}$ ) we can assume with no loss of generality that the maximum at $\left(t_{0}, X_{0}\right)$ is strict and global. Being a strict maximum implies that if $V\left(t_{i}, X_{i}\right)-\varphi\left(t_{i}, X_{i}\right) \rightarrow$ $V\left(t_{0}, X_{0}\right)-\varphi\left(t_{0}, X_{0}\right)$ then $\left(t_{i}, X_{i}\right) \rightarrow\left(t_{0}, X_{0}\right)$. Denote $P_{0}=D \varphi\left(t_{0}, X_{0}\right)$. For each $\epsilon>0$ let $X_{\epsilon}, P_{\epsilon} \in E$ be such that $X_{\epsilon}, P_{\epsilon}$ are continuous on $[0,1]$ and $\left|X_{0}-X_{\epsilon}\right|+\left|P_{0}-P_{\epsilon}\right|<\epsilon$.

For every $n$ we denote $A_{i}^{n}=\left(\frac{i-1}{n}, \frac{i}{n}\right), i=1, \cdots, n$. We then consider the function $\varphi_{n}$ : $(0, T) \times\left(\mathbb{R}^{d}\right)^{n} \rightarrow \mathbb{R}$ defined by

$$
\varphi_{n}(t, \mathbf{x}):=\varphi\left(t, \sum_{i=1}^{n} x_{i} \mathbf{1}_{A_{i}^{n}}\right),
$$

where $\mathbf{1}_{A_{i}^{n}}$ is the characteristic function of the set $A_{i}^{n}$.
Since the original maximum at $\left(t_{0}, X_{0}\right)$ was strict it is easy to see that the functions $u_{n}-\varphi_{n}$ must have local maxima at points $\left(t_{n}, \mathrm{x}^{n}\right)=\left(t_{n}, x_{1}^{n}, \cdots, x_{n}^{n}\right)$ such that

$$
t_{n} \rightarrow t_{0} \quad \text { and } \quad X^{n}=\sum_{i=1}^{n} x_{i}^{n} \mathbf{1}_{A_{i}^{n}} \rightarrow X_{0}
$$

In particular for sufficiently big $n$,

$$
\left|t_{n}-t_{0}\right|+\left|X_{0}-X^{n}\right| \leq \epsilon .
$$

Now, by the chain rule,

$$
D_{x_{i}} \varphi_{n}\left(t_{n}, \mathbf{x}^{n}\right)=\left\langle D \varphi\left(t_{n}, \sum_{i=1}^{n} x_{i}^{n} \mathbf{1}_{A_{i}^{n}}\right), \mathbf{1}_{A_{i}^{n}} e\right\rangle_{d}=\int_{A_{i}^{n}} D \varphi\left(t_{n}, \sum_{i=1}^{n} x_{i}^{n} \mathbf{1}_{A_{i}^{n}}\right),
$$

where $e=(1, \cdots, 1) \in \mathbb{R}^{d}$. If $\mathbf{x}=\left(x_{1}, \cdots, x_{n}\right)$, we will denote $x_{i}=\left(x_{i 1}, \cdots, x_{i d}\right)$. Then

$$
\operatorname{Tr}\left(A_{n} D^{2} \varphi_{n}\right)=\sum_{i, j=1}^{n} \operatorname{Tr}\left(D_{x_{i} x_{j}}^{2} \varphi_{n}\right)=\sum_{i, j=1}^{n} \sum_{k=1}^{d} \frac{\partial^{2} \varphi_{n}}{\partial x_{i k} \partial x_{j k}} .
$$

Now

$$
\frac{\partial^{2} \varphi_{n}}{\partial x_{i k} \partial x_{j k}}\left(t_{n}, \mathbf{x}^{n}\right)=\int_{A_{i}^{n}}\left(D^{2} \varphi\left(t_{n}, \sum_{i=1}^{n} x_{i}^{n} \mathbf{1}_{A_{i}^{n}}\right) \mathbf{1}_{A_{j}^{n}} e_{k}\right) \cdot e_{k}
$$

so

$$
\begin{aligned}
\operatorname{Tr}\left(A_{n} D^{2} \varphi_{n}\right) & =\sum_{i, j=1}^{n} \sum_{k=1}^{d} \int_{A_{i}^{n}}\left(D^{2} \varphi\left(t_{n}, \sum_{i=1}^{n} x_{i}^{n} \mathbf{1}_{A_{i}^{n}}\right) \mathbf{1}_{A_{j}^{n}} e_{k}\right) \cdot e_{k} \\
& =\sum_{k=1}^{d}\left\langle D^{2} \varphi\left(t_{n}, \sum_{i=1}^{n} x_{i}^{n} \mathbf{1}_{A_{i}^{n}}\right) e_{k}, e_{k}\right\rangle .
\end{aligned}
$$

We note that if

$$
\mathbf{x}, \mathbf{y} \in\left(\mathbb{R}^{d}\right)^{n}, \quad X=\sum_{i=1}^{n} x_{i} \mathbf{1}_{A_{i}^{n}} \quad \text { and } \quad Y=\sum_{i=1}^{n} y_{i} \mathbf{1}_{A_{i}^{n}}
$$

then

$$
|X-Y|^{2}=n^{-1} \sum_{i=1}^{n}\left|x_{i}-y_{i}\right|^{2}, \quad|X-Y|_{r}^{r}=n^{-1} \sum_{i=1}^{n}\left|x_{i}-y_{i}\right|^{r} .
$$

Furthermore, choosing $\xi, \eta \in E$, we have

$$
\begin{aligned}
\left|\sum_{i=1}^{n} H\left(x_{i}, \mu_{\mathbf{x}}^{i}, n \int_{A_{i}^{n}} \xi\right)-\sum_{i=1}^{n} H\left(y_{i}, \mu_{\mathbf{y}}^{i}, n \int_{A_{i}^{n}} \eta\right)\right| & \leq \sum_{i=1}^{n}\left|H\left(x_{i}, \mu_{\mathbf{x}}^{i}, n \int_{A_{i}^{n}} \xi\right)-H\left(y_{i}, \mu_{\mathbf{y}}^{i}, n \int_{A_{i}^{n}} \eta\right)\right| \\
& \leq C \sum_{i=1}^{n}\left(1+\left|n \int_{A_{i}^{n}} \xi\right|+\left|n \int_{A_{i}^{n}} \eta\right|\right) \int_{A_{i}^{n}} n|\xi-\eta| \\
& +\sum_{i=1}^{n} \sigma\left(\left(\left|x_{i}-y_{i}\right|+d_{r}\left(\mu_{\mathbf{x}}^{i}, \mu_{\mathbf{y}}^{i}\right)\right)\left(1+\left|n \int_{A_{i}^{n}} \xi\right|\right)\right)
\end{aligned}
$$

Thus,

$$
\begin{aligned}
& n^{-1}\left|\sum_{i=1}^{n} H\left(x_{i}, \mu_{\mathbf{x}}^{i}, n \int_{A_{i}^{n}} \xi\right)-\sum_{i=1}^{n} H\left(y_{i}, \mu_{\mathbf{y}}^{i}, n \int_{A_{i}^{n}} \eta\right)\right| \\
& \leq C \int_{\Omega}|\xi-\eta|+C\left(\left(\sum_{i=1}^{n} n\left|\int_{A_{i}^{n}} \xi\right|^{2}\right)^{\frac{1}{2}}+\left(\sum_{i=1}^{n} n\left|\int_{A_{i}^{n}} \eta\right|^{2}\right)^{\frac{1}{2}}\right)\left(\sum_{i=1}^{n} n\left|\int_{A_{i}^{n}}\right| \xi-\eta| |^{2}\right)^{\frac{1}{2}} \\
& \quad+\sigma\left(\frac{1}{n} \sum_{i=1}^{n}\left(\left|x_{i}-y_{i}\right|+\left(\frac{n}{n-1}\right)^{\frac{1}{r}}|X-Y|_{r}\right)\left(1+\left|n \int_{A_{i}^{n}} \xi\right|\right)\right) .
\end{aligned}
$$

We conclude

$$
\begin{aligned}
& n^{-1}\left|\sum_{i=1}^{n} H\left(x_{i}, \mu_{\mathbf{x}}^{i}, n \int_{A_{i}^{n}} \xi\right)-\sum_{i=1}^{n} H\left(y_{i}, \mu_{\mathbf{y}}^{i}, n \int_{A_{i}^{n}} \eta\right)\right| \\
& \leq C\left(\int_{\Omega}|\xi-\eta|^{2}\right)^{\frac{1}{2}}+C\left(\left(\sum_{i=1}^{n} \int_{A_{i}^{n}}|\xi|^{2}\right)^{\frac{1}{2}}+\left(\sum_{i=1}^{n} \int_{A_{i}^{n}}|\eta|^{2}\right)^{\frac{1}{2}}\right)\left(\sum_{i=1}^{n} \int_{A_{i}^{n}}|\xi-\eta|^{2}\right)^{\frac{1}{2}} \\
& +\sigma\left(\left(\left(\frac{1}{n} \sum_{i=1}^{n}\left|x_{i}-y_{i}\right|^{2}\right)^{\frac{1}{2}}+2|X-Y|_{r}\right)\left(1+\left(\sum_{i=1}^{n} n\left|\int_{A_{i}^{n}} \eta\right|^{2}\right)^{\frac{1}{2}}\right)\right.
\end{aligned}
$$

Finally, we have

$$
\begin{gather*}
n^{-1}\left|\sum_{i=1}^{n} H\left(x_{i}, \mu_{\mathbf{x}}^{i}, n \int_{A_{i}^{n}} \xi\right) \sum_{i=1}^{n} H\left(y_{i}, \mu_{\mathbf{y}}^{i}, n \int_{A_{i}^{n}} \eta\right)\right| \leq C(1+|\xi|+|\eta|)|\xi-\eta|+\sigma(3|X-Y|(1+|\xi|)) .  \tag{4.2}\\
\text { Set } \\
X_{\epsilon}^{n}=\sum_{i=1}^{n} X_{\epsilon}\left(\frac{i}{n}\right) \mathbf{1}_{A_{i}^{n}}, \quad \mathbf{x}_{\epsilon}^{n}=\left(X_{\epsilon}(1 / n), X_{\epsilon}(2 / n), \ldots, X_{\epsilon}(1)\right) .
\end{gather*}
$$

Using the definition of viscosity subsolution and 4.2) we now have, for every $\epsilon>0$,
$0 \geq \partial_{t} \varphi\left(t_{n}, X^{n}\right)-\kappa \sum_{k=1}^{d}\left\langle D^{2} \varphi\left(t_{n}, \sum_{i=1}^{n} x_{i}^{n} \mathbf{1}_{A_{i}^{n}}\right) e_{k}, e_{k}\right\rangle+\frac{1}{n} \sum_{i=1}^{n} H\left(x_{i}^{n}, \mu_{\mathbf{x}^{n}}^{i}, n \int_{A_{i}^{n}} D \varphi\left(t_{n}, X^{n}\right)\right)+\mathcal{F}\left(\mu_{\mathbf{x}^{n}}\right)$.
Hence, for large $n$,

$$
\begin{aligned}
0 & \geq \partial_{t} \varphi\left(t_{0}, X_{0}\right)-\kappa \sum_{k=1}^{d}\left\langle D^{2} \varphi\left(t_{0}, X_{0}\right) e_{k}, e_{k}\right\rangle+\frac{1}{n} \sum_{i=1}^{n} H\left(X_{\epsilon}\left(\frac{i}{n}\right), \mu_{\mathbf{x}_{\epsilon}^{n}}^{i}, n \int_{A_{i}^{n}} P_{\epsilon}\right)+F\left(X^{n}\right) \\
& -\rho_{0}(n)-C\left(1+\left|D \varphi\left(t_{n}, X^{n}\right)\right|+\left|P_{\epsilon}\right|\right)\left|D \varphi\left(t_{n}, X^{n}\right)-P_{\epsilon}\right|-\sigma\left(3\left|X_{\epsilon}^{n}-X^{n}\right|\left(1+\left|P_{\epsilon}\right|\right)\right) \\
\geq & \partial_{t} \varphi\left(t_{0}, X_{0}\right)-\kappa \sum_{k=1}^{d}\left\langle D^{2} \varphi\left(t_{0}, X_{0}\right) e_{k}, e_{k}\right\rangle+\frac{1}{n} \sum_{i=1}^{n} H\left(X_{\epsilon}\left(\frac{i}{n}\right), \mu_{\mathbf{x}_{\epsilon}^{n}}^{i}, n \int_{A_{i}^{n}} P_{\epsilon}\right)+F\left(X_{0}\right) \\
& -m_{1}(\epsilon)-\rho_{0}(n)-C\left(1+2\left|P_{0}\right|+2 \epsilon\right) 2 \epsilon-\sigma\left(\left(3\left|X_{\epsilon}^{n}-X_{\epsilon}\right|+3 \epsilon\right) \mid\left(1+\left|P_{0}\right|+\epsilon\right)\right),
\end{aligned}
$$

where $\lim _{n \rightarrow \infty} \rho_{0}(n)=0$. Since $X_{\epsilon}, P_{\epsilon}$ are continuous on $\Omega$, it follows that

$$
\lim _{n \rightarrow \infty} \frac{1}{n} \sum_{i=1}^{n} H\left(X_{\epsilon}\left(\frac{i}{n}\right), \mu_{\mathbf{x}_{\epsilon}^{n}}^{i}, n \int_{A_{i}^{n}} P_{\epsilon}\right)=\tilde{H}\left(X_{\epsilon}, \operatorname{law}\left(X_{\epsilon}\right), P_{\epsilon}\right) .
$$

Thus, letting $n \rightarrow \infty$ above we obtain

$$
\begin{aligned}
& \partial_{t} \varphi\left(t_{0}, X_{0}\right)-\kappa \sum_{k=1}^{d}\left\langle D^{2} \varphi\left(t_{0}, X_{0}\right) e_{k}, e_{k}\right\rangle+\tilde{H}\left(X_{\epsilon}, \operatorname{law}\left(X_{\epsilon}\right), P_{\epsilon}\right) \\
& \leq m_{1}(\epsilon)+C\left(1+2\left|P_{0}\right|+2 \epsilon\right) 2 \epsilon+2 \sigma\left(3 \epsilon\left(1+\left|P_{0}\right|+\epsilon\right)\right) .
\end{aligned}
$$

Finally, letting $\epsilon \rightarrow 0$ we conclude that

$$
\partial_{t} \varphi\left(t_{0}, X_{0}\right)-\kappa \sum_{k=1}^{d}\left\langle D^{2} \varphi\left(t_{0}, X_{0}\right) e_{k}, e_{k}\right\rangle+\tilde{H}\left(X_{0}, \operatorname{law}\left(X_{0}\right), D \varphi\left(t_{0}, X_{0}\right)\right)+F\left(X_{0}\right) \leq 0 .
$$

Thus $V$ is a viscosity subsolution of (1.3). Reasoning in the same way, we can prove that $V$ is a supersolution of (1.3).
Example 4.1. The following is an example of a typical particle system that leads to simple equations of type (1.2). Let $G: \mathbb{R}^{d} \rightarrow \mathbb{R}^{d}$ be a bounded even function such that

$$
|G(x)-G(y)| \leq L|x-y| \quad \forall x, y \in \mathbb{R}^{d}
$$

and let $W$ be a standard Wiener process in $\mathbb{R}^{d}$. Let $T>0$. We consider a system of $n$ particles whose states are given by the SDE with common noise

$$
\left\{\begin{array}{l}
d X_{i}(s)=\frac{1}{n-1} \sum_{j \neq i} G\left(X_{i}(s)-X_{j}(s)\right) d s+\sqrt{2 \kappa} d W(s) \quad t \leq s \leq T, \\
X_{i}(t)=x_{i} .
\end{array}\right.
$$

If we define

$$
u_{n}(t, \mathbf{x})=\mathbb{E}\left[-\int_{t}^{T} \mathcal{F}\left(\frac{1}{n} \sum_{i=1}^{n} \delta_{X_{i}(s)}\right) d s+\mathcal{U}_{0}\left(\frac{1}{n} \sum_{i=1}^{n} \delta_{X_{i}(T)}\right)\right]
$$

(where $\mathbb{E}$ above is the expectation with respect to a probability measure on some reference probability space) then the function $u_{n}$ is the viscosity solution of the terminal value problem

$$
\left\{\begin{aligned}
&-\partial_{t} u_{n}-\kappa \operatorname{Tr}\left(A_{n} D^{2} u_{n}\right)-\frac{1}{n-1} \sum_{i=1}^{n} \sum_{j \neq i} G\left(x_{i}-x_{j}\right) \cdot D_{x_{i}} u_{n} \\
&+\mathcal{F}\left(\frac{1}{n} \sum_{i=1}^{n} \delta_{x_{i}}\right)=0 \\
& u_{n}\left(T, x_{1}, \cdots, x_{n}\right)=\mathcal{U}_{0}\left(\frac{1}{n} \sum_{i=1}^{n} \delta_{x_{i}}\right) \text { in }(0, T) \times\left(\mathbb{R}^{d}\right)^{n}, \\
& \text { on }\left(\mathbb{R}^{d}\right)^{n},
\end{aligned}\right.
$$

where $A_{n}$ is as in 1.2 . In this example the Hamiltonian $H$ is defined by

$$
H(x, \nu, p)=-p \cdot \int_{\mathbb{R}^{d}} G(x-y) \nu(d y)
$$

It is obvious that $H$ satisfies (2.1) and (2.3). We point out that the boundedness of $G$ is needed here to guarantee (2.3). Regarding (2.2), let $x, y, p \in \mathbb{R}^{d}, \mu, \nu \in \mathcal{P}_{r}\left(\mathbb{R}^{d}\right)$, and let $\gamma$ be a Borel probability measure on $\mathbb{R}^{d} \times \mathbb{R}^{d}$ with marginals $\mu, \nu$. Then

$$
\begin{aligned}
& |H(x, \mu, p)-H(y, \nu, p)| \leq L|x-y||p|+\left|\int_{\mathbb{R}^{d}} G(x-z) \mu(d y) \cdot p-\int_{\mathbb{R}^{d}} G(x-w) \nu(d w) \cdot p\right| \\
& \quad \leq L|x-y||p|+\left|\int_{\mathbb{R}^{d}}(G(x-z)-G(x-w)) \gamma(d z, d w)\right||p| \\
& \quad \leq L|x-y||p|+L|p| \int_{\mathbb{R}^{d}}|z-w| \gamma(d z, d w) \\
& \quad \leq L|p|\left(|x-y|+\left(\int_{\mathbb{R}^{d}}|z-w|^{r} \gamma(d z, d w)\right)^{\frac{1}{r}}\right)
\end{aligned}
$$

Since this holds for every $\gamma$ we thus obtain

$$
|H(x, \mu, p)-H(y, \nu, p)| \leq L|p|\left(|x-y|+d_{r}(\mu, \nu)\right)
$$

5. L-VISCOSity solutions Versus viscosity solutions on the Wasserstein space

In this section, we consider either

$$
\mathcal{U}:[0, T) \times \mathcal{P}_{2}\left(\mathbb{R}^{d}\right) \rightarrow \mathbb{R} \quad \text { or } \quad \mathcal{U}: \mathcal{P}_{2}\left(\mathbb{R}^{d}\right) \rightarrow \mathbb{R}
$$

and

$$
U:[0, T) \times E \rightarrow \mathbb{R} \quad \text { or } \quad U: E \rightarrow \mathbb{R}
$$

such that $U(t, X)=\mathcal{U}\left(t, X_{\sharp} \mathcal{L}_{1}\right)$ or $U(X)=\mathcal{U}\left(X_{\sharp} \mathcal{L}_{1}\right)$. Recall

$$
E:=L^{2}\left((0,1), \mathbb{R}^{d}\right)
$$

When $\mathcal{U}$ is differentiable at $\mu \in \mathcal{P}_{2}\left(\mathbb{R}^{d}\right), \nabla_{\mu} \mathcal{U}(\mu)$, the Wasserstein gradient of $\mathcal{U}$ at $\mu$ satisfies

$$
\begin{equation*}
\nabla_{\mu} \mathcal{U}(\mu) \in T_{\mu} \mathcal{P}_{2}\left(\mathbb{R}^{d}\right) \tag{5.3}
\end{equation*}
$$

Assume next that $\nabla_{\mu} \mathcal{U}$ is differentiable at $(q, \mu)$ in the sense of [37] and $\nabla_{\mu \mu} \mathcal{U}(\mu)(q, \cdot)$, the Wasserstein gradient of $\nabla_{\mu} \mathcal{U}(\mu)$ at $\mu$ belongs to $L_{\mu \otimes \mu}^{\infty}\left(\mathbb{R}^{2 d}, \mathbb{R}^{d \times d}\right)$. We have

$$
\begin{equation*}
\left(\nabla_{\mu \mu} \mathcal{U}(\mu)\right)^{T}(q, x)=\left(\nabla_{\mu \mu} \mathcal{U}(\mu)\right)(x, q) \quad \forall q, x \in \mathbb{R}^{d} \tag{5.4}
\end{equation*}
$$

If we assume that $\mathcal{U}$ is twice differentiable in the sense of [37], then the map $(q, \mu) \mapsto$ $\nabla_{\mu} \mathcal{U}(\mu)(q)$ has a first order Taylor expansion on an appropriate set [24]. Furthermore, $\nabla_{\mu} \mathcal{U}(\mu)$ is Lipschitz on the support of $\mu$ and there exists a symmetric matrix $A_{1}(\mu) \in L_{\mu}^{\infty}\left(\mathbb{R}^{d}, \mathbb{R}^{d \times d}\right)$ which coincides almost everywhere with $\nabla_{q}\left(\nabla_{\mu} \mathcal{U}(\mu)\right)$. In this manuscript, we adopt the notation and terminology of [37] by defining $\nabla_{w}^{2} \mathcal{U}(\mu)$, the Wasserstein second differential of $\mathcal{U}$ at $\mu$ as

$$
\nabla_{w}^{2} \mathcal{U}(\mu)\left(\xi, \xi_{*}\right)=\int_{\mathbb{R}^{d}} A_{1}(\mu)(q) \xi(q) \cdot \xi_{*}(q) \mu(d q)+\int_{\mathbb{R}^{2 d}} \nabla_{\mu \mu} \mathcal{U}(\mu)\left(q_{*}, q\right) \xi(q) \cdot \xi_{*}\left(q_{*}\right) \mu(d q) \mu\left(d q_{*}\right)
$$

if $\xi, \xi_{*} \in L_{\mu}^{2}\left(\mathbb{R}^{d}, \mathbb{R}^{d}\right)$. Note that the ordering $\left(q_{*}, q\right)$ in the last integral is not a typo. By abuse of notation, we identify the bilinear forms $\nabla_{w}^{2} \mathcal{U}(\mu)$ with the operators

$$
\xi \longmapsto A_{1}(\mu) \xi+\int_{\mathbb{R}^{d}} \nabla_{\mu \mu} \mathcal{U}(\mu)(\cdot, q) \xi(q) \mu(d q)
$$

which, by $(5.4)$ and the fact that $A_{1}$ is symmetric $\mu$-almost everywhere, turns out to be selfadjoint on $L_{\mu}^{2}\left(\mathbb{R}^{d}, \mathbb{R}^{d}\right)$.

The relation $U(X)=\mathcal{U}\left(X_{\sharp} \mathcal{L}_{1}\right)$ expresses the fact that $U$ is invariant under the set of maps which preserve Lebesgue measure. This is what imposes a special structure on the second differential of $U$ at $X$ when it exists. When $U$ is twice differentiable at $X$ then for any $\zeta \in E$, $D^{2} U(X)(\zeta)(\cdot)$ belongs to $E$ and

$$
D^{2} U(X)(\zeta)(\cdot)=A_{1}(X) \zeta+\int_{(0,1)} \nabla_{\mu \mu} \mathcal{U}(\mu)(X, X(\omega)) \zeta(\omega) d \omega
$$

Given an arbitrary orthonormal basis $\left\{e_{1}, \cdots, e_{d}\right\}$ of $\mathbb{R}^{d}$, we identify each $e_{k}$ with the constant function which assumes the value $e_{k}$ everywhere. Abusing notation we write $e_{k}: \mathbb{R}^{d} \rightarrow \mathbb{R}^{d}$. Note that if $X \in E$ then $e_{k} \circ X \equiv e_{k}$ and so, we may also consider $e_{k}$ to be the constant function $e_{k}:(0,1) \rightarrow \mathbb{R}^{d}$. If $E_{0}$ is the finite dimensional space spanned by $\left\{e_{1}, \cdots, e_{d}\right\}$, we have the orthogonal decomposition

$$
E=E_{0} \oplus E_{0}^{\perp} .
$$

The partial trace of the operator $\zeta \rightarrow D^{2} U(X)(\zeta)$ on $E_{0}$ is

$$
\triangle_{\mathbb{R}^{d}} U(X)=\sum_{k=1}^{d}\left\langle D^{2} U(X) e_{k}, e_{k}\right\rangle .
$$

We have the relation

$$
\triangle_{\mathbb{R}^{d}} U(X)=\Delta_{w} \mathcal{U}(\mu) \circ X
$$

where, $\Delta_{w}$ is the partial Wasserstein Laplacian [24]. This relation will allow to compare viscosity solutions on the Wasserstein space and the Hilbert space.
5.1. Domain of definition of Wasserstein Hessian. In this section, we denote by $\pi^{1}, \pi^{2}$ : $\mathbb{R}^{2 d} \rightarrow \mathbb{R}^{d}$ the coordinate projection maps of $\mathbb{R}^{2 d}$ onto $\mathbb{R}^{d}$. Given a positive integer $D$, we denote by $C_{b}^{3}\left(\mathbb{R}^{D}\right)$ the set of $f \in C^{3}\left(\mathbb{R}^{D}\right)$ that have bounded second and third order derivatives.

We start by recalling a few facts about the Wasserstein tangent spaces $T_{\mu} \mathcal{P}_{2}\left(\mathbb{R}^{d}\right)$. Let $\xi_{0} \in$ $L_{\mu}^{2}\left(\mathbb{R}^{d}, \mathbb{R}^{d}\right)$. Note that $\xi_{0} \in T_{\mu} \mathcal{P}_{2}\left(\mathbb{R}^{d}\right)$ if and only if

$$
\begin{aligned}
& \lim _{r \rightarrow 0^{+}} \inf _{\phi} \sup _{(\nu, \gamma)}\left\{\frac{\left|\int_{\mathbb{R}^{2 d}}\left(\xi_{0}\left(q_{1}\right)-\nabla \phi\left(q_{1}\right)\right) \cdot\left(q_{2}-q_{1}\right) \gamma\left(d q_{1}, q_{2}\right)\right|}{\left\|\pi^{2}-\pi^{1}\right\|_{L^{2}(\gamma)}}: 0<\left\|\pi^{2}-\pi^{1}\right\|_{L^{2}(\gamma)} \leq r, \gamma \in \Gamma(\mu, \nu)\right\} \\
& =0
\end{aligned}
$$

Here, the infimum is performed over the set $C_{c}^{\infty}\left(\mathbb{R}^{d}\right)$ or equivalently over the set $C_{b}^{3}\left(\mathbb{R}^{d}\right)$. The space $T_{\mu} \mathcal{P}_{2}\left(\mathbb{R}^{d}\right)$, being a Hilbert space, can be identified with the co-tangent space. The Wasserstein gradient of a function $\mathcal{U}: \mathcal{P}_{2}\left(\mathbb{R}^{d}\right) \rightarrow \mathbb{R}$ at $\mu$ is the element of minimal norm in the subdifferential of $\mathcal{U}$ at $\mu$ and so, it belongs to $T_{\mu} \mathcal{P}_{2}\left(\mathbb{R}^{d}\right)$.

We would like to propose an analogous characterization for all Wasserstein derivatives of order less than or equal to 2 . We say that $f: \mathbb{R}^{2 d} \rightarrow \mathbb{R}$ is symmetric if $f\left(q_{1}, q_{2}\right)=f\left(q_{2}, q_{1}\right)$ for all $q_{1}, q_{2} \in \mathbb{R}^{d}$.

Given $\psi \in C_{b}^{3}\left(\mathbb{R}^{2 d}\right)$ which is symmetric and $\phi \in C_{b}^{3}\left(\mathbb{R}^{d}\right)$, we define

$$
\mathcal{V}_{(\phi, \psi)}^{\mu}(\nu):=\int_{\mathbb{R}^{d}} \phi(q)(\nu-\mu)(d q)+\frac{1}{2} \int_{\mathbb{R}^{2 d}} \psi\left(q_{1}, q_{2}\right)(\nu-\mu)\left(d q_{1}\right)(\nu-\mu)\left(d q_{2}\right)
$$

for $\mu, \nu \in \mathcal{P}_{2}\left(\mathbb{R}^{d}\right)$. The function $\mathcal{V}_{(\phi, \psi)}^{\mu}$ is twice differentiable in the sense of 37],

$$
\nabla_{\mu} \mathcal{V}_{(\phi, \psi)}^{\mu}(\nu)\left(q_{1}\right)=\nabla \phi\left(q_{1}\right)+\int_{\mathbb{R}^{d}} \nabla_{q_{1}} \psi\left(q_{1}, b\right)(\nu-\mu)(d b)
$$

and so,

$$
\nabla_{q_{1}}\left(\nabla_{\mu} \mathcal{V}_{(\phi, \psi)}^{\mu}(\nu)\left(q_{1}\right)\right)=\nabla^{2} \phi\left(q_{1}\right)+\int_{\mathbb{R}^{d}} \nabla_{q_{1} q_{1}} \psi\left(q_{1}, b\right)(\nu-\mu)(d b)
$$

We conclude

$$
\nabla_{\mu \mu} \mathcal{V}_{(\phi, \psi)}^{\mu}(\nu)(q)=\nabla_{q_{2} q_{1}} \psi\left(q_{1}, q_{2}\right)
$$

Note that if $X, Y \in E$ are such that $X_{\sharp} \mathcal{L}_{(0,1)}^{1}=\mu$ and $Y_{\sharp} \mathcal{L}_{(0,1)}^{1}=\nu$, then

$$
\begin{equation*}
\mathcal{V}_{(\phi, \psi)}^{\mu}(\nu)=V_{(\phi, \psi)}^{X}(Y) \tag{5.5}
\end{equation*}
$$

where we have set

$$
\begin{aligned}
V_{(\phi, \psi)}^{X}(Y) & =\int_{(0,1)}(\phi(Y(\omega))-\phi(X(\omega))) d \omega \\
& +\frac{1}{2} \int_{(0,1)^{2}}(\psi(Y(\omega), Y(o))+\psi(X(\omega), X(o))-2 \psi(Y(\omega), X(o))) d \omega d o .
\end{aligned}
$$

Remark 5.1. Let $X \in E$, let $\phi \in C_{b}^{3}\left(\mathbb{R}^{d}\right)$ let $\psi \in C_{b}^{3}\left(\mathbb{R}^{2 d}\right)$ be symmetric and set $V:=V_{(\phi, \psi)}^{X}$.
(i) Note that if $X^{*}, Y, Y^{*} \in E$ are such that $X$ and $X^{*}$ have the same law and $Y$ and $Y^{*}$ have the same law then (5.5) implies $V_{(\phi, \psi)}^{X}(Y)=V_{(\phi, \psi)}^{X^{*}}\left(Y^{*}\right)$.
(ii) The function $V$ is Fréchet differentiable everywhere on $E$ and for any $Y \in E$ we have

$$
D V(Y)(\omega)=\nabla \phi(Y(\omega))+\int_{(0,1)}\left(\nabla_{q_{1}} \psi(Y(\omega), Y(o))-\nabla_{q_{1}} \psi(Y(\omega), X(o))\right) d o
$$

(iii) The subset of $E$ where the function $D V$ is Fréchet differentiable may not be $E$ (cf. [37]), unless $\phi$ and $\psi$ are polynomials of degree 2. However, $D V$ is Gateaux differentiable everywhere on $E$ and for any $Y \in E$ we have

$$
\begin{aligned}
D^{2} V(Y)(\zeta, \zeta) & =\int_{(0,1)} \nabla^{2} \phi(X) \zeta \cdot \zeta d \omega+\int_{(0,1)^{2}} \nabla_{q_{2} q_{1}} \psi(Y(\omega), Y(o)) \zeta(\omega) \cdot \zeta(o) d \omega d o \\
& +\int_{(0,1)^{2}}\left(\nabla_{q_{1} q_{1}} \psi(Y(\omega), Y(o))-\nabla_{q_{1} q_{1}} \psi(Y(\omega), X(o))\right) \zeta(\omega) \cdot \zeta(\omega) d \omega d o
\end{aligned}
$$

for any $\zeta \in E$. In particular, $D V(X)=\nabla \phi \circ X$ and the operator $D^{2} V(X)$ is given by

$$
D^{2} V(X) \zeta=\nabla^{2} \phi(X) \zeta+\int_{(0,1)} \nabla_{q_{2} q_{1}} \psi(X, X(\omega)) \zeta(\omega) d \omega
$$

for $\zeta \in E$.
Lemma 5.2. Let $X, X^{*} \in E$, let $\phi \in C_{b}^{3}\left(\mathbb{R}^{d}\right)$ and let $\psi \in C_{b}^{3}\left(\mathbb{R}^{2 d}\right)$ be symmetric. If $X_{\sharp} \mathcal{L}_{(0,1)}^{1}=$ $X_{\sharp}^{*} \mathcal{L}_{(0,1)}^{1}$ then $V^{X}:=V_{(\phi, \psi)}^{X}$ is twice Fréchet differentiable at $X$ if and only if $V^{X^{*}}:=V_{(\phi, \psi)}^{X^{*}}$ is twice Fréchet differentiable at $X^{*}$.

Proof. We only need to prove one direction of the Lemma since the converse direction could be obtained by symmetry. Assume $V^{X}$ is twice Fréchet differentiable at $X$. By assumption, there exists a function $\epsilon: \mathbb{R} \rightarrow \mathbb{R}$, monotone nondecreasing, continuous at 0 and such that $\epsilon(0)=0$ and there exists a function $\epsilon_{0}: E \rightarrow \mathbb{R}$ such that $\left|\epsilon_{0}(h)\right| \leq \epsilon(|h|)$ and

$$
\begin{align*}
V^{X}(X+h) & =\int_{(0,1)} \nabla \phi(X) \cdot h d \omega+\frac{1}{2} \int_{(0,1)} \nabla^{2} \phi(X) h \cdot h d \omega \\
& +\frac{1}{2} \int_{(0,1)^{2}} \nabla_{q_{2} q_{1}} \psi(X(o), X(\omega)) h(\omega) \cdot h(o) d \omega d o+|h|^{2} \epsilon_{0}(h) \tag{5.6}
\end{align*}
$$

Since $X$ and $X^{*}$ have the same laws, it is well-known that we can choose a sequence of Borel functions $S_{n}:[0,1] \rightarrow[0,1]$ which are one-to-one, onto, such that $S_{n \sharp} \mathcal{L}_{(0,1)}^{1}=\mathcal{L}_{(0,1)}^{1}$ and such that

$$
\lim _{n}\left|X^{*}-X \circ S_{n}\right|=0
$$

In light of Remark 5.1(i), (5.6) implies

$$
\begin{aligned}
V^{X \circ S_{n}}\left(X \circ S_{n}+h\right) & =V^{X}\left(X+h \circ S_{n}^{-1}\right) \\
& =\int_{(0,1)} \nabla \phi(X) \cdot h \circ S_{n}^{-1} d \omega+\frac{1}{2} \int_{(0,1)} \nabla^{2} \phi(X) h \circ S_{n}^{-1} \cdot h \circ S_{n}^{-1} d \omega \\
& +\frac{1}{2} \int_{(0,1)^{2}} \nabla_{q_{2} q_{1}} \psi(X(o), X(\omega)) h \circ S_{n}^{-1}(\omega) \cdot h \circ S_{n}^{-1}(o) d \omega d o+\left|h \circ S_{n}^{-1}\right|^{2} \epsilon_{0}\left(h \circ S_{n}^{-1}\right) \\
& =\int_{(0,1)} \nabla \phi\left(X \circ S_{n}\right) \cdot h d \omega+\frac{1}{2} \int_{(0,1)} \nabla^{2} \phi\left(X \circ S_{n}\right) h \cdot h d \omega \\
& +\frac{1}{2} \int_{(0,1)^{2}} \nabla_{q_{2} q_{1}} \psi\left(X \circ S_{n}(o), X \circ S_{n}(\omega)\right) h(\omega) \cdot h(o) d \omega d o+|h|^{2} \epsilon_{0}\left(h \circ S_{n}^{-1}\right)
\end{aligned}
$$

We use again Remark 5.1(i) to conclude that

$$
\begin{aligned}
V^{X^{*}}\left(X \circ S_{n}+h\right) & =\int_{(0,1)} \nabla \phi\left(X \circ S_{n}\right) \cdot h d \omega+\frac{1}{2} \int_{(0,1)} \nabla^{2} \phi\left(X \circ S_{n}\right) h \cdot h d \omega \\
& +\frac{1}{2} \int_{(0,1)^{2}} \nabla_{q_{2} q_{1}} \psi\left(X \circ S_{n}(o), X \circ S_{n}(\omega)\right) h(\omega) \cdot h(o) d \omega d o+|h|^{2} \epsilon_{0}\left(h \circ S_{n}^{-1}\right)
\end{aligned}
$$

We let $n$ tend to $\infty$ to obtain

$$
\begin{aligned}
V^{X^{*}}\left(X^{*}+h\right) & =\int_{(0,1)} \nabla \phi\left(X^{*}\right) \cdot h d \omega+\frac{1}{2} \int_{(0,1)} \nabla^{2} \phi\left(X^{*}\right) h \cdot h d \omega \\
& +\frac{1}{2} \int_{(0,1)^{2}} \nabla_{q_{2} q_{1}} \psi\left(X^{*}(o), X^{*}(\omega)\right) h(\omega) \cdot h(o) d \omega d o+|h|^{2} \lim _{n \rightarrow \infty} \epsilon_{0}\left(h \circ S_{n}^{-1}\right)
\end{aligned}
$$

We use the fact that $\lim _{n \rightarrow \infty}\left|\epsilon_{0}\left(h \circ S_{n}^{-1}\right)\right| \leq \epsilon(|h|)$ to conclude that $V^{X^{*}}$ is twice Fréchet differentiable at $X^{*}$.

Definition 5.3. Let $\mu \in \mathcal{P}_{2}\left(\mathbb{R}^{d}\right)$, $\xi_{0} \in T_{\mu} \mathcal{P}_{2}\left(\mathbb{R}^{d}\right)$, let $A_{1} \in L_{\mu}^{\infty}\left(\mathbb{R}^{d}, \mathbb{R}^{d \times d}\right)$ be symmetric $\mu$-almost everywhere and let $A_{2} \in L_{\mu \otimes \mu}^{\infty}\left(\mathbb{R}^{2 d}, \mathbb{R}^{d \times d}\right)$ be such that $A_{2}\left(q_{1}, q_{2}\right)=A_{2}^{T}\left(q_{2}, q_{1}\right)$, $\mu \otimes$-almost everywhere. We say that $\left(\xi_{0}, A_{1}, A_{2}\right)$ belongs to $T_{\mu}^{*, 2} \mathcal{P}_{2}\left(\mathbb{R}^{d}\right)$ if
$\lim _{r \rightarrow 0^{+}} \inf _{(\phi, \psi)} \sup _{(\nu, \gamma)}\left\{\frac{\left|e_{0}\left(\mu, \gamma, \xi_{0}-\nabla \phi, A_{1}-\nabla^{2} \phi, A_{2}-\nabla_{q_{2} q_{1}} \psi\right)\right|}{\left\|\pi^{2}-\pi^{1}\right\|_{L^{2}(\gamma)}^{2}}: 0<\left\|\pi^{2}-\pi^{1}\right\|_{L^{2}(\gamma)} \leq r, \gamma \in \Gamma(\mu, \nu)\right\}=0$,
where the infimum is performed over the set of pairs $(\phi, \psi)$ such that $\phi \in C_{b}^{3}\left(\mathbb{R}^{d}\right), \psi \in C_{b}^{3}\left(\mathbb{R}^{2 d}\right)$ is symmetric and $V_{(\phi, \psi)}^{X}$ is twice Fréchet differentiable at $X$ with $\operatorname{law}(X)=\mu$.

Here we have set

$$
\begin{aligned}
e_{0}\left(\mu, \gamma, \xi_{0}, A_{1}, A_{2}\right) & :=\int_{\mathbb{R}^{2 d}}\left(\xi_{0}\left(q_{1}\right)+\frac{1}{2} A_{1}\left(q_{1}\right)\left(q_{2}-q_{1}\right)\right) \cdot\left(q_{2}-q_{1}\right) \gamma\left(d q_{1}, q_{2}\right) \\
& +\frac{1}{2} \int_{\mathbb{R}^{2 d} \times \mathbb{R}^{2 d}} A_{2}\left(q_{1}, z\right)(w-z) \cdot\left(q_{2}-q_{1}\right) \gamma\left(d q_{1}, q_{2}\right) \gamma(d z, d w)
\end{aligned}
$$

We shall later use the expression

$$
E\left(r, \mu, \xi_{0}, A_{1}, A_{2}\right):=\sup _{\nu} \sup _{\gamma}\left\{\frac{\left|e_{0}\left(\mu, \gamma, \xi_{0}, A_{1}, A_{2}\right)\right|}{\left\|\pi^{2}-\pi^{1}\right\|_{L^{2}(\gamma)}^{2}}\right\}
$$

where the supremum is performed over the set of pairs $(\nu, \gamma)$ such that $0<W_{2}^{2}(\nu, \mu) \leq r$ and $\gamma \in \Gamma(\mu, \nu)$.
5.2. Specific expression for superjets/subjets. For $\mu, \nu, \gamma, \xi_{0}, A_{1}, A_{2}$ as in Definition 5.3 and $t, s \in(0, T), a \in \mathbb{R}$, we set

$$
\begin{aligned}
e\left(\mathcal{U}, s, t, a, \mu, \nu, \gamma, \xi_{0}, A_{1}, A_{2}\right) & :=\mathcal{U}(s, \nu)-\mathcal{U}(t, \mu)-a(s-t)-\int_{\mathbb{R}^{2 d}} \xi_{0}\left(q_{1}\right) \cdot\left(q_{2}-q_{1}\right) \gamma\left(d q_{1}, d q_{2}\right) \\
& -\frac{1}{2} \int_{\mathbb{R}^{2 d}} A_{1}\left(q_{1}\right)\left(q_{2}-q_{1}\right) \cdot\left(q_{2}-q_{1}\right) \gamma(d x, d y) \\
& -\frac{1}{2} \int_{\mathbb{R}^{2 d} \times \mathbb{R}^{2 d}} A_{2}\left(q_{1}, z\right)(w-z) \cdot\left(q_{2}-q_{1}\right) \gamma(d z, d w) \gamma\left(d q_{1}, d q_{2}\right) .
\end{aligned}
$$

Similarly, for

$$
t, s \in(0, T), a \in \mathbb{R}, \quad X, Y, \zeta_{0} \in E, \quad \bar{A}_{1} \in L^{\infty}\left((0,1), \mathbb{R}^{d \times d}\right), \quad \bar{A}_{2} \in L^{\infty}\left((0,1)^{2}, \mathbb{R}^{d \times d}\right)
$$

we set

$$
\begin{aligned}
\bar{e}\left(U, s, t, a, X, Y, \zeta_{0}, \bar{A}_{1}, \bar{A}_{2}\right) & :=U(s, Y)-U(t, X)-a(s-t)-\int_{(0,1)} \zeta_{0} \cdot(Y-X) d \omega \\
& -\frac{1}{2} \int_{(0,1)} \bar{A}_{1}(Y-X) \cdot(Y-X) d \omega \\
& -\frac{1}{2} \int_{(0,1)^{2}} \bar{A}_{2}(\omega, o)(Y(o)-X(o)) \cdot(Y(\omega)-X(\omega)) d o d \omega
\end{aligned}
$$

If the functions $\mathcal{U}$ and $U$ are independent of $t$, the right hand sides of the above expressions do not have the $a(s-t)$ term and we will write $e\left(\mathcal{U}, \mu, \nu, \gamma, \xi_{0}, A_{1}, A_{2}\right)$ and $\bar{e}\left(U, X, Y, \zeta_{0}, \bar{A}_{1}, \bar{A}_{2}\right)$.
Definition 5.4. let $t \in(0, T)$ and $\mu \in \mathcal{P}_{2}\left(\mathbb{R}^{d}\right)$ and let $\mathcal{U}:[0, T) \times \mathcal{P}_{2}\left(\mathbb{R}^{d}\right) \rightarrow \mathbb{R}$.
(i) We define the parabolic second order subjet of $\mathcal{U}$ at $(t, \mu)$ to be the set $\mathcal{P}^{2,-} \mathcal{U}(t, \mu)$, which consists of $\left(a, \xi_{0}, A_{1}, A_{2}\right) \in \mathbb{R} \times T_{\mu}^{*, 2} \mathcal{P}_{2}\left(\mathbb{R}^{d}\right)$ such that

$$
\liminf _{(s, \nu) \rightarrow(t, \mu)} \inf _{\gamma \in \Gamma(\mu, \nu)} \frac{e\left(\mathcal{U}, s, t, \mu, \nu, \gamma, \xi_{0}, A_{1}, A_{2}\right)}{|s-t|+W_{2}^{2}(\nu, \mu)} \geq 0
$$

(ii) We define the parabolic second order superjet of $\mathcal{U}$ at $(t, \mu)$ to be the set $\mathcal{P}^{2,+} \mathcal{U}(t, \mu)$, of $\left(a, \xi_{0}, A_{1}, A_{2}\right) \in \mathbb{R} \times T_{\mu}^{*, 2} \mathcal{P}_{2}\left(\mathbb{R}^{d}\right)$ such that

$$
\limsup _{(s, \nu) \rightarrow(t, \mu)} \inf _{\gamma \in \Gamma(\mu, \nu)} \frac{e\left(\mathcal{U}, s, t, \mu, \nu, \gamma, \xi_{0}, A_{1}, A_{2}\right)}{|s-t|+W_{2}^{2}(\nu, \mu)} \leq 0
$$

We set

$$
S_{\left(A_{1}, A_{2}\right)}(X)(h)=A_{1}(X) h+\int_{(0,1)} A_{2}(X, X(\omega)) h(\omega) d \omega, \quad \text { for } h \in E .
$$

Owing to the properties of $A_{1}$ and $A_{2}$ in Definition 5.3, $S_{\left(A_{1}, A_{2}\right)}(X)$ is a self-adjoint operator on E.

Lemma 5.5. Let $\mu \in \mathcal{P}_{2}\left(\mathbb{R}^{d}\right)$ and let $X \in E$ be such that $X_{\sharp} \mathcal{L}_{(0,1)}^{1}=\mu$.
(i) If $\left(a, \xi_{0}, A_{1}, A_{2}\right) \in \mathcal{P}^{2,-} \mathcal{U}(t, \mu)$ then

$$
\left(a, \xi_{0}(X), S_{\left(A_{1}, A_{2}\right)}(X)\right) \in \mathcal{P}^{2,-} U(t, X)
$$

(ii) If $\left(a, \xi_{0}, A_{1}, A_{2}\right) \in \mathcal{P}^{2,+} \mathcal{U}(t, \mu)$ then

$$
\left(a, \xi_{0}(X), S_{\left(A_{1}, A_{2}\right)}(X)\right) \in \mathcal{P}^{2,+} U(t, X)
$$

Proof. It suffices to prove (i). Let us assume that $\left(a, \xi_{0}, A_{1}, A_{2}\right) \in \mathcal{P}^{2,-} \mathcal{U}(t, \mu)$. Let $Y \in E$ and set $Y_{\sharp} \mathcal{L}_{(0,1)}^{1}=\nu$. Choose first $\gamma \in \Gamma_{0}(\mu, \nu)$ and then choose $X^{*}, Y^{*} \in E$ such that

$$
\gamma:=\left(X^{*} \times Y^{*}\right)_{\sharp} \mathcal{L}_{(0,1)}^{1} .
$$

Note

$$
p:=(X \times Y)_{\sharp} \mathcal{L}_{(0,1)}^{1} \in \Gamma(\mu, \nu)
$$

and

$$
\begin{aligned}
\bar{e}\left(U, s, t, a, X, Y, \xi_{0} \circ X, A_{1} \circ X, A_{2} \circ(X \times X)\right) & =U(s, Y)-U(t, X)-a(s-t) \\
& -e_{0}\left(\mu, p, \xi_{0}, A_{1}, A_{2}\right) .
\end{aligned}
$$

Hence,

$$
\begin{align*}
\bar{e}\left(U, s, t, a, X, Y, \xi_{0} \circ X, A_{1} \circ X, A_{2} \circ(X \times X)\right) & =\mathcal{U}(s, \nu)-\mathcal{U}(t, \mu)-a(s-t) \\
& -e_{0}\left(\mu, p, \xi_{0}, A_{1}, A_{2}\right), \tag{5.7}
\end{align*}
$$

where $A_{2} \circ(X \times X)$ denotes the function $A_{2} \circ(X \times X)(\omega, o)=A_{2}(X(\omega), X(o))$.
Fix for a moment a symmetric function $\psi \in C_{b}^{3}\left(\mathbb{R}^{2 d}\right)$ and $\phi \in C_{b}^{3}\left(\mathbb{R}^{d}\right)$ such that $V_{(\phi, \psi)}^{X}$ is twice Fréchet differentiable at $X$. Recall that by (5.6)

Since the first marginal of $\gamma$ is $\mu, X$ and $X^{*}$ have the same laws and so, by Lemma 5.2, $V_{(\phi, \psi)}^{X^{*}}$ is twice Fréchet differentiable at $X^{*}$ and so,

$$
\begin{equation*}
V_{(\phi, \psi)}^{X^{*}}\left(Y^{*}\right)-e_{0}\left(\mu, \gamma, \nabla \phi, \nabla^{2} \phi, \nabla_{q_{2} q_{1}} \psi\right)=o\left(\left\|Y^{*}-X^{*}\right\|^{2}\right) . \tag{5.9}
\end{equation*}
$$

Using (5.7), we have the decomposition

$$
\begin{aligned}
& \bar{e}\left(U, s, t, a, X, Y, \xi_{0} \circ X, A_{1} \circ X, A_{2} \circ(X \times X)\right) \\
& =\mathcal{U}(s, \nu)-\mathcal{U}(t, \mu)-a(s-t)-e_{0}\left(\mu, p, \xi_{0}-\nabla \phi, A_{1}-\nabla^{2} \phi, A_{2}-\nabla_{q_{2} q_{1}} \psi\right) \\
& -e_{0}\left(\mu, p, \nabla \phi, \nabla^{2} \phi, \nabla_{q_{2} q_{1}} \psi\right)
\end{aligned}
$$

Thus,

$$
\begin{aligned}
& \bar{e}\left(U, s, t, a, X, Y, \xi_{0} \circ X, A_{1} \circ X, A_{2} \circ(X \times X)\right) \\
& =\mathcal{U}(s, \nu)-\mathcal{U}(t, \mu)-a(s-t)-e_{0}\left(\mu, p, \xi_{0}-\nabla \phi, A_{1}-\nabla^{2} \phi, A_{2}-\nabla_{q_{2} q_{1}} \psi\right) \\
& -e_{0}\left(\mu, p, \nabla \phi, \nabla^{2} \phi, \nabla_{q_{2} q_{1}} \psi\right)-e_{0}\left(\mu, \gamma, \xi_{0}, A_{1}, A_{2}\right) \\
& +e_{0}\left(\mu, \gamma, \xi_{0}-\nabla \phi, A_{1}-\nabla^{2} \phi, A_{2}-\nabla_{q_{2} q_{1}} \psi\right) \\
& +e_{0}\left(\mu, \gamma, \nabla \phi, \nabla^{2} \phi, \nabla_{q_{2} q_{1}} \psi\right) .
\end{aligned}
$$

Rearranging, we obtain

$$
\begin{aligned}
\bar{e}\left(U, s, t, a, X, Y, \xi_{0} \circ X, A_{1} \circ X, A_{2} \circ(X \times X)\right) & =\mathcal{U}(s, \nu)-\mathcal{U}(t, \mu)-a(s-t)-e_{0}\left(\mu, \gamma, \xi_{0}, A_{1}, A_{2}\right) \\
& -e_{0}\left(\mu, p, \xi_{0}-\nabla \phi, A_{1}-\nabla^{2} \phi, A_{2}-\nabla_{q_{2} q_{1}} \psi\right) \\
& -e_{0}\left(\mu, p, \nabla \phi, \nabla^{2} \phi, \nabla_{q_{2} q_{1}} \psi\right) \\
& +e_{0}\left(\mu, \gamma, \nabla \phi, \nabla^{2} \phi, \nabla_{q_{2} q_{1}} \psi\right) \\
& +e_{0}\left(\mu, \gamma, \xi_{0}-\nabla \phi, A_{1}-\nabla^{2} \phi, A_{2}-\nabla_{q_{2} q_{1}} \psi\right) .
\end{aligned}
$$

Using Remark 5.1(i) we conclude that

$$
\begin{aligned}
\bar{e}\left(U, s, t, a, X, Y, \xi_{0} \circ X, A_{1} \circ X, A_{2} \circ(X \times X)\right) & =\mathcal{U}(s, \nu)-\mathcal{U}(t, \mu)-a(s-t)-e_{0}\left(\mu, \gamma, \xi_{0}, A_{1}, A_{2}\right) \\
& -e_{0}\left(\mu, p, \xi_{0}-\nabla \phi, A_{1}-\nabla^{2} \phi, A_{2}-\nabla_{q_{2} q_{1}} \psi\right) \\
& V_{(\phi, \psi)}^{X^{*}}\left(Y^{*}\right)-e_{0}\left(\mu, p, \nabla \phi, \nabla^{2} \phi, \nabla_{q_{2} q_{1}} \psi\right) \\
& +e_{0}\left(\mu, \gamma, \nabla \phi, \nabla^{2} \phi, \nabla_{q_{2} q_{1}} \psi\right)-V_{(\phi, \psi)}^{X}(Y) \\
& +e_{0}\left(\mu, \gamma, \xi_{0}-\nabla \phi, A_{1}-\nabla^{2} \phi, A_{2}-\nabla_{q_{2} q_{1}} \psi\right) .
\end{aligned}
$$

We first use the fact that $\left(a, \xi_{0}, A_{1}, A_{2}\right) \in \mathcal{P}^{-} \mathcal{U}(t, \mu)$, second use (5.8) and (5.9) and third use the fact that $\left(\xi_{0}, A_{1}, A_{2}\right) \in T_{\mu}^{*, 2} \mathcal{P}_{2}\left(\mathbb{R}^{d}\right)$ to conclude that

$$
\begin{aligned}
& \bar{e}\left(U, s, t, a, X, Y, \xi_{0} \circ X, A_{1} \circ X, A_{2} \circ(X \times X)\right) \\
& \geq o\left(W_{2}^{2}(\nu, \mu)\right)+o(|s-t|)-o\left(\|Y-X\|^{2}\right)-o\left(\left\|Y^{*}-X^{*}\right\|^{2}\right) \\
& -E\left(\left\|\pi^{1}-\pi^{2}\right\|_{L^{2}(\gamma)}, \xi_{0}-\nabla \phi, A_{1}-\nabla^{2} \phi, A_{2}-\nabla_{q_{2} q_{1}} \psi\right)\left\|\pi^{1}-\pi^{2}\right\|_{L^{2}(\gamma)}^{2} \\
& -E\left(\left\|\pi^{1}-\pi^{2}\right\|_{L^{2}(p)}, \xi_{0}-\nabla \phi, A_{1}-\nabla^{2} \phi, A_{2}-\nabla_{q_{2} q_{1}} \psi\right)\left\|\pi^{1}-\pi^{2}\right\|_{L^{2}(p)}^{2} .
\end{aligned}
$$

Since

$$
\|Y-X\|=\left\|\pi^{1}-\pi^{2}\right\|_{L^{2}(p)}^{2} \geq\left\|Y^{*}-X^{*}\right\|=\left\|\pi^{1}-\pi^{2}\right\|_{L^{2}(\gamma)}^{2}
$$

we conclude that for any $r>0$
$\liminf _{(s, Y) \rightarrow(t, X)} \frac{\bar{e}\left(U, s, t, a, X, Y, \xi_{0} \circ X, A_{1} \circ X, A_{2} \circ(X \times X)\right)}{|s-t|+\|Y-X\|^{2}} \geq-E\left(r, \xi_{0}-\nabla \phi, A_{1}-\nabla^{2} \phi, A_{2}-\nabla_{q_{2} q_{1}} \psi\right)$

$$
\begin{equation*}
-E\left(r, \xi_{0}-\nabla \phi, A_{1}-\nabla^{2} \phi, A_{2}-\nabla_{q_{2} q_{1}} \psi\right) \tag{5.10}
\end{equation*}
$$

Maximizing $-E$, which means minimizing $E$ over $(r, \phi, \psi)$ and using the fact that $\left(\xi_{0}, A_{1}, A_{2}\right) \in$ $T_{\mu}^{*, 2} \mathcal{P}_{2}\left(\mathbb{R}^{d}\right)$ we conclude

$$
\liminf _{(s, Y) \rightarrow(t, X)} \frac{\bar{e}\left(U, s, t, a, X, Y, \xi_{0} \circ X, A_{1} \circ X, A_{2} \circ(X \times X)\right)}{|s-t|+\|Y-X\|^{2}} \geq 0
$$

which proves (i).
Definition 5.6. Suppose $\kappa>0$. An upper semicontinuous function $\mathcal{U}:[0, T) \times \mathcal{P}_{2}\left(\mathbb{R}^{d}\right) \rightarrow \mathbb{R}$ is an intrinsic viscosity subsolution of (1.1) on the Wasserstein space if $\mathcal{U}(0, \cdot) \leq \mathcal{U}_{0}$ on $\mathcal{P}_{2}\left(\mathbb{R}^{d}\right)$ and

$$
a-\kappa\left(\int_{\mathbb{R}^{d}} \operatorname{Tr}\left(A_{1}(q)\right) \mu(d q)+\int_{\mathbb{R}^{2 d}} \operatorname{Tr}\left(A_{2}\left(q_{1}, q_{2}\right)\right) \mu\left(d q_{1}\right) \mu\left(d q_{2}\right)\right)+\mathcal{H}\left(\mu, \mu, \xi_{0}\right)+\mathcal{F}(\mu) \leq 0
$$

for all $(t, \mu) \in(0, T) \times \mathcal{P}_{2}\left(\mathbb{R}^{d}\right)$ and $\left(a, \xi_{0}, A_{1}, A_{2}\right) \in \mathcal{P}^{2,+} \mathcal{U}(t, \mu)$.
A lower semicontinuous function $\mathcal{U}:[0, T) \times \mathcal{P}_{2}\left(\mathbb{R}^{d}\right) \rightarrow \mathbb{R}$ is an intrinsic viscosity supersolution of (1.1) on the Wasserstein space if $\mathcal{U}(0, \cdot) \geq \mathcal{U}_{0}$ on $\mathcal{P}_{2}\left(\mathbb{R}^{d}\right)$ and

$$
\begin{equation*}
a-\kappa\left(\int_{\mathbb{R}^{d}} \operatorname{Tr}\left(A_{1}(q)\right) \mu(d q)+\int_{\mathbb{R}^{2 d}} \operatorname{Tr}\left(A_{2}\left(q_{1}, q_{2}\right)\right) \mu\left(d q_{1}\right) \mu\left(d q_{2}\right)\right)+\mathcal{H}\left(\mu, \mu, \xi_{0}\right)+\mathcal{F}(\mu) \geq 0 \tag{5.11}
\end{equation*}
$$

for all $(t, \mu) \in(0, T) \times \mathcal{P}_{2}\left(\mathbb{R}^{d}\right)$ and $\left(a, \xi_{0}, A_{1}, A_{2}\right) \in \mathcal{P}^{2,-} \mathcal{U}(t, \mu)$.
If $\mathcal{U}$ is both an intrinsic viscosity subsolution and an intrinsic viscosity supersolution of (1.1) on the Wasserstein space, we say it is an intrinsic viscosity solution of (1.1) on the Wasserstein space.

Theorem 5.7. Let $\mathcal{U}:[0, T) \times \mathcal{P}_{2}\left(\mathbb{R}^{d}\right) \rightarrow \mathbb{R}$.
(i) If $\mathcal{U}$ is an L-viscosity subsolution of (1.1) on the Wasserstein space then it is an intrinsic viscosity subsolution of (1.1).
(ii) If $\mathcal{U}$ is an L-viscosity supersolution of (1.1) on the Wasserstein space then it is an intrinsic viscosity supersolution of (1.1).
(iii) If $\mathcal{U}$ is an L-viscosity solution of (1.1) on the Wasserstein space then it is an intrinsic viscosity solution of (1.1).

Proof. It suffices to prove (ii). Assume $\mathcal{U}$ is an $L$-viscosity supersolution of (1.1) on the Wasserstein space.

Let $\mu \in \mathcal{P}_{2}\left(\mathbb{R}^{d}\right)$ and $t \in[0, T)$. Choose $X \in E$ such that $X_{\sharp} \mathcal{L}_{1}=\mu$. We have $\mathcal{U}(0, \mu)=$ $U(0, X) \geq U_{0}(X)=\mathcal{U}_{0}(\mu)$.

In order to show that $\mathcal{U}$ is upper semicontinuous at $(t, \mu)$, we choose an arbitrary sequence $\left(\mu_{n}\right)_{n} \subset \mathcal{P}_{2}\left(\mathbb{R}^{d}\right)$ converging to $\mu$ and an arbitrary sequence $\left(t_{n}\right)_{n} \subset[0, T)$ converging to $t$. Let $\left(X_{n}\right)_{n} \subset E$ such that $X_{n, \sharp} \mathcal{L}_{1}=\mu_{n}$ and $\left(X_{n}\right)_{n}$ converges to $X$. We have

$$
\lim _{n \rightarrow \infty} \mathcal{U}\left(t_{n}, \mu_{n}\right)=\lim _{n \rightarrow \infty} U\left(t_{n}, X_{n}\right) \geq U(t, X)=\mathcal{U}(t, \mu)
$$

Thus, $\mathcal{U}$ is lower semicontinuous at $(t, \mu)$.
Let now $t>0$ and $\left(a, \xi_{0}, A_{1}, A_{2}\right) \in \mathcal{P}^{2,-} \mathcal{U}(t, \mu)$. We would like to show that 5.11) holds. Let $X \in E$ be such that $X_{\sharp} \mathcal{L}_{(0,1)}^{1}=\mu$. By Lemma 5.5.

$$
\left(a, \xi_{0}(X), S_{\left(A_{1}, A_{2}\right)}\right) \in \mathcal{P}^{-} U(t, X)
$$

Since $U$ is an $L$-viscosity supersolution of (1.1) on the Wasserstein space, we use Proposition 2.3 to infer
$a-\kappa\left(\sum_{k=1}^{d} \int_{(0,1)} A_{1}(X(\omega)) e_{k} \cdot e_{k} d \omega+\int_{(0,1)^{2}} A_{2}(X(\omega), X(o)) e_{k} \cdot e_{k} d \omega d o\right)+\tilde{H}\left(X, X_{\sharp} \mathcal{L}_{1}, \xi_{0}(X)\right)+F(X) \geq 0$.
This gives 5.11.
Remark 5.8. Let $U:[0, T) \times E \rightarrow \mathbb{R}$ and let $\mathcal{U}:[0, T) \times \mathcal{P}_{2}\left(\mathbb{R}^{d}\right) \rightarrow \mathbb{R}$ be such that $\mathcal{U}(t, \mu)=$ $U(t, X)$ whenever $X \in E$ is the law of $\mu$. In [43], it was proved that if $U$ is a viscosity solution of the first order equation

$$
\left\{\begin{array}{l}
\partial_{t} U+\tilde{H}\left(X, X_{\sharp} \mathcal{L}_{1}, D U\right)+F(X)=0 \quad \text { in }(0, T) \times E  \tag{5.12}\\
U(0, X)=U_{0}(X) \quad \text { on } E,
\end{array}\right.
$$

then $\mathcal{U}$ is a intrinsic-viscosity solution of the first order equation

$$
\left\{\begin{array}{l}
\partial_{t} \mathcal{U}+\mathcal{H}\left(\mu, \mu, \nabla_{\mu} \mathcal{U}\right)+\mathcal{F}(\mu)=0 \quad \text { in }(0, T) \times \mathcal{P}_{2}\left(\mathbb{R}^{d}\right)  \tag{5.13}\\
\mathcal{U}(0, \mu)=\mathcal{U}_{0}(\mu) \text { on } \mathcal{P}_{2}\left(\mathbb{R}^{d}\right),
\end{array}\right.
$$

according to the definition proposed in [43]. Therefore, Theorem 5.7 is an extension of the results of [43] from the case $\kappa=0$ to the case $\kappa>0$.

## 6. First Order Convex HJB Equations and Value Functions

In this section we show that if $\kappa=0$ and $H$ does not depend on $\mu$ and is convex in the gradient variable then the solutions $u_{n}$ of $(1.2$, which are value functions of optimal control problems for $n$-particle systems, converge to the value function of a variational problem in $\mathcal{P}_{2}\left(\mathbb{R}^{d}\right)$. Thus we obtain a representation formula for the solution of (1.1).

Hypothesis 6.1. The function $H=H(x, p)$, in addition to satisfying Hypotheses 2.1 in the $x$ and $p$ variables, is convex in the $p$ variable and

$$
\begin{equation*}
H(x, p) \geq C_{1}+C_{2}|p|^{2} \quad \text { for all } x, p \in \mathbb{R}^{d} \tag{6.1}
\end{equation*}
$$

for some constants $C_{1}, C_{2}$, where $C_{2}>0$.
We define $L(x, v)$ to be the Legendre transform of $H(x, p)$, that is

$$
L(x, v):=\sup _{p \in \mathbb{R}^{d}}(p \cdot v-H(x, p)), \quad x, v \in \mathbb{R}^{d} .
$$

This implies, by (2.3) and 6.1),

$$
\begin{equation*}
C_{3}+C_{4}|v|^{2} \leq L(x, v) \leq C_{5}+C_{6}|v|^{2} \quad \text { for some } C_{3}, C_{4}, C_{5}, C_{6} \text { with } C_{4}, C_{6}>0 . \tag{6.2}
\end{equation*}
$$

Given $\mu \in \mathcal{P}_{2}\left(\mathbb{R}^{d}\right), \xi \in L_{\mu}^{2}\left(\mathbb{R}^{d} ; \mathbb{R}^{d}\right)$, we define

$$
\mathcal{L}(\mu, \xi):=\int_{\mathbb{R}^{d}} L(x, \xi(x)) \mu(d x)-\mathcal{F}(\mu) .
$$

For $0 \leq t \leq T$, we define the action

$$
\overline{\mathcal{A}}(\sigma, v):=\int_{0}^{t} \mathcal{L}\left(\sigma_{\tau}, v_{\tau}\right) d \tau+\mathcal{U}_{0}\left(\sigma_{0}\right)
$$

Let $u_{n}:[0, T] \times\left(\mathbb{R}^{d}\right)^{n}$ be, as before, the viscosity solution to 1.2 , for $n=1, \ldots$.
For $\mu \in \mathcal{P}_{2}\left(\mathbb{R}^{d}\right)$, let

$$
\begin{equation*}
\overline{\mathcal{U}}(t, \mu):=\inf _{(\sigma, v)}\left\{\overline{\mathcal{A}}(\sigma, v) \mid \sigma_{t}=\mu\right\}, \tag{6.3}
\end{equation*}
$$

with the infimum taken over all the pairs $(\sigma, v)$, wher ${ }^{1} \sigma=\sigma_{\tau} \in A C^{2}\left(0, t ; \mathcal{P}_{2}\left(\mathbb{R}^{d}\right)\right), v=v_{\tau}$ is a velocity vector field for $\sigma_{\tau}$ and $\sigma_{t}=\mu$. Here $A C^{2}\left(0, t ; \mathcal{P}_{2}\left(\mathbb{R}^{d}\right)\right)$ is the space of absolutely continuous curves in $\mathcal{P}_{2}\left(\mathbb{R}^{d}\right)$ with square-integrable metric derivative, see [3, Definition 1.1.1]. Define

$$
\bar{u}_{n}(t, \mathbf{x})=\overline{\mathcal{U}}\left(t, \frac{1}{n} \sum_{j=1}^{n} \delta_{x_{j}}\right), \quad \mathbf{x}=\left(x_{1}, \ldots, x_{n}\right) \in\left(\mathbb{R}^{d}\right)^{n}
$$

We want to investigate the asymptotic relationship between $\bar{u}_{n}$ and $u_{n}$.
Set

$$
f(\mathbf{x})=\mathcal{F}\left(\frac{1}{n} \sum_{j=1}^{n} \delta_{x_{j}}\right), \quad u_{0}(\mathbf{x})=\mathcal{U}_{0}\left(\frac{1}{n} \sum_{j=1}^{n} \delta_{x_{j}}\right) .
$$

and

$$
l_{n}(\mathbf{x}, \mathbf{v})=-f(\mathbf{x})+\frac{1}{n} \sum_{j=1}^{n} L\left(x_{j}, v_{j}\right)
$$

[^0]Define

$$
\mathcal{C}_{n}(t, \mathbf{x}):=\left\{x(\cdot) \in A C^{2}\left(0, t ;\left(\mathbb{R}^{d}\right)^{n}\right) \mid x(t)=\mathbf{x}\right\} .
$$

With the conditions on $H$ listed in Hypothesis [2.1, the solution $u_{n}$ to (1.2) has the value function representation

$$
\begin{equation*}
u_{n}(t, \mathbf{x})=\inf _{\xi}\left\{\int_{0}^{t} l_{n}(\xi(\tau), \dot{\xi}(\tau)) d \tau+u_{0}(\xi(0)) \mid \xi(\cdot) \in \mathcal{C}_{n}(t, \mathbf{x})\right\} \tag{6.4}
\end{equation*}
$$

Denote by $\mathcal{A}$ the functional that is minimized in (6.4), i.e.,

$$
\mathcal{A}(\xi(\cdot)):=\int_{0}^{t} l_{n}(\xi(\tau), \dot{\xi}(\tau)) d \tau+u_{0}(\xi(0))
$$

Observe that when $\sigma_{\tau}=\frac{1}{n} \sum_{j=1}^{n} \delta_{x_{j}(\tau)}$ for $x(\cdot) \in A C^{2}\left(0, t ;\left(\mathbb{R}^{d}\right)^{n}\right)$, then

$$
v_{\tau}(x)=\sum_{j=1}^{n} \mathbf{1}_{x_{j}(\tau)} \dot{x}_{j}(\tau) \quad \text { for a.e. } \tau \in(0, t) \quad \text { and } \quad \overline{\mathcal{A}}(\sigma, v)=\mathcal{A}(x(\cdot)) \text {. }
$$

We will make use of the following lemma.
Lemma 6.2. Let $\mu \in \mathcal{P}_{2}\left(\mathbb{R}^{d}\right)$, and let $\sigma \in A C^{2}\left(0, t ; \mathcal{P}_{2}\left(\mathbb{R}^{d}\right)\right)$ be a path of velocity $w$ such that $\sigma_{t}=\mu$. There exist sequences: $\left\{y^{m}\right\}_{m=1}^{\infty}, y^{m} \in\left(\mathbb{R}^{d}\right)^{m},\left\{\sigma^{m}\right\}_{m=1}^{\infty}, \sigma^{m} \in \mathcal{C}_{m}\left(t, y^{m}\right)$ with corresponding velocity vector fields $w^{m}$, and $\left\{r_{m}\right\}_{1}^{\infty}, r_{m} \searrow 0$, such that

$$
\begin{gather*}
\sup _{0 \leq \tau \leq t} d_{2}\left(\sigma_{\tau}, \sigma_{\tau}^{m}\right) \leq r_{m},  \tag{6.5}\\
\overline{\mathcal{A}}\left(\sigma^{m}, v^{m}\right) \leq \overline{\mathcal{A}}(\sigma, v)+r_{m} . \tag{6.6}
\end{gather*}
$$

Proof. We are first going to prove the existence of such sequences as in the statement, for which

$$
\begin{equation*}
\int_{0}^{t} \int_{\mathbb{R}^{d}} L\left(x, w_{\tau}^{m}(x)\right) \sigma_{\tau}^{m}(d x) d \tau \leq \int_{0}^{t} \int_{\mathbb{R}^{d}} L\left(x, w_{\tau}(x)\right) \sigma_{\tau}(d x) d \tau+r_{m} \tag{6.7}
\end{equation*}
$$

Step 1. We start with a standard mollification procedure by setting

$$
\begin{gathered}
\eta(x):=\frac{1}{(4 \pi)^{d / 2}} \exp \left(-|x|^{2} / 4\right), \quad \eta^{\varepsilon}(x):=\frac{1}{\varepsilon^{d}} \eta(x / \varepsilon), \quad \sigma_{\tau}^{\varepsilon}=\sigma_{\tau} * \eta^{\varepsilon} \\
w_{\tau}^{\varepsilon}=\frac{w_{\tau} \sigma_{\tau} * \eta^{\varepsilon}}{\sigma_{\tau}^{\varepsilon}}, \quad j^{\varepsilon}(x, y):=\frac{\eta^{\varepsilon}(x-y)}{\int_{\mathbb{R}^{d}} \eta^{\varepsilon}(x-y) \sigma_{\tau}(d y)} .
\end{gathered}
$$

By Lemma 7.1.10 of [3],

$$
\begin{equation*}
d_{2}^{2}\left(\sigma_{\tau}^{\varepsilon}, \sigma_{\tau}\right) \leq \varepsilon \int_{\mathbb{R}^{d}}|x|^{2} \eta(x) d x \quad \forall 0 \leq \tau \leq t \tag{6.8}
\end{equation*}
$$

Let us now prove that

$$
\begin{equation*}
\limsup _{\varepsilon \rightarrow 0} \int_{0}^{t} \int_{\mathbb{R}^{d}} L\left(x, w_{t}^{\varepsilon}\right) \sigma_{\tau}^{\varepsilon}(d x) d \tau \leq \int_{0}^{t} \int_{\mathbb{R}^{d}} L\left(x, w_{t}(x)\right) \sigma_{\tau}(d x) d t . \tag{6.9}
\end{equation*}
$$

Note for any arbitrary fixed $\tau$,

$$
w_{\tau}^{\varepsilon}(x)=\int_{\mathbb{R}^{d}} w_{\tau}(y) j^{\varepsilon}(x, y) \sigma_{\tau}(d y)
$$

Since for every $x \in \mathbb{R}^{d}, \int_{\mathbb{R}^{d}} j^{\varepsilon}(x, y) \sigma_{\tau}(d y)=1$ and $L(x, \cdot)$ is convex, we use Jensen's inequality to infer

$$
\begin{equation*}
L\left(\cdot, w_{\tau}^{\varepsilon}\right)=L\left(\cdot, \int_{\mathbb{R}^{d}} w_{\tau}(y) j^{\varepsilon}(\cdot, y) \sigma_{\tau}(d y)\right) \leq \int_{\mathbb{R}^{d}} L\left(\cdot, w_{\tau}(y)\right) j^{\varepsilon}(\cdot, y) \sigma_{\tau}(d y)<\infty \tag{6.10}
\end{equation*}
$$

We have obtained the finiteness of the expression at the right handside of 6.10 since

$$
\left|L\left(\cdot, w_{\tau}(y)\right)\right| \leq C\left(1+\left|w_{\tau}(y)\right|^{2}\right), \quad j^{\varepsilon} \in L^{\infty}\left(\mathbb{R}^{2 d}\right), \quad w_{\tau} \in L^{2}\left(\sigma_{\tau}\right)
$$

Observe that the function $x \mapsto \int_{\mathbb{R}^{d}} L\left(x, w_{\tau}(y)\right) j^{\varepsilon}(x, y) \sigma_{\tau}(d y)$ belongs to $L^{1}\left(\sigma_{\tau}^{\varepsilon}\right)$. Indeed,

$$
\begin{aligned}
\int_{\mathbb{R}^{d}} \int_{\mathbb{R}^{d}}\left|L\left(x, w_{\tau}(y)\right)\right| j^{\varepsilon}(x, y) \sigma_{\tau}(d y) \sigma_{\tau}^{\varepsilon}(d x) \leq & \int_{\mathbb{R}^{d}} \int_{\mathbb{R}^{d}} C\left(1+\left|w_{\tau}(y)\right|^{2}\right) j^{\varepsilon}(x, y) \sigma_{\tau}(d y) \sigma_{\tau}^{\varepsilon}(d x) \\
& =C+\int_{\mathbb{R}^{d}}\left|w_{\tau}(y)\right|^{2}\left(\int_{\mathbb{R}^{d}} \eta^{\varepsilon}(x-y) d x\right) \sigma_{\tau}(d y) \\
& \leq C\left(1+\left\|w_{\tau}\right\|_{L^{2}\left(\sigma_{\tau}\right)}^{2}\right)
\end{aligned}
$$

Similarly,

$$
\begin{aligned}
\int_{\mathbb{R}^{d}} \int_{\mathbb{R}^{d}} L\left(x, w_{\tau}(y)\right) j^{\varepsilon}(x, y) \sigma_{\tau}(d y) \sigma_{\tau}^{\varepsilon}(d x) & =\int_{\mathbb{R}^{d}} \int_{\mathbb{R}^{d}} L\left(x, w_{\tau}(y)\right) j^{\varepsilon}(x, y) \sigma_{\tau}^{\varepsilon}(d x) \sigma_{\tau}(d y) \\
& =\int_{\mathbb{R}^{d}} \int_{\mathbb{R}^{d}} L\left(x, w_{\tau}(y)\right) \eta^{\varepsilon}(x-y) d x \sigma_{\tau}(d y)
\end{aligned}
$$

Thus, integrating 6.10 on both sides with respect to $\sigma_{\tau}^{\varepsilon}$, we get

$$
\begin{equation*}
\int_{\mathbb{R}^{d}} L\left(x, w_{\tau}^{\varepsilon}(x)\right) \sigma_{\tau}^{\varepsilon}(d x) \leq \int_{\mathbb{R}^{d}} \int_{\mathbb{R}^{d}} L\left(x, w_{\tau}(y)\right) \eta^{\varepsilon}(x-y) d x \sigma_{\tau}(d y) \tag{6.11}
\end{equation*}
$$

Classic arguments show that

$$
\lim _{\varepsilon \rightarrow 0} \int_{\mathbb{R}^{d}} \int_{\mathbb{R}^{d}} L\left(x, w_{\tau}(y)\right) \eta^{\varepsilon}(x-y) d x \sigma_{\tau}(d y)=\int_{\mathbb{R}^{d}} L\left(y, w_{\tau}(y)\right) \sigma_{\tau}(d y)
$$

From this, together with 6.11, it follows that

$$
\limsup _{\varepsilon \rightarrow 0} \int_{\mathbb{R}^{d}} L\left(x, w_{\tau}^{\varepsilon}\right) \sigma_{\tau}^{\varepsilon}(d x) \leq \int_{\mathbb{R}^{d}} L\left(x, w_{\tau}(x)\right) \sigma_{\tau}(d x)
$$

An application of Fatou's lemma now yields 6.9).
Step 2. Notice that the constructed $\sigma_{\tau}^{\varepsilon}$ solve the continuity equation

$$
\partial_{\tau} \sigma_{\tau}^{\varepsilon}+\operatorname{div}\left(w_{\tau}^{\varepsilon} \sigma_{\tau}^{\varepsilon}\right)=0 \quad \text { in }(0, t) \times \mathbb{R}^{d}
$$

because

$$
\operatorname{div}\left(w_{\tau}^{\varepsilon} \sigma_{\tau}^{\varepsilon}\right)=\operatorname{div}\left(\left(w_{\tau} \sigma_{\tau}\right) * \eta^{\varepsilon}\right)=\left(\operatorname{div}\left(w_{\tau} \sigma_{\tau}\right)\right) * \eta^{\varepsilon} \quad \text { and } \quad \partial_{\tau} \sigma_{\tau}^{\varepsilon}=\left(\partial_{\tau} \sigma_{\tau}\right) * \eta^{\varepsilon}
$$

Since $w_{\tau}^{\varepsilon} \sigma_{\tau}^{\varepsilon}$ is smooth, for arbitrary $\delta>0$ we can find $\sigma_{\tau}^{\varepsilon, N} \in \mathcal{C}_{N}\left(t, y^{N}\right)$ for some $N \in \mathbb{N}$ and $y^{N} \in\left(\mathbb{R}^{d}\right)^{N}$, satisfying

$$
\begin{aligned}
\sup _{0 \leq \tau \leq t} d_{2}\left(\sigma_{\tau}^{\varepsilon}, \sigma_{\tau}^{\varepsilon, N}\right) & \leq \delta \\
\int_{0}^{t} \int_{\mathbb{R}^{d}} L\left(x, w_{\tau}^{\varepsilon, N}(x)\right) \sigma_{\tau}^{\varepsilon, N}(d x) d \tau & \leq \int_{0}^{t} \int_{\mathbb{R}^{d}} L\left(x, w_{\tau}^{\varepsilon}(x)\right) \sigma_{\tau}^{\varepsilon}(d x) d \tau+\delta
\end{aligned}
$$

It is clear that combining the latter inequalities, together with 6.9 and 6.8 , gives the desired sequence $r_{m}$, such that inequalities (6.7) and 6.5 hold.

To finish the proof, note that by (6.5) and the uniform continuity of $\mathcal{F}, \mathcal{U}_{0}$, there exist a sequence $s_{m} \searrow 0$, such that

$$
\mathcal{U}_{0}\left(\sigma_{0}^{m}\right) \leq \mathcal{U}_{0}\left(\sigma_{0}\right)+s_{m} \quad-\int_{0}^{t} \mathcal{F}\left(\sigma_{\tau}^{m}\right) d \tau \leq-\int_{0}^{t} \mathcal{F}\left(\sigma_{\tau}\right) d \tau+s_{m}
$$

Denoting $r_{m}+2 s_{m}$ still by $r_{m}$, we obtain (6.6).

Lemma 6.3. For any $0 \leq t \leq T$, the value function $\overline{\mathcal{U}}(t, \cdot)$ is lower semicontinuous on $\mathcal{P}_{2}\left(\mathbb{R}^{d}\right)$.
Proof. Let $\mu^{n} \rightarrow \mu$ be such that

$$
\lim _{n \rightarrow \infty} \overline{\mathcal{U}}\left(t, \mu^{n}\right)=\liminf _{\nu \rightarrow \mu} \overline{\mathcal{U}}(t, \nu) .
$$

Let $\varepsilon>0$ and let $\sigma^{n} \in A C^{2}\left(0, t ; \mathcal{P}_{2}\left(\mathbb{R}^{d}\right)\right)$ be paths of velocity $v^{n}$ such that $\sigma_{t}^{n}=\mu^{n}$ and

$$
\int_{0}^{t} \mathcal{L}\left(\sigma_{\tau}, v_{\tau}\right) d \tau+\mathcal{U}_{0}\left(\mu^{n}\right)<\overline{\mathcal{U}}\left(t, \mu^{n}\right)+\varepsilon
$$

It follows from (6.2) that

$$
\begin{equation*}
\int_{0}^{t}\left\|v_{\tau}^{n}\right\|_{L^{2}\left(\sigma_{\tau}^{n}\right)}^{2} d \tau<C \tag{6.12}
\end{equation*}
$$

for some $C$ independent of $n$. Therefore, by Proposition 7.1 in the Appendix, we have the existence of a subsequence (still denoted by $\sigma^{n}$ ) and $\sigma \in A C^{2}\left(0, t ; \mathcal{P}_{2}\left(\mathbb{R}^{d}\right)\right.$ ), with $\sigma_{t}=\mu$, such that for every $s \in[0, t], \sigma_{s}^{n}$ converges narrowly to $\sigma_{s}$. Denote the product measures on $\mathbb{R}^{d} \times[0, t]$ by $\sigma_{s}^{n} d s$. These converge narrowly to $\sigma_{s} d s$. Furthermore, denote by $v_{s}^{n} \sigma_{s}^{n} d s$ the vector measure whose density with respect to $\sigma_{s}^{n} d s$ is the time-dependent vector field $v_{s}^{n}=v^{n}(s, x)$.

We then obtain from (6.12) that there exists a subsequence of $\left(\sigma^{n}, v^{n}\right)$, still indexed by $n$, such that $\sigma_{s}^{n} d s$ converge narrowly to $\sigma_{s} d s$ while $v_{s}^{n} \sigma_{s}^{n} d s$ converge narrowly to a vector measure $w$ on $\mathbb{R}^{d} \times[0, t]$.

Let $\varphi \in C_{c}^{1}\left((0, t) \times \mathbb{R}^{d}\right)$. Then

$$
\begin{align*}
0 & =\lim _{n \rightarrow \infty}\left(\int_{0}^{t} \int_{\mathbb{R}^{d}} \partial_{s} \varphi(s, x) \sigma_{s}^{n}(d x) d s+\int_{0}^{t} \int_{\mathbb{R}^{d}} D \varphi(s, x) \cdot v_{s}^{n}(x) \sigma_{s}^{n}(d x) d s\right) \\
& =\int_{0}^{t} \int_{\mathbb{R}^{d}} \varphi(s, x) \sigma_{s}(d x) d s+\int_{[0, t] \times \mathbb{R}^{d}} D \varphi(s, x) \cdot w(d x, d s) . \tag{6.13}
\end{align*}
$$

By Proposition 7.2 in the Appendix, $w \ll \sigma_{s} d s$, so there is an $L^{1}\left(\sigma_{s} d s\right)$ vector field $v(s, x)$ such that $w=v_{s} \sigma_{s} d s$, and, by the same proposition,

$$
\int_{0}^{t} \int_{\mathbb{R}^{d}} L\left(x, v_{\tau}\right) \sigma_{\tau}(d x) d \tau \leq \liminf _{n \rightarrow \infty} \int_{0}^{t} \int_{\mathbb{R}^{d}} L\left(x, v_{\tau}^{n}\right) \sigma_{\tau}^{n}(d x) d \tau .
$$

However, by (6.2) and 6.12, we actually obtain

$$
\int_{0}^{t}\left\|v_{\tau}\right\|_{L^{2}\left(\sigma_{\tau}\right)}^{2} d \tau<C
$$

for some constant $C$ which, together with (6.13), means that $v$ is a velocity vector field for $\sigma$. Therefore, since $\mathcal{F}$ is narrowly continuous, it follows that

$$
\overline{\mathcal{U}}(t, \mu) \leq \liminf _{\nu \rightarrow \mu} \overline{\mathcal{U}}(t, \nu) .
$$

Theorem 6.4. Given $\mu \in \mathcal{P}_{2}\left(\mathbb{R}^{d}\right), 0 \leq t \leq T$, there exists a sequence $\{\mathbf{x}(n)\}_{n=1}^{\infty}, \mathbf{x}(n) \in\left(\mathbb{R}^{d}\right)^{n}$, such that $d_{2}\left(\frac{1}{n} \sum_{j=1}^{n} \delta_{x_{j}(n)}, \mu\right) \rightarrow 0$ as $n \rightarrow 0$ and

$$
\overline{\mathcal{U}}(t, \mu)=\lim _{n \rightarrow \infty} \inf _{\xi(\cdot) \in \mathcal{C}_{n}(t, \mathbf{x}(n))}\left\{\int_{0}^{t} l_{n}(\xi(\tau), \dot{\xi}(\tau)) d \tau+u_{0}(\xi(0)) \mid \xi(t)=\mathbf{x}(n)\right\},
$$

i.e.,

$$
\lim _{n \rightarrow \infty} u_{n}(t, \mathbf{x}(n))=\overline{\mathcal{U}}(t, \mu) .
$$

In particular, $\overline{\mathcal{U}}=\mathcal{U}$ from Theorem 1.2 and $\overline{\mathcal{U}}$ is continuous and satisfies the continuity estimate (4.1).

Proof. Let $\left\{\sigma_{k}, v_{k}\right\}_{k=1}^{\infty}$ be a minimizing sequence of paths and velocities for $\mathcal{U}(t, \mu)$ such that

$$
\begin{equation*}
\overline{\mathcal{A}}\left(\sigma_{k}, v_{k}\right) \leq \overline{\mathcal{U}}(t, \mu)+1 / k \tag{6.14}
\end{equation*}
$$

By Lemma 6.2, for each $k \in \mathbb{N}$ there exists a sequence $\left\{\sigma_{k}^{m}, v_{k}^{m}\right\}_{m=1}^{\infty}$, with the $m$-th term in $\mathcal{C}_{m}\left(t, \sigma_{k}^{m}(t)\right)$, such that

$$
\begin{equation*}
u_{m}\left(t, \sigma_{k}^{m}(t)\right) \leq \mathcal{A}\left(\sigma_{k}^{m}\right) \leq \overline{\mathcal{A}}\left(\sigma_{k}, v_{k}\right)+1 / m \quad \text { and } \quad \sigma_{k}^{m}(t) \underset{m \rightarrow \infty}{\longrightarrow} \sigma_{k}(t)=\mu \text { in } d_{2} \tag{6.15}
\end{equation*}
$$

Then,

$$
u_{m}\left(t, \sigma_{k}^{m}(t)\right) \leq \overline{\mathcal{U}}(t, \mu)+1 / k+1 / m ;
$$

consequently,

$$
\limsup _{m \rightarrow \infty} u_{m}\left(t, \sigma_{m}^{m}(t)\right) \leq \overline{\mathcal{U}}(t, \mu) .
$$

Hence, since $d_{2}\left(\sigma_{m}^{m}(t), \mu\right) \underset{m \rightarrow \infty}{\longrightarrow} 0$, this, together with the lower semicontinuity of $\overline{\mathcal{U}}(t, \cdot)$ proved in Lemma 6.3, gives

$$
\lim _{m \rightarrow \infty} u_{m}\left(t, \sigma_{m}^{m}(t)\right)=\overline{\mathcal{U}}(t, \mu)
$$

Putting $\mathbf{x}(n):=\sigma_{n}^{n}(t), n=1, \ldots$, proves the statement.

## 7. Appendix

Proposition 7.1. Let $\mu \in \mathcal{P}_{2}\left(\mathbb{R}^{d}\right)$ and let $\left(\sigma^{n}, v^{n}\right)_{n=1}^{\infty}$ be a sequence such that for each $n \in \mathbb{N}$, $\sigma^{n} \in A C^{2}\left(0, t ; \mathcal{P}_{2}\left(\mathbb{R}^{d}\right)\right)$ and $v^{n}$ is a velocity vector field for $\sigma^{n}$. If

$$
\lim _{n \rightarrow \infty} d_{2}\left(\sigma_{t}^{n}, \mu\right)=0 \quad \text { and } \quad \int_{0}^{t}\left\|v_{\tau}^{n}\right\|_{L^{2}\left(\sigma_{\tau}^{n}\right)}^{2}<C \quad \forall n \in \mathbb{N}
$$

then there exists a subsequence $\left(\sigma^{n_{k}}\right)_{k=1}^{\infty}$ and $\sigma \in A C^{2}\left(0, t ; \mathcal{P}_{2}\left(\mathbb{R}^{d}\right)\right)$, with $\sigma_{t}=\mu$, such that for a.e. $\tau \in[0, t], \sigma_{\tau}^{n_{k}} \rightarrow \sigma_{\tau}$ narrowly.

Proof. Note that if $0 \leq s_{1}<s_{2} \leq t$, by Hölder's inequality we get

$$
\begin{equation*}
d_{2}\left(\sigma_{s_{1}}^{n}, \sigma_{s_{2}}^{n}\right) \leq \int_{s_{1}}^{s_{2}}\left\|v_{\tau}^{n}\right\|_{L^{2}\left(\sigma_{\tau}^{n}\right)} d \tau \leq \sqrt{C} \sqrt{s_{2}-s_{1}} \tag{7.1}
\end{equation*}
$$

In particular, $\left(\sigma^{n}\right)_{n}$ is bounded and equicontinuous in $\mathcal{P}_{2}\left(\mathbb{R}^{d}\right)$. We apply the refined version of the Ascoli-Arzelà theorem in Proposition 3.3.1 of [3] to conclude the proof.
Proposition 7.2. Let L be as in Section [6. Consider a sequence $\left\{\nu_{n}\right\}_{1}^{\infty} \cup\{\nu\}$ of finite, positive Borel measures on $[0, T] \times \mathbb{R}^{d}$ that converges narrowly to $\nu$. Suppose we have a sequence $g_{n}$ : $[0, T] \times \mathbb{R}^{d} \rightarrow \mathbb{R}^{d}$ of vector fields such that

$$
\begin{equation*}
\int_{[0, T] \times \mathbb{R}^{d}}\left|g_{n}(t, x)\right|^{2} \nu_{n}(d t, d x)<\infty \tag{7.2}
\end{equation*}
$$

and $\left(g_{n} \nu_{n}\right)_{n}$ converges narrowly to a vector-valued Borel measure $\lambda$ on $[0, T] \times \mathbb{R}^{d}$. Then:
(i) There exists a Borel vector field $v:[0, T] \times \mathbb{R}^{d} \rightarrow \mathbb{R}^{d}$ such that $\lambda=v \nu$.
(ii) We have

$$
\int_{[0, T] \times \mathbb{R}^{d}} L(x, v(t, x)) \nu(d t, d x) \leq \liminf _{n \rightarrow \infty} \int_{[0, T] \times \mathbb{R}^{d}} L\left(x, g_{n}(t, x)\right) \nu_{n}(d t, d x) .
$$

Proof. We define on $[0, T] \times \mathbb{R}^{2 d}$ the measures $f_{n}$ by

$$
\begin{equation*}
\int_{[0, T] \times \mathbb{R}^{2 d}} \Phi(t, x, w) f_{n}(d x, d w)=\int_{[0, T] \times \mathbb{R}^{d}} \Phi\left(t, x, g_{n}(t, x)\right) \nu_{n}(t, d x), \tag{7.3}
\end{equation*}
$$

for $\Phi \in C_{b}\left([0, T] \times \mathbb{R}^{2 d}\right)$.
(i) We use $(7.2)$ and the fact that $\left(\nu_{n}\right)_{n}$ is precompact for the narrow convergence topology to conclude that $\left(f_{n}\right)_{n}$ is precompact for the narrow convergence topology. Therefore, without loss of generality, we may assume that $\left(f_{n}\right)_{n}$ converges narrowly to some Borel measure on $[0, T] \times \mathbb{R}^{2 d}$ which we denote by $f_{\infty}$. When $\Phi$ depends only on the $(t, x)$ variables, passing to the limit in (7.3), we conclude that the first marginal of $f_{\infty}$ is $\nu$. Hence, there exists a Borel map $(t, x) \rightarrow f_{\infty}^{(t, x)}$ of probability measures (cf. [3] Subsection 5.3) such that we have the disintegration

$$
\int_{[0, T] \times \mathbb{R}^{2 d}} \Phi(t, x, w) f_{\infty}(d t, d x, d w)=\int_{[0, T] \times \mathbb{R}^{d}}\left(\int_{\mathbb{R}^{d}} \Phi(t, x, w) f_{\infty}^{(t, x)}(d w)\right) \nu(d t, d x),
$$

for all $\Phi \in C_{b}\left([0, T] \times \mathbb{R}^{2 d}\right)$.
Let $\varphi \in C\left([0, T] \times \mathbb{R}^{d}, \mathbb{R}^{d}\right)$ be a bounded function. Although $(t, x, w) \in[0, T] \times \mathbb{R}^{2 d} \mapsto$ $\varphi(t, x) \cdot w$ is not bounded, (7.2) allows to assert that (setting $z=(t, x))$
$\int_{[0, T] \times \mathbb{R}^{2 d}} \varphi(z) \cdot w f_{\infty}(d z, d w)=\lim _{n \rightarrow \infty} \int_{[0, T] \times \mathbb{R}^{2 d}} \varphi(z) \cdot w f_{n}(d z, d w)=\lim _{n \rightarrow \infty} \int_{[0, T] \times \mathbb{R}^{d}} \varphi(z) \cdot g_{n}(z) \nu_{n}(d z)$.
We now use the fact that $\lambda$ is a point of accumulation of $\left(g_{n} \nu_{n}\right)$ to conclude that

$$
\int_{[0, T] \times \mathbb{R}^{d}}\left(\int_{\mathbb{R}^{d}} \varphi(t, x) \cdot w f_{\infty}^{(t, x)}(d w)\right) \nu(d t, d x)=\int_{[0, T] \times \mathbb{R}^{d}} \varphi(t, x) \cdot \lambda(d t, d x) .
$$

We conclude the proof of (i) by setting $v(t, x):=\int_{\mathbb{R}^{d}} w f_{\infty}^{(t, x)}(d w)$.
(ii) Since $L$ is bounded below by the hypotheses, we may suppose without loss of generality that $L \geq 0$. For each $r>0$ let $\Phi_{r} \in C\left(\mathbb{R}^{2 d}\right)$ be a function which is identically 1 on the ball of radius $r$, is zero outside of the ball of radius $r+1$, but remains between 0 and 1 everywhere. We have

$$
\begin{aligned}
\int_{[0, T] \times \mathbb{R}^{2 d}} L(x, w) \Phi_{r}(x, w) f_{\infty}(d t, d x, d w) & =\liminf _{n \rightarrow \infty} \int_{[0, T] \times \mathbb{R}^{2 d}} L(x, w) \Phi_{r}(x, w) f_{n}(d t, d x, d w) \\
& \leq \liminf _{n \rightarrow \infty} \int_{[0, T] \times \mathbb{R}^{2 d}} L(x, w) f_{n}(d t, d x, d w)
\end{aligned}
$$

and so, letting $r \rightarrow \infty$ we conclude

$$
\int_{[0, T] \times \mathbb{R}^{2 d}} L(x, w) f_{\infty}(d t, d x, d w) \leq \liminf _{n \rightarrow \infty} \int_{[0, T] \times \mathbb{R}^{2 d}} L(x, w) f_{n}(d t, d x, d w) .
$$

Thus,

$$
\int_{[0, T] \times \mathbb{R}^{d}}\left(\int_{\mathbb{R}^{d}} L(x, w) f_{\infty}^{(t, x)}(d w)\right) \nu(d t, d x) \leq \lim _{n \rightarrow \infty} \int_{[0, T] \times \mathbb{R}^{2 d}} L(x, w) f_{n}(d t, d x, d w) .
$$

Since $L(x, \cdot)$ is convex, we apply Jensen's inequality and use the fact that $f_{\infty}^{(t, x)}$ is a Borel probability measure to conclude the proof.

We conclude the appendix with a proof of comparison for viscosity solutions of a class of equations that includes (1.3). Let $W$ be a real separable Hilbert space. We assume the following hypothesis.

## Hypothesis 7.3.

(i) The function $\hat{H}: W \times W \rightarrow \mathbb{R}$ satisfies

$$
\begin{equation*}
|\hat{H}(X, P)-\hat{H}(X, Q)| \leq C(1+|P|+|Q|)|P-Q| \quad \text { for all } P, Q, X \in W \tag{7.4}
\end{equation*}
$$

and

$$
\begin{equation*}
|\hat{H}(X, P)-\hat{H}(Y, P)| \leq \sigma(|X-Y|(1+|P|)) \quad \text { for all } P, X, Y \in W \tag{7.5}
\end{equation*}
$$

for some modulus of continuity $\sigma$.
(ii) The function $U_{0}: W \rightarrow \mathbb{R}$ is such that

$$
\begin{equation*}
\left|U_{0}(X)-U_{0}(Y)\right| \leq m_{1}(|X-Y|) \quad \text { for all } X, Y \in W \tag{7.6}
\end{equation*}
$$

for some modulus of continuity $m_{1}$.
We note that if $\hat{H}(X, P)=\tilde{H}(X, \operatorname{law}(X), P)+F(X)$ and Hypothesis 2.1 is satisfied then Hypothesis 7.3 holds. Thus comparison for viscosity solutions of 1.3 follows from the more general theorem below.
Theorem 7.4. Let Hypothesis 7.3 be satisfied and let $\kappa \geq 0$. Let $u$ be a viscosity subsolution of

$$
\left\{\begin{array}{l}
\partial_{t} u-\kappa \sum_{k=1}^{d}\left\langle D^{2} u e_{k}, e_{k}\right\rangle+\hat{H}(X, D u)=0 \quad \text { in }(0, T) \times W  \tag{7.7}\\
u(0, X)=U_{0}(X) \quad \text { on } W
\end{array}\right.
$$

$v$ be a bounded viscosity supersolution of (7.7) and suppose that there exists $M \geq 0$ such that

$$
\sup _{(t, X) \in[0, T) \times W} u(t, X) \leq \overline{M,} \sup _{(t, X) \in[0, T) \times W}-v(t, X) \leq M .
$$

Then $u \leq v$ on $[0, T) \times W$.
Proof. The proof is similar to the proof of Theorem 3.3. The main difference is that we have to use a Hilbert space version of the maximum principle for semicontinuous functions, Theorem 3.2 of [27], instead of Theorem 8.3 of [26]. For $\delta>0$, let $\varphi_{\delta}$ be the function from Lemma 3.2 applied to the modulus $\sigma_{1}(s)=(1+T) \sigma(s)+m_{1}(s)+(2 M+1) s$. In particular we have

$$
\begin{equation*}
\varphi_{\delta}(1) \geq 2 M+1, \quad \varphi_{\delta}(s) \geq m_{1}(s) \tag{7.8}
\end{equation*}
$$

If $u \not \leq v$ then there is $\nu>0$ such that

$$
\begin{equation*}
\sup _{(t, X) \in[0, T) \times E}(u-v) \geq \nu . \tag{7.9}
\end{equation*}
$$

We will show that this leads to a contradiction.
Let $\left\{\eta_{1}, \eta_{2}, \cdots\right\}$ be an orthonormal basis of $W$. For $N=1,2, \cdots$, we denote by $P_{N}$ the orthogonal projection in $W$ onto $\operatorname{span}\left\{\eta_{1}, \cdots, \eta_{N}\right\}$, and we set $Q_{N}=I-P_{N}$. Denote $h(X):=$ $\left(1+|X|^{2}\right)^{1 / 2}$. If (7.9) is true then for sufficiently small $\mu, \gamma, \alpha>0$

$$
\sup _{X, Y \in W, t \in[0, T]}\left(u(t, X)-u_{n}(t, Y)-\frac{\mu}{T-t}-\varphi_{\delta}\left(\left(\gamma+|X-Y|^{2}\right)^{\frac{1}{2}}\right)(1+t)-\alpha(h(X)+h(Y))\right)>0 .
$$

We also notice that the expression above goes to $-\infty$ as $|X|+|Y| \rightarrow+\infty$. Therefore, by the perturbed optimization result of Ekeland-Lebourg (see for instance [29], Theorem 3.25), for every $n \geq 1$ there exist $a_{n} \in \mathbb{R}, p_{n}, q_{n} \in W$ such that $\left|a_{n}\right|+\left|p_{n}\right|+\left|q_{n}\right|<\frac{1}{n}$ and

$$
\begin{aligned}
u(t, X)-u_{n}(t, Y) & -\frac{\mu}{T-t}-\varphi_{\delta}\left(\left(\gamma+|X-Y|^{2}\right)^{\frac{1}{2}}\right)(1+t)-\alpha(h(X)+h(Y)) \\
& +a_{n} t+\left\langle p_{n}, X\right\rangle+\left\langle q_{n}, Y\right\rangle
\end{aligned}
$$

attains a strict maximum at some point $(\bar{t}, \bar{X}, \bar{Y})$. By the construction of $\varphi_{\delta}$ we have $0<\bar{t}<T$ and $|\bar{X}-\bar{Y}|<1$. It now follows from Theorem 3.2 of [27], together with Remarks 2.3 and 3.1
there, that for every $N \geq 1$ there exist $b_{1}, b_{2} \in \mathbb{R}, S_{N}, R_{N} \in S(W)$ and $C>0$ independent of $N$, such that $S_{N}=P_{N} S_{N} P_{N}, R_{N}=P_{N} R_{N} P_{N}, S_{N} \leq R_{N}$ and such that, denoting $\bar{s}=$ $\left(\gamma+|\bar{X}-\bar{Y}|^{2}\right)^{\frac{1}{2}}$,

$$
\begin{gathered}
\left(b_{1}, \varphi_{\delta}^{\prime}(\bar{s}) \frac{\bar{X}-\bar{Y}}{\bar{s}}(1+\bar{t})+\alpha D h(\bar{X})-p_{n}, S_{N}+C Q_{N}+\alpha D^{2} h(\bar{X})\right) \in \overline{\mathcal{P}}^{2,+} u(\bar{t}, \bar{X}) \\
\left(b_{2}, \varphi_{\delta}^{\prime}(\bar{s}) \frac{\bar{X}-\bar{Y}}{\bar{s}}(1+\bar{t})-\alpha D h(\bar{X})+p_{n}, R_{N}-C Q_{N}-\alpha D^{2} h(\bar{X})\right) \in \overline{\mathcal{P}}^{2,-} v(\bar{t}, \bar{Y}) \\
b_{1}-b_{2}=\varphi_{\delta}(\bar{s})+\frac{\mu}{(T-\bar{t})^{2}}-a_{n}
\end{gathered}
$$

Using the definition of viscosity subsolution we now have

$$
\begin{aligned}
b_{1} & -\kappa \sum_{k=1}^{d}\left\langle\left(S_{N}+C Q_{N}+\alpha D^{2} h(\bar{X})\right) e_{k}, e_{k}\right\rangle \\
& +\hat{H}\left(\bar{X}, \varphi_{\delta}^{\prime}(\bar{s}) \frac{\bar{X}-\bar{Y}}{\bar{s}}(1+\bar{t})+\alpha D h(\bar{X})-p_{n}\right) \leq 0
\end{aligned}
$$

which implies, by (7.4),

$$
\begin{align*}
b_{1}-\kappa \sum_{k=1}^{d}\left\langle S_{N} e_{k}, e_{k}\right\rangle & +\hat{H}\left(\bar{X}, \varphi_{\delta}^{\prime}(\bar{s}) \frac{\bar{X}-\bar{Y}}{\bar{s}}(1+\bar{t})\right)  \tag{7.10}\\
& \leq \sigma_{2}\left(\frac{1}{N}\right)+\sigma_{3}\left(\frac{1}{n}\right)+\sigma_{4}(\alpha)
\end{align*}
$$

for some moduli $\sigma_{2}, \sigma_{3}, \sigma_{4}$. Similarly we have

$$
\begin{align*}
b_{2}-\kappa \sum_{k=1}^{d}\left\langle R_{N} e_{k}, e_{k}\right\rangle & +\hat{H}\left(\bar{Y}, \varphi_{\delta}^{\prime}(\bar{s}) \frac{\bar{X}-\bar{Y}}{\bar{s}}(1+\bar{t})\right)  \tag{7.11}\\
& \geq \sigma_{2}\left(\frac{1}{N}\right)+\sigma_{3}\left(\frac{1}{n}\right)+\sigma_{4}(\alpha)
\end{align*}
$$

Subtracting (7.11) from (7.10) and using $S_{N} \leq R_{N}$, 7.5), we obtain

$$
\begin{align*}
& \sigma_{2}\left(\frac{1}{N}\right)+\sigma_{3}\left(\frac{1}{n}\right)+\sigma_{4}(\alpha) \geq \varphi_{\delta}(\bar{s})+\frac{\mu}{(T-\bar{t})^{2}} \\
& +\hat{H}\left(\bar{X}, \varphi_{\delta}^{\prime}(\bar{s}) \frac{\bar{X}-\bar{Y}}{\bar{s}}(1+\bar{t})\right)-\hat{H}\left(\bar{Y}, \varphi_{\delta}^{\prime}(\bar{s}) \frac{\bar{X}-\bar{Y}}{\bar{s}}(1+\bar{t})\right)  \tag{7.12}\\
& \geq \varphi_{\delta}(\bar{s})+\frac{\mu}{(T-\bar{t})^{2}}-\sigma\left(\bar{s}\left(1+\varphi_{\delta}^{\prime}(\bar{s})(1+T)\right)\right) \\
& \geq \frac{\mu}{(T-\bar{t})^{2}}+\varphi_{\delta}(\bar{s})-\sigma_{1}\left(\varphi_{\delta}^{\prime}(\bar{s}) \bar{s}+\bar{s}\right) \geq \frac{\mu}{T^{2}}
\end{align*}
$$

where we have used the definition of $\sigma_{1}$ and Lemma 3.2 to justify the last two inequalities. Inequality (7.12) yields a contradiction after we send $N \rightarrow+\infty$, then $n \rightarrow+\infty$ and finally $\alpha \rightarrow 0$.

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[^0]:    ${ }^{1}$ We use the subindex notation $\sigma_{\tau}$ or $\sigma(\tau)$ interchangeably to mean the value of the path at time $\tau$.

