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A MODIFIED LEAST ACTION PRINCIPLE ALLOWING MASS CONCENTRATIONS FOR THE EARLY UNIVERSE RECONSTRUCTION PROBLEM

YANN BRENIER

ABSTRACT

We address the early universe reconstruction (EUR) problem (as considered by Frisch and coauthors in [24]), and the related Zeldovich approximate model [39]. By substituting the fully nonlinear Monge-Ampère equation for the linear Poisson equation to model gravitation, we introduce a modified mathematical model ("Monge-Ampère gravitation/MAG"), for which the Zeldovich approximation becomes exact. The MAG model enjoys a least action principle in which we can input mass concentration effects in a canonical way, based on the theory of gradient flows with convex potentials and somewhat related to the concept of self-dual Lagrangians developed by Ghoussoub [27]. A fully discrete algorithm is introduced for the EUR problem in one space dimension.

INTRODUCTION

This paper addresses the early universe reconstruction (EUR) problem discussed by Frisch and coauthors in [24, 17]. We refer to these papers for the detailed physical background of this important problem in cosmology. Here is a simplified mathematical formulation.

We consider (for simplicity) a smooth closed bounded 3D domain $D \subset \mathbb{R}^3$ and denote the space variable by $x \in D$. We are given two times $t_1 > t_0 > 0$, two probability measures $\rho_0(dx)$, $\rho_1(dx)$ on $D$. We look for a time-dependent family of probability measures $\rho(t, dx)$ on $D$ (depending continuously on the time variable $t$, with respect to the weak convergence of measures), interpolating $\rho_0$ and $\rho_1$, at $t = t_0$ and $t = t_1$ respectively, and minimizing the following action

$$\int_{t_0}^{t_1} dt \int_D \alpha(t)^2 \rho(t, dx) |v(t, x)|^2 + \beta(t)^2 |\nabla \varphi(t, x)|^2 \, dx,$$

where $v = v(t, x) \in \mathbb{R}^3$ is a vector-field and $\varphi = \varphi(t, x)$ a scalar field, respectively subject to:

$$\partial_t \rho + \nabla \cdot (\rho v) = 0, \quad \rho = 1 + \Delta \varphi.$$

Here $\Delta = \nabla^2$ is the Laplace operator. In the definition of the action, $\alpha$ and $\beta$ are time-dependent scaling parameters related to general relativity (GR). Following [24, 17] (case of an Einstein-de Sitter universe), we set:

$$\alpha(t) = t^{3/4}, \quad \beta(t) = t^{-1/4} \sqrt{3/2}.$$
The formal optimality conditions read

\[(0.4)\quad v = v(t, x) = \frac{\nabla \theta(t, x)}{\alpha(t)^2}, \quad \partial_t \theta + \frac{|\nabla \theta|^2}{2\alpha^2} + \beta^2 \varphi = 0,\]

which can be also (still formally) written:

\[(0.5)\quad \partial_t (\alpha^2 v) + \alpha^2 (v \cdot \nabla) v = -\beta^2 \nabla \varphi, \quad \nabla \times v = 0,\]

or

\[(0.6)\quad \partial_t (\alpha^2 \rho v) + \nabla \cdot (\alpha^2 \rho v \otimes v) = -\beta^2 \rho \nabla \varphi, \quad \nabla \times v = 0.\]

These equations, namely (0.2,0.6) are (up to the scaling factors \(\alpha, \beta\) which come from general relativity) nothing but the Euler-Poisson equations for a pressure-less, curl-free, self-gravitating gas subject to classical Newton gravitation. These equations can also be written, using “material coordinates”,

\[(0.7)\quad \partial_t (\alpha(t)^2 \partial_t X(t, a)) = -\beta(t)^2 (\nabla \varphi)(t, X(t, a)),\]

where \(a\) denotes the material coordinate, \(X(t, a)\) the position at time \(t\) of the mass particle with label \(a\). In the case of coefficients (0.3), we find

\[(0.8)\quad \frac{2t}{3} \partial_{tt} X(t, a) + \partial_t X(t, a) = -(\nabla \tilde{\varphi})(t, X(t, a)),\]

where \(\tilde{\varphi}(t, x) = \varphi(t, x)/t\) satisfies

\[\rho = 1 + t \Delta \tilde{\varphi}.\]

Remarkably enough, at early stage \(t \ll 1\), the density field must be uniformly equal to 1 (otherwise solutions are unbounded) and, even more surprisingly, the acceleration term is dominated by the velocity term, due to general relativity! (In some sense, Newton modified by Einstein returns to Aristoteles.) As a consequence, an amazingly simple approximate formula was proposed by Zeldovich [39]:

\[(0.9)\quad X(t, a) = a - t \nabla \tilde{\varphi}_0(a),\]

where \(\tilde{\varphi}_0\) is related to the behavior of the density field \(\rho(t, \cdot)\) at early stages

\[\rho(0, x) = 1, \quad \Delta \tilde{\varphi}_0 = \lim_{t \downarrow 0} \frac{\rho(t, x) - 1}{t}.\]

The Zeldovich approximate formula suggests possible mass concentrations in finite time. Indeed, denoting by \(\Lambda\) the largest eigenvalue of the Hessian matrix \(D^2 \tilde{\varphi}_0(a)\), for all \(a\), we see that, whenever \(\Lambda > 0\), the map \(a \rightarrow X(t, a)\) is no longer invertible at \(t = \Lambda^{-1}\). Beyond the concentration time, there are many possibilities of extending the formula and this is still a controversial issue from the physical viewpoint. It depends very much on whether or not we want to prevent interpenetration of particles. If we do so, we are naturally lead to the model of adhesion dynamics, where particles merge after collisions, which is the most possible dissipative behavior beyond concentrations. (See [5, 23, 24].) This issue can be simply addressed in terms of nonlinear hyperbolic PDEs [22]. Indeed, given a Zeldovich solution \(X\) defined by (0.9), let us introduce the field \(u(t, x)\) implicitly defined by:

\[(0.10)\quad u(t, X(t, a)) = a - X(t, a) = t \nabla \tilde{\varphi}_0(a),\]
as long as $a \to X(t,a)$ stays smooth and invertible. Then, we see that $u$ solves the multidimensional "inviscid Burgers" equation

\begin{equation}
\partial_t u + (u \cdot \nabla) u = 0.
\end{equation}

In one space dimension, if we want a global solution for all times, the monotonicity condition $\partial_t X(t,a) \geq 0$ exactly corresponds to “Oleinik’s entropy condition" $\partial_x u \leq 1/t$, which guarantees both global existence and uniqueness for solutions of the inviscid Burgers equation (0.11), written in “conservation form”

\begin{equation}
\partial_t u + \partial_x u = 0.
\end{equation}

Going back to the least action principle, it is remarkable that action (0.1) is (strictly) convex in the variables $(\rho, \rho v, \varphi)$, while the differential constraints (0.2) are linear in these variables. With classical tools of convex and functional analysis, Loeper [30] was able to prove the existence of a unique minimizer in the class of probability measures $\rho(t,dx)$ which have no singular part with respect to the Lebesgue measure, i.e. without concentrations, provided the data $\rho_0$ and $\rho_1$ are themselves concentration-free. [He also rigorously derived optimality conditions (0.6).] This existence and uniqueness result is quite remarkable. However, it provides only solutions to the Euler-Poisson system without concentrations. This is a major defect, since the initial value problem (as both $\rho$ and $v$ are prescribed at initial time $t_0$) is expected to have solutions with concentrations appearing in finite time, for most initial conditions. So, minimizing action (0.1) under constraints (0.2) is not a fully satisfactory way to recover solutions of the EP system from the knowledge of $\rho_0$ and $\rho_1$. The goal of the present paper is to investigate how the action can be modified so that its minimizers are not necessarily concentration-free. A similar problem, in the framework of adhesion-fragmentation processes, has been recently solved by Wolansky [38]. (See also the pioneering work of Shnirelman [35] for sticky particles and adhesion dynamics.) Our approach is different and more reminiscent of the recent theory of self-dual lagrangians by Ghoussoub [27]. Unfortunately, it does not apply to the desired Euler-Poisson system but rather to the modified system obtained by substituting the fully nonlinear Monge-Ampère equation $\rho = \det(I + tD^2_x \tilde{\varphi})$, for the Poisson equation $\rho = 1 + t\Delta \tilde{\varphi}$ (in case of coefficients (0.3)). Notice that these systems agree for solutions depending only on one spatial coordinate (i.e. with sheet structure) and are asymptotically close for $t << 1$. Of course, changing the model is not a satisfactory approach, without further justification. Our main argument is the following remarkable property of the resulting Euler-Monge-Ampère system: it admits as exact solutions the Zeldovich approximate solutions (0.9). As a secondary justification, let us recall that the Euler-Poisson system is, after all, itself an approximation of the full Einstein equations and it might be, from this viewpoint, equally good to use the Monge-Ampère equation and the Poisson equation. [A similar situation occurs in fluid mechanics when comparing the quasi-geostrophic and the semi-geostrophic approximations of the Euler equations for ocean and atmosphere dynamics, as discussed in [21]. See also [18].] However, there will be no attempt in the present paper to justify this last statement.

The structure of the paper is as follows:
In the first section, we introduce a family of infinite dimensional dynamical systems, in a rather abstract (but straightforward) framework, requiring a limited amount of classical functional analysis (Hilbert spaces) and elementary convex analysis (sub-differential calculus). Then, the ”Monge-Ampère” gravitation (MAG) model is introduced as a particular case. The model is first described in ”material variables” and then translated into ”Eulerian variables”, leading to a pressure-less Euler-Monge-Ampère system of PDEs. 

In section 2, we return to the abstract framework and observe that the potential part of the action has the very special property to be a squared distance function. This allows a rewriting of the action as an exact square and we find as special minimizers all the solutions of the gradient flow equation associated to the potential. [This special solutions play more or less the same role as ”instantons” in Yang-Mills theory [27].] It turns out that this gradient flow belongs to a very well-studied class of evolution equations, linked to the theory of ”maximal monotone operators” [19]. This suggests a somewhat canonical modification of the action.

In section 3, we adapt the previous discussion to Monge-Ampère gravitation and get a modified action that takes into account concentrations. Then, we observe that the solutions of the gradient flow equation exactly are the Zeldovich approximate solutions (0.9) of the Euler-Poisson system, which, in some sense, justifies the MAG model as a reasonable approximate model.

In section 4, we introduce a fully discrete algorithm to compute the EUR problem with Monge-Ampère gravitation.

In section 5, we introduce a numerical scheme for the initial value problem and, finally, in section 6, we provide numerical results in the very special case of one space variable.

1. The abstract framework

Let $H$ be a (separable) Hilbert space $H$ equipped with its norm denoted $\| \cdot \|$ and the corresponding inner product $(\cdot, \cdot)$. We first consider the general dynamical system

$$\frac{d^2 X}{dt^2} = (\nabla_H \Phi)[X],$$

where $t \to X(t)$ is valued in $H$, $\nabla_H$ denotes the gradient operator in $H$, and $\Phi$ is a given ”potential” defined on $H$. (Observe that we do not follow the usual sign convention for the potential, for notational convenience.) As well known, such a system admits a least action principle, at least at a formal level. Indeed, for a curve $t \to X(t)$ valued in the Hilbert space $H$, we may define its action between times $t_0$ and $t_1$, $t_1 > t_0$ by:

$$A_{[t_0, t_1]}[X] = \int_{t_0}^{t_1} \frac{1}{2} || \frac{dX}{dt} ||^2 + \Phi[X(t)] \ dt.$$

Then, the dynamical equation (1.13) can be seen as the formal optimality equation obtained by minimizing the action (1.14) as the end points $X(t_0)$ and $X(t_1)$ are fixed. 

Next, we crucially assume the potential to be of form:

$$\Phi[X] = \inf \left\{ \frac{||X - s||^2}{2} ; s \in S \right\},$$

where $S$ is a given bounded subset of $H$. Then, when it makes sense, $X - \nabla_H \Phi[X]$ is just the closest point $\pi[X]$ to $X$ in the set $S$. (Clearly this definition is ambiguous whenever
$X$ has several distinct closest points, which may happen unless $S$ is a convex set. In some cases, $X$ may have no closest point in $S$! As a consequence, (1.13) formally means:

\[(1.16) \quad \frac{d^2X}{dt^2} = X - \pi[X],\]

where $\pi[X]$ is the closest point to $X$ on $S$. With this formulation, we can guess a large class of explicit solutions. Indeed, let us assume that $X(0) = X_0$ has a unique closest point $\pi[X_0] = \pi_0$ on $S$. Then the linear (but not convex) combination of $X_0$ and $\pi_0$ given by:

\[(1.17) \quad X(t) = \pi_0 + e^t(X_0 - \pi_0)\]

solves (1.16) as long as $\pi_0$ stays the unique projection of $X(t)$. Intuitively, $X(t)$ gets repelled from its initial position in the opposite direction of its closest point on $S$, keeping for a while $\pi_0$ as its closest point on $S$ until a new point in $S$ gets even closer. Whenever $S$ is a convex set, this repulsion mechanism provides an obvious global solution. Indeed, all points contained in the infinite segment $\{\pi_0 + r(X_0 - \pi_0), \quad r \geq 0\}$ admits $\pi_0$ as their unique closest point on $S$. In the case of a non-convex set $S$, this is not true in general and formula (1.17) is able to provide no more than a local solution. The situation is very clear in the elementary case when $S$ is the unit sphere in $H$. Then, 0 is the unique point where $\Phi$ is not differentiable. We get as special solution

\[X(t) = r_0^{-1}(1 + (r_0 - 1)e^t)X_0,\]

where $X_0 \neq 0$ and $r_0 = ||X_0||$. We see that, if $r_0 < 1$, then the solution reaches 0 at time $T = -\log(1 - r_0)$ and its continuation beyond $T$ gets ambiguous.

**Miscellaneous mathematical remarks.** 1) The potential $\Phi$ given by (1.15) is a smooth perturbation of a Lipschitz concave function; indeed:

\[(1.18) \quad \Phi[X] = \frac{||X||^2}{2} - \Pi[X],\]

where $\Pi$ is the Lipschitz convex functional defined by:

\[(1.19) \quad \Pi[X] = \sup\{(x, s) - \frac{||s||^2}{2}; \quad s \in S\}.\]

A classical result of convex analysis [4] asserts that, for every Lipschitz convex function defined on a Hilbert space, the set $\tilde{H}$ where the function is differentiable is always "fat" in the topological sense of Baire: namely $\tilde{H}$ is dense and contains a countable intersection of dense open subsets of $H$ [4]. In the particular case of $\Pi$, the set of differentiability $\tilde{H}$ is contained in the set of all points $X$ in $H$ for which there is a unique closest point $s = \pi[X]$ on $S$. Thus, the potential $\Phi$ defined by (1.15) is everywhere differentiable on $\tilde{H}$ and its gradient in $H$ is given by:

\[(1.20) \quad \nabla_{\tilde{H}}\Phi[X] = X - \pi[X], \quad \forall X \in \tilde{H}.\]

2) In the case when $H$ is the finite-dimensional Hilbert space $\mathbb{R}^n$, for such a potential (namely a smooth perturbation of a concave Lipschitz function), the dynamical system (1.13) has a unique global solution for Lebesgue almost every initial condition $(X(0), X'(0)) \in \mathbb{R}^{2n}$ and is, therefore, well-posed in the sense of Bouchut and Ambrosio.
To the best of our knowledge, there is no similar theory in infinite dimension and the well-posedness of (1.13) is then a challenging open. (A somewhat related attempt is the theory developed by Ambrosio and Gangbo for some infinite dimensional hamiltonian systems [2]. See also [18, 25, 26].)

1.1. Monge-Ampère gravitation.

Definition. Since the dimension 3 does not matter in the definition of the MAG model, we consider a smooth bounded closed domain $D \subset \mathbb{R}^d$. We assume $D$ to be of unit Lebesgue measure. The MAG model is defined by choosing for $H$ the Hilbert space of all Lebesgue square-integrable maps from $D$ to $\mathbb{R}^d$,

$$H = L^2(D, R^d),$$

and for $S$ the subset of all Lebesgue measure-preserving maps $s$ of $D$:

$$S = \{ s \in H, \int_D f(s(a)) da = \int_D f(a) da, \forall f \in C(R^d) \}.$$

In addition, with respect to the abstract framework, we input coefficients $\alpha, \beta$ given by (0.3) and substitute

$$\beta^{-2}(t) \frac{d}{dt}(\alpha^2(t) \frac{dX}{dt}) = (\nabla H \Phi)[X] = X - \pi[X],$$

for (1.13).

Eulerian formulation. We now want to write the MAG model in "Eulerian coordinates". Given a solution $t \rightarrow X(t) \in H$, we introduce the pair of measures $(\rho, \rho v)$, respectively valued in $R$ and $R^d$ and defined by:

$$\rho(t, dx) = \int_D \delta(x - X(t, a)) da,$$

$$v(t, x) \rho(t, dx) = \int_D \partial_t X(t, a) \delta(x - X(t, a)) da.$$

Using the tools of "optimal transportation theory" [11, 12, 37] (cf. Appendix), we (formally) get for $(\rho, v)$ the following system of partial differential equations:

$$\partial_t \rho + \nabla \cdot (\rho v) = 0,$$

$$\partial_t (\alpha^2 \rho v) + \nabla \cdot (\alpha^2 \rho v \otimes v) = -\beta^2 \rho \nabla \varphi,$$

$$\rho = \det(I + D^2_x \varphi),$$

where $(\rho, v, \varphi)(t, x) \in R_+ \times R^d \times R$, $(t, x) \in R \times R^d$,

$$I = (\delta_{ij}), \quad D^2_x = (\partial^2_{x_i x_j}).$$

Thus our model can be seen as a nonlinear correction of the classical pressure-less Euler-Poisson system (EP), obtained by substituting the fully nonlinear Monge-Ampère equation

$$\rho = \det(I + D^2_x \varphi).$$

for the linear Poisson equation

$$\rho = 1 + \Delta \varphi.$$
For this reason, we call Monge-Ampère gravitation (MAG) the model defined by either (1.16,1.21,1.22) or (1.25). Notice that the MAG model coincides with the conventional Euler-Poisson model in dimension \( d = 1 \).

2. Modified action in the abstract framework

In this section, we go back to the abstract framework of a potential \( \Phi \) defined as the squared distance to a bounded subset \( S \) of a general Hilbert space \( H \), according to (1.16,1.18,1.19).

2.1. The self-dual form of the action. Since potential \( \Phi \) is a squared distance to some subset \( S \) inside \( H \), it solves, at least formally, the stationary Hamilton-Jacobi equation:

\[
\Phi = \frac{||\nabla H \Phi||^2}{2},
\]

where \( \nabla H \) denotes the gradient operator in \( H \). This suggest to rewrite, at least formally, the action (1.14) as

\[
\int_{t_0}^{t_1} \frac{1}{2} \left( \frac{dX}{dt} \right)^2 + \frac{||\nabla H \Phi[X(t)]||^2}{2} \, dt
\]

\[
= \int_{t_0}^{t_1} \frac{1}{2} \left( \frac{dX}{dt} \right)^2 - \nabla H \Phi[X(t)] + \left( \frac{dX}{dt}, \nabla H \Phi[X(t)] \right) \, dt
\]

\[
= \Phi[X(t_1)] - \Phi[X(t_0)] + \int_{t_0}^{t_1} \frac{1}{2} \left( \frac{dX}{dt} \right)^2 - \nabla H \Phi[X(t)] \, dt.
\]

Under this "self-dual" form (see [27] for a systematic study of "self-dual lagrangians"), it is obvious that any solution of

\[
\frac{dX}{dt} = (\nabla H \Phi)[X] = X - (\nabla H \Pi)[X]
\]

is always a minimizer of the action as \( X(t_0) \) and \( X(t_1) \) are fixed (just like instantons in euclidean Yang-Mills theory, cf.[27]).

2.2. A gradient flow equation. As already mentioned, in spite of the rather nice structure (1.18) of \( \Phi \), as a quadratic perturbation of a convex Lipschitz function, the corresponding second-order equation (1.16) is not so well understood. In sharp contrast, the first-order equation (2.30) is a standard "gradient flow" equation (GF), that can be solved by classical "maximal monotone operator" theory [19].

In the framework of maximal monotone operator theory, equation (2.30) is usually written as a sub-differential inclusion:

\[
-\frac{dX}{dt} + X \in \partial \Pi[X],
\]

which is well-posed in \( H \) since \( \Pi \) is Lipschitz and convex. Here, we use standard notations of convex analysis, for which \( \partial \) denotes the sub-differential of a convex function [19]:

\[
\partial \Pi[X] = \{ Z \in H; \quad \Pi[Y] \geq \Pi[X] + ((Z,Y - X)), \quad \forall Y \in H \}.
\]
A remarkable property \cite{19} of each solution $X(t) \in H$ is to be not only a Lipschitz continuous function of $t$ but also right-differentiable at each $t$ with

\begin{equation}
\frac{dX(t + 0)}{dt} + X(t) = d^0\Pi[X(t)], \quad \forall t
\end{equation}

where $d^0\Pi[X]$, following \cite{3}, denotes the element of $\partial\Pi[X]$ with minimal norm (which is uniquely defined):

\begin{equation}
||d^0\Pi[X]|| = \min\{||s|| \mid s \in \partial\Pi[X]\}.
\end{equation}

Finally, notice that $X(t)$ is a locally Lipschitz function of $t$ with values in the separable Hilbert space $H$, $X(t)$ is therefore almost everywhere differentiable in $t$ by Rademacher theorem. Since $X(t)$ is right-differentiable everywhere, we conclude that:

\begin{equation}
\frac{dX}{dt} + X(t) = d^0\Pi[X(t)],
\end{equation}

holds true both in the almost everywhere sense and in the sense of distributions.

2.3. **Proposals for a modified action.** Our main point is now to introduce a modified action. There are two possible ways to do it. First, we may introduce the modified potential $\tilde{\Phi}$:

\begin{equation}
\tilde{\Phi}[X] = \frac{1}{2}||X - d^0\Pi[X]||^2
\end{equation}

and the corresponding modified action

\begin{equation}
\tilde{A}_{[t_0,t_1]}[X] = \int_{t_0}^{t_1} \frac{1}{2}||\frac{dX}{dt}||^2 + \tilde{\Phi}[X(t)] \; dt = \int_{t_0}^{t_1} \frac{1}{2}||\frac{dX}{dt}||^2 + \frac{1}{2}||X - d^0\Pi[X]||^2 \; dt.
\end{equation}

Alternately, sticking more closely to the self-dual formulation, we may directly modify the action by setting

\begin{equation}
\hat{A}_{[t_0,t_1]}[X] = \int_{t_0}^{t_1} \frac{1}{2}||\frac{dX}{dt} - X + d^0\Pi[X]||^2 \; dt.
\end{equation}

It is not clear to us that these modified actions coincide (up to boundary terms). Nevertheless, we will take the second option, mostly for numerical purposes, because it leads to simpler algorithms.

3. **Application to the EUR problem with Monge-Ampère gravitation**

We now address the Monge-Ampère gravitation (MAG) model and the EUR problem. This means, with respect to the abstract framework, that $H$ and $S$ are now defined by (1.21,1.22) and (1.23) substitutes for (1.13).
3.1. A modified action for the EUR problem. In order to take into account coefficients \((\alpha, \beta)\) (given by (0.3)), we first rewrite the action as:

\[
A = \int_{t_0}^{t_1} \alpha(t)^2 \left\| \frac{dX}{dt} \right\|^2 + \beta(t)^2 \left\| \nabla_H \Phi[X(t)] \right\|^2 \; dt.
\]  

As in the homogeneous case \(\alpha = \beta = 1\), we keep in mind that

\[
\frac{1}{2} \left\| \nabla_H \Phi[X(t)] \right\|^2 = \Phi[X(t)]
\]

and look at the cross-term:

\[
J = \int_{t_0}^{t_1} \alpha(t)\beta(t) \left( \frac{dX}{dt}, \nabla_H \Phi[X(t)] \right) dt = \int_{t_0}^{t_1} \alpha(t)\beta(t) \frac{d}{dt} \left( \Phi[X(t)] \right) dt
\]

By integration by part, we get

\[
J - \alpha(t_1)\beta(t_1) \Phi[X(t_1)] + \alpha(t_0)\beta(t_0) \Phi[X(t_0)] = - \int_{t_0}^{t_1} \Phi[X(t)] \frac{d}{dt} \left( \alpha(t)\beta(t) \right) dt
\]

\[
= - \frac{1}{2} \int_{t_0}^{t_1} \left\| \nabla_H \Phi[X(t)] \right\|^2 \frac{d}{dt} \left( \alpha(t)\beta(t) \right) dt = - \lambda \frac{1}{2} \int_{t_0}^{t_1} \beta(t)^2 \left\| \nabla_H \Phi[X(t)] \right\|^2 dt,
\]

provided we assume

\[
\frac{d}{dt} \left( \alpha(t)\beta(t) \right) = \lambda\beta(t)^2,
\]

for some constant \(\lambda\), which is is consistent with data (0.3) if we choose \(\lambda = 1/\sqrt{6}\). From this calculation of the cross-term \(J\), we deduce that the action \(A\) defined by (3.39) can be written:

\[
A = BT + \int_{t_0}^{t_1} \left\| \alpha(t) \frac{dX}{dt} - \mu\beta(t) \nabla_H \Phi[X(t)] \right\|^2 \; dt.
\]

where \(BT\) is a boundary term depending only on \(X(t_1)\) and \(X(t_0)\), provided \(\mu^2 + \mu \lambda = 1\). For data (0.3), we get \(\lambda = 1/\sqrt{6}\) and \(\mu = \sqrt{2/3}\). Therefore, all solutions of the gradient-flow equation

\[
\alpha(t) \frac{dX}{dt} = \mu\beta(t) \nabla_H \Phi[X(t)]
\]

automatically are minimizers of the action (3.39). For data (0.3), this gradient-flow equation reduces to:

\[
t \frac{dX}{dt} = \nabla_H \Phi[X(t)] = X(t) - \nabla_H \Pi[X(t)].
\]

The gradient-flow equation should be understood in the more precise sense:

\[
t \frac{dX(t + 0)}{dt} = X(t) - d^0 \Pi[X(t)],
\]

which takes concentration into account, globally in time. In some sense, formulation (3.43) not only allows concentrations but guarantees the largest possible dissipation of kinetic energy during the concentration process (which is of course questionable from the
physical viewpoint.) Accordingly, we suggest, for the EUR problem with Monge-Ampère gravitation, the following modified action:

\[ \dot{A} = \int_{t_0}^{t_1} t^{-1/2} \left| \frac{dX}{dt} - X + d\Pi[X] \right|^2 \, dt. \]

3.2. Zeldovich solutions. Special solutions of (3.42) can be obtained thanks to the concept of "rearrangements with convex potential" as follows. By definition, the MAG model relies on the set \( S \) of all Lebesgue measure-preserving maps \( 1.22 \). This set contains the identity map \( Id \) as an obvious element. The set \( K \subset H \) of all points \( X \) which admits \( Id \) as a closest point on \( S \) plays a crucial role. It can be characterized (cf. Appendix on optimal transportation theory), as the convex cone of all maps \( X \in H \) with a convex potential, which means that there is a convex function \( \psi \) defined on \( \mathbb{R}^d \) and valued in \( [\infty, +\infty] \) which is almost everywhere differentiable on \( D \) with \( \nabla \psi(a) = X(a) \), a.e. on \( D \). It turns out that any map \( X \in H \) has a unique rearrangement \( X^2 \) in \( K \), which means

\[ \int_D \delta(x - X^2(a)) \, da = \int_D \delta(x - X(a)) \, da \]

(cf. Appendix).

Therefore, special solutions of (3.42) can be obtained, by looking for solutions \( X(t) \) valued in the convex cone \( K \) of all maps with convex potential. Indeed, for such solutions, we have:

\[ \nabla_H \Phi[X(t)] = X(t) - \nabla_H \Pi[X(t)] = X(t) - \pi[X(t)] = X(t) - Id, \]

and (3.42) reduces to the linear ODE

\[ t \frac{dX}{dt} = X - Id \]

as long as \( X(t) \) belongs to \( K \), i.e. \( X(t,a) = \nabla \psi(t,a) \), with \( \psi(t,a) \) convex in \( a \). This leads to the explicit formula:

\[ X(t,a) = \nabla \psi(t,a) = a + \frac{t}{t_0} (\nabla \psi(t_0,a) - a) = a + \frac{t}{t_0} (X(t_0,a) - a), \]

as long as \( \psi \) stays convex in \( a \). This exactly coincides with Zeldovich formula (0.9) discussed in the introduction. Remarkably enough, for Monge-Ampère gravitation, Zeldovich approximation (0.9) is just exact!

3.3. Modified action in one space dimension. Let us focus on the one space dimension case when: \( D = [-1/2, 1/2] \). Then, the modified potential \( \Phi \) can be explicitly computed in the case of a piecewise smooth map \( Y \) valued in \( K \). Indeed, in one space dimension, maps in \( K \), with convex potential are just increasing maps. So, there is a finite number of plateaux \([a_j, b_j] \) on which \( Y \) is constant with values \( Y_j \) and outside of which \( Y \) is a piecewise smooth strictly increasing function.

Notice that the corresponding image-measure \( \rho(dx) \) defined by

\[ \rho(dx) = \int_D \delta(x - Y(a)) \, da \]
has a singular part $\rho_s$ given by:

$$\rho_s(dx) = \sum_j (b_j - a_j) \delta(x - Y_j).$$

Then $d^0\Pi[Y]$ (the element of the sub-differential $\partial\Pi[Y]$ with minimal $L^2$ norm) coincides with the identity map outside of the plateaux and takes value $(a_j + b_j)/2$ inside $[a_j, b_j]$. After elementary calculations, we find

$$||Y - d^0\Pi[Y]||^2 = ||Y||^2 - 2((Y, Id)) + ||Id||^2 - \sum_j \int_{a_j}^{b_j} (a - \frac{a_j + b_j}{2})^2 da$$

$$= ||Y - Id||^2 - \frac{1}{12} \sum_j (b_j - a_j)^3.$$

Here we very clearly see the discrepancy between the original potential $\Phi$ and the modified potential $\tilde{\Phi}$:

$$\tilde{\Phi}[Y] = \Phi[Y] - \frac{1}{24} \sum_j (b_j - a_j)^3.$$

Remark. Specialists of nonlinear hyperbolic conservation laws will recognize in the second term of this expression the very expression of the so-called “entropy production” term for the inviscid Burgers equation (0.12), written in material coordinates [6, 22].

3.4. Eulerian formulation of the gradient flow equation. The gradient flow equation (3.41) has an Eulerian version. Indeed, the corresponding measures $(\rho, \rho v)$, defined by (1.24), are (formal) solutions of the following system of PDE:

$$\partial_t \rho - \nabla \cdot (\rho \nabla \tilde{\varphi}) = 0, \quad \rho = \det(I + tD^2_x \tilde{\varphi}),$$

where $\tilde{\varphi}(t, x) = \varphi(t, x)/t$. This model can be seen as a fully nonlinear counterpart of various models popular in biology (chemotaxis) or astronomy, involving the Poisson equation-or other linear equations involving a singular Green function– rather than the Monge-Ampère equation. A common feature of all these models is their ability at describing concentration phenomena [29, 28, 33, 20].

4. The discrete EUR problem with Monge-Ampère gravitation

4.1. A time-discrete scheme for the gradient flow equation. In view of numerical calculations, our first step is to get a time-discrete version of the modified action. Instead of directly getting a discrete version of (3.44), it seems wiser to us to start from a time-discrete version of the gradient flow equation (3.42).

A natural candidate is:

$$X_{n+1} = X_n(1 + \theta_n) - \nabla_H \Pi[X_n] \theta_n + \eta_n$$

where $X_n$ is an approximation of $X(t)$ at the $n^{th}$ time-step $T_n$, for $n = 0, \ldots, N$, $T_0 = t_0$, $T_N = t_1$, with $\theta_n = \frac{T_{n+1}}{T_n} - 1 \downarrow 0$. Here $\eta_n$ is a small perturbation added to the discrete solution so that, for every $n$, $X_n$ is a point of differentiability of $\Pi$. Thus, $\nabla_H \Pi[X_n]$ is well defined and is also the closest point $\pi[X_n]$ to $X_n$ in $S$. Indeed, as a smooth perturbation of a Lipschitz concave function on $H$, $\Pi$ is differentiable on a “fat” dense subset $\tilde{H}$ of $H$ (i.e. containing a countable intersection of dense open sets). Thus, we may choose a
perturbation term \( \eta_n \), arbitrarily small, so that \( X_n \) falls in the "good" set \( \tilde{H} \) where \( \nabla_H \Pi \) is well defined. By doing so, we do not generate a big error. Indeed, we keep control on the cumulated error thanks to the following stability estimate for two distinct solution \( X_n, \tilde{X}_n \) of (4.49) (where we neglect the perturbation terms \( \eta_n, \tilde{\eta}_n \) for notational simplicity),

\[
||X_{n+1} - \tilde{X}_{n+1}||^2 \leq (1 + \theta_n)^2||X_n - \tilde{X}_n||^2 + c_0 \theta_n^2,
\]

where \( c_0 \) is the squared diameter of \( S \). [Indeed, we get form (4.49)

\[
||X_{n+1} - \tilde{X}_{n+1}||^2 = (1 + \theta_n)^2||X_n - \tilde{X}_n||^2 - 2(1 + \theta_n)\theta_n((X_n - \tilde{X}_n, \nabla_H \Pi[X_n] - \nabla_H \Pi[\tilde{X}_n]))
\]

\[
+ \theta_n^2||\pi[X_n] - \pi[\tilde{X}_n]||^2
\]

and observe that the second term in the right-hand side is less than zero since \( \Pi \) is convex, and the third one is dominated by \( c_0 \theta_n^2 \).]

As a matter of fact, this stability estimate is also essentially sufficient to prove the convergence of the scheme as \( \theta_n \downarrow 0 \) to the continuous model (3.43), for the uniform convergence in time with respect to the strong topology of \( H \). (See [13, 14] for examples of similar results for various nonlinear hyperbolic conservation laws.) Notice that concentration phenomena, which are present at the continuous level, are correctly taken into account by the time discrete scheme, in spite of the fact that the discrete scheme never involves the computation of \( d^0 \Pi \), which is a big advantage in practice!

4.2. The time-discrete EUR problem with Monge-Ampère gravitation. From the time-discrete (4.49) for the gradient-flow equation, we define a time discrete version of modified action (3.44) just by setting:

\[
(4.51) \quad \sum_{n=0}^{N-1} r_n ||X_{n+1} - X_n(1 + \theta_n) + \pi[X_n]\theta_n||^2,
\]

with \( r_n = \frac{T_n^{3/2}}{T_{n+1} - T_n} \). Nevertheless, minimizing the time-discrete action (4.51), as both \( X_0 \) and \( X_N \) are fixed, is not really consistent with the original EUR problem, where the data are not \( X_0 \) and \( X_N \) but rather the corresponding probability measures \( \rho_0 \) and \( \rho_N \) defined by:

\[
\rho_0(dx) = \int_D \delta(x - X_0(a))da, \quad \rho_N(dx) = \int_D \delta(x - X_N(a))da.
\]

So there is a big loss of information (since the same probability measure can be generated by a continuum of maps). This problem can be addressed in terms of rearrangements with convex potentials. As a matter of fact, fixing \( \rho_0 \) and \( \rho_N \) is equivalent to fixing the rearrangements with convex potentials \( X_0^* \) and \( X_N^* \), rather than \( X_0 \) and \( X_N \) themselves. It is very fortunate that, one can rewrite the discrete scheme (4.49) as a self-consistent scheme for the rearrangement \( Y_n = X_n^* \) with convex potential. Indeed, let us assume, for simplicity, that, at each \( n \), the solution of the scheme \( X_n \) has a polar factorization \( X_n = Y_n \circ s_n \) (cf. Appendix), where \( Y_n = X_n^* \in K \) is the unique rearrangement with convex potential of \( X_n \) and \( s_n = \pi[X_n] \in S \) is the closest point in \( S \) to \( X_n \). Then, we can rewrite (4.49) as:

\[
Y_{n+1} \circ s_{n+1} = (Y_n(1 + \theta_n) - \theta_n Id) \circ s_n.
\]
But, this implies that $Y_{n+1}$ is the unique rearrangement of $Y_n(1 + \theta_n) - \theta_n Id$ with a convex potential. In other words, we have a well defined self-consistent scheme for $Y_n \in K$, namely:

\[(4.52) \quad Y_{n+1} = (Y_n(1 + \theta_n) - \theta_n Id)^\sharp.\]

Accordingly, the time discrete EUR problem with Monge-Ampère gravitation can be seen, in "polar coordinates" $(Y_n, s_n) \in K \times S$, as the minimization of

\[(4.53) \quad \sum_{n=0}^{N-1} r_n ||Y_{n+1} \circ s_{n+1} - (Y_n(1 + \theta_n) - \theta_n Id) \circ s_n||^2,\]

with $r_n = \frac{T_n^{3/2}}{T_{n+1} - T_n}$, as $Y_0$ and $Y_N$ are fixed in $K$.

Following "optimal transport" theory, we may introduce on $H$ the quadratic Monge-Kantorovich (MK2) (or "Wasserstein") distance,

\[(4.54) \quad d_{MK2}(X, \tilde{X}) = \inf\{||X \circ s - \tilde{X} \circ \tilde{s}||, \quad s, \tilde{s} \in S\},\]

which is nothing but the quotient distance in $H$ with respect to the action of the semi-group $S$. Then the discrete EUR problem is just the minimization in $Y_n \in K$ of

\[(4.55) \quad \sum_{n=0}^{N-1} r_n d_{MK2}(Y_{n+1}, Y_n(1 + \theta_n) - \theta_n Id)^2,\]

with $r_n = \frac{T_n^{3/2}}{T_{n+1} - T_n}$, as $Y_0$ and $Y_N$ are fixed in $K$.

(Equivalently, we could work on the so-called "Wasserstein" or "MK2" space as, for instance, in [34, 3, 2].)

### 4.3. The fully discrete least action principle.

Let us now introduce a fully discrete scheme, for which not only the time variable but also the space variable is discrete. The domain is divided into $L$ disjoints subdomains $D_i$ of Lebesgue measure $1/L$, for $i = 1, \cdots, L$, with barycenter $a_i$ and vanishing diameter as $L \to \infty$. In our abstract framework, it is enough to substitute for the spatial domain $D$, the discrete set $\{a_i, \quad i = 1, \cdots, L\}$. Accordingly, $H$ can be seen as the euclidean space $(\mathbb{R}^d)^L$ of all finite sequences of $L$ points in $\mathbb{R}^d \{X = (X_i \in \mathbb{R}^d)_{i=1,L}\}$ with the natural euclidean norm $|| \cdot ||$ induced by $\mathbb{R}^d$.

Meanwhile the set $S$ can be viewed as the set of all permutations $s$ of the $L$ first integers and $K$ is the corresponding cone of all sequences $Y_i$ such that

$$\sum_i Y_i \cdot (a_i - a_{s_i}) \geq 0,$$

for all permutations $s$. In one space dimension, $K$ is just the convex cone of all increasing sequences of $L$ real numbers.

The time-discrete EUR problem (4.53) makes sense at the fully discrete level without modification. In this discrete setting, $S$ is a group (which is untrue at the continuous level) and each $s$ can be inverted in $S$. Thanks to the group property of $S$ and the invariance of $|| \cdot ||$ with respect to $S$, the minimization problem can be further reduced to the
minimization of
\[
(4.56) \quad \sum_{n=0}^{N-1} r_n ||Y_{n+1} - (Y_n(1 + \theta_n) - \theta_n Id) \circ \sigma_{n+1}||^2,
\]
in $Y_n \in K$, $\sigma_n \in S$, as $Y_0$ and $Y_N$ are fixed in $K$ (just by setting $\sigma_{n+1} = s_n \circ s_{n+1}^{-1} \in S$).
To solve this minimization problem, a crude strategy is to use Gauss-Seidel type iterations.
We denote by $(Y^k_n, \sigma^k_n)$ the approximation of $(Y_n, \sigma_n)$ at iteration $k$ and time step $n$. Let
us fix $k$ and $n$. To get the updated values $\sigma^{k+1}_n$ and $Y^{k+1}_n$, we inductively suppose that we
already know $(Y^j_m, \sigma^j_m)$ for all $m$ if $j \leq k$ and for all $m < n$ if $j = k + 1$. Then, we
perform the two following steps:
i) First step: we get $s = \sigma^{k+1}_n$ by solving the combinatorial optimization problem
\[
(4.57) \quad \inf_{s \in S} ||Y^k_{n+1} - (Y^k_n(1 + \theta_n) - \theta_n Id) \circ s||.
\]
This step is particularly simple in one space dimension and just amounts to sorting in
increasing order the finite sequence $(Y^k_n(1 + \theta_n) - \theta_n Id)_i$, \(i = 1, \ldots, L\). It is much more
challenging in higher dimensions. The best known optimization methods need $O(L^3)$
operations, which is not satisfactory (see a related discussion in [17]).
ii) Second step: we get $Y = Y^{k+1}_n$ by minimizing in $Y \in K$:
\[
r_n ||Y^k_{n+1} - (Y(1 + \theta_n) - \theta_n Id) \circ \sigma^k_{n+1}||^2 + r_{n-1} ||Y - (Y^{k+1}_n(1 + \theta_{n-1}) - \theta_{n-1} Id) \circ \sigma^{k+1}_n||^2,
\]
where the first term can also be written
\[
r_n ||Y^k_{n+1} \circ (\sigma^k_{n+1})^{-1} - (Y(1 + \theta_n) - \theta_n Id)||^2,
\]
using the inverse permutation $(\sigma^k_{n+1})^{-1}$ and the invariance of $|| \cdot ||$ with respect to per-
mutations. After reorganizing squares, we see that $Y$ is just the least-square projection $H \rightarrow K$ of:
\[
V = \frac{r_n(1 + \theta_n) W + r_{n-1} Z}{r_n(1 + \theta_n)^2 + r_{n-1}},
\]
\[
W = Y^k_{n+1} \circ (\sigma^k_{n+1})^{-1} + \theta_n Id, \quad Z = (Y^{k+1}_n(1 + \theta_{n-1}) - \theta_{n-1} Id) \circ \sigma^{k+1}_n.
\]
So, we have obtained an effective algorithm. It is particularly simple in one space dimen-
sion (and much more challenging in higher dimensions!). Let us observe that, in one space
dimension, computing the least-square projection $Y = P_K[V]$ is different from sorting the
sequence $V$ in increasing order. However, still in one space dimension, this projection can
be approximately computed after a sequence of sorting steps, according to the asymptotic
formula (which is a byproduct of the ”transport-collapse method” [10]):
\[
(4.58) \quad P_K[V] = \lim_{M \rightarrow \infty} V^M_M, \quad V^M_m = (V^M_{m-1} + \frac{1}{M} V)^\sharp, \quad V^M_0 = 0, \quad m = 1, \cdots M.
\]
In practice, we already get a good accuracy for moderate values of $M$ (say $M = 10$).

5. Solution of the initial value problem

In order to validate the reconstruction scheme, we would like to solve the initial value
problem (IVP) consistently with the modified least action problem, and get a discrete
scheme for the IVP. Ideally, such a scheme should be derived directly from the modified
discrete least action principle. Unfortunately, we have not been able to do so, and we are
just going to suggest a simple scheme for the IVP which seems, in practice, consistent with the modified action, at least in one space dimension.

5.1. **A time-discrete scheme for the IVP.** Our idea to get a time-discrete solution of the IVP is to alternate the solution of the linear ODE

\[
\frac{d}{dt}(\alpha^2(t)V) = \beta^2(t)(Y - Id), \quad \frac{dY}{dt} = V, \quad Y(T_n) = Y_n, \quad V(T_n) = V_n,
\]

with coefficients \((\alpha, \beta)\) given by (0.3), on each time interval \([T_n, T_{n+1}[^2\] and the rearrangement of the result at time step \(T_{n+1}[^2\):

\[
Y_{n+1} = Y(T_{n+1}), \quad V_{n+1} = V(T_{n+1}).
\]

Using a plain explicit discretization of (5.59), we get:

\[
Y_{n+1} = (Y_n + (T_{n+1} - T_n)V_n)^2, \quad \alpha^2(T_{n+1})V_{n+1} = \alpha^2(T_n)V_n + \beta^2(T_n)(T_{n+1} - T_n)(Y_n - Id)
\]

The convergence analysis of this time-discrete scheme can be done in two different ways.

5.2. **The multidimensional case.** In the multidimensional case, our strategy for the convergence analysis of scheme (5.60) is inspired by our recent work [15], where a similar scheme is analyzed. We essentially use the fact that all maps with convex potential are of locally bounded variations, which provides enough compactness with respect to space variables. Time compactness is, as usual, directly obtained from the evolution scheme. We notice that \(V\) can be easily integrated out from \(Y\) by ODE (5.59).

**Theorem 5.1.** For every fixed initial condition \((Y_0, V_0) \in Y \times H\), the approximate solution \(Y_n\) admits at least a limit \(t \to Y(t) \in K\) valued in \(C^0([t_0, +\infty[, H)\) as the time step goes to zero. This limit satisfies:

\[
V(t, a)\alpha^2(t) = V_0(a)\alpha^2(t_0) + \int_{t_0}^t \beta^2(\tau)(Y(\tau, a) - a)d\tau,
\]

\[
\frac{d}{dt} \int_D f(Y(t, a))da = \int_D (\nabla f)(Y(t, a)) \cdot V(t, a)da,
\]

for all \(C^1\) function \(f\) on \(R^d\) (with \(|\nabla f|\) growing at most linearly at infinity).

Notice that, since \(Y\) is valued in \(K\), the knowledge of ”observables” \(\int_D f(Y(t, a))da\) for all test-functions \(f\) is enough to determine \(Y(t)\) which makes formulation (5.61) self-consistent. (However this does not guarantee uniqueness of solutions to the IVP.) So, we have a proposal to solve the IVP, and a corresponding discrete scheme, but we are not able to prove that formulation (5.61) is actually consistent with our modified least action principle.

5.3. **The one-dimensional case.** In the special case of one space variable, we get a much more precise information, following the analysis developed in [13] for similar problems (see also [14]):
Theorem 5.2. For every fixed initial condition \((Y_0, V_0) \in Y \times H\), as the time step goes to zero, the approximate solution \((Y_n)\) converges to the unique solution \(t \to Y(t) \in K\), valued in \(C^0([t_0, +\infty, H]\), of the mixed integral-sub-differential system:

\[-\partial_t Y + V \in \partial \Theta[Y],\]

where \(\Theta[Y] = 0\) whenever \(Y = Y(t, a)\) is monotonically increasing in \(a\) and \(\Theta[Y] = +\infty\) otherwise.

System (5.62) is well-posed in the \(L^2\) sense and can be shown (as in [13]) to be the limit (in the sense of maximal monotone operator theory) as \(\epsilon \downarrow 0\) of the perturbed system

\[-\partial_t Y + V = -\epsilon \partial_a (\log(\partial_a Y)), \quad \partial_t (\alpha^2(t)V) = \beta(t)(Y - Id),\]

which, in Eulerian variables (1.24), reduces to:

\[
\begin{align*}
\partial_t \rho + \partial_x (\rho v) &= 0, \\
\partial_t (\alpha^2(t) \rho v) + \partial_x (\alpha^2(t) \rho v^2) &= -\beta^2(t)\rho \partial_x \varphi + \epsilon \partial_{xx} v, \\
\rho &= 1 + \partial_x^2 \varphi,
\end{align*}
\]

and is just a pressure-less Navier-Stokes-Poisson system with vanishing viscosity (as in [9, 36]). As mentioned above, an interesting open question is to show that this approach (vanishing viscosity (5.64), subdifferential formulation (5.62) or scheme (5.60)) is actually consistent with the modified least action principle!

6. Numerical simulations in one space dimension

Our data are

\[t_0 = 1/2, \quad t_1 = 5/2, \quad N = 60, \quad L = 51, \quad a_i = -1 + (2i - 1)/L, \quad i = 1, \ldots, L,\]

\[Y_0 = X_0^\sharp, \quad (X_0)_i = a_i \omega_i, \quad ,\]

where \(\omega_i\) is a random number uniformly distributed between 1 and 2. Thus, \(Y_0\) looks like a devil’s staircase.

Concerning the final data \(Y_N \in K\), either:

(Case 1) the associate probability \(\rho_N\) is the barycenter of four Dirac’s measures:

\[\rho_N(dx) = \frac{\delta(x + 0.7) + 4 \delta(x - 0.2) + 3 \delta(x - 0.9) + \delta(x - 1.1)}{9}\]

or (Case 2) \(Y_N\) is the solution at time \(t_1\) of the initial value problem generated by the discrete gradient flow equation starting from \(Y_0 = X_0^\sharp\) at time \(t_0\).

In our plots, we draw the trajectories of the 51 particles during the 60 time steps of the time interval (the vertical axis corresponding to time and the horizontal one to space).

Case 1: we first plot the reconstructed solution (fig. 1). Then, with the reconstructed initial velocity we solve the initial value problem for the MAG equations with scheme (5.60) and plot the result (fig.2). We observe a nearly perfect match between figures 1 and 2.
Case 2: we first solve the IVP for the gradient flow equation with scheme (5.60) (fig. 3). Then, we reconstruct the solution from the initial and final data of the gradient flow solution (fig 4). (Here we observe some limited discrepancy.) Finally, with the reconstructed initial velocity we solve the initial value problem for the MAG equations (fig.5), again with scheme (5.60) and get a nearly perfect match.
7. Discussion

We have revisited the early universe reconstruction problem and suggested a modification of classical Newton gravitation by what we called Monge-Ampère gravitation. The main drawback of our approach is the lack of physical justification for such a modification. The main mathematical advantage is the obtention of a modified least action principle in which we can easily include mass concentration effects in an almost canonical way, using ideas from gradient flow theory. In addition, the well-known Zeldovich approximate solutions turn out to be exact solutions of the modified model, which provides an indirect validation of the model as a reasonable approximation for the early universe reconstruction (EUR) problem. According to these ideas, an algorithm has been designed in the 1D case. Our plan for the future includes: i) analysis of the initial value problem, consistently with the modified least action principle; ii) design of an efficient multidimensional algorithm; iii) study of the relative accuracy of the Newton and Monge-Ampère gravitation models with respect to general relativity.
Figure 3. Case 2/gradient flow solution -IVP before reconstruction

Figure 4. Case 2/gradient flow solution -reconstruction

Figure 5. Case 2/gradient flow solution -IVP after reconstruction
8. Appendix

8.1. Some useful results from optimal transport theory. The set $S$ defined by (1.22) has a semi-group structure for the composition rule and has the identity map $Id$ as neutral element. It is, in some sense, in duality with its "polar cone" $K \subset H$:

\[(8.65)\quad K = \{Y \in H : ((Y, Id - s)) \geq 0, \ \forall s \in S\} .\]

Let us recall few basic results of optimal transport theory [11, 12, 37] concerning $S$ and $K$. First, the set $K$ can be characterized as the closed convex cone of all maps $Y$ with a convex potential, which means that there is a convex function $\psi$ defined on $\mathbb{R}^d$ and valued in $]-\infty, +\infty]$ which is almost everywhere differentiable on $D$ with $\nabla \psi(x) = Y(x)$, a.e. on $D$.

Next, every map admits a unique rearrangement in $K$. More precisely

**Theorem 8.1.** ([11]) Every $X \in H$ admits a unique "rearrangement" $X^\sharp$ in $K$, which means:

\[\int_D f(X^\sharp(a))da = \int_D f(X(a))da, \ \forall f \in C(\mathbb{R}^d), \ \sup_x |f(x)|(1 + |x|^2)^{-1} < +\infty.\]

In addition, $X \rightarrow X^\sharp$ is continuous in $H$ (for the strong topology).

Moreover, there is a "polar factorization" of the Hilbert space $H$ by $S$ and $K$. More precisely:

**Theorem 8.2.** ([11]) Let $X \in H$ be a non degenerate map, in the sense that the measure $\rho(dx) = \int_D \delta(x - X(a))da$ has no singular part with respect to the Lebesgue measure. Then $X$ admits a unique "polar factorization"

\[(8.66)\quad X = Y \circ s, \ Y \in K, \ s \in S.\]

In addition, the second factor $s$ is characterized as the unique closest point $\pi[X]$ to $X$ in $S$ and can be written

\[(8.67)\quad \pi[X] = \nabla \psi \circ X,\]

where $\nabla \psi$ is the unique map $T : \mathbb{R}^d \rightarrow D$ with convex Lipschitz potential such that the Lebesgue measure restricted to $D$ is the image of $\rho$ by $T$:

\[(8.68)\quad \int_{\mathbb{R}^d} f(\nabla \psi(x))\rho(dx) = \int_D f(a)da, \ \forall f \in C(\mathbb{R}^d).\]

Let us finally observe as in [11, 12] that (8.68) can be seen as a "weak formulation" (not in the sense of distributions!) of the Monge-Ampère problem on $\mathbb{R}^d$ with range condition:

\[(8.69)\quad \rho = \det(D^2 \psi), \quad (\nabla \psi)(\mathbb{R}^d) = D.\]

8.2. Formal derivation of the MAG system. Let us first notice that the measures $(\rho, \rho v)$ defined by (1.24) automatically satisfy the mass conservation equation in the sense of distributions:

\[(8.70)\quad \partial_t \rho + \nabla \cdot (\rho v) = 0.\]
Next, let us assume that $X$ is non degenerate at time $t$ (which exactly means that $\rho(t, dx)$ has no singular part with respect to the Lebegseue measure).

Then

$$\nabla_H \Phi[X(t)] = X(t) - \pi[X(t)] = (Id - \nabla \psi(t, \cdot)) \circ X(t),$$

as in Theorem 8.1 and we can write the MAG equation (1.23):

$$\beta^{-2}(t) \frac{d}{dt}(\alpha^2(t) \frac{dX}{dt}) = (Id - \nabla \psi) \circ X,$$

where $\nabla \psi(t, \cdot)$ is the unique ”weak solution” of the Monge-Ampère problem (8.69).

In terms of $(\rho, v)$, (8.71) formally means

$$\partial_t (\alpha^2 \rho v) + \nabla \cdot (\alpha^2 \rho v \otimes v) = -\rho \beta^2(t) (x - \nabla \psi),$$

which, combined with (8.69), coincides with the Eulerian formulation (1.25) of the MAG model, with $\varphi(t, x) = \psi(t, x) - |x|^2/2$. Of course, this derivation is purely formal: not only concentrations have not been taken into account, but the derivation of (8.72) from (8.71) is far from being legitimate.

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References


