REMARKS ON OLDROYD-B AND RELATED COMPLEX FLUID MODELS

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To Dave Levermore on his 60th birthday with friendship and appreciation

Abstract. We prove global existence and uniqueness of solutions of Oldroyd-B systems, with relatively small data in \( \mathbb{R}^d \), in a large functional setting \( (C^\alpha \cap L^1) \). This is a stability result; solutions select an equilibrium and converge exponentially to it. Large spatial derivatives of the initial density and stress are allowed, provided the \( L^\infty \) norm of the density and stress are small enough. We prove global regularity for large data for a model in which the potential responds to high rates of strain in the fluid. We also prove global existence for a class of large data for a didactic scalar model which attempts to capture, in the simplest way, the essence of the dissipative nature of the coupling to fluid. This latter model has an unexpected cone invariance in function space that is crucial for the global existence.

Key words. Oldroyd-B, complex fluids, Fokker-Planck equations, blow up, global existence, Euler equations, Navier-Stokes equations, kinetic equations.

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1. Introduction

The complex fluids models we consider treat the interaction of a normal fluid (incompressible in our case) with insoluble matter suspended in it. Models have been devised to deal with microscopic elastic thread-like objects such as polymers [12, 26]. The complicated hydrodynamic interactions are simplified using the separation of scales, replacing the many degrees of freedom due to them by a few representative ones. In the models we consider here, an end-to-end vector in \( \mathbb{R}^3 \) (\( m \) in this paper) represents the orientation of the polymer and it is the sole variable retained to describe the microscopic object. The models that are most studied with this degree of simplification include the kinetic description of the Oldroyd-B model, the FENE model and variants. The probability distribution of the vectors \( m \) depends on time and physical space, and it is assumed to be absolutely continuous with respect to the usual Lebesgue measure \( dm \), so the chance of finding \( m \) at time \( t \) and location \( x \in \mathbb{R}^3 \) in the volume element \( dm \) is \( f(x,m,t)dm \). The aim is to describe the evolution of \( f \) in space and time. When the polymer concentration is so dilute that the polymers do not interact, and when the fluid is smooth, then the equation of \( f \) is a linear Fokker-Planck equation. The effect of the fluid can be understood perturbatively: because the objects suspended are assumed to have faster time scales and shorter length scales than the scales on which solvent varies, it is then justified to treat the fluid as varying little as far as the microscopic suspended objects are concerned. This macro-micro interaction can be rationally discussed, and leads to a kinetic description of the particles, given a fluid flow. On the other hand, the collective effect that the particles have on the flow itself is a macroscopic effect of microscopic insertions: a micro-macro interaction. These interactions are much more mysterious. There have been attempts to produce systematic upscalings based on non-equilibrium thermodynamics formalisms [2, 13]. When the microscopic insertions have a larger number

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of degrees of freedom, \( m \in M \), where \( M \) is a Riemannian manifold representing a finite number of degrees of freedom with constraints, then the kinetic equation is a Smoluchowski equation on the manifold. The upscaling principle advocated in [5, 6] is easiest formulated as the requirement that the sum of the energy of the fluid and the free energy of the microscopic quantities be a Lyapunov functional for the coupled system. In the known examples, this requirement leads to familiar rules of determining the added polymeric stress from the micro-micro and the macro-micro interactions [16]. The mathematical treatment of the coupled systems is far from being complete but has recently received increased attention. Early work [14, 28] established local existence results for Oldroyd-B and FENE type equations. The problem of existence of weak solutions is rather open. Global existence of weak solutions in the presence of spatial diffusion of the polymers was proved in a sequence of papers, the most recent of which is [1]. Global existence of weak solutions via propagation of compactness was proved under the corotational assumption [22, 23] and, very recently, for the full FENE model [25]. There is no such result for the Oldroyd-B model. The global existence of smooth solutions for small data for Oldroyd-B type models was established in [17, 21].

Global existence of smooth solutions for large data in 2D was established for Smoluchowski equations on compact manifolds [6, 7, 9, 10, 11, 27]. Global regularity for large data in the FENE case, under the corotational assumption, was proved in [19, 24]. An approach based on Lagrangian particle dynamics was described in [20]. Sufficient conditions for regularity in terms of bounds on the added stress tensor were established in [4, 15] and further refined in [18]. Numerical evidence for singularities was provided in [30].

In this paper we address issues related to global existence of smooth solutions in simple kinetic models. In the presence of a quiescent solvent, the polymer distribution is the unique stable time independent solution of the linear Fokker-Planck equation

\[
\partial_t f = \epsilon \text{div}_m (f \nabla_m (\log f + U)),
\]

i.e. \( f = Z^{-1} e^{-U} \) with \( Z \) a normalizing constant. The simplest form of \( U \) is harmonic, \( U(m) = \frac{|m|^2}{2R^2} \) and then the solution is Gaussian. The constant \( R^2 \) represents the expected value \( \langle |m|^2 \rangle \) of the square of the length of the end-to-end vector. In the presence of smooth fluids the kinetic equation changes and acquires a dependence on the macroscopic independent variable \( x \).

In the second section of the paper we provide a priori bounds for linear Fokker-Planck equations with Oldroyd-B type potentials. We do this by essentially giving a formula for the solution of the Fokker-Planck equation in terms of the particle trajectories of the underlying flow.

In the third section we derive a classical estimate for the velocity gradient in terms of the added polymeric stress, when the coupling to fluid is done via time-independent Stokes equations.

In the fourth section we prove global existence results for small data in \( C^\alpha \) at arbitrary Deborah numbers. The previously known global existence results for small data are done using energy estimates and require derivatives of the stress. The main difficulty in obtaining bounds for small data is the fact that quantities are not monotonically decreasing, and, in particular, the spatial gradient of the density of particles can grow in time but nevertheless saturates. The system formed by the kinetic equation for the particle distribution and the velocity of the fluid can be reduced to a closed system for a reduced added stress \( \tau \) and the particle density. This system looks like a
damped and driven Euler equation where \( \tau \) (a matrix) plays the role of vorticity and \( \rho \), the transported density, is part of the driver. The global existence of solutions is proved under the assumption that the \( L^\infty \) norm of the initial data is small enough (4.82). The class of initial data that lead to global existence includes densities and stresses with arbitrarily large gradients, provided the \( L^\infty \) norms are suitably small. While the limit added stress \( \tau \) and limit velocity vanish, the limit particle density is a rearranged initial particle density — a deformation of the initial density along the limit back-to-labels map. The class of solutions we discuss is in fact rather wide, and in particular the gradients need to be understood in the distributional sense. The proof of uniqueness of solutions is done using Lagrangian transformations. An Eulerian proof is not readily available at this moment. In order to explain this we pursue the analogy with Euler equations. Both the damped Euler vorticity equation and the equation for \( \tau \) have the form

\[
D_t \tau = -c \tau + N(\tau),
\]

where \( c \geq 0 \) is the damping and the nonlinearity \( N \) is quadratic and non-local, obtained from products of \( \tau \) and principal value singular integrals. Here \( D_t \) is the material derivative along the divergence-free velocity \( u \). The local (and global for small data) existence results are in Banach algebras in which \( N \) is a continuously differentiable nonlinear function. One of the largest and most natural such Banach algebras is \( C^\alpha \), \( 0 < \alpha < 1 \). The difference \( \delta \) of two solutions \( \tau_1, \tau_2 \) obeys the equation

\[
\mathcal{D}_t \delta = -c \delta + N'(\tau) \delta - v \cdot \nabla_x \tau,
\]

where \( \mathcal{D}_t \) is transport along the average velocity \( \tau = \frac{1}{2}(\tau_1 + \tau_2) \), \( \tau \) is the average of the solutions, \( \tau = \frac{1}{2}(\tau_1 + \tau_2) \), \( N'(\tau) \) is the linearization of \( N \) obtained using polarization, and \( v \) is the velocity difference \( v = u_1 - u_2 \). The last term in the equation is present because the velocities depend on \( \tau \) and \( u_1 \) is obtained from \( \tau_1 \) and \( u_2 \) from \( \tau_2 \). In both the (damped-driven) Euler case and in the Stokes-Fokker-Planck system we study, the dependence of \( \tau \) on \( \tau \) is linear. The term \( v \) depends therefore linearly on \( \delta \). Uniqueness proofs based on energy estimates proceed by estimating the growth of \( \delta \) in some norm. As the initial data vanishes, any closed, super-linear estimate of the norm is sufficient to prove uniqueness. (We could allow logarithmically sublinear estimates, but that is as far as that can go, in general). There are two sources of difficulty in obtaining these estimates. The transport term \( \bar{\tau} \cdot \nabla_x \delta \) can be dealt with by using integration by parts, if the norm of \( \delta \) we consider is an \( L^p \) norm. The other source of difficulty is \( \nabla_x \tau \). If we are in the framework of \( C^\alpha \cap L^1 \), adopted in this paper, then \( v \) is very nice but in order to deal with \( \nabla_x \tau \) we need to work in \( H^{-1} \). In general, the presence of both \( v \cdot \nabla_x \tau \) and \( \bar{\tau} \cdot \nabla_x \delta \) together is lethal for an Eulerian energy method approach unless we work in a space in which \( \nabla_x \tau \) is tame, like for instance \( W^{1,p}(\mathbb{R}^d) \) with \( p > d \). The reason there is a successful Eulerian proof of uniqueness for Euler equations for the vorticity in \( C^\alpha \) is that, for Euler equations, the vorticity equation comes from a velocity equation with a good cancellation property. Thus we can “retreat” to one less derivative in the equation for \( \delta \). In terms of the vorticity equation this means integrating against a vectorial stream function, i.e. working in \( H^{-1} \). In the Euler equations the transport term together with the term \( N'(\tau) \delta \) can be integrated by parts one more time, revealing only \( H^{-1} \) terms in \( \delta \). In the absence of such algebraic results, the retreat to one less derivative does not appear to work, and the fix for one difficult term does not work for the other, and vice-versa. The Lagrangian approach introduces commutators that are well behaved, and is successful.
In the fifth section we show that the system is regularized for large data if the potential is allowed to respond to excessively high rates of strain \( S \) in the fluid. The Oldroyd-B potential is harmonic, \( U = \frac{|m|^2}{2R^2} \), although it is not generated by a spring force that is related to a material property of the polymers. In fact, the representation of the polymer by \( m \) is over-simplified, and \( R \) has something to do with an average restoring force in the ensemble and not with the maximum extension allowed for an individual molecule. In the case there is feedback coupling to a fluid, the interaction might have properties that depend of the local properties of the fluid. We show that if we allow \( R \) to depend on the local rate of strain in the fluid, when this exceeds a threshold, and to grow with it (\( \frac{\partial R}{\partial S} = \delta(S) \)), then the equation has global smooth solutions for all large data. Here \( \delta(S) \) can vanish if the rate of strain \( S \) is not too high, and becomes asymptotically linear in \( S \) for high values of (norms of) \( S \). Allowing \( R \) to grow in response to high rate of strain is like allowing more entropic “slack” in the molecules at high rates of strain, and this turns out to be a mathematical regularizing mechanism.

In the sixth section we return to the Oldroyd-B large data problem, and give an example of blow up in one dimension for a crude model that does not respect the principle of free energy decay. We motivate and give then a simple scalar example of an equation that has some of the features of the type of equations arising in the large data Oldroyd B equation, and which has an unbounded set of initial data that lead to global smooth behavior. The reason for the global existence is an invariance of a cone in function space under the nonlinear evolution. The consequence of this invariance is an a priori bound that is sufficient for the persistence of smoothness of solutions.

### 2. Linear Fokker-Planck equation: kinematic observations

We take a vector field \( u(x, t) \) and a scalar \( f(x, m, t) \) representing a two or three-dimensional incompressible velocity and a particle distribution. We start by describing the simplest particle distribution equation:

\[
D_t f + ((\nabla_x u)m) \cdot \nabla_m f = \epsilon \text{div}_m (f \nabla_m (\log f + U(m))).
\]

(2.1)

Here \( D_t = \partial_t + u \cdot \nabla_x \) is the material derivative. The potential is given by

\[
U(m) = \frac{|m|^2}{2R^2}.
\]

(2.2)

In this section \( R \) is a positive constant. We consider here the case \( m \in \mathbb{R}^d \). We associate to the particle distribution \( f \) an added stress tensor:

\[
\sigma(x, t) = \int_{m \in \mathbb{R}^d} (m \otimes \nabla_m U) f(x, m, t) dm.
\]

(2.3)

In this normalization \( R \) has units of length and \( \epsilon/R^2 \) is an inverse time. The potential \( U \) and the stress \( \sigma \) are nondimensional. We study first (2.1) without regard to the coupling to the fluid. In what follows the velocity field could be quite arbitrary, but it is assumed to be smooth enough for the calculations below: the variable \( x \in \mathbb{R}^d \), and the velocity field \( u(x, t) \) is in \( C^{1,\alpha}(\mathbb{R}^d) \), is divergence-free,

\[
\text{div}_x u = 0,
\]

and decays at infinity. Before we start our computations, let us make a few general comments regarding the matrix \( \sigma \). First, because the potential \( U \) is radially symmetric
(2.2) and from its definition (2.3), it follows that $\sigma$ is symmetric and non-negative. Moreover its off-diagonal entries are bounded by the trace

$$|\sigma^{ij}(x,t)| \leq \frac{1}{2} Tr(\sigma(x,t)). \tag{2.4}$$

We take now the Fourier transform in $m$:

$$\hat{f}(x,t,\xi) = \int_{\mathbb{R}^d} e^{-im \cdot \xi} f(x,m,t) dm. \tag{2.5}$$

In view of (2.2), $\hat{f}$ obeys

$$D_t \hat{f}(x,\xi,t) + \left[ \frac{\epsilon R^2}{I} - (\nabla_x u)^T \right] \xi \cdot \nabla_\xi \hat{f} = -\epsilon |\xi|^2 \hat{f}(x,\xi,t). \tag{2.6}$$

We used the fact that $u$ is divergence-free and denoted by $(\nabla_x u)^T$ the transposed matrix. We solve (2.6) on characteristics. The connection between the Lagrangian paths of $u$ and the Oldroyd-B equation was exploited in [20] in the less complicated case $\epsilon = 0$, $\rho = 1$. Let $X(a,t)$ be the particle paths

$$\partial_t X(a,t) = u(X(a,t),t), \tag{2.7}$$

with $X(a,0) = a$, and let

$$g(a,t) = (\nabla_x u)(X(a,t),t). \tag{2.8}$$

Let now $\xi = \xi(a,\eta,t)$ solve the ODE

$$\frac{d}{dt} \xi = \frac{\epsilon}{R^2} \xi - g(a,t)^T \xi \tag{2.9}$$

with initial data $\xi(a,\eta,0) = \eta$. We take the fundamental matrix $\Phi(a,t)$, the solution of the linear ODE system

$$\frac{d}{dt} \Phi(a,t) = -g(a,t)^T \Phi(a,t) \tag{2.10}$$

with initial data $\Phi(a,0) = I$, and we then have

$$\xi(a,\eta,t) = e^{\frac{\epsilon t}{R^2}} \Phi(a,t) \eta, \tag{2.11}$$

We write

$$F(a,\eta,t) = \hat{f}(X(a,t),\xi(a,\eta,t),t). \tag{2.12}$$

Then (2.6) implies

$$\frac{d}{dt} F(a,\eta,t) = -\epsilon |\xi(a,\eta,t)|^2 F(a,\eta,t),$$

and by integrating we obtain

$$F(a,\eta,t) = e^{-\epsilon \int_0^t |\xi(a,\eta,s)|^2 ds} \hat{f}(a,\eta,0). \tag{2.13}$$
Invert the linear map $\eta \mapsto \xi(a,\eta,t)$ and write, from (2.11),
\[
\eta(a,\xi,t) = e^{-\frac{\eta}{H^2}} \Psi(a,t)\xi,
\]
where $\Psi(a,t) = \Phi(a,t)^{-1}$ obeys
\[
\frac{d}{dt} \Psi(a,t) = \Psi(a,t)g(a,t)^T
\]
with initial data $\Psi(a,0) = I$. Reading (2.12) at $\eta = \eta(a,\xi,t)$ for a fixed $\xi$ and substituting in (2.13) we obtain
\[
\hat{f}(X(a,t),\xi,t) = e^{-\epsilon \int_0^t |\xi(a,\eta(a,\xi,t),s)|^2 ds} \hat{f}_0(a,e^{-\frac{\eta}{H^2}} \Psi(a,t)\xi),
\]
where $\hat{f}_0(x,\xi)$ is the Fourier transform in $m$ of the initial data $f_0(x,m) = f(x,m,0)$. Let us consider now the “back-to-labels” map $A(x,t)$, the inverse of $X(a,t)$. Let
\[
M(x,t) = \Psi(A(x,t),t)
\]
and note that it obeys the transport equation
\[
D_t M = M(\nabla_x u)^T
\]
with initial data $M(x,0) = I$. Reading (2.16) at $a = A(x,t)$, we deduce
\[
\hat{f}(x,\xi,t) = e^{-\epsilon \int_0^t |\xi(A(x,t),\eta(A(x,t),\xi,t),s)|^2 ds} \hat{f}_0(A(x,t),e^{-\frac{\eta}{H^2}} M(x,t)\xi).
\]
We compute, directly from (2.11) and (2.14),
\[
\xi(A(x,t),\eta(A(x,t),\xi,t),t,s) = e^{-\frac{\eta(a,s)}{H^2}} Q(x,t,s)\xi
\]
with
\[
Q(x,t,s) = q(A(x,t),t,s), \quad \text{for } t \geq s,
\]
\[
q(a,t,s) = \Phi(a,s)\Psi(a,t) \quad \text{for } t \geq s,
\]
and deduce thus from (2.19)
\[
\hat{f}(x,\xi,t) = e^{-\epsilon \int_0^t e^{-2\frac{\eta(a,s)}{H^2}} |Q(x,t,s)|^2 ds} \hat{f}_0(A(x,t),e^{-\frac{\eta}{H^2}} M(x,t)\xi).
\]
Let us note now that the matrix $\sigma^{ij}(x,t)$ is computed from the Hessian of $\hat{f}$ at $\xi = 0$,
\[
\sigma^{ij}(x,t) = -\frac{1}{H^2} \frac{\partial^2 \hat{f}}{\partial \xi_i \partial \xi_j} (x,\xi,t)|_{\xi=0}.
\]
Using (2.23) and noting that the cross-terms vanish, we deduce
\[
\sigma^{ij}(x,t) = \frac{2\epsilon}{H^2} \rho(x,t) \int_0^t e^{-\frac{2\epsilon|\xi|^2}{H^2}} [Q(x,t,s)Q(x,t,s)]_{ij} ds
\]
\[
- \frac{2\epsilon}{H^2} M_{ki}(x,t) M_{lj}(x,t) \frac{\partial^2 f_0}{\partial \xi_i \partial \xi_j} (A(x,t),\xi)|_{\xi=0}.
\]
We used the notation \( \rho(x, t) := \hat{f}(x, 0, t) = \int_{\mathbb{R}^d} f(x, m, t) dm \). In view of (2.1) this obeys the transport equation
\[
D_t \rho = 0, \tag{2.26}
\]
and therefore it is given by
\[
\rho(x, t) = \hat{f}_0(A(x, t)) = \rho_0(A(x, t)) \tag{2.27}
\]
in terms of the initial particle density at \( x \). We used this in (2.25) as well as the summation convention. In terms of the initial stress, the expression is
\[
\sigma^{ij}(x, t) = \frac{2\varepsilon}{R^2} \rho(x, t) \int_0^t e^{-\frac{2\varepsilon(t-s)}{R^2}} \left[ Q^T(x, t, s)Q(x, t, s) \right]_{ij} ds
+ e^{-\frac{2\varepsilon t}{R^2}} M_{ki}(x, t)M_{lj}(x, t)\sigma^{kl}(A(x, t), 0). \tag{2.28}
\]
The solution of (2.28) solves the equation
\[
D_t \sigma = -\frac{2\varepsilon}{R^2} \sigma + (\nabla_x u)\sigma + \sigma(\nabla_x u)^T + \frac{2\varepsilon}{R^2} \rho(x, t)I, \tag{2.29}
\]
which can be easily derived from the Equation (2.1) by multiplying by \( \frac{1}{R^2} m_i m_j \) and integrating with respect to \( dm \). Let \( \theta(x, t) \) be a passive scalar, i.e., a solution of
\[
D_t \theta = 0. \tag{2.30}
\]
If \( M \) solves (2.18) then \( M\nabla_x \theta \) is again a passive scalar, i.e.,
\[
D_t (M \nabla_x \theta) = 0. \tag{2.31}
\]
This can be easily checked because
\[
D_t (\partial_j \theta) = -\partial_l (u_l \partial_j \theta),
\]
and so
\[
D_t (M_{ij} \partial_j \theta) = M_{ik}(\partial_k u_j)\partial_j \theta - M_{ij}(\partial_j u_l)\partial_l \theta = 0.
\]
This means that if \( \theta(x, t) = \theta_0(A(x, t)) \) with arbitrary smooth \( \theta_0(a) \) then
\[
M(x, t)\nabla_x \theta(x, t) = (\nabla_a \theta_0)(A(x, t)).
\]
In fact, because the initial datum of \( \Psi \) is the identity matrix, and because of uniqueness of ODEs, it follows from (2.15) that
\[
\Psi(a, t) = (\nabla_a X)^T(a, t), \tag{2.32}
\]
and consequently
\[
M_{ki}(x, t) = \frac{\partial X^i}{\partial a_k}(A(x, t), t). \tag{2.33}
\]
From (2.28) and (2.33) it follows that
\[
\sigma(x, t) = \frac{2\varepsilon}{R^2} \rho(x, t) \int_0^t e^{-\frac{2\varepsilon(t-s)}{R^2}} Q^T(x, t, s)Q(x, t, s) ds
+ e^{-\frac{2\varepsilon t}{R^2}} (\nabla_a X(A(x, t), t))\sigma_0(A(x, t))(\nabla_a X(A(x, t), t))^T. \tag{2.34}
\]
Let us introduce the notations
\[ e^i(x, t) := (\nabla X^i)(A(x, t), t), \] (2.35)
\[ e^{ij}(x, t) := \sigma_0(A(x, t)) e^i(x, t) \cdot e^j(x, t). \] (2.36)

Note that
\[ \sigma(x, t) = \frac{2\epsilon}{R^2} \rho(x, t) \int_0^t e^{-\frac{2\epsilon(t-s)}{R^2}} Q^T(x, t, s)Q(x, t, s)ds + e^{-\frac{2\epsilon t}{R^2}} (e^{ij}(x, t))_{ij}. \] (2.37)

Moreover,
\[ D_t e^i = (\nabla_x u) e^i \] (2.38)
with initial data that are constant in space and equal the canonical basis of \( \mathbb{R}^d \),
\[ e^i(x, 0) = (\delta_{ij})_j. \] Note that \( q \) given in (2.22) is given in terms of the gradient \( \nabla_a X \) by
\[ q(a, t, s) = [(\nabla_a X(a, s))^\top]^{-1}(\nabla_a X(a, t))^\top, \] (2.39)
solving
\[ \partial_s q(a, t, s) = -g^T(a, s)q(a, t, s) \] (2.40)
and
\[ q(a, t, t) = I, \quad q(a, t, 0) = (\nabla_a X(a, t))^\top. \] (2.41)

Passing to Lagrangian variables in (2.37) we obtain
\[ \sigma(X(a, t), t) = \frac{2\epsilon}{R^2} \rho_0(a) \int_0^t e^{-\frac{2\epsilon(t-s)}{R^2}} q^T(a, t, s)q(a, t, s)ds \]
\[ + e^{-\frac{2\epsilon t}{R^2}} (\nabla_a X(a, t))\sigma_0(a)(\nabla_a X(a, t))^\top, \] (2.42)
and by integrating by parts we deduce
\[ \sigma(X(a, t), t) - \rho_0(a)I = 2\rho_0(a) \int_0^t e^{-\frac{2\epsilon(t-s)}{R^2}} q^T(a, t, s)S(X(a, s), s)q(a, t, s)ds \]
\[ + e^{-\frac{2\epsilon t}{R^2}} (\partial_a X(a, t))\tau_0(a)(\partial_a X(a, t))^\top, \] (2.43)
where we introduced the reduced stress
\[ \tau(x, t) = \sigma(x, t) - \rho(x, t)I \] (2.44)
and where
\[ S(x, t) = \frac{1}{2} \left[ (\nabla_x u(x, t)) + (\nabla_x u(x, t))^\top \right] \] (2.45)
is the rate of strain. Returning to Eulerian variables, we have
\[ \tau(x, t) = 2\rho(x, t) \int_0^t e^{-\frac{2\epsilon(t-s)}{R^2}} Q^T(x, t, s)S(X(A(x, t), s), s)Q(x, t, s)ds \]
\[ + e^{-\frac{2\epsilon t}{R^2}} (\nabla_a X(A(x, t), t))\tau_0(A(x, t))(\nabla_a X(A(x, t), t))^\top. \] (2.46)
It is clear that in general, once the velocity \( u(x, t) \) is given, we can compute everything we need to know about \( \tau \) (and hence \( \sigma \)) from (2.46). In particular, we will be interested in the relationship between various norms of \( u \) and norms of \( \tau \), as these relationships will serve in establishing bounds for solutions of nonlinear equations. It is clear that we should start with the fields \( e_i \), and in order to understand their norms we perform some standard Lagrangian estimates. We take the finite difference

\[
\delta_h \left( \partial_a X^i(a, t) \right) = \partial_a X^i(a + h, t) - \partial_a X^i(a, t).
\]

Its equation follows from (2.7) by differentiation,

\[
\partial_t (\nabla_a X^i) = g(a, t)(\nabla_a X^i),
\]

and then by taking the finite difference:

\[
\frac{d}{dt} \delta_h \left( \partial_a X^i(a, t) \right) = \frac{g(a, t) + g(a + h, t)}{2} \delta_h \left( \partial_a X^i(a, t) \right) + \delta_h g(a, t) \frac{\partial_a X^i(a + h, t) + \partial_a X^i(a, t)}{2},
\]

where \( g(a, t) \) is the velocity gradient in Lagrangian coordinates (2.8) and

\[
\delta_h g(a, t) = g(a + h, t) - g(a, t).
\]

For matrices \( L \) we use the notation \( |L| \) for the Euclidean norm of the matrix \( |L| = \sqrt{\text{Tr} L^* L} \). We denote by

\[
\gamma(t) = \sup_a |g(a, s)| = \| \nabla_x u(\cdot, t) \|_{L^\infty(dx)}
\]

a quantity of some importance in the sequel. From (2.48), and the fact that the initial data for \( \partial_a X^i \) is \((\delta_{ij})_{j=1, \ldots, d}\), we have that

\[
\sup_a |\partial_a X^i(a, t)| \leq e^{\int_0^t \gamma(s) ds}.
\]

Because the initial data for \( \delta_h \left( \partial_a X^i \right) \) vanishes, we have, from Gronwall’s inequality and (2.52),

\[
|\delta_h \left( \partial_a X^i(a, t) \right)| \leq e^{\int_0^t \gamma(s) ds}\int_0^t |\delta_h g(a, s)| ds.
\]

We are interested in quantities in their Eulerian form. We consider the Hölder seminorm

\[
[\phi]_\alpha = \sup_{x \neq y} \frac{|\phi(x) - \phi(y)|}{|x - y|^\alpha}
\]

with \( 0 < \alpha < 1 \). We note that the back-to-labels maps are Lipschitz and

\[
\sup_{x \neq y} \frac{|A(x, t) - A(y, t)|}{|x - y|} \leq e^{\int_0^t \gamma(s) ds}
\]

because

\[
D_t A = 0
\]
with initial data $A(x, 0) = x$, and consequently
\[ D_t \nabla_x A = -(\nabla_x A)(\nabla_x u). \tag{2.57} \]

In view of (2.55) and the fact that Hölder seminorms behave nicely with respect to compositions with Lipschitz functions,

\[ [\phi \circ A]_\alpha \leq [\phi]_\alpha \lambda^\alpha \]

if $\lambda$ is a bound on the Lipschitz seminorm of $A$, it follows that

\[ [e^i(\cdot, t)]_\alpha \leq [\nabla X^i(\cdot, t)]_\alpha e^\alpha \int_0^t \gamma(s) ds \tag{2.58} \]

and, similarly

\[ [g(\cdot, t)]_\alpha \leq [\nabla x u(\cdot, t)]_\alpha e^\alpha \int_0^t \gamma(s) ds \tag{2.59} \]

From these considerations and from (2.53) it follows that

\[ [e^i(\cdot, t)]_\alpha \leq e^{(1+2\alpha) \int_0^t \gamma(s) ds} \int_0^t [\nabla x u(\cdot, s)]_\alpha ds \tag{2.60} \]

The reason for the $2\alpha$ “loss” in comparison to the Lagrangian estimate (2.53) is that we lose twice, once because the objective of the estimate is composed, and another because what we estimate it with is also composed with a Lipschitz function. In the same manner we also obtain

\[ \|\nabla_x e^i(\cdot, t)\|_{L^p(\mathbb{R}^d)} \leq e^{3\int_0^t \gamma(s) ds} \int_0^t \|\nabla_x \nabla_x u(\cdot, s)\|_{L^p(\mathbb{R}^d)} ds \tag{2.61} \]

We perform similar calculations for the matrix $Q(x, t, s)$, taking into account (2.21, 2.22, 2.40, 2.41). We deduce

\[ \|Q(\cdot, t, s)\|_{L^\infty(\mathbb{R}^d)} \leq Ce^{\int_s^t \gamma(z) dz}, \tag{2.62} \]

\[ |\delta_k q(a, t, s)| \leq Ce^{\int_s^t \gamma(z) dz} \int_s^t |\delta_k g(a, z)| dz, \tag{2.63} \]

and consequently

\[ [Q(\cdot, t, s)]_\alpha \leq Ce^{(1+2\alpha) \int_0^t \gamma(s) dz} \int_0^t [\nabla x u(\cdot, z)]_\alpha dz. \tag{2.64} \]

with $C = \sqrt{d}$ depending only on the dimension of space $d$. We are ready to state the estimates on $\tau$ in terms of $u$.

**Proposition 2.1.** Let $u$ be a divergence-free function belonging to $L^1(0, T; C^{1,\alpha}([\mathbb{R}^d]))$. Let $\sigma_0(a)$ be a Hölder continuous, $L^1(\mathbb{R}^d)$, positive symmetric matrix and let $\rho_0$ be a positive, Hölder continuous, $L^1(\mathbb{R}^d)$ function. Then $\tau = \sigma - \rho_\xi$, given in the expression (2.46), obeys

\[ \|\tau(\cdot, t)\|_{L^\infty(\mathbb{R}^d)} \leq C\|\rho_0\|_{L^\infty(\mathbb{R}^d)} \int_0^t \exp \left\{ -\frac{2t(s - t)}{R^2} + 2 \int_s^t \gamma(z) dz \right\} \gamma(s) ds \]

\[ + C\|\tau_0\|_{L^\infty(\mathbb{R}^d)} \exp \left\{ -\frac{2tR^2}{R^2} + 2 \int_0^t \gamma(s) ds \right\}, \tag{2.65} \]
Differentiating (2.38), we can write

\[ \tau(\cdot, t) \leq C[\rho_0]_0 \exp \int_0^t \gamma(z)dz \int_0^t \exp \left\{ -\frac{2\alpha(t-s)}{\mathcal{R}^2} + 2 \int_s^t \gamma(z)dz \right\} \gamma(s)ds \]

and

\[ + C\|\rho_0\|_{L^\infty(\mathbb{R}^d)} \int_0^t \exp \left\{ -\frac{2\alpha(t-s)}{\mathcal{R}^2} + (2 + 2\alpha) \int_s^t \gamma(z)dz \right\} \gamma(s) \int_s^t |\nabla_x u(\cdot, z)|_\alpha dz ds \]

\[ + C\|\rho_0\|_{L^\infty(\mathbb{R}^d)} \int_0^t \exp \left\{ -\frac{2\alpha(t-s)}{\mathcal{R}^2} + 2 \int_s^t \gamma(z)dz + 2\alpha \int_0^s \gamma(s)dz \right\} |\nabla_x u(\cdot, s)|_\alpha dz \]

\[ + C \left[ \tau_0 \alpha + \|\tau_0\|_{L^\infty(\mathbb{R}^d)} \int_0^t |\nabla_x u(\cdot, s)|_\alpha dt \right] \exp \left\{ -\frac{2\alpha t}{\mathcal{R}^2} + (2 + 2\alpha) \int_0^t \gamma(s)ds \right\} \]

(2.66)

with \( \gamma(t) = \sup_z |\nabla_x u(x, t)| \) and the constant \( C > 0 \) depending only on the dimension \( d \) of space.

**Proof.** The \( L^\infty \) estimate follows from (2.62) and

\[ \|e^t(\cdot, t)\|_{L^\infty(dx)} \leq e^{L_0^t \gamma(s)ds}, \]

(2.67)

which follows immediately from (2.38) and the fact that the initial data are of unit norm. The Hölder seminorm estimate follows from the algebra inequality

\[ [\phi]\alpha \leq [\phi]\alpha \|\psi\|_{L^\infty} + \|\phi\|_{L^\infty} |\psi|_\alpha, \]

(2.62) and (2.67) above, (2.60), (2.64), and the behavior of the Hölder continuous function under composition with \( A \) and (2.55). This concludes the proof of the proposition. \( \square \)

Before we conclude this section, we comment on the case of \( d = 2 \). In this case the velocity gradient matrix is

\[ \nabla_x u(x, t) = \begin{pmatrix} \lambda(x, t) & \mu(x, t) - \frac{\omega(x, t)}{2} \\ \mu(x, t) + \frac{\omega(x, t)}{2} & -\lambda(x, t) \end{pmatrix}, \]

(2.68)

where

\[ \lambda(x, t) = \frac{1}{2} \left( \partial_1 u^1(x, t) - \partial_2 u^2(x, t) \right), \]

(2.69)

\[ \mu(x, t) = \frac{1}{2} \left( \partial_1 u^2(x, t) + \partial_2 u^1(x, t) \right), \]

(2.70)

and

\[ \omega(x, t) = \partial_1 u^2(x, t) - \partial_2 u^1(x, t). \]

(2.71)

The symmetric part of the gradient (rate of strain matrix) is

\[ S(x, t) = \begin{pmatrix} \lambda(x, t) & \mu(x, t) \\ \mu(x, t) & -\lambda(x, t) \end{pmatrix}. \]

(2.72)

Differentiating (2.38), we can write

\[ D^2_t e^i = (|\nabla_x u|^2 + D_i (\nabla_x u)) e^i. \]

(2.73)
This equation holds in any number of dimensions, but because of the Hamilton-Cayley theorem, in two dimensions we have
\[
(\nabla_x u(x,t))^2 = \delta(x,t) I
\]  
(2.74)
with
\[
\delta(x,t) = \lambda^2(x,t) + \mu^2(x,t) - \frac{1}{4} \omega^2(x,t) = -\det(\nabla_x u)(x,t).
\]  
(2.75)
Therefore the Equation (2.73) becomes
\[
D_t^2 e^i = (\delta I + D_t(\nabla_x u)) e^i.
\]  
(2.76)
Similarly, we have
\[
d_s^2 q = - (\delta I + \partial_s (g(a,s))^{T}) q.
\]  
(2.77)
We can easily integrate (2.76, 2.77) when the velocity gradient is constant in space and time (so that \(D_t(\nabla x u) = 0\)). In this case \(\delta\) is a constant and its sign dictates the size of solutions \(e^i\) and \(Q\). Consequently, when \(\delta \leq 0\) or when \(\delta > 0\) but \(\sqrt{\delta} \leq \frac{\epsilon R}{2}\), then \(\tau\) is bounded, and when \(\delta > (\frac{\epsilon R}{2})^2\), then \(\tau\) may grow exponentially in time. Integrating (2.76, 2.77) when \(D_t(\nabla x u) = 0\) we see that if \(\delta \leq 0\) the functions \(e^i\) and \(Q\) are bounded in time, and when \(\delta > 0\) they grow like \(\exp \sqrt{\delta t}\) in time. If the initial distribution \(f_0\) is radially symmetric in \(m\) then the matrix \(\sigma_0(a)\) is a multiple of the identity matrix (a scalar matrix, or an isotropic stress),
\[
\sigma_0(A(x,t)) = c\rho_0(A(x,t)) I.
\]
If \(\delta\) is a positive constant, integrating (2.76) we have
\[
e^1(x,t) = \begin{pmatrix} \cosh(t\sqrt{\delta}) \\ \sinh(t\sqrt{\delta}) \end{pmatrix}, \quad e^2(x,t) = \begin{pmatrix} \sinh(t\sqrt{\delta}) \\ \cosh(t\sqrt{\delta}) \end{pmatrix}.
\]
In this case
\[
e^{11}(x,t) = e^{22}(x,t) = c\rho_0(A(x,t)) \cosh(2t\sqrt{\delta}),
\]
\[
e^{12}(x,t) = c\rho_0(A(x,t)) \sinh(2t\sqrt{\delta}).
\]
The case of a uniform gradient is used in rheology, but if the velocity gradient is square-integrable in space, then \(\delta\) has to have average zero in space because
\[
(\nabla_x u)^2_{ij} = \partial_k(u^i\partial_j u^k).
\]  
(2.78)
Therefore, for moderately varying \(\nabla x u\) we can expect \(\tau\) to be bounded.

3. Coupling to the Stokes system

The Stokes system is
\[
-\Delta_x u + \nabla_x p = k\text{div}_x \sigma
\]  
(3.1)
coupled with
\[
\text{div}_x u = 0.
\]  
(3.2)
The coefficient $k$ has units of inverse time, and represents the ratio between an energy (kT) per unit mass and kinematic viscosity. The fluid density is taken to be 1. We discuss the situation in which the fluid occupies all $\mathbb{R}^d$ and is at rest at infinity. The velocity and pressure decay in space and are obtained from $\sigma$ by classical singular integrals. We can equally consider the case in which the velocity and pressure are periodic in space. Let us denote by

$$R_j = \Lambda^{-1} \partial_j,$$

for $j = 1, \ldots, d$, in $\mathbb{R}^d$ the Riesz operators, where

$$\Lambda = (\Delta)^{\frac{1}{2}}$$

is the Zygmund operator. By modifying the pressure we may write

$$-\Delta_x u + \nabla_x \tilde{p} = k \text{div}_x (\tau),$$

and using (3.2) to solve for $\tilde{p}$,

$$\tilde{p} = -k R_m R_n (\tau^{mn}),$$

we deduce

$$u^i = k \Lambda^{-1} \left( R_i (\tau^{il}) + R_l R_m R_n (\tau^{mn}) \right)$$

and

$$\partial_j u^i = k R_j \left( R_i (\tau^{il}) + R_l R_m R_n (\tau^{mn}) \right),$$

which we abbreviate as

$$\nabla_x u = k \mathcal{R} \tau.$$  

The system formed by (2.39, 2.43) together with (3.7) is closed.

**3.1. Energetics.** The system formed by (2.1) coupled to (3.1) has a Liapunov functional. We take (2.1), multiply by $k (\log f + U)$, integrate, take (3.1) and multiply by $\frac{1}{2} k u$ and integrate, and add the two: we deduce

$$\frac{d}{dt} \int_{\mathbb{R}^d} \int_{\mathbb{R}^d} f \log f dm dx + \frac{1}{2} \int_{\mathbb{R}^d} Tr[\sigma(x, t)] dx$$

$$= -\int_{\mathbb{R}^d} |\nabla_x u(x, t)|^2 dx - k \epsilon \int_{\mathbb{R}^d} f |\nabla_m (\log f + U)|^2 dmdx.$$  

This is a general property of the class of Fokker-Planck equations in which we are interested, but in the case of the Oldroyd-B system we are considering here we can obtain the energetics directly from the Equation (2.29) by taking the trace, adding it to the Stokes Equation (3.1), multiplying by $\frac{1}{2} k u$ and integrating:

$$\partial_t \int_{\mathbb{R}^d} Tr \sigma dx + \frac{2}{k} \int_{\mathbb{R}^d} |\nabla_x u|^2 dx = -\frac{2 \epsilon}{R^2} \int_{\mathbb{R}^d} Tr \sigma dx + \frac{2 \epsilon}{R^2} \int_{\mathbb{R}^d} \rho_0 dx.$$  

Using the fact that $Tr \sigma$ together with $\rho$ bound the entries in $\tau$, we derive

$$\|\tau(t, \cdot)\|_{L^1(\mathbb{R}^d)} \leq C[\|\tau_0\|_{L^1(\mathbb{R}^d)} + \|\rho_0\|_{L^1(\mathbb{R}^d)}],$$

with a constant $C$ depending only on $d$. 

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3.2. Gradient bounds. We are going to investigate bounds on the time integral of the maximum gradient of velocity,

$$\int_0^T \| \nabla_x u(\cdot, t) \|_{L^\infty(dx)} dt.$$ 

In order to do so we recall a classical fact ([29], Chapter 3, Section 5), namely that operators given at Fourier level by multiplication by functions that are homogeneous of degree zero and are smooth on the unit sphere can be represented as sums of multiples of the identity and singular integral operators of a classical type.

**Proposition 3.1.** There exists a constant $C$ depending only on $d$ and $\alpha$ such that,

$$\| R \tilde{\sigma} \|_{L^\infty(dx)} \leq C \| \tilde{\sigma} \|_{L^\infty(dx)} \left\{ 1 + \log \left[ 1 + \frac{\| \tilde{\sigma} \|_{L^1(\mathbb{R}^d)} [\tilde{\sigma}]_\alpha^{\frac{d}{d-1}}}{{\| \tilde{\sigma} \|_{L^\infty(dx)}}} \right] \right\} ,$$

where

$$[\tilde{\sigma}]_\alpha = \max_{m,n} \sup_{x \neq y} \frac{\tilde{\sigma}^{mn}(x) - \tilde{\sigma}^{mn}(y)}{|x - y|^{\alpha}}$$

(3.14)

and $0 < \alpha < 1$. There exists a constant $C$ depending only on $d$ and $\alpha$ such that

$$[R \tilde{\sigma}]_\alpha \leq C [\tilde{\sigma}]_\alpha.$$  

(3.15)

**Proof.** In view of the fact that $R$ is made from operators of the form $R_j R_k$ and $R_j R_k R_l$, we have the representation

$$(R \tilde{\sigma})_{ij}(x) = C_{ijmn} \tilde{\sigma}^{mn}(x) + P.V. \int_{\mathbb{R}^d} k_{ijmn}(y) \tilde{\sigma}^{mn}(x - y) dy,$$

where $C_{ijmn}$ are constants and each $k_{ijmn}$ is a smooth function (actually a harmonic polynomial of degree no more than three) of $\hat{y} = y/|y|$ with mean zero on the unit sphere $\mathbb{S}^{d-1}$. We need to prove the inequality for each entry, and we reduce the proof to the scalar case. The proof then follows along very classical lines. We break the integral

$$K \tilde{\sigma}(x) = P.V. \int_{\mathbb{R}^d} \frac{k(\hat{y})}{|y|^d} \tilde{\sigma}(x - y) dy$$

into three pieces. We choose two numbers $m < M$ and write

$$I_1(x) = P.V. \int_{|y| \leq m} \frac{k(\hat{y})}{|y|^d} \tilde{\sigma}(x - y) dy,$$

$$I_2(x) = \int_{m \leq |y| \leq M} \frac{k(\hat{y})}{|y|^d} \tilde{\sigma}(x - y) dy,$$

and

$$I_3(x) = \int_{|y| \geq M} \frac{k(\hat{y})}{|y|^d} \tilde{\sigma}(x - y) dy.$$
We have, with the cancellation property of $k(\hat{y})$,

$$|I_1(x)| \leq C m^\alpha \|\tilde{\sigma}\|_\alpha,$$

and then clearly

$$|I_2(x)| \leq C \log \left( \frac{M}{m} \right) \|\tilde{\sigma}\|_{L^\infty},$$

and

$$|I_3(x)| \leq C M^{-d} \|\tilde{\sigma}\|_{L^1(dx)}.$$ 

Now we choose the two length scales,

$$M = \left[ \frac{\|\tilde{\sigma}\|_{L^1(\mathbb{R}^d)}}{\|\tilde{\sigma}\|_{L^\infty(dx)}} \right]^{\frac{1}{d}},$$

and

$$m = \left[ \frac{\|\tilde{\sigma}\|_{L^\infty(dx)}}{|\tilde{\sigma}|_\alpha} \right]^{\frac{1}{d}}.$$ 

If $M \geq m$ we use both in the inequalities above and we obtain

$$|K\tilde{\sigma}(x)| \leq C \|\tilde{\sigma}\|_{L^\infty(dx)} \left[ 2 + \log \left( \frac{M}{m} \right) \right].$$

If $M \leq m$ then we forgo the division above and split the integrals into only two pieces, up to $M$ and from $M$. We obtain then

$$|K\tilde{\sigma}(x)| \leq 2C \|\tilde{\sigma}\|_{L^\infty(dx)}.$$ 

This ends the proof of (3.13). The inequality (3.15) is a classical inequality for singular integrals of the type above. This ends the proof of the proposition.

4. Small data

The approach we choose to control $|\tau|_\alpha$ is via the explicit formula (2.46). The equation obeyed by $\tau$ can be derived directly from (2.29):

$$D_t \tau = -\frac{2\epsilon}{R^2} \tau + (\nabla_x u)\tau + \tau (\nabla_x \bar{u})^T + 2\rho S.$$ (4.1)

Because $2S = kR\tau + k(R\tau)^T$ depends linearly on $\tau$, it is clear that strong enough damping will prevail and send the solution to zero. It is clear however that

$$\frac{\epsilon}{kR^2} \geq C \|\rho_0\|_{L^\infty(\mathbb{R}^d)}$$ (4.2)

will have to be part of the requirement for this stability result. The system of (2.1), with (2.2) and (2.3) coupled with (3.1), is equivalent to the system formed by (4.1), the Stokes system (3.5) and (2.26). We use Proposition (2.1). We first prove a good local existence result. Let us denote

$$M_1 := \|\rho_0\|_{L^1(\mathbb{R}^d)} + \|\tau_0\|_{L^1(\mathbb{R}^d)}.$$ (4.3)
\[ M_\infty := \| \rho_0 \|_{L^\infty(\mathbb{R}^d)} + \| \tau_0 \|_{L^\infty(\mathbb{R}^d)}, \]  

and

\[ M_\alpha := [\rho_0]_\alpha + [\tau_0]_\alpha. \]  

Let us introduce the Deborah number

\[ D := \frac{k}{\kappa_0} = \frac{k R^2}{\epsilon}, \]  

the nondimensional ratio of the particle time scale \( R^2 \) to the viscous fluid’s response time to the stress added by the particles, \( k^{-1} \). We denoted half the damping rate by

\[ \kappa_0 = \frac{\epsilon}{R^2}. \]  

**Theorem 4.1.** Let \((\tau_0, \rho_0) \in L^1(\mathbb{R}^d)^2 \cap C^\alpha(\mathbb{R}^d)^2\). There exist constants \( \varepsilon > 0 \), \( \Gamma \geq 2 \) and there exists a time \( T_0 > 0 \) and a weak solution of \((2.26, 3.7, 4.1)\) \((\rho, \tau) \in C([0, T_0], W^{-1}(\mathbb{R}^d)^2) \cap L^\infty([0, T_0]; (C^\alpha(\mathbb{R}^d) \cap L^1(\mathbb{R}^d))^2), 1 < p < \infty,\) satisfying the equations in the weak sense, and such that

\[ Dk_0 T_0 M_\infty \left( 1 + \log \left( 1 + M_\alpha \frac{\tau}{\tau} M_\infty \right) \right) \geq \varepsilon \]  

and

\[ \| \rho(\cdot, t) \|_{L^1(\mathbb{R}^d)} + \| \tau(\cdot, t) \|_{L^1(\mathbb{R}^d)} \leq \Gamma M_1, \]  

and

\[ \| \rho(\cdot, t) \|_{L^\infty(\mathbb{R}^d)} + \| \tau(\cdot, t) \|_{L^\infty(\mathbb{R}^d)} \leq \Gamma M_\infty, \]  

and

\[ [\rho(\cdot, t)]_\alpha + [\tau(\cdot, t)]_\alpha \leq \Gamma M_\alpha \]  

hold on \([0, T_0]\).

**Proof.** We consider the system \((4.1, 3.7, 2.26)\) with initial data that have been regularized by convolution with a standard mollifier \( \psi_\delta \). We obtain uniform bounds for quantities of interest and pass to the limit \( \delta \to 0 \), removing the mollifier. In order to simplify the exposition, we will denote the solutions \( \tau(\delta), \rho(\delta) \) by \( \tau, \rho \) as if they did not depend on \( \delta \). The bounds and constants below will indeed be \( \delta \)-independent. We use the following notations

\[ \pi(t) := (Dk_0)^{-1} \sup_{0 \leq s \leq t} \| \nabla_x u(\cdot, s) \|_{L^\infty(\mathbb{R}^d)}, \]  

\[ b(t) := (Dk_0)^{-1} \sup_{0 \leq s \leq t} [\nabla_x u(\cdot, s)]_\alpha, \]  

\[ n(t) := \sup_{0 \leq s \leq t} \{ \| \tau(\cdot, s) \|_{L^\infty(\mathbb{R}^d)} + \| \rho(\cdot, s) \|_{L^\infty(\mathbb{R}^d)} \}, \]  

and

\[ m(t) := \sup_{0 \leq s \leq t} \{ [\tau(\cdot, s)]_\alpha + [\rho(\cdot, s)]_\alpha \}, \]  

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We prove first uniform bounds, and then discuss their use. We use (3.12) to obtain
\[
\sup_{0 \leq s \leq t} \| \tau(\cdot, s) \|_{L^1} + \| \rho(\cdot, s) \|_{L^1} \leq CM_1,
\] (4.16)
which is based on the fundamental dissipative property of the system. We note that the function \( (x, a, b) \mapsto x \log (1 + \frac{ab}{x}) \) is increasing in each of its arguments \( x > 0, a > 0, b > 0 \). Therefore, we may use upper bounds in the right-hand side of (3.13). Using (3.13) we deduce
\[
\gamma(t) \leq C n(t) \{ 1 + \log \left( 1 + \frac{m(t)^{\frac{1}{2-\alpha}} M_1^{\frac{1}{2-\alpha}}}{n(t)} \right) \},
\] (4.17)
and using (3.15)
\[
b(t) \leq C m(t).
\] (4.18)
We proceed to estimate in (2.65):
\[
n(t) \leq C M_\infty [1 + D_{K_0} t \tau(t)] e^{C D_{K_0} t \tau(t)},
\] (4.19)
and (2.66):
\[
m(t) \leq C M_\alpha [1 + D_{K_0} t \tau(t)] e^{C D_{K_0} t \tau(t)} + C (M_\infty D_{K_0} t)[1 + D_{K_0} t \tau(t)] e^{C D_{K_0} t \tau(t)} m(t).
\] (4.20)
The term \( D_{K_0} t \tau(t) \) can be bound from (4.17):
\[
D_{K_0} t \tau(t) \leq (CD_{K_0} t)n(t) \left\{ 1 + \log \left( 1 + \frac{m(t)^{\frac{1}{2-\alpha}} M_1^{\frac{1}{2-\alpha}}}{n(t)} \right) \right\}.
\] (4.21)
The structure of the inequalities is the following: if we denote the group
\[
x(t) := D_{K_0} t \tau(t),
\] (4.22)
we have from (4.19)
\[
n(t) \leq C M_\infty (1 + x(t)) e^{C x(t)},
\] (4.23)
and therefore from (4.21) we deduce
\[
x(t) \leq (CD_{K_0} t) M_\infty (1 + x(t)) e^{C x(t)} \left( 1 + \log \left( 1 + \frac{1}{C (1 + x(t)) e^{C x(t)}} \right) \right) \times \left\{ 1 + \log \left( 1 + m(t)^{\frac{1}{2-\alpha}} M_1^{\frac{1}{2-\alpha}} M_\infty^{-1} \right) \right\},
\] (4.24)
so
\[
F(x(t)) \leq C (M_\infty D_{K_0} t) \left\{ 1 + \log \left( 1 + m(t)^{\frac{1}{2-\alpha}} M_1^{\frac{1}{2-\alpha}} M_\infty^{-1} \right) \right\},
\] (4.25)
where
\[
F(x) = x e^{-C x} (1 + x)^{-1} \left( 1 + \log \left( 1 + \frac{1}{C (1 + x) e^{C x}} \right) \right)^{-1}.
\] (4.26)
Let
\[ a(t) := M_\infty D\kappa_0 t. \] (4.27)

Now (4.20) can be written as
\[ m(t) \leq C(1 + (x(t)))e^{C_2x(t)}[M_\alpha + a(t)m(t)] \] (4.28)
and
\[ F(x) \leq C_0 \left\{ 1 + \log \left( 1 + m(t) \frac{1}{M_\alpha} M_1 \frac{M_\alpha}{M_\infty} M_\infty^{-1} \right) \right\}. \] (4.29)

We fix \( C_0 \) a large enough absolute constant, larger than the \( C \) we encountered so far. Without loss of generality we may assume \( 1 + x \leq e^{C_0 x} \) and
\[ 1 + \log \left( 1 + \frac{1}{C(1 + x)e^{C_0 x}} \right) \leq 1 + \log(1 + C^{-1}) \leq C_0, \]
so that
\[ F(x) \geq xe^{-2C_0 x} \frac{1}{C_0} \]
and therefore, as long as \( x(t) \leq 1 \), we have from (4.29)
\[ x(t) \leq C_0^2 e^{2C_0 t} a(t) \left\{ 1 + \log \left( 1 + m(t) \frac{1}{M_\alpha} M_1 \frac{M_\alpha}{M_\infty} M_\infty^{-1} \right) \right\}. \] (4.30)
and from (4.28)
\[ \frac{m(t)}{M_\alpha} \leq C_0 e^{2C_0 x(t)} \left[ 1 + a(t) \frac{m(t)}{M_\alpha} \right]. \] (4.31)
Using \( x(t) \leq 1 \) and requiring
\[ a(t) \leq \frac{1}{2} C_0^{-1} e^{-2C_0} \] (4.32)
we deduce from (4.31) that, as long as \( x(t) \leq 1 \), it follows that
\[ \frac{m(t)}{M_\alpha} \leq 2C_0 e^{2C_0}. \] (4.33)

Returning to \( x(t) \), as long as \( x(t) \leq 1 \), we have from (4.30) and (4.33) that
\[ x(t) \leq C_0^2 e^{2C_0} a(t) \left\{ 1 + \log \left( 1 + (2C_0 e^{2C_0} M_\alpha) \frac{1}{M_1} M_1 \frac{M_\alpha}{M_\infty} M_\infty^{-1} \right) \right\}. \] (4.34)
We require therefore also
\[ a(t) \leq \frac{1}{2} \left\{ C_0^2 e^{2C_0} \left\{ 1 + \log \left( 1 + (2C_0 e^{2C_0} M_\alpha) \frac{1}{M_1} M_1 \frac{M_\alpha}{M_\infty} M_\infty^{-1} \right) \right\} \right\}^{-1} \] (4.35)
and thus, as long as \( x(t) \leq 1 \), it follows that actually \( x \) is bounded by \( x(t) \leq \frac{3}{4} \). The initial datum for \( x(t) \) is zero, and \( a(t) \) is explicitly proportional to \( t \) and starts from zero. Therefore, because of continuity, if the time \( t \) is taken small enough so that \( a(t) \) satisfies the requirements (4.32, 4.35) it follows that \( x(t) < \frac{3}{4} \) on that interval;
we can reason by contradiction that the first time it would reach \( \frac{3}{4} \), it would have to be not larger than \( \frac{1}{2} \), which is absurd. Consequently, it follows from (4.33) that \( m(t) \leq 2C_0 e^{2C_0 M_\alpha} \) and from (4.23) that \( n(t) \leq 2C_0 e^{2C_0 M_\infty} \) on the time interval, providing the required short time uniform bound. The rest of the proof of the theorem is based on the uniform bound. It is well known that regular solutions exist as long as the \( L^\infty \) bounds persist (see [4, 18], also next section). Therefore classical solutions with regularized initial data exist for this interval of time. Indeed, by regularizing the initial data we do not hamper the bounds \( M_1, M_\infty, M_\alpha \). The most straightforward proof of local existence for smooth data yields a very short time of existence, but the a priori bounds permit the solution to be extended as long as (4.32) and (4.35) hold.

Then we remove the regularization of the initial data, and pass to limit in the equation in distributional sense. We conclude that the limit obeys the equation because the quadratic nonlinearity is strongly continuous in \( L^2 \) (for the right-hand side) and in \( W^{-1,2} \) for the advective derivative. Further details are left for the interested reader.

We will pursue now the stability issue. Let us start by assuming that

\[
\frac{1}{\kappa_0} \| \nabla_x u(\cdot, t) \|_{L^\infty(\mathbb{R}^d)} \leq Ge^{-\kappa_0 t}
\]

(4.36)
on the time interval \( t \in [0, T] \). We require that

\[
G \leq \frac{1}{4}.
\]

(4.37)

Similarly, we assume

\[
\frac{1}{\kappa_0} [\nabla_x u(\cdot, t)]_\alpha \leq He^{-\kappa_0 t}
\]

(4.38)
on the same time interval. We will show that we can choose \( G, H \) so that these assumptions are invariant in time: once initiated they cannot break down. Before we embark on the calculation, we should point out the simple strategy: we obtain from (4.36), (4.38) a priori bounds on \( \tau \). We then deduce a posteriori bounds for \( \nabla_x u \) that are strictly better than the assumptions (4.36, 4.38). Because we work in a class in which the equation has unique solutions, if the initial data for \( \tau \) are small we deduce that, if initiated, the bounds continue indefinitely.

We start by estimating the terms in (2.66) using (4.36) and (4.38). We will use

\[
\int_0^t \gamma(z) dz \leq G,
\]

(4.39)

which follows from (4.36). We have four terms, and in all of them, as well as in (2.65), we bound

\[
\int_0^t \exp \left\{ -2\kappa_0 (t - s) + 2 \int_s^t \gamma(z) \right\} \gamma(s) d\gamma ds
\]

\[
\leq e^{2G} \exp \left\{ -2\kappa_0 t \right\} \int_0^t \exp \left\{ 2\kappa_0 s \right\} \kappa_0 G \exp \left\{ -\kappa_0 s \right\} ds \leq Ge^{-\kappa_0 t} e^{2G}
\]

\[
\leq \sqrt{\tau} Ge^{-\kappa_0 t},
\]

(4.40)

using (4.39) and the fact that \( 4G \leq 1 \). From this calculation and (2.65) it follows that

\[
\| \tau(\cdot, t) \|_{L^\infty(\mathbb{R}^d)} \leq Ce^{-\kappa_0 t} \left[ \| \rho_0 \|_{L^\infty(\mathbb{R}^d)} + \| \tau_0 \|_{L^\infty(\mathbb{R}^d)} \right] = Ce^{-\kappa_0 t} M_\infty,
\]

(4.41)
with $C$ depending on $d$ alone. Using (2.62) and (2.67) in (2.46) and integrating we obtain in the same way

$$
\|\tau(\cdot, t)\|_{L^1(\mathbb{R}^d)} \leq C e^{-\kappa_0 t} \left[ \|\rho_0\|_{L^1(\mathbb{R}^d)} + \|\tau_0\|_{L^1(\mathbb{R}^d)} \right] = C e^{-\kappa_0 t} M_1. \tag{4.42}
$$

We write the inequality (2.66) as

$$
[\tau(\cdot, t)]_\alpha \leq I + II + III + IV,
$$

and we start with

$$
I \leq C[\rho_0]_\alpha e^{\alpha G} \left( \frac{1}{2} \right) e^{-\kappa_0 t} = C[\rho_0]_\alpha e^{-\kappa_0 t},
$$

in which we used (4.39) and (4.40) above. This illustrates our use of the generic constant $C$: we used $\alpha G \leq \frac{1}{2}$. For $II$ we obtain

$$
II \leq C[\rho_0]_\alpha \|\nabla_x u(\cdot, s)\|_\alpha ds
\leq C[\rho_0]_\alpha \|\nabla_x u(\cdot, s)\|_\alpha ds.
$$

Finally, we use again (4.39) and (2 + 2\alpha)G \leq 1. Collecting terms, we have proved that (4.36, 4.37) and (4.38) imply that

$$
[\tau(\cdot, t)]_\alpha
\leq C e^{-\kappa_0 t} \left\{ [\rho_0]_\alpha + [\tau_0]_\alpha + H \|\rho_0\|_{L^\infty(\mathbb{R}^d)} + \|\nabla_x u(\cdot, s)\|_\alpha ds \right\}, \tag{4.48}
$$

Now from (3.9), (3.15) and a Gronwall-like inequality we derive

$$
[\tau(\cdot, t)]_\alpha
\leq C \left[ [\rho_0]_\alpha + [\tau_0]_\alpha + H \|\rho_0\|_{L^\infty(\mathbb{R}^d)} \right] \exp\left\{ -t\kappa_0 + \frac{Ck\|\rho_0\|_{L^\infty(\mathbb{R}^d)} + [\tau_0]_\alpha}{\kappa_0} \right\}.
$$

Indeed, (4.48) together with (3.9, 3.15) imply an inequality of the type

$$
y(t) \leq C e^{-t\kappa_0} \left[ a + b \int_0^t y(s) ds \right] = F(t)
$$
for the positive quantity \( y(t) = [\tau(\cdot, t)]_\alpha \), with

\[
a = [\rho_0]_\alpha + [\tau_0]_\alpha + H\|\rho_0\|_{L^\infty(\mathbb{R}^d)}
\]

and

\[
b = k\|\rho_0\|_{L^\infty(\mathbb{R}^d)} + \|\tau_0\|_{L^\infty(\mathbb{R}^d)}.
\]

Differentiating the right hand side \( F(t) \) and using the fact that \( y(t) \leq F(t) \) we deduce

\[
\frac{dF}{dt} \leq -\kappa_0 F + Cbe^{-t\kappa_0} F,
\]

and dividing by \( F \) and integrating we obtain

\[
\log \left( \frac{F(t)}{Ca} \right) \leq -\kappa_0 t + Cb \kappa_0.
\]

Substituting back, we deduce (4.49). Let us consider now the right-hand side of (3.13) with \( \tilde{\sigma} = \tau \). We recall that (4.41) is

\[
\|\tau(\cdot, t)\|_{L^\infty(\mathbb{R}^d)} \leq Ce^{-\kappa_0 t} M_\infty,
\]

(4.50)

(4.42) is

\[
\|\tau(\cdot, t)\|_{L^1(\mathbb{R}^d)} \leq Ce^{-\kappa_0 t} M_1,
\]

(4.51)

and (4.49) becomes

\[
[\tau(\cdot, t)]_\alpha \leq C[M_\alpha + H\|\rho_0\|_{L^\infty(\mathbb{R}^d)}]e^{-\kappa_0 t} \exp \{CDM_\infty \}.
\]

Assembling the right hand side of (3.13) we deduce

\[
\|\mathcal{R}\tau\|_{L^\infty(\mathbb{R}^d)} \leq CM_\infty e^{-\kappa_0 t} \left\{ 1 + \log \left[ 1 + \frac{M_1}{M_\infty} \frac{[M_\alpha + H M_\infty]}{\kappa_0} \right] + DM_\infty \right\}.
\]

(4.53)

Note that the term \( e^{-\kappa_0 t} \) canceled in the expression for the argument of the logarithm. Now we use (3.9) to obtain a posteriori estimates. From (4.53) we obtain

\[
\frac{1}{\kappa_0}\|\nabla_x u(\cdot, t)\|_{L^\infty(\mathbb{R}^d)} \leq CDM_\infty e^{-\kappa_0 t} \left\{ 1 + \log \left[ 1 + \frac{M_1}{M_\infty} \frac{[M_\alpha + H M_\infty] \frac{d}{\kappa_0}}{M_\infty} \right] + DM_\infty \right\}.
\]

(4.54)

Also, from (4.52) and (3.9) we have

\[
\frac{1}{\kappa_0}[\nabla_x u(\cdot, t)]_\alpha \leq CD[M_\alpha + H M_\infty] \exp \{-\kappa_0 t + CDM_\infty \}.
\]

(4.55)

Let us choose \( G = 1/4 \) and

\[
H = \frac{M_\alpha}{M_\infty}.
\]

(4.56)
Then (4.55) becomes
\[
\frac{1}{\kappa_0}[(\nabla_x u)(\cdot,t)]_\alpha \leq CDM_\alpha \exp \{-\kappa_0 t + CDM_\infty\}. \tag{4.57}
\]
The bound (4.52) becomes
\[
[\tau(\cdot,t)]_\alpha \leq CMe^{-\kappa_0 t} \exp \{CDM_\infty\}. \tag{4.58}
\]
The upper bound (4.54) implies the bound
\[
\frac{1}{\kappa_0} \| (\nabla_x u)(\cdot,t) \|_{L^\infty(\mathbb{R}^d)} \leq B_1 Ge^{-\kappa_0 t}, \tag{4.59}
\]
where
\[
B_1 = B_0 + CD^2 M_\infty^2 \tag{4.60}
\]
and \( B_0 = B_0(M_1, M_\alpha, M_\infty, D) \) is given by
\[
B_0(M_1, M_\alpha, M_\infty, D) = CDM_\infty \left\{ 1 + \log \left[ 1 + M_\infty^{-1} M_\alpha^d M_1^{d-1} \right] \right\}, \tag{4.61}
\]
with appropriate constants \( C \) depending on \( d \) only. The functions \( B_0(M_1, M_\alpha, M_\infty, D) \) and \( B_1(M_1, M_\alpha, M_\infty, D) \) are continuous in their arguments \( M_1 > 0, M_\alpha > 0, M_\infty > 0, D > 0 \), and hence locally bounded, and vanish when \( M_\infty \to 0 \) while the rest are held fixed, or when \( D \to 0 \) and the rest are fixed.

The bound (4.57) implies the bound
\[
\frac{1}{\kappa_0} [(\nabla_x u)(\cdot,t)]_\alpha \leq B_2 He^{-\kappa_0 t}, \tag{4.62}
\]
where the function \( B_2 = B_2(DM_\infty) \) is given by
\[
B_2 = CDM_\infty \exp \{CDM_\infty\} \tag{4.63}
\]
with an appropriate constant \( C \). We have proved therefore

**Proposition 4.2.** Let \( \tau(x,t) \) be the solution of (4.1) given in (2.46) where \( \rho(x,t) \) solves (2.26) and assume that \( u \) obeys (4.36), (4.38) on the time interval \([0,T]\). Then \( \tau \) and \( u \) obey the bounds (4.50, 4.51, 4.58, 4.59, 4.62) on the same time interval, where the constants \( M_1, M_\infty, M_\alpha \) are given in (4.3, 4.4, 4.5), the constants \( B_0, B_1, B_2 \) are given in (4.60, 4.61, 4.63) and the constants \( C \) are independent of solutions, fixed, and depend only on \( d \).

We combine now this proposition with the good local existence result, Theorem 4.1. We start with initial data that are small. To fix the notation, we denote by \( C_1 \) an upper bound for the constants \( C \) appearing in Proposition 4.2 above. We can arrange the constants so that, from
\[
[\tau(\cdot,t)]_\alpha \leq \Gamma M_\alpha, \tag{4.64}
\]
and
\[
\| \tau(\cdot,t) \|_{L^1(\mathbb{R}^d)} \leq \Gamma M_1, \tag{4.65}
\]
and
\[
\| \tau(\cdot,t) \|_{L^\infty(\mathbb{R}^d)} \leq \Gamma M_\infty, \tag{4.66}
\]
it follows that
\[ \frac{1}{\kappa_0} \| \nabla_x u(\cdot, t) \|_{L^\infty(\mathbb{R}^d)} \leq \frac{G}{8C_1} \] (4.67)
and
\[ \frac{1}{\kappa_0} [\nabla_x u(\cdot, t)]_\alpha \leq \frac{H}{8C_1} \] (4.68)
because of the constitutive relations between \( u \) and \( \tau \). Note that the size of \( \tau \) alone determines that of \( u \). Indeed, in view of (3.13)
\[ \frac{1}{\kappa_0} \| \nabla_x u(\cdot, t) \|_{L^\infty(\mathbb{R}^d)} \leq CB_0(M_1, M_\alpha, M_\infty, D), \] and in view of (3.15)
\[ \frac{1}{\kappa_0} [\nabla_x u(\cdot, t)]_\alpha \leq CD_\Gamma M_\alpha = CD_\Gamma M_\infty H. \]
So, the condition
\[ B_0(M_1, M_\alpha, M_\infty, D) \leq \frac{1}{32CC_1} \] (4.69)
implies (4.67), and the condition
\[ D\Gamma M_\infty \leq \frac{1}{8CC_1} \] (4.70)
implies (4.68). Then, we automatically see that (4.36) and (4.38) hold on the interval of time \([0, T_0]\). In fact, if \( e^{\kappa_0 T_0} \leq 2C_1 \) we deduce that
\[ \frac{1}{\kappa_0} \| \nabla_x u(\cdot, t) \|_{L^\infty(\mathbb{R}^d)} < \frac{1}{2} Ge^{-\kappa_0 t} \] (4.71)
and
\[ \frac{1}{\kappa_0} [\nabla_x u(\cdot, t)]_\alpha < \frac{1}{2} He^{-\kappa_0 t} \] (4.72)
hold on \([0, T_0]\). On the other hand, as long as (4.36, 4.38) hold, the decay estimates (4.50, 4.51, 4.58) imply that
\[ \| \tau(\cdot, t) \|_{L^\infty(\mathbb{R}^d)} \leq \frac{2}{3} M_\infty, \] (4.73)
\[ \| \tau(\cdot, t) \|_{L^1(\mathbb{R}^d)} \leq \frac{2}{3} M_1, \] (4.74)
and
\[ [\tau(\cdot, t)]_\alpha \leq \frac{2}{3} M_\alpha, \] (4.75)
provided \( e^{\kappa_0 t} \geq \frac{3}{2} C_1 e^{C_1 D M_\infty} \). The lower bound (4.8) on the life-span of the solution is
\[ \kappa_0 T_0 B_0(M_1, M_\alpha, M_\infty, D) \geq \varepsilon, \] (4.76)
and it shows that if
\[ B_0(M_1, M_\alpha, M_\infty, D) \left[ \log \left( \frac{3C_1}{2} \right) + C_1 DM_\infty \right] \leq \varepsilon \]
then the solution exists for a period of time that is long enough so that \( \tau \) is reduced from its original size (4.73, 4.74, 4.75) at the end \( T_0 \) of the interval. Requiring that
\[ C_1 DM_\infty \leq \log \left( \frac{4}{3} \right) \]
allows to arrange for \( T_0 \) so that
\[ \frac{3}{2}C_1 e^{C_1 DM_\infty} \leq e^{\varepsilon T_0} \leq 2C_1. \]
Therefore, if the initial data satisfy (4.69, 4.70, 4.77, 4.78) then the assumptions of Proposition 4.2 are satisfied with \( T = T_0 \). We can repeat now argument taking the same time step. Therefore we find that the solution, if unique, exists for all time and obeys the bounds (4.64, 4.65, 4.66). In order to prove decay it is useful to consider first the case of smooth initial data. The previous considerations of course apply and so the solutions are global and obey the bounds (4.64, 4.65, 4.66) for all time, and therefore obey (4.67, 4.68) for all time, and consequently (4.71, 4.72) for \( t \leq T_0 \) with \( T_0 \) fixed above. Now we consider the function
\[ \phi(T) := \max \left\{ \sup_{0 \leq t \leq T} e^{\kappa t} \frac{1}{G_{\kappa_0}} \| \nabla_x u(\cdot, t) \|_{L^\infty(\mathbb{R}^d)} ; \sup_{0 \leq t \leq T} e^{\kappa t} \frac{1}{H_{\kappa_0}} \| \nabla_x u(\cdot, t) \|_{\alpha} \right\}. \]
(4.79)
This is a continuous nondecreasing function of \( T \in \mathbb{R}_+ \) (it is here that we use the fact that we mollified the initial data). Also, (4.71, 4.72) imply that \( \phi(T) \leq \frac{1}{2} \), for \( T \leq T_0 \). Now we claim that
\[ \sup_{T \geq 0} \phi(T) \leq \frac{3}{4}. \]
(4.80)
Indeed, assume not. Then there would exist a first time \( T_1 \) so that \( \phi(T_1) = \frac{3}{4} \). In particular, because \( \phi(T_1) \leq 1 \), Proposition 4.2 applies on \([0, T_1] \), and hence (4.59, 4.62) hold on \([0, T_1] \). But these inequalities imply \( \phi(T_1) \leq \frac{1}{2} \), which is absurd. Thus, (4.80) holds for all time. Then, using again Proposition 4.2, (4.50, 4.51, 4.58) hold for all \( t \), and these are uniform decay bounds. The constants depend only on the norms indicated, so by removing the mollifiers we see that the upper bounds are inherited by the weak solutions in \( C^\alpha \cap L^1 \).

We note that the smallness conditions (4.69, 4.70, 4.77, 4.78) all follow from a single condition:
\[ B_0(M_1, M_\alpha, M_\infty, D) \leq \varepsilon_1. \]
(4.81)

THEOREM 4.3. Let the initial data \( \tau_0, \rho_0 \) satisfy \( (\tau_0, \rho_0) \in (L^1(\mathbb{R}^d))^2 \cap (C^\alpha(\mathbb{R}^d))^2 \)
\( 0 < \alpha < 1 \). Let \( M_1, M_\infty, \) and \( M_\alpha \) defined in (4.3, 4.4, 4.5) denote the size of the initial data. There exists a constant \( \varepsilon_1 \) such that, if
\[ DM_\infty \left\{ 1 + \log \left( 1 + M_\infty^{-1} M_\alpha^{\frac{\alpha}{\alpha-1}} M_1^{\frac{\alpha}{\alpha-1}} \right) \right\} \leq \varepsilon_1, \]
then there exists a unique global weak solution
\[ (\tau, \rho) \in L^\infty([0, \infty)), (L^1(\mathbb{R}^d) \cap C^\alpha(\mathbb{R}^d))^2) \cap C([0, \infty), (W^{-1,p}(\mathbb{R}^d))^2), \]
The norm $\|\tau\|_{L^p(\mathbb{R}^d)}$ decay exponentially:

\[
\|\tau(\cdot,t)\|_{L^p(\mathbb{R}^d)} \leq Ke^{-\kappa t},
\]

for $1 \leq p < \infty$ with explicit constants $K$ that depend on $d$, $D$, $M_\infty$, $M_\alpha$, $p$ and $\alpha > 0$ alone. The norms of $\rho$ remain bounded:

\[
\|\rho(\cdot,t)\|_{L^p(\mathbb{R}^d)} = \|\rho_0\|_{L^p(\mathbb{R}^d)}
\]

and

\[
\|\rho(\cdot,t)\|_{L^\infty(\mathbb{R}^d)} \leq K.
\]

**Remark 4.4.** The condition (4.82) is satisfied if the Deborah number is arbitrary, the initial $\tau_0$ is sufficiently small in $C^\alpha(\mathbb{R}^d)$ and the initial $\rho_0$ is small in $C^\alpha$. Nontrivial initial data with small $L^\infty$ norm and large spatial derivatives are allowed as well because of the logarithmic dependence on $M_\alpha$. The conditions are also satisfied if the initial data are of order one but the Deborah number is small. If the initial data are smoother, the smoothness propagates.

**Remark 4.5.** The proof of uniqueness, given below, uses Lagrangian transformations. Unlike the case of Euler equations, a proof in Eulerian coordinates seems difficult to obtain.

**Proof.** The previous argument is complete, modulo the uniqueness of the solution. We present here the proof of uniqueness. We consider the equation

\[
\frac{d}{dt} X = F[X],
\]

with $X(a,0) = a$, where $X = X(a,t)$ is viewed as an element of $\mathcal{X} := C([0,T], C^1(\mathbb{R}^d))$. We fix $T$. The function $F$ depends on the whole path $X$ not only on the value of $X$ at some point, and is obtained as follows. First we construct, using the initial data $\rho_0 \in L^1(\mathbb{R}^d) \cap C^\alpha(\mathbb{R}^d)$, $\sigma_0(a) = \rho_0(a)I + \tau_0(a) \in (L^1(\mathbb{R}^d) \cap C^\alpha(\mathbb{R}^d))^d$, and using (2.42), the map $X \mapsto \sigma[X]$ given by

\[
\sigma[X](a,t) = 2\kappa_0 \rho_0(a) \int_0^t e^{-2\kappa_0 (t-s)} q(a,t,s)q(a,t,s)^T ds + e^{-2\kappa_0 t}(\nabla a X(a,t))\sigma_0(a)(\nabla a X(a,t))^T
\]

where $q(a,t,s) = q[X](a,t,s) = (\nabla a X(a,t))(\nabla a X(a,s))^{-1}$ is given in (2.39).

We take

\[
\tau[X](a,t) = \sigma[X](a,t) - \rho_0(a)I,
\]
consider $A(x, t) = X^{-1}(x, t)$, compose $\tau[X]$ with it,

$$\tau_X(x, t) = \tau[X](X^{-1}(x, t), t),$$  \hfill (4.92)

and solve the Stokes system, resulting in (3.7)

$$u^i_X(x, t) = k\Lambda^{-1}(R_i\tau^j_X + R_iR_mR_n(\tau^{mn}_X)).$$  \hfill (4.93)

We write symbolically

$$u_X = k\Lambda^{-1}(\mathbb{H}(\tau[X] \circ X^{-1}) = k\Lambda^{-1}(\tau X),$$  \hfill (4.94)

where $\mathbb{H}$ stands for the combinations of Riesz transforms that appear in (4.93),

$$\mathbb{H}_{imn} = \delta_{im}R_n + R_iR_mR_n,$$  \hfill (4.95)

and thus $k\Lambda^{-1}\mathbb{H}$ is the inverse of the Stokes system (3.1). Finally, we compose with $X(a, t)$ to obtain

$$F^i[X](a, t) = u^i_X(X(a, t), t).$$  \hfill (4.96)

Thus $F[X]$ is obtained via the succession of compositions

$$X \mapsto \tau[X] \mapsto \tau_X = \tau[X] \circ X^{-1} \mapsto u_X = k\Lambda^{-1}\mathbb{H}(\tau_X) \mapsto F = u_X \circ X.$$  \hfill (4.97)

The norm in $\mathcal{X}$ is

$$\|X\|_{\mathcal{X}} := \sup_{0 \leq t \leq T} \|X(\cdot, t)\|_{C^{1,\alpha}(\mathbb{R}^d)}.$$  \hfill (4.98)

We consider a fixed constant $M$ and the set

$$\mathcal{D} := \{X \in \mathcal{X} | X(0, a) = a, \frac{1}{2} \leq \det \nabla_a X(a, t) \leq \frac{3}{2}, \|X\|_{\mathcal{X}} \leq M\}.$$  \hfill (4.99)

The initial data for the PDE serve as parameters in the definition of $F$. We wish to show that two solutions $X_1 \in \mathcal{D}$ and $X_2 \in \mathcal{D}$ of the Equation (4.89) corresponding to the same $\rho_0, \tau_0$ are identical. In order to do so we establish

$$\|DF[X]|Y]\|_{\mathcal{X}} \leq C\|Y\|_{\mathcal{X}}$$  \hfill (4.100)

with a uniform constant $C$ that depends on $M$. We have to be careful to avoid taking derivatives of $\sigma_0$ and $\tau_0$. We start by noting that the map

$$X \mapsto \sigma[X]$$  

is Fréchet differentiable at $X \in \mathcal{D}$ as a map from the Banach space $\mathcal{X}$ to the Banach space $\Sigma = C(0, T; [C^0(\mathbb{R}^d) \cap L^1(\mathbb{R}^d)])$ of time-continuous maps with values in the space of matrices with spatially Hölder continuous and integrable coefficients. The derivative is a bounded linear map in $L(\mathcal{X}, \Sigma)$,

$$Y \mapsto (D\sigma[X])Y.$$  

The derivative has a complicated expression that depends on history but it can be easily obtained; the derivative of $q[X]$ is

$$((Dq[X])Y)(a, t, s) = (\partial_a Y(a, t))(\nabla_a X(a, s))^{-1} - q(a, t, s)(\nabla_a Y(a, s))(\nabla_a X(a, s))^{-1}.$$  \hfill (4.101)
Then the derivative of $\sigma[X]$ is
\[
((D\sigma[X])(Y))(a,t)
\]
\[
= 2\kappa_0\rho_0(a) \int_0^t e^{-2\kappa_0(t-s)}(Dq[X]Y)(a,t,s)q(a,t,s)^T ds
\]
\[+ 2\kappa_0\rho_0(a) \int_0^t e^{-2\kappa_0(t-s)}q(a,t,s)(Dq[X]Y)^T(a,t,s)ds
\]
\[+ e^{-2\kappa_0t} \left[ (\nabla_a Y(a,t))\sigma_0(a)(\nabla_a X(a,t))^T + (\nabla_a X(a,t))\sigma_0(a)(\nabla_a Y(a,t))^T \right].
\]

(4.102)

It is clear that
\[
\|D\sigma[X]Y\|_\Sigma \leq C\|Y\|_X,
\]
where
\[
\|\tau\|_\Sigma = \sup_{0 \leq t \leq T} \left\{ \|\tau(\cdot, t)\|_{L^\infty(\mathbb{R}^d)} + \|\tau(\cdot, t)\|_{L^1(\mathbb{R}^d)} + [\tau(\cdot, t)]_\alpha \right\}.
\]

(4.104)

The map
\[
X \mapsto \tau[X]
\]
is just a translation of $\sigma[X]$ by an $X$-independent amount, so it has the same derivative
\[
(D\tau[X])Y = (D\sigma[X])Y.
\]

(4.105)

Now, the map
\[
X(a,t) \mapsto \tau_X(x,t)
\]
is obtained from the map $\tau[X]$ by composition with $X^{-1}$. We cannot afford to take the derivative of $\tau_X$. However, $u_X$ is obtained from $\tau_X$ by using a linear smoothing operator of degree $-1$, so we can differentiate it. We take a path $X_\varepsilon(a,t) = X(a,t) + \varepsilon Y(a,t)$, differentiate with respect to $\varepsilon$ and then set $\varepsilon = 0$:
\[
((Du_X)Y)(x,t) = k\Lambda^{-1}H((D\tau[X]Y) \circ X^{-1}) - k\Lambda^{-1}H((\nabla_x \tau_X)(Y \circ X^{-1})).
\]

(4.106)

The first term appears when we differentiate $\tau[X] \circ X^{-1}$ with respect to $X$, keeping $X^{-1}$ fixed. The second term arises when we differentiate $X^{-1}$ using the fact that
\[
\frac{d}{d\varepsilon} X_{\varepsilon}^{-1}|_{\varepsilon=0}(x,t) = - (\nabla_x X^{-1})Y(X^{-1}(x,t),t).
\]

This fact is obtained by differentiating $\frac{d}{d\varepsilon}$.
\[
X_{\varepsilon}^{-1}(X(a,t) + \varepsilon Y(a,t)) = a,
\]
setting $\varepsilon = 0$, and reading at $a = X^{-1}(x,t)$. Then the derivative of $\tau[X] \circ X_{\varepsilon}^{-1}$ with respect to $\varepsilon$ is obtained using the chain rule
\[
\frac{d}{d\varepsilon} (\tau[X])(X_{\varepsilon}^{-1})|_{\varepsilon=0} = ((\nabla_a \tau[X]) \circ X^{-1})(\frac{d}{d\varepsilon} X_{\varepsilon}^{-1})|_{\varepsilon=0}
\]
\[= -((\nabla_a \tau[X]) \circ X^{-1})(\nabla_x X^{-1})(Y \circ X^{-1}) = -((\nabla_x \tau_X)(Y \circ X^{-1})).
\]
The second term in (4.106) is strictly speaking a distribution, as it involves derivatives of \( \tau_X \), and is only Hölder continuous. At this stage, we can view this as a formal calculation that will be justified in the end. Finally, we need to compose back with \( X(a, t) \).

\[
(DF(X)Y)(a, t) = ((Du_X)Y)(X(a, t), t) + (\nabla u_X)(X(a, t), t)Y(a, t). \tag{4.107}
\]

The second term in the expression above is \textit{unbounded} as a linear operator from \( C^{1,\alpha} \) to itself, simply because the coefficient \( \nabla u_X \) is not differentiable. But this term combined with the second problematic term in (4.106) produces a commutator that is better behaved, \( K[X]Y \). Thus, we have

\[
(DF(X)Y)(a, t) = k \left\{ \Lambda^{-1} \mathbb{H}((D\tau[X]Y) \circ A) \right\} (X(a, t), t) + (K[X]Y)(a, t), \tag{4.108}
\]

where

\[
(K[X]Y)(a, t) = (\nabla u_X)(X(a, t), t)Y(a, t) - k(\Lambda^{-1} \mathbb{H}((\nabla \tau[X])(Y \circ X^{-1}))(X(a, t), t). \tag{4.109}
\]

We would like to show that the map \( Y \mapsto K[X]Y \) is a bounded linear operator in \( C(0, T; (C^{1,\alpha})^d) \) with norm uniformly bounded for \( X \in D \). Because composition with \( X^{-1} \) and composition with \( X \) are both bounded linear operators \( C^{1,\alpha} \to C^{1,\alpha} \), with norms controlled by \( M \), the boundedness of \( K[X]Y \) is equivalent to the boundedness of the map

\[
\phi \mapsto L[X]\phi,
\]

where

\[
\phi(x, t) = Y(X^{-1}(x, t), t) \tag{4.110}
\]

and

\[
(L[X]\phi)(x, t) = k(\nabla_x \Lambda^{-1} \mathbb{H}_{imn} \tau^m_X)(x, t)\phi(x, t) - k(\Lambda^{-1} \mathbb{H}((\nabla_x \tau_X)\phi))(x, t). \tag{4.111}
\]

It is important to specify the tensorial nature of this commutator, as not all such expressions are better behaved than their individual terms. In our case, \( \tau_X \) is a fixed symmetric matrix in \( \Sigma \), \( k \) is a number, and the commutator is

\[
(L[X]\phi)^i = k \left[ \phi^p \partial_p \Lambda^{-1} \mathbb{H}_{imn} \tau^m_X - \Lambda^{-1} \mathbb{H}_{imn}(\phi_p \partial_p \tau^m_X) \right]. \tag{4.112}
\]

We can write the commutator as

\[
L[X]\phi = k \left[ \phi^p (\partial_p \Lambda^{-1} \mathbb{H}_{imn})(\tau^m_X) - (\partial_p \Lambda^{-1} \mathbb{H}_{imn})(\phi^p \tau^m_X) \right] - k\Lambda^{-1} \mathbb{H}_{imn}((\partial_p \phi^p) \tau^m_X). \tag{4.113}
\]

For incompressible \( X \) we could only consider divergence-free \( \phi \),

\[
\partial_p \phi^p = 0, \tag{4.114}
\]

but that would force us to work tangent to volume-preserving maps and would make the proof a little more complicated; we do not need to use that because the map

\[
\phi \mapsto \Lambda^{-1} \mathbb{H}_{imn}((\partial_p \phi^p) \tau^m_X)
\]
is bounded as a linear map from $C(0, T; [C^{1, \alpha}]^d)$ to itself. Clearly, because $\tau_X$ is H"older continuous, and because classical Calderon-Zygmund operators are bounded in H"older spaces, there is no difficulty in bounding H"older norms of derivatives of the expression $\Lambda^{-1}H_{\text{imn}}(\partial_p \phi^t \tau_X^m)$. Proving that the undifferentiated quantity is bounded is done using the fact that $\tau_X \in L^1 \cap L^\infty$, and therefore $(\partial_p \phi^t) \tau_X \in L^1 \cap L^\infty$. The operator $\Lambda^{-1}H$ continuously maps $L^1 \cap L^\infty$ to $L^\infty$. The operators $R_{\text{imn}} = \partial_p \Lambda^{-1}H_{\text{imn}}$ are sums of classical Calderon-Zygmund operators and multiples of the identity. The commutators

$$
\phi \mapsto \phi^t R_{\text{imn}}(\tau_X^m) - R_{\text{imn}}(\phi^t \tau_X^m)
$$

are bounded as operators from $C(0, T; [C^{1, \alpha}]^d)$ to itself. This is quite obvious for smooth $\tau_X$ but a little less obvious for $\tau_X \in \Sigma$. Let us write the kernel of $R_{\text{imn}}$ as $K_{\text{imn}}$, so the commutator is

$$
\int_{\mathbb{R}^d} K(x-y)(\phi(x) - \phi(y)) \tau_X(y) dy,
$$

where we did not write all the indices and the time dependence for ease of notation. The kernel $K$ is smooth away from the origin and is homogeneous of order $-d$. Differentiating in some direction and writing $K'$ for the singular (of order $d + 1$) kernel obtained by differentiating $K$, we have

$$
T[X] \phi = P.V. \int K'(x-y)(\phi(x) - \phi(y)) \tau_X(y) dy
$$

plus a nice operator $(\nabla_x \phi^t) R_{\text{imn}}(\tau_X^m)$. This last operator is clearly bounded in $C^\alpha$, with bound controlled by $M$, so we concentrate our attention on $T[X]$. Now we write

$$
T[X] \phi = \int_0^1 d\lambda \left[ P.V. \int_{\mathbb{R}^d} (x-y)K'(x-y)\nabla \phi(x + \lambda(y - x)) \tau_X(y) dy \right].
$$

The kernel $(x-y)K'(x-y)$ is homogeneous of order $-d$. It might have nonzero average on the unit sphere. Nevertheless, we subtract the value $\nabla_x \phi(x)$:

$$
T_1[X] \phi = \int_0^1 d\lambda \left\{ P.V. \int_{\mathbb{R}^d} (x-y)K'(x-y) [\nabla_x \phi(x + \lambda(y - x))] - \nabla_x \phi(x)] \tau_X(y) dy \right\}.
$$

The contributions left from the average on the unit sphere, if nonzero, are a constant multiple of $(\nabla_x \phi(x)) \tau_X(x)$ and $\nabla_x \phi(x) T_2 \tau_X(x)$,

$$
T_2(\tau_X)(x) = \int_{\mathbb{R}^d} (x-y)K'(x-y)(\tau_X(x) - \tau_X(y)) dy,
$$

both bounded with values in $C^\alpha$. The fact that $T_1[X] \phi$ is bounded in $C^\alpha$, and similarly that $T_2(\tau_X)$ is a H"older continuous function, are classical. A proof can be found in [3]. We have one more term in $DF[X]$, namely

$$
Y \mapsto k(\Lambda^{-1}H((D\tau[X]Y) \circ X^{-1})) \circ X.
$$
Its boundedness is equivalent to the boundedness of the maps of the type
\[ \phi \mapsto \Lambda^{-1}H(g X \nabla \phi) \] (4.122)
in \( C^{1,\alpha} \), where \( g_X \) is in \( \Sigma \). These are easily bounded because when we take spatial derivatives we arrive at
\[ \phi \mapsto R(g X \nabla \phi), \] (4.123)
which are bounded in \( C^{1,\alpha} \), and if we do not take derivatives the \( L^\infty \) boundedness follows as above from the fact that \( g_X \nabla \phi \in L^1 \cap L^\infty \). We have therefore verified the fact that \( DF[X] \) is bounded in \( X \) uniformly for \( X \in D \). The function \( F \) is locally Lipschitz and because solutions of (4.89) start from the identity, they coincide for short time. The same argument does not need the initial data to be the identity, but rather the same invertible \( C^{1,\alpha} \) transformation so uniqueness propagates because \( F \) is locally Lipschitz. This concludes the proof of the theorem. 

5. A regularization

We consider here the system formed by the (2.1) coupled with (3.1), (3.2) via (2.2). In this section, we consider the same potential (2.2) but we allow \( R \) to be a function of \( x \) and \( t \) (but of course, not of \( m \)). The properties that \( \sigma \) is symmetric, non-negative, and the bound (2.4) remain valid. In order to see what dependences on \( R \) would be allowed by the energy considerations, we repeat the calculation leading to (3.10): we take (2.1), multiplied by \( k (\log f + U) \), integrate, take (3.1) multiply by \( u \), integrate, and add the two: we deduce

\[
\frac{d}{dt} \left[ \int_{\mathbb{R}^d} f \log f \, dm \, dx \right] + \frac{1}{2} \int_{\mathbb{R}^d} \text{Tr}[\sigma(x, t)] \, dx + k \int_{\mathbb{R}^d} \frac{D_t R}{R} \left( \int_{\mathbb{R}^d} U \, dm \right) \, dx
\]

\[= -\int_{\mathbb{R}^d} |\nabla_x u(x, t)|^2 \, dx - k\epsilon \int_{\mathbb{R}^d} \int_{\mathbb{R}^d} f |\nabla_m (\log f + U)|^2 \, dm \, dx. \] (5.1)

It is thus clear that \( D_t R \geq 0 \) is energetically favorable. The solution on characteristics might be less explicit, however the equation obeyed by \( \sigma \) is easily obtained by multiplying (2.1) by \( (m \otimes m) / R^2 \) and integrating. The result is very similar to (2.29):

\[ D_t \sigma = (\nabla_x u)\sigma + \sigma(\nabla_x u)^T - \frac{2\epsilon}{R^2} \sigma + \frac{2\epsilon}{R^2} \sigma - \frac{2D_t R}{R} \sigma. \] (5.2)

We again see that \( D_t R \geq 0 \) has the effect of an additional damping. In fact, if \( \frac{D_t R}{R} \) is a constant, then the effect is precisely one of enhanced damping and is similar to the situation covered previously in Theorem 4.3, but in a better regime. More interesting, perhaps, is a damping that responds locally to very high rate of strain in the fluid. Let us consider a coupled system in which, in addition to (2.1) coupled to (3.1), (3.2) via (2.2), \( R \) evolves according to

\[ D_t R = \delta(|\nabla_x u(x, t)|) \, R \] (5.3)

and \( \delta(g) \) is a smooth nonnegative function of one nonnegative variable \( g \), that vanishes for \( g \leq \frac{\kappa}{2} \),

\[
\delta(g) = \begin{cases} 
0, & \text{if } g \leq \frac{\kappa}{2}, \\
C_0 \sqrt{\kappa^2 + g^2}, & \text{if } g \geq \kappa,
\end{cases}
\] (5.4)
and satisfies
\[ |\delta'(g)| \leq 2C_0 \] (5.5)
for all \( g \geq 0 \). In particular \( \delta(|\nabla_x u|) \) satisfies
\[ \delta = C_0 \sqrt{\kappa^2 + |\nabla_x u|^2} \] (5.6)
if \( |\nabla_x u| \geq \kappa > 0 \). The constant \( \kappa > 0 \) is fixed, and represents the order of magnitude of the largest permissible temporal growth rate. The constant \( C_0 > 0 \) is a function of the dimension of space (this is needed because of the tensorial nature of the calculation), so that
\[ 3c|\nabla_x u| - 2\delta \leq 0 \] (5.7)
holds, if \( |\nabla_x u| \geq \kappa \), where \( c \) is the norm equivalence constant that bounds \( \max_u |\partial u| \leq c|\nabla u| \). Thus \( c \) is a norm equivalence constant in \( \mathbb{R}^d \), and we can choose \( C_0 = \frac{3}{2}c \).

Because the problem evolves in time, the initial data \( f_0 \) and \( R_0 \) need to be specified. We assume that \( f_0 \) is non-negative and smooth enough in \( x \), decaying fast enough in \( m \), \( (f_0 \in W^{1,p}(\mathbb{R}^d; L^1((1 + |m|^2)dm))) \) and \( R_0(a) \geq R_{\text{min}} > 0 \) is in \( W^{1,p}(\mathbb{R}^d) \) with \( p > d \). Then from (5.2) and from the fact that \( \sigma^{ii} \) are nonnegative for each fixed \( i \), it follows that
\[ D_t \text{Tr}(\sigma) \leq 2c\kappa \text{Tr}(\sigma) + \frac{2de}{R^2} \] (5.8)
which, after integration on characteristics, results in
\[ \sup_x \text{Tr}(\sigma(x,t)) \leq e^{2c\kappa t} \left[ \sup_x \text{Tr}(\sigma_0(x)) + \frac{de}{c\kappa R^2_{\text{min}}} \|\rho_0\|_{L^\infty(\mathbb{R}^d)} \right] = N_0(t). \] (5.9)
In view of (2.4), we thus have control of the \( L^\infty \) norm of \( \sigma \) in time. Once this is achieved, based on previous results for similar models, we may expect to prove regularity. The method of proof of regularity that we employ here, given boundedness of \( \sigma \), is the simplest and most explicit. The idea is to differentiate the equation, pay the price of differentiating the advective term, and do an Eulerian calculation in \( L^p \) spaces with \( p > d \). Working in Lebesgue spaces makes it easy to take advantage of incompressibility. The approach used in the proof of global existence for Smoluchowski equations on compact manifolds coupled with time-dependent Stokes equations in [6], and in the proof of global existence [11] for Smoluchowski equations coupled to Navier-Stokes equations in \( d = 2 \), can be adapted to the present situation as well. In other words, we can consider the time dependent \( R \) evolving according to (5.3), a Fokker-Planck equation (2.1) with potential given in (2.2) coupled via (2.3) to a fluid velocity evolving according to the time-dependent Stokes equation in \( d = 3 \) or Navier-Stokes equation in \( d = 2 \). The method of [4], used in [9] to prove the global regularity for \( d = 2 \) Smoluchowski equations on compact manifolds coupled with Navier-Stokes equations, requires a little less smoothness on the Navier-Stokes initial data. That method is also Eulerian, but uses commutation and a penalty that results in a controlled loss of regularity. The proof is a bit more technical as it employs Besov spaces and paradifferential calculus, but in principle it can be adapted to the noncompact particle phase space case with spatially dependent \( R \). We will not
pursue these matters here, but rather content ourselves with the simplest proof, in the simplest nontrivial case. Differentiating (5.2) we obtain

\[
D_t(\partial_k \sigma^{ij}) = -(\partial_k u^j)(\partial_k \sigma^{ij}) + (\partial_k u^i)(\partial_k \sigma^{ij}) + (\partial_k u^j)(\partial_k \sigma^{ij}) - 2\delta(\partial_k \sigma^{ij})
\]

\[
-\frac{2\nu}{R^2}(\partial_k \sigma^{ij}) + (\partial_k^2 u^i)\sigma^{ij} + (\partial_k^2 u^j)\sigma^{ij} - (\partial_k \frac{2\nu}{R^2})\sigma^{ij}
\]

(5.10)

We multiply by \(\partial_k(\sigma^{ij})\) and sum. The terms explicitly involving first derivatives of \(u\) are bounded using (5.7):

\[
[(-\partial_k u^j)(\partial_k \sigma^{ij}) + (\partial_k u^i)(\partial_k \sigma^{ij}) + (\partial_k u^j)(\partial_k \sigma^{ij}) - 2\delta(\partial_k \sigma^{ij}) - \frac{2\nu}{R^2}(\partial_k \sigma^{ij})] (\partial_k \sigma^{ij})
\]

\[
\leq 3\epsilon|\nabla_x \sigma|^2.
\]

(5.11)

The terms explicitly involving second derivatives of \(u\) are bounded using (5.9):

\[
[(\partial_k^2 u^i)\sigma^{ij} + (\partial_k^2 u^j)\sigma^{ij}] (\partial_k \sigma^{ij}) \leq CN_0(t)|\nabla_x \sigma|^2 |\nabla_x \sigma|.
\]

(5.12)

The term containing a derivative of \(\delta\) is bounded using (5.5), which implies

\[
|\nabla_x \delta(\nabla_x u)| \leq 2C_0|\nabla_x \sigma|
\]

(5.13)

and therefore

\[
-2\partial_k(\delta)(\sigma^{ij})(\partial_k \sigma^{ij}) \leq CN_0(t)|\nabla_x \sigma|^2 |\nabla_x \sigma|.
\]

(5.14)

Summarizing, we have so far, pointwise:

\[
\frac{1}{2} D_t|\nabla_x \sigma|^2
\]

\[
\leq 3\epsilon|\nabla_x \sigma|^2 + CN_0(t)|\nabla_x \sigma| |\nabla_x \sigma| + CN_0(t)\epsilon \left|\nabla_x \sigma\right| + C\epsilon \left|\nabla_x \rho\right| \left|\nabla_x \sigma\right|.
\]

(5.15)

We have to bound the terms involving \(\nabla_x \rho\) and \(\nabla_x \frac{1}{R^2}\). In view of the fact that \(\rho(x, t) = \rho_0(A(x, t))\), we have

\[
\nabla_x \rho = \left(\frac{\nabla_x A}{R^2}\right)^T (\nabla_x \rho_0)(A).
\]

(5.16)

Now, because of (2.57) and (5.3) we have

\[
D_t \left(\frac{\nabla_x A}{R^2}\right) = -\left(\frac{\nabla_x A}{R^2}\right) \left(\nabla_x u + 2\delta \mathbf{I}\right)
\]

(5.17)

and, in view of (5.7), we deduce that

\[
|\nabla_x A(x, t)| \leq \frac{e^{xt}}{R_{min}},
\]

(5.18)

and consequently

\[
\frac{|\nabla_x \rho(x, t)|}{R(x, t)^2} \leq C \frac{e^{xt}}{R_{min}} |\nabla_x \rho_0(A(x, t))| \leq C \frac{e^{xt}}{R_{min}} \|\nabla_x \rho_0\|_{L^\infty(\mathbb{R}^d)}.
\]

(5.19)
The term involving $\nabla_x \left( \frac{1}{R^2} \right)$ is treated using (5.3):

$$
D_t \frac{\nabla_x R}{R^3} = - (\nabla_x u)^T - 2\delta \left( \frac{\nabla_x R}{R^3} \right) + \frac{\nabla_x \delta}{R^2},
$$

and therefore in view of (5.7) and (5.13) we deduce

$$
D_t \left( \frac{\nabla_x R}{R^3} \right) \leq c_k \left( \frac{\nabla_x R}{R^3} \right) + C' \frac{|\nabla_x \nabla_x u|}{R^2},
$$

Integrating on characteristics, we obtain

$$
\frac{|\nabla_x R(x,t)|}{R^3(x,t)} \leq C e^{c_k t} \left( \frac{|\nabla_a R_0(A(x,t))|}{R_{min}^3} + \frac{1}{R_{min}^2} \int_0^t |\nabla_x \nabla_x u(X(A(x,t),s),s)| ds \right).
$$

Now we collect the terms in (5.15), divide by $|\nabla_x \sigma|$ and use (5.19) and (5.22):

$$
D_t |\nabla_x \sigma(x,t)| \leq 3c_k |\nabla_x \sigma(x,t)| + CN_0(t)|\nabla_x \nabla_x u(x,t)| + C e^{c_k t} \left( \frac{\nabla_a \rho_0(A(x,t))}{R_{min}} \right)
$$

$$
\times e^{c_k t} \left( \frac{|\nabla_x R_0(A(x,t))|}{R_{min}^3} + \frac{1}{R_{min}^2} \int_0^t |\nabla_x \nabla_x u(X(A(x,t),s),s)| ds \right).
$$

The inequality (5.23) has the form

$$
D_t (y(x,t)) \leq 3c_k y(x,t) + C(t)z(x,t) + D(t) \int_0^t z(X(A(x,t),s),s) ds + E(x,t),
$$

where

$$
\begin{align*}
    y(x,t) &= |\nabla_x \sigma(x,t)|, \\
    z(x,t) &= |\nabla_x \nabla_x u(x,t)|, \\
    C(t) &= CN_0(t) = C e^{2c_k t} \left[ \sup_x Tr(\sigma_0(x)) + \frac{d_0}{c_k R_{min}} \|\rho_0\|_{L^\infty(\mathbb{R}^d)} \right], \\
    D(t) &= C e^{c_k t} \left( N_0(t) + \|\rho_0\|_{L^\infty(\mathbb{R}^d)} \right) e^{c_k t}, \\
    E(x,t) &= C e^{c_k t} e^{c_k t} \left\{ |\nabla_a \rho_0(A(x,t))| + (N_0(t) + \|\rho_0\|_{L^\infty(\mathbb{R}^d)}) \frac{|\nabla_x R_0(A(x,t))|}{R_{min}^3} \right\}.
\end{align*}
$$

The interested reader can check that (5.23) is dimensionally balanced. Now it is time to start looking at $L^p$ norms. We start by noting that

$$
\left\| \int_0^t z(X(A(x,t),s),s) ds \right\|_{L^p(dx)} \leq \int_0^t \|z(X(A(x,t),s),s)\|_{L^p(dx)} ds = \int_0^t \|z(x,s)\|_{L^p(dx)} ds
$$

because of incompressibility. We integrate (5.24) on characteristics and take the $L^p$ norm. In order simplify the answer we use the fact that $C(t), D(t)$ are non-decreasing
functions of time. We obtain
\[
\|y(\cdot, t)\|_{L^p(\mathbb{R}^d)} \\
\leq e^{3c_\kappa t} \left\{ \|y_0\|_{L^p(\mathbb{R}^d)} + \|E(\cdot, t)\|_{L^p(\mathbb{R}^d)} + \int_0^t (C(s) + D(t)(t-s)) \|z(\cdot, s)\|_{L^p(\mathbb{R}^d)} ds \right\}.
\]
(5.27)

We recall (3.9), which implies
\[
\nabla_x \nabla_x u = k \tilde{R}(\nabla_x \tilde{\sigma}).
\]
(5.28)

In view of the well-known boundedness of Riesz transforms in \(L^p\) spaces, we deduce that
\[
\|z(\cdot, t)\|_{L^p(\mathbb{R}^d)} \leq C k \|y(\cdot, t)\|_{L^p(\mathbb{R}^d)},
\]
(5.29)

We note also that \(\|E(\cdot, t)\|_{L^p(\mathbb{R}^d)}\) is explicitly a sum of norms of initial data multiplied by exponentials of time. Now a simple Gronwall argument provides an a priori bound for \(y\) in \(L^p\),
\[
\|y(\cdot, t)\|_{L^p(\mathbb{R}^d)} \leq F_p(t),
\]
(5.30)

with \(F_p(t)\) an explicit function of time with exponential growth and depending only on norms of initial data \(\|\sigma_0\|_{L^\infty(\mathbb{R}^d)}, \|\rho_0\|_{L^\infty(\mathbb{R}^d)}, \|\nabla_x \sigma_0\|_{L^p(\mathbb{R}^d)}, \) and \(R_{\text{min}}\).

**Theorem 5.1.** Let \(f\) solve (2.1) with \(U\) given by (2.2) and \(R\) evolving according to (5.3) with smooth \(\delta\) satisfying (5.5) and (5.7). Let \(u\) be obtained by solving (3.1), (3.2), with \(\sigma\) given by (2.3). Assume that the initial distribution \(f_0\) and \(R_0\) satisfy
\[
\left\{ \begin{array}{l}
\sup_{x \in \mathbb{R}^d} \int_{\mathbb{R}^d} f_0(x, m)(1 + |m|^2)dm < \infty,
R_0(x) \geq R_{\text{min}} > 0, \\
\int_{\mathbb{R}^d} |\nabla_x R_0(x)|^p dx < \infty, \\
\int_{\mathbb{R}^d} \left( \int_{\mathbb{R}^d} (1 + |m|^2)|\nabla_x f_0(x, m)|dm \right)^p dx < \infty,
\end{array} \right.
\]
(5.31)

with \(p > d\). Then the solution exists for all time, is unique and obeys the a priori bounds (5.9) and (5.30).

**Proof.** The space \(W^{1,p}(\mathbb{R}^d)\) is a space of local existence and uniqueness of solutions. The condition \(\inf_x R(x, t) \geq R_{\text{min}}\) is invariant in time, because \(R(x, t) \geq R_0(A(x, t)) \geq R_{\text{min}}\). Then the a priori bounds are enough to finish the proof. \(\square\)

**Remark 5.2.** The a priori bounds hold in any \(L^p\), \(1 < p < \infty\).

**Remark 5.3.** Similar theorems hold if we replace steady-state Stokes equation with time-dependent Stokes equation in \(d = 2, 3\) and with Navier-Stokes equation in \(d = 2\).

**Remark 5.4.** Higher regularity of solutions can be obtained without difficulty from higher regularity of the initial data.
6. Large data

The problem of global regularity for arbitrary smooth initial data is open. The system formed by (2.1) and (3.1) has potentially finite-time blowup. Indeed, consider the system in $d = 1$:

\[
\begin{aligned}
D_t f(x, m, t) + u_x \partial_m (mf(x, m, t)) &= \epsilon \partial_m (f \partial_m (\log f(x, m, t) + U(m))) \\
U(m) &= \frac{m^2}{2R^2}, \\
\sigma(x, t) &= \int_{-\infty}^{\infty} \frac{m^2}{R^2} f(x, m, t) \, dm, \\
u_x &= kH\sigma,
\end{aligned}
\]

where $H$ is the Hilbert transform. This system is a 1-d analogue of (2.1) with (2.2), (2.3), and (3.9). This blows up in finite time. Indeed, we multiply the linear Fokker-Planck equation by $\frac{m^2}{R^2}$ and integrate $dm$ to obtain the analogue of (2.29):

\[
D_t \sigma = 2u_x \sigma - \frac{2\epsilon}{R^2} \sigma + \frac{2\epsilon}{R^2},
\]

which then, in view of (6.1), is

\[
D_t \sigma = 2k\sigma H\sigma - \frac{2\epsilon}{R^2} \sigma + \frac{2\epsilon}{R^2}
\]

and resembles the baby vorticity equation [8]. The blow up argument of [8] works here also, notwithstanding the fact that the present equation is computed on characteristics. We form $z = H\sigma + i\sigma$ and use the fact that

\[
H(\sigma H\sigma) = \frac{1}{2} ((H\sigma)^2 - \sigma^2)
\]

to deduce

\[
D_t z = k z^2 - \frac{2\epsilon}{R^2} z + i \frac{2\epsilon}{R^2},
\]

which blows up in finite time.

This simple example does not capture incompressibility and the tensorial nature of the problem, just like the baby vorticity equation of [8] does not capture them. The main quadratic nonlocal nonlinearity of (2.29) is modeled by $\sigma H\sigma$, and this is the only available model in one dimension. Unlike the incompressible Euler equation however, the system we study has an additional dissipative structure, and we believe that this is of some significance and represents the main fallacy of the one-dimensional model. We describe below a simple scalar model that addresses this issue. We take henceforth $d = 2$. We consider the variables

\[
\begin{aligned}
a(x, t) &= \frac{1}{2} (\sigma^{11}(x, t) - \sigma^{22}(x, t)), \\
b(x, t) &= \sigma^{12}(x, t) = \sigma^{21}(x, t), \\
c(x, t) &= \sigma^{11}(x, t) + \sigma^{22}(x, t) = Tr (\sigma(x, t)).
\end{aligned}
\]

The gradient of velocity is represented by $\lambda(x, t)$, $\mu(x, t)$, and $\omega(x, t)$ given in (2.69), (2.70), and (2.71). The Equation (2.29) can be written as the system

\[
\begin{aligned}
D_t a &= -\omega b + \lambda c - \frac{2\epsilon}{R^2} a, \\
D_t b &= \omega a + \mu c - \frac{2\epsilon}{R^2} b, \\
D_t c &= 4\lambda a + 4\mu b - \frac{2\epsilon}{R^2} c + 4\frac{\epsilon}{R^2} \rho, \\
D_t \rho &= 0.
\end{aligned}
\]
As we saw before, if we couple this with an equation for the velocity (steady or unsteady Stokes, or Navier-Stokes), then the regularity issue is decided by whether or not we can bound \( c \) in \( L^\infty(dx) \). We note in passing that the co-rotational system corresponds to \( \lambda = \mu = 0 \) in the system above, and the bound for \( c \) follows immediately.

Let us multiply the \( c \) equation by \( c^2/4 \), the \( a \) equation by \( 2a \), the \( b \) equation by \( 2b \), and subtract the last two from the first. We obtain

\[
D_t \left( \frac{c^2}{4} - a^2 - b^2 \right) = -\frac{4\epsilon}{R^2} \left( \frac{c^2}{4} - a^2 - b^2 \right) + \frac{2\epsilon}{R^2} \rho c.
\]

This cancellation of nonlinearity is not surprising because

\[
\frac{c^2}{4} - a^2 - b^2 = \det(\sigma),
\]

and the determinant is conserved along particle trajectories if \( \epsilon = 0 \). The matrix \( \sigma \) is symmetric and positive by construction, and is given in terms of \( a, b, c \) by

\[
\sigma = \begin{pmatrix}
\frac{c}{2} + a & b \\
b & \frac{c}{2} - a
\end{pmatrix}.
\]

The positivity of the matrix is equivalent (in this case) to the positivity of the determinant, i.e. to

\[
\frac{c^2}{4} - a^2 - b^2 > 0.
\]

The two eigenvalues of \( \sigma \)

\[
z_{1,2} = \frac{c}{2} \pm \sqrt{a^2 + b^2}
\]

are both positive. We have of course \( c = z_1 + z_2 > 0 \) and \( z_1 - z_2 = 2\sqrt{a^2 + b^2} \). Because of (6.9), \( c \) controls \( \sqrt{a^2 + b^2} \). But, on the other hand, unbounded growth of \( c \) on any trajectory cannot happen without unbounded growth of \( \sqrt{a^2 + b^2} \) on the same trajectory. Indeed, if a particle path would be such that \( \sqrt{a^2 + b^2} \) is bounded on it, but \( c \) grows without bound or blows up in finite time, then for large enough time we would have

\[
-\frac{\epsilon}{R^2} \left( \frac{c^2}{4} - a^2 - b^2 \right) + \frac{2\epsilon}{R^2} \rho c \leq 0
\]

on the particle path (because \( \rho \) is bounded), and from (6.6) we would arrive at the contradiction that \( c \) remains bounded. From (6.5) we can write

\[
c(x,t) = e^{-\frac{\epsilon}{R^2} t} c_0(A(x,t)) + \int_0^t e^{-\frac{\epsilon}{R^2} (t-s)} \left( \lambda a + \mu b - \frac{\epsilon}{R^2} \rho_0 \right) (X(A(x,t),s),s) ds.
\]

Thus, exponential growth or blowup of \( c \) is possible only if \( (\lambda a + \mu b) \) has time integrals on particle paths that are positive and grow at least exponentially and without bound.

In two dimensions we can express the velocity in terms of a stream function \( \psi \) and write

\( u^1 = -\partial_2 \psi \), \( u^2 = \partial_1 \psi \). Then \( \omega = \Delta \psi \), and therefore

\[
\lambda = \partial_1 \partial_2 (-\Delta)^{-1} \omega = B \omega
\]

(6.12)
and

\[ \mu = -\frac{1}{2} (\partial_1^2 - \partial_2^2) (-\Delta)^{-1} \omega = -A \omega, \]  

(6.13)

where

\[ A = \frac{1}{2} (\partial_1^2 - \partial_2^2) (-\Delta)^{-1} = \frac{1}{2} (R_1^2 - R_2^2) \]  

(6.14)

and

\[ B = \partial_1 \partial_2 (-\Delta)^{-1} = R_1 R_2 \]  

(6.15)

are bounded operators in \( L^p(\mathbb{R}^2) \) spaces. They are also bounded, selfadjoint in \( L^2(\mathbb{R}^2) \), they commute \( AB = BA \), and each is given by a multiplier at Fourier level:

\[ \hat{A}h(\xi) = \frac{\xi_1^2 - \xi_2^2}{2|\xi|^2} \hat{h}(\xi) \]  

(6.16)

and

\[ \hat{B}h(\xi) = -\frac{\xi_1 \xi_2}{|\xi|^2} \hat{h}(\xi). \]  

(6.17)

Note that

\[ 4(A^2 + B^2) = 1. \]  

(6.18)

Let us consider now the time independent Stokes system (3.1). Taking the curl of (3.1) and inverting the Laplacian we obtain

\[ \omega = 2k (Ab - Ba). \]  

(6.19)

Consequently, from (6.12) we have

\[ \lambda = 2k (-B^2 a + A B b), \]  

(6.20)

and from (6.13) we have

\[ \mu = 2k (A B a - A^2 b). \]  

(6.21)

It is convenient to measure time in units of \( \frac{1}{2k} \). The system (6.5) coupled with (6.19), (6.12), and (6.13) is

\[
\begin{align*}
D_t a &= -(A(b) - B(a))b + [-B^2(a) + AB(b)]c - \varepsilon^{-1} a, \\
D_t b &= (A(b) - B(a))a + [AB(a) - A^2(b)]c - \varepsilon^{-1} b, \\
D_t c &= 4 \left( [-B^2(a) + AB(b)]a + [AB(a) - A^2(b)]b \right) - \varepsilon^{-1} c + 2\varepsilon^{-1} \rho, \\
D_t \rho &= 0,
\end{align*}
\]  

(6.22)

where \( \varepsilon \) is given by (4.6). The all-important term \( \lambda a + \mu b \) is given by

\[ \frac{1}{2k}(\lambda a + \mu b) = AB(\omega) - bA(\omega) = -aB(B(a) - A(b)) + bA(B(a) - A(b)). \]  

(6.23)
This expression is quadratic, nonlocal, and has negative spatial integral. While blow-up requires the pointwise positivity of this expression (6.11), its spatial average is negative. Integrating the third equation in (6.22) we obtain, using the selfadjointness of $A$ and $B$, \begin{equation}
abla \int_{\mathbb{R}^2} c(x,t)dx + 4 \int_0^t e^{-\frac{\nu}{2}} \int_{\mathbb{R}^2} |B(a)(x,s) - A(b)(x,s)|^2 dx 
abla = e^{-\frac{\nu}{2}} \int_{\mathbb{R}^2} c_0(x)dx + 2(1 - e^{-\frac{\nu}{2}}) \int_{\mathbb{R}^2} \rho_0(x)dx. \tag{6.24} \end{equation}
This is just the energetic bound on $\int Tr \sigma dx$. We saw earlier that if the size of initial data is of order 1, then small enough $\varepsilon$ leads to global existence of solutions. The case of large initial data and moderate $\varepsilon$ is wide open. In order to clarify the issues, we will take $\varepsilon = \infty$. In this case the system is simplified somewhat, because in view of (6.6) we have \begin{equation} c = 2\sqrt{a^2 + b^2 + d_0(x,t)}, \tag{6.25} \end{equation}
where \begin{equation} d_0(x,t) = \text{Det} (\sigma_0(A(x,t))), \tag{6.26} \end{equation}
and so the first two equations of (6.5) become just \begin{equation} \begin{cases} D_t a = -b(Ab - Ba) + 2\sqrt{(a^2 + b^2 + d_0)} B(Ab - Ba), \\ D_t b = a(Ab - Ba) - 2\sqrt{(a^2 + b^2 + d_0)} A(Ab - Ba), \end{cases} \tag{6.27} \end{equation}
and consequently \begin{equation} D_t \sqrt{(a^2 + b^2 + d_0)} = 2aB(Ab - Ba) - 2bA(Ab - Ba). \tag{6.28} \end{equation}
Integrating in space we again get the bound on the trace. Blowup for the system (6.27) is still too difficult to analyse. The simplest possible didactic model of a scalar equation exhibiting this kind of quadratic nonlocal structure with an $L^1$ dissipation is a scalar equation of the form \begin{equation} \partial_\tau \tau = -\tau(A^2 \tau), \tag{6.29} \end{equation}
where $A$ is a bounded selfadjoint operator given in Fourier representation by multiplication by a function that is homogeneous of degree zero, \begin{equation} \hat{A}h(\xi) = \alpha(\xi) \hat{h}(\xi), \tag{6.30} \end{equation}
with $\alpha(\lambda \xi) = \alpha(\xi)$, $\alpha(\xi) \in \mathbb{R}$, $\int_{\mathbb{S}^2} \alpha(\xi)d\xi = 0$, $\alpha$ smooth on the unit sphere. The unknown $\tau(x,t)$ (representing the trace of $\sigma$) is a positive scalar. We note the salient features of this model; if $\tau_0$ is smooth and positive, (let us consider for instance $\tau(\cdot,0) \in L^1 \cap C^{0,s}$ with $s > 0$) then the solution exists and is unique for some time, and remains positive as long as it exists. If a bound on $\|\tau\|_{L^\infty}$ is given, then higher regularity of the solution follows. Integrating the equation, we have the dissipative bound \begin{equation} \int \tau(x,t)dx + \int_0^t \|A\tau\|^2_{L^2} dt \leq \int \tau(x,0)dx, \tag{6.31} \end{equation}
and multiplying by $A^2\tau$ and integrating we have

$$
\int |A\tau(x,t)|^2 \, dx \leq \int |A\tau(x,0)|^2 \, dx. \tag{6.32}
$$

We are going to give an example of global regularity for some large data for such an equation, in the spatially periodic case. For simplicity we take the period to be $2\pi$ in each direction. We write

$$
\tau(x,t) = \tau_0(t) + \sum_{k \in \mathbb{Z}^2 \setminus \{0\}} \tau_k(t) e^{ik \cdot x}. \tag{6.33}
$$

The requirement that $\tau$ be real is implemented by $\tau_k = \tau_{-k}$. We will consider the symbol $\alpha(k)$, and assume that it is real valued, even, $\alpha(-k) = \alpha(k)$, and bounded $|\alpha(k)| \leq \Gamma$. We also assume $\alpha(0) = 0$. The Equation (6.29) is the infinite system of ODEs

$$
\frac{d\tau_l}{dt} = -\sum_{k+j=l} \tau_k \alpha^2(j) \tau_j. \tag{6.34}
$$

**Proposition 6.1.** Let $\tau(x,0) = \sum_k \tau_k(0) e^{ik \cdot x}$ satisfy

$$
\tau_{-k}(0) = \overline{\tau_k(0)},
$$

$$
\sum_k (1 + |k|)^s |\tau_k(0)| \leq C_s(0) < \infty, \tag{6.36}
$$

for some $s > 0$, and

$$
\tau_0(0) \geq \sum_{k \neq 0} |\tau_k(0)|. \tag{6.37}
$$

Then the solution of (6.34) exists for all time and obeys

$$
\tau_{-k}(t) = \overline{\tau_k(t)},
$$

$$
\sum_k (1 + |k|)^s |\tau_k(t)| \leq C_s(0) e^{2s+1} \tau_0(0) \Gamma^2 t < \infty, \tag{6.39}
$$

and

$$
\tau_0(t) \geq \sum_{k \neq 0} |\tau_k(t)| \tag{6.40}
$$

which of course implies that $\tau(x,t)$ remains smooth and positive.

**Proof.** We start with the dissipation, the equation at $l = 0$, which reads

$$
\frac{d\tau_0}{dt} = -\sum_{k \neq 0} \alpha(k)^2 |\tau_k(t)|^2. \tag{6.41}
$$

Now for $l \neq 0$ we have

$$
\frac{d\tau_l}{dt} = -\sum_{k+j=l, k \neq 0, j \neq 0} \tau_k \alpha^2(j) \tau_j - \tau_0 \alpha^2(l) \tau_l,
$$
and therefore
\[
\frac{d|\tau_l|}{dt} \leq -\tau_0 \alpha^2(l)|\tau_l| + \sum_{k+j=l, \ k \neq 0, \ j \neq 0} |\tau_k| \alpha^2(j)|\tau_j|.
\] (6.42)

Summing in \(l \neq 0\), we obtain
\[
\frac{d}{dt} \sum_{l \neq 0} |\tau_l| \leq -\tau_0 \sum_{l \neq 0} \alpha^2(l)|\tau_l| + \sum_{l \neq 0} \sum_{k+j=l, \ k \neq 0, \ j \neq 0} |\tau_k| \alpha^2(j)|\tau_j|
= -\tau_0 \sum_{l \neq 0} \alpha^2(l)|\tau_l| + \sum_{j \neq 0} \sum_{k \neq 0} \alpha^2(j)|\tau_j||\tau_k| - \sum_{j \neq 0} \alpha^2(j)|\tau_j||\tau_l|
= \left( \sum_{j \neq 0} \alpha^2(j)|\tau_j| \right) \left( -\tau_0 + \sum_{k \neq 0} |\tau_k| \right) + \frac{d}{dt}\tau_0.
\]

This results in the inequality
\[
\frac{d}{dt} \left( \sum_{l \neq 0} |\tau_l| - \tau_0 \right) \leq \left( \sum_{j \neq 0} \alpha^2(j)|\tau_j| \right) \left( \sum_{l \neq 0} |\tau_l| - \tau_0 \right) .
\] (6.43)

If the initial data is non-positive, the quantity remains non positive. Thus we have the invariance of this cone in function space. Once this is achieved, we know
\[
\sum_{l \neq 0} |\tau_l(t)| \leq \tau_0(t) \leq \tau_0(0),
\] (6.44)

which implies an \(L^\infty\) bound on \(\tau\). We know this should be sufficient for regularity. In the present situation, the proof of regularity is quite straightforward: we multiply (6.42) by \((1 + |l|)^s\) and sum. Using \(|l| \leq 2 \max(|k|, |j|)\) and \(|\alpha(j)| \leq \Gamma\), we obtain
\[
\frac{d}{dt} \sum_{l \neq 0} (1 + |l|)^s|\tau_l| \leq 2^{s+1} \tau_0(0) \Gamma^2 \left( \sum_{l \neq 0} (1 + |l|)^s|\tau_l| \right).
\]

Using (6.44) we deduce
\[
\frac{d}{dt} \sum_{l \neq 0} (1 + |l|)^s|\tau_l| \leq 2^{s+1} \tau_0(0) \Gamma^2 \left( \sum_{l \neq 0} (1 + |l|)^s|\tau_l| \right),
\] (6.45)

which finishes the proof.

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