

# Dissipative solutions of the Euler equations

Camillo De Lellis

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# The Euler and Navier-Stokes equations

They describe the motion of an incompressible fluid under some assumptions.

The unknowns are the pressure (a scalar field) and the velocity.

$$\partial_t \mathbf{v} + \operatorname{div}(\mathbf{v} \otimes \mathbf{v}) + \nabla p = \nu \Delta \mathbf{v}$$

$$\operatorname{div} \mathbf{v} = 0$$

$\nu > 0$  Navier Stokes

$\nu = 0$  Euler

The  $i$ -th component of the advective term  $\operatorname{div}(\mathbf{v} \otimes \mathbf{v})$  is given by

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# The Euler and Navier-Stokes equations II

In this talk we will consider solutions which are defined on the entire 3-dimensional (resp. 2-dimensional) space and over some time interval  $I$ .

In several occasions we will consider **periodic solutions** and thus the domain of definition will be  $\mathbb{T}^3 \times I$  (or  $\mathbb{T}^2 \times I$ ).  
 $I$  might be

- ▶ a bounded interval or a half line; in this case the left endpoint will be 0 and the equations will be complemented with an initial condition (**Cauchy problem**):

$$v(0, \cdot) = v_0$$

- ▶ the entire real line (**ancient solutions**)

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The Euler equations were derived more than 250 years ago (by Euler!)  
The Navier-Stokes equations date back to the middle of the 19th century.

Nonetheless several fundamental and outstanding open questions are still open: the most famous one is the [blow-up problem](#) for 3-dimensional solutions of the Cauchy problem.

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# The conservation of energy

If  $(v, p)$  is a  $C^1$  solution, we can scalar multiply the first equation by  $v$ :

$$\partial_t \frac{|v|^2}{2} + \operatorname{div} \left( \left( \frac{|v|^2}{2} + p \right) v \right) = \nu \left( \Delta \frac{|v|^2}{2} - |Dv|^2 \right)$$

(NB:  $C^1$  obviously enough for  $\nu = 0$  (Euler). Less obvious when  $\nu > 0$ : use for instance the regularity theory for Navier-Stokes.)

Integrate in space (and by parts!) to derive the dissipation law for the kinetic energy:

$$\frac{d}{dt} \int |v|^2(x, t) dx = -\nu \int |Dv|^2(x, t) dx \quad (1)$$

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# Weak solutions of the Euler equations

Three possible definitions of **generalized solutions**:

1: use the theory of distributions to define derivatives. Assume square summability of  $v$  ( $v \in L^2$ ) to safely define  $v \otimes v$ .

2: use Fourier series (periodic setting) in space and reduce the PDE to an (**infinite-dimensional**) system of ODEs for the Fourier coefficients. The minimal assumption to give a meaning:  $v(t, \cdot) \in L^2$  (with some uniformity in  $t$ ).

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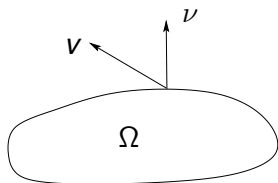
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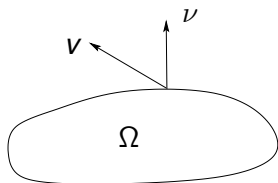


Conclusion: 
$$\int_{\partial\Omega} v \cdot \nu = 0$$

Balancing the momentum: 
$$\frac{d}{dt} \int_{\Omega} v = \int_{\partial\Omega} v(v \cdot \nu) + \int_{\partial\Omega} p\nu$$

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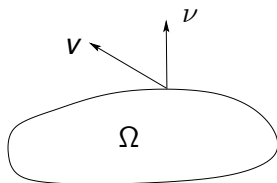


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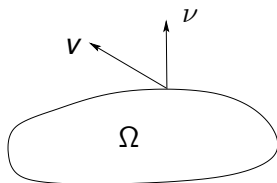


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# The Scheffer-Shnirelman nonuniqueness

All these notions are equivalent and from now on: **weak = generalized**.

Theorem (Scheffer 1993)

*There are compactly supported nonzero weak solutions in  $L^2(\mathbb{R}^2 \times \mathbb{R})$ .*

A different proof in the periodic setting given by Shnirelman in 1998.

Obviously these solutions do not preserve the total kinetic energy.

Theorem (Shnirelman 2000)

*There are weak solutions in 3-space dimension with total kinetic energy which is **strictly decreasing**.*

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# Bounded "bad" weak solutions

## Theorem (D-Székelyhidi 2007)

*There are compactly supported nontrivial **bounded** weak solutions in any space dimension.*

D-S 2009: There are weak solutions which “behave in all possible ways” in terms of local/global energy conservation. **None of the criteria proposed so far in the literature restores uniqueness of weak solutions to the Cauchy problem if the initial data are discontinuous.**

Székelyhidi 2011: **this remains true even for fairly mild discontinuities.**

Similar conclusions hold for other equations of fluid dynamics, where analogous methods can be used: Cordoba-Faraco-Gancedo, Shvidkoy, Wiedemann, Chiodaroli, ...

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# Differential inclusions

Our 2007 paper plunged Scheffer's nonuniqueness Theorem in a long tradition of counterintuitive examples in [differential inclusions](#) and in [differential geometry](#).

In the theory of differential inclusions you are looking at problems of the following type.

## Problem

Given a set  $K$  of  $k \times n$  matrices study maps  $u : \mathbb{R}^n$  (or  $\Omega \subset \mathbb{R}^n$ )  $\rightarrow \mathbb{R}^k$  such that

$$\nabla u(x) \in K \quad \text{for all } x \in \Omega. \quad (2)$$

It happens in several interesting situations that  $C^1$  solutions are not so interesting because they are forced to be [affine](#). In these cases we can look at [Lipschitz](#) solutions (which are differentiable a.e.!) and we turn (2) into

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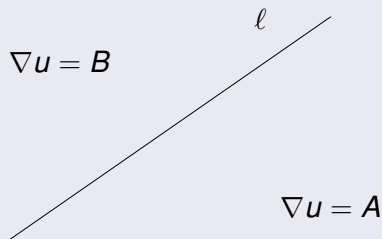
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Let us look at two “cousins” of the D-S Theorem.

## Exercise

Consider two  $2 \times 2$  matrices  $A$  and  $B$ : is there a Lipschitz planar map  $u : \mathbb{R}^2 \rightarrow \mathbb{R}^2$  with  $\nabla u = A$  “on the left” and  $\nabla u = B$  “on the right”?



Solution: It exists if and only if the direction of  $\ell$  is in the kernel of  $A - B$ .

However...

Theorem (Kirchheim 2003)

*There are  $2 \times 2$  matrices  $A_1, A_2, \dots, A_5$  and a Lipschitz map  $u : \mathbb{R}^2 \rightarrow \mathbb{R}^2$  such that*

- ▶  $\text{rank}(A_i - A_j) = 2$ ;
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Another cousin: [Müller - Šverak 1999] there are solutions of nonlinear strongly elliptic systems of PDEs which are **nowhere differentiable**.

Connection between differential inclusions and elliptic systems??  
I.e.: the Cauchy-Riemann equations are a differential inclusion!

The techniques used by Kirchheim and Müller-Šverak have a long tradition: Cellina, Bressan, Bressan-Flores, Dacorogna-Marcellini, Sychev, Székelyhidi (Tartar, DiPerna).

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# A toy example

$\Omega \subset \mathbb{R}^2$  smooth bounded open set. We look for planar (Lipschitz!) real-valued maps  $\alpha : \Omega \rightarrow \mathbb{R}$  such that

$$|\nabla\alpha| = 1 \quad (4)$$

(+ maybe some boundary conditions...).

**PLAN:** Start from some smooth map  $\varphi_0$  with  $|\nabla\varphi_0| < 1$ .

Set up an iteration scheme producing  $\varphi_0 \rightarrow \varphi_1 \rightarrow \varphi_2 \rightarrow \dots$

such that

▶  $|\nabla\varphi_k| < 1$ ;

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$$\int_{\Omega} (1 - |\nabla\varphi_k|) \leq \beta \int_{\Omega} (1 - |\nabla\varphi_{k-1}|)$$

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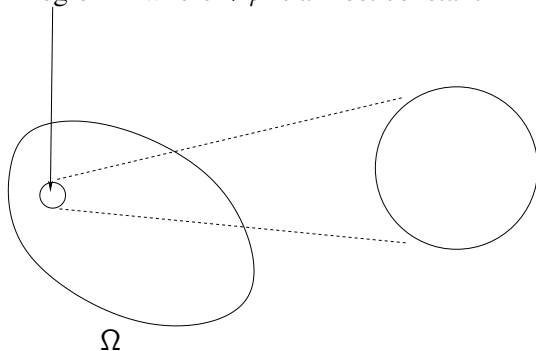
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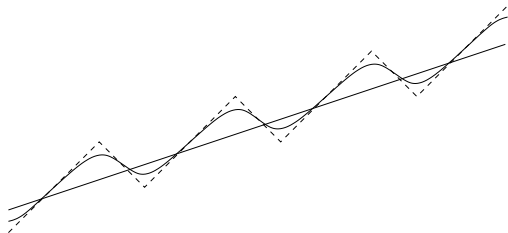
The iteration: from  $\varphi_k = \varphi$  to  $\varphi_{k+1} = \psi$

A region  $R$  where  $\nabla\varphi$  is almost constant



## A toy example III

We make the slope of  $\varphi$  “steeper in average” in  $R$  by adding a periodic function which oscillates rapidly (in the direction of  $\nabla\varphi$ ): we see below a cross section of  $\varphi$  and of the perturbed function  $x \mapsto \varphi(x) + \frac{1}{\lambda}p(\lambda x)$



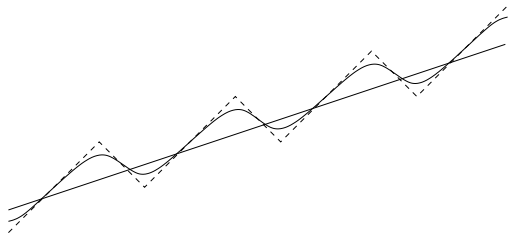
Next, **cut off** the perturbation to make it compactly supported in the region  $R$ :

$$\psi(x) = \varphi(x) + \frac{1}{\lambda}p(\lambda x)c(x)$$

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We are now ready for the key computation:

$$\nabla\psi(x) = \nabla\varphi(x) + \underbrace{c(x)\nabla\rho(\lambda x)}_{\text{Improvement}} + \underbrace{\frac{1}{\lambda}\rho(\lambda x)\nabla c(x)}_{\text{Error}}$$

The Improvement “pushes” the slope towards 1 (at least in most places!).

The error can be made as small as we wish if  $\lambda$  is very large: this is not destroying what we gained with the Improvement.

Take care, do not get immediately to slope 1 (or above!) with the Improvement: for  $\lambda$  large we will keep the inequality  $|\nabla\psi| < 1$ .



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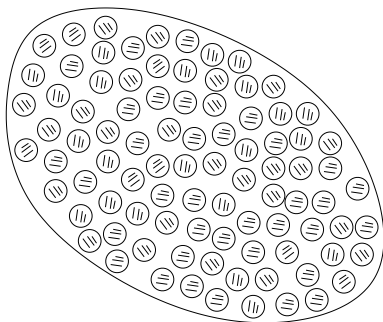
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# A toy example V

Repeat now this in many many small balls which cover a substantial portion of the region where  $|\nabla\varphi|$  is “far” from 1.



$\Omega$

# From analysis to geometry

The upshot is: in all these “crazy” constructions the final (more or less counterintuitive) map is achieved through the addition of very fine oscillations to some underlying “subsolution”: the oscillations “pile up” and we reach the desired map only after infinitely many steps.

The Müller-Šverak paper is a landmark result also because the authors realized that similar ideas had already been used in geometry.

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An older tradition of counterintuitive construction indeed exists in differential geometry (Nash-Kuiper, Smale's paradox, Gromov, Eliashberg, ...).

Rather than trying to introduce the  $h$ -principle let me give an example, (maybe the "mother" of all these constructions?).

Consider a (smooth) Riemannian manifold  $(M, g)$ : **an isometry  $u : M \rightarrow \mathbb{R}^N$  is a map preserving the length of curves.**

In what follows we deal with  $C^1$  maps which are also embeddings: **isometric embeddings.**

## Corollary

*Consider the standard sphere  $(\mathbb{S}^2, \sigma)$  or the flat square  $([0, 1]^2, f)$ . For any given  $\varepsilon > 0$  there are  $C^1$  isometric embeddings of these manifolds in a euclidean three-dimensional ball of radius  $\varepsilon$ ,  $B_\varepsilon(0) \subset \mathbb{R}^3$ !*

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Indeed the Theorem of Nash-Kuiper is much more general and much more precise: every **short embedding** (i.e. which shrinks the length of curves) of a compact Riemannian manifold can be uniformly approximated by  $C^1$  isometric embeddings.

In the framework introduced by Gromov this can be translated into a “ $C^0$ -dense  $h$ -principle” (combining Nash-Kuiper with the Hirsch-Smale  $h$ -principle).

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Something like that holds for all the results mentioned in the theory of differential inclusions... **With a big caveat:**

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Let us go back to the Navier-Stokes equations and assume that the “viscosity” is small:

$$\partial_t v + \operatorname{div} (v \otimes v) + \nabla p = \varepsilon \Delta v$$

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The Kolmogorov's theory of fully developed turbulence (K41) predicts that for most solutions  $Q(t)$  is independent of  $\varepsilon$ . More precisely

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## Conjecture (Onsager 1949)

(A) Assume  $v$  is a (periodic) weak solution of the Euler equations satisfying an Hölder condition with exponent  $\alpha > \frac{1}{3}$ :

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Then the total kinetic energy of  $v$  is conserved:

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(B) Let  $\alpha < \frac{1}{3}$ . Then there are weak solutions satisfying the Hölder condition with exponent  $\alpha$  such that

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The solutions considered by Onsager are not differentiable... A physicist (with a Nobel prize in chemistry) considers indeed weak solutions in our modern sense in 1949.

And he gave a very rigorous definition, following “road 2” from some slides ago, i.e. via a Fourier series expansion.

Part (A) of the Conjecture has been proved by Eyink and Constantin-E-Titi in 1993. (see also Cheskidov-Constantin-Friedlander-Shvidkoy 2008)

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# Continuous dissipative solutions

## Theorem (D-Székelyhidi 2011)

Let  $I$  be a compact interval and  $e : I \rightarrow \mathbb{R}$  any given smooth positive function. Then there is a continuous solution  $(v, p)$  of the Euler equations in  $\mathbb{T}^3 \times I$  such that

$$\int |v|^2(x, t) dx = e(t) \quad \forall t \in I.$$

To our knowledge:

- ▶ This is the first example of its kind in the “analysis of PDEs” which lies in the “ $C^0$  category” (which is instead customary in geometry);
- ▶ It relies upon “more complicated elliptic operators” ( $\Delta$  and  $\Delta^2$ ) whereas all other “convex integration” constructions rely on “integrating”  $\frac{d^k}{dr^k}$  (cp. with a question of Gromov).

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[Borisov 1963, 2004] The Nash-Kuiper Theorem for isometric embeddings holds for  $C^{1,\alpha}$  maps if  $\alpha$  is sufficiently small (recall:  $C^{1,\alpha}$  means  $C^1$  + the differential of the map is  $\alpha$ -Hölder). For instance

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*There is an  $\alpha_0 > 0$  with the following property. Consider the standard sphere  $(\mathbb{S}^2, \sigma)$ . For any given  $\varepsilon$  and any  $\alpha < \alpha_0$  there are  $C^{1,\alpha}$  isometric embeddings of these manifolds in a euclidean three-dimensional ball of radius  $\varepsilon$ ,  $B_\varepsilon(0) \subset \mathbb{R}^3$ .*

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There is therefore a striking analogy with the Onsager’s conjecture:

- ▶ Large Hölder exponents  $\Rightarrow$  “rigidity”.
- ▶ Small Hölder exponents allows for “flexibility” of the solutions.

[Conti-D-Székelyhidi 2009] It is possible to give a much shorter proof of Borisov’s Theorem exploiting the key computations of the Constantin-E-Titi proof of the “rigidity part” of Onsager’s conjecture.

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