\textbf{A-free Rigidity and Applications to the Compressible Euler System}

Elisabetta Chiodaroli, Eduard Feireisl, Ondřej Kreml, and Emil Wiedemann

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\textbf{Abstract}

Can every measure-valued solution to the compressible Euler equations be approximated by a sequence of weak solutions? We prove that the answer is negative: Generalizing a well-known rigidity result of Ball and James to a more general situation, we construct an explicit measure-valued solution for the compressible Euler equations which can not be generated by a sequence of distributional solutions. We also give an abstract necessary condition for measure-valued solutions to be generated by weak solutions, relying on work of Fonseca and Müller.

While a priori it is not unexpected that not every measure-valued solution arises from a sequence of weak solutions, it is noteworthy that this observation in the compressible case is in contrast with the incompressible situation, where every measure-valued solution can be approximated by weak solutions, as shown by Székelyhidi and Wiedemann.

1 Introduction

The compressible Euler equations

\begin{align*}
\partial_t \rho + \text{div}_x (\rho v) &= 0 \\
\partial_t (\rho v) + \text{div}_x (\rho v \otimes v) + \nabla_x p(\rho) &= 0
\end{align*}

\text{(1.1)}

describe the motion of a perfect fluid whose state is determined by the distribution of the fluid velocity \(v\) and of the mass density \(\rho\). The pressure \(p\) is given as a function of the density and satisfies \(p(\rho) \geq 0, \ p'(\rho) > 0\). These equations were first formulated by Euler in 1757, but we are still far from a complete understanding of the phenomena described by the Euler equations.

If we consider the Cauchy problem for the compressible Euler equations (1.1) on \(Q = [0,T] \times \Omega\) (for a, say, bounded smooth domain \(\Omega \subset \mathbb{R}^3\)) with bounded initial data

\((\rho, v)(0, \cdot) = (\rho_0, v_0)\)

then we say that a pair \((\rho, v) \in L^1([0,T] \times \Omega)\) is a weak solution to (1.1) (together with a natural boundary condition \(v \cdot n = 0\) on \(\partial \Omega\), \(n\) being the unit outer normal to \(\partial \Omega\)) if it
satisfies the following equations:

\[
\int_0^T \int_\Omega \partial_t \psi \rho + \nabla_x \psi \cdot \rho v \, dx \, dt + \int_\Omega \psi(0, x) \rho_0(x) \, dx = 0,
\]

\[
\int_0^T \int_\Omega \partial_t \phi \cdot \rho v + \nabla_x \phi : \rho v \otimes v + \text{div}_x \phi \rho \, dx \, dt + \int_\Omega \phi(0, x) \cdot \rho_0(x) v_0(x) \, dx = 0
\]

for all \( \psi \in C^\infty_c([0, T] \times \Omega) \) and all \( \phi \in C^\infty_c([0, T] \times \Omega; \mathbb{R}^3) \) (in particular, it is part of the definition that all the integrals exist, which is certainly the case for solutions with finite energy). The concept of weak solutions allows to deal with shock formation which may occur even for smooth data; however weak solutions are well known to be non-unique. In the literature many admissibility criteria have been proposed to restore uniqueness and the most classical one relies on suitable entropy conditions.

Entropy solutions have been widely accepted as the suitable solution framework for systems of conservation laws in several space dimensions, for which the compressible Euler equations (1.1) are the most paradigmatic example. However, in space dimensions greater than one the existence of entropy weak solutions is not known, so that the weaker concept of measure-valued solution was introduced by Neustupa [19], who also proved the existence of such solutions. Previously, the notion of measure-valued solutions for systems of conservation laws had been introduced by DiPerna in the seminal paper [7].

Measure-valued solutions are only parametrized measures (often also called Young measures), instead of integrable functions. A parametrized measure is a family

\[\{\nu_{t,x}\}_{(t,x) \in [0,T] \times \Omega}\]

of probability measures on the configuration space \( \mathbb{R}^+ \times \mathbb{R}^3 \), and it is called a measure-valued solution of (1.1) if it satisfies the equations in an average sense (see Definition 1 below). Let us emphasize that, although measure-valued solutions are typically conceived as arising from a sequence of approximate solutions, it is not at all clear if every measure-valued solution in the sense of Definition 1 is generated by such a sequence of (approximate) weak solutions. In fact it is the main result of this paper that, for the compressible Euler system, this is not the case.

On the one hand, parametrized measures allow to characterize complex phenomena such as oscillations and concentrations. On the other hand, measure-valued solutions have been criticized for being too weak, as is apparent from their obvious non-uniqueness (but see the weak-strong uniqueness results in [2, 14]). However, recent results by the first three authors jointly with De Lellis [3, 4, 5] have demonstrated that for the compressible Euler equations (1.1) even entropy solutions (weak solutions satisfying a suitable entropy inequality) may not be unique, thus raising anew the problem of a correct notion of solutions for (1.1).

In the recent paper [11], the authors show that current numerical schemes may not necessarily converge to an entropy solution of systems of conservation laws in several space dimensions, when the mesh is refined. Instead, entropy measure-valued solutions
are suggested as an appropriate solution paradigm for systems of conservation laws, on
the base of a large number of numerical experiments.

In the context of hyperbolic conservation laws in one space dimension, when the equa-
tion is complemented with a suitable entropy condition, the Tartar-Murat theory of com-
pensated compactness applies (see [8]) and oscillations in weakly converging sequences of
solutions can be excluded. The situation is very different for the incompressible Euler
equations in several space dimensions: In [21], Székelyhidi and the last author showed
that any measure-valued solution can be generated by a sequence of exact solutions. This
means that, in any space dimension, measure-valued solutions and weak solutions are
substantially the same, thus leading to a very large set of weak solutions.

Motivated by this series of results, we aimed at understanding the corresponding prob-
lem for the compressible Euler equations (1.1) (we treat here only the three-dimensional
case). This issue may prove interesting both in view of the solution theory of the com-
pressible Euler system (1.1) itself and for possible selection criteria among measure-valued
solutions of incompressible Euler based on low Mach number limits. In this paper we prove
(Theorem 3) that any measure-valued solution to the compressible Euler system which
arises as a limit of a sequence of weak solutions must satisfy a Jensen-type inequality,
which is reminiscent of analogous inequalities in the so called $A^*$-free setting (cf. [13] and
[15, 16]). To our knowledge this is the first time that such a Jensen inequality appears in
the context of fluid dynamics.

The example of the incompressible Euler equations [21] shows that in general it cannot
be excluded that the Jensen inequality is vacuous, i.e. it is automatically satisfied and thus
poses no actual constraint. In order to directly show that the condition is non-empty, one
would have to exhibit a measure-valued solution and an $A$-quasiconvex function violating
the Jensen condition. Due to the abstract and ill-understood nature of quasiconvexity,
this appears difficult.

Instead, we follow a different strategy in order to construct an explicit measure-valued
solution for the compressible Euler equations which can not be generated by a sequence
of distributional solutions (Theorem 4): our construction relies on a rigidity result for so-
called $A^*$-free sequences (Theorem 2 below) in the spirit of compensated compactness [22,
18], which generalizes a well-known result of Ball and James (Proposition 2 in [1], see
also Lemma 2.2 in [17]) on sequences of gradients. Rigidity results of this type play an
important role in nonlinear elasticity, but have not yet been used in fluid mechanics.
Our generalization to the $A^*$-free setting will not come as a surprise to experts, but we
have not been able to find it explicitly stated in the literature. Under the constant rank
assumption, our Theorem 2 follows from Theorem 4 in [20], but we do not require this
assumption and our proof is simple and self-contained. We thus hope it will be found
useful independently of our specific application in the context of the Euler equations.

The substantial difference between weak and measure-valued solutions for the com-
pressible Euler equations (1.1) in contrast with the situation for incompressible Euler [21]
indicates that the different role of the pressure in the compressible and incompressible
theory has an important effect and confers additional rigidity to weak solutions of the
compressible system.

Finally, let us briefly comment on the converse of Theorem 3: is the Jensen inequality
sufficient for a measure-valued solution to be generated by weak solutions? In the incom-
pressible case, this was shown to be true [21], but the proof relies on convex integration
techniques which do not seem to be flexible enough to transfer to the compressible case.

The paper is organized as follows: Sections 2, 3.1 and 3.2 contain preliminary material
on measure-valued solutions and linear differential constraints. In Section 3.3 we state
and prove our \( A \)-free rigidity theorem. Section 4 then gives the mentioned applications
of the general theory to the compressible Euler equations.

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2 Measure-valued solutions and subsolutions

We denote by \( \mathcal{M}(\mathbb{R}^d) \) the space of finite Radon measures on \( \mathbb{R}^d \) and by \( \mathcal{M}^1(\mathbb{R}^d) \) the
subset of probability measures. For a subset \( \Omega \subset \mathbb{R}^N \), we also denote by \( L^\infty(\Omega; \mathcal{M}^1(\mathbb{R}^d)) \)
the space of weakly*-measurable maps \( \nu \) from \( \Omega \) into \( \mathcal{M}^1(\mathbb{R}^d) \), meaning that for each
Borel function \( f : \mathbb{R}^d \rightarrow \mathbb{R} \), the map 

\[
x \mapsto \langle \nu_x, f \rangle := \int_{\mathbb{R}^d} f(z) d\nu_x(z)
\]

is measurable (with respect to Lebesgue measure).

A \textit{parametrized measure} or \textit{Young measure} on a set \( \Omega \subset \mathbb{R}^N \) is a map \( \nu \in L^\infty(\Omega; \mathcal{M}^1(\mathbb{R}^d)) \)
which assigns to almost every point \( x \in \Omega \) a probability measure \( \nu_x \in \mathcal{M}^1(\mathbb{R}^d) \) on the
phase space \( \mathbb{R}^d \). We note in passing that \( \langle \nu_x, f \rangle := \int_{\mathbb{R}^d} f(z) d\nu_x(z) \) is the expectation of \( f \)
with respect to the probability measure \( \nu_x \).

In the context of measure-valued solutions to evolution equations, one of the variables
is interpreted as time and the domain takes the form \( \Omega = [0, T] \times \bar{\Omega} \) for some space domain
\( \bar{\Omega} \) (we will omit the tilde as long as no confusion is to be expected). We assume in this
work that the space domain is bounded. In particular, for the three-dimensional Euler
equations, we have \( N = 1 + 3 \) and \( d = 4 \). We will then denote the state variables by
\( \xi \in \mathbb{R}^4 \) and more precisely we introduce the following notation in order to formalize the
definition of measure-valued solution to (1.1):

\[
\xi = [\xi_0, \xi_1] = [\xi_0, \xi_1, \xi_2, \xi_3] \in \mathbb{R}^+ \times \mathbb{R}^3
\]

\[
\langle \nu_{t,x}, \xi_0 \rangle = \int_{\mathbb{R}^4} \xi_0 d\nu_{t,x}(\xi) = \bar{\rho}
\]

\[
\langle \nu_{t,x}, \sqrt{\xi_0} \xi' \rangle = \int_{\mathbb{R}^4} \sqrt{\xi_0} \xi' d\nu_{t,x}(\xi) = \bar{\rho} \bar{v}
\]

\[
\langle \nu_{t,x}, \xi' \otimes \xi' \rangle = \int_{\mathbb{R}^4} \xi' \otimes \xi' d\nu_{t,x}(\xi) = \bar{\rho} \bar{v} \otimes \bar{v}
\]

\[
\langle \nu_{t,x}, p(\xi_0) \rangle = \int_{\mathbb{R}^4} p(\xi_0) d\nu_{t,x}(\xi) = \bar{p}(\bar{\rho}).
\]
One should think of $\xi_0$ as the state of the density $\rho$ and $\xi'$ as the state of $\sqrt{\rho}v$.

**Definition 1 (Measure-valued solution)** A measure-valued solution to the compressible Euler equations (1.1) is a Young measure $\nu_{t,x}$ on $\mathbb{R}^+ \times \mathbb{R}^3$ with parameters in $[0,T] \times \Omega$ which satisfies the Euler equations in an average sense, i.e.

$$
\int_0^T \int_\Omega \partial_t \psi \rho + \nabla_x \psi \cdot \bar{\rho}v \, dx \, dt + \int_\Omega \psi(0,x) \rho_0(x) \, dx = 0
$$

$$
\int_0^T \int_\Omega \partial_t \varphi \cdot \bar{\rho}v + \nabla_x \varphi \cdot \bar{\rho}v \otimes \bar{v} + \text{div}_x \bar{\varphi} \rho \, dx \, dt + \int_\Omega \varphi(0,x) \cdot \rho_0(x) v_0(x) \, dx = 0
$$

for all $\psi \in C^\infty_c([0,T) \times \Omega)$ and all $\varphi \in C^\infty_c([0,T) \times \Omega; \mathbb{R}^3)$.

Again, it is understood that all the integrals have to be well-defined and finite. Note that this definition involves only the “classical” Young measure and ignores effects of concentration (confer the notion of “generalized” Young measure as established in [9]). This is mainly for reasons of simplicity, and because our counterexample (Theorem 4) excludes an approximating sequence even if it forms concentrations (cf. Remark 5 below).

Observe that every measurable function $u : \Omega \to \mathbb{R}^d$ naturally gives rise to a Young measure by defining $\nu_x := \delta_{u(x)}$: such Young measures are called atomic. Thus every weak solution of (1.1) defines an atomic measure-valued solution.

If $\{z_n\}$ is a sequence of measurable functions $\Omega \to \mathbb{R}^d$ we say that $\{z_n\}$ generates the Young measure $\nu$ if for all bounded Carathéodory functions $f : \Omega \times \mathbb{R}^d \to \mathbb{R}$

$$
\lim_{n \to \infty} \int_\Omega f(x,z_n(x)) \varphi(x) \, dx = \int_\Omega \langle \nu_x, f(x,\cdot) \rangle \varphi(x) \, dx
$$

for all $\varphi \in L^1(\Omega)$. By the fundamental theorem of Young measures any sequence of functions bounded in $L^p(\Omega)$ (for any $p \geq 1$) generates, up to a subsequence, some Young measure.

Our first result gives a necessity condition, in the form of a Jensen-type inequality, for a measure-valued solution to (1.1) in order to be generated by a sequence of weak solutions. In order to formulate this Jensen inequality and to apply the linear $\mathcal{A}$-free theory of Fonseca and Müller [13] to the Euler equations, we need to relax our problem to a linear setting. Thus we introduce a suitable notion of subsolutions for our system (1.1). More precisely we deal with solutions to the related linear system

$$
\begin{align*}
\partial_t \rho + \text{div}_x m &= 0 \\
\partial_t m + \text{div}_x U + \nabla_x q &= 0,
\end{align*}
$$

which we call subsolutions. The linear system (2.1) is obtained from (1.1) by first changing the state variable from $(\rho, v)$ to $(\rho, m := \rho v)$, where $m$ is the linear momentum, and then by replacing every nonlinearity appearing in the original compressible Euler equations
(1.1) by a new variable: in particular \( U \in S_0^3 \) is a symmetric trace-free \( 3 \times 3 \) matrix which replaces the traceless part of the matrix \( \rho v \otimes v = \frac{m \otimes m}{\rho} \), i.e.

\[
U = \frac{m \otimes m}{\rho} - \frac{|m|^2}{3\rho} \mathbf{I},
\]

and \( q \) encodes the pressure plus the term coming from the trace of \( \rho v \otimes v \), more explicitly

\[
q = p(\rho) + \frac{|m|^2}{3\rho}
\]

(see [6, 4] for similar definitions). Weak solutions to (2.1) are functions \((\rho, m, U, q)\) which satisfy (2.1) in the sense of distributions.

Accordingly we have to define measure-valued subsolutions. Similarly as above, we use the following notation:

\[
\begin{align*}
[\xi_0, \xi', \mathbf{Z}, \tilde{\xi}] & \in \mathbb{R}^+ \times \mathbb{R}^3 \times S_0^3 \times \mathbb{R}^+ \\
\langle \mu_{t,x}, \xi_0 \rangle &= \bar{p} \\
\langle \mu_{t,x}, \xi' \rangle &= \bar{m} \\
\langle \mu_{t,x}, \mathbf{Z} \rangle &= \bar{U} \\
\langle \mu_{t,x}, \tilde{\xi} \rangle &= \bar{q}
\end{align*}
\]

**Definition 2 (Measure valued subsolution)** A measure-valued solution to the linear system (2.1) is a Young measure \( \mu_{t,x} \) on \( \mathbb{R}^+ \times \mathbb{R}^3 \times S_0^3 \times \mathbb{R}^+ \) with parameters in \([0,T] \times \Omega\) which satisfies the linear system (2.1) in an average sense, i.e.

\[
\begin{align*}
\int_0^T \int_{\Omega} \partial_t \psi \bar{p} + \nabla_x \psi \cdot \bar{m} \, dx \, dt + \int_{\Omega} \psi(0, x) \rho_0(x) \, dx &= 0 \\
\int_0^T \int_{\Omega} \partial_t \varphi - \nabla_x \varphi \cdot \bar{U} + \text{div}_x \varphi \bar{q} \, dx \, dt \\
&+ \int_{\Omega} \varphi(0, x) \cdot m_0(x) \, dx = 0
\end{align*}
\]

for all \( \psi \in C_c^\infty([0,T] \times \Omega) \) and all \( \varphi \in C_c^\infty([0,T] \times \Omega; \mathbb{R}^3) \).

We also define, in analogy with [21], the lift of a Young measure from the space of solutions to the space of subsolutions. Indeed once we are given a measure-valued solution to the compressible Euler equations (1.1) this defines naturally a corresponding measure–valued subsolution to the linear system (2.1).

**Definition 3 (Lift)** Let \( \nu_{t,x} \) be a measure valued solution to the Euler equations. Denote \( Q : \mathbb{R}^+ \times \mathbb{R}^3 \mapsto \mathbb{R}^+ \times \mathbb{R}^3 \times S_0^3 \times \mathbb{R}^+ \)

\[
Q(\xi) := (\xi_0, \sqrt{\xi_0} \xi', \xi' \otimes \xi' - \frac{1}{3} |\xi'|^2 \mathbf{I}, p(\xi_0) + \frac{1}{3} |\xi'|^2)
\]

We define the lifted measure \( \tilde{\nu}_{t,x} \) as

\[
\langle \tilde{\nu}_{t,x}, f \rangle := \langle \nu_{t,x}, f \circ Q \rangle
\]

for \( f \in C_0(\mathbb{R}^+ \times \mathbb{R}^3 \times S_0^3 \times \mathbb{R}^+) \) and a.e. \((t, x) \in [0,T] \times \Omega\).
Remark 1 Let us note that the lifted measure $\tilde{\nu}_{t,x}$ as introduced in Definition 3 is, in fact, a measure valued solution to the linear system (2.1), or, equivalently, it is a measure valued subsolution.

3 Linearity constrained sequences

3.1 $\mathcal{A}$–free setting and constant rank property

As detailed in Section 4.1, the linear system (2.1) associated to the compressible Euler equations fits nicely into the so-called $\mathcal{A}$-free framework for linear partial differential constraints, introduced by Tartar [22]. Consider a general linear system of $l$ differential equations in $\mathbb{R}^N$ written as

\[ \mathcal{A}z := \sum_{i=1}^{N} A^{(i)} \frac{\partial z}{\partial x_i} = 0, \]

where $A^{(i)} (i = 1, ..., N)$ are $l \times d$ matrices and $z : \mathbb{R}^N \to \mathbb{R}^d$ is a vector-valued function.

Next, we define the $l \times d$ matrix $A(w) := \sum_{i=1}^{N} w_i A^{(i)}$ for $w \in \mathbb{R}^N$.

Definition 4 (Constant rank) We say that $\mathcal{A}$ has the constant rank property if there exists $r \in \mathbb{N}$ such that

\[ \text{rank} A(w) = r \]

for all $w \in S^{N-1}$.

The following theorem from [13] will be the cornerstone of our proof of Theorem 3 below. In order to state it, we first recall the general definition of $\mathcal{A}$-quasiconvex functions.

Definition 5 ($\mathcal{A}$-Quasiconvexity) A function $f : \mathbb{R}^d \to \mathbb{R}$ is said to be $\mathcal{A}$-quasiconvex if

\[ f(z) \leq \int_{(0,1)^N} f(z + w(x)) dx \]

for all $z \in \mathbb{R}^d$ and all $w \in C^\infty_{\text{per}}((0,1)^N; \mathbb{R}^d)$ such that $Aw = 0$ and $\int_{(0,1)^N} w(x) dx = 0$.

Recall that a sequence $\{z_n\}$ is called $p$-equiintegrable if the sequence $\{|z_n|^p\}$ is equiintegrable in the usual sense.

Theorem 1 (Theorem 4.1 in [13]) Let $1 \leq p < \infty$ and let $\{\nu_x\}_{x \in \Omega}$ be a weakly measurable family of probability measures on $\mathbb{R}^d$. Let $\mathcal{A}$ have the constant rank property. There exists a $p$-equiintegrable sequence $\{z_n\}$ in $L^p(\Omega; \mathbb{R}^d)$ that generates the Young measure $\nu$ and satisfies $\mathcal{A}z_n = 0$ in $\Omega$ if and only if the following conditions hold:
(i) there exists \( z \in L^p(\Omega; \mathbb{R}^d) \) such that \( A z = 0 \) and \( z(x) = \langle \nu_x, \text{id} \rangle \) a.e. \( x \in \Omega \);

(ii) \[
\int_{\Omega} \int_{\mathbb{R}^d} |w|^p \, d\nu_x(w) \, dx < \infty;
\]

(iii) for a.e. \( x \in \Omega \) and all \( A \)-quasiconvex functions \( g \) that satisfy \( |g(w)| \leq C(1 + |w|^p) \) for some \( C > 0 \) and all \( w \in \mathbb{R}^d \) one has

\[
(3.3) \quad \langle \nu_x, g \rangle \geq g(\langle \nu_x, \text{id} \rangle).
\]

### 3.2 Wave analysis

When considering first order linear PDEs in the form (3.1), a special class of solutions are plane waves, i.e. solutions of the form \( z(x) = h(x \cdot \xi) \bar{z} \) with \( h : \mathbb{R} \to \mathbb{R} \). In order to identify such solutions, one has to solve the relation \( \sum_{i=1}^N \xi_i A^{(i)} \bar{z} = 0 \) which gives rise to the following definition.

**Definition 6** Consider a linear differential operator \( A \) as in (3.1). Its wave cone \( \Lambda \) is defined as the set of all \( \bar{z} \in \mathbb{R}^d \setminus \{0\} \) for which there exists \( \xi \in \mathbb{R}^N \setminus \{0\} \) such that

\[
z(x) = h(x \cdot \xi) \bar{z}
\]

satisfies \( A z = 0 \) for any choice of profile function \( h : \mathbb{R} \to \mathbb{R} \).

Equivalently, \( \bar{z} \in \Lambda \) if and only if \( \bar{z} \neq 0 \) and there exists \( \xi \in \mathbb{R}^N \setminus \{0\} \) such that \( A(\xi) \bar{z} = 0 \).

The claimed equivalence can be easily verified by taking the Fourier transform of \( A z = 0 \). It can be seen immediately that \( \Lambda \) is relatively closed in \( \mathbb{R}^d \setminus \{0\} \). In other words, the wave cone \( \Lambda \) characterizes the directions of one dimensional oscillations compatible with (3.1).

Let us give an equivalent reformulation of the constraint (3.1) and accordingly yet another characterization of the wave cone. Observe that

\[
\sum_{i=1}^N A^{(i)} \frac{\partial z}{\partial x_i} = \left( \sum_{i=1}^N \sum_{k=1}^d A_{jk}^{(i)} \frac{\partial z_k}{\partial x_i} \right)_{j=1,\ldots,l} = \left( \sum_{i=1}^N \frac{\partial}{\partial x_i} \sum_{k=1}^d A_{jk}^{(i)} z_k \right)_{j=1,\ldots,l}.
\]

Therefore, if we define the \( l \times N \)-matrix \( Z_A \) by

\[
(Z_A)_{ji} = \sum_{k=1}^d A_{jk}^{(i)} z_k, \quad j = 1, \ldots, l, \quad i = 1, \ldots, N,
\]

then (3.1) can be rewritten as

\[
\text{div} Z_A = 0.
\]

Moreover, the condition \( A(\xi) \bar{z} = 0 \) from the definition of the wave cone translates to \( \bar{Z}_A \xi = 0 \) (where \( \bar{Z}_A \) is obtained from \( \bar{z} \) via (3.4)), so that the following are equivalent:
1. $\bar{z} \in \Lambda$;
2. $\bar{z} \neq 0$ and rank $\bar{Z}_A < N$.

It follows immediately that $\Lambda = \mathbb{R}^d \setminus \{0\}$ if $l < N$.

### 3.3 $A$-free rigidity

In this section we prove a generalization of the well-known rigidity result of Ball and James for sequences of gradients [1, Proposition 2] to the $A$-free framework. Besides its possible independent interest, its application to the system (3.1) will guarantee a form of compensated compactness that allows us to construct a measure-valued solution to the compressible Euler system which is not generated by a sequence of weak solutions.

**Theorem 2** Let $\Omega \subset \mathbb{R}^N$ be a Lipschitz and bounded domain, $A$ a linear operator of the form (3.1), and $1 < p < \infty$. Let $\bar{z}_1, \bar{z}_2 \in \mathbb{R}^d$, $\bar{z}_1 \neq \bar{z}_2$ be two constant states such that $\bar{z}_2 - \bar{z}_1 \not\in \Lambda$.

Then if $z_n : \Omega \rightarrow \mathbb{R}^d$ is a family of functions with

$$\|z_n\|_{L^p(\Omega; \mathbb{R}^d)} \leq c,$$

(3.6)

$Az_n = 0$ in $\mathcal{D}'(\Omega)$,

which generates a compactly supported Young measure

$$\nu_x \in \mathcal{M}^1(\mathbb{R}^d)$$

such that

$$\text{supp}[\nu_x] \subset \{\lambda \bar{z}_1 + (1 - \lambda) \bar{z}_2, \ \lambda \in [0, 1]\} \text{ for a.a. } x \in \Omega,$$

then

$$z_n \rightarrow z_\infty \text{ in } L^p(\Omega),$$

which implies that

$$\nu_x = \delta_{z_\infty(x)}, \ z_\infty(x) \in \{\lambda \bar{z}_1 + (1 - \lambda) \bar{z}_2, \ \lambda \in [0, 1]\} \text{ for a.a. } x \in \Omega.$$

More specifically, $z_\infty$ is a constant function of the form

$$z_\infty = \bar{x} \bar{z}_1 + (1 - \bar{x}) \bar{z}_2,$$

for some fixed $\bar{x} \in [0, 1]$.

**Remark 2** Note that we are assuming neither $p$-equiintegrability of the sequence nor the constant rank property. If, however, the constant rank property holds, then hypothesis (3.6) can be relaxed to

$$Az_n \rightarrow 0 \text{ in } W^{-1,p}(\Omega)$$

using Lemma 2.14 from [13].

The rest of the Section is devoted to the proof of Theorem 2 which will be split in two steps corresponding to as many subsections.
3.3.1 The specific form of $z_n$

Let us start with a preliminary remark. Since we already know that
\[ z_n \to z_\infty \text{ weakly in } L^p(\Omega; \mathbb{R}^d), \]
it is enough to show that \( \{z_n\} \) contains a subsequence converging in \( L^p(\Omega) \).
Moreover, as the result is local, it is enough to show the claim on any set
\[ \Omega_\delta = \{x \in \Omega : \text{dist}[x, \partial \Omega] > \delta\}, \ \delta > 0. \]
In other words, we may assume, without loss of generality, that (3.6) holds in an open set containing \( \overline{\Omega} \).

Now, we aim to prove the following claim.

CLAIM: The functions \( z_n \) have the following specific form

(3.8) \[ z_n(x) = e_n(x) + \lambda_n(x)\bar{z}_1 + (1 - \lambda_n(x))\bar{z}_2, \]
where \( e_n \to 0 \) in \( L^p(\Omega; \mathbb{R}^d) \) as \( n \to \infty \), and \( \lambda_n \) are bounded measurable functions
\[ 0 \leq \lambda_n(x) \leq 1 \text{ a.e. } x \in \Omega. \]

Proof. [Proof of the CLAIM] Let \( \delta > 0 \) be an arbitrary positive constant. Then there is a function
\[ F_\delta \in C^\infty(\mathbb{R}^d), \ 0 \leq F_\delta \leq 1 \]
such that
\[ F_\delta = 0 \text{ in an open neighborhood of the segment } \{\lambda \bar{z}_1 + (1 - \lambda)\bar{z}_2, \ \lambda \in [0, 1]\} \]
\[ F_\delta(z) = 1 \text{ if dist}[z; \{\lambda \bar{z}_1 + (1 - \lambda)\bar{z}_2, \ \lambda \in [0, 1]\}] > \delta. \]
Next, we can simply rewrite \( z_n \) as
\[ z_n = F_\delta(z_n)z_n + (1 - F_\delta(z_n))z_n. \]
Of course, we have
\[ \|F_\delta(z_n)z_n\|_{L^p(\Omega; \mathbb{R}^d)} \leq \|z_n\|_{L^p(\Omega; \mathbb{R}^d)}, \]
and, by virtue of hypothesis (3.7), we also have
\[ \|F_\delta(z_n)z_n\|_{L^p(\Omega; \mathbb{R}^d)} \to 0 \text{ as } n \to \infty, \text{ for any } \delta > 0. \]
Finally, we can write
\[ (1 - F_\delta(z_n))z_n = (1 - F_\delta(z_n))(z_n - \lambda_n\bar{z}_1 - (1 - \lambda_n)\bar{z}_2) - F_\delta(z_n)(\lambda_n\bar{z}_1 + (1 - \lambda_n)\bar{z}_2) \]
for certain functions \( \lambda_n \) with

\[
0 \leq \lambda_n(x) \leq 1 \quad \text{for a.a. } x \in \Omega,
\]

where

\[
|(1 - F_\delta(z_n)) (z_n - \lambda_n \tilde{z}_1 - (1 - \lambda_n) \tilde{z}_2)| \leq \delta
\]

and

\[
\|F_\delta(z_n)) (\lambda_n \tilde{z}_1 + (1 - \lambda_n) \tilde{z}_2)\|_{L^p(\Omega; \mathbb{R}^d)} \to 0 \quad \text{as } n \to \infty
\]

for any fixed \( \delta > 0 \). (In fact \( \lambda_n \) are chosen so that \( \lambda_n \tilde{z}_1 + (1 - \lambda_n) \tilde{z}_2 \) is the projection of \( z_n \) onto the segment joining \( \tilde{z}_1 \) with \( \tilde{z}_2 \).)

As \( \delta > 0 \) can be taken arbitrarily small, (3.8) and the claim follow.

\[ \Box \]

### 3.3.2 Constraint imposed by equations

We use hypothesis (3.6) in the form (3.5) taking advantage of the specific form of \( z_n \) proved in (3.8). We employ here the map (3.4) from vectors in \( \mathbb{R}^d \) to matrices in \( \mathbb{R}^{l \times N} \) and denote by a capital letter the value of the map at the corresponding lowercase letter, e.g. \( e_n \) is mapped to \( E_n \) according to (3.4). We notice that the strong convergence of \( e_n \) implies that \( E_n \to 0 \) strongly in \( L^p \).

Hence by (3.6) we have

\[
0 = \text{div} E_n(x) + \text{div} [\lambda_n(x)(\tilde{Z}_1 - \tilde{Z}_2)] = \text{div} E_n(x) - (\tilde{Z}_2 - \tilde{Z}_1) \nabla \lambda_n(x) \quad \text{in } D'(\Omega).
\]

Now, using our preliminary remark that the equations are satisfied in an open neighborhood of \( \overline{\Omega} \), we may use a regularization \( v \mapsto (v)_\varepsilon \) by means of a suitable family of convolution kernels to obtain

\[
(\tilde{Z}_2 - \tilde{Z}_1) \nabla (\lambda_n)_\varepsilon = \text{div} (E_n)_\varepsilon \equiv \chi_{n,\varepsilon},
\]

where

\[
\|\chi_{n,\varepsilon}\|_{W^{-1,p}(\Omega; \mathbb{R}^d)} \to 0 \quad \text{as } n \to \infty \quad \text{uniformly in } \varepsilon.
\]

Next we make use of the fact that \( \tilde{z}_2 - \tilde{z}_1 \notin \Lambda \), and hence \( \tilde{Z}_2 - \tilde{Z}_1 \) has full rank (recall the discussion in Section 3.2) and thus possesses a left inverse. Then it follows from (3.9)–(3.10) that

\[
\|\nabla (\lambda_n - \langle \lambda_n \rangle)_\varepsilon\|_{W^{-1,p}(\Omega; \mathbb{R}^d)} \to 0 \quad \text{uniformly in } \varepsilon;
\]

whence

\[
\|\nabla (\lambda_n - \langle \lambda_n \rangle)\|_{W^{-1,p}(\Omega; \mathbb{R}^d)} \to 0;
\]

and finally by [10, Lemma 10.10] we can conclude that

\[
\|\lambda_n - \langle \lambda_n \rangle\|_{L^p(\Omega; \mathbb{R}^d)} \to 0 \quad \text{as } n \to \infty
\]

where

\[
\langle \lambda_n \rangle := \frac{1}{|\Omega|} \int_{\Omega} \lambda_n(x) \, dx.
\]

Going back to (3.8) we obtain the desired strong convergence of \( \{z_n\} \) to a constant function. We have proved Theorem 2.
4 Application to the compressible Euler equations

4.1 A necessary condition for generability by weak solutions

With \( t = x_0 \) the linear system (2.1) related to subsolutions of the compressible Euler equations can be rewritten in the form (3.1). More precisely, keeping in mind that \( U \) is a symmetric traceless \( 3 \times 3 \) matrix, we define the state vector

\[
z := (\rho, m_1, m_2, m_3, U_{11}, U_{12}, U_{13}, U_{22}, U_{23}, q) \in \mathbb{R}^{10}.
\]

Accordingly, the \( 4 \times 10 \) matrices \( A_{L}^{(i)} \) for \( i = 0, ..., 3 \) have the following form

\[
A_{L}^{(0)} = \begin{pmatrix}
1 & 0 & 0 & 0 & 0 & 0 & 0 & 0 & 0 & 0 \\
0 & 1 & 0 & 0 & 0 & 0 & 0 & 0 & 0 & 0 \\
0 & 0 & 1 & 0 & 0 & 0 & 0 & 0 & 0 & 0 \\
0 & 0 & 0 & 1 & 0 & 0 & 0 & 0 & 0 & 0
\end{pmatrix}
\]

\[
A_{L}^{(1)} = \begin{pmatrix}
0 & 1 & 0 & 0 & 0 & 0 & 0 & 0 & 0 & 0 \\
0 & 0 & 0 & 1 & 0 & 0 & 0 & 0 & 1 & 0 \\
0 & 0 & 0 & 0 & 1 & 0 & 0 & 0 & 0 & 0 \\
0 & 0 & 0 & 0 & 0 & 1 & 0 & 0 & 0 & 0
\end{pmatrix}
\]

\[
A_{L}^{(2)} = \begin{pmatrix}
0 & 0 & 1 & 0 & 0 & 0 & 0 & 0 & 0 & 0 \\
0 & 0 & 0 & 0 & 1 & 0 & 0 & 0 & 0 & 0 \\
0 & 0 & 0 & 0 & 0 & 1 & 0 & 0 & 0 & 0 \\
0 & 0 & 0 & 0 & 0 & 0 & 1 & 0 & 0 & 0
\end{pmatrix}
\]

\[
A_{L}^{(3)} = \begin{pmatrix}
0 & 0 & 0 & 1 & 0 & 0 & 0 & 0 & 0 & 0 \\
0 & 0 & 0 & 0 & 1 & 0 & 0 & 0 & 0 & 0 \\
0 & 0 & 0 & 0 & 0 & 0 & 1 & 0 & 0 & 0 \\
0 & 0 & 0 & 0 & 0 & 0 & 0 & 1 & 0 & 0
\end{pmatrix}
\]

With this notation, the system (2.1) takes the form \( A_{L}z := \sum_{i=0}^{3} A_{L}^{(i)} \frac{\partial z}{\partial x_i} = 0 \). We can prove the following lemma.

**Lemma 1** The operator \( A_{L} \) defined as \( A_{L}z := \sum_{i=0}^{3} A_{L}^{(i)} \frac{\partial z}{\partial x_i} \), with the choice of matrices \( A_{L}^{(i)} \) given by (4.1)–(4.4), has the constant rank property with \( r = 4 \).

**Proof.** This follows from an easy linear algebra computation. \( \square \)

The constant rank property is crucial to us since it allows to exploit the theory of Fonseca and Müller for our linear system (2.1). Recall once more that a sequence \( \{z_n\} \) is called \( p \)-equiintegrable if the sequence \( \{|z_n|^p\} \) is equiintegrable. Here, \( p \)-equi-integrability refers to both variables \( t \) and \( x \). Then, by virtue of Theorem 1 and Lemma 1, we can finally prove the following theorem, which is one of our main results:
Step 2: The sequence \{(\rho_n, v_n)\} is a sequence of weak solutions to the compressible Euler system (1.1) such that \{\rho_n\} is \(\gamma\)-equiintegrable and \{\sqrt{\rho_n} v_n\} is 2-equiintegrable. Suppose moreover \{(\rho_n, \sqrt{\rho_n} v_n)\} generates a Young measure \(\nu\) on \(\mathbb{R}^+ \times \mathbb{R}^3\). Then \(\nu\) is a measure-valued solution to the compressible Euler system (1.1) and the lifted measure \(\tilde{\nu}\) on \(\mathbb{R}^+ \times \mathbb{R}^3 \times \mathcal{S}_0^3 \times \mathbb{R}^+\) satisfies

\[
\langle \tilde{\nu}_{t,x}, g \rangle \geq g((\tilde{\nu}_{t,x}, \text{id}))
\]

for all \(\mathcal{A}_L\)-quasiconvex functions \(g\) with at most linear growth (i.e. \(|g(w)| \leq C(1 + |w|)\)).

**Remark 3**

The assumption on the pressure is very natural. Together with the equiintegrability assumptions it implies that the nonlinear terms \(\{\rho_n v_n \otimes v_n\}\) and \(\{p(\rho_n)\}\) do not concentrate, so that our classical Young measure framework is applicable. The equiintegrability assumption, however, is more difficult to justify, since it is not implied by the usual apriori estimates for (1.1) related to the energy inequality. We use this assumption here for convenience in order to not overly complicate the presentation; the interested reader is referred to the paper [12], which extends the Fonseca-Müller theory to non-equidintegrable sequences. An application of this extension to the Euler system then implies a result similar to Theorem 3 including concentration effects. We omit details.

**Remark 4**

Condition (4.5) is rather abstract, as it involves the little understood concept of \(\mathcal{A}_L\)-quasiconvexity. It is not even obvious whether this condition is vacuous or not. Indeed, for the incompressible Euler equations, one can obtain a similar Jensen inequality, although the wave cone is so large that the corresponding \(\mathcal{A}\)-quasiconvexity already implies convexity, so that (4.5) is tautologically satisfied. The main point of this paper is to demonstrate that this is not the case for the compressible system, which is accomplished in Subsection 4.2 below.

**Proof.** The proof proceeds through the following steps:

Step 1: The fact that the Young measure \(\nu\) is a measure valued solution to the compressible Euler equations (1.1) is a direct consequence of the Fundamental Theorem of Young measures and the lack of concentrations implied by the equiintegrability assumptions.

Step 2: The sequence \{(\rho_n, v_n)\} naturally gives rise to a sequence of weak solutions \(\{z_n\} = \{\rho_n, m_n, U_n, q_n\}\) to the linear system (2.1), by defining

\[
m_n = \rho_n v_n, \quad U_n := \rho_n v_n \otimes v_n - \frac{1}{3} \rho_n |v_n|^2 I, \quad q_n := p(\rho_n) + \frac{3}{2} \rho_n |v_n|^2,
\]

where \(I\) is the \(3 \times 3\) identity matrix. The assumptions of \(\gamma\)-equiintegrability of \(\{\rho_n\}\) (with \(\gamma \geq 1\)) and of 2-equiintegrability of \(\{\sqrt{\rho_n} v_n\}\) guarantee that the sequence of subsolutions \(\{z_n\}\) is equiintegrable. Indeed, since \(\{\sqrt{\rho_n} v_n\}\) is 2-equiintegrable, the sequences \(\rho_n v_n \otimes v_n\) and \(\rho_n |v_n|^2\) are equiintegrable; the equiintegrability of \(\{p(\rho_n)\}\) follows from \(\gamma\)-equiintegrability of \(\{\rho_n\}\) and the growth assumption on \(p\); and the equiintegrability of \(\{m_n\}\) follows from the 2-equiintegrability of \(\{\sqrt{\rho_n} v_n\}\) and the 2-equiintegrability of \(\{\sqrt{\rho_n}\}\), which in turn follows from the imbedding \(L^\gamma((0, T) \times \Omega) \subset L^1((0, T) \times \Omega)\) on our bounded domain.
Step 3: We show that the lifted measure \( \check{\nu} \) is generated by the sequence \( \{z_n\} \). This is an easy consequence of the definition of the lifted measure and the fact that \( \nu \) is generated by \( \{(\rho_n, \sqrt{\rho_n} v_n)\} \):

\[
\int_0^T \int_\Omega \langle \check{\nu}_{t,x}, (g(t, x, \cdot)) \varphi(x, t) \rangle dx dt = \int_0^T \int_\Omega \langle \nu_{t,x}, (g \circ Q)(t, x, \cdot) \rangle \varphi(x, t) dx dt
\]

\[
= \lim_{n \to \infty} \int_0^T \int_\Omega g(t, x, Q(\rho_n, \sqrt{\rho_n} v_n))(t, x) \varphi(x, t) dx dt
\]

\[
= \lim_{n \to \infty} \int_0^T \int_\Omega g(t, x, z_n(t, x)) \varphi(x, t) dx dt
\]

for all test functions \( \varphi \).

Step 4: Since the sequence of subsolutions \( \{z_n\} \) is equiintegrable, we can apply Theorem 1 with the choice \( \mu = \check{\nu} \) and \( p = 1 \) to conclude that \( \check{\nu} \) is a measure valued solution to the linear system (2.1) and it satisfies the Jensen inequality (4.5).

4.2 An explicit example

In this final section we aim to show that there exists a measure-valued solution of (1.1) which is not generated by a sequence of weak solutions, which also shows that the Jensen condition (4.5) is not vacuous. To this end, let us study the wave cone for our linear system (2.1). By virtue of (3.4), to each state vector

\[
z := (\rho, m_1, m_2, m_3, U_{11}, U_{12}, U_{13}, U_{22}, U_{23}, q) \in \mathbb{R}^{10}
\]

we associate the \( 4 \times 4 \) matrix \( Z_{A_L} \) given as

\[
Z_{A_L} = \begin{bmatrix}
\rho & m_1 & m_2 & m_3 \\
m_1 & U_{11} + q & U_{12} & U_{13} \\
m_2 & U_{12} & U_{22} + q & U_{23} \\
m_3 & U_{13} & U_{23} & -U_{11} - U_{22} + q
\end{bmatrix}.
\]

Hence, the wave cone for the operator \( A_L \) is equal to

\[
\Lambda_L = \left\{ \bar{z} \in \mathbb{R}^{10} \text{ such that } \det(Z_{A_L}) = 0 \right\}
\]

Moreover, by (3.5), the linear system (2.1) can be written

\[
\text{div}_{t,x} Z_{A_L} = 0.
\]

Finally, using Theorem 2 and an explicit construction, we can prove the existence of a measure-valued solution to (1.1) which is not generated by weak solutions:

**Theorem 4** There exists a measure-valued solution of the compressible Euler system (1.1) which is not generated by any sequence of \( L^p \)-bounded weak solutions to (1.1) (for any choice of \( p > 1 \)).
Remark 5  
1. Any reasonable sequence of approximate solutions of (1.1) will satisfy some uniform energy bound, so that the assumption of \(L^p\)-boundedness will always be met.

2. As Theorem 2 did not require any equiintegrability, the statement of Theorem 4 is true even when the potential generating sequence is allowed to concentrate. In the language of generalized Young measures, this means that there exists a generalized measure-valued solution which can not be generated by a sequence of weak solutions (take the measure from Theorem 4 as the oscillation part and choose the concentration part arbitrarily).

Proof. The idea is to choose suitable \(z^1 := (\rho^1, m^1, U^1, q^1)\) and \(z^2 := (\rho^2, m^2, U^2, q^2)\) such that the homogeneous Young measure

\[
\tilde{\nu} = \frac{1}{2} \delta_{z^1} + \frac{1}{2} \delta_{z^2}
\]

cannot be a limit of bounded weak solutions to the linear system (2.1).

Set \(\rho^1 = 1\), \(m^1 = e_1\), \(U^1 = \text{diag}(2/3, -1/3, -1/3)\), \(q^1 = p(1) + \frac{1}{3}\) and, for some \(\gamma > 0\), \(\rho^2 = \gamma\), \(m^2 = e_1\), \(U^2 = U^1/\gamma\), \(q^2 = p(\gamma) + \frac{1}{3\gamma}\). Notice that with this choice of \(z^1\) and \(z^2\), \(\tilde{\nu}\) arises as the lifted Young measure of some measure-valued solution \(\nu\) to (1.1); indeed, this follows from

\[
U^i = \frac{m^i \otimes m^i}{\rho^i} - \frac{|m^i|^2}{3\rho^i} I, \quad q^i = p(\rho^i) + \frac{|m^i|^2}{3\rho^i}, \quad i = 1, 2.
\]

Now, observe that \(\tilde{z} := z^2 - z^1\) is not in the wave cone for the operator \(A_L\). Indeed, forming the corresponding \(4 \times 4\)-matrix \(\tilde{Z}\) according to (4.6), a direct calculation yields that, with our choice of \(z^1\) and \(z^2\), \(\tilde{Z} = \text{diag}(\gamma - 1, \frac{1}{\gamma} - 1 + p(\gamma) - p(1), p(\gamma) - p(1), p(\gamma) - p(1))\) so that the determinant of the matrix \(\tilde{Z}\) is

\[
(\gamma - 1) \left( \frac{1}{\gamma} - 1 + p(\gamma) - p(1) \right) (p(\gamma) - p(1))^2,
\]

and by the assumptions on the pressure (recall \(p \geq 0\), \(p' > 0\)) it is readily seen that \(\gamma\) can be chosen in such a way that this is nonzero, i.e. \(\tilde{z} \notin \Lambda_L\).

Finally, by virtue of Theorem 2 applied with the choice \(\tilde{z}_1 := z^1\) and \(\tilde{z}_2 := z^2\), we easily see that \(\tilde{\nu}\) cannot be generated by any sequence of subsolutions. Moreover, as noted before, \(\tilde{\nu}\) arises as the lifting of a measure valued solution \(\nu\) to the original compressible Euler equations (1.1) of the form

\[
\nu = \frac{1}{2} \delta \left( \rho^1, \frac{m^1}{\sqrt{\rho^1}} \right) + \frac{1}{2} \delta \left( \rho^2, \frac{m^2}{\sqrt{\rho^2}} \right)
\]

which, as a consequence, cannot be generated by any sequence of weak solutions to (1.1), since this would contradict what we have shown at the level of subsolutions. This proves Theorem 4.
Remark 6 The construction of an explicit measure valued solution for the compressible Euler system (1.1) which is not generated by any sequence of $L^p$-bounded weak solutions (for any $p > 1$) hinges crucially on the role of the density–dependent pressure: this makes the construction clearly not possible for the incompressible Euler equations.

References


