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Oscillations,

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## Introduction

Les équations de l'écoulement d'un fluide parfait incompressible remontent aux travaux d'Euler (1707-1783). A supposer que le fluide est isolé et s'écoule dans un espace  $\mathbb{R}^d$  avec d > 1, ces équations s'écrivent

(0.0.1) 
$$\partial_t u + (u \cdot \nabla)u + \nabla p \equiv 0, \quad \text{div } u \equiv 0$$

Ici  $u \in \mathbb{R}^d$  et  $p \in \mathbb{R}$  représentent respectivement la vitesse et la pression du fluide. On complète (0.0.1) en fixant les valeurs de u à l'instant initial

$$(0.0.2) u(0,x) \equiv u_0(x) \,.$$

La théorie d'existence et d'unicité pour le problème de Cauchy (0.0.1)-(0.0.2)associé à une condition initiale  $u_0(x)$  régulière est maintenant bien établie [3, 19, 21]. En revanche, de nombreuses questions restent ouvertes en ce qui concerne les phénomènes susceptibles de se produire lorsque la régularité de  $u_0$  se dégrade. L'objectif de ce document est précisément d'apporter un éclairage sur ces aspects.

Une perte d'informations sur la régularité de  $u_0$  peut se modéliser par une explosion de la norme  $L^{\infty}$  de certaines dérivées de  $u_0$ . Cela se produit en particulier lorsque l'expression  $u_0$  se met à varier rapidement, disons sur des longueurs de taille  $\varepsilon$  avec  $\varepsilon \in [0, 1]$  qui tend vers 0, dans les directions transverses à un feuilletage prescrit de  $\mathbb{R}^d$ . On a alors affaire à une famille de conditions initiales  $u_0^{\varepsilon}$  pouvant être représentée sous la forme d'une oscillation monophase. On a typiquement

(0.0.3) 
$$u_0^{\varepsilon}(x) = h^{\varepsilon}(x) = \mathbb{H}^{\varepsilon}\left(x, \frac{\varphi(x)}{\varepsilon}\right), \quad \varepsilon \in [0, 1]$$

pour une phase  $\varphi \in \mathcal{C}^2(\mathbb{R}^d; \mathbb{R})$  et un profil  $\mathbb{H}^{\varepsilon}(x, \theta) \in \mathcal{C}^2(\mathbb{R}^d \times \mathbb{T}; \mathbb{R}^d)$  qui dépend de la variable *lente* x choisie dans l'espace physique  $\mathbb{R}^d$  et de la variable *rapide*  $\theta$  décrivant le tore  $\mathbb{T} := \mathbb{R}/\mathbb{Z}$ . On note  $\mathbb{H} := \mathbb{H}^0$  le profil principal que l'on suppose non trivial au sens où  $\partial_{\theta}\mathbb{H} \neq 0$ . On dit alors de l'oscillation  $u_0^{\varepsilon}$  qu'elle est de grande amplitude. En l'absence d'hypothèses supplémentaires, le problème de Cauchy oscillant (0.0.1)-(0.0.3) est réputé mal posé [24]. Une part des difficultés [9] provient de la façon dont les oscillations sont gérées par la contrainte de divergence nulle. Dans ce qui suit, on s'intéresse au cas de données initiales  $h^{\varepsilon}$  dont la matrice Jacobienne  $D_x h^{\varepsilon}$  est nilpotente. Plus précisément, on suppose qu'il existe un entier  $r \in \{2, \dots, d\}$  tel que

(0.0.4) 
$$D_x h^{\varepsilon}(x)^r = 0, \quad \forall (\varepsilon, x) \in ]0, 1] \times \mathbb{R}^d$$

Le Théorème 2.6 énoncé dans [6] garantit alors qu'il suffit de résoudre le système dit des gaz sans pression [12, 14, 22], à savoir

$$(0.0.5) \qquad \qquad \partial_t u^{\varepsilon} + (u^{\varepsilon} \cdot \nabla) u^{\varepsilon} \equiv 0$$

en vue de récupérer des solutions  $u^{\varepsilon}$  de (0.0.1). C'est ce point de vue qui sera adopté. Avant toute chose, la question se pose de savoir s'il est possible de construire des familles  $\{h^{\varepsilon}\}_{\varepsilon}$  vérifiant simultanément les deux contraintes (0.0.3) et (0.0.4). Il s'agit là d'un problème d'*optique géométrique* assez atypique, posé dans un contexte non encore répertorié, et qui ne se ramène pas à des situations connues.

Le système d'équations aux dérivées partielles (0.0.4) contient l'équation de Monge-Ampère "det  $D_x h^{\varepsilon}(x) \equiv 0$ " ainsi que d'autres contraintes. Il présente un caractère non linéaire évident. Sa complexité fait que, à l'exception de quelques cas très spécifiques signalés dans [6], son étude générale sous couvert de (0.0.3) n'a pas encore été complètement traitée.

Une fois les structures des fonctions  $h^{\varepsilon}$  identifiées, on souhaiterait déterminer la façon dont elles se trouvent propagées via l'équation d'évolution (0.0.5). Le régime envisagé étant *sur-critique* [8, 11, 15], de violentes instabilités sont susceptibles de survenir.

En particulier, pour d = 3, lorsque le profil  $\mathbb{H}^{\varepsilon}(x, \theta)$  incorpore des oscillations dans des directions  $\nabla \psi(x)$  qui sont transverses à  $\nabla \varphi(x)$ , il peut se produire une *superposition d'oscillations*. En l'occurrence, la phase  $\psi$  se met à osciller selon  $\varphi$  à la fréquence  $\varepsilon^{-1}$ . Il s'ensuit que les dérivés de  $u^{\varepsilon}$  à un instant t > 0sont de taille  $\varepsilon^{-2}$  alors qu'elles sont en  $\varepsilon^{-1}$  à l'instant initial t = 0.

Ainsi, le passage de t = 0 à tout instant t > 0 s'accompagne d'un changement qualitatif brutal dans le comportement asymptotique de la famille  $\{u^{\varepsilon}\}_{\varepsilon}$ . Cette monté instantanée vers les hautes fréquences traduit une complexité soudainement croissante des mouvements du fluide, renvoyant (pour  $\varepsilon \to 0$ ) à l'image de ce que seraient les *turbulences*. La justification de tels phénomènes a déjà été entreprise dans [5], mais c'était en dehors du contexte de divergence nulle considéré ici, qui complique beaucoup la discussion. Au delà de ces motivations physiques, l'objectif de cette thèse est surtout de fournir une *analyse BKW* complète de (0.0.4) dans le cadre fourni par (0.0.3). Ce faisant, on aura besoin de dégager les structures géométriques (feuilletages) qui s'avèrent génériquement compatibles avec la propagation d'oscillations de grande amplitude.

Les deux paragraphes qui suivent font un bref descriptif des résultats qui sont obtenus dans ce mémoire de thèse.

### 0.1 Résumé de la première partie.

La partie 1 fait l'objet d'un article déjà publié [7]. La discussion porte au départ sur un système qui généralise quelque peu (0.0.5). Etant donnée une fonction  $\mathbf{V} \in \mathcal{C}^{\infty}(\mathbb{R}^d; \mathbb{R}^d)$ , on y considère les équations

$$(0.1.1) \qquad \partial_t u^{\varepsilon} + (\mathbf{V} \circ u^{\varepsilon} \cdot \nabla_x) \ u^{\varepsilon} = 0, \qquad \varepsilon \in ]0,1], \qquad d \in \mathbb{N}, \qquad d > 2.$$

On complète (0.1.1) à l'aide d'une famille de données initiales  $\{u_0^{\varepsilon}\}_{\varepsilon}$  ajustée comme en (0.0.3). On fixe r > 0 et on travaille sur des domaines localisés en espace, du type

$$\Omega_r^T := \{(t, x) \in [0, T] \times \mathbb{R}^d; \quad | \ x | +t \ V < r \}, \quad (V, T, r) \in (\mathbb{R}_+^*)^3.$$

On pose  $W := \mathbf{V} \circ H$ . On suppose que les ingrédients  $\mathbf{V}$ ,  $\mathbb{H}^{\varepsilon}$  et  $\varphi$  sont des fonctions régulières, disons de classe  $\mathcal{C}^1$ . On impose de plus  $\partial_{\theta} W \neq 0$  et le caractère non stationnaire de la phase

$$(0.1.2) \qquad \nabla \varphi(x) \neq 0, \qquad \forall x \in \Omega_r^0$$

La contrainte (0.1.1) forme un système hyperbolique quasilinéaire qui admet une vitesse de propagation finie contrôlée par

$$V := Sup \left\{ \left| \mathbf{V} \circ \mathbb{H}^{\varepsilon}(x, \theta) \right|; \ (\varepsilon, x, \theta) \in [0, 1] \times \Omega_{r}^{0} \times \mathbb{T} \right\}.$$

Le paramètre  $\varepsilon \in [0, 1]$  étant fixé, des résultats standards [20] garantissent l'existence de  $T^{\varepsilon} > 0$  tel que le problème de Cauchy (0.1.1)-(0.0.3) ait une solution  $u^{\varepsilon}(t, x)$  de classe  $\mathcal{C}^1$  sur  $\Omega_r^{T^{\varepsilon}}$ . Par contre, on a généralement

(0.1.3) 
$$\limsup_{\varepsilon \to 0} T^{\varepsilon} = 0$$

Les obstructions viennent de la formation éventuelle de chocs. Notant

$$\mathbb{X}^{\varepsilon}(t,x) := x + t W^{\varepsilon}\left(x, \frac{\varphi(x)}{\varepsilon}\right), \qquad W^{\varepsilon} := \mathbf{V} \circ \mathbb{H}^{\varepsilon},$$

le croisement des droites (dites caractéristiques)  $\{\mathbb{X}^{\varepsilon}(t,x), t \in \mathbb{R}_+\}$  peut en effet venir contredire la persistence de la régularité  $\mathcal{C}^1$  de  $u^{\varepsilon}$ . La définition 1.2.1 qui suit met de coté les situations pour lesquelles cela ne se produit pas.

**Définition 0.1.1.** La famille  $\{h^{\varepsilon}\}_{\varepsilon}$  est dite <u>compatible</u> s'il existe un instant T > 0 et une constante c > 0 tels que

$$(0.1.4) \qquad det \ D_x \mathbb{X}^{\varepsilon}(t,x) \ge c > 0, \quad \forall (t,x,\varepsilon) \in [0,T] \times \Omega_r^0 \times ]0,1].$$

Il est possible de traduire le critère (0.1.4) sous la forme des conditions nécessaires et suffisantes contenues dans (1.2.13), portant sur le couple ( $\varphi$ , W). C'est ce qui est fait au niveau de la Proposition 1.2.3. Toutefois, la contrainte (1.2.13) se prète difficilement à une étude complète. C'est pourquoi, dans une première approche, on a eu recours à des conditions plus restrictives (qui restent cependant intrinsèques) sous la forme de la notion suivante.

**Définition 0.1.2.** Le couples  $(\varphi, W) \in C^2(\Omega^0_r; \mathbb{R}) \times C^1(\Omega^0_r \times \mathbb{T}; \mathbb{R}^d)$  est dit <u>bien préparé</u> s'il satisfait le système de contraintes suivant

$$(0.1.5) \qquad \begin{cases} \partial_{\theta} W(x,\theta) \subset \nabla \varphi^{\perp} \\ \Pi_{\partial_{\theta} W(x,\theta)^{\perp}} D_{x} W(x,\theta) \Pi_{\nabla \varphi(x)^{\perp}} = 0 \end{cases}, \qquad \forall (x,\theta) \in \Omega^{0}_{r} \times \mathbb{T}$$

où  $\Pi_u$  désigne le projecteur orthogonal dans la direction  $u \in \mathbb{R}^d$ .

Le chapitre 1.3 est consacré à la discussion du système non linéaire (0.1.5). Cela requiert en premier lieu la succession des Lemmes 1.3.1, 1.3.2 et 1.3.3 en vue de mettre en valeur les conditions satisfaites par la phase  $\varphi$ . Celle-ci doit être localement constante suivant un champ d'espaces vectoriels noté  $\mathbb{E}$ . Une fois la structure de  $\varphi$  dégagée, il devient possible d'identifier celle du profil W via un travail de factorisation, voir la Proposition 1.3.1.

La partie 1.4 aborde le problème de l'évolution en temps. Par construction, on sait que toute famille  $\{h^{\varepsilon}\}_{\varepsilon}$  issue d'un couple  $(\varphi, W)$  bien préparé donne lieu à l'existence d'une suite  $\{u^{\varepsilon}\}_{\epsilon}$  composée de solutions de (0.1.1) sur un domaine  $\Omega_r^T$  indépendant de  $\varepsilon \in [0, 1]$ . On peut même établir la propagation de (0.1.5) au travers de (0.1.1).

**Théorème 1.** On se donne un couple  $(\varphi, W)$  <u>bien</u> <u>préparé</u>. Alors, le problème de Cauchy formé par le système (à priori surdéterminé)

(0.1.6) 
$$\begin{cases} \partial_t \mathbf{H} + \mathbf{V} \circ \mathbf{H} \cdot \nabla_x \mathbf{H} = 0, \\ \partial_t \Phi + (\mathbf{V} \circ \mathbf{H}) \cdot \nabla_x \Phi = 0, \\ \mathbf{V} \circ \mathbf{H}^* \cdot \nabla_x \Phi = 0, \end{cases}$$

associé à la donnée initiale

(0.1.7) 
$$\mathbf{H}(0, x, \theta) = H(x, \theta), \qquad \Phi(0, x) = \varphi(x)$$

admet une solution unique sur  $\Omega_r^T \times \mathbb{T}$  pour un certain T > 0. L'onde simple  $u^{\varepsilon}(t,x) := \mathbf{H}(t,x,\Phi(t,x)/\varepsilon)$  ainsi récupérée est solution de (0.1.1) sur  $\Omega_r^T$ . De plus, pour tout  $t \in [0,T]$ , la trace  $(\Phi(t,.),\mathbf{H}(t,.))$  vérifie (0.1.5).

#### 0.2 Résumé de la second partie.

La seconde partie est consacrée uniquement à (0.0.5) dans le seul cas de la dimension trois (d = 3). Il s'agit cette fois-ci d'aboutir à une discussion aussi complète que possible de ce que contient le système (0.0.3)-(0.0.4). La prise en compte des situations non encore traitées dans la Partie 1 s'avère en fait délicate. Elle fournit sa matière aux longs développements de la Partie 2 et requiert une analyse fine qui, comme on le verra, est en lien avec des questions de géométrie.

Le problème en question admet à première vue une formulation simple. On cherche toutes les familles  $\{h^{\varepsilon}\}_{\varepsilon}$  impliquant des données initiales  $h^{\varepsilon}$  qui se mettent sous la forme

(0.2.1) 
$$h^{\varepsilon}(x) = w\left(x, \frac{\varphi(x)}{\varepsilon}\right), \qquad \partial_{\theta}w \neq 0, \qquad \forall (\varepsilon, x) \in ]0, 1] \times \Omega^{0}_{r}$$

pour une phase  $\varphi$  non stationnaire, et vérifiant

(0.2.2) 
$$(D_x h^{\varepsilon}(x))^r = 0, \quad r \in \{2,3\}, \quad \forall (x,\varepsilon) \in \Omega_r^0 \times ]0,1].$$

On peut formaliser cette double propriété (0.2.1)-(0.2.2) au niveau du couple  $(\varphi, w)$  suivant la définition suivante.

**Définition 0.2.1.** Soit  $\varphi \in C^1(\Omega^0_r; \mathbb{R})$  et  $w \in C^1(\Omega^0_r \times \mathbb{T}; \mathbb{R}^3)$  tels que

(0.2.3) 
$$\partial_{\theta} w(x,\theta) \neq 0, \qquad \nabla \varphi(x) \neq 0.$$

Le couple  $(\varphi, w)$  est dit <u>compatible</u> sur le domaine  $\Omega_r^0$  si la famille  $\{h^{\varepsilon}\}_{\varepsilon}$  qui est construite à partir de  $(\varphi, w)$  via (0.2.1) vérifie (0.2.2).

Une étape préliminaire (Proposition 2.2.1) consiste à identifier via des calculs formels les conditions à imposer sur  $\varphi$  et w. Lorsque

$$rg(D_xw(x,\theta)) = dim(Im(D_xw)(x,\theta)) = 1, \quad \forall (x,\theta) \in \Omega^0_r \times \mathbb{T},$$

on peut écrire w sous la forme

$$w(x,\theta) = \mathbf{W}(\psi(x,\theta),\theta), \qquad \forall (x,\theta) \in \Omega^0_r \times \mathbb{T}$$

puis extraire le système d'équations (2.2.8)-(2.2.10)-(2.2.9) portant sur les ingrédients que sont  $\varphi$ ,  $\psi$  et **W** (voir le Lemme 2.2.1). Il devient alors assez facile de conclure.

En revanche, sous l'hypothèse

$$rg(D_xw(x,\theta)) = dim(Im(D_xw)(x,\theta)) = 2, \qquad \forall (x,\theta) \in \Omega^0_r \times \mathbb{T},$$

le travail de factorisation de la fonction w s'avère plus compliqué. Il faut la Proposition 2.2.2 pour établir qu'on peut obtenir

$$w(x,\theta) = \mathbf{W}(\varphi(x),\psi(x,\theta),\theta), \qquad \forall (x,\theta) \in \Omega^0_r \times \mathbb{T}$$

tandis que la Proposition 2.2.3 met en valeur les contraintes (2.2.17), (2.2.18), (2.2.19) et (2.2.20) qu'il convient de retenir en ce qui concerne  $\varphi$ ,  $\psi$  et **W**. Ces deux derniers énoncés 2.2.2 et 2.2.3 sont prouvés respectivement au niveau des chapitres 2.2.2 et 2.2.3.

On remarque alors que l'annulation ou non de la quantité  $\nabla \varphi \cdot \partial_{\psi} \mathbf{W}$  joue un rôle intrinsèque, au sens où elle n'est pas modifiée par les changements d'inconnues naturellement autorisés. On distingue donc deux situations, celle correspondant à  $\nabla \varphi \cdot \partial_{\psi} \mathbf{W} \equiv 0$ , et celle pour laquelle  $\nabla \varphi \cdot \partial_{\psi} \mathbf{W} \not\equiv 0$ .

Lorsque  $\nabla \varphi \cdot \partial_{\psi} \mathbf{W} \equiv 0$ , on retrouve des structures similaires à celles observées dans la Partie 1. Par exemple, il existe deux fonctions f et g donnant lieu à

$$\nabla \varphi \equiv {}^t (f(\varphi), 1, g(\varphi)) \partial_2 \varphi$$

Pour autant, en vue d'aboutir à une description complète de ce cas, il convient d'incorporer quelques aspects inédits (en comparaison de la Partie 1).

Le traitement de l'autre situation, lorsque  $\nabla \varphi \cdot \partial_{\psi} \mathbf{W} \neq 0$ , réclame une tout autre démarche qui occupe l'essentiel de ce qui suit. Le point de départ est fourni par le système de contraintes du Lemme 2.4.1. Il s'agit maintenant de trouver des fonctions  $\varphi, \psi$  et  $\mathbf{W}$  non triviales dans le sens où

$$\nabla \varphi \wedge \nabla \psi \neq 0, \qquad \partial_{\varphi} \mathbf{W} \wedge \partial_{\psi} \mathbf{W} \neq 0, \qquad \partial_{\theta} w \neq 0$$

et ajustées de manière à ce qu'il existe une fonction  $k(x, \theta)$  telle que

$$(0.2.4) \begin{cases} \nabla \varphi \cdot \partial_{\theta} w = 0, \\ \nabla \psi \cdot \partial_{\theta} w = 0, \\ \nabla \varphi \cdot (\partial_{\varphi} \mathbf{W} - k \ \partial_{\psi} \mathbf{W}) = 0, \\ \nabla \psi \cdot (\partial_{\varphi} \mathbf{W} - k \ \partial_{\psi} \mathbf{W}) = 0, \\ (k \ \nabla \varphi + \nabla \psi) \cdot \partial_{\varphi} \mathbf{W} = 0, \\ (k \ \nabla \varphi + \nabla \psi) \cdot \partial_{\psi} \mathbf{W} = 0. \end{cases}$$

Un travail de reformulation mené au Paragraphe 2.4.1.1 (voir la Proposition 2.4.1) permet dans un premier temps de remplacer (0.2.4) par

$$(0.2.5) \begin{cases} \partial_{1}\varphi + \partial_{v}\mathfrak{L}(\psi, v) \,\partial_{3}\varphi + \frac{\partial_{\theta}\psi}{\partial_{\theta}v} \left[\partial_{2}\varphi + \partial_{\psi}\mathfrak{L}(\psi, v) \,\partial_{3}\varphi\right] \equiv 0, \\ \partial_{1}\psi + \partial_{v}\mathfrak{L}(\psi, v) \,\partial_{3}\psi + \frac{\partial_{\theta}\psi}{\partial_{\theta}v} \left[\partial_{2}\psi + \partial_{\psi}\mathfrak{L}(\psi, v) \,\partial_{3}\psi\right] \equiv 0, \\ \partial_{1}\psi + \partial_{v}\mathfrak{L}(\psi, v) \,\partial_{3}\psi - \frac{\partial_{\theta}\psi}{\partial_{\theta}\mathbf{V}} \,\partial_{\varphi}\mathbf{V} \left[\partial_{1}\varphi + \partial_{v}\mathfrak{L}(\psi, v) \,\partial_{3}\varphi\right] \equiv 0. \end{cases}$$

Quelques précisions sont nécessaires pour la lecture de (0.2.5). Il faut dire que les expressions  $\partial_{\theta} \mathbf{V}$  et  $\partial_{\varphi} \mathbf{V}$  sont évaluées au point  $(\varphi(x), \psi(x, \theta), \theta)$ , que la fonction  $\mathfrak{L}$  induit de la non linéarité, et que le symbole v désigne

$$v(x,\theta) := \mathbf{V}(\varphi(x), \psi(x,\theta), \theta)$$

Bien que déjà pas mal décanté, le système (0.2.5) n'est pas pour autant directement exploitable. Il faut encore effectuer plusieurs changements de variables (voir le Paragraphe 2.4.1.2) afin d'isoler la clé sur laquelle repose l'ensemble de l'édifice.

Les systèmes (0.2.4) et (0.2.5) sont tout à fait spécifiques à la problématique qui a été mise en place. Visiblement, ils n'ont jamais été étudiés. Du coup, la méthode employée pour les résoudre est complètement originale. Elle ne s'appuie pas du tout sur une approche classique en optique géométrique.

La stratégie, a priori peu naturelle mais a fortiori incontournable, consiste à éclater la phase  $\varphi(x)$  en une fonction  $\Phi(x_1, x_2, u, v)$  qui au lieu de dépendre des trois variables  $x_1, x_2$  et  $x_3$  met en jeu quatre variables notées  $x_1, x_2, u$  et v. On passe de  $\mathbb{R}^3$  à  $\mathbb{R}^4$ . L'expression  $\Phi$  ainsi obtenue doit vérifier deux équations de transport dans  $\mathbb{R}^4$ , en l'occurrence

$$(0.2.6) X \Phi \equiv 0, X := \partial_1 + R(x_1, x_2, u, v) \partial_2,$$

(0.2.7) 
$$Y \Phi \equiv 0, \qquad Y := R(x_1, x_2, u, v) \ \partial_u + \partial_v.$$

Intuitivement, la première contrainte (0.2.6) peut être assimilée à la condition  $\partial_{\theta}\varphi \equiv 0$ . Elle est destinée à faire en sorte que la non linéarité n'induise pas une auto-oscillation de la phase  $\varphi$ : on veut que la fonction  $\varphi$  ne se mette pas à dépendre de  $\theta$  ce qui ne va pas de soi compte tenu de la force des effets non linéaires en présence. Quant à la seconde restriction (0.2.7), elle provient de (0.2.2) à l'issu du long travail de ré-interprétation qui a été effectué.

Le point crucial maintenant, c'est que toute fonction  $\Phi$  compatible avec (0.2.6) et (0.2.7) doit aussi vérifier  $Z \Phi \equiv 0$  pour tout champ de vecteurs Z appartenant à l'algèbre de Poisson  $\mathcal{A}$  engendrée par X et Y. Comme la construction impose que le gradient de  $\varphi$  (et donc de  $\Phi$ ) soit non trivial, il faut nécessairement que la dimension de  $\mathcal{A}$  soit inférieure ou égale à 3. Ce critère induit le système d'équations aux dérivées partielles (2.4.38) dans le cas  $\dim \mathcal{A} = 2$ . Par contre, lorsque  $\dim \mathcal{A} = 3$ , on doit gérer :

- (0.2.8)  $(XR) YXR 2 (XR) XYR + (YR) X^2R \equiv 0,$
- (0.2.9)  $(YR) XYR 2 (YR) YXR + (XR) Y^2R \equiv 0.$

Par ailleurs, la fonction R étant issue du procédé de réduction, elle doit être soumise aux contraintes attenantes. Plus précisément, on doit avoir

$$R(x_1, x_2, u(x, v), v) = \partial_v u(x, v)$$

pour une certaine fonction  $u(x, v) \in C^1(\mathbb{R}; \mathbb{R})$  qui est astreinte à la double loi de conservation scalaire

$$(0.2.10) \qquad \partial_2 u + \partial_u \mathfrak{L}(u,v) \ \partial_3 u = 0, \qquad \partial_1 u + \partial_v \mathfrak{L}(u,v) \ \partial_3 u = 0.$$

Le chapitre 2.4.3 montre que toute fonction R obtenue par un tel procédé s'écrit sous la forme

(0.2.11) 
$$R = -\frac{\partial_v \alpha}{\partial_u \alpha}, \qquad \alpha := \mathfrak{K}(u,v) + \partial_v \mathfrak{L}(u,v) \ x_1 + \partial_u \mathfrak{L}(u,v) \ x_2$$

où  $\mathfrak{K}$  et  $\mathfrak{L}$  sont deux fonctions quelconques prises dans  $\mathcal{C}^1(\mathbb{R}^2; \mathbb{R})$ . Arrivé à ce stade, l'enjeu consiste à pouvoir résoudre (0.2.8)-(0.2.9) en s'appuyant sur la liberté dont on dispose quant au choix des fonctions  $\mathfrak{K}$  et  $\mathfrak{L}$ . Autrement dit, il s'agit de tester les conditions d'intégrabilité (0.2.8) et (0.2.9) dans le contexte offert par (0.2.11).

Ce programme est ce qui fournit la matière du chapitre 2.4.4. Les différentes situations possibles sont triées selon la dimension de l'algèbre  $\mathcal{A}$  (c'est à dire 2 ou 3) puis, plus finement, selon l'annulation ou non des quantités XR ou YR. On obtient ainsi une une classification (presque exhaustive) de tous les coefficients R autorisés. Vient ensuite un travail de reconstruction permettant de remonter de la connaissance de R à l'identification de  $\varphi$ ,  $\psi$ et  $\mathbf{W}$ . C'est ce qui est fait au niveau du chapitre 2.4.5. Quelques exemples venant illustrer la façon dont la procédure peut se concrétiser sont apportés à l'occasion du chapitre 2.4.6.

Le chapitre 2.5 est consacré au problème de l'évolution en temps. Il met en oeuvre les objets  $\varphi$ ,  $\psi$  et **W** extraits ci-dessus de la manière suivante.

**Théorème 2.** Il existe un instant T > 0 tel que le système suivant (qui a priori est surdéterminé)

(0.2.12) 
$$\begin{cases} \partial_t \Phi + (\mathbf{W}(\Phi, \Psi, \theta) \cdot \nabla) \Phi = 0, \\ \partial_t \Psi + (\mathbf{W}(\Phi, \Psi, \theta) \cdot \nabla) \Psi = 0, \\ (\partial_\theta \mathbf{W}(\Phi, \Psi, \theta) + \partial_\theta \Psi \partial_\psi \mathbf{W}(\Phi, \Psi, \theta)) \cdot \nabla \Phi = 0, \end{cases}$$

associé aux données initiales  $\Phi(0, x) = \varphi(x)$  et  $\Psi(0, x, \theta) = \psi(x, \theta)$  admette une unique solution sur  $\Omega_r^T \times \mathbb{T}$ . Pour tout  $\varepsilon \in [0, 1]$ , l'oscillation

$$u^{\varepsilon}(t,x) := \mathbf{W}\Big(\Phi(t,x), \Psi\big(t,x,\frac{\Phi(t,x)}{\varepsilon}\big), \frac{\Phi(t,x)}{\varepsilon}\Big), \qquad \varepsilon \in \left]0,1\right]$$

est solution de (0.0.1) et (0.0.5) sur  $\Omega_r^T \times \mathbb{T}$ . De plus, pour tout  $t \in [0,T]$ le couple  $(\Phi(t,.), \widetilde{\mathbf{W}}(t,.))$  où  $\widetilde{\mathbf{W}}(t,x,\theta) = \mathbf{W}(\Phi(t,x), \Psi(t,x,\theta), \theta)$  est encore compatible sur le domaine  $B(0, r - tV[\times\mathbb{T}]$ . Le paragraphe 2.5.2 revient sur le phénomène de superposition des phases qui a été évoqué précédemment, en vue de l'illustrer au travers d'un exemple concret. Tout compte fait, son apparition apparait comme un sous-produit assez anecdotique de l'analyse.

Il y a aussi une annexe qui se positionne comme suit. La Proposition 2.4.4 traite du problème soulevé par l'étude du système (0.2.8)-(0.2.9) complété de (0.2.11), ceci dans le cas le plus complexe à savoir lorsque  $\dim \mathcal{A} = 3$ ,  $XR \neq 0$  et  $YR \neq 0$ . Pour simplifier la présentation, l'énoncé 2.4.4 fournit directement des formules possibles pour  $\mathfrak{K}$  et  $\mathfrak{L}$ , puis se contente de vérifier qu'elles conviennent. Montrer que la liste ainsi obtenue est exhaustive est loin d'être facile. Ce seul point réclame les raisonnements subtils et les lourds calculs placés en Appendice 2.6.

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## Chapter 1

# Compatibility conditions to allow some large amplitude WKB analysis for Burger's type systems.

Abstract. In this article, we discuss the problem of finding large amplitude asymptotic expansions for monophase oscillating solutions of the following multidimensional (d > 1) Burger's type system

 $(\diamondsuit) \qquad \partial_t \mathbf{u} + (V \circ \mathbf{u} \cdot \nabla_x) \mathbf{u} = 0, \quad \mathbf{u} \in \mathbb{R}^d, \quad (t, x) \in \mathbb{R} \times \mathbb{R}^d, \quad V \in \mathcal{C}^1(\mathbb{R}^d; \mathbb{R}^d).$ More precisely, we are concerned with families  $\{\mathbf{u}^\varepsilon\}_{\varepsilon \in ]0,1]}$  made of solutions to  $(\diamondsuit)$  and having a development of the form  $\mathbf{u}^\varepsilon(t, x) = \mathbf{H}(t, x, \frac{\Phi(t, x)}{\varepsilon}) + O(\varepsilon)$  where the function  $\theta \longmapsto \mathbf{H}(t, x, \theta)$  is periodic. In general, due to the formation of shocks, such a construction is not possible on a domain  $\Omega$  which does not depend on  $\varepsilon \in ]0, 1]$ . In this article, we perform a detailed analysis of the restrictions to impose on the profile  $\mathbf{H}$  and on the phase  $\Phi$  in order to remedy this. Among these compatibility conditions, we can isolate some new interesting system of nonlinear partial differential equations : see (1.1.11). We explain how to solve it and we describe how the underlying structure is propagated through the evolution equation.

### 1.1 Introduction.

Note  $x = (x_1, \cdots, x_d) \in \mathbb{R}^d$  and

$$|x| := \left(\sum_{j=1}^{d} x_j^2\right)^{\frac{1}{2}}, \qquad \partial_j := \frac{\partial}{\partial x_j}, \qquad \partial_\theta := \frac{\partial}{\partial \theta}.$$

Let  $(T, \mathbf{V}, r) \in (\mathbb{R}^*_+)^3$ . Work on the domain

 $\Omega^T := \left\{ (t, x) \in [0, T] \times \mathbb{R}^d; |x| + \mathbf{V} t \le r \right\}, \qquad d \in \mathbb{N} \setminus \{0, 1\}.$ 

Select a function  $V \in \mathcal{C}^1(\mathbb{R}^d; \mathbb{R}^d)$  and consider the Burger's type system

(1.1.1) 
$$\partial_t \mathbf{u} + (V \circ \mathbf{u} \cdot \nabla_x) \mathbf{u} = 0, \quad \mathbf{u} \in \mathbb{R}^d, \quad (t, x) \in \Omega^T$$

Associate (1.1.1) with a family of initial datas

(1.1.2) 
$$\mathbf{u}^{\varepsilon}(0,x) = h^{\varepsilon}(x) = H\left(x,\frac{\varphi(x)}{\varepsilon}\right) + O(\varepsilon), \quad \varepsilon \in ]0,1]$$

defined on the ball  $B(0,r] := \{x \in \mathbb{R}^d; |x| \le r\}$ , built with

$$H(x,\theta) \in \mathcal{C}^1(B(0,r] \times \mathbb{T}; \mathbb{R}^d), \qquad \varphi(x) \in \mathcal{C}^1(B(0,r]; \mathbb{R}), \qquad \mathbb{T} := \mathbb{R}/\mathbb{Z}$$

and consisting of large amplitude high frequency monophase oscillating waves, which means to require a non trivial (main) profile

$$(1.1.3) \qquad \exists (x,\theta) \in B(0,r] \times \mathbb{T} \, ; \qquad \partial_{\theta} W(x,\theta) \neq 0 \, , \qquad W \, := \, V \circ H$$

and a non stationary phase

(1.1.4) 
$$\nabla_x \varphi(x) \neq 0, \quad \forall x \in B(0,r].$$

To describe more precisely the expressions involved in (1.1.2), select a function

$$\begin{array}{rcl} H \,:\, [0,1] \times B(0,r] \times \mathbb{T} & \longrightarrow & \mathbb{R}^d \\ & (\varepsilon, x, \theta) & \longmapsto & H^{\varepsilon}(x, \theta) \end{array}$$

which is smooth with respect to the parameter  $\varepsilon \in [0, 1]$ 

$$H \in \mathcal{C}^{\infty}([0,1]; \mathcal{C}^1(B(0,r] \times \mathbb{T}; \mathbb{R}^d))$$

and whose Taylor expansion near  $\varepsilon = 0$  is noted

(1.1.5) 
$$H^{\varepsilon}(x,\theta) := H(x,\theta) + \sum_{j=1}^{m} \varepsilon^{j} H^{j}(x,\theta) + O(\varepsilon^{m+1}), \qquad m \in \mathbb{N}^{*}.$$

Define :

(1.1.6) 
$$h^{\varepsilon}(x) := H^{\varepsilon}\left(x, \frac{\varphi(x)}{\varepsilon}\right), \qquad W^{\varepsilon}(x, \theta) := V \circ H^{\varepsilon}(x, \theta).$$

Associate (1.1.1) with the family of initial datas  $\{h^{\varepsilon}\}_{\varepsilon \in [0,1]}$ . The evolution equation (1.1.1) is a quasilinear (diagonal) system of hyperbolic equations. The speed of propagation is finite. More precisely, it can be uniformly controlled by

$$\mathbb{R} \ni \mathbf{V} := \left\{ \sup |V \circ H^{\varepsilon}(x,\theta)| \, ; \, (\varepsilon, x, \theta) \in [0,1] \times B(0,r] \times \mathbb{T} \right\}.$$

Standard results [20] guarantee the existence of  $T^{\varepsilon} > 0$  such that the Cauchy problem (1.1.1)-(1.1.2) has a local  $C^1$  solution  $\mathbf{u}^{\varepsilon}(t, x)$  on the truncated cone  $\Omega^{T^{\varepsilon}}$ . In the context of (1.1.1), the limitations on  $T^{\varepsilon}$  are due to the formation of shocks. The size of  $T^{\varepsilon}$  can be estimated very precisely [1, 5, 22] in terms of  $h^{\varepsilon}$ . In general, this yields

(1.1.7) 
$$\limsup_{\varepsilon \to 0} T^{\varepsilon} = 0$$

In this article, we exhibit solutions  $\mathbf{u}^{\varepsilon}$  on a fixed domain  $\Omega^{T}$  (with T > 0) having the asymptotic description

(1.1.8) 
$$\mathbf{u}^{\varepsilon}(t,x) = \mathbf{H}\left(t,x,\frac{\Phi(t,x)}{\varepsilon}\right) + O(\varepsilon), \quad \varepsilon \in ]0,1].$$

The main novelty in comparison with usual works [18] in WKB analysis is the size of the involved oscillations. Indeed, in a quasilinear context such as (1.1.1), the standard regime is given by *weakly nonlinear geometric optics* [16] which means to consider expansions of the following form

(1.1.9) 
$$\mathbf{u}^{\varepsilon}(t,x) = \mathbf{u}(t,x) + \varepsilon \mathbf{H}^{1}\left(t,x,\frac{\Phi(t,x)}{\varepsilon}\right) + O(\varepsilon), \quad \varepsilon \in \left]0,1\right].$$

Of course, to deal with (1.1.8) in place of (1.1.9) requires to manage much stronger nonlinear phenomena. In particular, the interplay between the phase  $\Phi$  and the profile **H** is reinforced.

In fact, the construction can be faced only if the expressions  $\varphi := \Phi_{|t=0}$  and  $H := \mathbf{H}_{|t=0}$  satisfy very special restrictions. The corresponding constraints in the case of the dimension d = 2 are brought out in the recent contribution [5]. The aim of the present paper is precisely to generalize the discussion when d > 2 and to study more deeply the structure to impose on  $\varphi$  and H.

• In the Section 2, we exhibit (Proposition 1.2.2) necessary and sufficient compatibility conditions on  $\varphi(x)$  and  $W(x,\theta) := V \circ H(x,\theta)$  in order to guarantee that

(1.1.10) 
$$\liminf_{\varepsilon \to 0} T^{\varepsilon} = \tilde{T} > 0.$$

From these compatibility conditions, we can isolate some interesting system of nonlinear partial differential equations which we introduce below.

Let  $\mathbf{u} = {}^t(\mathbf{u}_1, \cdots, \mathbf{u}_d) \in \mathbb{R}^d$ . Note  $\mathbf{u}^{\perp}$  or  ${}^t\mathbf{u}^{\perp}$  the hyperplane of  $\mathbb{R}^d$  composed with the directions orthogonal to the vector  $\mathbf{u}$ , that is

 $\mathbf{u}^{\perp} \equiv {}^{t}\mathbf{u}^{\perp} := \left\{ \mathbf{v} = {}^{t}(\mathbf{v}_{1}, \cdots, \mathbf{v}_{d}) \in \mathbb{R}^{d} ; {}^{t}\mathbf{v} \cdot \mathbf{u} = \sum_{j=1}^{d} \mathbf{v}_{j} \mathbf{u}_{j} = 0 \right\}.$ 

Consider the orthogonal projector  $\Pi_F$  from  $\mathbb{R}^d$  onto the vector space F, that is the operator  $\Pi_F$  defined by the conditions

 $\mathbf{u} = \Pi_F \mathbf{u} + (I - \Pi_F) \mathbf{u}, \qquad \Pi_F \mathbf{u} \in F, \qquad (I - \Pi_F) \mathbf{u} \in F^{\perp}.$ 

Select  $W \in \mathcal{C}^1(B(0,r] \times \mathbb{T}; \mathbb{R}^d)$ . The symbol  $D_x W(x,\theta)$  is for the Jacobian matrix

$$D_x W(x,\theta) = \left(\partial_j W_i(x,\theta)\right)_{1 \le i, j \le d}, \qquad W(x,\theta) = {}^t (W_1, \cdots, W_d).$$

Définition 1.1.1. The couple

 $(\varphi, W) \in \mathcal{C}^2(B(0, r]; \mathbb{R}) \times \mathcal{C}^2(B(0, r] \times \mathbb{T}; \mathbb{R}^d)$ 

is said to be <u>well</u> <u>prepared</u> if it satisfies the following system

(1.1.11) 
$$\begin{cases} \partial_{\theta} W(x,\theta) \subset \nabla \varphi(x)^{\perp} \\ \Pi_{\partial_{\theta} W(x,\theta)^{\perp}} D_{x} W(x,\theta) \Pi_{\nabla \varphi(x)^{\perp}} = 0 \end{cases}, \qquad \forall (x,\theta) \in B(0,r] \times \mathbb{T}.$$

As explained before, the study of (1.1.11) is the main motivation of the present article. Indeed, the structure of (1.1.11) is new and interesting. It is not a usual quasilinear system because it is made of fully nonlinear constraints on the derivatives  $\partial_j W_i$ ,  $\partial_\theta W_i$  and  $\partial_j \varphi$ . It extends to the case  $d \geq 3$  preliminary conditions which have been emphasized (only when d = 2) in the recent contribution [5].

• In the Section 3, we work under natural assumptions on  $\varphi$  and W. In this framework, we succeed in classifying all the solutions of (1.1.11). The fact that such a complete discussion is available is very surprising. At all events, this confirms the coherence of (1.1.11).

The first observation is that any phase  $\varphi$  involved in (1.1.11) inherits some affine structure. Its level surfaces must be spanned by pieces of vector spaces (see Lemmas 1.3.2 and 1.3.3). This geometrical particularity seems to always play an important part when dealing with phase involved in a supercritical WKB calculus, as here. Once  $\varphi$  is determined, it becomes possible to identify all the profiles  $W(x, \theta)$ which are subjected to (1.1.11). This is done in Proposition 1.3.1. Quite a lot freedom is available in the construction of  $W(x, \theta)$ .

The function  $W(x, \theta)$  can be put in the form

 $W(x,\theta) = W_{\parallel} \big( \varphi(x), \psi(x,\theta) \big) + W_{\perp} \big( \varphi(x) \big)$ 

where  $W_{\parallel} \in \mathcal{C}^1(\mathbb{R}^2; \mathbb{R}^d)$  and  $W_{\perp} \in \mathcal{C}^1(\mathbb{R}; \mathbb{R}^d)$  are conveniently well-polarized vector fields whereas  $\psi \in \mathcal{C}^1(B(0, r] \times \mathbb{T}; \mathbb{R})$  is any scalar function. Define

$$\langle W \rangle(x) \equiv \bar{W}(x) := \int_{\mathbb{T}} W(x,\theta) \, d\theta \,, \qquad W^*(x,\theta) := W(x,\theta) - \bar{W}(x) \,.$$

The construction of large amplitude oscillating solutions to system (1.1.1) - or to variants of system (1.1.1) - is a delicate problem which has recently called some

attention. The article [14] and the related contributions are mainly concerned with *time* oscillations. In the continuity of the works [1, 10, 5], we are faced here with the case of *spatial* oscillations.

According to Section 2, any family  $\{h^{\varepsilon}\}_{\varepsilon} \in \mathcal{C}^1(B(0,r];\mathbb{R}^d)^{[0,1]}$  issued from a well prepared couple  $(\varphi, W)$  leads to a family  $\{\mathbf{u}^{\varepsilon}\}_{\varepsilon}$  which is composed with  $\mathcal{C}^1$  solutions  $\mathbf{u}^{\varepsilon}$  of (1.1.1) on  $\Omega^{\tilde{T}}$ . Now, the question is to determine the asymptotic behaviour of  $\{\mathbf{u}^{\varepsilon}\}_{\varepsilon}$  when  $\varepsilon$  goes to 0. In particular, we want to understand how the constraint (1.1.1) is propagated through the evolution equation (1.1.1).

• Explanations are given in the Section 4. They can be obtained just by looking at the simple wave solutions of (1.1.1).

**Théorème 3.** Suppose that the couple

$$(\varphi, W) \in \mathcal{C}^2(B(0, r]; \mathbb{R}) \times \mathcal{C}^1(B(0, r] \times \mathbb{T}; \mathbb{R}^d), \qquad W := V \circ H$$

is well prepared. Then, the Cauchy problem consisting in the (apparently overdetermined) system

(1.1.12) 
$$\begin{cases} \partial_t \mathbf{H} + V \circ \mathbf{H} \cdot \nabla_x \mathbf{H} = 0, \\ \partial_t \Phi + \langle V \circ \mathbf{H} \rangle \cdot \nabla_x \Phi = 0, \\ (V \circ \mathbf{H})^* \cdot \nabla_x \Phi = 0, \end{cases}$$

associated with the initial datas

(1.1.13) 
$$\mathbf{H}(0, x, \theta) = H(x, \theta), \qquad \Phi(0, x) = \varphi(x)$$

has a unique solution on  $\Omega^T \times \mathbb{T}$  for some T > 0. For all  $\varepsilon \in [0,1]$ , the simple wave  $\mathbf{u}^{\varepsilon}(t,x) := \mathbf{H}(t,x,\frac{\Phi(t,x)}{\varepsilon})$  is a solution of (1.1.1) on  $\Omega^T$ . Moreover, for all  $t \in [0,T]$ , the trace  $(\Phi(t,\cdot), \mathbf{H}(t,\cdot))$  is still subjected to (1.1.11).

At the time t = 0, it is also possible to take into account (small) perturbations of  $H(x, \frac{\varphi(x)}{\varepsilon})$ . For instance, we can select

$$h^{\varepsilon}(x) = H^{\varepsilon}\left(x, \frac{\varphi(x)}{\varepsilon}\right), \qquad \varepsilon \in \left]0, 1\right]$$

where  $H^{\varepsilon}(x,\theta)$  is like in (1.1.5). Again, the discussion of the Section 2 indicates that corresponding  $C^1$  solutions  $\mathbf{u}^{\varepsilon}(t,x)$  of (1.1.1) are still available on  $\Omega^T$ . When  $\varepsilon$  goes to 0, the expression  $\mathbf{u}^{\varepsilon}(t,x)$  remains close (in a convenient sense) to the simple wave  $\mathbf{H}(t,x,\frac{\Phi(t,x)}{\varepsilon})$ . This result can be proved by adapting and extending the method presented in [5]. The related analysis will not be developed here.

#### **1.2** Analysis of the compatibility conditions

Introduce the curves  $t \mapsto (X(t; x, \lambda), \Lambda(t; x, \lambda))$  associated with the integration of (1.1.1) along the characteristics. They are defined by the ordinary differential

equations

(1.2.1) 
$$\begin{cases} \frac{d}{dt}X = V(\Lambda), & X(0; x, \lambda) = x, \\ \\ \frac{d}{dt}\Lambda = 0, & \Lambda(0; x, \lambda) = \lambda. \end{cases}$$

The corresponding solutions are

(1.2.2) 
$$X(t;x,\lambda) = x + t V(\lambda), \qquad \Lambda(t;x,\lambda) = \lambda.$$

Define

(1.2.3) 
$$\mathbb{X}^{\varepsilon}(t,x) := X(t;x,h^{\varepsilon}(x)) = x + t W^{\varepsilon}(x,\frac{\varphi(x)}{\varepsilon}), \qquad W^{\varepsilon} := V \circ H^{\varepsilon}.$$

Any smooth  $C^1$  solution of (1.1.1)-(1.1.2) must be subjected to the relation

(1.2.4) 
$$\mathbf{u}^{\varepsilon}(t, \mathbb{X}^{\varepsilon}(t, x)) = \mathbf{u}^{\varepsilon}(t, x + t \ V \circ h^{\varepsilon}(x)) = h^{\varepsilon}(x).$$

Fix  $\varepsilon \in [0,1]$ . For t small enough, say for  $t \in [0, \tilde{T}^{\varepsilon}]$  with  $\tilde{T}^{\varepsilon} > 0$ , the implicit theorem guarantees that the application

$$\begin{array}{rcl} \mathbb{X}_t^{\varepsilon} : B(0,r] & \longrightarrow & \mathbb{X}^{\varepsilon} \bigl( t, B(0,r] \bigr) \\ & x & \longmapsto & \mathbb{X}^{\varepsilon} (t,x) \end{array}$$

is a  $\mathcal{C}^1$  diffeomorphism. Then, due to the definition of the maximal speed of propagation **V**, any point (t, x) contained in  $\Omega^{\tilde{T}^{\varepsilon}}$  is sure to be realized as  $(t, x) = (t, \mathbb{X}^{\varepsilon}(t, y))$  for some unique  $y \in B(0, r]$ . We can define

(1.2.5) 
$$\mathbf{u}^{\varepsilon}(t,x) := h^{\varepsilon} \circ (\mathbb{X}_t^{\varepsilon})^{-1}(x), \qquad (t,x) \in \Omega^{\tilde{T}^{\varepsilon}}$$

which yields a  $C^1$  solution on  $\Omega^{\tilde{T}^{\varepsilon}}$  of the Cauchy problem (1.1.1)-(1.1.2). The relation (1.2.5) implies that

(1.2.6) 
$$D_x \mathbf{u}^{\varepsilon}(t,x) := D_x h^{\varepsilon} \circ (\mathbb{X}_t^{\varepsilon})^{-1}(x) \operatorname{Co} \left[ D_x \mathbb{X}^{\varepsilon}(t,x) \right] / \det D_x \mathbb{X}^{\varepsilon}(t,x)$$

where  $\operatorname{Co}[M]$  stands for the co-matrix of M. We have (1.2.7)

$$D_{x}\mathbb{X}^{\varepsilon}(t,x) = \varepsilon^{-1} t \partial_{\theta}W^{\varepsilon}\left(x,\frac{\varphi(x)}{\varepsilon}\right) \otimes {}^{t}\nabla\varphi(x) + \mathbf{I} + t D_{x}W^{\varepsilon}\left(x,\frac{\varphi(x)}{\varepsilon}\right), \qquad W^{\varepsilon} := V \circ H^{\varepsilon}$$

where we adopt the following convention

$$u \otimes v = (u_i v_j)_{1 \le i,j \le d}, \qquad u = {}^t(u_1, \cdots, u_d), \qquad v = {}^t(v_1, \cdots, v_d).$$

Classical results - see for instance [20] - assert that a  $C^1$  solution  $\mathbf{u}^{\varepsilon}(t, x)$  on  $\Omega^T$  can be extended in time as long as the matrix  $D_x \mathbf{u}^{\varepsilon}(t, x)$  is bounded. In view of the formula (1.2.6), to recover a  $C^1$  solution  $\mathbf{u}^{\varepsilon}(t, x)$  on  $\Omega^T$ , it is necessary and sufficient to have

 $\det D_x \mathbb{X}^{\varepsilon}(t,x) > 0, \qquad \forall (t,x) \in \Omega^T.$ 

Therefore, the life span of a  $\mathcal{C}^1$  solution on a domain of propagation is bounded below by

$$T^{\varepsilon} := \sup \left\{ T > 0; \det D_x \mathbb{X}^{\varepsilon}(t, x) > 0, \forall (t, x) \in [0, T] \times B(0, r] \right\}.$$

In general, due to the presence in (1.2.7) of the (singular) term with  $\varepsilon^{-1}$  in factor, only (1.1.7) can be asserted. Now, the opposite situation is still possible providing that the family  $\{h^{\varepsilon}\}_{\varepsilon}$  is conveniently adjusted. This situation is distinguished below.

**Définition 1.2.1.** - see (1.1.6) and (1.2.3) for the definitions of  $h^{\varepsilon}$  and  $\mathbb{X}^{\varepsilon}$  - The family  $\{h^{\varepsilon}\}_{\varepsilon}$  is said to be <u>compatible</u> if there exists T > 0 and c > 0 such that

(1.2.8)  $\det D_x \mathbb{X}^{\varepsilon}(t,x) \ge c > 0, \qquad \forall (t,x,\varepsilon) \in [0,T] \times B(0,r] \times ]0,1].$ 

The preceding discussion can be summarized by the following statement.

**Proposition 1.2.1.** - see (1.1.6) for the definition of  $h^{\varepsilon}$  - Suppose that the family  $\{h^{\varepsilon}\}_{\varepsilon}$  is compatible. Then, for all  $\varepsilon \in [0, 1]$ , the expression  $\mathbf{u}^{\varepsilon}(t, x)$  defined through (1.2.5) is a  $\mathcal{C}^1$  solution on  $\Omega^T$  of the Cauchy problem (1.1.1)-(1.1.2).

Our aim now is to transcribe (1.2.8) in terms of constraints to impose on  $\varphi(x)$  and  $W(x, \theta)$ . To this end, introduce

(1.2.9) 
$$\mathcal{V} := \left\{ (x,\theta) \in B(0,r] \times \mathbb{T}; \ \partial_{\theta} W(x,\theta) \neq 0 \right\}, \qquad W := V \circ H.$$

We assume (1.1.3), that is  $\mathcal{V} \neq \emptyset$ .

**Proposition 1.2.2.** - see (1.1.6) for the definitions of  $h^{\varepsilon}$  and  $W^{\varepsilon}$  - The family  $\{h^{\varepsilon}\}_{\varepsilon}$  can be compatible only if :

(1.2.10) 
$${}^t\nabla\varphi(x)\cdot\partial_\theta W(x,\theta) = 0, \quad \forall (x,\theta)\in B(0,r]\times\mathbb{T}$$

where we recall that

$$W(x,\theta) = W^0(x,\theta) = V \circ H(x,\theta).$$

**Proof.** The reasoning is based on the identity (1.2.7) which can be formulated as

$$\varepsilon \ D_x \mathbb{X}^{\varepsilon}(t,x) \ = \ M^0\big(t,x,\frac{\varphi(x)}{\varepsilon}\big) \ + \ \varepsilon \ M^1\big(t,x,\frac{\varphi(x)}{\varepsilon}\big) \ + \ \varepsilon^2 \ t \ R^{\varepsilon}\big(t,x,\frac{\varphi(x)}{\varepsilon}\big)$$

where

$$\begin{split} M^{0}(t,x,\theta) &:= t \,\partial_{\theta} W(x,\theta) \otimes {}^{t} \nabla \varphi(x) \,, \\ M^{1}(t,x,\theta) &:= \mathrm{I} + t \, D_{x} W(x,\theta) + t \,\partial_{\theta} \big[ D_{\mathbf{u}} V \big( H^{0}(x,\theta) \big) H^{1}(x,\theta) \big] \otimes {}^{t} \nabla \varphi(x) \,, \end{split}$$

whereas the matrix  $R^{\varepsilon}(t, x, \theta)$  is a continuous function of the variables  $(\varepsilon, t, x, \theta) \in [0, 1] \times \mathbb{R} \times B(0, r] \times \mathbb{T}$ . Assume that the restriction (1.2.8) is satisfied for some T > 0 and some c > 0. We start by showing

(1.2.11) 
$${}^t \nabla \varphi(x) \cdot \partial_\theta W(x,\theta) \ge 0, \quad \forall (x,\theta) \in \mathcal{V}$$

To this end, we argue by contradiction. We suppose that we can find a point  $(\bar{x}, \bar{\theta}) \in \mathcal{V}$  such that

(1.2.12) 
$${}^{t}\nabla\varphi(\bar{x})\cdot\partial_{\theta}W(\bar{x},\bar{\theta}) < 0$$

This information allows to express the matrices  $M^0(t, \bar{x}, \bar{\theta})$  and  $M^1(t, \bar{x}, \bar{\theta})$  in a basis of  $\mathbb{R}^d$  having the form  $(e_1, e_2, \cdots, e_d)$  where  $e_1 := \partial_{\theta} W(\bar{x}, \bar{\theta})$  and where  $(e_2, \cdots, e_d)$  is a basis of  $\nabla \varphi(\bar{x})^{\perp}$ .

In this special basis, the matrices  $M^0$  and  $M^1$  look like

$$M^{0} = \begin{pmatrix} t \ {}^{t} \nabla \varphi \cdot \partial_{\theta} W & 0 & \dots & 0 \\ 0 & 0 & \cdots & 0 \\ \vdots & \vdots & \ddots & \vdots \\ 0 & 0 & \cdots & 0 \end{pmatrix}, \qquad M^{1} = \begin{pmatrix} m_{11}^{1} & \dots & m_{1d}^{1} \\ m_{21}^{1} & \dots & m_{2d}^{1} \\ \vdots & & \vdots \\ m_{d1}^{1} & \dots & m_{dd}^{1} \end{pmatrix}.$$

It follows that

$$\det D_x \mathbb{X}^{\varepsilon}(t,\bar{x}) = \varepsilon^{-d} \det \left[ M^0(t,\bar{x},\frac{\varphi(\bar{x})}{\varepsilon}) + \varepsilon M^1(t,\bar{x},\frac{\varphi(\bar{x})}{\varepsilon}) + O(\varepsilon^2) \right]$$
$$= \varepsilon^{-1} t \, {}^t \nabla \varphi(\bar{x}) \cdot \partial_\theta W(\bar{x},\frac{\varphi(\bar{x})}{\varepsilon}) \, \det M^\flat(t,\bar{x},\frac{\varphi(\bar{x})}{\varepsilon}) + O(1)$$

with

<sup>h</sup>  
$$M^{\flat} = M^{\flat}(t, \bar{x}, \bar{\theta}) = \begin{pmatrix} m_{22}^{1} & \dots & m_{2d}^{1} \\ \vdots & & \vdots \\ m_{d2}^{1} & \dots & m_{dd}^{1} \end{pmatrix}$$

When t = 0, we have  $M^1(0, \bar{x}, \bar{\theta}) = I$  so that  $M^{\flat} = I_{\mathbb{R}^{d-1}}$  and det  $M^{\flat} = 1$ . By continuity, for t small enough (say  $t \in [0, \tilde{T}]$  with  $\tilde{T} > 0$ ), it remains

$$\det M^{\flat}\big(t, \bar{x}, \frac{\varphi(\bar{x})}{\varepsilon}\big) \geq \frac{1}{2}, \qquad \forall (t, \varepsilon) \in [0, \tilde{T}] \times ]0, 1].$$

Choose  $t \in [0, \tilde{T}]$  and a sequence  $\{\varepsilon_n\}_n \in [0, 1]^{\mathbb{N}}$  tending to 0 and such that

$$\forall n \in \mathbb{N}, \quad \exists k_n \in \mathbb{Z}; \quad \varphi(\bar{x}) = \varepsilon_n \left(\theta + 2 k_n \pi\right).$$

Then, by construction, we have

$$\exists C \in \mathbb{R} \, ; \quad \det D_x \mathbb{X}^{\varepsilon_n}(t, \bar{x}) \, \leq \, \frac{t}{2 \, \varepsilon_n} \, {}^t \nabla \varphi(\bar{x}) \cdot \partial_\theta W(\bar{x}, \bar{\theta}) \, + \, C \, , \quad \forall \, n \in \mathbb{N} \, .$$

For *n* large enough, the right hand side becomes negative. This is not compatible with (1.2.8). This means that the case (1.2.12) must be excluded. Now, because the function  $\theta \longmapsto W(x, \theta)$  is periodic, we have

$$\int_0^1 {}^t \nabla \varphi(x) \cdot \partial_\theta W(x,\theta) \ d\theta = {}^t \nabla \varphi(x) \cdot W(x,1) - {}^t \nabla \varphi(x) \cdot W(x,0) = 0.$$

Combining this with (1.2.11), we see that the restriction (1.2.10) is necessary.

Below, up to the end of the proof of Proposition 1.2.3, we select  $(x, \theta) \in \mathcal{V}$  such that  ${}^t \nabla \varphi(x) \cdot \partial_{\theta} W(x, \theta) = 0$ . Introduce the notations

 $\tilde{e}_1 := \partial_{\theta} W(x, \theta), \qquad \tilde{e}_d := {}^t \nabla \varphi(x), \qquad {}^t \tilde{e}_1 \cdot \tilde{e}_d = 0.$ 

We can complete  $\tilde{e}_1$  and  $\tilde{e}_d$  into some orthonormal basis  $(\tilde{e}_1, \tilde{e}_2, \cdots, \tilde{e}_{d-1}, \tilde{e}_d)$  of  $\mathbb{R}^d$ . In this special basis, the matrix  $\mathbf{I} + t D_x W(x, \theta)$  looks like :

.

$$\mathbf{I} + t \ D_x W(x, \theta) = \begin{pmatrix} \tilde{m}_{11}^1 & \dots & \tilde{m}_{1(d-1)}^1 & \tilde{m}_{1d}^1 \\ \tilde{m}_{21}^1 & \dots & \tilde{m}_{2(d-1)}^1 & \tilde{m}_{2d}^1 \\ \vdots & & \vdots & \\ \tilde{m}_{d1}^1 & \dots & \tilde{m}_{d(d-1)}^1 & \tilde{m}_{dd}^1 \end{pmatrix}$$

We can extract the  $(d-1) \times (d-1)$  matrix :

$$\mathcal{L}(t, x, \theta) = \begin{pmatrix} \tilde{m}_{21}^1 & \dots & \tilde{m}_{2(d-1)}^1 \\ \vdots & & \vdots \\ \tilde{m}_{d1}^1 & \dots & \tilde{m}_{d(d-1)}^1 \end{pmatrix}$$

Observe that  $\mathcal{L}$  is the realisation (in some specific basis) of the linear application :

$$\begin{aligned} \mathcal{L} \, : \, \nabla \varphi(x)^{\perp} & \longrightarrow \quad \partial_{\theta} W(x,\theta)^{\perp} \\ \mathbf{u} & \longmapsto \quad \Pi_{\partial_{\theta} W(x,\theta)^{\perp}} \left( \mathbf{I} \, + \, t \, D_{x} W(x,\theta) \right) \, \mathbf{u} \, . \end{aligned}$$

**Proposition 1.2.3.** The family  $\{h^{\varepsilon}\}_{\varepsilon}$  can be compatible only if there is T > 0 such that for all  $t \in [0, T]$ , we have :

$$(1.2.13) \qquad \qquad (-1)^d \det \mathcal{L}(t, x, \theta) \ge 0.$$

**Proof.** Assume again that the restriction (1.2.8) is satisfied for some T > 0 and some c > 0. We already know that (1.2.10) is verified. In the basis  $(\tilde{e}_1, \dots, \tilde{e}_d)$  of  $\mathbb{R}^d$ , the matrices  $M^0$  and  $M^1$  take the form

$$M^{0} = \begin{pmatrix} 0 & \cdots & 0 & t \ |\nabla\varphi|^{2} \\ 0 & \cdots & 0 & 0 \\ \vdots & & \vdots & \vdots \\ 0 & \cdots & 0 & 0 \end{pmatrix}, \qquad M^{1} = \begin{pmatrix} m_{11}^{1} & \cdots & m_{1d}^{1} \\ m_{21}^{1} & \cdots & m_{2d}^{1} \\ \vdots & & \vdots \\ m_{d1}^{1} & \cdots & m_{dd}^{1} \end{pmatrix}.$$

It follows that

$$\det D_x \mathbb{X}^{\varepsilon}(t,x) = \varepsilon^{-1} (-1)^d t |\nabla \varphi(x)|^2 \det M^{\sharp}\left(t,x,\frac{\varphi(x)}{\varepsilon}\right) \\ +1 + t g^{\varepsilon}\left(t,x,\frac{\varphi(x)}{\varepsilon}\right)$$

with

$$M^{\sharp}(t,x,\theta) = \begin{pmatrix} m_{21}^{1} & \dots & m_{2(d-1)}^{1} \\ \vdots & & \vdots \\ m_{d1}^{1} & \dots & m_{d(d-1)}^{1} \end{pmatrix} \equiv \Pi_{\partial_{\theta}W(x,\theta)^{\perp}} M^{1} \Pi_{\nabla\varphi(x)^{\perp}}$$

whereas the scalar application  $g^{\varepsilon}(t, x, \theta)$  is a continuous function of all the variables  $(\varepsilon, t, x, \theta) \in [0, 1] \times \mathbb{R} \times B(0, r] \times \mathbb{T}$ . Observe that

$$\begin{bmatrix} \mathbf{u} \otimes {}^t \nabla \varphi(x) \end{bmatrix} \mathbf{v} = 0, \qquad \forall (\mathbf{u}, \mathbf{v}) \in \mathbb{R}^d \times \nabla \varphi(x)^{\perp}.$$

Therefore, the expression of  $M^{\sharp}$  can be simplified according to

$$M^{\sharp}(t,x,\theta) = \mathcal{L}(t,x,\theta) \equiv \Pi_{\partial_{\theta}W(x,\theta)^{\perp}} \left( \mathbf{I} + t \ D_{x}W(x,\theta) \right) \Pi_{\nabla\varphi(x)^{\perp}}$$

Follow the argument of the preceding proof, using a well adjusted sequence  $\{\varepsilon_n\}_n$ , in order to extract the necessary condition

$$(-1)^d \det M^{\sharp}(t, x, \theta) \ge 0, \qquad \forall (t, x, \theta) \in [0, T] \times B(0, r] \times \mathbb{T}$$

which is exactly (1.2.13).

Remark 2.1. In the basis  $(\tilde{e}_1, \cdots, \tilde{e}_d)$ , we can get the decomposition

$$\mathcal{L}(t,x,\theta) = \mathcal{L}_0 + t \,\tilde{\mathcal{L}}(x,\theta), \qquad \mathcal{L}_0 := \begin{pmatrix} 0 & 1 & 0 & \cdots & 0 & 0 \\ 0 & 0 & 1 & \cdots & 0 & 0 \\ \vdots & \vdots & \vdots & \ddots & \vdots & \vdots \\ 0 & 0 & 0 & \cdots & 0 & 1 \\ 0 & 0 & 0 & \cdots & 0 & 0 \end{pmatrix}$$

with  $\tilde{\mathcal{L}}(x,\theta) \equiv \prod_{\partial_{\theta}W(x,\theta)^{\perp}} D_x W(x,\theta) \prod_{\nabla \varphi(x)^{\perp}}$ . This special structure implies the existence of coefficients  $\alpha_j(x,\theta)$  such that

$$(-1)^d \det \mathcal{L}(t, x, \theta) = \sum_{j=1}^{d-1} \alpha_j(x, \theta) t^j.$$

Noting

$$J(x,\theta) := \left\{ \begin{array}{ll} \min \mathcal{J} & \text{if} \quad \mathcal{J} := \left\{ j \; ; \; \alpha_j(x,\theta) \neq 0 \right\} \neq \emptyset \, , \\ d-1 & \text{if} \quad \mathcal{J} = \emptyset \, , \end{array} \right.$$

the condition (1.2.13) is equivalent to the restriction

(1.2.14) 
$$\alpha_{J(x,\theta)}(x,\theta) \ge 0, \qquad \forall (x,\theta) \in B(0,r] \times \mathbb{T}.$$

On the one hand, especially when  $d \gg 1$ , the conditions (1.2.14) can be technically difficult to deal with. On the other hand, nothing guarantees that they are intrinsic. Instead of looking at (1.2.14), we will consider

(1.2.15) 
$$\Pi_{\partial_{\theta}W(x,\theta)^{\perp}} D_x W(x,\theta) \Pi_{\nabla\varphi(x)^{\perp}} = 0, \qquad \forall (x,\theta) \in \mathcal{V}.$$

This relation is clearly intrinsic and, if it is satisfied, we are sure that

 $\det \mathcal{L}(t, x, \theta) = \det \mathcal{L}_0 = 0.$ 

We can summarize the preceding discussion by :

**Proposition 1.2.4.** Suppose that the relations (1.2.10) and (1.2.15) are verified. Then, the family  $\{h^{\varepsilon}\}_{\varepsilon}$  is compatible.

**Proof.** Under conditions (1.2.10) and (1.2.15), it remains

 $\det D_x \mathbb{X}^{\varepsilon}(t,x) = 1 + t \ g^{\varepsilon}\left(t,x,\frac{\varphi(x)}{\varepsilon}\right), \qquad g^{\varepsilon} \in C^0\left([0,1] \times \mathbb{R} \times B(0,r] \times \mathbb{T}; \mathbb{R}\right).$ 

In particular, we get :

$$\det D_x \mathbb{X}^{\varepsilon}(t,x) \ge 1 - C(T) t, \qquad \forall (t,x,\varepsilon) \in [0,T] \times B(0,r] \times ]0,1]$$

where the function  $T \mapsto C(T)$  is increasing. Now, it suffices to choose T > 0 small enough to recover (1.2.8).

Remark 2.2. Suppose that  $V : \mathbb{R}^d \longrightarrow \mathbb{R}^d$  is a  $\mathcal{C}^1$  diffeomorphism. Then, it is equivalent to solve (1.1.1) or

(1.2.16) 
$$\partial_t \mathbf{w} + (\mathbf{w} \cdot \nabla_x) \mathbf{w} = 0, \quad \mathbf{w} := V \circ \mathbf{u}$$

completed with the initial data

(1.2.17) 
$$\mathbf{w}(0,x) = W\left(x,\frac{\varphi(x)}{\varepsilon}\right), \qquad \varepsilon \in \left]0,1\right].$$

The system (1.1.11) can also be interpreted as a compatibility condition in order to solve the Cauchy problem (1.2.16)-(1.2.17) in the class of  $\mathcal{C}^1$  solutions, locally in time, on some domain  $\Omega^T$  with T > 0 independent of  $\varepsilon \in [0, 1]$ . This interpretation explains why the relevant constraint is concerned with  $V \circ H$  instead of dealing separately with V and H.  $\bigtriangleup$ 

From now on, we consider functions  $\varphi$  and W satisfying (1.2.10) and (1.2.15). In other words, we will concentrate on well prepared couples ( $\varphi$ , W).

### **1.3** Existence of compatible families

The goal of this subsection is to show through a constructive proof that the system (1.1.11) actually admits non trivial solutions. We want also to understand the structure of its generic solutions.

Of course, to face (1.1.11), preliminary assumptions are needed. We select some phase  $\varphi \in C^2(B(0, r]; \mathbb{R})$  with no critical point in B(0, r]. Without loss of generality (relabelling the coordinates and diminishing r if necessary) we can adjust  $\varphi$  so that

(1.3.1) 
$$\partial_d \varphi(x) \neq 0, \quad \forall x \in B(0,r], \quad r > 0.$$

 $\triangle$ 

We take  $W = V \circ H \in \mathcal{C}^2(B(0, r] \times \mathbb{T}, \mathbb{R}^d)$ . Introduce the linear subspace of  $\mathbb{R}^d$  spanned by the vectors  $\partial_{\theta} W(x, \theta)$  with  $\theta \in \mathbb{T}$ , that is

(1.3.2) 
$$\mathbf{E}(x) := \left\{ \sum_{j=1}^{N} \mu_j \ \partial_{\theta} W(x, \theta_j); \quad (\mu_1, \cdots, \mu_N) \in \mathbb{R}^N, \\ (\theta_1, \cdots, \theta_N) \in \mathbb{T}^N, \ N \in \mathbb{N} \right\}$$

Choose N = 1,  $\mu_1 = 1$  and  $\theta_1 = \theta$  in this definition to see that

 $\partial_{\theta} W(x,\theta) \in \mathbf{E}(x) \subset \mathbb{R}^d, \qquad \forall (x,\theta) \in B(0,r] \times \mathbb{T}.$ 

Because  $\mathbf{E}(x)$  is of finite dimension, we can find  $J^x$  numbers  $\theta_1^x, \dots, \theta_{J^x}^x$  such that

$$\mathbf{E}(x) = \left\{ \sum_{j=1}^{J_x} \mu_j \ \partial_\theta W(x, \theta_j^x); \ (\mu_1, \cdots, \mu_{J^x}) \in \mathbb{R}^{J^x} \right\}, \qquad J^x := \dim \mathbf{E}(x).$$

Then, in view of the first line of (1.1.11), we must have

 $\mathbf{E}(x) \subset \nabla \varphi(x)^{\perp} \,, \qquad \forall \, (x,\theta) \in B(0,r] \times \mathbb{T} \,.$ 

On the one hand, the case  $J^x = \dim \mathbf{E}(x) = 0$  is not interesting because this situation corresponds to the absence of oscillations. On the other hand, we have necessarily

$$J^x \leq \dim \nabla \varphi(x)^{\perp} = d - 1, \qquad \forall x \in B(0, r].$$

Due to the continuity of  $\partial_{\theta} W$ , the application  $x \mapsto \dim \mathbf{E}(x)$  is lower semicontinuous. In particular, the set

$$\left\{ x \in B(0, r[; J^x > d - \frac{3}{2}) \right\} = \left\{ x \in B(0, r[; J^x = d - 1) \right\}$$

is open. Now, suppose that  $J^0 = d - 1$ . By restricting r > 0 if necessary, we can always suppose that  $J^x = d - 1$  for all  $x \in B(0, r[$ . More generally, in what follows, we will retain the generic case where the application  $x \mapsto J^x = \dim \mathbf{E}(x)$  is constant (not necessarily equal to d - 1) on B(0, r]:

(1.3.3) 
$$\exists J \in \{1, \cdots, d-1\}; \quad \dim \mathbf{E}(x) = J, \quad \forall x \in B(0, r].$$

Denote by the symbol  $\mathcal{G}_d^J$  the *Grassmanian* manifold of linear subspaces of  $\mathbb{R}^d$  with dimension J.

**Lemme 1.3.1.** Assume that  $W \in C^2(B(0, r[\times \mathbb{T}, \mathbb{R}^d) \text{ and } (1.3.3))$ . Then  $\mathbf{E} \in C^1(B(0, r[, \mathcal{G}^J_d))$ .

**Proof.** Let  $x_0 \in B(0, r]$ . By hypothesis, we can find  $\theta_1^{x_0}, \dots, \theta_J^{x_0}$  in  $\mathbb{T}$  such that  $(\partial_{\theta} W(x_0, \theta_1^{x_0}), \dots, \partial_{\theta} W(x_0, \theta_J^{x_0}))$  is a basis of  $\mathbf{E}(x_0)$ . Hence, we can extract a  $J \times J$  determinant

$$\delta(x_0) := \det \left( \partial_{\theta} W_{i_j}(x_0, \theta_k^{x_0}) \right)_{1 \le j, k \le J}, \qquad i_j \in \llbracket 1, d \rrbracket$$

such that  $\delta(x_0) \neq 0$ . Since  $\partial_{\theta} W$  is continuous, the function  $x \mapsto \delta(x)$  is continuous. Therefore, we can isolate some small open neighborhood  $\Omega$  of  $x_0$  such that

 $\delta(x) \neq 0, \quad \forall x \in \Omega, \quad x_0 \in \Omega.$ 

For  $x \in \Omega$ , the family  $(\partial_{\theta} W(x, \theta_1^{x_0}), \cdots, \partial_{\theta} W(x, \theta_J^{x_0}))$  is still linearly independent and it is built with J vectors of  $\mathbf{E}(x)$ . Since by hypothesis  $\mathbf{E}(x)$  is of dimension J, this is in fact a basis of  $\mathbf{E}(x)$ . Obviously, the application

$$x \longmapsto \left(\partial_{\theta} W(x, \theta_1^{x_0}), \cdots, \partial_{\theta} W(x, \theta_J^{x_0})\right)$$

is of class  $C^1$  in  $\Omega$ . This remark gives the expected local regularity of **E**. Finally, since  $x_0 \in B(0, r]$  can be chosen arbitrarily, the Lemma 1.3.1 is proved.

Recall that

(1.3.4) 
$$\partial_{\theta} W(x,\theta) \in \mathbf{E}(x) \subset \nabla \varphi(x)^{\perp}, \quad \forall (x,\theta) \in B(0,r] \times \mathbb{T}.$$

The second line of (1.1.11) implies that

$$\Pi_{\mathbf{E}(x)^{\perp}} D_x W(x,\theta) \Pi_{\nabla \varphi(x)^{\perp}} = 0, \qquad \forall (x,\theta) \in B(0,r] \times \mathbb{T}.$$

Observe that, in this formulation, the two projectors (on the left and on the right) do not depend any more on the variable  $\theta \in \mathbb{T}$ . This allows to extract the mean value to get

(1.3.5) 
$$\Pi_{\mathbf{E}(x)^{\perp}} D_x W^*(x,\theta) \Pi_{\nabla \varphi(x)^{\perp}} = 0, \qquad \forall (x,\theta) \in B(0,r] \times \mathbb{T}.$$

**Lemme 1.3.2.** Let  $\varphi \in C^2(B(0,r],\mathbb{R})$  and  $W \in C^1(B(0,r] \times \mathbb{T},\mathbb{R}^d)$  satisfying respectively the conditions (1.3.1) and (1.3.3). Suppose that the relations (1.3.4) and (1.3.5) are satisfied. Then, the application  $x \mapsto \mathbf{E}(x)$  is constant on the level surfaces of  $\varphi$ . More precisely

(1.3.6) 
$$\exists \mathbb{E} \in \mathcal{C}^1(\mathbb{R}, \mathcal{G}_d^J); \qquad \mathbf{E}(x) = \mathbb{E} \circ \varphi(x), \quad \forall x \in B(0, r].$$

**Proof.** Let us denote  $\delta_{ij}$  the usual Dirichlet symbol, and  $\delta^{(k)}$  the vector of  $\mathbb{R}^d$  whose components are  $(\delta_{ik})_{1 \leq i \leq d}$ . The d-1 vectors

$$v_k(x) = -\delta^{(k)} + \partial_k \varphi(x) / \partial_d \varphi(x) \,\delta^{(d)}, \qquad 1 \le k \le d-1$$

form a  $\mathcal{C}^1$  basis of  $\nabla \varphi(x)^{\perp}$ . By permutting the components of  $\mathbb{R}^d$  and by diminishing r if necessary, we can always arrange the datas so that

 $\mathbf{E}(x) \oplus \langle v_1(x), \cdots, v_{d-J-1}(x) \rangle = \nabla \varphi(x)^{\perp}, \qquad \forall x \in B(0, r].$ 

Therefore, for all  $j \in [\![1, J]\!]$ , the vector  $v_{d-j}(x) \in \nabla \varphi(x)^{\perp}$  can be decomposed according to

$$v_{d-j}(x) = e_j(x) - \sum_{k=1}^{d-J-1} \alpha_j^k(x) \ v_k(x), \qquad e_j(x) \in \mathbf{E}(x)$$

where, due to the assumptions related to the regularity of  $\varphi$  and **E**, we have

$$e_j = (e_j^1, \cdots, e_j^d) \in \mathcal{C}^1(B(0, r]; \mathbb{R}^d), \qquad \alpha_j^k \in \mathcal{C}^1(B(0, r]; \mathbb{R}).$$

The vectors  $e_j$  with  $j \in [\![1, J]\!]$  are necessarily independent. They form a basis of  $\mathbf{E}(x)$ . Besides, we have the general formula

$$W(x,\theta) = \bar{W}(x) + \int_0^\theta \partial_\theta W(x,\tilde{\theta}) \, d\tilde{\theta} - \int_{\mathbb{T}} \left( \int_0^\theta \partial_\theta W(x,\tilde{\theta}) \, d\tilde{\theta} \right) \, d\theta$$

that, in view of (1.3.4), implies

$$W(x,\theta) = \overline{W}(x) + \sum_{j=1}^{J} w_j^*(x,\theta) e_j(x), \qquad w_j^* \in \mathcal{C}^1(B(0,r] \times \mathbb{T};\mathbb{R}).$$

Now, the relation (1.3.5) reads

$$\sum_{j=1}^{J} w_j^*(x,\theta) \Pi_{\mathbf{E}(x)^{\perp}} D_x e_j(x) \Pi_{\nabla \varphi(x)^{\perp}} = 0, \qquad \forall (x,\theta) \in B(0,r] \times \mathbb{T}.$$

Recall that the dimension of  $\mathbf{E}(x)$  is J. This implies that

$$\exists \left(\theta_{1}^{x}, \cdots, \theta_{J}^{x}\right) \in \mathbb{T}^{J}; \qquad \det \left[w_{i}^{*}\left(x, \theta_{j}^{x}\right)\right]_{1 \leq i, j \leq J} \neq 0$$

Combining the preceding informations, we see that (1.3.5) is equivalent to

(1.3.7) 
$$\Pi_{\mathbf{E}(x)^{\perp}} D_x e_j(x) \Pi_{\nabla \varphi(x)^{\perp}} = 0, \qquad \forall (j, x) \in \llbracket 1, J \rrbracket \times B(0, r].$$

The vector space  $\mathbf{E}(x)^{\perp}$  is of dimension d - J. It is generated by the vector  $e_d(x) := \nabla \varphi(x)$  and the d - J - 1 vectors

$$e_j(x) = -\delta^{(j-J)} + \sum_{k=1}^J \alpha_k^{j-J}(x) \,\,\delta^{(d-k)}, \qquad j \in [\![J+1, d-1]\!]$$

Therefore (1.3.7) is exactly the same as

(1.3.8) 
$${}^{t}e_{l}(x) D_{x}e_{j}(x) \Pi_{\nabla\varphi(x)^{\perp}} = 0, \quad \forall (l, j, x) \in [\![J+1, d]\!] \times [\![1, J]\!] \times B(0, r] .$$

For  $j \in [\![1, J]\!]$ , compute

$$D_{x}e_{j}(x) = \sum_{k=1}^{d-J-1} \nabla_{x}\alpha_{j}^{k}(x) v_{k}(x) + \left[\sum_{k=1}^{d-J-1} \alpha_{j}^{k}(x) \nabla_{x} \left(\partial_{k}\varphi(x)/\partial_{d}\varphi(x)\right) + \nabla_{x} \left(\partial_{d-j}\varphi(x)/\partial_{d}\varphi(x)\right)\right] \delta^{(d)}.$$

Applying on the left  ${}^{t}e_{l}(x)$  with  $l \in [[J+1, d-1]]$ , yields

$${}^{t}e_{l}(x) D_{x}e_{j}(x) = \nabla_{x}\alpha_{j}^{l-J}(x), \qquad 1 \le j \le J < l \le d-1.$$

We can extract from (1.3.8) that

(1.3.9) 
$$\nabla_x \alpha_j^{l-J}(x) \Pi_{\nabla \varphi(x)^{\perp}} = 0, \quad \forall (l, j, x) \in \llbracket J + 1, d - 1 \rrbracket \times \llbracket 1, J \rrbracket \times B(0, r].$$

**Independent statement.** Let  $\varphi \in C^1(B(0,r],\mathbb{R})$  satisfying (1.3.1). Let  $\alpha \in C^1(B(0,r],\mathbb{R})$  a function which is subjected to the relation (1.3.9). Then, restricting r > 0 if necessary, we can always find some function  $Z \in C^1(\mathbb{R},\mathbb{R})$  such that

(1.3.10) 
$$\alpha(x) = Z \circ \varphi(x), \quad \forall x \in B(0, r].$$

**Proof of the independent statement.** The geometric reason of (1.3.10) is the following. The relation (1.3.9) means that the vectors  $\nabla_x \alpha(x)$  and  $\nabla \varphi(x)$  are parallel or that the tangent spaces at x to the level surfaces associated with the scalar functions  $\alpha$  and  $\varphi$  coincide. Since the level surfaces associated with  $\alpha$  and  $\varphi$  are spanned by these tangent spaces, we can deduce that  $\alpha$  and  $\varphi$  have common level surfaces. Moreover, the relation (1.3.1) allows to characterize (locally near 0) these level surfaces through the different values of  $\varphi$ . This is why we have (1.3.10).

Now, we can also proceed as follows. Due to (1.3.1), the functions  $x_1, x_2, \dots, x_{d-1}$  and  $\varphi(x)$  form locally (near 0) a system of coordinates. Therefore, we can find  $Z \in \mathcal{C}^1(\mathbb{R}^d, \mathbb{R})$  such that

$$\alpha(x) = Z(\hat{x}, \varphi(x)), \qquad \hat{x} := (x_1, x_2, \cdots, x_{d-1}), \qquad \forall x \in B(0, r].$$

Decompose  $\nabla \varphi(x)$  according to

$$\nabla\varphi(x) = \left(\nabla_{\hat{x}}\varphi(x), \partial_d\varphi(x)\right), \qquad \nabla_{\hat{x}}\varphi(x) = \left(\partial_1\varphi(x), \cdots, \partial_{d-1}\varphi(x)\right) \in \mathbb{R}^{d-1}.$$

Given  $\hat{h} \in \mathbb{R}^{d-1}$ , define

$$h_d(x, \hat{h}) := -\partial_d \varphi(x)^{-1} \nabla_{\hat{x}} \varphi(x) \cdot \hat{h}.$$

Observe that

$$(\hat{h}, h_d(x, \hat{h})) \in \nabla \varphi(x)^{\perp}, \quad \forall \hat{h} \in \mathbb{R}^{d-1}$$

Testing (1.3.9) with such choices gives rise to

$$abla_{\hat{x}} Zig(x_1, x_2, \cdots, x_{d-1}, arphi(x)ig) \cdot \hat{h} \,=\, 0\,, \qquad orall \, \hat{h} \in \mathbb{R}^{d-1}\,.$$

This information clearly implies that the function Z does not depend on its d-1 first variables. We have (1.3.10).

Applying the independent statement to the functions  $\alpha_j^{l-J}$ , we see that we can exhibit functions

$$Z_j^k \in \mathcal{C}^1(\mathbb{R}, \mathbb{R}), \qquad (k, j) \in \llbracket 1, d - J - 1 \rrbracket \times \llbracket 1, J \rrbracket$$

such that, for all  $(k, j) \in [\![1, d - J - 1]\!] \times [\![1, J]\!]$ , we have

(1.3.11) 
$$e_j^k(x) = -\alpha_j^k(x) = Z_j^k \circ \varphi(x), \qquad \forall x \in B(0,r].$$

This construction of the  $Z_j^k$  is not classical and it is one of the main difficulties in the proof of Lemma 1.3.2. Finally, the remaining conditions to consider are obtained by taking  $j \in [\![1, J]\!]$  and l = d. Namely

$$\nabla \varphi(x) \ D_x e_j(x) \ \Pi_{\nabla \varphi(x)^\perp} = \ 0 \,, \qquad \forall \, (j,x) \in \llbracket 1, J \rrbracket \times B(0,r] \,.$$

Use (1.3.1) and (1.3.11) to simplify this into

$$\nabla_x e_j^d(x) \ \Pi_{\nabla \varphi(x)^{\perp}} = 0, \qquad \forall (j, x) \in \llbracket 1, J \rrbracket \times B(0, r]$$

where we recall that

$$e_j^d(x) = -\sum_{k=1}^{d-J-1} Z_j^k \circ \varphi(x) \,\partial_k \varphi(x) \,/ \,\partial_d \varphi(x) + \,\partial_{d-j} \varphi(x) \,/ \,\partial_d \varphi(x) \,.$$

Again, this means the existence of  $Z_i^d \in \mathcal{C}^1(\mathbb{R}, \mathbb{R})$  such that

$$e_j^d(x) \,=\, Z_j^d \circ \varphi(x)\,, \qquad \forall\, (j,x) \in [\![1,J]\!] \times B(0,r]\,.$$

Briefly, we have obtained, for all  $j \in [\![1, J]\!]$ , that

$$e_j(x) = Z_j \circ \varphi(x), \qquad Z_j = {}^t(Z_j^1, \cdots, Z_j^{d-J-1}, 0, \cdots, 0, -1, 0, \cdots, 0, Z_j^d).$$

The vector space  $\mathbf{E}$  is spanned by the  $e_j$  with  $j \in [\![1, J]\!]$ . Therefore, it depends only on  $\varphi$ , in a  $\mathcal{C}^1$  way. This gives rise to (1.3.6). In particular,  $\mathbf{E}$  is constant on the level surfaces of  $\varphi$ .

Combining (1.3.4) and (1.3.6), we can produce the necessary condition

(1.3.12) 
$$\nabla \varphi(x) \in \mathbb{E} \circ \varphi(x)^{\perp} = \mathbf{E}(x)^{\perp}, \quad \forall x \in B(0, r].$$

The condition (1.3.12) is a geometrical constraint on  $\varphi$  underlying the resolution of (1.1.11). We explain below how to solve it.

Lemme 1.3.3. Select :

- a curve  $\mathbb{E} \in \mathcal{C}^2(\mathbb{R}, \mathcal{G}_d^J)$  of J-dimensional vector spaces of  $\mathbb{R}^d$ ,

- a  $\mathcal{C}^2$  submanifold  $\mathcal{S} \subset \mathbb{R}^d$  of dimension d - J, containing  $0 \in \mathbb{R}^d$ ,

- a  $\mathcal{C}^2$  scalar function  $\chi: \mathcal{S} \longrightarrow \mathbb{R}$ .

Note  $T_0 S$  the tangent space of S at the point  $0 \in \mathbb{R}^d$ . We suppose that

(1.3.13) 
$$T_0 \mathcal{S} + \mathbb{E}(\chi(0)) = \mathbb{R}^d$$

Then, we can find r > 0 such that the nonlinear equation (1.3.12) completed with  $\varphi_{|S \cap B(0,r]} \equiv \chi$  has a unique  $C^2$  solution on B(0,r]. We will say that the phase  $\varphi$  is <u>generated</u> by  $(\mathbb{E}, S, \chi)$ .

**Proof.** Select  $\delta > 0$  and J functions

 $Z_j \in \mathcal{C}^2(]\chi(0) - \delta, \chi(0) + \delta[; \mathbb{R}), \qquad j \in [\![1, J]\!]$ 

adjusted such that, for all  $t \in ]\chi(0) - \delta, \chi(0) + \delta[, (Z_1(t), \dots, Z_J(t))]$  is a basis of  $\mathbb{E}(t)$ . Note

$$\Omega^{\delta}_{\mathcal{S}} := \chi^{-1} \big( ]\chi(0) - \delta, \chi(0) + \delta[ \big) \subset \mathcal{S}, \qquad z = {}^t (z^1, \dots, z^J) \in \mathbb{R}^J.$$

Consider the  $\mathcal{C}^2$  application

$$\begin{split} \Xi \, : \, \Omega^{\delta}_{\mathcal{S}} \times \mathbb{R}^J & \longrightarrow \quad \mathbb{R}^d \\ (y,z) & \longmapsto \quad \Xi(y,z) \, := \, y \, + \, \sum_{j=1}^J \, z^j \, Z_j \circ \chi(y) \, . \end{split}$$

Because of (1.3.13), the linear operator

$$D_x \Xi(0,0) : T_0 \mathcal{S} \times \mathbb{R}^J \longrightarrow \mathbb{R}^d$$
  
(h,k)  $\longmapsto h + \sum_{j=1}^J k^j Z_j \circ \chi(y)$ 

is invertible. The inverse mapping Theorem can be applied at the point  $(0,0) \in S \times \mathbb{R}^J$ . It guarantees the existence of r > 0 such that  $\Xi$  is a  $\mathcal{C}^2$  diffeomorphism from a neighbourhood of  $(0,0) \in S \times \mathbb{R}^J$  onto B(0,r]. Introduce the projection

$$\begin{array}{cccc} \Gamma \,:\, \mathcal{S} \times \mathbb{R}^J & \longrightarrow & \mathcal{S} \\ & (y,z) & \longmapsto & \Gamma(y,z) := \end{array}$$

Now, we can define

$$\varphi := \chi \circ \Gamma \circ \Xi^{-1} \in \mathcal{C}^2(B(0,r];\mathbb{R}))$$

Since  $(\Gamma \circ \Xi^{-1})_{|S \cap B(0,r]} = Id$ , we have  $\varphi_{|S \cap B(0,r]} \equiv \chi_{|S \cap B(0,r]}$ . Moreover, the function  $\varphi$  is constant on the set

$$\mathcal{F}_y := \left( y + \langle Z_1 \circ \chi(y), \cdots, Z_J \circ \chi(y) \rangle \right) \cap B(0, r], \qquad y \in \mathcal{S} \cap B(0, r].$$

y .

More precisely,  $\mathcal{F}_y$  is a piece of affine manifold with direction  $\mathbb{E} \circ \chi(y)$ , on which  $\varphi$  takes the value  $\chi(y)$ . In particular

$$\nabla \varphi(x) \in (T_x \mathcal{F}_y)^{\perp} = \mathbb{E} \circ \chi(y)^{\perp} = \mathbb{E} \circ \varphi(x)^{\perp}, \quad \forall x \in \mathcal{F}_y.$$

Since the  $\mathcal{F}_y$  with  $y \in S \cap B(0, r]$  form a foliation of B(0, r], we have obtained the expected relation (1.3.12).

**Proposition 1.3.1.** Let  $\varphi$  be generated by  $(\mathbb{E}, S, \chi)$ . The couple  $(\varphi, W)$  is well prepared if and only if there exist two functions  $W_{\parallel} \in C^1(\mathbb{R}^2, \mathbb{R}^d)$  and  $W_{\perp} \in C^1(\mathbb{R}, \mathbb{R}^d)$ satisfying

(1.3.14) 
$$W_{\parallel}(t,s) \in \mathbb{E}(t), \qquad W_{\perp}(t) \in \mathbb{E}(t)^{\perp}, \qquad \forall (t,s) \in \mathbb{R}^2$$

and a scalar function  $\psi \in \mathcal{C}^1(B(0,r] \times \mathbb{T};\mathbb{R})$  such that

(1.3.15) 
$$W(x,\theta) = W_{\parallel}(\varphi(x),\psi(x,\theta)) + W_{\perp}(\varphi(x)), \quad \forall (x,\theta) \in B(0,r] \times \mathbb{T}.$$

**Proof.** Note  $(Z_1(t), \dots, Z_J(t))$  some orthonormal basis of  $\mathbb{E}(t)$  with a  $\mathcal{C}^1$  regularity with respect to  $t \in \mathbb{R}$ . Complete it with some  $\mathcal{C}^1$  orthonormal basis  $(e_{J+1}(t), \dots, e_d(t))$  of  $\mathbb{E}(t)^{\perp}$ , again of class  $\mathcal{C}^1$ . In view of (1.3.4), the definition of  $\mathbf{E}(x)$  and Lemma 1.3.2, the profile  $W(x, \theta)$  can be decomposed according to

$$W(x,\theta) = \sum_{k=1}^{J} w_j(x,\theta) Z_j \circ \varphi(x) + \sum_{k=J+1}^{d} w_j(x) e_j \circ \varphi(x)$$

with

$$w_j \in \mathcal{C}^1(B(0,r] \times \mathbb{T};\mathbb{R}), \quad \forall j \in \llbracket 1, J \rrbracket,$$

$$w_j \in \mathcal{C}^1(B(0,r];\mathbb{R}), \qquad \forall j \in \llbracket J+1,d \rrbracket.$$

Compute the derivative of  $W(x,\theta)$  with respect to the variable x and compose on the right with  $\Pi_{\nabla\varphi(x)^{\perp}}$ . It remains

$$D_{x}W(x,\theta) \Pi_{\nabla\varphi(x)^{\perp}} = \sum_{k=1}^{J} \nabla_{x}w_{j}(x,\theta) \cdot \Pi_{\nabla\varphi(x)^{\perp}} \times Z_{j} \circ \varphi(x) + \sum_{k=J+1}^{d} \nabla_{x}w_{j}(x) \cdot \Pi_{\nabla\varphi(x)^{\perp}} \times e_{j} \circ \varphi(x).$$

Select a point  $(x, \theta) \in \mathcal{V}$  which means that  $\partial_{\theta} W(x, \theta) \neq 0$ . Without loss of generality, we can suppose that  $\partial_{\theta} W_J(x, \theta) \neq 0$ . Otherwise, just permute the components of  $\mathbb{R}^d$  to obtain this condition. By construction, the hyperplane  $\partial_{\theta} W(x, \theta)^{\perp}$  is generated by the d - J vectors  $e_j \circ \varphi(x)$  with  $j \in [J + 1, d]$  and the J - 1 vectors

$$\partial_{\theta} w_J(x,\theta) \ Z_j \circ \varphi(x) \ - \ \partial_{\theta} w_j(x,\theta) \ Z_J \circ \varphi(x) \,, \qquad j \in \llbracket 1, J - 1 \rrbracket \,.$$

The requirement (1.2.15) is equivalent to the conditions

(1.3.16) 
$$\nabla_x w_j(x) \cdot \Pi_{\nabla \varphi(x)^{\perp}} = 0, \qquad \forall j \in \llbracket J + 1, d \rrbracket,$$

(1.3.17) 
$$(\partial_{\theta} w_J \nabla_x w_j - \partial_{\theta} w_j \nabla_x w_J)(x,\theta) = 0, \qquad \forall j \in [\![1, J-1]\!].$$

On the one hand, from (1.3.16), we deduce that

$$\exists \tilde{w}_j \in \mathcal{C}^1(\mathbb{R}, \mathbb{R}); \qquad w_j(x) = \tilde{w}_j \circ \varphi(x), \qquad \forall j \in \llbracket J + 1, d \rrbracket.$$

On the other hand, it follows from the relations (1.3.17) that the mappings  $\Upsilon_t$  parameterized by  $t \in \mathbb{R}$  and defined on the level sets

$$\mathcal{G}_t := \left\{ x \in B(0, r] \, ; \, \varphi(x) = t \right\}$$

by the formulas

$$\begin{array}{cccc} \Upsilon_t &: \ \mathcal{G}_t \times \mathbb{T} & \longrightarrow & \mathbb{R}^J \\ & (x, \theta) & \longmapsto & {}^t(w_1, \cdots, w_J) \end{array}$$

have rank one. Thus, to each  $\Upsilon_t$  corresponds a foliation of  $\mathcal{G}_t \times \mathbb{T}$  by submanifolds of dimension d-1. This foliation depends on the parameter t. It can be described by using a function  $\psi \in \mathcal{C}^1(B(0,r] \times \mathbb{T}, \mathbb{R})$  so that

$$w_j(x,\theta) = \tilde{w}_j(\varphi(x),\psi(x,\theta)), \quad \forall j \in \llbracket 1, J \rrbracket.$$

Define

$$W_{\perp}(t) := \sum_{j=J+1}^{d} \tilde{w}_{j}(t) e_{j}(t), \qquad W_{\parallel}(t,s) := \sum_{j=1}^{J} \tilde{w}_{j}(t,s) Z_{j}(t).$$

By construction, we have both (1.3.14) and (1.3.15).

Conversely, suppose that  $W(x, \theta)$  has the form (1.3.15) with  $W_{\parallel}(x, \theta)$  and  $W_{\perp}(x, \theta)$  as in (1.3.14). Then

$$\partial_{\theta} W(x,\theta) = \partial_{\theta} \psi(x,\theta) \times \partial_{s} W_{\parallel} \big( \varphi(x), \psi(x,\theta) \big) \in \mathbb{E} \big( \varphi(x) \big) \equiv \mathbf{E}(x)$$

which is (1.3.4) and gives rise to the first part of (1.1.11). Moreover

$$D_x W(x,\theta) \Pi_{\nabla \varphi(x)^{\perp}} = \nabla_x \psi(x,\theta) \cdot \Pi_{\nabla \varphi(x)^{\perp}} \times \partial_s W_{\parallel} (\varphi(x), \psi(x,\theta)) \,.$$

Since  $\partial_{\theta} W$  and  $\partial_s W_{\parallel}$  are collinear, we get the second equation of (1.1.11).

#### **1.4** Simple wave solutions

The aim of this last part is to explain how the initial oscillating data  $h^{\varepsilon}(x)$  is transformed through the evolution equation (1.1.1). Below, we consider this question in a simplified context, by looking only on simple wave solutions.

**Définition 1.4.1.** Let  $\varepsilon \in [0, 1]$ . We say that  $\mathbf{u}^{\varepsilon} \in \mathcal{C}^1(\Omega^T; \mathbb{R})$  is a <u>simple</u> wave if it can be put in the following form

$$\mathbf{u}^{\varepsilon}(t,x) = \mathbf{H}(t,x,\frac{\Phi(t,x)}{\varepsilon}), \qquad \mathbf{H} \in \mathcal{C}^{1}(\Omega^{T} \times \mathbb{T};\mathbb{R}^{d}), \quad \Phi \in \mathcal{C}^{1}(\Omega^{T};\mathbb{R}).$$

The Theorem 3 explains how to associate with a well prepared couple  $(\varphi, W)$  a simple wave  $\mathbf{u}^{\varepsilon}(t, x)$  which is a solution on  $\Omega^{T}$  of the Burger's type system (1.1.1). It remains to show this statement 3.

**Proof of Theorem 3.** Compose the first equation of (1.1.12) with  $D_{\mathbf{u}}V \circ \mathbf{H}$  in order to extract

(1.4.1) 
$$\begin{cases} \partial_t \mathbf{W} + (\mathbf{W} \cdot \nabla_x) \mathbf{W} = 0, \\ \partial_t \Phi + (\bar{\mathbf{W}} \cdot \nabla_x) \Phi = 0, \\ \mathbf{W}^* \cdot \nabla_x \Phi = 0, \end{cases} \quad \mathbf{W} := V \circ \mathbf{H}.$$

This must be associated with the initial data

(1.4.2) 
$$\mathbf{W}(0, x, \theta) = W(x, \theta), \qquad \Phi(0, x) = \varphi(x).$$

First, we discuss about (1.4.1)-(1.4.2). From Proposition 1.3.1 we can write

$$W(x,\theta) = W_{\parallel}(\varphi(x),\psi(x,\theta)) + W_{\perp}(\varphi(x))$$

Solve locally in time, say on  $\Omega^T$  for some T > 0, the scalar conservation law

(1.4.3) 
$$\partial_t \Phi + W_{\perp}(\Phi) \cdot \nabla_x \Phi = 0, \qquad \Phi(0, x) = \varphi(x).$$

Recall that  $\mathbf{E}(x) = \mathbb{E} \circ \varphi(x)$  is spanned by the *J* vectors  $e_j(x) = Z_j \circ \varphi(x)$  where the  $Z_j$  are defined at the end of the proof of Lemma 1.3.2. Now, fix any  $j \in [\![1, J]\!]$ and compute

 $\left[\partial_t + W_{\perp}(\Phi) \cdot \nabla_x\right] \left(Z_j \circ \Phi \cdot \nabla_x \Phi\right) = -\left(\nabla_x \Phi \cdot W'_{\perp} \circ \Phi\right) \times \left(Z_j \circ \Phi \cdot \nabla_x \Phi\right).$ 

Combining (1.1.13) and (1.3.12), we can extract

$$(Z_j \circ \Phi \cdot \nabla_x \Phi)(0, x) = 0, \qquad \forall (j, x) \in \llbracket 1, J \rrbracket \times B(0, r].$$

In view of the preceding equation, this polarization identity is propagated in time which means that

$$Z_j \circ \Phi(t, x) \cdot \nabla_x \Phi(t, x) = 0, \qquad \forall (t, x) \in [0, T] \times B(0, r]$$

or equivalently that

(1.4.4) 
$$\nabla_x \Phi(t,x) \subset \mathbb{E} \circ \Phi(t,x)^{\perp}, \qquad \forall (t,x) \in [0,T] \times B(0,r].$$

Now, introduce the function

$$\tilde{W}(t,s)\,:=\,W_{\parallel}(t,s)\,+\,W_{\perp}(t)\,,\qquad(t,s)\in\mathbb{R}^2$$

and the scalar conservation law

(1.4.5) 
$$\partial_t \Psi + \tilde{W}(\Phi(t,x),\Psi) \cdot \nabla_x \Psi = 0.$$

Complete (1.4.5) with the initial data

(1.4.6) 
$$\Psi(0, x, \theta) = \psi(x, \theta), \qquad \psi \in \mathcal{C}^1(B(0, r] \times \mathbb{T}; \mathbb{R}).$$

In (1.4.5), the variable  $\theta \in \mathbb{T}$  plays the part of a parameter. For T > 0 small enough, the Cauchy problem (1.4.5)-(1.4.6) has a local solution on  $\Omega^T$ . Finally, define the profile **W** through

$$\mathbf{W}(t,x,\theta) := \tilde{W}\big(\Phi(t,x), \Psi(t,x,\theta)\big), \qquad \mathbf{W}(0,x,\theta) = W(x,\theta).$$

By construction, we have

$$\mathbf{W}^*(t, x, \theta) = W_{\parallel} \big( \Phi(t, x), \Psi(t, x, \theta) \big)^*.$$

The informations (1.3.14) and (1.4.4) imply that

$$\mathbf{W}^*(t, x, \theta) \cdot \nabla_x \Phi(t, x) = 0, \qquad \forall (t, x) \in \Omega^T.$$

Taking into account (1.3.14) and (1.4.3), we have also

 $\partial_t \Phi \,+\, {\bf W} \cdot \nabla_x \Phi \,=\, \partial_t \Phi \,+\, W_\perp \circ \Phi \cdot \nabla_x \Phi \,=\, 0\,.$ 

Then, with (1.4.5), we can deduce that

(1.4.7) 
$$\partial_t \mathbf{W} + \mathbf{W} \cdot \nabla_x \mathbf{W} = \partial_s \tilde{W} \left( \partial_t \Psi + \mathbf{W} \cdot \nabla_x \Psi \right) = 0, \qquad \mathbf{W}(0, x, \theta) = W(x, \theta).$$

To sum up, we have constructed functions  $\Phi$  and **W** satisfying (1.4.1).

Now, we concentrate on (1.1.12). First, solve separately (on some domain  $\Omega^T$  with T > 0) the Cauchy problem

(1.4.8) 
$$\partial_t \mathbf{H} + V \circ \mathbf{H} \cdot \nabla_x \mathbf{H} = 0, \qquad \mathbf{H}(0, x, \theta) = H(x, \theta).$$

Observe that the expression  $\tilde{\mathbf{W}} := V \circ \mathbf{H}$  is by construction subjected to

(1.4.9) 
$$\partial_t \tilde{\mathbf{W}} + \tilde{\mathbf{W}} \cdot \nabla_x \tilde{\mathbf{W}} = 0, \qquad \tilde{\mathbf{W}}(0, x, \theta) = W(x, \theta).$$

The Cauchy problems (1.4.7) and (1.4.9) are made of the same quasilinear constraints and the same initial data. Since the corresponding  $C^1$  solutions must coincide, we have necessarily  $\tilde{\mathbf{W}} = V \circ \mathbf{H} \equiv \mathbf{W}$  on  $\Omega^T$ .

Briefly, the first equation of (1.1.12) is verified because this is precisely (1.4.8) whereas the two other conditions of (1.1.12) are satisfied because they correspond exactly to the two last conditions in (1.4.1). This explains why the apparently overdetermined system (1.1.12)-(1.1.13) has a unique solution on  $\Omega^T \times \mathbb{T}$  for some T > 0.

Finally, define the simple wave  $\mathbf{u}^{\varepsilon}(t,x) := \mathbf{H}(t,x,\frac{\Phi(t,x)}{\varepsilon})$ . Compute

$$\partial_{t}\mathbf{u}^{\varepsilon} + V(\mathbf{u}^{\varepsilon}) \cdot \nabla_{x}\mathbf{u}^{\varepsilon} = \left(\partial_{t}\mathbf{H} + V \circ \mathbf{H} \cdot \nabla_{x}\mathbf{H}\right)\left(t, x, \frac{\varphi(x)}{\varepsilon}\right) \\ + \frac{1}{\varepsilon} \left[\left(\partial_{t}\Phi + V \circ \mathbf{H} \cdot \nabla_{x}\Phi\right) \partial_{\theta}\mathbf{H}\right]\left(t, x, \frac{\varphi(x)}{\varepsilon}\right).$$

The fact that  $\mathbf{u}^{\varepsilon}(t, x)$  is a solution of (1.1.1) becomes a direct consequence of the equations inside (1.1.12). Moreover, the definition of  $\mathbf{W}$  indicates clearly that the structure (1.3.15) is conserved for  $t \in [0, T]$ . Therefore (see the end of the proof of Proposition 1.3.1), for all  $t \in [0, T]$ , the trace  $(\Phi(t, \cdot), \mathbf{W}(t, \cdot))$  is still well prepared. This last remark concludes the proof of Theorem 3.

# Chapter 2

# Large amplitude Oscillating solutions For incompressible Euler Equations in space dimension 3.

Abstract. In this article, we construct large amplitude oscillating waves  $(u^{\varepsilon})_{\varepsilon \in ]0,1]}$ which are *local solutions* on some open domain of the time-space  $\mathbb{R}_+ \times \mathbb{R}^3$  of both the three dimensional Burger equations (without source term) and the incompressible Euler equations (without pressure). The functions  $u^{\varepsilon}(t, x)$  are mainly characterized by the fact that the corresponding Jacobian matrices  $D_x u^{\varepsilon}(t, x)$  are nilpotent of rank one or two. Our purpose here is to describe the interesting geometrical features of the expressions  $u^{\varepsilon}(t, x)$  obtained by this way.

# 2.1 Detailed introduction.

# 2.1.1 Presentation of the framework.

Let  $(T, V, r) \in (\mathbb{R}^*_+)^3$  with  $T V \leq r$ . We work on a domain of determination having the form of a truncated cone like

$$\Omega_r^T := \left\{ (t, x) \in [0, T] \times \mathbb{R}^3; \ |x| + t \, V \le r \right\}, \qquad |x| := \sqrt{x_1^2 + x_2^2 + x_3^2}.$$

We are looking at expressions  $u^{\varepsilon}(t, x)$ , with  $\varepsilon \in ]0, 1]$ , which are special solutions of three dimensional Burger equations without source term, namely

(2.1.1) 
$$\partial_t u^{\varepsilon} + (u^{\varepsilon} \cdot \nabla) u^{\varepsilon} = 0, \qquad (t, x) \in \Omega_r^T \subset \mathbb{R} \times \mathbb{R}^3.$$

We complete (2.1.1) with a family of oscillating initial data

(2.1.2) 
$$u^{\varepsilon}(0,x) = h^{\varepsilon}(x) = \begin{pmatrix} h_1^{\varepsilon}(x) \\ h_2^{\varepsilon}(x) \\ h_3^{\varepsilon}(x) \end{pmatrix} = w\left(x, \frac{\varphi(x)}{\varepsilon}\right), \quad (x,\varepsilon) \in \Omega_r^0 \times ]0,1].$$

The function  $h^{\varepsilon}(x)$  is defined on the closed ball  $\Omega^0_r$  (having center zero and radius r) by using a bounded profile  $w(x,\theta) \in \mathcal{C}^1_b(\Omega^0_r \times \mathbb{T}; \mathbb{R}^3)$  satisfying

(2.1.3) 
$$\exists (x,\theta) \in \Omega^0_r \times \mathbb{T}; \qquad \partial_\theta w(x,\theta) \neq 0, \qquad \mathbb{T} := \mathbb{R}/\mathbb{Z}.$$

We use also a phase  $\varphi \in \mathcal{C}^1(\Omega^0_r; \mathbb{R})$  which is assumed to be not stationary

(2.1.4) 
$$\nabla\varphi(x) := {}^t \big(\partial_1\varphi(x), \partial_2\varphi(x), \partial_3\varphi(x)\big) \neq 0, \qquad \forall \ x \in \Omega^0_r$$

The equation (2.1.1) is the prototype of a quasilinear hyperbolic system. Thus, the solution  $u^{\varepsilon}(t, x)$  of (2.1.1) which is issued from the bounded initial data  $h^{\varepsilon}(x)$  inherits a finite speed of propagation V. In view of (2.1.2), noting  $w = {}^{t}(w_1, w_2, w_3) \in \mathbb{R}^3$ , we can take

$$V := \sup\left\{\left(\sum_{i=1}^{3} w_i(x,\theta)^2\right)^{1/2}; \ (x,\theta) \in \Omega^0_r \times \mathbb{T}\right\} < \infty.$$

**Example 1.** Choose T = V = r = 1. Select any non constant function  $w_3^e \in C^{\infty}(\mathbb{T};\mathbb{R})$  which is bounded by 1 and define

$$\varphi^{e}(x) := x_{1}, \quad w^{e}(x,\theta) := {}^{t}(0,0,w_{3}^{e}(\theta)), \quad u^{e\varepsilon}(x) := {}^{t}(0,0,w_{3}^{e}(\varphi^{e}(x)/\varepsilon)).$$

Observe that

(2.1.5) 
$$\partial_t u^{e\varepsilon} + (u^{e\varepsilon} \cdot \nabla) u^{e\varepsilon} \equiv 0, \quad \operatorname{div} u^{e\varepsilon} \equiv 0, \quad \left( D_x u^{e\varepsilon}(x) \right)^2 \equiv 0.$$

The expression  $u^{e\varepsilon}(x)$  is a very basic example of a contact discontinuity solution of (2.1.1). More elaborated patterns are proposed in [5, 6, 7, 13, 14]. Extensions can be obtained either by considering nonlinear phases  $\varphi$  or by adding some dependence in other variables than  $\varphi$ . In this article, we explain what can be done in these two directions. More precisely, we construct and classify all functions  $\varphi(x)$  and  $w(x, \theta)$  (if need be, the profile w can also depend in a smooth way on  $\varepsilon \in [0, 1]$ ) allowing to solve the oscillating Cauchy problem (2.1.1)-(2.1.2) in the class of  $\mathcal{C}^1$ -functions on a domain of determination  $\Omega_r^T$  with  $(T, r) \in \mathbb{R}^*_+ \times \mathbb{R}^*_+$  independent of  $\varepsilon \in [0, 1]$ . Note  $D_x h^{\varepsilon}$  the Jacobian matrix of  $h^{\varepsilon}$ , that is

$$D_x h^{\varepsilon}(x) = \begin{pmatrix} \partial_1 h_1^{\varepsilon}(x) & \partial_2 h_1^{\varepsilon}(x) & \partial_3 h_1^{\varepsilon}(x) \\ \partial_1 h_2^{\varepsilon}(x) & \partial_2 h_2^{\varepsilon}(x) & \partial_3 h_2^{\varepsilon}(x) \\ \partial_1 h_3^{\varepsilon}(x) & \partial_2 h_3^{\varepsilon}(x) & \partial_3 h_3^{\varepsilon}(x) \end{pmatrix} \in \mathcal{M}_3(\mathbb{R}^3).$$

Our starting point is the Theorem 2.6 of [6]. To find on  $\Omega_r^T \to \mathcal{C}^1$ -solution of the Cauchy problem (2.1.1)-(2.1.2), it suffices to look at what happens at the initial time t = 0. A necessary and sufficient condition is to impose

(2.1.6) 
$$(D_x h^{\varepsilon}(x))^3 = 0, \quad \forall (x, \varepsilon) \in \Omega^0_r \times ]0, 1].$$

Then (see [6]), the solution of (2.1.1)-(2.1.2) satisfies div  $u^{\varepsilon} = 0$  in  $\Omega_r^T$ . It means that the solutions of (2.1.1) under study are also (local) solutions of the incompressible Euler equations (with constant pressure) :

(2.1.7) 
$$\partial_t u^{\varepsilon} + (u^{\varepsilon} \cdot \nabla) u^{\varepsilon} + \nabla p_{\varepsilon} = 0, \quad \text{div} u^{\varepsilon} = 0, \quad p_{\varepsilon} = c.$$

In what follows, we work with the conditions (2.1.1), (2.1.2) and (2.1.6). We seek simple wave solutions meaning that we want to solve directly (2.1.1)-(2.1.6) through a construction relying on the special form (2.1.2). Taking into account (2.1.7), this can be viewed as a preliminary step towards a more general (large amplitude) WKB calculus concerning incompressible or compressible Euler equations. The long-term perspective is indeed to incorporate at the level of (2.1.1) the influence of extra terms (like pressure, viscosity,  $\cdots$ ) and the presence of complete expansions for the profile such as

(2.1.8) 
$$w_{\varepsilon}(x,\theta) = w(x,\theta) + \sum_{j=1}^{\infty} \varepsilon^{\kappa j} w^{j}(x,\theta), \qquad \kappa \in ]0,1] \cap \mathbb{Q}.$$

Let us recall here what has yet been obtained concerning (2.1.7) when the initial data are adjusted as in (2.1.2) and (2.1.8). The case  $\kappa = 1$  with a profile  $w(x, \theta) = w(x)$  independent of the fast variable  $\theta \in \mathbb{T}$  is well-known. It is a variant of standard results in weakly nonlinear geometric optics [17, 23]. The case  $\kappa \in [0, 1] \cap \mathbb{Q}$  with still  $w(x, \theta) = w(x)$  is fully discussed in [9]. The case  $\kappa = 1$  associated with (2.1.3) corresponds to a more singular situation. It is much more delicate. It is what here holds our attention.

In the case  $\kappa = 1$  together with (2.1.3), the WKB analysis of incompressible Euler equations is supposed to be not well-posed [24]. This is due to a strong coupling

between the profile  $w(x, \theta)$  and the phase  $\varphi(x)$ . In such a regime, many unstable phenomena (see for instance [14, 21]) can occur. Therefore, any progress in this direction requires to work in a very specific context, like here (2.1.1)-(2.1.2)-(2.1.6), with adapted tools.

The study of (2.1.2)-(2.1.6) is not so easy to achieve. In [5, 6, 10], some very special examples are proposed implying functions  $h^{\varepsilon}(x)$  which are adjusted such that the matrix  $D_x h^{\varepsilon}(x)$  is of rank 1. These preliminary advancements are partially completed in [7] by exploring (without restriction on the space dimension  $d \in \mathbb{N}^*_+$ ) some *necessary* condition on  $\varphi(x)$  and  $w(x,\theta)$  giving rise to matrices  $D_x h^{\varepsilon}(x)$  which are nilpotent, as in (2.1.6).

# 2.1.2 The main results.

In this paper, we restrict our attention to the case d = 3 but, this time, we seek *necessary and sufficient* conditions on  $(\varphi, w)$  to have (2.1.2)-(2.1.6). This approach leads to the notion of *compatible couple* given below.

**Définition 2.1.1.** Let  $\varphi \in C^1(\Omega^0_r; \mathbb{R})$  and  $w \in C^1(\Omega^0_r \times \mathbb{T}; \mathbb{R}^3)$  two functions satisfying the preliminary assumptions

(2.1.9)  $\partial_{\theta} w(x,\theta) \neq 0, \quad \nabla \varphi(x) \neq 0.$ 

The couple  $(\varphi, w)$  is said to be <u>compatible</u> on  $\Omega^0_r \times \mathbb{T}$  if the family  $\{h^{\varepsilon}\}_{\varepsilon}$  which is associated to  $(\varphi, w)$  through (2.1.2) satisfies (2.1.6).

It is possible to derive an exhaustive description of all compatible couples. In the statement below, for the sake of brevity, we express this remarkable fact in a rather imprecise form.

**Théorème 4.** There is a whole class of compatible couples  $(\varphi, w)$ .

The interesting aspects will appear in the text when precising the structure of the functions  $\varphi$  and w such involved, and especially when describing the geometrical features of  $\varphi$  and how to get them.

Retain here that we can perform a complete WKB analysis of the constraints (2.1.2), (2.1.6) and (2.1.9). Then, applying Theorem 2.6 of [6], we are sure to recover by this way the existence of *large amplitude high-frequency waves*  $u^{\varepsilon}(t,x)$  which are special solutions of (2.1.7) on  $\Omega_r^T$ . Now, the structure of the expressions  $u^{\varepsilon}(t,x)$  can be precised as follows.

**Théorème 5.** Let  $(\varphi, w)$  be a couple which is compatible on  $\Omega_r^0 \times \mathbb{T}$ . There are functions  $\mathbf{W}(\varphi, \psi, \theta) \in \mathcal{C}^1(\mathbb{R}^2 \times \mathbb{T}; \mathbb{R})$  and  $\psi(x, \theta) \in \mathcal{C}^1(\Omega_r^0 \times \mathbb{T}; \mathbb{R})$  such that the profile  $w(x, \theta)$  can be factorized through

(2.1.10) 
$$w(x,\theta) = \mathbf{W}(\varphi(x),\psi(x,\theta),\theta), \quad \nabla\varphi \wedge \nabla\psi \neq 0.$$

#### CHAPTER 2. OSCILLATING SOLUTIONS

There is also some T > 0 such that the Cauchy problem

(2.1.11) 
$$\begin{cases} \partial_t \Phi + (\mathbf{W}(\Phi, \Psi, \theta) \cdot \nabla) \Phi = 0, & \Phi(0, x) = \varphi(x), \\ \partial_t \Psi + (\mathbf{W}(\Phi, \Psi, \theta) \cdot \nabla) \Psi = 0, & \Psi(0, x, \theta) = \psi(x, \theta), \end{cases}$$

has a solution  $(\Phi, \Psi)(t, x, \theta)$  on the domain  $\Omega_r^T \times \mathbb{T}$ . We have  $\partial_{\theta} \Phi \equiv 0$  and, for all  $\varepsilon \in [0, 1]$ , the oscillation

(2.1.12) 
$$u^{\varepsilon}(t,x) = \mathbf{W}\big(\Phi(t,x), \Psi(t,x,\Phi(t,x)/\varepsilon), \Phi(t,x)/\varepsilon\big), \quad \varepsilon \in ]0,1]$$

is a solution of (2.1.1) on the domain  $\Omega_r^T$  with initial data  $u^{\varepsilon}(0, \cdot)$  as in (2.1.2). Moreover, for all  $t \in [0, T]$  the couple  $(\Phi(t, \cdot), \widetilde{\mathbf{W}}(t, \cdot))$  where

$$\mathbf{W}(t, x, \theta) := \mathbf{W}(\Phi(t, x), \Psi(t, x, \theta), \theta)$$

is still compatible on  $B(0, r - tV[\times \mathbb{T}])$ . More precisely, for all  $t \in [0, T]$ , we must have

(2.1.13)  $\nabla \Phi \cdot \partial_{\theta} \mathbf{W} + \partial_{\theta} \Psi \ \nabla \Phi \cdot \partial_{\Psi} \mathbf{W} \equiv 0,$ 

$$(2.1.14) \qquad (\nabla \Phi \cdot \partial_{\Psi} \mathbf{W}) (\nabla \Psi \cdot \partial_{\theta} \mathbf{W} + \partial_{\theta} \Psi \nabla \Psi \cdot \partial_{\Psi} \mathbf{W}) \equiv 0,$$

(2.1.15) 
$$(\nabla \Phi \cdot \partial_{\varphi} \mathbf{W})^2 + (\nabla \Phi \cdot \partial_{\Psi} \mathbf{W}) (\nabla \Psi \cdot \partial_{\varphi} \mathbf{W}) \equiv 0,$$

(2.1.16)  $\nabla \Phi \cdot \partial_{\varphi} \mathbf{W} + \nabla \Psi \cdot \partial_{\Psi} \mathbf{W} \equiv 0.$ 

In comparison with preceding works [5, 6, 7, 10, 14], this second result 5 includes various situations which have not yet been studied. It allows to exhibit many new phenomena with respect to both the propagation and the interaction of oscillations.

# 2.1.3 Plan of the article.

We present here the plan of the present article. We take this opportunity to make some clarifications and to indicate ideas of proof.

• In Chapter 2.2, we discuss the notion of *compatible couple*. More precisely, the Proposition 2.2.1 of Section 2.2.1 says that any compatible couple  $(\varphi, w)$  must verify a list S, namely (2.2.1)-(2.2.2)-(2.2.3)-(2.2.4), of conditions which are independent of the parameter  $\varepsilon \in [0, 1]$ .

Then, in the Proposition 2.2.2 which is proved in Section 2.2.2, we observe that there exists a scalar function  $\psi \in C^1(\Omega_r^0 \times \mathbb{T}; \mathbb{R})$  leading to a factorization of the involved profiles  $w(x, \theta)$  in the form (2.1.10). It follows simplifications when dealing with the system S. It remains (see the Proposition 2.2.3 proved in Section 2.2.3) some necessary and sufficient conditions to impose on the three ingredients  $\varphi$ ,  $\psi$ and  $\mathbf{W}$ . In fact, the matter is to work with the relations (2.1.13), (2.1.14), (2.1.15) and (2.1.16) at the time t = 0.

• Chapter 2.3 consider the simplest case, when  $\nabla \varphi \cdot \partial_{\psi} \mathbf{W} \equiv 0$ . Then, as it is explained in Section 2.3.1, the level surfaces of the phase  $\varphi$  can be associated with

some foliated structure of  $\mathbb{R}^3$  by planes. This information is a crucial key which, in Section 2.3.2, enables progress leading to a complete description of  $(\varphi, \psi, \mathbf{W})$ , and therefore  $(\varphi, w)$ .

• Chapter 2.4 is devoted to the case  $\nabla \varphi \cdot \partial_{\psi} \mathbf{W} \neq 0$ . Then, without loss of generality, the profile  $w(x, \theta)$  can be assumed to be of the form

$$w(x,\theta) = {}^{t} (v,\psi, \mathfrak{L}(\psi,v))(x,\theta), \qquad v(x,\theta) = \mathbf{V}(\varphi(x),\psi(x,\theta),\theta)$$

where  $\mathfrak{L}(\psi, v)$  and  $\mathbf{V}(\varphi, \psi, \theta)$  are auxiliary functions. On the other hand, the expression  $\psi(x, \theta)$  can always be factorized according to

$$\psi(x, \theta) = u(x, v(x, \theta)), \qquad \partial_v u(x, v) \not\equiv 0.$$

In Section 2.4.1, see the Proposition 2.4.1, the information  $\nabla \varphi \cdot \partial_{\psi} \mathbf{W} \neq 0$  is exploited in order to rephrase the conditions (2.1.13), (2.1.14), (2.1.15) and (2.1.16), written at the time t = 0 on  $\varphi \equiv \Phi(0, \cdot)$ ,  $\psi \equiv \Psi(0, \cdot)$  and  $\mathbf{W}$ , in terms of the more convenient conditions (2.4.13), (2.4.14) and (2.4.15) which concern only  $\varphi$  and  $\psi$  (as well as  $\mathfrak{L}$  and  $\mathbf{V}$ ).

After eliminating the special case  $\partial_3 u \equiv 0$ , we concentrate on the remaining situation  $\partial_3 u \neq 0$ . At this stage, the question becomes the following (see also the remark 2.4.3.1 for a functional analysis viewpoint).

The intermediate problem under study. The question is to find smooth <u>non</u> <u>constant</u> functions  $\Phi(x_1, x_2, u, v)$ , locally defined in  $\mathbb{R}^4$ , satisfying the two transport equations

(2.1.17) 
$$X \Phi \equiv 0, \quad Y \Phi \equiv 0, \quad X := \partial_1 + R \partial_2, \quad Y := R \partial_u + \partial_v$$

and involving a variable coefficient  $R(x_1, x_2, u, v)$  which can be identified through the implicit relation

$$(2.1.18) \qquad \qquad \partial_v u(x,v) = R(x_1, x_2, u(x,v), v)$$

where the function u(x, v) must satisfy the two conservation laws

(2.1.19) 
$$\partial_1 u + \partial_v \mathfrak{L}(u, v) \ \partial_3 u = 0, \qquad \partial_2 u + \partial_u \mathfrak{L}(u, v) \ \partial_3 u = 0.$$

At the level of (2.1.19), the variable v plays the part of a parameter. When solving (2.1.19), there are degrees of freedom related to the choices of  $u(0, 0, x_3)$ and  $\mathfrak{L}(u, v)$ . Once the function u (and therefore R) is known, the difficulty is to find solutions  $\Phi$  of (2.1.17) satisfying  $\nabla \Phi \neq 0$ . Let us say a few words about the origin of the conditions (2.1.17) and  $\nabla \Phi \neq 0$ . In fact, the expression  $\Phi$  is issued from  $\varphi$  after a blowing-up procedure. Indeed, one has

$$\varphi(x) = \Phi(x_1, x_2, u(x, v), v), \qquad R = \partial_v u \neq 0, \qquad v(x, \theta).$$

In this context, the condition  $Y \Phi \equiv 0$  means simply that  $\varphi$  does not depend on v. Since the letter v is aimed to be replaced by a function  $v(x,\theta)$  of the variables  $(x,\theta) \in \mathbb{R}^3 \times \mathbb{T}$ , this is equivalent to say that  $\partial_{\theta} \varphi \equiv 0$ . This is a natural requirement. Despite the strength of the nonlinearity, we do not want that the phase  $\varphi$  starts to oscillate with respect to itself. The other restrictions  $X \Phi \equiv 0$ , (2.1.18) and (2.1.19) are coming from (2.1.6) after the reduction procedure.

Recall that the phase  $\varphi$  is supposed to be not stationary, see (2.1.4). This is possible if and only if the Poisson algebra  $\mathcal{A}$  generated by the two vector fields X and Y is of dimension strictly less than four (dim  $\mathcal{A} < 4$ ). The corresponding integrability criterion (of Frobenius type) can be traducted in terms of conditions on R. Actually, the Proposition 2.4.2 in Section 2.4.2 exhibits the relevant nonlinear PDE's to impose on R. In the case dim  $\mathcal{A} = 2$ , we find (2.4.38). When dim  $\mathcal{A} = 2$ , we have to deal with (2.4.39)-(2.4.40).

Note that the construction of phases  $\varphi$  (through  $\Phi$ ) is associated with the production of special foliations of  $\mathbb{R}^4$ . The related subtle informations would be out of reach when working with functions  $\varphi$  depending only on  $x \in \mathbb{R}^3$ . Now, the difficulty is that the coefficient R must also be issued from (2.1.18) after solving the two conservation laws given line (2.1.19). It follows that the expression R inherits some special structure described at the level of Proposition 2.4.3 in Section 2.4.3. Given smooth functions  $\mathfrak{K}$  and  $\mathfrak{L}$ , introduce

(2.1.20) 
$$\alpha(x_1, x_2, u, v) := \mathfrak{K}(u, v) + \partial_v \mathfrak{L}(u, v) x_1 + \partial_u \mathfrak{L}(u, v) x_2.$$

The function R must be in the form

(2.1.21) 
$$R(x_1, x_2, u, v) = -\partial_v \alpha(x_1, x_2, u, v) / \partial_u \alpha(x_1, x_2, u, v).$$

In Section 2.4.4, we test the integrability conditions (2.4.39) and (2.4.40) in the framework of (2.1.20) and (2.1.21). Surprisingly, all requirements are met for many choices of the functions  $\mathfrak{K}$  and  $\mathfrak{L}$  leading in Section 2.4.5 to a complete classification of all compatible couples ( $\varphi, w$ ).

To our knowledge, the preceding approach and the corresponding analysis is completely original and new. In the end, it furnishes a good description of the *class* of functions  $\varphi(x)$  and  $w(x, \theta)$  mentioned in the Theorem 4.

We conclude the chapter 2.4 by producing in the paragraph 2.4.6 illustrative examples of compatible couples  $(\varphi, w)$ .

• In Section 2.5, we study the time evolution problem. We show Theorem 5. This result is proved in the paragraph 2.5.1. It furnishes, in the context of the equation (2.1.1), a complete description of what can happen in terms of smooth large amplitude oscillations. The formula (2.1.12) generalizes previous examples exhibited in [7, 10, 13, 24].

The families  $\{u^{\varepsilon}\}_{\varepsilon\in]0,1]}$  exhibited in (2.1.12) belong to a regime which, in non linear geometric optics, is called supercritical (because one order derivatives of  $u^{\varepsilon}$ 

explode when  $\varepsilon$  goes to 0). Expressions like  $u^{\varepsilon}$  are very unstable objects [10] unless some small viscosity is added [4]. Their asymptotic behaviours (always as  $\varepsilon \to 0$ ) can involve interesting features.

For instance, in Section 2.5.2, we can exhibit a phenomenon of *superposition* of oscillations. It is obtained by selecting compatible couples  $(\varphi, w_{\varepsilon})$  where, in contrast to (2.1.2), the profiles  $w_{\varepsilon}$  depend on  $\varepsilon \in [0, 1]$ . More precisely, the expression  $w_{\varepsilon}$  is built with functions **W** and  $\psi$  through the formula

$$w_{\varepsilon}(x,\theta) = \mathbf{W}(\varphi(x),\psi(x)/\varepsilon,\theta), \qquad \mathbf{W}(\varphi,\cdot,\theta) \in \mathcal{C}^{\infty}(\mathbb{T};\mathbb{R}^3).$$

At the time t = 0, we are faced with a large amplitude multiphase oscillation

(2.1.22) 
$$u^{\varepsilon}(0,x) = \mathbf{W}(\varphi(x),\psi(x)/\varepsilon,\varphi(x)/\varepsilon), \qquad \nabla\varphi \wedge \nabla\psi \neq 0.$$

On the other hand, at any time  $t \in [0, T]$ , the function  $\Psi(t, \cdot)$  starts to really depend on  $\theta \in \mathbb{T}$  giving rise to

(2.1.23) 
$$u^{\varepsilon}(t,x) = \mathbf{W}\left(\Phi(t,x), \frac{\Psi(t,x,\Phi(t,x)/\varepsilon)}{\varepsilon}, \frac{\Phi(t,x)}{\varepsilon}\right)$$

Thus, the interaction of large amplitude waves oscillating in transversal directions at the frequency  $\varepsilon^{-1}$  can produce oscillations with frequency  $\varepsilon^{-2}$ . Such a *turbulent* effect was already mentioned in [5] in the context of the system (2.1.1) when d = 2. On the contrary, we were not able to prove the same effect in the case of (2.1.7) when d = 2. It seems that, when adding the divergence free condition, it is a specificity of the space dimension d = 3.

• The aim of Appendix 6 is to check that the list of situations enumerated at the level of Proposition 2.4.4 is exhaustive. The corresponding work of verification is quite long and technical. The difficulties are due to the fact that it is delicate to interpret the integrability conditions to impose on R into convenient constraints on the functions  $\mathfrak{K}$  and  $\mathfrak{L}$  appearing at the level of (2.1.20). This will be done step by step, from paragraph 2.6.1 up to 2.6.5.

# 2.2 Compatible couples.

From now on, we write  $f \equiv 0$  and  $f \not\equiv 0$  to mean respectively that f is identically zero on its domain of definition or that it is a non-zero function.

# 2.2.1 The notion of compatible couples.

Given two vectors  $u = {}^t(u_1, u_2, u_3) \in \mathbb{R}^3$  and  $v = {}^t(v_1, v_2, v_3) \in \mathbb{R}^3$ , we note

 $u \cdot v := u_1 v_1 + u_2 v_2 + u_3 v_3 \,,$ 

$$u \otimes v := \begin{pmatrix} u_1v_1 & u_1v_2 & u_1v_3 \\ u_2v_1 & u_2v_2 & u_2v_3 \\ u_3v_1 & u_3v_2 & u_3v_3 \end{pmatrix}, \qquad u \wedge v := \begin{pmatrix} u_2v_3 - u_3v_2 \\ u_3v_1 - u_1v_3 \\ u_1v_2 - u_2v_1 \end{pmatrix}.$$

We can interpret (2.1.6) in the form of conditions on  $(w, \varphi)$ .

**Proposition 2.2.1.** Let  $\varphi \in C^1(\Omega^0_r; \mathbb{R})$  and  $w \in C^1(\Omega^0_r \times \mathbb{T}; \mathbb{R}^3)$  satisfying the preliminary assumptions (2.1.9). The couple  $(\varphi, w)$  is compatible on  $\Omega^0_r \times \mathbb{T}$  if and ony if it is a solution on  $\Omega^0_r \times \mathbb{T}$  of the system S made of

(2.2.1) 
$$\nabla \varphi \cdot \partial_{\theta} w \equiv 0,$$

(2.2.2) 
$$\nabla \varphi \cdot (D_x w \ \partial_\theta w) \equiv 0,$$

$$(2.2.3) (D_x w)^3 \equiv 0,$$

(2.2.4) 
$$M (D_x w)^2 + D_x w M D_x w + (D_x w)^2 M \equiv 0, \quad M := \partial_\theta w \otimes \nabla \varphi$$

Proof of Proposition 2.2.1. We find

$$D_x h^{\varepsilon}(x) = (D_x w) \left( x, \frac{\varphi(x)}{\varepsilon} \right) + \frac{1}{\varepsilon} \, \partial_{\theta} w \left( x, \frac{\varphi(x)}{\varepsilon} \right) \otimes \nabla \varphi(x) \,.$$

The constraint (2.1.6) can also be formulated as

$$\sum_{j=0}^{3} \varepsilon^{-j} \Xi_j \left( x, \frac{\varphi(x)}{\varepsilon} \right) \equiv 0, \qquad \Xi_j(x, \theta) \in C^0 \left( \Omega_r^0 \times \mathbb{T}, \mathcal{M}_3(\mathbb{R}^3) \right)$$

where

$$\begin{aligned} \Xi_0 &= (D_x w)^3 , \qquad \Xi_1 &= (D_x w)^2 M + D_x w M D_x w + M (D_x w)^2 , \\ \Xi_3 &= M^3 , \qquad \Xi_2 &= M^2 D_x w + D_x w M^2 + M D_x w M . \end{aligned}$$

To guarantee (2.1.6) for all  $\varepsilon \in [0, 1]$ , it is necessary and sufficient to impose

(2.2.5) 
$$\Xi_j \equiv 0, \qquad \forall (x,\theta) \in \Omega_r^0 \times \mathbb{T}, \qquad \forall j \in \{0,1,2,3\}.$$

Our aim is to solve (2.2.5) for some  $r \in \mathbb{R}^*_+$ . The constraints  $\Xi_0 \equiv 0$  and  $\Xi_1 \equiv 0$  are repetitions of respectively (2.2.3) and (2.2.4). Since

$$M^3 = (
abla arphi \cdot \partial_{ heta} w)^2 \; \partial_{ heta} w \otimes 
abla arphi = 0 \,, \qquad \partial_{ heta} w \otimes 
abla arphi \not\equiv 0 \,,$$

the examination of  $\Xi_3$  leads to (2.2.1). In view of (2.2.1), we have also  $M^2 \equiv 0$ . Thus, the condition  $\Xi_2 \equiv 0$  reduces to  $M D_x w M \equiv 0$ , that is (2.2.2).

The system S, as presented above, is not yet exploitable. The purpose of this chapter 2.2 is to put it in a suitable form. In view of (2.2.3), the rank of the matrix  $D_x w$  is either one or two (the zero case being trivial). The next paragraphs 2.2.1.1 and 2.2.1.2 deal separately with these two situations.

#### 2.2.1.1 The case of rank one.

In this paragraph, we suppose that

(2.2.6) 
$$rg(D_xw(x,\theta)) = dim(Im(D_xw)(x,\theta)) = 1, \qquad \forall (x,\theta) \in \Omega_r^0 \times \mathbb{T}.$$

By the constant rank theorem [2] and due to the compacity of the torus  $\mathbb{T}$ , by restricting  $r \in \mathbb{R}^*_+$  if necessary, we can find two functions  $\psi \in \mathcal{C}^1(\Omega^0_r \times \mathbb{T}; \mathbb{R})$  and  $\mathbf{W} \in \mathcal{C}^1(\mathbb{R} \times \mathbb{T}; \mathbb{R}^3)$  with  $\nabla \psi \neq 0$  and  $\partial_{\psi} \mathbf{W} \neq 0$  such that

(2.2.7) 
$$w(x,\theta) = \mathbf{W}(\psi(x,\theta),\theta), \quad \forall (x,\theta) \in \Omega_r^0 \times \mathbb{T}.$$

**Lemme 2.2.1.** Assume (2.1.9) and (2.2.7). Then, the couple  $(\varphi, w)$  is compatible on the domain  $\Omega_r^0 \times \mathbb{T}$  if and only if the following conditions are verified :

(2.2.8) 
$$\nabla \varphi \cdot \partial_{\theta} w \equiv 0$$

(2.2.9) 
$$(\nabla \varphi \cdot \partial_{\psi} \mathbf{W}) (\nabla \psi \cdot \partial_{\theta} w) \equiv 0$$

(2.2.10) 
$$\nabla \psi \cdot \partial_{\psi} \mathbf{W} \equiv 0.$$

**Proof of Lemma 2.2.1.** The condition (2.2.8) is the same as (2.2.1). Taking into account (2.2.7), we find  $D_x w = \partial_{\psi} \mathbf{W} \otimes \nabla \psi$  so that (2.2.3) becomes

 $(\nabla \psi \cdot \partial_{\psi} \mathbf{W})^2 \ \partial_{\psi} \mathbf{W} \otimes \nabla \psi \equiv 0 \,, \qquad \partial_{\psi} \mathbf{W} \otimes \nabla \psi \neq 0$ 

which implies (2.2.10). Knowing (2.2.10), the constraint (2.2.4) reduces to

$$D_x w M D_x w = (\nabla \varphi \cdot \partial_{\psi} \mathbf{W}) (\nabla \psi \cdot \partial_{\theta} w) \partial_{\psi} \mathbf{W} \otimes \nabla \psi \equiv 0.$$

We recover here (2.2.9) which also allows to guarantee (2.2.2).

In this paragraph, we suppose that

(2.2.11) 
$$rg(D_xw(x,\theta)) = dim(Im(D_xw)(x,\theta)) = 2, \qquad \forall (x,\theta) \in \Omega_r^0 \times \mathbb{T}.$$

As before, we can apply the constant rank theorem [2] in order to find three functions  $\psi \in \mathcal{C}^1(\Omega^0_r \times \mathbb{T}; \mathbb{R}), \ \tilde{\psi} \in \mathcal{C}^1(\Omega^0_r \times \mathbb{T}; \mathbb{R})$  and  $\mathbf{W} \in \mathcal{C}^1(\mathbb{R} \times \mathbb{T}; \mathbb{R}^3)$  with  $\nabla \psi \neq 0, \ \nabla \tilde{\psi} \neq 0, \ \nabla \psi \wedge \nabla \tilde{\psi} \neq 0$  and  $\partial_{\psi} \mathbf{W} \neq 0$  such that

(2.2.12) 
$$w(x,\theta) = \mathbf{W}\big(\tilde{\psi}(x,\theta),\psi(x,\theta),\theta\big), \qquad \forall (x,\theta) \in \Omega^0_r \times \mathbb{T}.$$

In the Section 2.2.2, we will show that we can take  $\tilde{\psi} \equiv \varphi$ . The precise statement is the following.

**Proposition 2.2.2.** Let  $(\varphi, w)$  be a compatible couple on the domain  $\Omega_r^0 \times \mathbb{T}$ . By restricting  $r \in \mathbb{R}^*_+$  if necessary, we can find a function  $\psi \in \mathcal{C}^1(\Omega_r^0 \times \mathbb{T}; \mathbb{R})$  satisfying  $\nabla \varphi \wedge \nabla \psi \neq 0$  and a vector function  $\mathbf{W} \in \mathcal{C}^1(\mathbb{R}^2 \times \mathbb{T}; \mathbb{R}^3)$  such as

(2.2.13) 
$$w(x,\theta) = \mathbf{W}(\varphi(x),\psi(x,\theta),\theta), \quad \forall (x,\theta) \in \Omega^0_r \times \mathbb{T}.$$

Assuming (2.2.13), we can compute

(2.2.14) 
$$D_x w(x,\theta) = \partial_{\varphi} \mathbf{W} \otimes \nabla \varphi + \partial_{\psi} \mathbf{W} \otimes \nabla \psi.$$

In view of (2.2.11), the two vectors  $\nabla \varphi$  and  $\nabla \psi$ , as well as  $\partial_{\varphi} \mathbf{W}$  and  $\partial_{\psi} \mathbf{W}$ , must be independent. In other words :

(2.2.15) 
$$\nabla \varphi \wedge \nabla \psi \neq 0, \qquad \partial_{\varphi} \mathbf{W} \wedge \partial_{\psi} \mathbf{W} \neq 0.$$

On the other hand, the condition (2.1.9) amounts to the same thing as

(2.2.16) 
$$\partial_{\theta}\psi \ \partial_{\psi}\mathbf{W} + \partial_{\theta}\mathbf{W} \neq 0.$$

In the Section 2.2.3, we will further exploit the information (2.2.13) in order to interpret the system S differently. Just retain here that :

**Proposition 2.2.3.** Assume (2.2.11) and (2.2.13) together with the preliminary hypothesis (2.2.16). Then, the couple  $(\varphi, w)$  is compatible on the domain  $\Omega_r^0 \times \mathbb{T}$  if and only if we have (2.2.15) and the following conditions :

(2.2.17)  $\nabla \varphi \cdot \partial_{\theta} \mathbf{W} + \partial_{\theta} \psi \, \nabla \varphi \cdot \partial_{\psi} \mathbf{W} \equiv 0,$ 

(2.2.18) 
$$(\nabla \varphi \cdot \partial_{\psi} \mathbf{W}) (\nabla \psi \cdot \partial_{\theta} \mathbf{W} + \partial_{\theta} \psi \nabla \psi \cdot \partial_{\psi} \mathbf{W}) \equiv 0,$$

(2.2.19) 
$$(\nabla \varphi \cdot \partial_{\varphi} \mathbf{W})^2 + (\nabla \varphi \cdot \partial_{\psi} \mathbf{W}) (\nabla \psi \cdot \partial_{\varphi} \mathbf{W}) \equiv 0,$$

(2.2.20) 
$$\nabla \varphi \cdot \partial_{\varphi} \mathbf{W} + \nabla \psi \cdot \partial_{\psi} \mathbf{W} \equiv 0$$

Comparing the two Propositions 2.2.1 and 2.2.3, we see that (2.2.8)-(2.2.9)-(2.2.10) can be handled as a special case of (2.2.17)-(2.2.18)-(2.2.19)-(2.2.20). It suffices to work with  $\partial_{\varphi} \mathbf{W} \equiv 0$ . Thus, in the chapters 2.3 and 2.4, we can concentrate on the system (2.2.17)-(2.2.18)-(2.2.19)-(2.2.20). We will examine separately what happens when respectively  $\nabla \varphi \cdot \partial_{\psi} \mathbf{W} \equiv 0$  and  $\nabla \varphi \cdot \partial_{\psi} \mathbf{W} \not\equiv 0$ .

# 2.2.2 Factorization of compatible couples.

Suppose (2.2.11). To obtain (2.2.13), we proceed in two steps. First, in the paragraph 2.2.2.1, we produce a local version of the Proposition 2.2.2. Then, in the paragraph 2.2.2.2, we complete the proof of the Proposition 2.2.2.

#### 2.2.2.1 The local version of the Proposition 2.2.2.

Note  $\vec{0} = (0,0,0) \in \Omega_r^0 \subset \mathbb{R}^3$ . In this paragraph, we work locally, near a point  $(\vec{0}, \tilde{\theta}) \in \Omega_r^0 \times \mathbb{T}$ . We select some open connected neighbourhood  $\Gamma$  satisfying  $(\vec{0}, \tilde{\theta}) \in \Gamma \subset \Omega_r^0 \times \mathbb{T}$ . Typically, we can take

$$\Gamma \equiv \Gamma^{\theta}_{r,\tilde{r}} := \Omega^{0}_{r} \times \left] \theta - \tilde{r} , \theta + \tilde{r} \right[, \qquad (r,\tilde{r},\theta) \in \mathbb{R}^{*}_{+} \times \left] 0, 1 \right[ \times \mathbb{T}]$$

Let  $(\varphi, w)$  be a couple which is compatible on  $\Gamma_{r,\tilde{r}}^{\tilde{\theta}}$ . By exchanging  $w(x,\theta)$  into  $w(x,\theta-\tilde{\theta})$ , we can always suppose that  $\tilde{\theta}=0$ . In what follows, we will argue on

 $\Gamma^0_{r,\tilde{r}}$ . Note *i*, *j* and *k* three distinct elements chosen among the set  $\{1, 2, 3\}$ . The constraint (2.2.11) means that there is *k* giving rise to

(2.2.21) 
$$\nabla w_k(x,\theta) \in Vec \langle \nabla w_i(x,\theta), \nabla w_j(x,\theta) \rangle$$
,  $\forall (x,\theta) \in \Gamma^0_{r,\tilde{r}}$ ,

(2.2.22) 
$$\nabla w_i(x,\theta) \wedge \nabla w_j(x,\theta) \neq 0 \quad , \qquad \forall \ (x,\theta) \in \Gamma^0_{r,\tilde{r}} \, .$$

The direction  $\nabla \varphi(\vec{0})$  cannot be simultaneously colinear to the two vectors  $\nabla w_i(\vec{0}, 0)$ and  $\nabla w_j(\vec{0}, 0)$ . Pick the indice  $l \in \{i, j\}$  in such a way that  $\nabla \varphi(\vec{0}) \wedge \nabla w_l(\vec{0}, 0) \neq 0$ . Then, do a permutation on the three directions  $x_1, x_2$  and  $x_3$  (with the corresponding permutation on the components  $w_1, w_2$  and  $w_3$ ) in order to have l = 1and k = 3. Then, by restricting  $r \in \mathbb{R}^*_+$  and  $\tilde{r} \in ]0, 1[$ , we can obtain

(2.2.23) 
$$\nabla \varphi \wedge \nabla w_1 \neq 0, \quad \forall (x,\theta) \in \Gamma^0_{r,\tilde{r}}$$

while the conditions (2.2.21) and (2.2.22) become

$$(2.2.24) \qquad \nabla w_3(x,\theta) \in Vec \left\langle \nabla w_1(x,\theta), \nabla w_2(x,\theta) \right\rangle \quad , \qquad \forall \ (x,\theta) \in \Gamma^0_{r,\tilde{r}},$$

(2.2.25) 
$$\nabla w_1(x,\theta) \wedge \nabla w_2(x,\theta) \neq 0 \quad , \qquad \forall \ (x,\theta) \in \Gamma^0_{r,\tilde{r}}$$

The constraint (2.2.24) allows to deduce the existence of a scalar function  $\mathbb{W}_3$  in  $C^1(\mathbb{R}^2 \times ] - \tilde{r}, \tilde{r}[;\mathbb{R})$  such that

(2.2.26) 
$$w_3(x,\theta) = \mathbb{W}_3(w_1(x,\theta), w_2(x,\theta), \theta), \quad \forall (x,\theta) \in \Gamma^0_{r,\tilde{r}}.$$

Then, using the convention

$$\mathbb{W}(w_1, w_2, \theta) = \begin{pmatrix} \mathbb{W}_1(w_1, w_2, \theta) \\ \mathbb{W}_2(w_1, w_2, \theta) \\ \mathbb{W}_3(w_1, w_2, \theta) \end{pmatrix} := \begin{pmatrix} w_1 \\ w_2 \\ \mathbb{W}_3(w_1, w_2, \theta) \end{pmatrix},$$
we can get

we can get

(2.2.27) 
$$w(x,\theta) = \mathbb{W}(w_1(x,\theta), w_2(x,\theta), \theta), \quad \forall \ (x,\theta) \in \Gamma^0_{r,\tilde{r}}.$$

**Lemme 2.2.2.** Select a couple  $(\varphi, w)$  which is compatible on  $\Gamma^0_{r,\tilde{r}}$  and which satisfies (2.2.24) together with (2.2.25). Then, there exists a scalar function  $W_2 \in C^1(\mathbb{R}^2 \times ] - \tilde{r}, \tilde{r}[;\mathbb{R})$  such that the component  $w_2$  can be put in the form

(2.2.28) 
$$w_2(x,\theta) = W_2(\varphi(x), w_1(x,\theta), \theta), \quad \forall \ (x,\theta) \in \Gamma^0_{r,\tilde{r}}.$$

Proof of Lemma 2.2.2. To obtain (2.2.28), it suffices to show that

(2.2.29) 
$$\nabla w_2(x,\theta) \in Vec \left\langle \nabla \varphi(x), \nabla w_1(x,\theta) \right\rangle, \quad \forall \ (x,\theta) \in \Gamma^0_{r,\tilde{r}}.$$

The proof is by contradiction. Suppose that (2.2.29) is not verified :

(2.2.30) 
$$\exists (x_0, \theta_0) \in \Gamma^0_{r, \tilde{r}}, \qquad \nabla w_2(x_0, \theta_0) \notin Vec \left\langle \nabla \varphi(x_0), \nabla w_1(x_0, \theta_0) \right\rangle.$$

Combining (2.2.23) and (2.2.30), we see that the vectors  $\nabla \varphi(x_0)$ ,  $\nabla w_1(x_0, \theta_0)$  and  $\nabla w_2(x_0, \theta_0)$  give rise to a basis of  $\mathbb{R}^3$ . In addition, by using the definition of the  $\Xi_j$  and the restrictions (2.2.1), (2.2.2), (2.2.3) and (2.2.4), we can get

$$(D_x w + \partial_\theta w \otimes \nabla \varphi)^3 = \sum_{j=0}^3 \Xi_j = 0.$$

Thus, the matrice

$$D_x w + \partial_\theta w \otimes \nabla \varphi = \begin{pmatrix} {}^t \nabla w_1 + \partial_\theta w_1 {}^t \nabla \varphi \\ {}^t \nabla w_2 + \partial_\theta w_2 {}^t \nabla \varphi \\ {}^t \nabla w_3 + \partial_\theta w_3 {}^t \nabla \varphi \end{pmatrix}$$

is at most of rank two. In view of (2.2.26), the third row vector is

$$\nabla w_3 + \partial_\theta w_3 \nabla \varphi = \partial_{w_1} \mathbb{W}_3 \left( \nabla w_1 + \partial_\theta w_1 \nabla \varphi \right) + \partial_{w_2} \mathbb{W}_3 \left( \nabla w_2 + \partial_\theta w_2 \nabla \varphi \right) + \partial_\theta \mathbb{W}_3 \nabla \varphi .$$

It must be a linear combination of the two first row vectors so that

(2.2.31) 
$$(\partial_{\theta} \mathbb{W}_3)(w_1(x_0, \theta_0), w_2(x_0, \theta_0), \theta_0) = 0$$

In what follows, the functions will be (unless stated otherwise) computed at the point  $(x, \theta) = (x_0, \theta_0)$ . The information (2.2.31) implies that

$$\partial_{\theta} w_3 = \partial_{w_1} \mathbb{W}_3(w_1, w_2, \theta_0) \ \partial_{\theta} w_1 + \partial_{w_2} \mathbb{W}_3(w_1, w_2, \theta_0) \ \partial_{\theta} w_2 \, .$$

Looking at (2.1.9), we note that either  $\partial_{\theta} w_1(x_0, \theta_0) \neq 0$  or  $\partial_{\theta} w_2(x_0, \theta_0) \neq 0$ . We will below consider the case  $\partial_{\theta} w_2(x_0, \theta_0) \neq 0$ . The other situation (that is  $\partial_{\theta} w_1 \neq 0$ ) can be dealt in a similar way.

The constraint (2.2.31) allows simplifications when writing (2.2.1), (2.2.2), (2.2.3) and (2.2.4). For example, the condition (2.2.1) reduces to

(2.2.32) 
$$\nabla \varphi \cdot \partial_{w_2} \mathbb{W} = -\frac{\partial_{\theta} w_1}{\partial_{\theta} w_2} \nabla \varphi \cdot \partial_{w_1} \mathbb{W}.$$

The condition (2.2.3) is nothing other than

$$(D_x w)^3 = \left[ (D_x w)^2 \,\partial_{w_1} \mathbb{W} \right] \otimes \nabla w_1 + \left[ (D_x w)^2 \,\partial_{w_2} \mathbb{W} \right] \otimes \nabla w_2 \equiv 0 \,.$$

Taking into account (2.2.25), this identity is possible only if

(2.2.33) 
$$(D_x w)^2 \,\partial_{w_1} \mathbb{W} \equiv 0, \qquad (D_x w)^2 \,\partial_{w_2} \mathbb{W} \equiv 0.$$

Defining  $\alpha := {}^t \nabla w_1 D_x w \partial_{w_1} \mathbb{W}$  and  $\beta := {}^t \nabla w_2 D_x w \partial_{w_1} \mathbb{W}$ , we find

$$(D_x w)^2 \partial_{w_1} \mathbb{W} = \alpha^t (1, 0, \partial_{w_1} \mathbb{W}_3) + \beta^t (0, 1, \partial_{w_2} \mathbb{W}_3).$$

The first constraint of (2.2.33) means that the two coefficients  $\alpha$  and  $\beta$  are zero, yielding

- $(2.2.34) \qquad (\nabla w_1 \cdot \partial_{w_1} \mathbb{W})^2 + (\nabla w_1 \cdot \partial_{w_2} \mathbb{W}) (\nabla w_2 \cdot \partial_{w_1} \mathbb{W}) = 0,$
- (2.2.35)  $(\nabla w_2 \cdot \partial_{w_1} \mathbb{W}) (\nabla w_1 \cdot \partial_{w_1} \mathbb{W} + \nabla w_2 \cdot \partial_{w_2} \mathbb{W}) = 0.$

By the same method followed this time at the level of the second condition, we can extract the necessary and sufficient conditions

$$(2.2.36) \qquad (\nabla w_1 \cdot \partial_{w_2} \mathbb{W}) (\nabla w_1 \cdot \partial_{w_1} \mathbb{W} + \nabla w_2 \cdot \partial_{w_2} \mathbb{W}) = 0$$

$$(2.2.37) \qquad (\nabla w_2 \cdot \partial_{w_2} \mathbb{W})^2 + (\nabla w_1 \cdot \partial_{w_2} \mathbb{W}) (\nabla w_2 \cdot \partial_{w_1} \mathbb{W}) = 0.$$

We claim that it is not possible to have

(2.2.38) 
$$\nabla w_1 \cdot \partial_{w_1} \mathbb{W} + \nabla w_2 \cdot \partial_{w_2} \mathbb{W} \neq 0.$$

Indeed, suppose that (2.2.38) is true. Then, the relations (2.2.35) and (2.2.36) imply that  $\nabla w_2 \cdot \partial_{w_1} \mathbb{W} = 0$  and that  $\nabla w_1 \cdot \partial_{w_2} \mathbb{W} = 0$ . Using these informations, the relations (2.2.34) and (2.2.37) lead to  $\nabla w_1 \cdot \partial_{w_1} \mathbb{W} = 0$  and  $\nabla w_2 \cdot \partial_{w_2} \mathbb{W} = 0$ . Now, these two last informations are in contradiction with (2.2.38). Therefore, we are sure that

(2.2.39) 
$$\nabla w_1 \cdot \partial_{w_1} \mathbb{W} + \nabla w_2 \cdot \partial_{w_2} \mathbb{W} = 0.$$

The condition (2.2.39) induces (2.2.35) and (2.2.36). It is also adjusted in such a way that (2.2.37) is equivalent to (2.2.34). Thus, the analysis of (2.2.3) is the same as the one of (2.2.34) and (2.2.39). These two constraints (2.2.34) and (2.2.39) say in particular that the two vectors

$$(\nabla w_1 \cdot \partial_{w_1} \mathbb{W}, \nabla w_2 \cdot \partial_{w_1} \mathbb{W}) \in \mathbb{R}^2, \qquad (\nabla w_1 \cdot \partial_{w_2} \mathbb{W}, \nabla w_2 \cdot \partial_{w_2} \mathbb{W}) \in \mathbb{R}^2$$

are colinear. In other words, we can find  $(\tilde{\alpha}, \tilde{\beta}) \in \mathbb{R}^2 \setminus (0, 0)$  such that

(2.2.40) 
$$\nabla w_1 \cdot (\tilde{\alpha} \ \partial_{w_1} \mathbb{W} + \tilde{\beta} \ \partial_{w_2} \mathbb{W}) = 0, \qquad \nabla w_2 \cdot (\tilde{\alpha} \ \partial_{w_1} \mathbb{W} + \tilde{\beta} \ \partial_{w_2} \mathbb{W}) = 0.$$

Now, we consider (2.2.2) computed at  $(x_0, \theta_0)$ . Exploiting the informations (2.2.31), (2.2.32) and (2.2.39), we can formulate (2.2.2) according to

(2.2.41) 
$$\begin{bmatrix} 2 \ \partial_{\theta} w_1 \ (\nabla w_1 \cdot \partial_{w_1} \mathbb{W}) + \partial_{\theta} w_2 \ (\nabla w_1 \cdot \partial_{w_2} \mathbb{W}) \\ - \frac{(\partial_{\theta} w_1)^2}{\partial_{\theta} w_2} \ (\nabla w_2 \cdot \partial_{w_1} \mathbb{W}) \end{bmatrix} (\nabla \varphi \cdot \partial_{w_1} \mathbb{W}) = 0.$$

Multiply (2.2.41) by  $\partial_{\theta} w_2$  ( $\nabla w_2 \cdot \partial_{w_1} W$ ). Then, use (2.2.34) and (2.2.39) to obtain

(2.2.42) 
$$(\nabla \varphi \cdot \partial_{w_1} \mathbb{W}) (\nabla w_2 \cdot \partial_{\theta} w)^2 = 0.$$

In the same way, multiply (2.2.41) by  $\partial_{\theta} w_2$  ( $\nabla w_1 \cdot \partial_{w_2} W$ ). Then, use (2.2.34) in order to extract

(2.2.43) 
$$(\nabla \varphi \cdot \partial_{w_1} \mathbb{W}) (\nabla w_1 \cdot \partial_{\theta} w)^2 = 0.$$

Two situations can happen :

- $ightarrow 1 
  ightarrow \frac{The\ case}{\nabla \varphi} \nabla \varphi \cdot \partial_{w_1} \mathbb{W} \neq 0$ . The equations (2.2.1), (2.2.42) and (2.2.43) imply that  $\nabla \varphi$ ,  $\nabla w_1$  and  $\nabla w_2$  belong to the same plane  $(\partial_{\theta} w^{\perp})$ . It follows that these vectors are linearly dependent, in contradiction with (2.2.30).
- $\triangleright 2 \triangleleft \underline{The \ case} \nabla \varphi \cdot \partial_{w_1} \mathbb{W} = 0.$  From (2.2.32), we deduce that  $\nabla \varphi \cdot \partial_{w_2} \mathbb{W} = 0.$  It follows that

(2.2.44) 
$$\nabla \varphi \cdot (\alpha' \partial_{w_1} \mathbb{W} + \beta' \partial_{w_2} \mathbb{W}) = 0, \quad \forall (\alpha', \beta') \in \mathbb{R}^2.$$

We choose  $\alpha' = \tilde{\alpha}$  and  $\beta' = \tilde{\beta}$ . According to the definition of the function  $\mathbb{W}$  and since  $(\tilde{\alpha}, \tilde{\beta}) \neq (0, 0)$ , we have

$$\tilde{\alpha} \partial_{w_1} \mathbb{W} + \beta \partial_{w_2} \mathbb{W} = {}^t (\tilde{\alpha}, \beta, \star) \neq (0, 0, 0).$$

The informations (2.2.40) and (2.2.44) (where  $\alpha' = \tilde{\alpha}$  and  $\beta' = \tilde{\beta}$ ) indicate that the vectors  $\nabla \varphi$ ,  $\nabla w_1$  and  $\nabla w_2$  belong to the same plane, namely  $(\tilde{\alpha} \partial_{w_1} \mathbb{W} + \tilde{\beta} \partial_{w_2} \mathbb{W})^{\perp}$ . It follows that these three vectors are linearly dependent. This is in contradiction with (2.2.30).

In conclusion, we have (2.2.29), as expected.

**Proposition 2.2.4** (local version of the Proposition 2.2.2). Assume (2.2.11) and select any  $\tilde{\theta} \in \mathbb{T}$ . Let  $(\varphi, w)$  be a couple which is compatible on  $\Gamma^{\tilde{\theta}}_{r,\tilde{r}}$ . Then, by selecting  $r \in \mathbb{R}^*_+$  and  $\tilde{r} \in ]0,1[$  conveniently and by permuting the directions  $x_1, x_2$ and  $x_3$  (with accordingly the components  $w_1$ ,  $w_2$  and  $w_3$  of w), it is possible to obtain (2.2.23) and to write the profile  $w(x, \theta)$  in the form

(2.2.45) 
$$w(x,\theta) = W(\varphi(x), w_1(x,\theta), \theta), \qquad (x,\theta) \in \Gamma^{\theta}_{r,\tilde{r}}$$

with a function  $W = {}^t(W_1, W_2, W_3) \in C^1(\mathbb{R}^2 \times ] - \tilde{r}, \tilde{r}[;\mathbb{R}^3)$  whose two first components  $W_1$  and  $W_2$  satisfy

(2.2.46) 
$$W_1(\varphi, w_1, \theta) = w_1, \quad \forall (\varphi, w_1, \theta) \in \mathbb{R}^2 \times ] - \tilde{r}, \tilde{r}[,$$

(2.2.47) 
$$\partial_{\varphi} W_2(\varphi, w_1, \theta) \neq 0, \quad \forall (\varphi, w_1, \theta) \in \mathbb{R}^2 \times ] - \tilde{r}, \tilde{r}[.$$

**Proof of proposition 2.2.4.** Without loss of generality, we can suppose that  $\tilde{\theta} = 0$ . Taking into account (2.2.28), the function  $w_3$  can be put in the form

$$w_3(x,\theta) = W_3(\varphi(x), w_1(x,\theta), \theta), \qquad \forall \ (x,\theta) \in \Gamma^0_{r,i}$$

with  $W_3(\varphi, w_1, \theta) := \mathbb{W}_3(w_1, W_2(\varphi, w_1, \theta), \theta)$ . In addition, we can define

$$W_1(\varphi, w_1, \theta) := w_1, \qquad \forall \ (\varphi, w_1, \theta) \in \mathbb{R}^2 \times ] - \tilde{r}, \tilde{r}[.$$

With these conventions, we recover both (2.2.45) and (2.2.46). Recalling (2.2.23), to have (2.2.11), the vector  $\partial_{\varphi} W \wedge \partial_{w_1} W$  must not vanish on  $\mathbb{R}^2 \times ] - \tilde{r}, \tilde{r}[$ . This amounts to saying that the function  $\partial_{\varphi} W_2$  does not vanish on  $\mathbb{R}^2 \times ] - \tilde{r}, \tilde{r}[$ . This is exactly what requires the condition (2.2.47).

#### 2.2.2.2 The proof of the Proposition 2.2.2.

Select a compatible couple  $(\varphi, w)$ . The condition (2.2.11) implies that

(2.2.48) 
$$\dim \operatorname{Vec} \langle \nabla w_1, \nabla w_2, \nabla w_3 \rangle = 2, \qquad \forall \ (x,\theta) \in \ \Omega^0_r \times \mathbb{T}$$

Locally, by permuting the directions  $x_1, x_2$  and  $x_3$  as it is made in the Proposition 2.2.4, we can get  $\nabla \varphi \in Vec \langle \nabla w_1, \nabla w_2 \rangle$ . It follows that the direction  $\nabla \varphi$  belongs to the vector space  $Vec \langle \nabla w_1, \nabla w_2, \nabla w_3 \rangle$ . Observe that this property does not depend on the choice of the coordinates. Thus, it remains to be true in all the domain under study. We must have

(2.2.49) 
$$\nabla \varphi \in Vec \langle \nabla w_1, \nabla w_2, \nabla w_3 \rangle, \quad \forall (x, \theta) \in \Omega^0_r \times \mathbb{T}$$

Fix  $\theta \in \mathbb{T}$ . Given a function  $\Psi_{\theta} \in C^1(\mathbb{R}^3; \mathbb{R})$  and  $r_{\theta} \in ]0, r[$ , introduce

$$\psi_{\theta}(x,\tilde{\theta}) := \Psi_{\theta}(w_1, w_2, w_3)(x, \tilde{\theta}), \qquad \forall (x, \tilde{\theta}) \in \Omega^0_{r_{\theta}} \times ]\theta - r_{\theta}, \theta + r_{\theta}[.$$

We can deduce from (2.2.48) and (2.2.49) the existence of  $\Psi_{\theta} \in C^1(\mathbb{R}^3; \mathbb{R})$  and  $r_{\theta} \in ]0, r[$  such that  $\nabla \varphi$  is not collinear to  $\nabla \psi_{\theta}$ , namely that the first component of  $\nabla \varphi \wedge \nabla \psi_{\theta}$  is positive

$$(\nabla \varphi \wedge \nabla \psi_{\theta})_{1} > 0, \qquad \forall (x, \theta) \in \Omega^{0}_{r_{\theta}} \times ]\theta - r_{\theta}, \theta + r_{\theta}[$$

whereas

 $\operatorname{Vec}\left\langle \nabla\varphi,\nabla\psi_{\theta}\right\rangle \equiv \operatorname{Vec}\left\langle \nabla w_{1},\nabla w_{2},\nabla w_{3}\right\rangle, \quad \forall \left(x,\theta\right)\in\Omega_{r_{\theta}}^{0}\times]\theta - r_{\theta},\theta + r_{\theta}[\,.$ 

The family of intervals  $]\theta - r_{\theta}, \theta + r_{\theta}[$  with  $\theta \in \mathbb{T}$  is an open cover of  $\mathbb{T}$ . Since  $\mathbb{T}$  is compact, there is a finite subcover  $\mathbb{T} \subset \bigcup_{i=1}^{N} ]\theta_i - r_{\theta_i}, \theta_i + r_{\theta_i}[$ . Now, consider some associated partition of unity  $\{\chi_i\}_{i=1}^{N}$  where the functions  $\chi_i \in C^{\infty}(\mathbb{T}; \mathbb{R}_+)$  are adjusted such that  $supp \chi_i \subset ]\theta_i - r_{\theta_i}, \theta_i + r_{\theta_i}[$  and  $\sum_{i=1}^{N} \chi_i \equiv 1$ . We replace  $r \in \mathbb{R}^*_+$  by the minimum of the numbers  $r_{\theta_i}$  (with  $i \in \{1, \dots, N\}$ ). Then, we can introduce  $\psi(x, \theta) := \sum_{i=1}^{N} \psi_{\theta_i}(x, \theta) \chi_i(\theta)$ . The preceding construction yields (2.2.15) as well as

$$(2.2.50) \qquad Vec \left\langle \nabla \varphi, \nabla \psi \right\rangle \equiv Vec \left\langle \nabla w_1, \nabla w_2, \nabla w_3 \right\rangle, \qquad \forall \left( x, \theta \right) \in \Omega^0_r \times \mathbb{T}.$$

The restriction (2.2.50) means that the three components  $w_i$  can be expressed as functions of  $\varphi$ ,  $\psi$  and  $\theta$ . At this level, we recover (2.2.13).

## **2.2.3** Necessary and sufficient constraints on $(\varphi, \psi, \mathbf{W})$ .

In this Section 2.2.3, we first show the Proposition 2.2.3, see the paragraph 2.2.3.1. Then, in the paragraph 2.2.3.2, we exclude the situations already examined in [7] and we precise the assumptions on  $(\varphi, \psi, \mathbf{W})$  to be retained.

#### 2.2.3.1 The proof of the Proposition 2.2.3.

The restriction (2.2.17) is just a repetition of (2.2.1). Concerning (2.2.18), it comes from the constraint (2.2.2) in which the matrix  $D_x w(x, \theta)$  is replaced as in (2.2.14). The relation (2.2.1) induces simplifications leading to (2.2.18). With (2.2.14), we can formulate (2.2.3) according to

$$(D_x w)^2 \partial_{\varphi} \mathbf{W} \otimes \nabla \varphi + (D_x w)^2 \partial_{\psi} \mathbf{W} \otimes \nabla \psi = 0.$$

Recall (2.2.15). The two vectors  $\nabla \varphi$  and  $\nabla \psi$  being independent, the above identity is equivalent to

- $(2.2.51) \qquad (D_x w)^2 \,\partial_{\varphi} \mathbf{W} = 0,$
- (2.2.52)  $(D_x w)^2 \partial_{\psi} \mathbf{W} = 0.$

Plug (2.2.14) into (2.2.51). Then, exploit (2.2.15) in order to extract

(2.2.53)  $(\nabla \varphi \cdot \partial_{\varphi} \mathbf{W})^{2} + (\nabla \psi \cdot \partial_{\varphi} \mathbf{W}) (\nabla \varphi \cdot \partial_{\psi} \mathbf{W}) = 0,$ 

(2.2.54) 
$$(\nabla \psi \cdot \partial_{\varphi} \mathbf{W}) (\nabla \psi \cdot \partial_{\psi} \mathbf{W} + \nabla \varphi \cdot \partial_{\varphi} \mathbf{W}) = 0.$$

We do the same with (2.2.52). This time, we get

(2.2.55)  $(\nabla \psi \cdot \partial_{\psi} \mathbf{W})^{2} + (\nabla \psi \cdot \partial_{\varphi} \mathbf{W}) (\nabla \varphi \cdot \partial_{\psi} \mathbf{W}) = 0,$ 

(2.2.56) 
$$(\nabla \varphi \cdot \partial_{\psi} \mathbf{W}) (\nabla \psi \cdot \partial_{\psi} \mathbf{W} + \nabla \varphi \cdot \partial_{\varphi} \mathbf{W}) = 0.$$

The relations (2.2.19) and (2.2.53) are similar. Observe that we cannot have  $\nabla \psi \cdot \partial_{\psi} \mathbf{W} + \nabla \varphi \cdot \partial_{\varphi} \mathbf{W} \neq 0$ . Indeed, in such a case, (2.2.53), (2.2.54), (2.2.55) and (2.2.56) would provide

 $\partial_{\varphi} \mathbf{W} \in \operatorname{Vec} \langle \nabla \varphi, \nabla \psi \rangle^{\perp} , \qquad \partial_{\psi} \mathbf{W} \in \operatorname{Vec} \langle \nabla \varphi, \nabla \psi \rangle^{\perp} .$ 

In other words, because of (2.2.15), the two vectors  $\partial_{\varphi} \mathbf{W}$  and  $\partial_{\psi} \mathbf{W}$  of  $\mathbb{R}^3$  would be collinear. This is clearly not coherent with (2.2.15). Therefore, we are sure to have (2.2.20).

Now, we have to show the opposite implication, that is the "only if" part of the Proposition 2.2.3. Using (2.2.13) and (2.2.15), the relations (2.2.1), (2.2.2), (2.2.11) and (2.1.9) are respectively equivalent to (2.2.17), (2.2.18), (2.2.15) and (2.2.16).

In addition, we have seen that looking at (2.2.3) is the same as imposing (2.2.53), (2.2.54), (2.2.55) and (2.2.56). The three conditions (2.2.53), (2.2.54) and (2.2.56) are taken into account at the level of (2.2.19) and (2.2.20). In view of (2.2.20), the remaining condition (2.2.55) reduces to (2.2.53).

It remains to check that the relation (2.2.4) is indeed a consequence of the constraints of the Proposition 2.2.3. To this end, use (2.2.14) in order to identify the different terms of (2.2.4). With (2.2.19) et (2.2.20), we can easily recover  $M(D_xw)^2 \equiv 0$ . Then, we can exploit (2.2.17) and (2.2.18) to obtain

$$D_x w M D_x w + (D_x w)^2 M = (\nabla \psi \cdot \partial_\theta w) (\nabla \varphi \cdot \partial_\varphi \mathbf{W} + \nabla \psi \cdot \partial_\psi \mathbf{W}) \partial_\psi \mathbf{W} \otimes \nabla \varphi.$$

In view of (2.2.20), we have (2.2.4). The proof of the Proposition 2.2.3 is finished.

#### 2.2.3.2 Further adjustments.

Before going further in the analysis, we must take care to deal with situations which are not considered in [7]. Noting  $\widetilde{\mathbf{W}}(x,\theta) := \mathbf{W}(\varphi(x),\psi(x,\theta),\theta)$ , the article [7] is based on the following condition, see (35) of [7]:

(2.2.57) 
$$\Pi_{\partial_{\theta}\widetilde{\mathbf{W}}^{\perp}} D_{x}\widetilde{\mathbf{W}} \Pi_{\nabla \varphi^{\perp}} \equiv \Pi_{(\partial_{\theta}\psi \, \partial_{\psi}\mathbf{W} + \partial_{\theta}\mathbf{W})^{\perp}} (\partial_{\psi}\mathbf{W} \otimes \nabla \psi) \Pi_{\nabla \varphi^{\perp}} \equiv 0.$$

Therefore, in order not to repeat what is made in [7], we have to work with  $\mathbf{W}$ ,  $\varphi$  and  $\psi$  adjusted such that  $\partial_{\psi} \mathbf{W} \wedge \partial_{\theta} \mathbf{W} \neq 0$  and  $\nabla \psi \wedge \nabla \varphi \neq 0$ .

There are different ways to factorize the profile  $w(x,\theta)$  as it is proposed in (2.2.13). Indeed, if  $\chi(\varphi,\psi,\theta) \in \mathcal{C}^{\infty}(\mathbb{R}^2 \times \mathbb{T};\mathbb{R})$  is any function such that  $\partial_{\psi}\chi \not\equiv 0$ , noting  $\tilde{\psi} := \chi(\varphi,\psi,\theta)$ , we find  $w \equiv \mathbf{W}(\varphi,\psi,\theta) \equiv \widetilde{\mathbf{W}}(\varphi,\tilde{\psi},\theta)$  with  $\mathbf{W}(\varphi,\psi,\theta) \equiv \widetilde{\mathbf{W}}(\varphi,\chi(\varphi,\psi,\theta),\theta)$ . Then, we find  $\partial_{\psi}\mathbf{W} \equiv \partial_{\psi}\chi \ \partial_{\tilde{\psi}}\widetilde{\mathbf{W}} \not\equiv 0$  together with

$$\partial_{\varphi} \mathbf{W} \equiv \partial_{\varphi} \widetilde{\mathbf{W}} + \partial_{\varphi} \chi \ \partial_{\psi} \mathbf{W} / \partial_{\psi} \chi \,, \qquad \partial_{\theta} \widetilde{\psi} \equiv \partial_{\theta} \psi \ \partial_{\psi} \chi + \partial_{\theta} \chi \,.$$

In this transformation, the conditions  $\partial_{\psi} \mathbf{W} \neq 0$  and  $\partial_{\varphi} \mathbf{W} \neq 0$  are preserved. On the other hand, we have some freedom concerning  $\partial_{\theta}\psi$ . By adjusting  $\chi$  conveniently, we can make sure that  $\partial_{\theta}\psi \neq 0$  or  $\partial_{\theta}\psi \equiv 0$ . According to circumstances, we will use one or other of these two conditions. In preparation for what follows, we put aside the framework (2.2.58) given below

(2.2.58)  $\partial_{\theta}\psi \neq 0, \quad \partial_{\varphi}\mathbf{W} \neq 0, \quad \partial_{\psi}\mathbf{W} \wedge \partial_{\theta}\mathbf{W} \neq 0, \quad \nabla\psi \wedge \nabla\varphi \neq 0.$ 

# **2.3** Compatible couples when $\nabla \varphi \cdot \partial_{\psi} \mathbf{W} \equiv 0$ .

We discuss here the system (2.2.17)-(2.2.18)-(2.2.19)-(2.2.20) under the restriction (2.2.58) and when  $\nabla \varphi \cdot \partial_{\psi} \mathbf{W} \equiv 0$ . In other words, we have to deal with (2.2.13), (2.2.15) and (2.2.58) combined with

(2.3.1)  $\nabla \varphi \cdot \partial_{\theta} \mathbf{W} = 0,$ 

(2.3.2) 
$$\nabla \varphi \cdot \partial_{\varphi} \mathbf{W} = 0$$

$$(2.3.3) \qquad \qquad \nabla \psi \cdot \partial_{\psi} \mathbf{W} = 0$$

(2.3.4) 
$$\nabla \varphi \cdot \partial_{\eta} \mathbf{W} = 0$$

# **2.3.1** The foliated structure associated to the phase $\varphi$ .

The phase  $\varphi$  must here inherit some special structure.

**Lemme 2.3.1.** Assume (2.1.9), (2.2.15) and (2.2.58) as well as (2.3.1), (2.3.2), (2.3.3) and (2.3.4). By restricting  $r \in \mathbb{R}^*_+$  and by permuting the coordinates  $x_1, x_2, x_3$  and the components  $\partial_1 \varphi, \partial_2 \varphi, \partial_3 \varphi$ , we can find two scalar functions  $f \in C^1(\mathbb{R}; \mathbb{R})$  and  $g \in C^1(\mathbb{R}; \mathbb{R})$  adjusted such that

(2.3.5) 
$$\nabla \varphi(x) \equiv {}^t (f \circ \varphi(x), 1, g \circ \varphi(x)) \,\partial_2 \varphi(x), \qquad \forall x \in \Omega^0_r.$$

**Proof of the Lemma 2.3.1.** The conditions (2.3.1) and (2.3.4) say that the direction  $\nabla \varphi$  is parallel to  $\partial_{\theta} \mathbf{W} \wedge \partial_{\psi} \mathbf{W} \neq 0$ . It follows that the direction  $\nabla \varphi$  can be viewed as a function of only  $(\varphi, \psi, \theta)$ . By restricting  $r \in \mathbb{R}^*_+$  and by permuting the coordinates  $x_1, x_2, x_3$  and the components  $\partial_1 \varphi, \partial_2 \varphi, \partial_3 \varphi$ , we can always recover

$$\nabla \varphi = E(\varphi, \psi, \theta) \ \partial_2 \varphi, \qquad E(\varphi, \psi, \theta) := {}^t \big( f(\varphi, \psi, \theta), 1, g(\varphi, \psi, \theta) \big)$$

Since the function  $\varphi$  does not depend on  $\theta$ , we must have  $\partial_{\theta}\psi \partial_{\psi}E + \partial_{\theta}E \equiv 0$ . When  $\partial_{\psi}E \equiv 0$ , we find also  $\partial_{\theta}E \equiv 0$  so that (2.3.5) is verified. From now on, we suppose that  $\partial_{\psi}E \neq 0$ .

The application  $\partial_{\theta}\psi$  can be represented as a function of only the variables  $(\varphi, \psi, \theta)$ , say  $\partial_{\theta}\psi = k(\varphi, \psi, \theta)$  with  $k \in C^1(\mathbb{R}^2 \times \mathbb{T}; \mathbb{R})$ . Consider any function  $\chi(\varphi, \psi, \theta)$ satisfying  $\partial_{\psi}\chi \neq 0$  and  $k \partial_{\psi}\chi + \partial_{\theta}\chi \equiv 0$ . Define  $\tilde{\psi} := \chi(\varphi, \psi, \theta)$ . We can change the set of independent variables  $(\varphi, \psi, \theta)$  into  $(\varphi, \tilde{\psi}, \theta)$  to find  $E(\varphi, \psi, \theta) \equiv \tilde{E}(\varphi, \tilde{\psi}, \theta)$ . Observe that

$$\partial_{\theta} \left[ E(\varphi, \psi, \theta) \right] \equiv \partial_{\theta} \psi \, \partial_{\psi} E + \partial_{\theta} E \equiv 0 \equiv \partial_{\theta} \left[ \tilde{E}(\varphi, \tilde{\psi}, \theta) \right] \equiv \partial_{\theta} \tilde{\psi} \, \partial_{\tilde{\psi}} \tilde{E} + \partial_{\theta} \tilde{E} \, .$$

By construction, we have  $\partial_{\theta} \tilde{\psi} \equiv 0$ . It follows that the function  $\tilde{E}$  does not depend on  $\theta$ . Retain that

(2.3.6) 
$$\nabla \varphi = \tilde{E}(\varphi, \tilde{\psi}) \ \partial_2 \varphi, \quad \tilde{E}(\varphi, \tilde{\psi}) := {}^t \left( \tilde{f}(\varphi, \tilde{\psi}), 1, \tilde{g}(\varphi, \tilde{\psi}) \right).$$

Since  $\partial_{\psi} E \equiv 0$ , we must have  $\partial_{\tilde{\psi}} \tilde{E} \neq 0$ . Writing  $\mathbf{W}(\varphi, \psi, \theta) \equiv \widetilde{\mathbf{W}}(\varphi, \tilde{\psi}, \theta)$ , we still have to deal with (2.3.1)-(2.3.2)-(2.3.3)-(2.3.4) but this time with  $\widetilde{\mathbf{W}}$  and  $\tilde{\psi}$  in place of  $\mathbf{W}$  and  $\psi$ . We decompose  $\widetilde{\mathbf{W}}$  into

(2.3.7) 
$$\widetilde{\mathbf{W}}(\varphi, \tilde{\psi}, \theta) = \alpha^{t}(0, -\tilde{g}, 1) + \beta^{t}(1, -\tilde{f}, 0) + \gamma^{t}(\tilde{f}, 1, \tilde{g})$$

where the three functions  $\alpha$ ,  $\beta$  and  $\gamma$  depend on  $\varphi$ ,  $\tilde{\psi}$  and  $\theta$ . The condition (2.3.1) yields  $\partial_{\theta}\gamma \equiv 0$ . On the other hand, the restriction (2.3.4) leads to

$$(2.3.8) \qquad \partial_{\tilde{\psi}}\gamma\left(\tilde{f}^2+1+\tilde{g}^2\right)-\alpha\,\partial_{\tilde{\psi}}\tilde{g}-\beta\,\partial_{\tilde{\psi}}\tilde{f}+\gamma\left(\tilde{f}\,\partial_{\tilde{\psi}}\tilde{f}+\tilde{g}\,\partial_{\tilde{\psi}}\tilde{g}\right) \equiv 0\,.$$

Taking the derivative of (2.3.8) with respect to  $\theta$ , we find

(2.3.9) 
$$\partial_{\theta} \alpha \, \partial_{\tilde{\psi}} \tilde{g} + \partial_{\theta} \beta \, \partial_{\tilde{\psi}} f \equiv 0 \, .$$

The symmetry of second derivatives expressed in the form  $\partial_{13}^2 \varphi \equiv \partial_{31}^2 \varphi$  can be traducted according to

(2.3.10) 
$$(-\partial_{\tilde{\psi}}\tilde{g},\tilde{f}\,\partial_{\tilde{\psi}}\tilde{g}-\tilde{g}\,\partial_{\tilde{\psi}}\tilde{f},\partial_{\tilde{\psi}}\tilde{f})\cdot{}^{t}(\partial_{1}\tilde{\psi},\partial_{2}\tilde{\psi},\partial_{3}\tilde{\psi}) \equiv 0.$$

Combining (2.3.9) and (2.3.10) with  $\partial_{\tilde{\psi}} \tilde{E} \neq 0$ , we can deduce that

(2.3.11) 
$$\nabla \tilde{\psi} \cdot \partial_{\theta} \widetilde{\mathbf{W}} \equiv \partial_{\theta} \beta \, \partial_1 \tilde{\psi} - (\tilde{f} \, \partial_{\theta} \beta + \tilde{g} \, \partial_{\theta} \alpha) \, \partial_2 \tilde{\psi} + \partial_{\theta} \alpha \, \partial_3 \tilde{\psi} \equiv 0 \, .$$

Recall that  $\nabla \varphi \wedge \nabla \tilde{\psi} \neq 0$ . Thus, the relations (2.3.1), (2.3.3), (2.3.4) and (2.3.11) indicate that the two vectors  $\partial_{\theta} \widetilde{\mathbf{W}}$  and  $\partial_{\tilde{\psi}} \widetilde{\mathbf{W}}$  are collinear. It follows that  $\partial_{\theta} \mathbf{W} \wedge \partial_{\psi} \widetilde{\mathbf{W}} = \partial_{\psi} \chi \ \partial_{\theta} \widetilde{\mathbf{W}} \wedge \partial_{\tilde{\psi}} \widetilde{\mathbf{W}} \equiv 0$ . This last information is clearly in contradiction with (2.2.58).

Recall here a basic result (see also [5, 7]) concerning (2.3.5).

**Lemme 2.3.2.** Select three functions  $f(\varphi)$ ,  $g(\varphi)$  and  $\varphi_{00}(x_2)$  in  $C^1(\mathbb{R};\mathbb{R})$ . Then, for  $r \in \mathbb{R}^*_+$  small enough, there is a unique expression  $\varphi(x) \in C^1(\Omega^0_r;\mathbb{R})$  satisfying (2.3.5), that is

(2.3.12) 
$$\partial_1 \varphi - f \circ \varphi \ \partial_2 \varphi = 0, \qquad \partial_3 \varphi - g \circ \varphi(x) \ \partial_2 \varphi = 0, \qquad \forall x \in \Omega^0_r$$

together with the initial data  $\varphi(0, x_2, 0) = \varphi_{00}(x_2)$  for all  $x_2 \in ] - r, r[$ .

**Proof of the Lemma 2.3.2.** The Cauchy problem for the first conservation law involved at the level of (2.3.12), namely

(2.3.13) 
$$\partial_1 \varphi_0 - f \circ \varphi_0 \ \partial_2 \varphi_0 = 0, \qquad \varphi_0(0, x_2) = \varphi_{00}(x_2)$$

has a local  $\mathcal{C}^1$  solution  $\varphi_0(x_1, x_2)$  near the point  $(0, 0) \in \mathbb{R}^2$ . Then, consider the local  $\mathcal{C}^1$  solution  $\varphi(x)$  of

(2.3.14) 
$$\partial_3 \varphi - g \circ \varphi(x) \ \partial_2 \varphi = 0, \qquad \varphi(x_1, x_2, 0) = \varphi_0(x_1, x_2).$$

To verify (2.3.5), it suffices now to check that  $\Xi := \partial_1 \varphi - f \circ \varphi(x) \partial_2 \varphi \equiv 0$  also when  $x_3 \neq 0$ . This property is in fact a consequence of the preceding construction which implies that

$$\partial_3 \Xi - g \circ \varphi(x) \ \partial_2 \Xi = g' \circ \varphi \ \partial_2 \varphi \ \Xi, \qquad \Xi(x_1, x_2, 0) = 0.$$

# **2.3.2** The description of $(\varphi, w)$ .

In this paragraph 2.3.2, the starting point is the description (2.3.7) which is based on some auxiliary function  $\psi(x)$  (not depending on  $\theta$ ). At this stage, we know that w can be put in the form

(2.3.15)  

$$w(x,\theta) = \mathbf{W}(\varphi(x),\psi(x),\theta) = \alpha(\varphi(x),\psi(x),\theta) \begin{pmatrix} 0\\ -g \circ \varphi(x)\\ 1 \end{pmatrix} + \beta(\varphi(x),\psi(x),\theta) \begin{pmatrix} 1\\ -f \circ \varphi(x)\\ 0 \end{pmatrix} + \gamma(\varphi(x),\psi(x),\theta) \begin{pmatrix} f \circ \varphi(x)\\ 1\\ g \circ \varphi(x) \end{pmatrix}$$

with a phase  $\varphi$  satisfying (2.3.12). It remains to adjust the ingredients  $\varphi$ ,  $\psi$  and **W** according to (2.3.1)-...-(2.3.4). We have already observed that the constraint

(2.3.1) is the same as  $\partial_{\theta}\gamma \equiv 0$ . In the same way, using again (2.3.12), the condition (2.3.4) is equivalent to  $\partial_{\psi}\gamma \equiv 0$ . Thus, the function  $\gamma$  depends only on the variable  $\varphi$ . Retain that  $\gamma(\varphi, \psi, \theta) \equiv \gamma(\varphi)$ .

Now, we can interpret the two remaining restrictions (2.3.2) and (2.3.3) into

- (2.3.16)  $-\alpha g' \beta f' + \gamma' (f^2 + 1 + g^2) + \gamma (f f' + g g') = 0,$
- (2.3.17)  $\partial_{\psi} \alpha \, \nabla \psi \cdot {}^t(0, -g, 1) + \partial_{\psi} \beta \, \nabla \psi \cdot {}^t(1, -f, 0) = 0.$

From (2.3.16), it is easy to extract

(2.3.18) 
$$\partial_{\theta} \alpha g' + \partial_{\theta} \beta f' \equiv 0, \qquad \partial_{\psi} \alpha g' + \partial_{\psi} \beta f' = 0.$$

The discussion about (2.3.16)-(2.3.17) is separated in two cases.

## **2.3.2.1** The case $f' \equiv g' \equiv 0$ .

By hypothesis, we have  $f \equiv a$  and  $g \equiv b$  with  $(a, b) \in \mathbb{R}^2$ . It follows that

(2.3.19) 
$$\varphi(x) = \varphi_{00}(a x_1 + x_2 + b x_3), \qquad \varphi_{00} \in \mathcal{C}^1(\mathbb{R}; \mathbb{R})$$

In view of (2.3.16), we have also  $\gamma \equiv c$  for some  $c \in \mathbb{R}$ . On the other hand, the function  $\psi(x)$  can always be put in the form

(2.3.20) 
$$\psi(x) = \Psi(x_1, x_3, a x_1 + x_2 + b x_3), \quad \Psi(X, Y, Z) \in \mathcal{C}^1(\mathbb{R}^3; \mathbb{R}).$$

Then, the condition (2.3.17) becomes the following scalar conservation law (implying Z and  $\theta$  as parameters)

(2.3.21) 
$$\partial_{\psi}\beta(\varphi_{00}(Z),\Psi,\theta) \ \partial_{X}\Psi + \partial_{\psi}\alpha(\varphi_{00}(Z),\Psi,\theta) \ \partial_{Y}\Psi \equiv 0.$$

At the level of (2.3.21), the variables Z and  $\theta$  play the part of parameters. Since  $\Psi(X, Y, Z)$  does not depend on  $\theta \in \mathbb{T}$ , we must have (when  $\partial_{\psi} \alpha \neq 0$ )

(2.3.22) 
$$\partial_{\psi}\beta = \chi(\varphi,\psi) \; \partial_{\psi}\alpha, \qquad \chi \in \mathcal{C}^1(\mathbb{R}^2;\mathbb{R}).$$

The equation (2.3.21) reduces to

(2.3.23) 
$$\chi(\varphi_{00}(Z),\Psi) \ \partial_X \Psi + \partial_Y \Psi \equiv 0.$$

We can sum up the situation when  $\nabla \varphi \cdot \partial_{\psi} \mathbf{W} \equiv 0$  and  $f' \equiv g' \equiv 0$  through the following result.

**Proposition 2.3.1.** Select any constants  $(a, b, c) \in \mathbb{R}^3$ . Select any smooth functions  $\varphi_{00}(Z)$ ,  $\chi(\varphi, \psi)$  and  $\alpha(\varphi, \psi, \theta)$ , any solutions  $\beta(\varphi, \psi, \theta)$  and  $\Psi(X, Y, Z)$  satisfying respectively (2.3.22) and (2.3.23). Define  $\varphi(x)$  and  $\psi(x)$  according to (2.3.19) and (2.3.20). Consider the function  $w(x, \theta)$  given by

(2.3.24) 
$$w = \alpha(\varphi, \psi, \theta) \begin{pmatrix} 0 \\ -b \\ 1 \end{pmatrix} + \beta(\varphi, \psi, \theta) \begin{pmatrix} 1 \\ -a \\ 0 \end{pmatrix} + c \begin{pmatrix} a \\ 1 \\ b \end{pmatrix}.$$

Then, the couple  $(\varphi, w)$  is compatible.

Take  $\varphi$  as indicated in (2.3.19). Given any function  $m \in \mathcal{C}^1(\mathbb{R} \times \mathbb{T}; \mathbb{R})$ , define

$$\beta(\varphi,\psi,\theta) := m(\varphi,\theta) + \varphi \int_0^{\psi} s \; (\partial_{\psi}\alpha)(\varphi,s,\theta) \; ds \, .$$

Then, we recover (2.3.21) with  $\psi(x) = x_1/(1+x_3\varphi(x))$ . The vectors  $\nabla\varphi$  and  $\nabla\psi$  are not collinear. By choosing *m* and  $\alpha$  conveniently, we can obtain

$$\partial_{\psi} \mathbf{W} \wedge \partial_{\theta} \mathbf{W} = (\partial_{\psi} \alpha \ \partial_{\theta} \beta - \partial_{\psi} \beta \ \partial_{\theta} \alpha)^{t} (a, 1, b)$$
  
=  $\partial_{\psi} \alpha \left( \partial_{\theta} m - \varphi \int_{0}^{\psi} \partial_{\theta} \alpha(\varphi, s, \theta) \ ds \right)^{t} (a, 1, b) \neq 0.$ 

The relation (2.2.57) is not satisfied. This example shows that the situations considered in Proposition 2.3.1 may not fall under the scope of [7].

Note that the support in (X, Y) of any non trivial solution  $\Psi \neq 0$  of (2.3.23) cannot be compact. Moreover, when  $\chi$  depends in a non linear way on  $\psi$ , due to the formation of singularities, the construction is valid only *locally*.

## **2.3.2.2** The case $f' \not\equiv 0$ or $g' \not\equiv 0$ .

In view of (2.3.18), we must have  $\partial_{\theta} \mathbf{W} \wedge \partial_{\psi} \mathbf{W} \equiv 0$  that implies (2.2.57). This situation is excluded at the level of (2.2.58) because it has been treated in [7].

Still, for the sake of completeness, we explain below what happens. We deal with the case  $f' \neq 0$ , the other situation  $(g' \neq 0)$  being similar. This time, seek the function  $\psi(x)$  in the form

(2.3.25) 
$$\psi(x) = \Psi(x_1, x_3, \varphi(x)), \qquad \Psi(X, Y, Z) \in \mathcal{C}^1(\mathbb{R}^3; \mathbb{R}).$$

From (2.3.18), extract  $\partial_{\psi}\beta$  in function of  $\partial_{\psi}\alpha$ . Plug the result into (2.3.17). Due to (2.2.58), we must have  $\partial_{\psi}\alpha \neq 0$ . Thus, it remains

(2.3.26) 
$$\Psi(X,Y,\varphi) = \Psi_0(g'(\varphi)Y + f'(\varphi)X), \qquad \Psi_0 \in \mathcal{C}^1(\mathbb{R};\mathbb{R}).$$

Thus, the variable  $\varphi$  being fixed, the function  $\Psi$  is constant on lines. Again, its support cannot be compact.

**Proposition 2.3.2.** Select functions f, g,  $\gamma$  and  $\Psi_0$  in  $C^1(\mathbb{R};\mathbb{R})$  with  $f' \neq 0$ . By applying the Lemma 2.3.2, we can construct a phase  $\varphi(x)$  which is solution of (2.3.12). Define the function  $\psi(x)$  as it is indicated in (2.3.25) and (2.3.26). Given any  $\alpha \in C^1(\mathbb{R}^2 \times \mathbb{T};\mathbb{R})$  with  $\partial_{\psi}\alpha \neq 0$ , define  $\beta \in C^1(\mathbb{R}^2 \times \mathbb{T};\mathbb{R})$  through the relation (2.3.16). Finally, consider the expression  $w(x,\theta)$  which is given by (2.3.15) where  $\gamma(\varphi, \psi, \theta) \equiv \gamma(\varphi)$ .

Then, the couple  $(\varphi, w)$  is compatible.

To illustrate the situation under study, we produce some example. Just take  $f(\varphi) = \varphi, g(\varphi) = \varphi^{-1}$  and  $\gamma(\varphi) \equiv 0$ . As a solution of (2.3.5), we can choose

$$\varphi(x) = \frac{1 - x_2}{2x_1} + \sqrt{\left(\frac{1 - x_2}{2x_1}\right)^2 - \frac{x_3}{x_1}}.$$

Concerning  $\psi$ , given any function  $\Xi \in \mathcal{C}^1(\mathbb{R}^3; \mathbb{R})$ , we can take

$$\psi(x) \equiv \psi(x, heta) = \Xi ig( arphi(x), x_2 + 2 \, arphi(x) \, x_1, x_3 - arphi(x)^2 \, x_1 ig) \, .$$

# **2.4** Compatible couples when $\nabla \varphi \cdot \partial_{\psi} \mathbf{W} \neq 0$ .

We discuss here the system (2.2.17)-(2.2.18)-(2.2.19)-(2.2.20) under the restriction (2.2.58) and when  $\nabla \varphi \cdot \partial_{\psi} \mathbf{W} \neq 0$ .

**Lemme 2.4.1.** Assume  $\nabla \varphi \cdot \partial_{\psi} \mathbf{W} \neq 0$ . The couple  $(\varphi, w)$  with w given by (2.2.13) is compatible if and only if there exists a function  $k(x, \theta)$  such that

(2.4.1) 
$$\nabla \varphi \cdot \partial_{\theta} w \equiv 0,$$

(2.4.2) 
$$\nabla \psi \cdot \partial_{\theta} w \equiv 0,$$

(2.4.3)  $\nabla \varphi \cdot \left( \partial_{\varphi} \mathbf{W} - k \; \partial_{\psi} \mathbf{W} \right) \equiv 0 \,,$ 

(2.4.4)  $\nabla \psi \cdot \left(\partial_{\varphi} \mathbf{W} - k \,\partial_{\psi} \mathbf{W}\right) \equiv 0,$ (2.4.5)  $\left(k \,\nabla (2 + \nabla \psi) - \partial_{\psi} \mathbf{W}\right) = 0,$ 

(2.4.5) 
$$(k \nabla \varphi + \nabla \psi) \cdot \partial_{\varphi} \mathbf{W} \equiv 0,$$

(2.4.6)  $(k \nabla \varphi + \nabla \psi) \cdot \partial_{\psi} \mathbf{W} \equiv 0.$ 

**Proof of the Lemma 2.4.1.** The relation (2.4.1) is a repetition of (2.2.17). When  $\nabla \varphi \cdot \partial_{\psi} \mathbf{W} \neq 0$ , the condition (2.2.18) amounts to the same thing as (2.4.2). On the other hand, from (2.2.19) and (2.2.20), we can extract

$$\left(\nabla\varphi\cdot\partial_{\psi}\mathbf{W}\right)\left(\nabla\psi\cdot\partial_{\varphi}\mathbf{W}\right) - \left(\nabla\psi\cdot\partial_{\psi}\mathbf{W}\right)\left(\nabla\varphi\cdot\partial_{\varphi}\mathbf{W}\right) \equiv 0$$

meaning that the vectors  ${}^{t}(\nabla \varphi \cdot \partial_{\psi} \mathbf{W}, \nabla \psi \cdot \partial_{\psi} \mathbf{W})$  and  ${}^{t}(\nabla \varphi \cdot \partial_{\varphi} \mathbf{W}, \nabla \psi \cdot \partial_{\varphi} \mathbf{W})$ are colinear. The second one can be obtained by multiplying the first one (which by hypothesis is not equal to zero) by a factor k. This is precisely (2.4.3) and (2.4.4). From (2.2.19) or (2.2.20) with (2.4.3) and (2.4.4), we can extract (2.4.5) and (2.4.6). Reciprocally, from the informations (2.4.3), (2.4.4), (2.4.5) and (2.4.6), it is easy to deduce (2.2.19) and (2.2.20).

# 2.4.1 Reduction of the problem : preliminaries.

The system (2.4.1)-...-(2.4.6) is not yet in a suitable form.

#### 2.4.1.1 Restatement of the problem.

Since  $\partial_{\psi} \mathbf{W} \neq 0$ , by permuting the coordinates, we can always suppose that  $\partial_{\psi} \mathbf{W}_2 \neq 0$ , allowing to exchange the variable  $\psi$  into  $\mathbf{W}_2(\varphi, \psi, \theta)$ . After this modification, we have to deal with

(2.4.7) 
$$\mathbf{W}(\varphi,\psi,\theta) = {}^{t} \left( \mathbf{V}(\varphi,\psi,\theta),\psi,\mathbf{W}_{3}(\varphi,\psi,\theta) \right), \qquad \mathbf{V} := \mathbf{W}_{1}.$$

Recall (2.2.58) which says in paticular that  $\partial_{\varphi} \mathbf{W} \neq 0$ . After permuting the two indices 1 and 3 (if necessary), we can suppose that  $\partial_{\varphi} \mathbf{W}_1 \equiv \partial_{\varphi} \mathbf{V} \neq 0$ . It follows that we can regard  $\mathbf{W}_3$  as a function of  $(\psi, \mathbf{V}, \theta)$ . In other words, we can find some function  $\mathfrak{L}(\psi, \mathbf{V}, \theta) \in \mathcal{C}^1(\mathbb{R}^2 \times \mathbb{T}; \mathbb{R})$  such that

(2.4.8) 
$$\mathbf{W}(\varphi,\psi,\theta) = {}^{t} \left( \mathbf{V}(\varphi,\psi,\theta), \psi, \mathfrak{L}(\psi,\mathbf{V}(\varphi,\psi,\theta),\theta) \right).$$

Using (2.2.58) together with (2.4.1), (2.4.2), (2.4.3) and (2.4.4), we can see that the two vectors  $\partial_{\theta} w$  and  $\partial_{\varphi} \mathbf{W} - k \partial_{\psi} \mathbf{W}$  are collinear meaning that there is a function  $\beta(x, \theta)$  which is adjusted such that

(2.4.9) 
$$\partial_{\varphi} \mathbf{W} - \tilde{k} \, \partial_{\psi} \mathbf{W} = \beta \, \partial_{\theta} \mathbf{W}, \qquad \tilde{k} := k + \beta \, \partial_{\theta} \psi.$$

Knowing (2.4.8), this information (2.4.9) becomes

(2.4.10) 
$$\beta \ \partial_{\theta} \mathfrak{L} \equiv 0, \qquad \partial_{\varphi} \mathbf{V} = \beta \ \partial_{\theta} \mathbf{V}, \qquad k = -\beta \ \partial_{\theta} \psi.$$

Since  $\partial_{\varphi} \mathbf{V} \neq 0$ , we must have  $\beta \neq 0$  and  $\partial_{\theta} \mathbf{V} \neq 0$ , this last condition being also a consequence of (2.2.58). Necessarily, we must have  $\partial_{\theta} \mathfrak{L} \equiv 0$ . Introduce

(2.4.11) 
$$v(x,\theta) := \mathbf{V}(\varphi(x),\psi(x,\theta),\theta), \quad v \in \mathcal{C}^1(\Omega^0_r \times \mathbb{T};\mathbb{R}).$$

Note simply

$$\begin{split} \mathfrak{L}(\psi, v) &= \mathfrak{L}(\psi, v)(x, \theta) &:= \mathfrak{L}\big(\psi(x, \theta), v(x, \theta)\big), \\ \partial_{\psi} \mathfrak{L}(\psi, v) &:= \partial_{\psi} \mathfrak{L}\big(\psi(x, \theta), v(x, \theta)\big), \\ \partial_{v} \mathfrak{L}(\psi, v) &:= \partial_{\mathbf{V}} \mathfrak{L}\big(\psi(x, \theta), v(x, \theta)\big). \end{split}$$

Observe that  $\nabla \varphi \cdot {}^{t}(1, 0, \partial_{v} \mathfrak{L}) \partial_{\theta} \mathbf{V} = \nabla \varphi \cdot \partial_{\theta} w - \nabla \varphi \cdot \partial_{\psi} \mathbf{W} \partial_{\theta} \psi$ . In view of the restriction (2.4.1), the condition  $\partial_{\theta} \psi \neq 0$  of (2.2.58) and the hypothesis  $\nabla \varphi \cdot \partial_{\psi} \mathbf{W} \neq 0$ , we are sure that  $\partial_{\theta} \mathbf{V} \neq 0$ . Retain that

(2.4.12) 
$$\partial_{\varphi} \mathbf{V} \neq 0, \qquad \partial_{\theta} \mathbf{V} \neq 0, \qquad \lambda \equiv -\partial_{\theta} \psi \partial_{\varphi} \mathbf{V} / \partial_{\theta} \mathbf{V}.$$

**Proposition 2.4.1.** Assume (2.2.58), (2.4.8) and  $\nabla \varphi \cdot \partial_{\psi} \mathbf{W} \neq 0$ . Then, the function  $\mathfrak{L}$  does not depend on the variable  $\theta \in \mathbb{T}$  and we have (2.4.12). Moreover, the system (2.4.1)-...-(2.4.6) is equivalent to

(2.4.13) 
$$\partial_{\theta} v \left[ \partial_{1} \varphi + \partial_{v} \mathfrak{L}(\psi, v) \partial_{3} \varphi \right] + \partial_{\theta} \psi \left[ \partial_{2} \varphi + \partial_{\psi} \mathfrak{L}(\psi, v) \partial_{3} \varphi \right] \equiv 0,$$

(2.4.14) 
$$\partial_{\theta} v \left[ \partial_{1} \psi + \partial_{v} \mathfrak{L}(\psi, v) \partial_{3} \psi \right] + \partial_{\theta} \psi \left[ \partial_{2} \psi + \partial_{\psi} \mathfrak{L}(\psi, v) \partial_{3} \psi \right] \equiv 0,$$

(2.4.15) 
$$\partial_1 \psi + \partial_v \mathfrak{L}(\psi, v) \,\partial_3 \psi - \partial_\theta \psi \, \frac{\partial_\varphi \mathbf{V}}{\partial_\theta \mathbf{V}} \left[ \partial_1 \varphi + \partial_v \mathfrak{L}(\psi, v) \,\partial_3 \varphi \right] \equiv 0 \,,$$

where v is given by (2.4.11) whereas  $\partial_{\varphi} \mathbf{V}$  and  $\partial_{\theta} \mathbf{V}$  are computed at  $(\varphi, \psi, \theta)$ .

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In view of (2.4.11), from (2.4.13) and (2.4.14), we can easily deduce that

$$(2.4.16) \qquad \partial_{\theta} v \left[ \partial_{1} v + \partial_{v} \mathfrak{L}(\psi, v) \partial_{3} v \right] + \partial_{\theta} \psi \left[ \partial_{2} v + \partial_{\psi} \mathfrak{L}(\psi, v) \partial_{3} v \right] \equiv 0.$$

**Proof of the Proposition 2.4.1.** We have already seen that the function  $\mathfrak{L}$  does not depend on the variable  $\theta \in \mathbb{T}$  and that the conditions inside (2.4.12) are verified. By construction, we know also that

(2.4.17) 
$$w(x,\theta) = {}^{t} \left( v(x,\theta), \psi(x,\theta), \mathfrak{L}(\psi,v)(x,\theta) \right), \qquad \partial_{\theta} w \neq 0.$$

Taking into account (2.2.58) and (2.4.8), the two constraints (2.4.1) and (2.4.2) are equivalent to the existence of some (nonzero) function  $\alpha(x, \theta)$  such that

(2.4.18) 
$$\partial_{\theta} v \begin{pmatrix} 1\\ 0\\ \partial_{v} \mathfrak{L} \end{pmatrix} + \partial_{\theta} \psi \begin{pmatrix} 0\\ 1\\ \partial_{\psi} \mathfrak{L} \end{pmatrix} = \alpha \nabla \varphi \wedge \nabla \psi.$$

Since  $\partial_{\varphi} \mathbf{V} \neq 0$ , we have  $\partial_{\varphi} \mathbf{W} \wedge \partial_{\psi} \mathbf{W} = \partial_{\varphi} \mathbf{V}^{t} (-\partial_{v} \mathfrak{L}, -\partial_{\psi} \mathfrak{L}, 1) \neq 0$ . Combining this with (2.4.5), (2.4.6) and (2.4.12) yields the existence of some (nonzero) scalar function  $\gamma(x, \theta)$  such that

(2.4.19) 
$$-\partial_{\theta}\psi \ \frac{\partial_{\varphi}\mathbf{V}(\varphi(x),\psi(x,\theta),\theta)}{\partial_{\theta}\mathbf{V}(\varphi(x),\psi(x,\theta),\theta)} \ \nabla\varphi + \nabla\psi = \gamma \left(\begin{array}{c} \partial_{v}\mathfrak{L} \\ \partial_{\psi}\mathfrak{L} \\ -1 \end{array}\right).$$

Plug the expression  $\nabla \psi$  given by (2.4.19) into (2.4.18) in order to extract

(2.4.20) 
$$\partial_{\theta} v = -\alpha \gamma \left( \partial_2 \varphi + \partial_{\psi} \mathfrak{L} \, \partial_3 \varphi \right),$$

(2.4.21) 
$$\partial_{\theta}\psi = +\alpha \gamma \left(\partial_{1}\varphi + \partial_{v}\mathfrak{L} \ \partial_{3}\varphi\right),$$

(2.4.22) 
$$\partial_{\theta} v \, \partial_{v} \mathfrak{L} + \partial_{\theta} \psi \, \partial_{\psi} \mathfrak{L} = + \alpha \, \gamma \left( \partial_{\psi} \mathfrak{L} \, \partial_{1} \varphi - \partial_{v} \mathfrak{L} \, \partial_{2} \varphi \right).$$

Since  $\alpha \gamma \neq 0$  (because  $\partial_{\theta} \psi \neq 0$ ), from (2.4.20) and (2.4.21), we can deduce (2.4.13). On the other hand, the relation (2.4.22) provides no new information because it is a linear combination of (2.4.20) and (2.4.21). Observe that

$$(1,0,\partial_{\nu}\mathfrak{L}) \cdot {}^{t}(\partial_{\nu}\mathfrak{L},\partial_{\psi}\mathfrak{L},-1) \equiv 0, \qquad (0,1,\partial_{\psi}\mathfrak{L}) \cdot {}^{t}(\partial_{\nu}\mathfrak{L},\partial_{\psi}\mathfrak{L},-1) \equiv 0.$$

Using these identities and (2.4.13), coming back to (2.4.19) multiplied by the non zero vector valued function  $\partial_{\theta} w$ , we can obtain (2.4.14). The last condition (2.4.15) is just the product of (2.4.19) with the vector  ${}^{t}(1, 0, \partial_{v}\mathfrak{L})$ .

Conversely, suppose that  $\varphi(x)$  and  $\psi(x,\theta)$  are such that  $\nabla \varphi \wedge \nabla \psi \neq 0$  and satisfy (locally) the system (2.4.13)-(2.4.14)-(2.4.15) for some functions  $\mathfrak{L}(\psi, v)$  and  $\mathbf{V}(\varphi, \psi, \theta)$ . Define v and w as in (2.4.11) and (2.4.17).

Both (2.4.1) and (2.4.2) become a direct consequence of (2.4.13) and (2.4.14). We can obtain (2.4.9), that is (2.4.3) and (2.4.4), through (2.4.10) by adjusting the coefficient  $\beta$  (and then k) conveniently. At this stage, the interpretation of (2.4.13)-(2.4.14)-(2.4.15) is that the vector on the left of (2.4.19) is orthogonal to the direction  ${}^{t}(1, 0, \partial_{v}\mathfrak{L})$  and  ${}^{t}(0, 1, \partial_{\psi}\mathfrak{L})$ . Thus, we must have (2.4.19) for some coefficient  $\gamma$ . This is exactly (2.4.5) and (2.4.6).

Since  $\partial_{\theta} w \neq 0$ , by a small rotation in the space variable  $x \in \mathbb{R}^3$ , we can always obtain that  $\partial_{\theta} v \neq 0$ . All the restrictions in (2.4.12) are stable under such a modification (if it is small enough). In what follows, we work locally in  $(x, \theta)$  under the assumptions  $\partial_{\theta} v \neq 0$  and (2.4.12). We will exploit these informations in order to perform different changes of variables which are crucial when discussing the content of (2.4.13)-(2.4.14)-(2.4.15).

#### 2.4.1.2 Various changes of variables.

Subtract (2.4.15) from (2.4.14), use (2.4.13) to replace  $\partial_1 \varphi + \partial_v \mathfrak{L}(\psi, v) \partial_3 \varphi$ , and then exploit (2.2.58) to make simplifications in order to extract the identity

(2.4.23) 
$$(\partial_2 \psi + \partial_{\psi} \mathfrak{L} \ \partial_3 \psi) / \partial_{\theta} \psi = \partial_{\varphi} \mathbf{V} \ (\partial_2 \varphi + \partial_{\psi} \mathfrak{L} \ \partial_3 \varphi) / \partial_{\theta} \mathbf{V} .$$

The identity a/b = c/d implies that  $a/b = (c + \gamma a)/(d + \gamma b)$  for all  $\gamma \in \mathbb{R}$ . This implication applied to (2.4.23) with  $\gamma = \partial_{\psi} \mathbf{V}$  furnishes

$$\frac{\partial_2 \psi + \partial_\psi \mathfrak{L} \ \partial_3 \psi}{\partial_\theta \psi} = \frac{\partial_\varphi \mathbf{V} \ (\partial_2 \varphi + \partial_\psi \mathfrak{L} \ \partial_3 \varphi) + \partial_\psi \mathbf{V} \ (\partial_2 \psi + \partial_\psi \mathfrak{L} \ \partial_3 \psi)}{\partial_\theta \mathbf{V} + \partial_\psi \mathbf{V} \ \partial_\theta \psi}$$

Recalling (2.4.11), this is the same as

(2.4.24) 
$$(\partial_2 \psi + \partial_{\psi} \mathfrak{L} \ \partial_3 \psi) / \partial_{\theta} \psi = (\partial_2 v + \partial_{\psi} \mathfrak{L} \ \partial_3 v) / \partial_{\theta} v \,.$$

Since  $\partial_{\theta} v \neq 0$ , we can work (locally) with the variables (x, v) in place of  $(x, \theta)$ . The function  $\varphi$  does not depend on  $\theta \in \mathbb{T}$  and therefore it does not depend on v. On the contrary, the function  $\psi$  can be put in the form

(2.4.25) 
$$\psi(x,\theta) = u(x,v(x,\theta)), \quad \partial_v u \neq 0.$$

Formulating (2.4.24) at the level of u(x, v) yields

(2.4.26) 
$$\partial_2 u + \partial_u \mathfrak{L}(u, v) \ \partial_3 u = 0.$$

Recalling (2.4.16) and exploiting (2.4.26), the constraint (2.4.14) becomes

(2.4.27) 
$$\partial_1 u + \partial_v \mathfrak{L}(u, v) \ \partial_3 u = 0.$$

Knowing what is the function u(x, v), it is not complicated to obtain  $v(x, \theta)$ . To this end, it suffices to consider the scalar conservation law

(2.4.28) 
$$\begin{array}{l} \partial_1 v + \partial_v \mathfrak{L}(u(x,v),v) \ \partial_3 v \\ + \partial_v u(x,v) \ \left[ \partial_2 v + \partial_u \mathfrak{L}(u(x,v),v) \ \partial_3 v \right] \equiv 0 \,. \end{array}$$

To sum up, the system (2.4.13)-(2.4.14)-(2.4.15) amounts to the same thing as to identify the two expressions u(x, v) and  $v(x, \theta)$  as it is explained above and then to focus on the remaining constraint, namely

(2.4.29) 
$$\begin{array}{l} \partial_1 \varphi + \partial_v \mathfrak{L}(u(x,v),v) \ \partial_3 \varphi \\ + \partial_v u(x,v) \ \left[ \partial_2 \varphi + \partial_u \mathfrak{L}(u(x,v),v) \ \partial_3 \varphi \right] \ \equiv \ 0 \,. \end{array}$$

Recall that  $v \in K \subset \mathbb{R}$  must be seen here, at the level of (2.4.29), as a parameter. Thus, all the difficulty is to solve (2.4.29) with a phase  $\varphi(x)$  which does not depend on v. We first explain what happens when  $\partial_3 u \equiv 0$ . Then, we present the problematic when  $\partial_3 u \neq 0$ .

• <u>The</u> case  $\left\lfloor \partial_3 u \equiv 0 \right\rfloor$ . In view of (2.4.26) and (2.4.27), we have  $\nabla_x u \equiv 0$ . It follows that u(x,v) = U(v) with a function  $U \in \mathcal{C}^1(K;\mathbb{R})$  such that  $U' \neq 0$ . Necessarily, the function **V** depends only on the variable  $\psi$ . This is a contradiction with (2.4.12). For the sake of completeness, we still describe below what happens when  $\partial_3 u \equiv 0$  and  $U' \neq 0$ . Noting  $\tilde{\mathfrak{L}}(v) := \mathfrak{L}(U(v), v)$ , we can see that (2.4.29) becomes

(2.4.30) 
$$\partial_1 \varphi(x) + U'(v) \ \partial_2 \varphi(x) + \hat{\mathcal{L}}'(v) \ \partial_3 \varphi(x) \equiv 0.$$

Recall that the variables x and v are independent. Thus, the relation (2.4.30) implies that  $\tilde{\mathfrak{L}}(v) \equiv \mathfrak{L}(U(v), v) = a v + b U(v) + c$  for some  $(a, b, c) \in \mathbb{R}^3$ . We have to deal with

(2.4.31) 
$$\partial_1 \varphi + a \, \partial_3 \varphi + U'(v) \, (\partial_2 \varphi + b \, \partial_3 \varphi) \equiv 0.$$

This is possible only if  $U' \equiv c \in \mathbb{R}$  and

$$\varphi(x) = \Phi(c x_1 - x_2, (a + b c) x_2 - c x_3), \qquad \Phi \in \mathcal{C}^1(\mathbb{R} \times \mathbb{R}; \mathbb{R}).$$

On the other hand, the function v can be obtained through

(2.4.32) 
$$\partial_1 v + a \,\partial_3 v + U'(v) \,(\partial_2 v + b \,\partial_3 v) = 0, \qquad v(0, \cdot) = v_0.$$

By varying the ingredients  $a, b, U, \Phi$  and  $v_0$ , we can obtain a whole class of solutions to the system (2.4.13)-(2.4.14)-(2.4.15).

• <u>The</u> case  $\partial_3 u \neq 0$ . Since  $\partial_3 u \neq 0$ , we can exchange the variables (x, v) into  $(x_1, x_2, u, v)$ . In particular, the applications  $\varphi$ ,  $\partial_v u$  and  $\partial_3 u$  can be regarded as functions of  $(x_1, x_2, u, v)$  instead of (x, v). Taking this point of view into account, we adopt the following conventions

(2.4.33) 
$$\varphi(x) = \Phi(x_1, x_2, u(x, v), v), \qquad \forall (x, v),$$

(2.4.34) 
$$\partial_v u(x,v) = R(x_1, x_2, u(x,v), v), \qquad \forall (x,v),$$

(2.4.35) 
$$\partial_3 u(x,v) = S(x_1, x_2, u(x,v), v), \qquad \forall (x,v).$$

Recall that  $R \neq 0$  and  $S \neq 0$ . The constraint (2.4.29) becomes

(2.4.36)  $X \Phi \equiv 0, \qquad X := \partial_1 + R \partial_2,$ 

whereas the fact that  $\varphi$  does not depend on v amounts to the same thing as

(2.4.37)  $Y \Phi \equiv 0, \qquad Y := R \partial_u + \partial_v.$ 

The rest of this chapter 2.4 is devoted to the case  $\partial_3 u \neq 0$ . Thus, it should be clearly noted here what the current matter is.

**Remaining work.** When  $\partial_3 u \neq 0$ , the problem is to find a <u>non</u> <u>constant</u> function  $\Phi(x_1, x_2, u, v)$  satisfying the transport equations (2.4.36) and (2.4.37) with a coefficient  $R(x_1, x_2, u, v)$  issued from (2.4.26), (2.4.27) and (2.4.34).

Forcing the presence of u and v at the level of  $\varphi$  and passing through (2.4.37) to express that  $\partial_{\theta}\varphi \equiv 0$  may seem unnatural. However, this process allows to simplify the equation (2.4.29). It leads to the above problem which, to our knowledge, is original. The strategy to solve it is the following.

In the Section 2.4.2, we extract from (2.4.36)-(2.4.37) the necessary and sufficient conditions (2.4.39) and (2.4.40) to impose on R. In the Section 2.4.3, we exhibit the special form (2.4.51) of a coefficient R coming from (2.4.26), (2.4.27) and (2.4.34). In the Section 2.4.4, we test our criteria (2.4.39) and (2.4.40) on the functions R which conform to (2.4.51). All requirements are met in different cases leading to a classification of all compatible couples (when  $\nabla \varphi \cdot \partial_{\psi} \mathbf{W} \neq 0$  and  $\partial_3 u \neq 0$ ). Illustrative examples are proposed in the Section 2.4.6.

# 2.4.2 Reduction of the problem : geometrical step.

The existence of a non constant solution to (2.4.36)-(2.4.37) relies deeply on the geometrical properties of the two vector fields X and Y. Introduce the Lie algebra  $\mathcal{A}$  generated by the successive Poisson brackets of X and Y. The dimension being 4, we find here

 $\mathcal{A} \equiv \langle X, Y, [X;Y], [X;[X;Y]], [Y;[X;Y]] \rangle.$ 

**Proposition 2.4.2.** The system (2.4.36)-(2.4.37) has a non constant solution  $\Phi$  if and only if the dimension of A is strictly less than 4. Two different situations may occur :

i) dim  $\mathcal{A} = 2$ . The function  $\Phi(x_1, x_2, u, v)$  depends on two independent variables. Then, the coefficient R must satisfy :

 $(2.4.38) \qquad \partial_1 R + R \ \partial_2 R \equiv X \ R \equiv 0, \qquad R \ \partial_u R + \partial_v R \equiv Y \ R \equiv 0.$ 

ii) dim  $\mathcal{A} = 3$ . The function  $\Phi(x_1, x_2, u, v)$  depends on one variable. The coefficient R must satisfy  $XR \neq 0$  or  $YR \neq 0$  together with

- (2.4.39) (XR) YXR 2 (XR) XYR + (YR) X<sup>2</sup>R = 0,
- (2.4.40) (YR) XYR 2 (YR) YXR + (XR) Y<sup>2</sup>R = 0.

**Proof of Proposition 2.4.2.** Recall that the Poisson bracket of the vector fields X and Y is the vector field [X; Y] which is adjusted such that

 $[X;Y] f = -YR \partial_2 f + XR \partial_u f = XYf - YXf, \qquad \forall f \in \mathcal{C}^{\infty}(\mathbb{R}^4;\mathbb{R}).$ 

From (2.4.36) and (2.4.37), it is easy to infer that  $Z \Phi \equiv 0$  for all  $Z \in \mathcal{A}$ . Thus, when  $\dim \mathcal{A} = 4$ , the function  $\Phi$  is constant. It means that the phase  $\varphi$  is stationary, in contradiction with (2.1.4). We examine the other situations.

i) dim  $\mathcal{A} = 2$ . This situation can occur if and only if [X; Y] is a linear combination of X and Y, giving rise to (2.4.38). By applying the Frobenius Theorem [2], we see that the field of planes  $Vec \langle X, Y \rangle$  is associated with a foliated structure of  $\mathbb{R}^4$  by submanifolds of dimension 2 along which  $\Phi$  must be constant. Clearly, the function  $\Phi$  inherits two degrees of freedom. In particular, it can be a non constant solution of the system (2.4.36)-(2.4.37).

i) dim  $\mathcal{A} = 3$ . To avoid (2.4.38), we have to require  $XR \neq 0$  or  $YR \neq 0$ . Then, to obtain  $\dim \mathcal{A} = 3$ , it is necessary to impose

(2.4.41) 
$$[X; [X; Y]] \in Vec \langle X, Y, [X; Y] \rangle,$$

$$(2.4.42) \qquad \qquad \left[Y; [X;Y]\right] \in Vec\left\langle X, Y, [X;Y]\right\rangle$$

Under the conditions (2.4.41) and (2.4.42), we find that  $\dim \mathcal{A} = 3$ . By applying again the Frobenius Theorem [2], we can see that the field of hyperplanes  $\operatorname{Vec} \langle X, Y, [X; Y] \rangle$  is associated with a foliated structure of the space  $\mathbb{R}^4$  by hypersurfaces along which  $\Phi$  must be constant. On the other hand, the function  $\Phi$ can actually vary in the directions which are transversal to these hypersurfaces. Now, it remains to convert (2.4.41)-(2.4.42) in the form of constraints implying the coefficient R. To this end, compute

$$\begin{bmatrix} X; [X;Y] \end{bmatrix} f = (-2XYR + YXR) \partial_2 f + X^2R \partial_u f, \\ \begin{bmatrix} Y; [X;Y] \end{bmatrix} f = -Y^2R \partial_2 f + (2YXR - XYR) \partial_u f.$$

Taking into acount these informations combined with the specific forms of X, Y and [X;Y], we can deduce that the two constraints (2.4.41) and (2.4.42) can be verified on condition that [X;Y] is collinear to both [X;[X;Y]] and [Y;[X;Y]]. This remark, leads directly to (2.4.39) and (2.4.40).

Given some initial data  $R(0, x_2, u, 0) := R_{00}(x_2, u)$ , we can solve the system of two conservation laws (2.4.38) in the same way as in the Lemma 2.3.2. Then, to recover  $\Phi$ , it suffices to fix any function  $\Phi_{00}(x_2, u)$  satisfying  $\nabla_{x_2,u} \Phi_{00} \neq 0$  and to integrate the two equations

$$(2.4.43) \qquad \qquad \partial_1 \Phi + R \ \partial_2 \Phi \equiv 0, \qquad R \ \partial_u \Phi + \partial_v \Phi \equiv 0.$$

The discussion about (2.4.39)-(2.4.40) is delicate. We explain below how to construct R and  $\Phi$  in the more general situation (when  $XR YR \neq 0$ ).

**Lemme 2.4.2.** Fix any (non zero) function  $\mathcal{Q}(y, R, \Phi) \in \mathcal{C}^3(\mathbb{R}^3; \mathbb{R}^*)$ . Select any couple of functions  $R_{00}(x_1, x_2) \in \mathcal{C}^1(\mathbb{R}^2; \mathbb{R})$  and  $\Phi_{00}(x_1, x_2) \in \mathcal{C}^1(\mathbb{R}^2; \mathbb{R})$  satisfying  $\nabla_{x_1, x_2} \Phi_{00} \neq 0$  as well as

$$(2.4.44) \qquad \qquad \partial_1 R_{00} + R_{00} \ \partial_2 R_{00} \neq 0, \qquad \partial_1 \Phi_{00} + R_{00} \ \partial_2 \Phi_{00} \equiv 0.$$

Then, the system (2.4.39)-(2.4.40) has a solution  $R(x_1, x_2, u, v)$  such that

(2.4.45) 
$$R(x_1, x_2, 0, 0) = R_{00}(x_1, x_2), \qquad X R \neq 0, \qquad Y R \neq 0.$$

Moreover, there exists a non constant solution  $\Phi$  of (2.4.36)-(2.4.37) such that

(2.4.46) 
$$\Phi(x_1, x_2, 0, 0) = \Phi_{00}(x_1, x_2), \qquad \frac{YR}{XR} \equiv \mathcal{Q}(Rx_1 - x_2, R, \Phi).$$

**Proof of Lemma 2.4.2.** We start by studying a little more the structure of the system (2.4.39)-(2.4.40). Since  $XR \neq 0$ , we can introduce the quantity Q := YR/XR. In fact, the restrictions (2.4.39) and (2.4.40) are equivalent to

$$(2.4.47) YQ - Q XQ = 0,$$

$$(2.4.48) -[Y;X]R + (XQ)(XR) = 0.$$

Since  $YR \neq 0$  whereas  $X\Phi \equiv 0$ , we can always consider that Q is a function of the variables  $(x_1, x_2, R, \Phi)$ , namely

$$Q(x_1, x_2, u, v) = \mathfrak{Q}(x_1, x_2, R(x_1, x_2, u, v), \Phi(x_1, x_2, u, v)).$$

In view of the definition of Q and knowing that  $X\Phi \equiv 0$  and  $Y\Phi \equiv 0$ , the equation (2.4.47) reduces to  $X\mathfrak{Q} = 0$ , meaning that  $\mathfrak{Q} \equiv \mathcal{Q}(Rx_1 - x_2, R, \Phi)$  for some function  $\mathcal{Q}(T, R, \Phi) \in \mathcal{C}^1(\mathbb{R}^3; \mathbb{R})$ .

The conditions (2.4.47) and (2.4.48) become the two scalar conservation laws

(2.4.49) 
$$\partial_v R + R \,\partial_u R - \mathcal{Q}(R \, x_1 - x_2, R, \Phi) \,\partial_1 R \\ - \mathcal{Q}(R \, x_1 - x_2, R, \Phi) \,R \,\partial_2 R \equiv 0,$$

(2.4.50) 
$$\begin{aligned} \partial_u R &- \mathcal{Q}(R \, x_1 - x_2, R, \Phi) \ \partial_2 R \\ &+ (x_1 \, \partial_T \mathcal{Q} + \partial_R \mathcal{Q})(R \, x_1 - x_2, R, \Phi) \ (\partial_1 R + R \, \partial_2 R) \equiv 0. \end{aligned}$$

Consider the equation (2.4.50) written for  $R_0(x_1, x_2, u)$  and associated with the initial data  $R_0(x_1, x_2, 0) = R_{00}(x_1, x_2)$ . At first sight, the access to  $R_0$  (and R) requires the knowledge of  $\Phi_0(x_1, x_2, u) := \Phi(x_1, x_2, u, 0)$  (and  $\Phi$ ). Nevertheless, by construction, the function  $\Phi$  is constant along the characteristics associated with (2.4.49) and (2.4.50). Thus, in doing so, it suffices to know who is  $\Phi_{00}(x_1, x_2) := \Phi_0(x_1, x_2, 0)$ .

Look at (2.4.49) as an evolution equation in v associated with the initial data  $R_0$ . For the same reasons as above, we can solve this Cauchy problem knowing only who is  $\Phi_{00}$ . There is still a difficulty coming from a problem of compatibility between (2.4.49) and (2.4.50). We must check that the expression R thus obtained is still a solution of (2.4.50). To this end, it suffices to show that (2.4.50) is propagated (in the direction v). This is due to the identity

$$(Y - \alpha X) \left\{ -[Y; X] R + (XQ) (XR) \right\} = \frac{2 \left\{ -[Y; X] R + (XQ) (XR) \right\}^2}{XR}$$

Note that  $Y R \neq 0$  as a consequence of (2.4.49) and  $Q \neq 0$ . The function  $\Phi$  can be obtained by the same procedure, by first integrating (2.4.50) and then by looking at (2.4.49). Geometrically, we have

$$\nabla \Phi = \lambda^{t} (-RXR, XR, YR, -RYR), \qquad \lambda \neq 0$$

implying that the level surfaces of  $\Phi$  intersect the plane  $\{u = v = 0\} \subset \mathbb{R}^4$  transversally. Thus, there is a unique function  $\Phi$  satisfying (2.4.46).

## 2.4.3 Reduction of the problem : analytical step.

In the preceding paragraph 2.4.2, we have developped only the aspects of R related to (2.4.38) or (2.4.39) and (2.4.40). However, the coefficient  $R(x_1, x_2, u, v)$  is also linked through the implicit relation (2.4.34) to the selection of a function u(x, v) satisfying (2.4.26) and (2.4.27).

At the level of (2.4.26) and (2.4.27), the variable v plays the part of a parameter. The situation is the same as in the Lemma 2.3.2. It suffices to select some data  $u_{00}(x_3, v) \equiv u(0, 0, x_3, v)$  such that  $\partial_3 u_{00} \neq 0$  in order to obtain (locally in  $\mathbb{R}^4$ ) some solution u of (2.4.26) and (2.4.27) satisfying  $\partial_3 u \neq 0$ .

**Proposition 2.4.3.** Let u(x, v) be any (local) solution of (2.4.26) and (2.4.27) satisfying  $\partial_3 u \neq 0$ . Define  $R(x_1, x_2, u, v)$  through (2.4.34). Then, there is a function  $\mathfrak{K} \in \mathcal{C}^1(\mathbb{R}^2; \mathbb{R})$  such that R can be put in the form

(2.4.51) 
$$R(x_1, x_2, u, v) = -\partial_v \alpha(x_1, x_2, u, v) / \partial_u \alpha(x_1, x_2, u, v)$$

where the scalar function  $\alpha$  is given by

(2.4.52) 
$$\alpha(x_1, x_2, u, v) := \mathfrak{K}(u, v) + \partial_v \mathfrak{L}(u, v) x_1 + \partial_u \mathfrak{L}(u, v) x_2.$$

In this context, the two restrictions  $R \neq 0$  and  $S \neq 0$  which are prerequisites in the analysis, see after (2.4.35), become  $\partial_u \mathfrak{K} \neq 0$  and  $\partial_v \mathfrak{K} \neq 0$ .

Proof of Proposition 2.4.3. From (2.4.26) and (2.4.27), it is easy to deduce

$$(2.4.53) \qquad \qquad \partial_1(\partial_v u) + \partial_v \mathfrak{L} \ \partial_3(\partial_v u) + (\partial_{uv}^2 \mathfrak{L} + \partial_{uv}^2 \mathfrak{L} \ \partial_v u) \ \partial_3 u = 0.$$

(2.4.54) 
$$\partial_1(\partial_3 u) + \partial_v \mathfrak{L} \ \partial_3(\partial_3 u) + \partial_{uv}^2 \mathfrak{L} \ (\partial_3 u)^2 = 0$$

(2.4.55) 
$$\partial_2(\partial_v u) + \partial_u \mathfrak{L} \, \partial_3(\partial_v u) + (\partial_{uv}^2 \mathfrak{L} + \partial_{uu}^2 \mathfrak{L} \, \partial_v u) \, \partial_3 u = 0,$$

(2.4.56) 
$$\partial_2(\partial_3 u) + \partial_u \mathfrak{L} \ \partial_3(\partial_3 u) + \partial_u^2 \mathfrak{L} \ (\partial_3 u)^2 = 0.$$

Since  $\partial_3 u \neq 0$ , these equations can be interpreted in the variables  $x_1, x_2, u$  and v. Then, it remains the following ODEs (with respect to  $x_1$  and  $x_2$ ):

(2.4.57) 
$$\partial_1(R/S) = -\partial_{vv}^2 \mathfrak{L}, \qquad \partial_1(1/S) = \partial_{uv}^2 \mathfrak{L},$$

(2.4.58) 
$$\partial_2(R/S) = -\partial_{uv}^2 \mathfrak{L}, \qquad \partial_2(1/S) = \partial_{uu}^2 \mathfrak{L}.$$

Observe that u and v play the part of parameters. It is easy to integrate (2.4.57) and (2.4.58). There are functions k(u, v) and h(u, v) such that

(2.4.59) 
$$R/S = k(u,v) - \partial_{vv}^2 \mathfrak{L}(u,v) x_1 - \partial_{uv}^2 \mathfrak{L}(u,v) x_2,$$

(2.4.60) 
$$1/S = h(u,v) + \partial_{uv}^2 \mathfrak{L}(u,v) \ x_1 + \partial_{uu}^2 \mathfrak{L}(u,v) \ x_2 .$$

In fact, the two functions k and h are linked together. This is due to the equality of the mixed partials derivatives  $\partial_v(\partial_3 u)$  and  $\partial_3(\partial_v u)$ :

$$\partial_v [\partial_3 u(x,v)] = \partial_v [S(x_1, x_2, u, v)] = \partial_u S R + \partial_v S = \partial_3 [\partial_v u(x,v)] = \partial_3 [R(x_1, x_2, u, v)] = \partial_u R S.$$

In other words, we must have

 $(R \ \partial_u S - S \ \partial_u R)/S^2 = -\partial_u (R/S) = -\partial_v S/S^2 = \partial_v (1/S).$ 

Apply this at the level of (2.4.59) and (2.4.60) to obtain  $-\partial_u k = \partial_v h$ . There is  $\mathfrak{K}(u,v)$  such that  $k = -\partial_v \mathfrak{K}$  and  $h = \partial_u \mathfrak{K}$ . Dividing (2.4.59) by (2.4.60) and replacing k and h as indicated previously, we get (2.4.51) and (2.4.52).

The explicit formulas (2.4.51) and (2.4.52) indicate that  $R = -\partial_1 \beta / \partial_2 \beta$  with

$$\beta(x_1, x_2, u, v) := \partial_v \mathfrak{K} x_1 + \partial_u \mathfrak{K} x_2 + \frac{1}{2} \ \partial_{vv}^2 \mathfrak{L} \ x_1^2 + \partial_{uv}^2 \mathfrak{L} \ x_1 \ x_2 + \frac{1}{2} \ \partial_{uu}^2 \mathfrak{L} \ x_2^2 \,.$$

Combining the informations obtained in this paragraph 2.4.3 with (2.4.36) and (2.4.37), we can observe that

(2.4.61) 
$$R = -\frac{\partial_1 \Phi}{\partial_2 \Phi} = -\frac{\partial_1 \beta}{\partial_2 \beta}, \qquad R = -\frac{\partial_v \alpha}{\partial_u \alpha} = -\frac{\partial_v \Phi}{\partial_u \Phi}.$$

Now, we can produce another interpretation of the intermediate problem under study which is emphasized in the introduction.

**Remark 2.4.3.1.** The question is to know if we can find two functions  $\mathfrak{K}(u, v)$ and  $\mathfrak{L}(u, v)$  allowing a simultaneous factorization of some  $\Phi$  in the form

$$\Phi = \mathcal{A}(x_1, x_2, lpha(x_1, x_2, u, v)) = \mathcal{B}(u, v, eta(x_1, x_2, u, v)), \qquad 
abla \Phi 
ot \equiv 0.$$

Since  $\nabla \Phi \neq 0$ , the two functions  $\mathcal{A}$  and  $\mathcal{B}$  cannot be constant. This is the source of the difficulty.

## 2.4.4 Test of the integrability conditions.

At this stage, we have to plug the coefficient R given by (2.4.51)-(2.4.52) into the integrability conditions (2.4.38) or (2.4.39)-(2.4.40). In this procedure, we have a little freedom coming from the choice of  $\mathfrak{L}$  and  $\mathfrak{K}$ . The matter is to check that the related constraints on  $\mathfrak{L}$  and  $\mathfrak{K}$  can indeed be realized for non trivials choices of  $\mathfrak{L}$  and  $\mathfrak{K}$ .

In the paragraph 2.4.4.1, we examine the case  $\dim \mathcal{A} = 2$ , that is (2.4.38). Then, in the paragraph 2.4.4.2, we consider the case  $\dim \mathcal{A} = 3$ , that is (2.4.39)-(2.4.40).

#### 2.4.4.1 The two-dimensional criterion.

This is when  $\dim \mathcal{A} = 2$ . We have to deal with (2.4.38).

**Lemme 2.4.3.** A function R given by (2.4.51) with  $\alpha$  as in (2.4.52) satisfies (2.4.38) if and only if one of the two distinct following conditions is verified :

**i.1.** We have  $\partial_{vv}^2 \mathfrak{L} \equiv 0$ . The function  $\mathfrak{L}$  is linear, say  $\mathfrak{L}(u, v) = a + bu + cv$  with  $(a, b, c) \in \mathbb{R}^3$ . Moreover, we can find  $\mathfrak{R} \in \mathcal{C}^1(\mathbb{R}; \mathbb{R})$  such that

(2.4.62) 
$$R(u,v) = \Re(\Re(u,v)), \qquad \Re(\Re) \ \partial_u \Re + \partial_v \Re = 0.$$

**i.2.** We have  $\partial_{nn}^2 \mathfrak{L} \neq 0$ . We can find  $\mathfrak{H} \in \mathcal{C}^1(\mathbb{R};\mathbb{R})$  such that

(2.4.63) 
$$\partial_u \mathfrak{L}(u,v) = \mathfrak{H}(\partial_v \mathfrak{L}(u,v)), \quad \partial_u \mathfrak{K} - \mathfrak{H}'(\partial_v \mathfrak{L}) \ \partial_v \mathfrak{K} = 0.$$

In the first case (2.4.62), we are faced with a scalar conservation law. In the second case (2.4.63), we have to solve some Hamilton-Jacobi equation. In those cases, the determination of  $\mathfrak{K}$  and  $\mathfrak{L}$  can be achieved once two functions in  $\mathcal{C}^1(\mathbb{R};\mathbb{R})$  are given, namely  $\mathfrak{R}(\cdot)$  and  $\mathfrak{K}(u,0)$  or  $\mathfrak{H}(\cdot)$  and  $\mathfrak{K}(u,0)$ .

**Proof of Lemma 2.4.3.** The calculation of XR gives rise to a polynomial fraction in x. More precisely, we find  $XR = -(\partial_v \alpha)^{-3} P(x)$  with

$$P(x) = a_{(0,0)} + \Xi(\mathfrak{L}) \sum a_{\beta} x^{\beta}, \qquad \Xi(\mathfrak{L}) := \partial_{uu}^{2} \mathfrak{L} \partial_{vv}^{2} \mathfrak{L} - (\partial_{uv}^{2} \mathfrak{L})^{2}.$$

The sum runs over all multi-indices  $\beta \in \mathbb{N}^2$  such that  $1 \leq |\beta| \leq 2$ . We find

$$\begin{split} a_{(0,0)} &= (\partial_u \mathfrak{K})^2 \ \partial_{vv}^2 \mathfrak{L} - 2 \ \partial_u \mathfrak{K} \ \partial_v \mathfrak{K} \ \partial_{uv}^2 \mathfrak{L} + (\partial_v \mathfrak{K})^2 \ \partial_{uu}^2 \mathfrak{L} \,, \quad a_{(1,0)} &= 2 \ \partial_v \mathfrak{K} \,, \\ a_{(0,1)} &= 2 \ \partial_u \mathfrak{K} \,, \qquad a_{(2,0)} &= \partial_{vv}^2 \mathfrak{L} \,, \qquad a_{(1,1)} &= 2 \ \partial_{uv}^2 \mathfrak{L} \,, \qquad a_{(0,2)} &= \partial_{uu}^2 \mathfrak{L} \,. \end{split}$$

Suppose that  $\Xi(\mathfrak{L}) \neq 0$ . Then, the condition  $XR \equiv 0$  requires that all the coefficients  $a_{\beta}$  with  $|\beta| \leq 2$  are equal to zero. In particular, it follows that  $\partial_u \mathfrak{K} \equiv 0$  and  $\partial_v \mathfrak{K} \equiv 0$ . This is not possible because this situation was excluded. Necessarily, we must impose  $\Xi(\mathfrak{L}) \equiv 0$ .

**i.1.** When  $\partial_{vv}^2 \mathfrak{L} \equiv 0$ , the condition  $\Xi(\mathfrak{L}) \equiv 0$  becomes  $\partial_{uv}^2 \mathfrak{L} \equiv 0$ . It remains  $a_{(0,0)} = (\partial_v \mathfrak{K})^2 \partial_{uu}^2 \mathfrak{L} \equiv 0$ . The function  $\mathfrak{L}$  must be linear in u and v, say  $\mathfrak{L}(u, v) = 0$ .

 $a+b\,u+c\,v$ . It follows that  $R \equiv -\partial_v \mathfrak{K}/\partial_u \mathfrak{K}$ . The other constraint  $YR \equiv 0$  amounts to the same thing as

 $\partial_u \mathfrak{K} (R \ \partial_u R + \partial_v R) \equiv -\partial_v \mathfrak{K} \ \partial_u R + \partial_u \mathfrak{K} \ \partial_v R \equiv 0$ implying that  $R \equiv \mathfrak{R}(\mathfrak{K})$  for some  $\mathfrak{R} \in \mathcal{C}^1(\mathbb{R}; \mathbb{R})$ . We have (2.4.62).

**i.2.** When  $\partial_{vv}^2 \mathfrak{L} \neq 0$ , the relation  $\Xi(\mathfrak{L}) \equiv 0$  is equivalent to  $\partial_u \mathfrak{L} = \mathfrak{H}(\partial_v \mathfrak{L})$ for some  $\mathfrak{H} \in \mathcal{C}^1(\mathbb{R};\mathbb{R})$ . Then, the condition  $a_{(0,0)} \equiv 0$  leads to the condition  $\partial_{vv}^2 \mathfrak{L} \left[ \partial_u \mathfrak{K} - \mathfrak{H}'(\partial_v \mathfrak{L}) \partial_v \mathfrak{K} \right]^2 \equiv 0$ . We recognize here the second part of (2.4.63). We find  $R \equiv -\mathfrak{H}'(\partial_v \mathfrak{L})^{-1}$  and, combining the preceding informations, it becomes easy to check that the relation  $YR \equiv 0$  is sure to be satisfied.

#### 2.4.4.2 The three-dimensional criterion.

This is when  $\dim \mathcal{A} = 3$ . We have to deal with (2.4.39) and (2.4.40), knowing that  $XR \neq 0$  or  $YR \neq 0$ . We consider separately the different situations which can happen concerning XR or YR.

**Lemme 2.4.4.** [Case  $XR \equiv 0$  and  $YR \not\equiv 0$ ]. A function R given by (2.4.51) with  $\alpha$  as in (2.4.52) satisfies  $XR \equiv 0$ , (2.4.39) and (2.4.40) without (2.4.38) if and only if the function  $\mathfrak{L}$  is linear, say  $\mathfrak{L}(u, v) = a + bu + cv$  with  $(a, b, c) \in \mathbb{R}^3$ , whereas  $R \equiv -\partial_v \mathfrak{K} / \partial_u \mathfrak{K}$  with  $\mathfrak{K}$  such that

$$(2.4.64) \qquad \qquad (\partial_v \mathfrak{K})^2 \ \partial_{uu}^2 \mathfrak{K} - 2 \ \partial_u \mathfrak{K} \ \partial_v \mathfrak{K} \ \partial_{uv}^2 \mathfrak{K} + (\partial_u \mathfrak{K})^2 \ \partial_{vv}^2 \mathfrak{K} \neq 0.$$

**Proof of Lemma 2.4.4.** The discussion is the same as in the proof of Lemma 2.4.3. The option **i.2** must be excluded because it leads to  $YR \equiv 0$ . Just go back to **i.1** where (2.4.62) must be exchanged with (2.4.64).

**Lemme 2.4.5.** [Case  $XR \neq 0$  and  $YR \equiv 0$ ]. A function R given by (2.4.51) with  $\alpha$  as in (2.4.52) satisfies  $YR \equiv 0$ , (2.4.39) and (2.4.40) without (2.4.38) if and only if the function  $\Re$  is linear in u and v, say  $\Re(u, v) = \alpha u + \beta v + \gamma$  with  $\alpha \neq 0$  and  $\beta \neq 0$ , whereas the function  $\mathfrak{L}(u, v)$  is polynomial in u and v with degree less or equal to 2. Moreover, the involved coefficients must be adjusted in order to have  $XR \neq 0$ .

**Proof of Lemma 2.4.5.** We have (2.4.40) and the condition (2.4.39) reduces to  $YXR \equiv 0$  yielding  $[X;Y]R = XYR - YXR \equiv 0 \equiv XR\partial_u R$ . It means that  $\partial_u R \equiv 0$  and therefore  $\partial_v R \equiv 0$ . The function R does not depend on (u, v). In particular, for  $(x_1, x_2) = (0, 0)$ , we find that  $\partial_v \Re/\partial_u \Re$  is constant. Since  $\partial_u \Re \neq 0$ and  $\partial_v \Re \neq 0$ , we must have

$$\mathfrak{K}(u,v) \equiv K(u-a\,v)\,, \qquad a \in \mathbb{R}^*\,, \qquad K \in \mathcal{C}^2(\mathbb{R};\mathbb{R})\,, \qquad K' \not\equiv 0\,.$$

Either  $K'' \equiv 0$  and all derivatives  $D^{\beta}\mathfrak{L}$  with  $|\beta| = 2$  are constant, leading to the description of Lemma 2.4.5. Or  $K'' \not\equiv 0$  and  $\mathfrak{L}(u, v) = F(u - av) + \beta v$  for some function  $F \in \mathcal{C}^{2}(\mathbb{R}; \mathbb{R})$  and some constant  $\beta \in \mathbb{R}$ . Nevertheless, this last case must be excluded. Indeed, it yields  $R \equiv -a$  so that  $XR \equiv 0$  (in contradiction with the hypothesis  $XR \not\equiv 0$ ).

The remaining case is when  $XR \neq 0$  and  $YR \neq 0$ .

**Proposition 2.4.4.** [Case  $XR \neq 0$  and  $YR \neq 0$ ]. A function R which is such that  $(XR)(YR) \neq 0$  and which is given by (2.4.51) with  $\alpha$  as in (2.4.52) satisfies (2.4.39) and (2.4.40) when the expressions  $\Re$  and  $\mathfrak{L}$  are adjusted according to one of the two following (distinct) situations :

**ii.1.** Both functions  $\mathfrak{L}(u, v)$  and  $\mathfrak{K}(u, v)$  are polynomial in u and v with degree less or equal to 2. More precisely, we have

(2.4.65)  $\mathfrak{L}(u,v) = a_{20} u^2 + 2 a_{11} u v + a_{02} v^2 + a_1 u + a_2 v + a_0, \quad a_{\star} \in \mathbb{R},$ 

(2.4.66) 
$$\mathfrak{K}(u,v) = k_{20} u^2 + 2 k_{11} u v + k_{02} v^2 + k_1 u + k_2 v + k_0, \quad k_\star \in \mathbb{R},$$

with coefficients  $a_{20}$ ,  $a_{11}$  and  $a_{02}$  (not all equal to zero) and coefficients  $k_{20}$ ,  $k_{11}$ and  $k_{02}$  (not all equal to zero) adjusted such that

$$(2.4.67) k_{11} a_{02} - k_{02} a_{11} = k_{20} a_{02} - k_{02} a_{20} = k_{20} a_{11} - k_{11} a_{20} = 0.$$

**ii.2.** The functions  $\mathfrak{L}(u, v)$  can be put in the form

(2.4.68) 
$$\mathfrak{L}(u,v) = a \ u + \mathbb{F}(b \ u + v), \qquad (a,b) \in \mathbb{R}^2$$

where the auxiliary function  $\mathbb{F} \in \mathcal{C}^3(\mathbb{R}; \mathbb{R})$  satisfies  $\mathbb{F}^{(3)} \neq 0$  and the ODE

(2.4.69) 
$$(\gamma s^2 + 2\beta s + \delta) \mathbb{F}^{(3)}(s) + 3(\gamma s + \beta) \mathbb{F}^{(2)}(s) = 0, \qquad s \in \mathbb{R}$$

with constants  $\gamma$ ,  $\beta$  and  $\delta$  not all equal to zero. The gradient of  $\mathfrak{K}(u, v)$  is adjusted as indicated at the level of (2.4.75) (with polynomial functions A and B which are defined in the proof).

**ii.3.** The function  $\mathfrak{L}(u, v)$  can be put in the form

(2.4.70) 
$$\mathfrak{L}(u,v) = u \mathbb{F}(u^{-1}(v+\alpha)) + \mathbb{G}(u), \qquad \alpha \in \mathbb{R}$$

where the auxiliary functions  $\mathbb{F} \in \mathcal{C}^2(\mathbb{R};\mathbb{R})$  and  $\mathbb{G} \in \mathcal{C}^2(\mathbb{R};\mathbb{R})$  satisfy

(2.4.71) 
$$\mathbb{F}^{(2)}(u) \neq 0, \qquad \delta \mathbb{F}^{(2)}(u) = u^3 \mathbb{G}^{(2)}(u), \qquad u \in \mathbb{R}$$

with  $\delta \in \mathbb{R}^*$ . Moreover  $\mathfrak{K}(u, v) = \partial_v \mathfrak{L}(u, v)$ .

**Proof of Proposition 2.4.4.** Below, we check that the different choices described in the paragraphs **ii.1**, **ii.2** and **ii.3** are convenient. Showing that there are no other possible situations is delicate. This aspect of the discussion is postponed to the Appendix 2.6. Recall that (2.4.39)-(2.4.40) is equivalent to (2.4.47)-(2.4.48) or to (2.4.49)-(2.4.50). We start by looking at the equation (2.4.49) which is the same as  $Z R \equiv 0$  where Z is the vector field

 $Z := Y - \mathcal{Q}(R x_1 - x_2, R, \Phi) X, \qquad X = \partial_1 + R \ \partial_2, \qquad Y = R \ \partial_u + \partial_v.$ 

By construction, we have also

(2.4.72) 
$$Z(Rx_1 - x_2) \equiv 0, \quad Z\Phi \equiv 0, \quad Z[Q(Rx_1 - x_2, R, \Phi)] \equiv 0.$$

Select  $\tilde{v} \in \mathbb{R}$  near 0. Given  $f \in \mathcal{C}_l^1(\mathbb{R}^4; \mathbb{R})$ , note  $f_{\tilde{v}}(x_1, x_2, u) := f(x_1, x_2, u, \tilde{v})$ . For instance, we have  $R_{\tilde{v}}(x_1, x_2, u) := R(x_1, x_2, u, \tilde{v})$  and

$$Q_{\tilde{v}}(x_1, x_2, u) := Q(x_1, x_2, u, \tilde{v}) = \mathcal{Q}(R_{\tilde{v}} x_1 - x_2, R_{\tilde{v}}, \Phi_{\tilde{v}}).$$

We also adopt the following conventions

$$\begin{aligned} d_v f &:= \partial_v \big[ f(x_1, x_2, u + R_{\tilde{v}} v, \tilde{v} + v) \big] \\ &= (R_{\tilde{v}} \, \partial_u f + \partial_v f)(x_1, x_2, u + R_{\tilde{v}} v, \tilde{v} + v) \,, \\ d_v^2 f &:= \partial_v (d_v f) \,= \, (R_{\tilde{v}}^2 \, \partial_{uu}^2 f + 2 \, R_{\tilde{v}} \, \partial_{uv}^2 f + \partial_{vv}^2 f)(x_1, x_2, u + R_{\tilde{v}} v, \tilde{v} + v) \,. \end{aligned}$$

To avoid confusions, retain that, in general, we have  $d_v^2 \not\equiv d_v \circ d_v$ . In view of (2.4.72), the characteristic associated with (2.4.49) and starting from the point  $(x_1, x_2, u, \tilde{v})$  is a straight line given by

$$(2.4.73) (X_1, X_2, U, V)(v) = (x_1 - Q_{\tilde{v}} v, x_2 - Q_{\tilde{v}} R_{\tilde{v}} v, u + R_{\tilde{v}} v, \tilde{v} + v).$$

The function R must be constant along the characteristics. Expressing this principle in connection with the definitions (2.4.51)-(2.4.52) yields

$$(2.4.74) d_v \mathfrak{K} + d_v (\partial_v \mathfrak{L}) x_1 + d_v (\partial_u \mathfrak{L}) x_2 - Q_{\tilde{v}} v d_v^2 \mathfrak{L} \equiv 0.$$

• The situation ii.1. Observe that, due to (2.4.65), the three quantities  $d_v(\partial_v \mathfrak{L})$ ,  $d_v(\partial_u \mathfrak{L})$  and  $d_v^2 \mathfrak{L}$  are constant functions. Thus, applying the second order derivative  $\partial_{vv}^2$  to the identity (2.4.74), we can extract

 $\partial^3_{uuu} \mathfrak{K}(U, V) R^3_{\tilde{v}} + 3 \partial^3_{uuv} \mathfrak{K}(U, V) R^2_{\tilde{v}} + 3 \partial^3_{uvv} \mathfrak{K}(U, V) R_{\tilde{v}} + \partial^3_{vvv} \mathfrak{K}(U, V) = 0.$ Since the three variables  $R_{\tilde{v}}$ , U and V are independent, we must have (2.4.66). Then, observe that

$$\partial_u R = -2 (\partial_u \alpha)^{-1} (k_{11} + k_{20} R), \qquad \partial_v R = -2 (\partial_u \alpha)^{-1} (k_{02} + k_{11} R), \partial_1 R = -2 (\partial_u \alpha)^{-1} (a_{02} + a_{11} R), \qquad \partial_2 R = -2 (\partial_u \alpha)^{-1} (a_{11} + a_{20} R).$$

It follows that

$$Q(x_1, x_2, u, v) \equiv Q(R) = \frac{YR}{XR} = \frac{k_{02} + 2k_{11}R + k_{20}R^2}{a_{02} + 2a_{11}R + a_{20}R^2}$$

We can work at the level of (2.4.47)-(2.4.48). By construction, the condition (2.4.47) is satisfied. On the other hand, (2.4.48) becomes

$$\partial_u R - \mathcal{Q}(R) \partial_2 R + \mathcal{Q}'(R) (\partial_1 R + R \partial_2 R) = 0.$$

This relation amounts to the same thing as

 $(k_{11}a_{02} - k_{02}a_{11}) + (k_{20}a_{02} - k_{02}a_{20}) R + (k_{20}a_{11} - k_{11}a_{20}) R^{2} = 0.$ 

This polynomial function of R is identically zero if and only if the restriction (2.4.67) is verified.

• The situation ii.2. Since  $\mathbb{F}^{(2)} \neq 0$ , we can introduce

(2.4.75) 
$$A(u,v) := \frac{\partial_u \mathfrak{K}(u,v)}{\mathbb{F}^{(2)}(b \ u+v)}, \qquad B(u,v) := \frac{\partial_v \mathfrak{K}(u,v)}{\mathbb{F}^{(2)}(b \ u+v)}.$$

With these conventions, the function R can be put in the form

$$R = -(B(u, v) + x_1 + b x_2) (A(u, v) + b x_1 + b^2 x_2)^{-1}$$

whereas

$$Q = \tilde{Q}(R, u, v) := \left(\partial_u A R^2 + (\partial_v A + \partial_u B) R + \partial_v B\right) (R b + 1)^{-2}.$$

The condition (2.4.47) reduces to  $\partial_v \hat{Q} + R \ \partial_u \hat{Q} = 0$ . Taking into account the above specific form of Q, we find a fraction in R whose coefficients must be zero. This criterion leads to

(2.4.76) 
$$\partial_{uu}^2 A = 2 \partial_{uv}^2 A + \partial_{uu}^2 B = \partial_{vv}^2 A + 2 \partial_{uv}^2 B = \partial_{vv}^2 B \equiv 0.$$

Exploiting (2.4.76), we can obtain

$$\begin{split} A(u,v) \, &=\, +\, \alpha\, u\, v - \gamma\, v^2 + a_0^1\, u + a_1^1\, v + a^0\,, \\ B(u,v) \, &=\, -\, \alpha\, u^2 + \gamma\, u\, v + b_0^1\, u + b_1^1\, v + b^0\,. \end{split}$$

Look at (2.4.48) which can also be formulated as  $\partial_u R - Q \ \partial_2 R + XQ = 0$ . Noting  $\mathfrak{D} := A + b \ x_1 + b^2 \ x_2$ , we find that

 $\mathfrak{D}(R b+1) XQ = 2 b \partial_v B - \partial_v A - \partial_u B + (b \partial_v A + b \partial_u B - 2 \partial_u A) R.$ 

Again, the condition (2.4.48) becomes a fraction in R whose coefficients must be zero. It follows that

$$(2.4.77) -3 \partial_u A + 2 b \partial_v A + b \partial_u B = 0,$$

$$(2.4.78) -2 \partial_u B + 3 b \partial_v B - \partial_v A = 0$$

It remains  $\alpha = -b\gamma$  and

$$A(u,v) = -b\gamma u v - \gamma v^{2} + b(-b_{0}^{1} + 2 b b_{1}^{1}) u + (-2 b_{0}^{1} + 3 b b_{1}^{1}) v + a^{0}.$$

Coming back to (2.4.75), we have to test the existence of  $\mathfrak{K}$  through Clairaut's Theorem. This is guaranteed by (2.4.69) if we choose  $\beta := b_0^1 - b b_1^1$  and  $\delta = b b^0 - a^0$ . The remaining restriction on  $\gamma$ ,  $\beta$  and  $\delta$  comes from the two conditions  $XR \neq 0$  and  $YR \neq 0$ .

• The situation ii.3. In this context, the definition of R gives rise to

(2.4.79) 
$$R = -\frac{R_1(x_1, x_2, u, v)}{R_2(x_1, x_2, u, v)} := -\frac{1 + x_1 + a(u, v) x_2}{a(u, v) (1 + x_1) + b(u, v) x_2},$$

where we have introduced

(2.4.80) 
$$a(u,v) := -u^{-1} (v+\alpha), \quad b(u,v) = u^{-2} [(v+\alpha)^2 + \delta].$$

We can use the formula given for R in (2.4.79) to compute

$$Q(x_1, x_2, u, v) \equiv \frac{YR}{XR} \equiv \frac{(YR_1)R_2 - (YR_2)R_1}{(XR_1)R_2 - (XR_2)R_1}$$

From (2.4.79), we can also extract

$$1 + x_1 = -\frac{h_1(u, v, R)}{h_2(u, v, R)} x_2 := -\frac{a(u, v) + R b(u, v)}{1 + R a(u, v)} x_2.$$

Then, replacing  $1 + x_1$  accordingly in the expression of Q, we can derive

$$Q(x_1, x_2, u, v) = \mathfrak{Q}(R, x_2, u, v) = \mathfrak{Q}_1(R, u, v) \ \mathfrak{Q}_2(R, u, v)^{-1} \ x_2$$

where  $\mathfrak{Q}_1$  and  $\mathfrak{Q}_2$  are only functions of R, u and v. We find

$$\begin{aligned} \mathfrak{Q}_1 &:= (b Y a - a Y b) h_2^2 + Y b h_1 h_2 - Y a h_1^2, \\ \mathfrak{Q}_2 &:= (b - a^2) h_2 (h_2 + R h_1). \end{aligned}$$

Many simplifications occur. It remains  $Q = -x_2 u^{-1}$  allowing to check (2.4.47) directly. On the other hand, the condition (2.4.48) reduces to

$$R_1 (u \partial_u + x_2 \partial_2) R_2 - R_2 (u \partial_u + x_2 \partial_2) R_1 + R_1 R_2 \equiv 0.$$

Taking into account the definitions of  $R_1$ ,  $R_2$ , a and b, this last relation becomes obvious to verify.

## 2.4.5 Discussion summary.

Up to now, we have described which conditions are needed in order to progress. Our aim here is to explain how to proceed concretely in order to build compatible couples  $(\varphi, w)$  in the case  $\nabla \varphi \cdot \partial_{\psi} \mathbf{W} \neq 0$ . Select two functions  $\mathfrak{L}$  and  $\mathfrak{K}$  as it is indicated in the paragraphs 2.4.4.1 or 2.4.4.2. In particular, we have  $\partial_u \mathfrak{K} \neq 0$ and  $\partial_v \mathfrak{K} \neq 0$ . Define the coefficient  $R(x_1, x_2, u, v)$  through (2.4.51) and (2.4.52). Knowing R, we have access to  $\Phi$ . More precisely, when  $\dim \mathcal{A} = 2$ , the function  $\Phi$  is entirely determined by prescribing

(2.4.81) 
$$\Phi_{00}(x_2, u) := \Phi(0, x_2, u, 0) \in \mathcal{C}^1(\mathbb{R}^2; \mathbb{R}), \qquad \nabla_{x_2, u} \Phi_{00} \neq 0.$$

On the other hand, when  $\dim \mathcal{A} = 3$  (implying that  $XR \neq 0$  and  $YR \neq 0$ ), the situation is more restricted. Then the function  $\Phi_{00}(x_2, u)$  must in fact depend only on one variable, say u. Indeed, it can be obtained by solving

(2.4.82) 
$$-YR \ \partial_2 \Phi_{00} + XR \ \partial_u \Phi_{00} = 0, \quad \Phi_{00}(0, u) = \Phi_{000}(u) \in \mathcal{C}^1(\mathbb{R}; \mathbb{R}).$$

Once the function  $\Phi_{00}(x_2, u)$  is fixed as it is indicated above, we can recover a non stationnary phase  $\Phi(x_1, x_2, u, v)$ . Now, select any function  $\chi \in \mathcal{C}^1(\mathbb{R}; \mathbb{R})$  such that  $\chi' \neq 0$  and consider the solution  $u_{00}(x_3, v)$  of the following ordinary differential equation (in the variable v)

(2.4.83) 
$$\partial_u \mathfrak{K}(u_{00}, v) \ \partial_v u_{00} + \partial_v \mathfrak{K}(u_{00}, v) = 0, \qquad u_{00}(x_3, 0) = \chi(x_3).$$

By construction, the expression  $u_{00}(x_3, v)$  satisfy (2.4.34). The resolution of the equations (2.4.26) and (2.4.27), where v belongs to some compact set  $K \subset \mathbb{R}$  and plays the part of a parameter, has already been discussed.

Given  $\mathfrak{L}$  and  $u_{00}$ , there is a unique expression  $u(x, v) \in \mathcal{C}^1(\Omega^0_r \times K; \mathbb{R})$  satisfying (2.4.26)-(2.4.27) together with the initial data

(2.4.84) 
$$u(0,0,x_3,v) = u_{00}(x_3,v) \in \mathcal{C}^1(\mathbb{R}^2;\mathbb{R}).$$

Moreover, as a consequence of the proof of Proposition 2.4.3, we have the relation (2.4.34) for all (x, v). Deducing the expression  $\varphi$  from  $\Phi(x_1, x_2, u, v)$  and u(x, v) through the formula (2.4.33), we obtain a function  $\varphi(x)$  which does not depend on the variable v.

The determination of the function  $v(x, \theta)$  is delicate. Combining (2.4.11) and (2.4.25), we can extract the functional identity

(2.4.85) 
$$v(x,\theta) = \mathbf{V}(\varphi(x), u(x,v(x,\theta)), \theta), \quad \forall (x,\theta) \in \Omega^0_r \times \mathbb{T}.$$

In particular, for x = 0 and for all  $\theta \in \mathbb{T}$ , we are faced with

$$v(0,\theta) = \mathbf{V}(\varphi(0), u_{00}(0, v(0,\theta)), \theta), \qquad \varphi(0) = \Phi_{00}(0, \chi(0)).$$

To simplify, we can seek a function  $v(x, \theta)$  such that  $v(0, \cdot) \equiv 0$ . It means that the function  $\mathbf{V}(\varphi, \psi, \theta)$  must be such that

(2.4.86) 
$$\mathbf{V}(\varphi(0), \chi(0), \theta) \equiv 0, \quad \forall \theta \in \mathbb{T}$$

In what follows, we select a function V satisfying (2.4.86). We suppose also that  $\partial_{\theta} \mathbf{V}$  is not the zero function and that

(2.4.87) 
$$\partial_{v}\mathfrak{K}(\chi(0),0) \ \partial_{\psi}\mathbf{V}(\varphi(0),\chi(0),\theta) + \partial_{u}\mathfrak{K}(\chi(0),0) \neq 0, \quad \forall \theta \in \mathbb{T}.$$

For each  $\theta \in \mathbb{T}$ , the informations (2.4.86) and (2.4.87) allow to apply the implicit Theorem at the point  $(0, \theta, 0)$  to the application

$$\begin{array}{cccc} \mathbb{R}^3 \times \mathbb{T} \times \mathbb{R} & \longrightarrow & \mathbb{R} \\ (x, \theta, v) & \longmapsto & v - \mathbf{V}\big(\varphi(x), u(x, v), \theta\big) \end{array}$$

It yields locally, near  $(0, \theta) \in \mathbb{R}^3 \times \mathbb{T}$ , a unique function  $v(x, \theta)$  satisfying the relation of (2.4.85). Due to the compactness of the torus  $\mathbb{T}$ , by adjusting the number  $r \in \mathbb{R}^+_*$ sufficiently small, we can recover (2.4.85). Note that the expression  $v(x, \theta)$  is (by construction) necessarily a solution of (2.4.28). Moreover, we do not have  $\partial_{\theta} v \equiv 0$ .

Define the function  $\psi(x,\theta)$  through (2.4.25). All the ingredients  $\varphi(x)$ ,  $\psi(x,\theta)$  and  $\mathbf{V}(\varphi,\psi,\theta)$  are determined. It means that the profile  $w(x,\theta)$  is known. Just use (2.2.13) and (2.4.8). By construction, the couple  $(\varphi,w)$  is compatible. Below, we sum up the preceding discussion by clearly precising the degrees of freedom at disposal in the construction of compatible couples  $(\varphi,w)$ .

**Proposition 2.4.5.** In the case  $\nabla \varphi \cdot \partial_{\psi} \mathbf{W} \neq 0$ , the class of compatible couples  $(\varphi, w)$  is entirely determined by giving locally

- functions  $\mathfrak{L}(u,v)$  and  $\mathfrak{K}(u,v)$  coming from the paragraphs 2.4.4.1 or 2.4.4.2;
- a function  $\Phi_{00}(x_2, u)$  which must satisfy (2.4.82) when dim  $\mathcal{A} = 3$ ;
- a function  $\chi(x_3)$ ;
- a function  $\mathbf{V}(\varphi, \psi, \theta)$  which is adjusted as in (2.4.86) and (2.4.87).

## 2.4.6 Illustrative examples.

The purpose here is to illustrate the various situations which can occur through corresponding examples. In practice, we select functions  $\mathfrak{L}$  and  $\mathfrak{K}$  resulting from the different cases classified in Section 2.4.4. In each case, we produce the corresponding phases  $\varphi(x)$ , and also the ingredients u and v allowing to recover the profile  $w(x, \theta)$  through (2.4.17) and (2.4.25).

To facilitate the presentation, we recall below the equations to deal with. Once  $\mathfrak{L}$  and  $\mathfrak{K}$  are fixed, the expression R is given by (2.4.51) and (2.4.52). By construction, there exist adequate functions  $\Phi$  such that

(2.4.88) 
$$\partial_1 \Phi + R \,\partial_2 \Phi \equiv 0, \qquad \partial_v \Phi + R \,\partial_u \Phi \equiv 0.$$

The function u must satisfy (2.4.26), (2.4.27) and (2.4.34), that is :

(2.4.89) 
$$\begin{cases} \partial_1 u + \partial_v \mathfrak{L}(u, v) \ \partial_3 u \equiv 0, \\ \partial_2 u + \partial_u \mathfrak{L}(u, v) \ \partial_3 u \equiv 0, \\ \partial_v u(x, v) = R(x_1, x_2, u(x, v), v) \end{cases}$$

The function  $v(x,\theta)$  is obtained through (2.4.28), that is

(2.4.90) 
$$\frac{\partial_1 v + \partial_v \mathfrak{L}(u(x,v),v) \ \partial_3 v}{+ \partial_v u(x,v) \ [\partial_2 v + \partial_u \mathfrak{L}(u(x,v),v) \ \partial_3 v] \equiv 0.$$

Then, it becomes possible to determine  $\varphi$  through (2.4.33). By construction, the function  $\varphi$  does not depend on  $\theta$  and it satisfies (2.4.29).

#### 2.4.6.1 Example in the case i.1 of Lemma 2.4.3.

By assumption, the function  $\mathfrak{L}$  is linear, say  $\mathfrak{L}(u,v) = a \ u + b \ v + c$  with  $(a,b,c) \in \mathbb{R}^3$ . The function  $R \equiv -\partial_v \mathfrak{K} / \partial_u \mathfrak{K}$  must be as indicated at the level of (2.4.62). To simplify, just take  $R \equiv 1$  so that  $\Phi = \Phi_{00}(x_2 - x_1, u - v)$ . From (2.4.89), we deduce that  $u(x,v) = \chi(x_3 - a \ x_2 - b \ x_1) + v$ . On the other hand, the function v can be written

$$v(x,\theta) = v_0 (x_1 - x_2, x_3 - (a+b) x_2, \theta), \qquad \partial_\theta v_0 \neq 0.$$
  
It remains to compute  $\varphi(x) = \Phi_{00} (x_2 - x_1, \chi(x_3 - a x_2 - b x_1)).$ 

#### 2.4.6.2 Example in the case i.2 of Lemma 2.4.3.

To simplify the discussion, we work with the choice  $\mathfrak{H}(t) = t$  implying that both  $\mathfrak{L}$  and  $\mathfrak{K}$  are functions of u + v. For instance, we have  $\mathfrak{L}(u, v) = L(u + v)$  for some function L satisfying  $L^{(2)} \neq 0$ . On the other hand, we find  $R \equiv -1$  and  $\Phi = \Phi_{00}(x_1 + x_2, u + v)$ . Looking at (2.4.89), we can infer that u(x, v) can be put in the form  $\tilde{u}(x_1 + x_2, x_3) - v$  where  $\tilde{u}(z, x_3)$  is obtained by solving the conservation law

 $\partial_z \tilde{u}(z, x_3) + L'(\tilde{u}(z, x_3)) \ \partial_3 \tilde{u}(z, x_3) = 0, \qquad \tilde{u}(0, x_3) = \chi(x_3).$ From (2.4.90), we deduce that  $v(x, \theta) = \tilde{v}(x_1 + x_2, x_3, \theta)$ . Observe also that  $\varphi(x) = \Phi_{00}(x_1 + x_2, \tilde{u}(x_1 + x_2, x_3)).$ 

#### 2.4.6.3 Example in the case of Lemma 2.4.4.

The function  $\mathfrak{L}$  is here linear, say  $\mathfrak{L}(u, v) = a \ u + b \ v + c$  with  $(a, b, c) \in \mathbb{R}^3$ . The function  $\mathfrak{K}$  must satisfy (2.4.64). We choose  $\mathfrak{K}(u, v) = -\frac{1}{2} \ v^2 + u$  in order to deal with  $R \equiv v$ . From (2.4.88), we can extract that  $\Phi = \Phi_{00}(2 \ u - v^2)$ . As expected, we see that  $\Phi$  depends this time on only one variable. Moreover

 $u(x,v) = 2^{-1} v^2 + \chi(x_3 - a x_2 - b x_1), \qquad \chi^{(1)} \neq 0.$ 

From (2.4.90), we obtain that

 $v(x,\theta) = \tilde{v}(b\,x_1 - x_3, a\,x_2 - x_3, \theta), \qquad \tilde{v}(y,z) \in \mathcal{C}^1(\mathbb{R}^2;\mathbb{R})$ 

where  $\tilde{v}(y, z)$  must satisfy the Burger's law  $b \partial_z \tilde{v} + a \tilde{v} \partial_y \tilde{v} \equiv 0$ . Finally :

 $\varphi(x) = (\Phi_{00} \circ \chi)(x_3 - a x_2 - b x_1).$ 

#### 2.4.6.4 Example in the case of Lemma 2.4.5.

The context is as in Lemma 2.4.5 with  $\Re = \alpha u + \beta v + \gamma$ . Choose  $\mathfrak{L} = 2^{-1} v^2$  so that  $R = -\alpha^{-1} (\beta + x_1)$  and  $\Phi \equiv \varphi = \Phi_{00} (x_2 + \alpha^{-1} \beta x_1 + (2\alpha)^{-1} x_1^2)$ . Note that  $u(x, v) = -\alpha^{-1} [(\beta + x_1) v - x_3] + c$  with  $c \in \mathbb{R}$ . On the other hand, the function v is obtained through

 $\partial_1 v + v \,\partial_3 v - \alpha^{-1} (x_1 + \beta) \,\partial_2 v \equiv 0.$ 

#### 2.4.6.5 Example in the case ii.1 of Proposition 2.4.4.

In agreement with (2.4.65) and (2.4.66), we can select  $\Re(u, v) = \mathfrak{L}(u, v) = v^2 + u$ so that  $R = -2(v + x_1)$  and  $\Phi \equiv \Phi_{00}(v^2 + u + x_1^2 + 2x_1v + x_2)$ . Moreover

 $u(x,v) = -2 v x_1 - x_2 + x_3 - v^2, \qquad \varphi(x) = \Phi_{00}(x_1^2 + x_3),$ 

whereas  $v(x, \theta)$  is any solution of

 $\partial_1 v - 2 (x_1 + v) \partial_2 v - 2 x_1 \partial_3 v = 0.$ 

#### 2.4.6.6 Example in the case ii.2 of Proposition 2.4.4.

Choose  $\mathfrak{K}(u,v) = uv^{-1}$  and  $\mathfrak{L}(u,v) = (2v)^{-1}$  so that  $R = uv^{-1} - x_1v^{-2}$ . We can take

$$\Phi(x_1, x_2, u, v) = \Phi_{00}\left(\frac{u}{v}x_1 - \frac{x_1^2}{2v^2} - x_2\right), \qquad u(x, v) = x_3 v + \frac{x_1}{2v} + \alpha v.$$

The function v is again solution of a suitable conservation law. On the other hand, we have again to deal with a phase  $\varphi$  which is some function of a quadratic expression in x, namely  $\varphi(x) = \Phi_{00}(x_1 x_3 + \alpha x_1 - x_2)$ .

#### 2.4.6.7 Example in the case ii.3 of Proposition 2.4.4.

In accordance with (2.4.71), select

$$\mathbb{F}(t) = \frac{t^2}{2}, \qquad \mathbb{G}(t) = \frac{1}{2t}, \qquad \delta = 1, \qquad \alpha = 0, \qquad \mathfrak{L}(u, v) = \frac{v^2 + 1}{2u}.$$

From (2.4.27), we can deduce the implicit relation

(2.4.91) 
$$u(x,v) = \tilde{U}(v x_1 - u(x,v) x_3, x_2, v), \quad \tilde{U}(X, x_2, v) \in C^1(\mathbb{R}^3; \mathbb{R})$$

From (2.4.26) together with (2.4.91), we can also derive

(2.4.92) 
$$\tilde{U}(X, x_2, v) = \underline{U}(2 X \tilde{U} - (v^2 + 1) x_2, v), \quad \underline{U}(Y, v) \in C^1(\mathbb{R}^2; \mathbb{R}).$$

Use (2.4.91) and (2.4.92) in order to extract respectively  $\partial_v u$  and  $\partial_X \tilde{U}$ . Replace  $x_3$  as indicated by a function of  $x_1$ ,  $x_2$ , Y, v and  $\underline{U}$ . By this way, we obtain a first expression for  $\partial_v u$ . It is compared below with the one coming directly from (2.4.51)-(2.4.52). We find :

$$R \equiv \partial_v u = \frac{\frac{\partial_v \underline{U}}{2 \underline{U} \partial_Y \underline{U}} + x_1 - \frac{v}{u} x_2}{\frac{v}{\underline{U}} \left[ \frac{\underline{U} - 2 Y \partial_Y \underline{U}}{2 v \underline{U} \partial_Y \underline{U}} + x_1 \right] - \frac{v^2 + 1}{u^2} x_2} = \frac{1 + x_1 - \frac{v}{u} x_2}{\frac{v}{\underline{U}} \left[ 1 + x_1 \right] - \frac{v^2 + 1}{u^2} x_2}.$$

It follows that  $\underline{U}(Y, v) = v \pm \sqrt{Y + v^2}$ . The function u(x, v) can now be deduced by just imposing u(0, 0) = 0 together with the implicit relation :

 $u(x,v) = \underline{U}(2 v x_1 u(x,v) - 2 x_3 u(x,v)^2 - (v^2 + 1) x_2, v).$ 

On the other hand, we seek  $\Phi(x_1, x_2, u, v)$  in the form

 $\Phi = \underline{\Phi} \left( u \left( 1 + x_1 \right) - v x_2, x_2, u, v \right), \qquad \underline{\Phi} (Y, x_2, u, v) \in \mathcal{C}^1(\mathbb{R}^4; \mathbb{R}).$ 

Taking into account the preceding definition of R, the condition (2.4.36) gives rise to the equation  $-x_2 \ \partial_Y \underline{\Phi} + Y \ \partial_2 \underline{\Phi} = 0$ . Thus, there is some function  $\underline{\Phi}_0(X, u, v) \in \mathcal{C}^1(\mathbb{R}^3; \mathbb{R})$  such that  $\underline{\Phi}(Y, x_2, u, v) = \underline{\Phi}_0(Y^2 + x_2^2, u, v)$ . From (2.4.37), we can then deduce that  $\partial_v \underline{\Phi}_0 \equiv 0$  and 2  $X \ \partial_X \underline{\Phi}_0 + u \ \partial_u \underline{\Phi}_0 = 0$ . In conclusion, the following choice is suitable :

$$\Phi(x_1, x_2, u, v) = \Phi_{00} \left( \frac{[u \ (1+x_1) - v \ x_2]^2 + x_2^2}{u^2} \right), \qquad \Phi_{00} \in \mathcal{C}^1(\mathbb{R}; \mathbb{R}).$$

With u(x, v) and  $\Phi(x_1, x_2, u, v)$  as above, we can deduce  $\varphi(x)$  through (2.4.33).

# 2.5 The time evolution problem.

Let  $(\varphi, w)$  be a compatible couple. We recall that the profile  $w(x, \theta)$  can be put in the form (2.2.13) with a triplet  $(\varphi, \psi, \mathbf{W})$  which is adjusted as it is indicated in (2.2.17)-(2.2.18)-(2.2.19)-(2.2.20) and which satisfies (2.2.15). The purpose of this last chapter is to explain what happens as time evolves.

## 2.5.1 Propagation of compatible datas.

The purpose of this paragraph 2.5.1 is to show the Theorem 5. Consider the system (2.1.11). Standard results (see for instance [25]) guarantee the existence locally in time, say on the domain  $\Omega_r^T \times \mathbb{T}$  with  $T \in \mathbb{R}^*_+$ , of a  $\mathcal{C}^1$ -solution to (2.1.11). Introduce  $\mathbf{U}(t, x, \theta) := \mathbf{W}(\Phi(t, x, \theta), \Psi(t, x, \theta), \theta)$ . From (2.1.11), we can easily deduce that

(2.5.1) 
$$\partial_t \mathbf{U} + (\mathbf{U} \cdot \nabla) \mathbf{U} = 0, \quad \mathbf{U}(0, x, \theta) = \mathbf{W}(\varphi(x), \psi(x, \theta), \theta) = w(x, \theta).$$

By integrating (2.1.11) along the associated characteristics (which are straight lines), we can exhibit the identities

(2.5.2) 
$$\Phi(t, x, \theta) = \varphi(x - t \mathbf{U}(t, x, \theta)) , \quad \forall (t, x, \theta) \in \Omega_r^T \times \mathbb{T},$$

(2.5.3) 
$$\Psi(t, x, \theta) = \psi(x - t \mathbf{U}(t, x, \theta), \theta), \quad \forall (t, x, \theta) \in \Omega_r^T \times \mathbb{T}.$$

**Lemme 2.5.1.** Assume that the three ingredients  $\varphi$ ,  $\psi$  and  $\mathbf{W}$  are adjusted according to (2.2.17)-(2.2.18)-(2.2.19)-(2.2.20). Then, the function  $\Phi(t, x, \theta)$  issued from (2.1.11) is such that  $\partial_{\theta} \Phi \equiv 0$ . Moreover, noting

$$y \equiv y(t,x) := x - t \mathbf{U}(t,x,\theta), \qquad \Xi(y,\theta) := \left(\varphi(y), \psi(y,\theta), \theta\right) \in \mathbb{R}^2 \times \mathbb{T},$$

the expression  $\Psi(t, x, \theta)$  coming from (2.1.11) satisfies

(2.5.4) 
$$\frac{\partial_{\theta} \Psi(t, x, \theta) \equiv \partial_{\theta} \psi(y, \theta)}{-t \nabla \psi(y, \theta) \cdot \left[\partial_{\theta} \mathbf{W}(\Xi(y, \theta)) + \partial_{\theta} \psi(y, \theta) \partial_{\psi} \mathbf{W}(\Xi(y, \theta))\right]}.$$

**Proof of the Lemma 2.5.1.** Use the relations (2.5.2) and (2.5.3) with the formula given for **U** to compute  $\partial_{\theta} \Phi$  and  $\partial_{\theta} \Psi$  according to

$$\mathcal{M} \left( \begin{array}{c} \partial_{\theta} \Phi(t, x, \theta) \\ \partial_{\theta} \Psi(t, x, \theta) \end{array} \right) = \left( \begin{array}{c} -t \, \nabla \varphi(y) \cdot \partial_{\theta} \mathbf{W} \big( \Xi(y, \theta) \big) \\ \partial_{\theta} \psi(y, \theta) - t \, \nabla \psi(y, \theta) \cdot \partial_{\theta} \mathbf{W} \big( \Xi(y, \theta) \big) \end{array} \right)$$

with a matrix  $\mathcal{M}$  given by

$$\mathcal{M}(t, y, \theta) := \begin{pmatrix} 1 + t \, \nabla \varphi(y) \cdot \partial_{\varphi} \mathbf{W} & t \, \nabla \varphi(y) \cdot \partial_{\psi} \mathbf{W} \\ t \, \nabla \psi(y, \theta) \cdot \partial_{\varphi} \mathbf{W} & 1 + t \, \nabla \psi(y, \theta) \cdot \partial_{\psi} \mathbf{W} \end{pmatrix}$$

In the preceding formula for the matrix  $\mathcal{M}$ , the functions  $\partial_{\star} \mathbf{W}$  are evaluated at the point  $\Xi(y,\theta)$ . A consequence of (2.2.19) and (2.2.20) is the information:  $\det \mathcal{M}(t, y, \theta) = 1$ . It follows that

$$\partial_{\theta} \Phi(t, x, \theta) = -t \nabla \varphi \cdot (\partial_{\theta} \mathbf{W} + \partial_{\theta} \psi \, \partial_{\psi} \mathbf{W}) + t^{2} \left[ (\nabla \psi \cdot \partial_{\theta} \mathbf{W}) (\nabla \varphi \cdot \partial_{\psi} \mathbf{W}) - (\nabla \varphi \cdot \partial_{\theta} \mathbf{W}) (\nabla \psi \cdot \partial_{\psi} \mathbf{W}) \right].$$

Note that the right hand term can be regarded as a function of  $(y, \theta)$ . The condition (2.2.17) is the same as

(2.5.5) 
$$\nabla \varphi(y) \cdot \left[ \partial_{\theta} \psi(y,\theta) \, \partial_{\psi} \mathbf{W} \big( \Xi(y,\theta) \big) + \partial_{\theta} \mathbf{W} \big( \Xi(y,\theta) \big) \right] = 0 \, .$$

Therefore, it remains

 $\partial_{\theta} \Phi(t, x, \theta) \,=\, t^2 \, \left( \nabla \varphi \cdot \partial_{\psi} \mathbf{W} \right) \, \left( \nabla \psi \cdot \partial_{\theta} \mathbf{W} + \partial_{\theta} \psi \, \nabla \psi \cdot \partial_{\psi} \mathbf{W} \right).$ 

Due to (2.2.18), this is  $\partial_{\theta} \Phi \equiv 0$ . Proceeding as above, we can obtain

$$\begin{aligned} \partial_{\theta}\Psi(t,x,\theta) &= \partial_{\theta}\psi + t \left(\partial_{\theta}\psi \ \nabla\varphi \cdot \partial_{\varphi}\mathbf{W} - \nabla\psi \cdot \partial_{\theta}\mathbf{W}\right) \\ &+ t^2 \left[ \left(\nabla\varphi \cdot \partial_{\theta}\mathbf{W}\right) \left(\nabla\psi \cdot \partial_{\varphi}\mathbf{W}\right) - \left(\nabla\varphi \cdot \partial_{\varphi}\mathbf{W}\right) \left(\nabla\psi \cdot \partial_{\theta}\mathbf{W}\right) \right]. \end{aligned}$$

Exploiting again (2.2.19), (2.2.20) and (2.5.5), this is equivalent to

$$\partial_{\theta}\Psi(t, x, \theta) = \partial_{\theta}\psi - t \left(\nabla\psi \cdot \partial_{\theta}\mathbf{W} + \partial_{\theta}\psi \nabla\psi \cdot \partial_{\psi}\mathbf{W}\right) - t^{2} \left(\nabla\varphi \cdot \partial_{\varphi}\mathbf{W}\right) \nabla\psi \cdot \left(\partial_{\theta}\mathbf{W} + \partial_{\theta}\psi \partial_{\psi}\mathbf{W}\right).$$

Due to (2.2.18) and (2.2.19), the term in factor of  $t^2$  is necessarily equal to zero. By this way, we can see how (2.5.4) appears.

Consider the expression  $u^{\varepsilon}$  which is defined on the domain  $\Omega_r^T$  through

(2.5.6) 
$$u^{\varepsilon}(t,x) := \mathbf{U}\left(t,x,\frac{\Phi(t,x)}{\varepsilon}\right) \\ = \mathbf{W}\left(\Phi(t,x),\Psi\left(t,x,\frac{\Phi(t,x)}{\varepsilon}\right),\frac{\Phi(t,x)}{\varepsilon}\right), \quad \varepsilon \in ]0,1].$$

By construction, we have  $u^{\varepsilon}(0, \cdot) \equiv h^{\varepsilon}(\cdot)$  with  $h^{\varepsilon}$  given by (2.1.2). A direct computation based on (2.1.11) indicates that  $u^{\varepsilon}(t, x)$  is indeed a solution of (2.1.1) on  $\Omega_r^T$ . By applying the Theorem 2.6 of [6], we obtain that  $(D_x u^{\varepsilon}(t, x))^3 \equiv 0$  on B(0, r - tV) for all  $t \in [0, T]$ . Repeating at the time  $t \in [0, T]$  the procedure of the Section 2.2, we can deduce that the constraints (2.2.17), (2.2.18), (2.2.19) and (2.2.20) are propagated. In other words :

**Lemme 2.5.2.** For all  $t \in [0,T]$ , the solutions  $\Phi(t,x)$  and  $\Psi(t,x,\theta)$  of (2.1.11) satisfy (2.1.13), (2.1.14), (2.1.15) and (2.1.16).

These identities can also be derived by using (2.2.17)-(2.2.18)-(2.2.19)-(2.2.20) as well as (2.5.2), (2.5.3) and the Lemma 2.5.1. The Theorem 5 is proved.

For the sake of completeness, we can also remark that the rank of the solution is a preserved quantity. In the case of rank 1, this is obvious. In the case of rank 2, this is a consequence of what follows.

**Lemme 2.5.3.** The solutions  $\Phi(t, x)$  and  $\Psi(t, x, \theta)$  of (2.1.11) satisfy

(2.5.7) 
$$(\nabla \Phi \wedge \nabla \Psi)(t, x, \theta) = (\nabla \varphi \wedge \nabla \psi)(y, \theta), \qquad \forall (t, x, \theta) \in \Omega_r^T \times \mathbb{T}.$$

Thus, the volume measure  $\nabla \Phi \wedge \nabla \Psi$  is constant along the characteristics. This is in fact a by-product of the divergence free relation.

**Proof of the Lemma 2.5.3.** By differentiating (2.5.2) and (2.5.3) with respect to  $x_i$ , we can extract

$$M(t, y, \theta) \left(\begin{array}{c} \partial_j \Phi(t, x, \theta) \\ \partial_j \Psi(t, x, \theta) \end{array}\right) = \left(\begin{array}{c} \partial_j \varphi(y) \\ \partial_j \psi(y) \end{array}\right), \qquad \forall j \in \{1, 2, 3\}.$$

It follows that

(2.5.8)  $\nabla \Phi = (1 + t \nabla \psi \cdot \partial_{\psi} \mathbf{W}) \nabla \varphi - t (\nabla \varphi \cdot \partial_{\psi} \mathbf{W}) \nabla \psi,$ 

(2.5.9) 
$$\nabla \Psi = -t \left( \nabla \psi \cdot \partial_{\varphi} \mathbf{W} \right) \nabla \varphi + \left( 1 + t \nabla \varphi \cdot \partial_{\varphi} \mathbf{W} \right) \nabla \psi.$$

Now, we can use (2.5.8) and (2.5.9) in order to compute the cross product of  $\nabla \Phi$  and  $\nabla \Psi$ . Due to (2.2.19) and (2.2.20), it remains (2.5.7).

### 2.5.2 Asymptotic phenomena.

Families like  $\{u^{\varepsilon}\}_{\varepsilon}$  give many informations on the complex phenomena which may occur at the level of (2.1.7) when passing to the limit (as  $\varepsilon \to 0$ ).

Noting  $\mathcal{S}(t)$  with  $t \in \mathbb{R}^*_+$  the semi-group operator which is associated with incompressible Euler equations, we can for instance use  $\{u^{\varepsilon}\}_{\varepsilon}$  to study the well-posedness (or not) of  $\mathcal{S}(t)$  in functional spaces (thus arising the delicate problem of the localization of the solutions, see [4]). We can also investigate the weak  $L^2$ -continuity (or not) of  $\mathcal{S}(t)$  (in the spirit of [5, 21]). These applications of our current approach will not be developed in these pages. Nevertheless, we will point out some related very specific aspect.

We want here to show that the phenomenon of superposition of oscillations already noted in [5] (only when d = 2 and in the absence of the divergence free constraint) can indeed occur at the level of (2.1.7) when d = 3.

The idea is to start at the initial time t = 0 with a function  $\psi(x)$  which does not see the variable  $\theta \in \mathbb{T}$  and with a function  $\mathbf{W}_{\varepsilon}(\cdot)$  which depends on the parameter  $\varepsilon \in ]0,1]$  and contains oscillations in the variable  $\psi$ , as it is indicated in (2.1.22). Then, in order to prove the mechanism (2.1.23), it suffices to exhibit some  $t \in ]0,T]$ such that  $\partial_{\theta} \Psi \neq 0$ .

In view of (2.5.1) and because  $\partial_{\theta} \psi \equiv 0$ , it suffices to test if

(2.5.10) 
$$\exists (x,\theta) \in \Omega^0_r \times \mathbb{T}; \qquad \partial_\theta \mathbf{W}(\varphi,\psi,\theta) \cdot \nabla \psi = \partial_\theta w \cdot \nabla \psi \neq 0.$$

In the framework  $\nabla \varphi \cdot \partial_{\psi} \mathbf{W} \neq 0$  of the Section 2.4.1, because of (2.4.2), it is not possible to obtain (2.5.10). When  $\nabla \varphi \cdot \partial_{\psi} \mathbf{W} \equiv 0$ , in the case  $f' \neq 0$  and  $g' \neq 0$ , we can see with (2.3.15), (2.3.23) and (2.3.24) that

$$\partial_{\theta}w \cdot \nabla \psi = \left\{ \partial_{\theta}\alpha \begin{pmatrix} 0 \\ -g \\ 1 \end{pmatrix} + \partial_{\theta}\beta \begin{pmatrix} 1 \\ -f \\ 0 \end{pmatrix} \right\} \cdot \left\{ \Psi_{0}' \begin{pmatrix} f'(\varphi) \\ 0 \\ g'(\varphi) \end{pmatrix} + \partial_{\varphi}\Psi \begin{pmatrix} \partial_{1}\varphi \\ \partial_{2}\varphi \\ \partial_{3}\varphi \end{pmatrix} \right\}.$$

Taking into account (2.3.12) and (2.3.18), we must necessarily have  $\partial_{\theta} w \cdot \nabla \psi \equiv 0$ . It remains to examine the situation of the paragraph 2.3.2.1. The context is the one of Proposition 2.3.1. Choose  $(a, b) \in (\mathbb{R}^*)^2$ , c = 0 and  $\varphi_{00} \in \mathcal{C}^1(\mathbb{R}; \mathbb{R})$ . Define  $\varphi(x)$  as in (2.3.19). Select some function  $\alpha_{\varepsilon}$  which is such that

$$\alpha_{\varepsilon}(\varphi,\psi,\theta) = A(\varphi,\psi,\psi/\varepsilon,\theta)\,, \qquad \partial_{\psi}A \neq 0\,, \qquad A \in \mathcal{C}^{\infty}(\mathbb{R}^2 \times \mathbb{T}^2;\mathbb{R})\,.$$

Choose two auxiliary functions  $\phi(\theta) \in \mathcal{C}^{\infty}(\mathbb{T};\mathbb{R})$  and  $\Psi_0(T,Z) \in \mathcal{C}^{\infty}(\mathbb{R}^2;\mathbb{R})$  satisfying  $\phi' \neq 0$  and  $\partial_T \psi_0 \neq 0$ . Take

$$\chi \equiv 1$$
,  $\beta_{\varepsilon}(\varphi, \psi, \theta) = \alpha_{\varepsilon}(\varphi, \psi, \theta) - \phi(\theta)$ ,  $\Psi(X, Y, Z) = \Psi_0(Y - X, Z)$ .

Obviously, we have (2.3.22) and (2.3.23). With  $w(x,\theta)$  defined according to (2.3.24), compute

$$\partial_{\theta} w \cdot \nabla \psi = \left\{ \partial_{\theta} \alpha \begin{pmatrix} 0 \\ -b \\ 1 \end{pmatrix} + \partial_{\theta} \beta \begin{pmatrix} 1 \\ -a \\ 0 \end{pmatrix} \right\} \cdot \left\{ \partial_{T} \Psi_{0} \begin{pmatrix} -1 \\ 0 \\ 1 \end{pmatrix} + \partial_{Z} \Psi_{0} \begin{pmatrix} a \\ 1 \\ b \end{pmatrix} \right\}$$
$$= \partial_{T} \Psi_{0} \phi' \neq 0.$$

In fact, the corresponding solution of (2.1.7) can be produced explicitly. It is

$$u^{\varepsilon}(t,x) = A_{\varepsilon}(t,x) \begin{pmatrix} 1\\ -a-b\\ 1 \end{pmatrix} - \phi\left(\frac{\varphi(x)}{\varepsilon}\right) \begin{pmatrix} 1\\ -a\\ 0 \end{pmatrix}$$

with

$$A_{\varepsilon}(t,x) := A\left(\varphi(x), \frac{\psi\left(x_3 - x_1 - t\,\phi\left(\frac{\varphi(x)}{\varepsilon}\right)\right)}{\varepsilon}, \frac{\varphi(x)}{\varepsilon}\right)$$

# 2.6 Appendix.

This appendix is concerned with the three-dimensional criterion which is studied at the level of Section 2.4.4.2. The matter is to consider the more complicated case, when  $XR \neq 0$  and  $YR \neq 0$ . The Proposition 2.4.4 gives sufficient conditions on  $\Re$ and  $\mathfrak{L}$  in order to solve the system (2.4.39)-(2.4.40). The aim of this Appendix is to explain (under suitable assumptions that will be precised later) why there are no other possible choices.

Thus, in all this Section 2.6, we deal with (2.4.39)-(2.4.40) or (2.4.47)-(2.4.48), in the case  $XR \neq 0$  and  $YR \neq 0$ . The starting point of our analysis is the equation (2.4.74). In a first approach, we assume that  $\partial^2_{vv} \mathfrak{L} \neq 0$ . We will see in the paragraph 2.6.5 that the case  $\partial^2_{vv} \mathfrak{L} \equiv 0$  can be dealt separately and that it does not produce other cases than (2.4.65).

# 2.6.1 Preliminary informations.

In what follows, we assume that  $\partial_{vv}^2 \mathfrak{L} \neq 0$ . In doing so, as seen below, no information is forgotten.

**Lemme 2.6.1.** The assumption  $\partial_{vv}^2 \mathfrak{L} \equiv 0$  is not compatible with solving the system (2.4.39)-(2.4.40) in the case  $(XR)(YR) \neq 0$ .

**Proof of Lemma 2.6.1.** First, consider the case when  $\partial_{vv}^2 \mathfrak{L} \equiv 0$  and also  $\partial_{uv}^2 \mathfrak{L} \not\equiv 0$ . To this end, introduce  $c(u, v) := \partial_{uu}^2 \mathfrak{L} / \partial_{uv}^2 \mathfrak{L}$ . The notations are as in (2.4.73). The starting point of the discussion is the identity (2.4.74) which can here be reduced to

$$\frac{d_v \mathfrak{K}}{\partial_{uv}^2 \mathfrak{L}} + R_{\tilde{v}} x_1 + \left[ c(U, V) R_{\tilde{v}} + 1 \right] x_2 - Q_{\tilde{v}} v \left[ c(U, V) R_{\tilde{v}}^2 + 2 R_{\tilde{v}} \right] \equiv 0.$$

Since  $\partial_{vv}^2 \mathfrak{L} \equiv 0$  and  $\partial_{uv}^2 \mathfrak{L} \not\equiv 0$ , we are sure that  $\partial_1 R_{\tilde{v}} \not\equiv 0$ . Thus, we can work with  $R_{\tilde{v}}, x_2, u, v, \tilde{v}$  instead of  $x_1, x_2, u, v, \tilde{v}$ . In particular, for  $R_{\tilde{v}} = 0$ , it remains

$$\partial_v \mathfrak{K}(u,V) \ \partial^2_{uv} \mathfrak{L}(u,V)^{-1} + x_2 \equiv 0, \qquad \forall (x_2,u,V), \qquad V = v + \tilde{v}.$$

This is not possible since the three variables  $x_2$ , u and V are independent. Thus, we have necessarily  $\partial^2_{uv} \mathfrak{L} \equiv 0$ . Knowing that  $\partial^2_{vv} \mathfrak{L} \equiv \partial^2_{uv} \mathfrak{L} \equiv 0$  and that  $XR \neq 0$ , the function  $\mathfrak{L}$  can be put in the form F(u) + bv for some constant  $b \in \mathbb{R}$  and some function  $F \in \mathcal{C}^2(\mathbb{R}; \mathbb{R})$  satisfying  $F^{(2)} \neq 0$ . This time, the identity (2.4.74) becomes

(2.6.1) 
$$d_v \mathfrak{K}(U,V) + x_2 F^{(2)}(U) R_{\tilde{v}} - v Q_{\tilde{v}} F^{(2)}(U) R_{\tilde{v}}^2 \equiv 0.$$

In particular, for  $R_{\tilde{v}} = 0$ , we find  $\partial_v \mathfrak{K} \equiv 0$ , that is  $\mathfrak{K}(u, v) \equiv \tilde{\mathfrak{K}}(u)$ . Then, dividing (2.6.1) by  $R_{\tilde{v}}$  and taking v = 0, we obtain  $\tilde{\mathfrak{K}}'(u) + x_2 F^{(2)}(u) \equiv 0$ . Since  $F^{(2)} \neq 0$ , this furnishes the expected contradiction.

From now on, assume that  $\partial_{uv}^2 \mathfrak{L} \neq 0$ . Introduce the two auxiliary functions

(2.6.2) 
$$a(u,v) := \partial_{uv}^2 \mathfrak{L} / \partial_{vv}^2 \mathfrak{L}, \qquad b(u,v) := \partial_{uu}^2 \mathfrak{L} / \partial_{vv}^2 \mathfrak{L}$$

From the informations (2.4.51)-(2.4.52) written with  $v = \tilde{v}$ , we obtain that

(2.6.3) 
$$\left[a(u,\tilde{v})R_{\tilde{v}}+1\right]\partial_2 R_{\tilde{v}} - \left[b(u,\tilde{v})R_{\tilde{v}}+a(u,\tilde{v})\right]\partial_1 R_{\tilde{v}} = 0.$$

Now, the idea is to manipulate (2.4.74) in order to eliminate the contribution

$$d_{v}\mathfrak{K} \equiv R_{\tilde{v}} \,\partial_{u}\mathfrak{K}(u+R_{\tilde{v}}\,v,\tilde{v}+v) + \partial_{v}\mathfrak{K}(u+R_{\tilde{v}}\,v,\tilde{v}+v) + \partial_{v}\mathfrak{K}(u+R_{v}\,v,\tilde{v}+v) + \partial_{v}\mathfrak{K}(u+$$

To this end, it suffices to apply the vector field  $\partial_2 R_{\tilde{v}} \partial_1 - \partial_1 R_{\tilde{v}} \partial_2$  to the equation (2.4.74). Then, use (2.6.3) in order to extract

(2.6.4) 
$$\begin{aligned} \Xi(u, v, \tilde{v}, R_{\tilde{v}}) &:= \left[a(u, \tilde{v}) - a(U, V)\right] + \left[b(u, \tilde{v}) - b(U, V)\right] R_{\tilde{v}} \\ &+ \left[a(U, V) b(u, \tilde{v}) - a(u, \tilde{v}) b(U, V)\right] R_{\tilde{v}}^2 \\ &- v \ \tilde{\chi}(x_1, x_2, u, \tilde{v}) \left[1 + 2 a(U, V) R_{\tilde{v}} + b(U, V) R_{\tilde{v}}^2\right] = 0 \end{aligned}$$

with

$$\tilde{\chi}(x_1, x_2, u, \tilde{v}) := \begin{bmatrix} b(u, \tilde{v}) R_{\tilde{v}} + a(u, \tilde{v}) \end{bmatrix} \partial_1 Q_{\tilde{v}} - \begin{bmatrix} a(u, \tilde{v}) R_{\tilde{v}} + 1 \end{bmatrix} \partial_2 Q_{\tilde{v}}$$

By definition, the function  $\tilde{\chi}$  does not depend on v. On the other hand, in view of (2.6.4), it is an expression of the four variables  $u, v, \tilde{v}$  and  $R_{\tilde{v}}$  (which are independent because  $XR \neq 0$ ). It means that

(2.6.5) 
$$\exists \chi \in \mathcal{C}^{\infty}(\mathbb{R}^3; \mathbb{R}); \qquad \tilde{\chi}(x_1, x_2, u, \tilde{v}) = \chi(u, \tilde{v}, R_{\tilde{v}}).$$

**Lemme 2.6.2.** There exist functions f, g, k and l in  $C^2(\mathbb{R}; \mathbb{R})$  such that

(2.6.6) 
$$a(u,v) = f(u)v + g(u), \quad b(u,v) = (2f^2 - f')(u)v^2 + k(u)v + l(u)$$

where the expression  $\mathbf{Z}(u) := 2 f(u)^2 - f'(u)$  must satisfy

(2.6.7) 
$$\mathbf{Z}'(u) = 2 f(u) \mathbf{Z}(u)$$

**Proof of Lemma 2.6.2.** As indicated line (2.6.5), we can replace  $\tilde{\chi}$  by  $\chi$  at the level of (2.6.4). Then, compute

 $\partial_{vv}^2 \Xi(u, v, \tilde{v}, 0) = -\partial_{vv}^2 a(u, \tilde{v} + v) = 0.$ 

Thus, we can find functions f and g such that a(u,v) = f(u)v + g(u). On the other hand  $\partial_v \Xi(u,v,\tilde{v},0) \equiv -\partial_v a(u,\tilde{v}+v) - \chi(u,\tilde{v},0) \equiv 0$  which implies that  $\chi(u,\tilde{v},0) = -f(u)$ . It follows that

(2.6.8)  
$$0 = \partial_{R_{\tilde{v}}} \Xi(u, v, \tilde{v}, 0) = b(u, \tilde{v}) - b(u, \tilde{v} + v) + v (\tilde{v} + v) [2 f(u)^2 - f'(u)] + v [2 f(u) g(u) - g'(u)] - v \partial_{R_{\tilde{v}}} \chi(u, \tilde{v}, 0).$$

In particular, we must have

 $\partial_{R_{\tilde{v}}vv}^3 \Xi(u, v, \tilde{v}, 0) = -\partial_{vv}^2 b(u, \tilde{v} + v) + 4 f(u)^2 - 2 f'(u) = 0.$ 

In other words, there are two functions k and l such that the second part of (2.6.6) is verified. Coming back to (2.6.8), we can see that

 $\partial_{R_{\tilde{v}}}\chi(u,\tilde{v},0) = -\mathbf{Z}(u) \,\tilde{v} - k(u) + 2 \,f(u) \,g(u) - g'(u) \,.$ 

Then look at the condition

(2.6.9) 
$$0 = (\partial_{R_{\tilde{v}}})^2 \Xi(u, v, \tilde{v}, 0) = -v^2 \partial_{uu}^2 a(u, V) - 2 v \partial_u b(u, V) + 2 \left[ a(u, V) b(u, \tilde{v}) - a(u, \tilde{v}) b(u, V) \right] - v (\partial_{R_{\tilde{v}}})^2 \chi(u, \tilde{v}, 0) - 4 v \partial_{R_{\tilde{v}}} \chi(u, \tilde{v}, 0) a(u, V) + 2 v f(u) \left[ 2 v \partial_u a(u, V) + b(u, V) \right].$$

Taking into account the preceding informations on a and b, then the expression  $(\partial_{R_{\tilde{v}}})^2 \Xi(u, v, \tilde{v}, 0)$  is a polynomial function with respect to v. In particular, the coefficient in factor of  $v^3$  must be zero. This criterion yields (2.6.7).

From (2.6.4), we can extract a formula for  $\chi$ , namely

$$\chi(u,\tilde{v},R_{\tilde{v}}) = \tilde{\mathfrak{P}}_1(u,\tilde{v},v,R_{\tilde{v}}) \ \tilde{\mathfrak{P}}_2(u,\tilde{v},v,R_{\tilde{v}})^{-1}$$

where  $\tilde{\mathfrak{P}}_1$  and  $\tilde{\mathfrak{P}}_2$  are polynomial functions in  $R_{\tilde{v}}$ . Work in a neighbourhood of  $\mathbb{R}^2$  where  $v \neq 0$ . Since  $\chi$  does not depend on v, we must have

(2.6.10) 
$$[(v^3 \partial_v \tilde{\mathfrak{P}}_1) (v \tilde{\mathfrak{P}}_2) - (v^2 \tilde{\mathfrak{P}}_1) (v^2 \partial_v \tilde{\mathfrak{P}}_2)] v^{-4} \tilde{\mathfrak{P}}_2^{-4} \equiv 0.$$

Replace  $R_{\tilde{v}}$  by (U - u)/v. Then, change the point of view by adopting  $u, U, \tilde{v}$  and v as being the new (independent) variables. Noting simply  $\tilde{a}$  and  $a^*$  when the function a is evaluated at the points  $(u, \tilde{v})$  and  $(U, \tilde{v})$ , the condition (2.6.10) becomes

$$\mathfrak{D}(u,\tilde{v},U,v) := (\mathfrak{P}_1 \mathfrak{P}_2 - \mathfrak{P}_3 \mathfrak{P}_4)(u,\tilde{v},U,v) = 0.$$

More precisely

$$\begin{aligned} \mathfrak{P}_{1} &= v^{3} \left[ -f(U) - (U-u) \, \mathbf{Z}(U) \right] \\ &+ v^{2} \left[ \tilde{a} - a^{*} - 2 \, (U-u) \, \mathbf{Z}(U) \, \tilde{v} - (U-u) \, k(U) - (U-u)^{2} \, \tilde{a} \, \mathbf{Z}(U) \right] + \\ &+ v \left[ (U-u) \, (\tilde{b} - b^{*}) + (U-u)^{2} \, \left( \tilde{b} \, f(U) - 2 \, \tilde{a} \, \mathbf{Z}(U) \, \tilde{v} - \tilde{a} \, k(U) \right) \right] + \\ &+ \left[ (U-u)^{2} \, \left( \tilde{b} \, a^{*} - \tilde{a} \, b^{*} \right) \right], \end{aligned}$$

$$\begin{aligned} \mathfrak{P}_{2} &= v^{2} \left[ 1 + 4 \left( U - u \right) f(U) + 2 \left( U - u \right)^{2} f^{(1)}(U) + 3 \left( U - u \right)^{2} \mathbf{Z}(U) \\ &+ (U - u)^{3} \mathbf{Z}^{(1)}(U) \right] \\ &+ v \left[ (U - u)^{3} \left( k^{(1)}(U) + 2 \mathbf{Z}^{(1)}(U) \tilde{v} \right) \\ &+ (U - u)^{2} \left( 2 \partial_{u} a^{*} + 4 \mathbf{Z}(U) \tilde{v} + 2 k(U) \right) + 2 \left( U - u \right) a^{*} \right] \\ &+ \left[ \partial_{u} b^{*}(U - u)^{3} + b^{*}(U - u)^{2} \right], \end{aligned}$$

$$\mathfrak{P}_{3} = v^{2} \left[ 1 + 2 \left( U - u \right) f(U) + (U - u)^{2} \mathbf{Z}(U) \right] + v \left[ 2 \left( U - u \right) a^{*} + 2 \left( U - u \right)^{2} \mathbf{Z}(U) \, \tilde{v} + (U - u)^{2} \, k(U) \right] + \left[ (U - u)^{2} \, b^{*} \right],$$

$$\begin{aligned} \mathfrak{P}_{4} &= v^{3} \left[ -f(U) - (U-u) \left( f^{(1)}(U) + 2 \mathbf{Z}(U) \right) - (U-u)^{2} \mathbf{Z}^{(1)}(U) \right] \\ &+ v^{2} \left[ (U-u) \left( -\partial_{u}a^{*} - 2 \mathbf{Z}(U) \tilde{v} - k(U) \right) \\ &+ (U-u)^{2} \left( -2 \mathbf{Z}^{(1)}(U) \tilde{v} - k^{(1)}(U) - 2 \tilde{a} \mathbf{Z}(U) \right) \\ &+ (U-u)^{3} \left( -\tilde{a} \mathbf{Z}^{(1)}(U) \right) \right] \\ &+ v \left[ (U-u)^{2} \left( -\partial_{u}b^{*} + \tilde{b} f(U) - 2 \tilde{a} \mathbf{Z}(U) \tilde{v} - \tilde{a} k(U) \right) \\ &+ (U-u)^{3} \left( -\tilde{a} k^{(1)}(U) + \tilde{b} f^{(1)}(U) - 2 \tilde{a} \mathbf{Z}^{(1)}(U) \tilde{v} \right) \right] \\ &+ \left[ \tilde{b} \partial_{u}a^{*} - \tilde{a} \partial_{u}b^{*} \right] (U-u)^{3}. \end{aligned}$$

All expressions  $\mathfrak{P}_j$  are polynomial functions in v and  $\tilde{v}$ , with degree at most 5 in v and 3 in  $\tilde{v}$ . It follows that  $\mathfrak{D}$  can be put in the form

$$\mathfrak{D}(u,\tilde{v},U,v) = \sum_{j=0}^{5} \mathfrak{D}_{j}(u,\tilde{v},U) \ v^{j} \equiv 0, \quad \mathfrak{D}_{j}(u,\tilde{v},U) = \sum_{k=0}^{3} \mathfrak{D}_{j}^{k}(u,U) \ \tilde{v}^{k} \,.$$

Of course, the condition  $\mathfrak{D} \equiv 0$  amounts to the same thing as

(2.6.11) 
$$\mathfrak{D}_j^k \equiv 0, \qquad \forall (j,k) \in \{0,\cdots,5\} \times \{0,\cdots,3\}.$$

**Lemme 2.6.3.** Concerning the structures of the functions  $\mathbf{Z}$  and f, there are three possible situations:

a)  $\mathbf{Z} \equiv 0$  and  $f \equiv 0$ ;

b) There exists a constant  $\delta \in \mathbb{R}$  such that  $\mathbf{Z} \equiv 0$  and  $f(u) = (\delta - 2u)^{-1}$ ;

c) There exist two constants  $(\delta_1, \delta_2) \in \mathbb{R}^2$  such that  $\mathbf{Z}(u) = (u^2 - 2\,\delta_1\,u + \delta_2)^{-1}$ and  $f(u) = (-u + \delta_1) \, (u^2 - 2\,\delta_1\,u + \delta_2)^{-1}$ .

**Proof of Lemma 2.6.3.** When  $\mathbf{Z} \equiv 0$ , the function f must satisfy  $2f^2 - f' \equiv 0$  and we are faced with situations a) or b).

From now on, suppose that  $\mathbf{Z} \neq 0$ . First, look at the coefficient  $\mathfrak{D}_0^5$  which is

$$\mathfrak{D}_{0}^{5} = (U-u)^{4} \times \left\{ \left[ \mathbf{Z}(u) \ f(U) - f(u) \ \mathbf{Z}(U) \right] \mathbf{Z}(U) + (U-u) \ \mathbf{Z}(u) \left[ \mathbf{Z}^{(1)}(U) \ f(U) - \mathbf{Z}(U) \ f^{(1)}(U) \right] \right\}.$$

Since  $\mathbf{Z} \neq 0$ , dividing the expression  $\mathfrak{D}_0^5$  by  $\mathbf{Z}(u) \ \mathbf{Z}(U)^2$ , the condition  $\mathfrak{D}_0^5 \equiv 0$  becomes

$$\frac{f(U)}{\mathbf{Z}(U)} - \frac{f(u)}{\mathbf{Z}(u)} - (U-u) \left[\frac{f(U)}{\mathbf{Z}(U)}\right]^{(1)} = 0.$$

It means that the function  $f \mathbb{Z}^{-1}$  is linear with respect to U. Combining this information with (2.6.7), it first remains

(2.6.12) 
$$\frac{f(U)}{\mathbf{Z}(U)} = -U + \delta_1, \qquad \delta_1 \in \mathbb{R}$$

and then, replacing f by  $\mathbf{Z}'/(2\mathbf{Z})$ , we have access to c).

# **2.6.2** Discussion of the case $\mathbf{Z} \neq 0$ .

This paragraph 2.6.2 is aimed to be the source of the situation **ii.3** described in Proposition 2.4.4. First, by changing u into  $u - \delta_1$  and defining  $\gamma := \delta_2 - \delta_1^2$ , we can always assume that

(2.6.13) 
$$f(u) = -u/(u^2 + \gamma)^{-1}, \quad \mathbf{Z}(u) = (u^2 + \gamma)^{-1}.$$

The functions a and b are given by (2.6.6). We have to determine g, k and l.

**Lemme 2.6.4.** Assume that  $Z \neq 0$ . Then, there are constants  $\alpha \in \mathbf{R}$ ,  $\beta \in \mathbf{R}$  and  $\delta \in \mathbf{R}$  such that

(2.6.14) 
$$g(u) = -\frac{\alpha \, u - \beta}{u^2 + \gamma}, \qquad k(u) = \frac{2 \, \alpha}{u^2 + \gamma}, \qquad l(u) = \frac{\delta}{u^2 + \gamma}.$$

Proof of Lemma 2.6.4. Complete (2.6.2) with the introduction of

(2.6.15) 
$$\tilde{c}(u,v) := \partial_v \mathfrak{K} / \partial_{vv}^2 \mathfrak{L}, \qquad \tilde{d}(u,v) := \partial_u \mathfrak{K} / \partial_{vv}^2 \mathfrak{L}.$$

We have here to work with

(2.6.16) 
$$R = -\frac{R_1}{R_2} := -\frac{\tilde{c}(u,v) + x_1 + a(u,v) x_2}{\tilde{d}(u,v) + a(u,v) x_1 + b(u,v) x_2}$$

It is easy to check that  $1 + a R \neq 0$ . Otherwise  $R = a(u, v)^{-1}$  so that  $XR \equiv 0$  which is in contradiction with the assumptions of Proposition 2.4.4. From (2.6.16), we can extract

$$x_1 = -\left\{ R \ \tilde{d}(u,v) + \tilde{c}(u,v) + x_2 \ \left[ R \ b(u,v) + a(u,v) \right] \right\} / \left[ R \ a(u,v) + 1 \right].$$

Use the relation (2.6.16) to compute  $Q(x_1, x_2, u, v) \equiv YR/XR$ . Then replace  $x_1$  as it is indicated above in order to obtain

$$Q(x_1, x_2, u, v) = \mathfrak{Q}(R, x_2, u, v) = \frac{\mathfrak{Q}_0(R, u, v)}{\mathfrak{Q}_2(R, u, v)} + \frac{\mathfrak{Q}_1(R, u, v)}{\mathfrak{Q}_2(R, u, v)} x_2$$

with

$$\mathfrak{Q}_{i}(R, u, v) = \sum_{j=0}^{3} \mathfrak{Q}_{i}^{j}(u, v) R^{j}, \qquad i \in \{0, 1, 2\}$$

and

$$\begin{array}{lll} \mathfrak{Q}_1^3 = \partial_u b \ a - \partial_u a \ b, & \mathfrak{Q}_1^2 = a \ \partial_v b + \partial_u b - \partial_v a \ b, \\ \mathfrak{Q}_1^1 = \partial_v b + \partial_u a \ , & \mathfrak{Q}_1^0 = \partial_v a \ , \\ \mathfrak{Q}_0^3 = \partial_u \tilde{d} \ a - \partial_u a \ \tilde{d} \ , & \mathfrak{Q}_0^2 = \partial_u \tilde{d} + a^2 \ \partial_v (\frac{d}{a}) + a^2 \ \partial_u (\frac{c}{a}) \ , \\ \mathfrak{Q}_2^3 = a \ b \ , & \mathfrak{Q}_2^2 = b + 2 \ a^2 \ , \\ \mathfrak{Q}_2^1 = 3 \ a \ , & \mathfrak{Q}_2^0 = 1 \ . \end{array}$$

The condition (2.4.48) is the same as  $\partial_u R - Q \partial_2 R + XQ \equiv 0$ . Compute

$$\begin{split} \partial_u R &= -\frac{p_u}{p_0} := -\frac{\partial_u \tilde{c} + \partial_u a \, x_2 + R \, \left[\partial_u d + \partial_u a \, x_1 + \partial_u b \, x_2\right]}{\tilde{d} + a \, x_1 + b \, x_2} \,, \\ \partial_2 R &= -\frac{p_2}{p_0} := -\frac{a + b \, R}{\tilde{d} + a \, x_1 + b \, x_2} \,, \\ XR &= -\frac{p_x}{p_0} := -\frac{1 + 2 \, a \, R + b \, R^2}{\tilde{d} + a \, x_1 + b \, x_2} \,. \end{split}$$

With these conventions, the condition (2.4.48) becomes

$$(2.6.17) - p_u + \mathfrak{Q} \ p_2 - \partial_R \mathfrak{Q} \ p_x - R_1 \ \partial_2 \mathfrak{Q} \equiv 0.$$

Observe that (2.6.17) is linear with respect to the variable  $x_2$ . In particular, the term in factor of  $x_2$  must be zero implying that

(2.6.18) 
$$-\mathfrak{Q}_2 (\partial_u a + \partial_u b R) - \partial_R \mathfrak{Q}_1 (R^2 b + 2 a R + 1) + \partial_u a R (R b + a) (R^2 b + 2 a R + 1) + 4 \mathfrak{Q}_1 (R b + a) \equiv 0.$$

Since  $\partial_1 R \neq 0$  (because  $1 + a R \neq 0$ ), the four variables  $R, x_2, u$  and v are independent. The left hand side of (2.6.18) is a polynomial function in R of degree 4 whose four coefficients must all be zero.

Applying this criterion, we can derive the two following conditions :

(2.6.19) 
$$2 b \partial_u a - 2 a \partial_u b + b \partial_v b \equiv 0, \qquad -\partial_u b + 2 b \partial_v a \equiv 0.$$

From (2.6.6), (2.6.13) and (2.6.19), we can derive the expressions of k and l given in (2.6.14). Then, it suffices to exploit (2.6.19) to obtain g.

At this stage, the functions a and b are determined. From the first part of (2.6.2), we can deduce the existence of functions  $\mathbb{F}$ ,  $\mathbb{G}$  and  $\mathfrak{g}$  such that

(2.6.20) 
$$\mathfrak{L}(u,v) = (u^2 + \gamma)^{1/2} \mathbb{F}((v+\alpha) (u^2 + \gamma)^{-1/2} + \mathfrak{g}(u)) + \mathbb{G}(u)$$

where  $\mathbb{F}^{(2)} \neq 0$  and  $\mathfrak{g}$  is some primitive of the function  $u \mapsto \beta (u^2 + \gamma)^{-3/2}$ . Then, testing the second part of (2.6.2) with the formula (2.6.20), we can see that necessarily  $\beta = \gamma = 0$ , yielding the form (2.4.70) for  $\mathfrak{L}$ . As already observed at the level of the proof of Proposition 2.4.4, once  $\mathfrak{L}$  is given by (2.4.70) with  $\mathbb{F}$  and  $\mathbb{G}$  as in (2.4.71), the choice  $\mathfrak{K} = \partial_v \mathfrak{L}$  is suitable. The existence of other convenient functions  $\mathfrak{K}$  will not be discussed.

# **2.6.3** Exclusion of the case $\mathbf{Z} \equiv 0$ and $f \neq 0$ .

In this paragraph, we consider the situation b) of Lemma 2.6.3. When  $\mathbf{Z} \equiv 0$ , we can compute the quantities  $\mathcal{D}_i^k$  of (2.6.11) to find the following list :

$$\begin{split} \mathfrak{D}_{5}^{0} &:= 2 \; (U-u) \; [f^{(1)}(U) - 2 \; f^{2}(U)] \,, \\ \mathfrak{D}_{4}^{1} &:= \; f(u) - f(U) \; + \; (U-u) \; [4 \; f(U) \; f(u) - 4 \; f^{2}(U) + f^{(1)}(U)] \\ &\; + \; (U-u)^{2} \; [2 \; f^{(1)}(U) \; f(u)] \,, \\ \mathfrak{D}_{4}^{0} &:= \; g(u) - g(U) \; + \; (U-u) \; [-4 \; f(U) \; g(U) + g^{(1)}(U) + 4 \; f(U) \; g(u)] \\ &\; + \; (U-u)^{2} \; [2 \; f^{(1)}(U) \; g(u) + k^{(1)}(U) - 3 \; f(U) \; k(U)] \\ &\; + \; (U-u)^{3} \; [k^{(1)}(U) \; f(U) - k(U) \; f^{(1)}(U)] \,, \\ \mathfrak{D}_{3}^{2} &:= \; 2 \; (U-u) \; [f(u) \; f(U) - f^{2}(U) + (U-u) \; f^{(1)}(U) \; f(u)] \,, \\ \mathfrak{D}_{3}^{1} &:= \; (U-u) \; [2 \; f(u) \; g(U) - 4 \; g(U) \; f(U) + 2 \; f(U) \; g(u) + k(u) - k(U)] \\ &\; + \; (U-u)^{2} \; [2 \; g^{(1)}(U) \; f(u) - 6 \; f(U) \; k(U) + 2 \; k(U) \; f(u) \\ &\; + \; 2 \; g(u) \; f^{(1)}(U) + 4 \; k(u) \; f(U) + 2 \; k(U) \; f(u) \\ &\; + \; 2 \; g(u) \; f^{(1)}(U) \; f(u) - 2 \; f^{(1)}(U) \; k(U) \\ &\; + \; 2 \; f^{2}(U) \; k(u) - \; 2 \; f(u) \; f(U) \; k(U) + f^{(1)}(U) \; k(u)] \\ &\; + \; (U-u)^{4} \; [2 \; f(U) \; f(u) \; k^{(1)}(U) - 2 \; f^{(1)}(U) \; k(U)] \,, \end{aligned}$$

$$\begin{split} \mathfrak{D}_{2}^{2} &:= (U-u)^{2} \left[ 3 \ k(u) \ f(U) - 3 \ k(U) \ f(U) \right] (U) + 4 \ k(u) \ f^{2}(U) \\ &\quad -4 \ k(U) \ f(U) \ f(u) + k(u) \ f^{(1)}(U) - k(U) \ f^{(1)}(U) \right] \\ &\quad + (U-u)^{4} \left[ 4 \ f(U) \ f(u) \ k^{(1)}(U) - 4 \ f(u) \ f^{(1)}(U) \ k^{(1)}(U) \right] \\ &\quad + (U-u)^{2} \left[ 2 \ g^{(1)}(U) \ g(u) - 2 \ g^{(2)}(U) + (u) - 1(U) \right] \\ &\quad + (U-u)^{2} \left[ 2 \ g^{(1)}(U) \ g(u) + 2 \ k(U) \ g(u) - 2 \ k(U) \ g(U) \\ &\quad + 4 \ f(U) \ l(u) - 4 \ f(U) \ l(U) + l^{(1)}(U) \right] \\ &\quad + (U-u)^{3} \left[ l^{(1)}(U) \ f(U) + 2 \ k^{(1)}(U) \ g(u) + k^{(1)}(U) \ g(u) \\ &\quad - g^{(1)}(U) \ k(U) - f^{(1)}(U) \ l(U) - 2 \ f(U) \ k(U) \ g(u) \\ &\quad + 2 \ l(u) \ f^{2}(U) + f^{(1)}(U) \ l(u) - k^{2}(U) \right] \\ &\quad + (U-u)^{4} \left[ 2 \ g(u) \ f(U) \ k^{(1)}(U) - 2 \ f^{(1)}(U) \ g(u) \ k(U) \right] \\ &\quad + (U-u)^{4} \left[ 2 \ g^{(1)}(U) \ f(u) + f(U) \ l^{(1)}(U) - 2 \ k^{(2)}(U) \ k(U) \ g(u) \\ &\quad + k^{(1)}(U) \ g^{(1)}(U) - k(U) \ g^{(1)}(U) + 2 \ k^{(1)}(U) \ g(u) \\ &\quad + k^{(1)}(U) \ g^{(1)}(U) - k(U) \ g^{(1)}(U) + 2 \ k^{(1)}(U) \ g^{(1)}(U) \\ &\quad + 4 \ g^{(1)}(U) \ l^{(1)}(U) - l^{(1)}(U) \ l^{(1)}(U) + 2 \ g^{(1)}(U) \ k(U) \ g^{(1)}(U) \\ &\quad + 4 \ g^{(1)}(U) \ l^{(1)}(U) + k(U) \ k^{(1)}(U) - 2 \ g^{(1)}(U) \ k^{(1)}(U) \ g^{(1)}(U) \\ &\quad + l^{(1)}(U) \ f^{(1)}(U) - l^{(1)}(U) \ g^{(1)}(U) \ k^{(1)}(U) \ g^{(1)}(U) \\ &\quad + l^{(1)}(U) \ f^{(1)}(U) \ l^{(1)}(U) \ g^{(1)}(U) \ + 4 \ f^{(1)}(U) \ g^{(1)}(U) \ h^{(1)}(U) \ g^{(1)}(U) \ h^{(1)}(U) \$$

2 5 4 2 5 4 4 - - 2

$$\begin{split} \mathfrak{D}_{1}^{1} &:= (U-u)^{3} \left[ k(u) \, l(U) - 2 \, k(U) \, l(U) + l(u) \, k(U) + 2 \, k(u) \, g^{2}(U) \right. \\ &\quad + 4 \, f(U) \, g(U) \, l(u) - 2 \, f(u) \, l(U) \, g(U) \\ &\quad - 2 \, f(U) \, g(u) \, l(U) - 2 \, g(u) \, g(U) \, k(U) \right] \\ &\quad + (U-u)^{4} \left[ k(u) \, l^{(1)}(U) - 2 \, f(u) \, k(U) \, l(U) + l(u) \, k^{(1)}(U) + 2 \, l(u) \, k(U) \, f(U) \\ &\quad - 2 \, k^{2}(U) \, g(u) + 2 \, k(u) \, k(U) \, g(U) - 2 \, f(u) \, l(U) \, g^{(1)}(U) \\ &\quad - 2 \, g^{(1)}(U) \, g(u) \, k(U) - 2 \, g(u) \, l(U) \, f^{(1)}(U) + 2 \, f(U) \, l^{(1)}(U) \, g(u) \\ &\quad + 2 \, g(U) \, f(u) \, l^{(1)}(U) + 2 \, g(U) \, g(u) \, k^{(1)}(U) \right] \\ &\quad + (U-u)^{5} \left[ k(u) \, f(U) \, l^{(1)}(U) + k^{(1)}(U) \, k(u) \, g(U) - k(U) \, k(u) \, g^{(1)}(U) \\ &\quad - 2 \, k(U) \, l(u) \, f^{(1)}(U) - l(U) \, k(u) \, f^{(1)}(U) + 2 \, l(u) \, f(U) \, k^{(1)}(U) \right] \\ &\quad + l(u) \, k(U) \, f(U) \, l(U) - 2 \, f(u) \, k(U) \, l(U) + k(u) \, k(U) \, g(U) \\ &\quad + l(u) \, k(U) \, f(U) - g(u) \, k^{2}(U) \right] \\ &\quad + (U-u)^{5} \left[ k(u) \, f(U) \, l^{(1)}(U) + k(u) \, k^{(1)}(U) \, g(U) + l(u) \, f(U) \, k^{(1)}(U) \right] \\ &\quad - k(u) \, k(U) g^{(1)}(U) - l(u) \, k(U) \, f^{(1)}(U) - k(u) \, l(U) \, f^{(1)}(U) \right] \\ &\quad - l(u) \, l(U) \, f^{(1)}(U) + l(u) \, g(U) \, k(U) \right] \\ &\quad + (U-u)^{5} \left[ k(u) \, l^{(1)}(U) \, g(U) + l(u) \, l^{(1)}(U) \, g^{(1)}(U) - l(u) \, k(U) \, g^{(1)}(U) \right] \\ &\quad - l(u) \, l(U) \, f^{(1)}(U) - g(u) \, l^{2}(U) \right] \\ &\quad + (U-u)^{5} \left[ l(u) \, l^{(1)}(U) \, g(U) - l(u) \, l(U) \, g^{(1)}(U) \right] . \end{split}$$

- - (---)

We start by obtaining preliminary informations on g, k and l.

**Lemme 2.6.5.** In the case  $\mathbf{Z} \equiv 0$  and  $f(u) = (\delta - 2u)^{-1}$ , we can find a constant  $\beta \in \mathbb{R}$  and a function  $\mathfrak{Q} \in \mathcal{C}^2(\mathbb{R};\mathbb{R})$  such that

- $g(u) = \mathfrak{Q}(u) f(u)^2, \qquad \mathfrak{Q}^{(2)}(u) = 2 \beta,$ (2.6.21)
- $k(u) = \beta f(u),$ (2.6.22)
- $l(u) = \beta g(u).$ (2.6.23)

**Proof of Lemma 2.6.5.** Since  $f \neq 0$ , the condition  $\mathfrak{D}_1^3 \equiv 0$  amounts to the same thing as

$$\frac{k(u)}{f(u)} - \frac{k(U)}{f(U)} + (U - u) \left(\frac{k(U)}{f(U)}\right)^{(1)} \equiv 0.$$

In other words  $k(u) = (\alpha u + \beta) f(u)$  with  $(\alpha, \beta) \in \mathbb{R}^2$ . On the other hand, the restriction  $\mathfrak{D}_0^3 \equiv 0$  is equivalent to

$$\frac{k(u)}{f(u)} \frac{k(U)}{f(U)} - \frac{k(U)^2}{f(U)^2} + (U-u) \frac{k(u)}{f(u)} \left(\frac{k(U)}{f(U)}\right)^{(1)} = 0.$$

This is possible only if  $\alpha = 0$ , yielding (2.6.22).

Knowing (2.6.22), the constraint  $\mathfrak{D}_4^0 \equiv 0$  allows to extract

$$\frac{g(u)}{f(u)^2} = \left[ \frac{g(U)}{f(U)^2} + 4 U \frac{g(U)}{f(U)} - U \frac{g^{(1)}(U)}{f(U)^2} + U^2 \beta \right] + \\ \left[ -4 \frac{g(U)}{f(U)} + \frac{g^{(1)}(U)}{f(U)^2} - 2 U \beta \right] u + \beta u^2 .$$

Since the left hand side of this identity does not depend on the variable U, the coefficients of the right hand side must be constants. We recognize here what is said at the level (2.6.21). In view of of the informations (2.6.21) and (2.6.22), the condition  $\mathfrak{D}_1^2 \equiv 0$  becomes

$$(2.6.24) \begin{array}{l} 2 \ l(u) \ f(U)^2 &= -\beta^2 \ f(u) \ f(U) + \beta^2 \ f(U)^2 \\ &- 2 \ \beta \ \mathfrak{Q}(U) \ f(u) \ f(U)^3 + 2 \ \beta \ \mathfrak{Q}(u) \ f(u)^2 \ f(U)^2 \\ &+ 2 \ f(u) \ f(U) \ l(U) + (U - u) \ \left\{ -2 \ \beta^2 \ f(u) \ f(U)^2 \\ &+ 2 \ \beta \ \mathfrak{Q}^{(1)}(U) \ f(u) \ f(U)^3 + 4 \ \beta \ \mathfrak{Q}(U) \ f(u) \ f(U)^4 \\ &+ 4 \ l(U) \ f(u) \ f(U)^2 - 2 \ f(u) \ f(U) \ l^{(1)}(U) \right\}. \end{array}$$

Dividing this identity by  $2 f(u)^2 f(U)^2$ , the left hand side is simply  $l(u) f(u)^{-2}$ whereas the right hand side is some polynomial function in u of degree 2. To better visualize the content of (2.6.24), we can work with the auxiliary variables  $X := \delta - 2 u$  and  $Y := \delta - 2 U$ . Observe that

$$\mathfrak{Q}(u) = \check{\mathfrak{Q}}(X) = \frac{\beta}{4} X^2 + \check{\mathfrak{Q}}_1 X + \check{\mathfrak{Q}}_0,$$
$$l(u) f(u)^{-2} = \widetilde{\mathfrak{Q}}(X) = \frac{\widetilde{\mathfrak{Q}}_2}{4} X^2 + \widetilde{\mathfrak{Q}}_1 X + \widetilde{\mathfrak{Q}}_0.$$

The relation (2.6.24) is the same as

$$\begin{split} 2\,\tilde{\mathfrak{Q}}(X)\,Y^2 &= -\beta^2\,X\,Y^3 + \beta^2\,X^2\,Y^2 - 2\,\beta\,\check{\mathfrak{Q}}(Y)\,X\,Y + 2\,\beta\,\check{\mathfrak{Q}}(X)\,Y^2 \\ &+ 2\,\check{\mathfrak{Q}}(Y)\,X\,Y + (X - Y)\,\left[ -\beta^2\,X\,Y^2 + \beta\,(\check{\mathfrak{Q}})^{(1)}(Y)\,X\,Y \\ &+ 2\,\beta\,\check{\mathfrak{Q}}(Y)\,X - (\check{\mathfrak{Q}})^{(1)}(Y)\,X\,Y - 2\,\check{\mathfrak{Q}}(Y)\,X \right]. \end{split}$$

For X = 0 and  $Y \neq 0$ , we obtain  $\tilde{\mathfrak{Q}}_0 = \beta \, \check{\mathfrak{Q}}_0$ . Examining the coefficients in factor of  $X^2 Y^2$  and  $X Y^2$ , we get respectively  $\tilde{\mathfrak{Q}}_2 = \beta^2$  and  $\tilde{\mathfrak{Q}}_1 = \beta \, \check{\mathfrak{Q}}_1$ . In other words we have  $\tilde{\mathfrak{Q}} = \beta \, \check{\mathfrak{Q}}$ , that is (2.6.23).

The above study does not exploit all the informations which are contained in the coefficients  $\mathfrak{D}_*^{\star}$ . We can go further

**Lemme 2.6.6.** In the case  $\mathbf{Z} \equiv 0$  and  $f(u) = (\delta - 2u)^{-1}$ , we must have  $k \equiv 0$  and  $l \equiv 0$ . On the other hand, we must have g(u) = df(u) with  $d \in \mathbb{R}$ .

**Proof of Lemma 2.6.6.** Knowing (2.6.22) and (2.6.23), the coefficient  $\mathfrak{D}_1^0$  can be simplified into  $\beta^2 \left[ g(u) \ g(U) - g^2(U) + (U-u) \ g(u) \ g^{(1)}(U) \right] = 0$ . If  $\beta \neq 0$ , exploiting (2.6.21), this is possible only if

$$2 \check{\mathfrak{Q}}(X) \check{\mathfrak{Q}}(Y) Y^2 - 2 \check{\mathfrak{Q}}(Y)^2 X^2 + (X - Y) \check{\mathfrak{Q}}(X) \left[\check{\mathfrak{Q}}^{(1)}(Y) Y^2 + 4 \check{\mathfrak{Q}}(Y) Y\right] = 0.$$

For  $X \neq 0$  and Y = 0, we obtain  $\check{\mathfrak{Q}}_0 = 0$ . Then, dividing by  $Y^2$  and taking again  $X \neq 0$  and Y = 0, we can see that  $\check{\mathfrak{Q}}_1 = 0$  and also  $\beta = 0$ . It means that  $k \equiv 0$  and  $l \equiv 0$ . Then, look at the condition  $\mathfrak{D}_3^0 \equiv 0$  which is

$$g(u) g(U) - g(U)^{2} + (U - u) g'(U) g(u) = 0.$$

Necessarily, we must have  $g(u) = d (e - 2u)^{-1}$  for some  $(d, e) \in \mathbb{R}^2$ . Taking into account (2.6.21), the only possible choice for the constant e is  $e = \delta$  giving rise to the expected result.

**Lemme 2.6.7.** The case  $\mathbf{Z} \equiv 0$  and  $f(u) = (\delta - 2u)^{-1}$  must be excluded.

**Proof of Lemma 2.6.7.** If  $F' \equiv 0$ , we have  $G^{(2)} \equiv 0$  and  $\mathfrak{L}$  is linear in u and v. This is in contradiction with the assumptions  $XR \neq 0$  and  $YR \neq 0$ . Therefore  $F' \neq 0$ ,  $F'(v) = c(v+d)^{-2}$  and  $G^{(2)}(v) = \delta c(v+d)^{-3}$  for some  $c \in \mathbb{R}^*$ . Now, consider the relation (2.4.74) which can be simplified into

(2.6.25) 
$$\begin{aligned} d_v \mathfrak{K} + x_1 \left[ F^{(1)}(V) \ R_{\tilde{v}} + F^{(2)}(V) \ U + G^{(2)}(V) \right] + x_2 \ F^{(1)}(V) \\ - v \ Q_{\tilde{v}} \left[ 2 \ R_{\tilde{v}} \ F^{(1)}(V) + F^{(2)}(V) \ U + G^{(2)}(V) \right] = 0 \,. \end{aligned}$$

Multiply this expression by  $(V + d)^3$  and then take V = -d. It remains

 $c (\delta - 2U) (x_1 + (d + \tilde{v}) Q_{\tilde{v}}) = 0, \qquad c \in \mathbb{R}^*.$ It means that  $Q_{\tilde{v}} = -x_1 (d + \tilde{v})^{-1}$ . Replace  $Q_{\tilde{v}}$  accordingly at the level of (2.6.25). Multiply the expression thus obtained by  $(d + \tilde{v}) F'(V)^{-1}$  and then take  $\tilde{v} = -d$ in order to obtain  $2R_{\tilde{v}}v - 2U + \delta = -2u + \delta = 0$  which cannot be satisfied for all u. This is again a contradiction showing finally that the case  $\mathbf{Z} \equiv 0$  and  $f \neq 0$ must indeed be excluded.

# **2.6.4** Discussion of the case $f \equiv 0$ .

When  $f \equiv 0$ , the structure of the coefficients  $\mathfrak{D}^*_*$  is simplified. It becomes easier to understand what contains the system (2.6.11).

**Lemme 2.6.8.** In the case  $f \equiv 0$ , the functions k, l, and g must satisfy one of the two distinct following restrictions:

a) We have  $k \equiv 0, g^{(1)} \equiv 0$  and  $l^{(1)} \equiv 0$ ;

b) There exists a constant  $\Theta \in \mathbb{R}^*$  such that  $k \equiv -2 \Theta$ ,  $g^{(1)} \equiv \Theta$  and  $l \equiv 0$ .

Proof of Lemma 2.6.8. We start by looking at the condition

$$\mathfrak{D}_4^0(u,U) = g(u) - g(U) + (U-u) \ g^{(1)}(U) + (U-u)^2 \ k^{(1)}(U) = 0 \,.$$

Apply  $\partial_{uu}^2$  to see that  $g^{(2)} = -2 k^{(1)} = \tilde{\delta}$  with  $\tilde{\delta} \in \mathbb{R}$ . Then, consider the restriction  $\mathfrak{D}_2^1 \equiv 0$  which can be rewritten

$$\begin{aligned} -3 \ k(U) \ g(U) + 3 \ k(u) \ g(U) \\ + (U - u) \ \left[ -2 \ k(U)^2 + 2 \ k^{(1)}(U) \ g(u) + k^{(1)}(U) \ g(U) \\ + k(u) \ g^{(1)}(U) + 2 \ k(u) \ k(U) - k(U) \ g^{(1)}(U) \right] \\ + (U - u)^2 \ \left[ k(u) \ k^{(1)}(U) \right] = 0 \,. \end{aligned}$$

Apply  $\partial_{uu}^2$  and replace U by u to deduce that  $g^{(2)} \equiv \tilde{\delta} = 0$ . The next step is to look at the relation  $\mathfrak{D}_3^0 \equiv 0$  which becomes

$$2 g(u) g(U) - 2 g(U)^{2} + l(u) - l(U) + (U - u) [2 g^{(1)}(U) g(u) + 2 k(U) g(u) - 2 k(U) g(U) + l^{(1)}(U)] + (U - u)^{2} [-g^{(1)}(U) k(U) - k(U)^{2}] = 0.$$

Apply again  $\partial_{uu}^2$  and replace U by u. Since  $\tilde{\delta} = 0$ , we obtain by this way that l is a polynomial function of degree at most 2, say  $l(u) = l_2 u^2 + l_1 u + l_0$  with  $(l_0, l_1, l_2) \in \mathbb{R}^3$ . Now, the condition  $\mathfrak{D}_0^0 \equiv 0$  says that

(2.6.26) 
$$\begin{array}{c} l(u) \ g(U) \ l(U) - g(u) \ l^2(U) \\ + (U - u) \ l(u) \ \left[ l^{(1)}(U) \ g(U) - l(U) \ g^{(1)}(U) \right] = 0 \,. \end{array}$$

Composing (2.6.26) on the left with  $\partial_{uuu}^3$  and replacing U by u gives rise to

$$l^{(2)} \left[ l^{(1)} g - l g^{(1)} \right] \equiv 0.$$

At this stage, we claim that  $l^{(2)} \equiv 2 l_2 \equiv 0$ . To see why this is true, first assume that  $l_2 \neq 0$ . Due to the above relation, we must have g = cl with  $c \in \mathbb{R}$ . We cannot have  $c \in \mathbb{R}^*$  because then  $0 \equiv g^{(2)} \equiv c l^{(2)}$  and therefore  $l^{(2)} \equiv 0$  which is a contradiction. Retain that c = 0 so that  $g \equiv 0$ . Looking at the condition  $\mathfrak{D}_1^0 \equiv 0$ , we get

$$l(u) \ l(U) - l^2(U) + (U - u) \ l(u) \ l^{(1)}(U) \equiv 0, \qquad l(u) = l_2 \ u^2 + l_1 \ u + l_0$$

which is not possible when  $l_2 \neq 0$ . We must have  $l_2 \equiv 0$  and  $l(u) = l_1 u + l_0$ .

Now, we can come back to (2.6.26). Apply  $\partial_{uu}^2$  to (2.6.26) and replace U by u. It remains  $l^{(1)} \left[ l^{(1)}g - lg^{(1)} \right] \equiv 0$ . We claim that  $l^{(1)} \equiv l_1 \equiv 0$ . Indeed, if  $l_1 \neq 0$ , we must have g = c l with  $c \in \mathbb{R}$  and the condition  $\mathfrak{D}_1^0 \equiv 0$  becomes

$$l(u) \ l(U) - l(U)^{2} + (U - u) \ l(u) \ l^{(1)}(U) - (U - u)^{2} \ c \ k(U) \ l(u) \ l^{(1)}(U) \equiv 0.$$

Since  $l^{(2)} \equiv 0$ , applying  $\partial_{uu}^2$  and replacing U by u gives rise to

$$l^{(1)}(u) \left[ l^{(1)}(u) + c \, k(u) \, l(u) \right] \equiv 0 \, .$$

Recalling that  $k^{(1)} \equiv 0$ , we can see that this is not possible when  $l^{(1)} \neq 0$ . In other words, the functions k and l are constants, say  $k = k \in \mathbb{R}$  and  $l = l \in \mathbb{R}$ . The condition  $\mathfrak{D}_1^1 \equiv 0$  can be rewritten

2 k 
$$g(U)^2 - 2$$
 k  $g(u)$   $g(U)$   
+  $(U - u) \left[ -2 k^2 g(u) + 2 k^2 g(U) - 2 k g^{(1)}(U) g(u) \right]$   
-  $(U - u)^2 \left[ k^2 g^{(1)}(U) \right] \equiv 0.$ 

Apply  $\partial_{uu}^2$  and replace U by u in order to extract

Now look at the restriction  $\mathfrak{D}_3^0 \equiv 0$  which amounts to the same thing as

$$\begin{bmatrix} 2 g(u) g(U) - 2 g(U)^2 \end{bmatrix} + (U - u) \begin{bmatrix} 2 g^{(1)}(U) g(u) + 2 k g(u) - 2 k g(U) \end{bmatrix} - (U - u)^2 \begin{bmatrix} k g^{(1)}(U) + k^2 \end{bmatrix} \equiv 0.$$

Combining this information with (2.6.27), we obtain that either  $g^{(1)} \equiv 0$  and  $k \equiv 0$ , or that  $k \equiv -2$   $g^{(1)} \equiv -2\Theta$  with  $\Theta \in \mathbb{R}^*$ . We recover here the two situations *a*) and *b*) described in Lemma 2.6.8. In the case *b*), we can even obtain more informations about *l*. Indeed, the condition  $\mathfrak{D}_1^0 \equiv 0$  says that

$$\left\{ g(U)^2 - g(u) \ g(U) + (U - u) \ \left[ k \ g(U) - k \ g(u) - g(u) \ g^{(1)}(U) \right] \right. \\ \left. + 2^{-1} \ (U - u)^2 \ \left[ k \ g^{(1)}(U) \right] \right\} \equiv 0 \,.$$

Apply  $\partial_{uu}^2$  and replace U by u to get  $l \left[ 2 g^{(1)}(u)^2 + 3 k g^{(1)}(u) \right] \equiv -4 l \Theta^2 \equiv 0$ . It follows that  $l \equiv 0$  as expected.

Solving the system (2.6.11) is necessary but not sufficient in order to get solutions of (2.4.39)-(2.4.40). A direct study is still required. We start by examining the case b of Lemma 2.6.8.

**Lemme 2.6.9.** The case  $k \equiv -2 \Theta$ ,  $g^{(1)} \equiv \Theta$ ,  $l \equiv 0$  with  $\Theta \in \mathbb{R}^*$  is excluded.

**Proof of Lemma 2.6.9.** Taking into account (2.6.2), we have

$$\partial^2_{uv} \mathfrak{L} = (\Theta \ u + \Theta_0) \ \partial^2_{vv} \mathfrak{L} \,, \quad \partial^2_{uu} \mathfrak{L} = -2 \ \Theta \ v \ \partial^2_{vv} \mathfrak{L} \,, \quad \Theta \neq 0 \,, \quad \Theta_0 \in \mathbb{R} \,.$$

To simplify the presentation, we will adopt in this paragraph the following conventions  $A := \partial_v \mathfrak{K} / \partial_{vv}^2 \mathfrak{L}$  and  $B := \partial_u \mathfrak{K} / \partial_{vv}^2 \mathfrak{L}$ . We have to deal with

$$R = -\frac{q_0}{p_0} = -\frac{A + x_1 + (\Theta \ u + \Theta_0) \ x_2}{B + (\Theta \ u + \Theta_0) \ x_1 - 2 \ \Theta \ v \ x_2}$$

We can compute

$$\begin{split} \partial_2 R &= -\frac{p_2}{p_0} = - \frac{\Theta \ u + \Theta_0 - 2 \ \Theta \ v \ R}{p_0} \,, \\ \partial_u R &= -\frac{p_u}{p_0} = - \frac{\partial_u A + \Theta \ x_2 + (\partial_u B + \Theta \ x_1) \ R}{p_0} \end{split}$$

and similar expressions for  $\partial_1 R$  and  $\partial_v R$ . Briefly, the quantity Q is given by

$$Q = \frac{p_y}{p_x} = \frac{\partial_v A + \left[\partial_u A + \partial_v B - \Theta x_2\right] R + \left[\partial_u B + \Theta x_1\right] R^2}{1 + 2 \left[\Theta u + \Theta_0\right] R - 2 \Theta v R^2}$$

Now, we can write  $XQ = -p_Q p_x^{-1} p_0^{-1}$  with

$$p_Q := [\partial_u A + \partial_v B - \Theta x_2 + 2 (\partial_u B + \Theta x_1) R] p_x$$
$$- [2(\Theta u + \Theta_0) - 4 \Theta v R] p_y.$$

In this context, the condition (2.4.48) becomes  $p_u p_x - p_y p_2 + p_Q \equiv 0$ . Now, the idea is to multiply this quantity by  $p_0^3$  and to replace everywhere  $p_0 R$  by  $-q_0$  in order to obtain that some polynomial function in  $x_1$  and  $x_2$  with coefficients depending on u and v must be zero. In particular, the coefficient in factor of  $x_1^4$  must be zero. This criterion yields  $\Theta \left[ (\Theta \ u + \Theta_0)^2 + 2 \Theta v \right] \equiv 0$ . This is not possible due to the assumption  $\Theta \in \mathbb{R}^*$ .

### 2.6.5 Necessary conditions.

We can summarize what has been obtained above and in the preceding Sections 2.6.1, 2.6.3 and 2.6.4 through the following statement.

**Proposition 2.6.1.** When  $(XR)(YR) \neq 0$ , the system (2.4.39)-(2.4.40) can be solved if and only if  $\partial_{vv}^2 \mathfrak{L} \neq 0$ . Then, there are two functions f and g in  $\mathcal{C}^2(\mathbb{R};\mathbb{R})$  such that  $\partial_{uu}^2 \mathfrak{L}(u,v) = [f(u)v + g(u)] \partial_{vv}^2 \mathfrak{L}$ . Defining  $\mathbf{Z} := 2f^2 - f'$ , we have either  $\mathbf{Z} \equiv 0$  or  $\mathbf{Z}(u) = (u^2 - 2\delta_1 u + \delta_2)^{-1}$  with  $(\delta_1, \delta_2) \in \mathbb{R}^2$ . In the case  $\mathbf{Z} \equiv 0$ , the function  $\mathfrak{L}$  is necessarily of the form (2.4.65) or (2.4.68).

**Proof of Proposition 2.6.1.** The first part is a repetition of Lemma 2.6.3. When  $\mathbf{Z} \equiv 0$ , due to Lemmas 2.6.7, 2.6.8 and 2.6.9, there are two constants  $b \in \mathbb{R}$  and  $c \in \mathbb{R}$  such that  $\partial_{uv}^2 \mathfrak{L} - b \ \partial_{vv}^2 \mathfrak{L} \equiv 0$  and  $\partial_{uu}^2 \mathfrak{L} - c \ \partial_{vv}^2 \mathfrak{L} \equiv 0$ .

The first relation indicates that

 $\mathfrak{L}(u,$ 

$$v) = \mathbb{G}(u) + \mathbb{F}(b \ u + v), \qquad (\mathbb{F}, \mathbb{G}) \in \mathcal{C}^2(\mathbb{R}; \mathbb{R})^2, \qquad \mathbb{F}^{(2)} \neq 0.$$

Plugging this expression of  $\mathfrak{L}$  in the second relation, we obtain

 $\mathbb{G}^{(2)}(u) + (b^2 - c) \mathbb{F}^{(2)}(b \, u + v) \equiv 0.$ 

When  $b^2 - c \neq 0$ , since the variables u and  $b \ u + v$  are independent, both functions  $\mathbb{F}^{(2)}$  and  $\mathbb{G}^{(2)}$  must be constants and we recover (2.4.65). On the contrary, when  $b^2 - c = 0$ , there is no restriction on  $\mathbb{F}$  (apart from  $\mathbb{F}^{(2)} \neq 0$ ) but the function  $\mathbb{G}$  must be linear in u. It follows that  $\mathfrak{L}(u, v)$  can be put in the form (2.4.68).

# Bibliography

- Andrei Biryuk. On multidimensional Burgers type equations with small viscosity. In *Contributions to current challenges in mathematical fluid mechanics*, Adv. Math. Fluid Mech., pages 1–30. Birkhäuser, Basel, 2004.
- [2] William M. Boothby. An introduction to differentiable manifolds and Riemannian geometry. Academic Press [A subsidiary of Harcourt Brace Jovanovich, Publishers], New York-London, 1975. Pure and Applied Mathematics, No. 63.
- [3] Jean-Yves Chemin. Perfect incompressible fluids, volume 14 of Oxford Lecture Series in Mathematics and its Applications. The Clarendon Press Oxford University Press, New York, 1998. Translated from the 1995 French original by Isabelle Gallagher and Dragos Iftimie.
- [4] C. Cheverry. A deterministic model for the propagation of turbulent oscillations. J. Differ. Equations, 247(9):2637-2679, 2009.
- [5] C. Cheverry and O. Guès. Counter-examples to concentration-cancellation. Arch. Ration. Mech. Anal., 189(3):363–424, 2008.
- [6] C. Cheverry, O. Guès, and G. Métivier. Large-amplitude high-frequency waves for quasilinear hyperbolic systems. Adv. Differential Equations, 9(7-8):829– 890, 2004.
- [7] C. Cheverry and M. Houbad. Compatibility conditions to allow some large amplitude WKB analysis for Burger's type systems. *Phys. D*, 237(10-12):1429–1443, 2008.
- [8] Christophe Cheverry. Sur la propagation de quasi-singularitiés. In Séminaire: Équations aux Dérivées Partielles. 2004–2005, pages Exp. No. VIII, 20. École Polytech., Palaiseau, 2005.
- Christophe Cheverry. Cascade of phases in turbulent flows. Bull. Soc. Math. France, 134(1):33–82, 2006.
- [10] Christophe Cheverry. Recent results in large amplitude monophase nonlinear geometric optics. In *Instability in models connected with fluid flows. I*,

volume 6 of Int. Math. Ser. (N. Y.), pages 267–288. Springer, New York, 2008.

- [11] Christophe Cheverry, Olivier Guès, and Guy Métivier. Oscillations fortes sur un champ linéairement dégénéré. Ann. Sci. École Norm. Sup. (4), 36(5):691– 745, 2003.
- [12] Jean-François Coulombel. From gas dynamics to pressureless gas dynamics. Proc. Amer. Math. Soc., 134(3):683–688 (electronic), 2006.
- [13] Ronald J. DiPerna and Andrew J. Majda. Oscillations and concentrations in weak solutions of the incompressible fluid equations. *Comm. Math. Phys.*, 108(4):667–689, 1987.
- [14] Isabelle Gallagher and Laure Saint-Raymond. On pressureless gases driven by a strong inhomogeneous magnetic field. SIAM J. Math. Anal., 36(4):1159– 1176 (electronic), 2005.
- [15] Emmanuel Grenier. On the nonlinear instability of Euler and Prandtl equations. Comm. Pure Appl. Math., 53(9):1067–1091, 2000.
- [16] Olivier Guès. Ondes multidimensionnelles  $\epsilon$ -stratifiées et oscillations. Duke Math. J., 68(3):401–446, 1992.
- [17] J.-L. Joly, G. Métivier, and J. Rauch. Generic rigorous asymptotic expansions for weakly nonlinear multidimensional oscillatory waves. *Duke Math.* J., 70(2):373–404, 1993.
- [18] Jean-Luc Joly, Guy Métivier, and Jeffrey Rauch. Several recent results in nonlinear geometric optics. In Partial differential equations and mathematical physics (Copenhagen, 1995; Lund, 1995), volume 21 of Progr. Nonlinear Differential Equations Appl., pages 181–206. Birkhäuser Boston, Boston, MA, 1996.
- [19] Pierre-Louis Lions. Mathematical topics in fluid mechanics. Vol. 1, volume 3 of Oxford Lecture Series in Mathematics and its Applications. The Clarendon Press Oxford University Press, New York, 1996. Incompressible models, Oxford Science Publications.
- [20] A. Majda. Compressible fluid flow and systems of conservation laws in several space variables, volume 53 of Applied Mathematical Sciences. Springer-Verlag, New York, 1984.
- [21] Andrew J. Majda and Andrea L. Bertozzi. Vorticity and incompressible flow, volume 27 of Cambridge Texts in Applied Mathematics. Cambridge University Press, Cambridge, 2002.

- [22] F. Poupaud. Global smooth solutions of some quasi-linear hyperbolic systems with large data. Ann. Fac. Sci. Toulouse Math. (6), 8(4):649–659, 1999.
- [23] Steven Schochet. Fast singular limits of hyperbolic PDEs. J. Differential Equations, 114(2):476–512, 1994.
- [24] D. Serre. Oscillations non-linéaires hyperboliques de grande amplitude; dim ≥
  2. In Nonlinear variational problems and partial differential equations (Isola d'Elba, 1990), volume 320 of Pitman Res. Notes Math. Ser., pages 245–294. Longman Sci. Tech., Harlow, 1995.
- [25] Denis Serre. Systems of conservation laws. 1. Cambridge University Press, Cambridge, 1999. Hyperbolicity, entropies, shock waves, Translated from the 1996 French original by I. N. Sneddon.