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The Transport Oka-Grauert principle for simple surfaces

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THE TRANSPORT OKA-GRAUERT PRINCIPLE FOR SIMPLE SURFACES

BY JAN BOHR & GABRIEL P. PATERNAIN

ABSTRACT. — This article considers the attenuated transport equation on Riemannian surfaces in the light of a novel twistor correspondence under which matrix attenuations correspond to holomorphic vector bundles on a complex surface. The main result is a transport version of the classical Oka-Grauert principle and states that the twistor space of a simple surface supports no nontrivial holomorphic vector bundles. This solves an open problem on the existence of matrix holomorphic integrating factors on simple surfaces and is applied to give a range characterisation for the non-Abelian X-ray transform. The main theorem is proved using the inverse function theorem of Nash and Moser and the required tame estimates are obtained from recent results on the injectivity of attenuated X-ray transforms and microlocal analysis of the associated normal operators.

RÉSUMÉ (Le principe de transport d'Oka-Grauert pour les surfaces simples)

Cet article étudie l'équation de transport atténuée sur les surfaces riemanniennes à la lumière d'une nouvelle correspondance de twisteurs dans laquelle les atténuations de matrice correspondent à des fibrés vectoriels holomorphes sur une surface complexe. Le résultat principal est une version de transport du principe classique d'Oka-Grauert et stipule que l'espace des twisteurs d'une surface simple ne supporte aucun fibré vectoriel holomorphe non trivial. Ceci résout un problème ouvert sur l'existence de facteurs intégrants holomorphes matriciels sur des surfaces simples et est appliqué pour donner une caractérisation du domaine pour la transformation en rayons X non abélienne. Le théorème principal est démontré en utilisant le théorème d'inversion locale de Nash et Moser, et les estimations nécessaires sont obtenues à partir de résultats récents sur l'injectivité des transformées en rayons X atténuées et l'analyse microlocale des opérateurs normaux associés.

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1. INTRODUCTION

Inverse problems play a central role in different parts of analysis and geometry. In these problems, there is often an underlying PDE of transport type involving the geodesic vector field of a Riemannian manifold that drives the behaviour of various X-ray transforms. In recent years, a series of papers has culminated in general injectivity results (modulo gauge transformations) for a fundamental class of *nonlinear* X-ray transforms on simple Riemannian surfaces. One goal of this paper is to give a characterisation of the *range* for this class of transforms via a theory of ‘holomorphic integrating factors’. The result is reminiscent of the Ward correspondence for anti-self-dual Yang-Mills fields, but without solitonic degrees of freedom. The range characterisation turns out to be equivalent, via a novel twistor correspondence, to a non-existence theorem for holomorphic vector bundles on certain complex surfaces, resembling the classical Oka-Grauert theorem. Remarkably, the proof of this complex geometric result uses essentially both the theory of transport equations and microlocal analysis.

We now describe the setting of the paper in more detail. Let (M, g) be a compact Riemannian surface with smooth boundary ∂M . Let $SM = \{(x, v) \in TM : g(v, v) = 1\}$ be the unit tangent bundle and X the geodesic vector field on SM . This paper addresses three aspects related to the transport equation

$$(1.1) \quad (X + \mathbb{A})R = 0 \quad \text{on } SM$$

with matrix attenuations $\mathbb{A} \in C^\infty(SM, \mathbb{C}^{n \times n})$: (1) The existence of special solutions to (1.1), called (matrix-)holomorphic integrating factors. (2) A twistor correspondence between attenuations \mathbb{A} and holomorphic vector bundles on a complex surface. (3) A range characterisation for the non-Abelian X-ray transform, which arises from boundary measurements of solutions to (1.1).

These considerations are closely related and are motivated by an inverse problem that we now describe. Assume that ∂M is strictly convex and that M is non-trapping, i.e., all geodesics in M reach ∂M in finite time. We denote with ν the inward pointing unit normal to ∂M and partition the boundary of SM into $\partial SM = \partial_+ SM \cup \partial_- SM$, where

$$\partial_\pm SM = \{(x, v) \in SM : x \in \partial M, \pm g(\nu(x), v) \geq 0\}.$$

Then, by standard ODE theory, equation (1.1) admits a unique continuous solution $R = R^0 : SM \rightarrow \text{GL}(n, \mathbb{C})$, differentiable along the geodesic flow, with $R^0 = \text{Id}$ on $\partial_- SM$ and we define the *scattering data* of $\mathbb{A} \in C^\infty(SM, \mathbb{C}^{n \times n})$ by

$$C_{\mathbb{A}} := R^0|_{\partial_+ SM} \in C^\infty(\partial_+ SM, \text{GL}(n, \mathbb{C})).$$

The nonlinear map $\mathbb{A} \mapsto C_{\mathbb{A}}$ is called the *non-Abelian X-ray transform*. The inverse problem of recovering an attenuation \mathbb{A} from measurements of its scattering data $C_{\mathbb{A}}$ has been subject of a number of recent papers [34, 33, 28] (with earlier in work [50, 46, 12, 30, 10]) and the question of injectivity is now well understood in the following setting: Let $G \subset \text{GL}(n, \mathbb{C})$ be a Lie group with Lie algebra \mathfrak{g} and suppose

that \mathbb{A} is given in terms of a 1-form $A \in \Omega^1(M, \mathfrak{g})$ and a matrix field $\Phi \in C^\infty(M, \mathfrak{g})$ as

$$\mathbb{A}(x, v) = A_x(v) + \Phi(x).$$

We then write $\mathbb{A} = (A, \Phi)$ (referred to as \mathfrak{g} -pair) and $C_{\mathbb{A}} = C_{A, \Phi}$ and note that the scattering data of a \mathfrak{g} -pair is a G -valued function. Two special cases are of particular importance: If $\Phi = 0$, then $C_A = C_{A, 0}$ describes parallel transport of the connection that A induces on the trivial bundle $M \times \mathbb{C}^n$. If $A = 0$ and $\mathfrak{g} = \mathfrak{so}(3)$, then $C_\Phi = C_{0, \Phi}$ arises as measurement data in a novel imaging method called *Polarimetric Neutron Tomography* [7, 28, 42, 17].

A surface (M, g) is called *simple*, if ∂M is strictly convex and M is non-trapping and free of conjugate points.

THEOREM 1.1 (Paternain, Salo, Uhlmann – 2012 & 2020). — *Let (M, g) be a simple surface and $G = \mathrm{U}(n)$ [34] or $G = \mathrm{GL}(n, \mathbb{C})$ [33]. Suppose that two \mathfrak{g} -pairs (A, Φ) and (B, Ψ) have the same scattering data, $C_{A, \Phi} = C_{B, \Psi} \in C^\infty(\partial_+ SM, G)$. Then*

$$(B, \Psi) = (A, \Phi) \triangleleft \varphi := (\varphi^{-1} d\varphi + \varphi^{-1} A \varphi, \varphi^{-1} \Phi \varphi)$$

for some gauge $\varphi \in C^\infty(M, G)$ with $\varphi = \mathrm{Id}$ on ∂M . □

Here $\mathrm{U}(n)$ is the unitary group, with Lie algebra $\mathfrak{u}(n) = \{T \in \mathbb{C}^{n \times n} : T^* = -T\}$ consisting of skew-Hermitian matrices. On manifolds of dimension ≥ 3 a similar result was obtained in [37], using the groundbreaking techniques of Uhlmann and Vasy [49] that also underpin the recent solution of the boundary rigidity problem [48]. For a more detailed account on the history and applications of the non-Abelian X-ray transform we refer to [33, 31] as well as the recent monograph [36].

1.1. HOLOMORPHIC INTEGRATING FACTORS. — Our first contribution concerns (*matrix-*) *holomorphic integrating factors (HIF)*, which were initially sought after as a tool to prove Theorem 1.1 and are now used to obtain the range characterisations in Section 1.3.

To define holomorphic integrating factors we use the fact that every smooth function $F : SM \rightarrow \mathbb{C}^{n \times n}$ has a unique decomposition $F = \sum_{k \in \mathbb{Z}} F_k$ in terms of its vertical Fourier modes (see Section 2 for more details). Then F is called *fibrewise holomorphic* iff $F_k = 0$ for $k < 0$ and we define

$$(1.2) \quad \mathbb{G} = \{F \in C^\infty(SM, \mathrm{GL}(n, \mathbb{C})) : F \text{ and } F^{-1} \text{ are fibrewise holomorphic}\}.$$

DEFINITION 1.2. — A function $F \in \mathbb{G}$ is called a holomorphic integrating factor for the attenuation $\mathbb{A} \in C^\infty(SM, \mathbb{C}^{n \times n})$, if it satisfies the equation $(X + \mathbb{A})F = 0$ on SM .

If an attenuation \mathbb{A} admits holomorphic integrating factors, then necessarily its Fourier modes vanish for $k < -1$, which is to say that \mathbb{A} is a member of the set

$$(1.3) \quad \mathbb{U} = \{\mathbb{A} \in C^\infty(SM, \mathbb{C}^{n \times n}) : \mathbb{A}_k = 0 \text{ for } k < -1\}.$$

Note that \mathbb{U} in particular contains all $\mathfrak{gl}(n, \mathbb{C})$ -pairs $\mathbb{A} = (A, \Phi)$, which have nonzero Fourier modes only for $|k| \leq 1$. We prove the following result:

THEOREM 1.3. — *On a simple surface (M, g) every attenuation $\mathbb{A} \in \mathcal{U}$ admits holomorphic integrating factors.*

In the Abelian case ($n = 1$) this theorem was established in [43] and has since become an indispensable tool in the treatment of attenuated and tensor tomography. The non-Abelian case is much harder and has so far only been addressed in a Euclidean setting. There a weak form of HIF (with F only being continuous) was constructed by Novikov [30] for \mathbb{A} sufficiently small and, without smallness assumption, by Eskin and Ralston [11]. The question of whether smooth matrix HIF exist on simple surfaces has since been open and we can now give an affirmative answer.

The idea behind the proof of Theorem 1.3 is conceptually quite simple. The set \mathbb{G} from (1.2) forms a group and acts on \mathcal{U} from the right by

$$(1.4) \quad \mathbb{A} \triangleleft F = F^{-1}XF + F^{-1}\mathbb{A}F,$$

such that the orbit of $0 \in \mathcal{U}$ contains precisely those attenuations \mathbb{A} that admit holomorphic integrating factors. Theorem 1.3 can thus be reformulated as transitivity of this group action. Using results on the attenuated X-ray transform from [34, 33] and microlocal analysis of the associated normals operators, we show that the derivative of $F \mapsto \mathbb{A} \triangleleft F$ at $\text{Id} \in \mathbb{G}$ is surjective for all $\mathbb{A} \in \mathcal{U}$. After establishing appropriate tame estimates, we use this together with the inverse function theorem of Nash and Moser to show that all orbits of \mathbb{G} are open. As \mathcal{U} is connected, the action must be transitive.

1.2. TWISTOR CORRESPONDENCE. — The second purpose of this article is to promote a novel viewpoint on transport equations as in (1.1) by relating them to holomorphic vector bundles on a *twistor space* Z associated to (M, g) . This is inspired by Penrose’s twistor programme [38] and the paradigm that solutions to integrable systems should be parametrised by complex geometric objects [3, 19, 18, 25].

The twistor space Z can be constructed for any oriented surface (M, g) and, as a smooth manifold, equals the unit disk bundle

$$Z = \{(x, v) \in TM : g(v, v) \leq 1\}.$$

We equip Z with a complex structure that turns Z^{int} into a classical complex surface, and that degenerates at $SM \subset \partial Z$. Postponing precise definitions to Section 4, we note that standard constructions from complex geometry can be carried out ‘smooth up to the boundary’, in particular there is a natural moduli space

$$\mathfrak{M} = \mathfrak{M}_n(Z) = \{\text{Holomorphic rank } n \text{ vector bundles on } Z\} / \sim,$$

where \sim denotes isomorphism of holomorphic vector bundles. We establish several correspondence principles (see Propositions 4.4 and 4.12) which relate the complex geometry on Z to transport problems on SM . In particular, we prove that there is an isomorphism

$$(1.5) \quad \mathfrak{M} \cong \mathcal{U}/\mathbb{G},$$

where the right hand side is the quotient space under the action defined in (1.4). In light of (1.5), we may reformulate Theorem 1.3 as follows (see also Theorem 4.13 where the result is stated in context):

THEOREM 1.4 (Transport Oka-Grauert principle). — *Let Z be the twistor space of a simple surface (M, g) . Then $\mathfrak{M}_n(Z) = 0$, that is, Z supports no nontrivial holomorphic vector bundles.*

This result is reminiscent of the Oka-Grauert principle in complex geometry (cf. [14, 13]), which states that on a Stein manifold the classification of continuous and holomorphic vector bundles coincide. This, amongst other similarities elaborated on in Section 4, suggests the following slogan:

Twistor spaces of simple surfaces behave like (contractible) Stein surfaces.

It is tempting to try and prove Theorem 1.4 by complex geometric methods, thus deriving Theorem 1.3 as corollary. However, there are several obstacles to this: First one would need to show that Z^{int} is indeed a Stein surface – this is easily seen if (M, g) is flat (see Lemma 4.10), but remains challenging for other geometries. Second, one has to deal with the degeneracy of the complex structure at ∂Z , which is a highly nontrivial task. The work of Eskin and Ralston [11] can be interpreted as such a ‘desingularisation’, similar to the one performed, albeit in a different setting, by LeBrun and Mason in [20]. We discuss this approach in more detail in Section 4.4.

For general simple surfaces it seems to be preferable to prove Theorem 1.4 using transport techniques, requiring however, the injectivity result in [34] a priori. It is curious to note that the techniques in [34] were in turn inspired by the Kodaira vanishing theorem from complex geometry.

In [5], we investigate the behaviour of holomorphic vector bundles in a non-simple setting, specifically for the twistor space Z of a closed Anosov surface of genus g . There we demonstrate that $\mathfrak{M}_1(Z, Z^{\text{int}})$, the moduli space of holomorphic line bundles on Z up to equivalence in the interior, is isomorphic to $\mathbb{C}^g / \mathbb{Z}^{2g} \times \mathbb{C}$.

1.3. RANGE CHARACTERISATION. — Finally, we provide a characterisation of the range of the non-Abelian X-ray transform $\mathbb{A} \mapsto C_{\mathbb{A}}$ in terms of boundary objects. This is inspired by the range characterisations for the linear X-ray transform by Pestov and Uhlmann [39] and the subsequent work for attenuated X-ray transforms in [35, 2].

Our main result concerns the range of $(A, \Phi) \mapsto C_{A, \Phi}$ for $\mathfrak{u}(n)$ -pairs and is formulated in terms of a ‘boundary operator’

$$P : C_{\alpha}^{\infty}(\partial_+ SM, \text{Her}_n^+) \longrightarrow C^{\infty}(\partial_+ SM, \text{U}(n)),$$

where $\text{Her}_n^+ \subset \mathbb{C}^{n \times n}$ denotes Hermitian positive definite matrices. Postponing precise definitions to Section 5, we note that the domain of P and the operator P itself are defined in terms of the following objects:

- The *scattering relation* $\alpha : \partial_+ SM \rightarrow \partial_- SM$ of (M, g) , sending starting point and direction of a geodesic to end point and direction.

– A nonlinear type of Hilbert transform

$$\mathcal{H}^+ : C^\infty(\partial SM, \text{Her}_n^+) \longrightarrow C^\infty(\partial SM, \text{GL}(n, \mathbb{C})),$$

defined in terms of the Birkhoff factorisation in loop groups [41], see Section 5.1.

THEOREM 1.5 (Range characterisation for $\mathfrak{u}(n)$ -pairs). — *Suppose that (M, g) is a simple surface (or more generally, that $\mathfrak{M} = 0$). Then an element $q \in C^\infty(\partial_+ SM, \text{U}(n))$ lies in the range of $\{\mathfrak{u}(n)\text{-pairs}\} \ni (A, \Phi) \mapsto C_{A, \Phi}$ if and only if*

$$q = h \cdot Pw \cdot (h^{-1} \circ \alpha)$$

for some $w \in C_\alpha^\infty(\partial_+ SM, \text{Her}_n^+)$ and a contractible map $h \in C^\infty(\partial M, \text{U}(n))$ (i.e., the element induced in $\pi_1(\text{U}(n))$ is trivial.)

Together with Theorem 1.1, we now have a complete understanding of injectivity and range properties of the (nonlinear) non-Abelian X-ray transform $(A, \Phi) \mapsto C_{A, \Phi}$ on simple surfaces.

The theorem is restated as Theorem 5.8, where it is complemented by a number of further characterisations concerning in particular the range of $\Phi \mapsto C_\Phi$ and $A \mapsto C_A$, as well as the case of $\mathfrak{gl}(n, \mathbb{C})$ -valued attenuations. Also a characterisation in the non-simple case is discussed. For precise statements we refer to Section 5.

Let us illustrate the idea behind our range characterisations with the case of the transform $C^\infty(M, \mathfrak{u}(n)) \ni \Phi \mapsto C_\Phi$ (cf. Theorem 5.11 below). To produce an element in the range, we start with some function $w \in C_\alpha^\infty(\partial_+ SM, \text{Her}_n^+)$. This can be extended to a smooth first integral $w^\sharp : SM \rightarrow \text{Her}_n^+$, constant along the geodesic flow. By Birkhoff's factorisation theorem, $w^\sharp = F^*F$ for some $F \in \mathbb{G}$. We now make the assumption that the 0th Fourier mode of F satisfies $F_0 = \text{Id}$, in which case the factorisation is unique. Consider $q := F(F^{-1} \circ \alpha)|_{\partial_+ SM} \in C^\infty(\partial_+ SM, \text{U}(n))$. Then q is given solely in terms of boundary data and in fact equals $q = C_\Phi$, the scattering data of the matrix field

$$(1.6) \quad \Phi = -(XF)F^{-1} \in C^\infty(M, \mathfrak{u}(n)).$$

In particular, q lies in the range of $\Phi \mapsto C_\Phi$. We prove that on a simple surface all elements in the range arise in this way by showing that *every* matrix field Φ is of the form (1.6). This in turn is a consequence of Theorem 1.3.

Theorem 1.5 bears a striking resemblance with the Ward correspondence for the anti-self-dual Yang-Mills (ASDYM) equation by Mason in [24]: there a one-to-one correspondence is set up between solutions to the ASDYM equation on $\tilde{\mathbb{M}} = S^2 \times S^2$ (with split signature) on the one hand and pairs (E, H) on the other hand, where E is a holomorphic vector bundle on a complex twistor space associated with $\tilde{\mathbb{M}}$ and H is a Hermitian metric on E , restricted to a real subspace. The two ‘parameters’ E and H are also referred to as *solitonic* and *radiative/dispersive* degrees of freedom, respectively. Back to Theorem 1.5, and ignoring the gauge h , we see that the range of the non-Abelian X-ray transform is also parametrised by a Hermitian metric, given

by $w \in C_\alpha^\infty(\partial_+ SM, \text{Her}_n^+)$. Notably, there are no solitonic degrees of freedom, which is in line with the Transport Oka-Grauert principle in Theorem 1.4.

At last, let us mention a potential application of our range characterisations. In the context of Polarimetric Neutron Tomography, it has been of recent interest to rigorously study statistical algorithms for recovering a matrix field Φ from noisy measurements of C_Φ [28, 29]. In particular it was shown in [6], that if M is the Euclidean unit disk, then Φ can be recovered by a statistical algorithm *in polynomial time*, provided there is a suitable initialiser. Knowing the range of $\Phi \mapsto C_\Phi$ is a possible starting point to construct a computable initialiser – we hope to address this in future work.

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2. PRELIMINARIES

Here we provide some well-known background material which may be found in [15, 47]; for a recent presentation and its relevance to geometric inverse problems in two dimensions we refer to [36]. Throughout, (M, g) is a compact, oriented two dimensional Riemannian manifold with smooth and possibly empty boundary ∂M .

The unit sphere bundle SM is a compact 3-manifold with boundary $\partial SM = \{(x, v) \in SM : x \in \partial M\}$, containing $\partial_0 SM := \partial_+ SM \cap \partial_- SM$ as submanifold. The geodesic vector field X is the infinitesimal generator of the geodesic flow φ_t on SM and for $(x, v) \in SM$ we denote with $\tau(x, v) \in [0, \infty]$ the first time $t \mapsto \varphi_t(x, v)$ exits SM . The vertical vector field V is defined as the infinitesimal generator of the circle action that the orientation of M induces on the fibres of SM . The pair X, V can be completed to a global frame of $T(SM)$ by considering the vector field $X_\perp := [X, V]$. There are two further structure equations given by $[V, X_\perp] = X$ and $[X, X_\perp] = -KV$, where K is the Gaussian curvature of M . The Sasaki metric on SM is the unique Riemannian metric for which $\{X, X_\perp, V\}$ is an orthonormal frame and the volume for for this metric is denoted by $d\Sigma^3$. The induced area form on ∂SM is denoted by $d\Sigma^2$.

If $x = (x_1, x_2)$ are isothermal coordinates in (M, g) so that the metric has the form $g = e^{2\lambda(x)} dx^2$ and if θ is the angle between v and ∂_{x_1} , then in the (x, θ) coordinates in SM the vector fields have the following explicit formulas:

$$\begin{aligned}
 (2.1) \quad X &= e^{-\lambda} \left(\cos \theta \frac{\partial}{\partial x_1} + \sin \theta \frac{\partial}{\partial x_2} + \left(-\frac{\partial \lambda}{\partial x_1} \sin \theta + \frac{\partial \lambda}{\partial x_2} \cos \theta \right) \frac{\partial}{\partial \theta} \right), \\
 X_\perp &= -e^{-\lambda} \left(-\sin \theta \frac{\partial}{\partial x_1} + \cos \theta \frac{\partial}{\partial x_2} - \left(\frac{\partial \lambda}{\partial x_1} \cos \theta + \frac{\partial \lambda}{\partial x_2} \sin \theta \right) \frac{\partial}{\partial \theta} \right), \\
 V &= \frac{\partial}{\partial \theta}.
 \end{aligned}$$

The space $L^2(SM, \mathbb{C}^n)$ is defined in terms of the measure $d\Sigma^3$ and the standard Hermitian inner product on \mathbb{C}^n . There is an orthogonal decomposition $L^2(SM, \mathbb{C}^n) =$

$\bigoplus_{k \in \mathbb{Z}} H_k$, where H_k is the eigenspace of $-iV$ corresponding to the eigenvalue k . A function $u \in L^2(SM, \mathbb{C}^n)$ has a Fourier series expansion

$$u = \sum_{k=-\infty}^{\infty} u_k,$$

where $u_k \in H_k$. For $k \in \mathbb{Z}$ and $I \subset \mathbb{Z}$ we define

$$(2.2) \quad \Omega_k = C^\infty(SM, \mathbb{C}^n) \cap H_k \quad \text{and} \quad \bigoplus_{k \in I} \Omega_k = C^\infty(SM, \mathbb{C}^n) \cap \left(\bigoplus_{k \in I} H_k \right).$$

DEFINITION 2.1. — Let $u \in L^2(SM, \mathbb{C}^n)$.

(i) u is called *fibrewise holomorphic*, iff $u_k = 0$ for $k < 0$. Similarly, u is called *fibrewise anti-holomorphic*, iff $u_k = 0$ for $k > 0$.

(ii) u is called *even* iff $u_k = 0$ for $k \in 2\mathbb{Z} + 1$, or equivalently iff $u(x, -v) = u(x, v)$ for all $(x, v) \in SM$. Similarly, u is called *odd* iff $u_k = 0$ for $k \in 2\mathbb{Z}$, or equivalently iff $u(x, -v) = -u(x, v)$ for all $(x, v) \in SM$.

We tacitly use these definitions also on ∂SM , noting that $u \in L^2(\partial SM, \mathbb{C}^n)$ has an analogous decomposition $u = \sum_{k \in \mathbb{Z}} u_k$ into Fourier modes.

As in [15] we introduce the first order operators

$$(2.3) \quad \begin{aligned} \eta_+, \eta_- &: C^\infty(SM, \mathbb{C}^n) \longrightarrow C^\infty(SM, \mathbb{C}^n), \\ \eta_+ &:= (X + iX_\perp)/2, \quad \eta_- := (X - iX_\perp)/2. \end{aligned}$$

Clearly $X = \eta_+ + \eta_-$. We have

$$\eta_+ : \Omega_m \longrightarrow \Omega_{m+1}, \quad \eta_- : \Omega_m \longrightarrow \Omega_{m-1}, \quad (\eta_+)^* = -\eta_-, \quad [\eta_\pm, V] = \mp i\eta_\pm.$$

In particular, X has the following important mapping property

$$X : \bigoplus_{k \geq 0} \Omega_k \longrightarrow \bigoplus_{k \geq -1} \Omega_k.$$

We will often use all of the above for smooth functions taking values in complex $n \times n$ matrices, which we indistinctly denote by $\mathbb{C}^{n \times n}$ or $\mathfrak{gl}(n, \mathbb{C})$, if we wish to think of them as Lie algebra of $\mathrm{GL}(n, \mathbb{C})$. (We also use the notation Ω_k in the matrix valued case.)

2.1. FACTORISATION THEOREMS. — For the range characterisations below it will be important to factor $\mathrm{GL}(n, \mathbb{C})$ -valued maps on SM in terms of the group \mathbb{G} from (1.2). This requires a bundle-version of two well-known factorisation theorems for loop groups that we now recall, following the notation and presentation in [41, §8].

Let us denote by $\mathrm{LGL}_n(\mathbb{C})$ the set of all smooth maps $\gamma : S^1 \rightarrow \mathrm{GL}(n, \mathbb{C})$. The set has a natural structure of an infinite dimensional Lie group as explained in [41, §3.2]. This group contains several subgroups which are relevant for us. We shall denote by $\mathrm{L}^+\mathrm{GL}_n(\mathbb{C})$ the subgroup consisting of those loops γ which are boundary values of holomorphic maps

$$\gamma : \{z \in \mathbb{C} : |z| < 1\} \longrightarrow \mathrm{GL}(n, \mathbb{C}).$$

We let ΩU_n denote the set of smooth loops $\gamma : S^1 \rightarrow U(n)$ such that $\gamma(1) = \text{Id}$. The first result we shall use is Theorems 8.1.1 in [41]:

THEOREM 2.2. — *Any loop $\gamma \in \text{LGL}_n(\mathbb{C})$ can be factored uniquely as $\gamma = \gamma_u \cdot \gamma_+$, with $\gamma_u \in \Omega U_n$ and $\gamma_+ \in L^+ \text{GL}_n(\mathbb{C})$. In fact, the product map*

$$\Omega U_n \times L^+ \text{GL}_n(\mathbb{C}) \longrightarrow \text{LGL}_n(\mathbb{C})$$

is a diffeomorphism.

The second result we shall need is the celebrated Birkhoff factorisation theorem (cf. [41, Th. 8.1.1]), stating that loops $\gamma \in \text{LGL}_n(\mathbb{C})$ can be factored as $\gamma = \gamma_- \cdot \Delta \cdot \gamma_+$, where $\gamma_-^*, \gamma_+ \in L^+ \text{GL}_n(\mathbb{C})$ and Δ is a group homomorphism from S^1 into the diagonal matrices in $\text{GL}(n, \mathbb{C})$. In fact, we require only a version for loops with values in the space of positive definite Hermitian matrices, denoted

$$\text{Her}_n^+ = \{H \in \mathbb{C}^{n \times n} : \xi^* H \xi > 0 \text{ for all } \xi \in \mathbb{C}^n \setminus \{0\}\}.$$

In this case, Δ always equals Id and the statement is equivalent to the preceding theorem. We postpone a precise formulation to Theorem 2.3 below.

Consider now a compact non-trapping surface (M, g) with strictly convex boundary. It is well known that such surfaces are diffeomorphic to a disc (cf. [36]) and thus there exists a section $\mathbf{1} : M \rightarrow SM$ which trivialises the bundle SM to $M \times S^1$. One can then perform loop group factorisations fibrewise to obtain:

THEOREM 2.3. — *Let (M, g) be a non-trapping surface with strictly convex boundary.*

(i) *Any $R \in C^\infty(SM, \text{GL}(n, \mathbb{C}))$ can be factored as $R = UF$ (or $R = FU$) where $F \in \mathbb{G}$ and $U \in C^\infty(SM, U(n))$. If R is even, then also U and F are even. Moreover, F is unique up to left (or right) multiplication by a function in $C^\infty(M, U(n))$.*

(ii) *Any $H \in C^\infty(SM, \text{Her}_n^+)$ can be factored as $H = F^*F$ with $F \in \mathbb{G}$. If H is even, then also F is even. Moreover, F is unique up to left multiplication by a function in $C^\infty(M, U(n))$.*

Proof. — Part (i), modulo the statement on even functions, follows from Theorem 2.2, applied to the loop $R(x, \cdot)$ for each $x \in M$. Normalising such that $U(x, \mathbf{1}(x)) = \text{Id}$, the resulting factors U and F are smooth on SM – we refer to Theorem 4.2 in [33] and its proof for more details. Now suppose that $R = UF$ is even. Denoting $a : SM \rightarrow SM$ the antipodal map, defined by $a(x, v) = (x, -v)$, we then have $UF = R = R \circ a = (U \circ a)(F \circ a)$. As the factorisation is unique up to gauge, there exists a function $h \in C^\infty(M, U(n))$ with $U = (U \circ a)h$ and $F = h^*(F \circ a)$. Consequently U and F must be even.

For (ii) note that any $H \in C^\infty(SM, \text{Her}_n^+)$ admits a square root, i.e., there exists an $R \in C^\infty(SM, \text{GL}(n, \mathbb{C}))$ with $H = R^*R$. Using (i), we may decompose $R = UF$, with U unitary and $F \in \mathbb{G}$ and thus $H = F^*U^*UF = F^*F$, as desired. The uniqueness claim follows from the observation that if we had two factorisations $F^*F = \tilde{F}^*\tilde{F}$, then $(\tilde{F}^*)^{-1}F^* = \tilde{F}F^{-1}$. It follows that $\tilde{F}F^{-1}$ is both fibrewise holomorphic and antiholomorphic and thus equal to $h \in C^\infty(M, U(n))$. \square

REMARK 2.4. — There is a boundary version of the theorem in terms of the group $\mathbb{H} = \{f = F|_{\partial SM} : F \in \mathbb{G}\} = \{f \in C^\infty_{\text{Id}}(\partial SM, \text{GL}(n, \mathbb{C})) : f \text{ is fibrewise holomorphic}\}$.

Here and below the subscript Id refers to maps that are homotopic to Id. Indeed, all of the following maps are surjective and injective up to a gauge in $C^\infty_{\text{Id}}(\partial M, \text{U}(n))$:

$$(2.4) \quad C^\infty_{\bullet}(\partial SM, \text{U}(n)) \times \mathbb{H} \longrightarrow C^\infty_{\bullet}(\partial SM, \text{GL}(n, \mathbb{C})), \quad (u, f) \longmapsto uf,$$

$$(2.5) \quad \mathbb{H} \times C^\infty_{\bullet}(\partial SM, \text{U}(n)) \longrightarrow C^\infty_{\bullet}(\partial SM, \text{GL}(n, \mathbb{C})), \quad (u, f) \longmapsto fu,$$

$$(2.6) \quad \mathbb{H} \longrightarrow C^\infty(\partial SM, \text{Her}_+^n), \quad f \longmapsto f^*f,$$

Here C^∞_{\bullet} stands for smooth maps r which have a $\text{GL}(n, \mathbb{C})$ -valued extension to all of SM , or equivalently for which the induced homomorphism $r_* : \mathbb{Z} \times \mathbb{Z} \rightarrow \mathbb{Z}$ between fundamental groups satisfies $r_*(1, 0) = 0$. (Since $\partial SM = \partial M \times S^1$ is a torus and $\text{GL}(n, \mathbb{C})$ has fundamental group \mathbb{Z} , the map r induces a homomorphism $r_* : \mathbb{Z} \times \mathbb{Z} \rightarrow \mathbb{Z}$.) To show (2.4)–(2.6), we extend $r \in C^\infty_{\bullet}(\partial SM, \text{GL}(n, \mathbb{C}))$ to a function $R \in C^\infty(SM, \text{GL}(n, \mathbb{C}))$ and apply Theorem 2.3 to R in order to find appropriate factors for r . We emphasise however, that the factors u and f in the previous display can be found pointwise for every $x \in \partial M$ by solving a Birkhoff factorisation problem in the fibre $S_x M$.

3. MATRIX HOLOMORPHIC INTEGRATING FACTORS

In this section we prove Theorem 1.3 on the existence of matrix holomorphic integrating factors on simple surfaces. Recall from the discussion below the theorem that this is equivalent to proving that the group \mathbb{G} from (1.2) acts transitively on \mathcal{U} from (1.3) via the rule (1.4). To show transitivity, we use the Nash-Moser inverse function theorem in the form of Theorem 2.4.1 in [16, §III], which requires that:

(a) \mathbb{G} is a tame Fréchet Lie group, \mathcal{U} is a connected, tame Fréchet manifold and the action of \mathbb{G} on \mathcal{U} is smooth tame;

(b) for all $\mathbb{A} \in \mathcal{U}$, the derivative of $F \mapsto \mathbb{A} \triangleleft F$ at Id has a tame right inverse.

Here tameness is understood with respect to the grading $(\|\cdot\|_{H^s} : s = 0, 1, \dots)$ by Sobolev norms and condition (a) is satisfied in view of standard estimates; for more details we refer to Appendix Section 6.1. The key condition is (b) and we claim that the derivative in question is given by

$$(3.1) \quad \mathcal{J}_{\mathbb{A}} : T_{\text{Id}}\mathbb{G} \longrightarrow \mathcal{U}, \quad \mathcal{J}_{\mathbb{A}}(H) = XH + [\mathbb{A}, H],$$

with $T_{\text{Id}}\mathbb{G} = \{G \in C^\infty(SM, \mathbb{C}^{n \times n}) : G \text{ fibrewise holomorphic}\}$ and $[\cdot, \cdot]$ denoting the commutator. To see this, fix $\mathbb{A} \in \mathcal{U}$ and consider $F_t = \text{Id} + tH \in \mathbb{G}$ for $H \in T_{\text{Id}}\mathbb{G}$ and small $t \in \mathbb{R}$. Let $s \geq 0$, then for $|t|$ sufficiently small, the Neumann series $\sum_{k \geq 0} (-tH)^k$ converges in the Sobolev space $H^s(SM)$ and one computes that

$$\mathbb{A} \triangleleft F_t = tXH + \mathbb{A} + t\mathbb{A}H - tH\mathbb{A} + o_{\|\cdot\|_{H^s}}(1) \quad \text{as } t \longrightarrow 0,$$

which yields the formula in (3.1).

The proof of Theorem 1.3 is complete, if we show that in the simple case the map $\mathcal{J}_{\mathbb{A}}$ in (3.1) has a tame right inverse for all $\mathbb{A} \in \mathcal{U}$. This is implied by the following

proposition, which is formulated in terms of \mathbb{C}^n -valued functions – the required right inverse for $\mathcal{T}_{\mathbb{A}}$ is obtained by going ‘one level higher’, i.e., viewing $\widehat{\mathbb{A}} = [\mathbb{A}, \cdot]$ as attenuation with values in $\text{End}(\mathbb{C}^{n \times n}) \cong \mathbb{C}^{n^2 \times n^2}$, acting on \mathbb{C}^{n^2} -valued functions.

PROPOSITION 3.1. — *Let (M, g) be a simple surface and $\mathbb{A} \in \mathcal{U}$. Then the map*

$$(X + \mathbb{A}) : \bigoplus_{k \geq 0} \Omega_k \longrightarrow \bigoplus_{k \geq -1} \Omega_k$$

is onto and admits a right inverse $L_{\mathbb{A}} : \bigoplus_{k \geq -1} \Omega_k \rightarrow \bigoplus_{k \geq 0} \Omega_k$ obeying the tame estimate

$$\|L_{\mathbb{A}}f\|_{H^s} \lesssim \|f\|_{H^{s+1}}, \quad f \in \bigoplus_{k \geq -1} \Omega_k, s \geq 0,$$

where \lesssim means up to a constant that depends only on (M, g) , \mathbb{A} and s .

The proposition relies on a number of lemmas that we will discuss first. The first lemma, modulo the tame estimates, appears as Proposition 4.5 in [1] and relies on the fact that the attenuated X-ray transform $I_{A, \Phi}$ is injective on $\Omega_0 \oplus \Omega_1$. Recall that $I_{\mathbb{A}} = I_{A, \Phi} : C^\infty(SM, \mathbb{C}^n) \rightarrow C^\infty(\partial_+ SM, \mathbb{C}^n)$ (for a $u(n)$ -pair $\mathbb{A} = (A, \Phi)$) is defined by $I_{\mathbb{A}}f = u^f|_{\partial_+ SM}$, where $u^f : SM \rightarrow \mathbb{C}^n$ is the unique continuous solution (differentiable along the geodesic flow) of $(X + \mathbb{A})u^f = -f$ on SM and $u^f = 0$ on $\partial_- SM$. The tame estimates can be traced back to mapping properties of the associated normal operator.

LEMMA 3.2. — *Let (M, g) be simple and $\mathbb{A} = (A, \Phi)$ a skew-Hermitian pair. Then for any $f_m + f_{m+1} \in \Omega_m \oplus \Omega_{m+1}$ there is a solution $u \in C^\infty(SM, \mathbb{C}^n)$ to*

$$(3.2) \quad (X + \mathbb{A})u = 0 \quad \text{and} \quad u_m = f_m, \quad u_{m+1} = f_{m+1}.$$

The solution operator $S_{m, \mathbb{A}} : \Omega_m \oplus \Omega_{m+1} \rightarrow C^\infty(SM, \mathbb{C}^n)$, sending $f_m + f_{m+1}$ to $u = S_{m, \mathbb{A}}(f_m + f_{m+1})$, may be chosen to satisfy the tame estimates

$$(3.3) \quad \|S_{m, \mathbb{A}}(f_m + f_{m+1})\|_{H^s} \lesssim \|f_m + f_{m+1}\|_{H^{s+1}}, \quad f_m + f_{m+1} \in \Omega_m \oplus \Omega_{m+1}, s \geq 0,$$

where \lesssim means up to a constant that depends only on (M, g) , \mathbb{A} , m and s .

Proof. — First consider the case $m = 0$. Write $I_{\mathbb{A}}^{0,1} : \Omega_0 \oplus \Omega_1 \rightarrow C^\infty(\partial_+ SM, \mathbb{C}^n)$ for the attenuated X-ray transform, restricted to $\Omega_0 \oplus \Omega_1$. This transform is injective, as the natural gauge from [34, Th.1.3] is fixed on $\Omega_0 \oplus \Omega_1$. Indeed, if $I_{\mathbb{A}}^{0,1}(f) = 0$ for $f \in \Omega_0 \oplus \Omega_1$, then there is a smooth $p : M \rightarrow \mathbb{C}^n$ with $p|_{\partial M} = 0$ such that $f = \Phi p + (X + A)p$. Since $f_{-1} = 0$ we see that $(\eta_- + A_{-1})p = 0$ and this gives $p = 0$ via Lemma 6.2 since any holomorphic function that vanishes on the boundary must be identically zero.

By means of Santaló’s formula the L^2 -adjoint $(I_{\mathbb{A}}^{0,1})^*$ with respect to the measure $\langle \nu(x), v \rangle d\Sigma^2(x, v)$ on $\partial_+ SM$ can be characterised by the equivalence

$$(3.4) \quad f_0 + f_1 = (I_{\mathbb{A}}^{0,1})^*h \iff f_0 = (h^\sharp)_0, \quad f_1 = (h^\sharp)_1.$$

This is valid for all $h \in \mathcal{S}_{\mathbb{A}}^\infty(\partial_+ SM, \mathbb{C}^n)$, the set of functions $h \in C^\infty(\partial_+ SM, \mathbb{C}^n)$ for which the solution h^\sharp to $(X + \mathbb{A})h^\sharp = 0$ with $h^\sharp|_{\partial_+ SM} = h$ is smooth on all of SM (cf. [35, §5], adding a matrix field is unproblematic). The first statement of the

lemma is then the assertion that $(I_{\mathbb{A}}^{0,1})^* : \mathcal{S}_{\mathbb{A}}^{\infty}(\partial_+ SM, \mathbb{C}^n) \rightarrow \Omega_0 \oplus \Omega_1$ is onto. Indeed, if $f_0 + f_1 = (I_{\mathbb{A}}^{0,1})^* h$, then $u = h^{\sharp}$ solves (3.2).

This assertion, together with the tame estimates, is proved with help of the associated normal operator, which was shown to be elliptic in [1]. More precisely, if we make the identification $\Omega_0 \oplus \Omega_1 \cong C^{\infty}(M, \mathbb{C}^{2n})$ (which is possible after trivialising SM), then

$$(I_{\mathbb{A}}^{0,1})^* I_{\mathbb{A}}^{0,1} : C_c^{\infty}(M^{\text{int}}, \mathbb{C}^{2n}) \longrightarrow C^{\infty}(M^{\text{int}}, \mathbb{C}^{2n})$$

is an elliptic pseudodifferential operator of order -1 [1, §4.2]. To proceed, embed M into a closed surface (N, g) and cover N by open subsets U_1, \dots, U_m such that $M \subset U_1$, and $M_j = \bar{U}_j$ is a simple surface for all j . Let $\psi_1, \dots, \psi_m \in C^{\infty}(N, \mathbb{R})$ be such that $\psi_1 \equiv 1$ on M , $\text{supp } \psi_j \subset M_j$ and $\sum_{j=1}^m \psi_j^2 = 1$ on N . Further, extend $\mathbb{A} = (A, \Phi)$ to a pair $\mathbb{A}_1 = (A_1, \Phi_1)$ with compact support in SM_1^{int} and set $\mathbb{A}_2 = \dots = \mathbb{A}_m = 0$. Following the template from [36, §8.2], we see that

$$P = \sum_{j=1}^m \psi_j (I_{\mathbb{A}_j}^{0,1})^* I_{\mathbb{A}_j}^{0,1} \psi_j : C^{\infty}(N, \mathbb{C}^{2n}) \longrightarrow C^{\infty}(N, \mathbb{C}^{2n})$$

is an elliptic pseudodifferential operator of order -1 on N , which is self-adjoint and thus has Fredholm index zero. As each of the operators $I_{\mathbb{A}_j}^{0,1}$ is injective, also P is injective and thus it defines a homeomorphism $P : C^{\infty}(N, \mathbb{C}^{2n}) \rightarrow C^{\infty}(N, \mathbb{C}^{2n})$. For $f = f_0 + f_1 \in \Omega_0 \oplus \Omega_1 \cong C^{\infty}(M, \mathbb{C}^{2n})$ we can now define

$$(3.5) \quad S_{0,\mathbb{A}}(f_0 + f_1) = (I_{\mathbb{A}_1}^{0,1} \psi_1 P^{-1}(Ef))^{\sharp_1}|_{SM},$$

where $E : C^{\infty}(M, \mathbb{C}^{2n}) \rightarrow C^{\infty}(N, \mathbb{C}^{2n})$ is an extension operator and the map $(\cdot)^{\sharp_1} : \mathcal{S}_{\mathbb{A}_1}^{\infty}(\partial_+ SM_1, \mathbb{C}^n) \rightarrow C^{\infty}(SM_1, \mathbb{C}^n)$ is defined similar as above, this time with respect to \mathbb{A}_1 . First note that $u = S_{0,A}(f_0 + f_1)$ indeed solves (3.2): Write $h_1 = I_{\mathbb{A}_1}^{0,1} \psi_1 P^{-1}(Ef)$ and $h = u|_{\partial_+ SM}$, then $h_1^{\sharp_1}|_{SM} = h^{\sharp}$ and thus

$$(I_{\mathbb{A}}^{0,1})^* h = (I_{\mathbb{A}_1}^{0,1})^* h_1 = \psi_1 (I_{\mathbb{A}_1}^{0,1})^* I_{\mathbb{A}_1}^{0,1} \psi_1 (P^{-1}(Ef)) = f \quad \text{on } M,$$

where we used the characterisation (3.4) and the fact that $\psi_1 \equiv 1$ on M , while all other ψ_j 's vanish. Consequently, $S_{0,\mathbb{A}}$ is the desired solution operator and it remains to check the tame estimates.

We check tameness of the operators in (3.5) separately. First note that the extension E , multiplication by ψ_1 and application of $I_{\mathbb{A}}^{0,1}$ satisfy the appropriate tame estimates in a Sobolev scale *without loss of derivatives*. For E this is the content of Seeley's classical article [44] and for $I_{\mathbb{A}}^{0,1}$ this is a standard forward estimate [45, Th. 4.2.1]. Further we have

$$\|P^{-1}g\|_{H^s(N)} \lesssim \|g\|_{H^{s+1}(N)}, \quad g \in C^{\infty}(N, \mathbb{C}^{2n}), \quad s \geq 0,$$

which follows from $P : H^s(N, \mathbb{C}^{2n}) \rightarrow H^{s+1}(N, \mathbb{C}^{2n})$ being injective with closed range. Next, note that $\text{supp } I_{\mathbb{A}_1}^{0,1}(\psi_1 g) \subset K$ for all $g \in C^{\infty}(M_1, \mathbb{C}^{2n})$ and a fixed compact set $K \subset \partial_+ SM_1$ with $K \cap \partial_0 SM_1 = \emptyset$. In order to obtain the tame estimates for $S_{0,\mathbb{A}}$, it thus remains to show

$$(3.6) \quad \|h^{\sharp_1}\|_{H^s(SM_1)} \lesssim \|h\|_{H^s(\partial_+ SM_1)} \quad h \in C_K^{\infty}(\partial_+ SM_1), \quad s \geq 0,$$

where the subscript indicates that $\text{supp } h \subset K$. Let $R : SM_1 \rightarrow U(n)$ be a smooth solution to $(X + \mathbb{A}_1)R = 0$ with $R = \text{Id}$ on $\partial_+ SM_1$ (this exists, because \mathbb{A}_1 has compact support). Further, write $\psi : SM_1 \rightarrow \partial_+ SM_1$ for the foot-point projection, sending (x, v) to the unique point on $\partial_+ SM_1$ on the same geodesic. Then $h^{\sharp_1} = R \cdot \psi^* h$ and we conclude (3.6) from the following mapping properties: Multiplication by R is bounded $H^s(SM_1, \mathbb{C}^n) \rightarrow H^s(SM_1, \mathbb{C}^n)$ and pull-back by ψ is bounded $H_K^s(\partial_+ SM_1, \mathbb{C}^n) \rightarrow H^s(SM_1, \mathbb{C}^n)$ (again subscript K indicates a support restriction). This concludes the Lemma for $m = 0$.

For general $m \in \mathbb{Z}$ the operator $S_{m, \mathbb{A}}$ is obtained by conjugation with $e^{im\theta}$, where the angle θ is chosen with respect to a trivialisation of SM . Indeed, given $f_m + f_{m+1} \in \Omega_m \oplus \Omega_{m+1}$, let $\tilde{f}_0 = e^{-im\theta} f_m$, $\tilde{f}_1 = e^{-im\theta} f_{m+1}$ and $\tilde{\mathbb{A}} = \mathbb{A} + e^{-im\theta} X(e^{im\theta})$. Put $\tilde{u} = S_{0, \tilde{\mathbb{A}}}(\tilde{f}_0 + \tilde{f}_1)$, then $S_{m, \mathbb{A}}(f_m + f_{m+1}) := u = e^{im\theta} \tilde{u}$ satisfies

$$\begin{aligned} Xu &= (X e^{im\theta}) \tilde{u} - e^{im\theta} (\tilde{\mathbb{A}} \tilde{u}) = -\mathbb{A}u, \\ u_m &= e^{im\theta} \tilde{u}_0 = f_m, \quad \text{and} \quad u_{m+1} = e^{i(m+1)\theta} \tilde{u}_1 = f_{m+1}, \end{aligned}$$

as desired. The tame estimates follow immediately from the case $m = 0$ and the proof is complete. \square

The next lemma is essentially a result in complex analysis:

LEMMA 3.3. — *Let (M, g) be non-trapping with strictly convex boundary and $A \in \Omega^1(M, \mathfrak{gl}(n, \mathbb{C}))$ a matrix valued 1-form. Then $\mu_{\pm} = \eta_{\pm} + A_{\pm 1} : \Omega_m \rightarrow \Omega_{m \pm 1}$ is onto and admits a right inverse $T_{A, \pm, m} : \Omega_{m \pm 1} \rightarrow \Omega_m$ obeying the tame estimates*

$$\|T_{A, \pm, m} q\|_{H^{s+1}} \lesssim \|q\|_{H^s}, \quad q \in \Omega_m, \quad s \geq 0,$$

where \lesssim means up to a constant that depends only on (M, g) , A , m and s .

Proof. — We only consider μ_- , the result for μ_+ follows by a similar method. Fix global isothermal coordinates, such that elements in Ω_m are given by $he^{im\theta}$ for a function $h \in C^\infty(\mathbb{D}, \mathbb{C}^n)$ (in particular $A_{\pm 1} = e^{\pm i\theta} a_{\pm}$) and the metric is $g = e^{2\lambda} g_{\text{Eucl}}$ for some conformal factor $\lambda(z)$. Here $\mathbb{D} \subset \mathbb{C}$ is the closed unit disk. Define $\alpha \in C^\infty(\mathbb{D}, \mathbb{C}^{n \times n})$ by $\alpha(z) = e^{\lambda(z)} a_-(z)$. Then a computation similar to [36, Lem. 6.1.8] yields

$$(3.7) \quad \mu_-(he^{im\theta}) = e^{-(m+1)\lambda} \bar{\partial}_\alpha (he^{m\lambda}) \cdot e^{i(m-1)\theta}, \quad h \in C^\infty(\mathbb{D}),$$

where $\bar{\partial}_\alpha : C^\infty(\mathbb{D}, \mathbb{C}^n) \rightarrow C^\infty(\mathbb{D}, \mathbb{C}^n)$ is defined by $\bar{\partial}_\alpha u = \partial_{\bar{z}} u + \alpha u$. As multiplication operators are tame (without loss of derivatives), it suffices to construct a tame right inverse for $\bar{\partial}_\alpha$. By Lemma 6.2 there is a solution $R \in C^\infty(\mathbb{D}, \text{GL}(n, \mathbb{C}))$ to $\bar{\partial}_\alpha R = 0$. Then $R^{-1} \bar{\partial}_\alpha (Ru) = \bar{\partial}_0 u$ for all $u \in C^\infty(\mathbb{D})$ and we have reduced the problem to finding a tame right inverse for $\bar{\partial}_0 \equiv \partial_{\bar{z}}$.

It is a basic result in complex analysis that the equation $\partial_{\bar{z}} u = h$ over \mathbb{C} , given some $h \in C_c^\infty(\mathbb{C}, \mathbb{C}^n)$, is solved by

$$u(z) = Ph(z) = -\frac{1}{2\pi i} \int_{\mathbb{C}} \frac{h(\zeta)}{\zeta - z} d\bar{\zeta} \wedge d\zeta, \quad z \in \mathbb{C}.$$

A right inverse for $\bar{\partial}_0$ on \mathbb{D} is thus given by

$$T : C^\infty(\mathbb{D}, \mathbb{C}^n) \longrightarrow C^\infty(\mathbb{D}, \mathbb{C}^n), \quad Tf = P(Ef)|_{\mathbb{D}},$$

where $E : C^\infty(\mathbb{D}, \mathbb{C}^n) \rightarrow C^\infty(\mathbb{C}, \mathbb{C}^n)$ is a Seeley extension operator, say chosen such that $\text{supp } Ef \subset 2\mathbb{D}$ for all $f \in C^\infty(\mathbb{D}, \mathbb{C}^n)$. Let $\chi \in C_c^\infty(\mathbb{C}, \mathbb{R})$ with $\chi \equiv 1$ on \mathbb{D} , then for all $s \geq 0$ we have

$$\|Tf\|_{H^{s+1}(\mathbb{D})} \leq \|\chi P(Ef)\|_{H^{s+1}(\mathbb{C})} \lesssim \|Ef\|_{H^s(2\mathbb{D})} \lesssim \|f\|_{H^s(\mathbb{D})},$$

where we have used that P is a pseudodifferential operator of order -1 and thus it has the mapping property $H_c^s(\mathbb{C}) \rightarrow H_{\text{loc}}^{s+1}(\mathbb{C})$.

The right-inverse $T_{A,-,m}$ of $\mu_- : \Omega_m \rightarrow \Omega_{m-1}$ is obtained from T by conjugating with R and multiplying with scalar factors as indicated in (3.7). In particular, the tame estimate (3.7) follows from the previous display and the proof is complete. \square

The final ingredient is a non-holomorphic version of Proposition 3.1 and follows from well-known solvability results and estimates concerning the attenuated transport equation over *smooth* functions.

LEMMA 3.4. — *Let (M, g) be non-trapping with strictly convex boundary and suppose $\mathbb{A} \in C^\infty(SM, \mathfrak{u}(n))$. Then $X + \mathbb{A} : C^\infty(SM, \mathbb{C}^n) \rightarrow C^\infty(SM, \mathbb{C}^n)$ has a right inverse $U_{\mathbb{A}}$ that obeys the tame estimates*

$$\|U_{\mathbb{A}}f\|_{H^s(SM)} \lesssim \|f\|_{H^s(SM)} \quad s \geq 0, f \in C^\infty(SM, \mathbb{C}^n),$$

where \lesssim means that the inequality holds up to a multiplicative constant that depends only on (M, g) , \mathbb{A} and s .

Proof. — First assume that both \mathbb{A} and f have compact support in SM^{int} . Then the unique continuous solution $g : SM \rightarrow \mathbb{C}^n$ to

$$(X + \mathbb{A})g = f \text{ on } SM \quad \text{and} \quad g = 0 \text{ on } \partial_- SM$$

vanishes near the glancing region $\partial_0 SM$ and consequently is smooth on SM . Following [28, Lem. 5.12], we have

$$(3.8) \quad \|g\|_{L^2} \leq \tau_\infty \cdot \|f\|_{L^2},$$

where $\tau_\infty = \sup_{SM} \tau$. Let P be a differential operator on SM of order $m \geq 0$ and with constant coefficients with respect to the commuting frame $\{X, P_T, P_V\}$ from [28, Lem. 5.1]. Then $\tilde{g} = Pg$ is the unique solution of

$$(X + \mathbb{A})\tilde{g} = Pf + [\mathbb{A}, P]g \text{ on } SM \quad \text{and} \quad \tilde{g} = 0 \text{ on } \partial_- SM,$$

where $[\cdot, \cdot]$ denotes the commutator. As $\tilde{f} = Pf + [\mathbb{A}, P]g$ has compact support, we may apply (3.8) to obtain

$$\|Pg\|_{L^2} \leq \tau_\infty (\|Pf\|_{L^2} + \|[\mathbb{A}, P]g\|_{L^2}) \lesssim \|f\|_{H^m} + \|g\|_{H^{m-1}},$$

where we used that $[\mathbb{A}, P]$ is a differential operator of order $m - 1$. The H^m -norm of g can be bounded in terms of $\|Pg\|_{L^2}$, if P is taken uniformly elliptic. By induction (and an interpolation argument to pass to non-integral regularities) it then follows

that $\|g\|_{H^s} \lesssim \|f\|_{H^s}$ for all $s \geq 0$, with implicit constant only depending on τ_∞ , s and \mathbb{A} .

The right inverse $U_{\mathbb{A}}$ for general \mathbb{A} and f can be obtained by a standard extension trick: Embed M in the interior of a slightly large manifold (M_1, g) which is also non-trapping and has strictly convex boundary and extend \mathbb{A} to a smooth attenuation $\mathbb{A}_1 : SM_1 \rightarrow \mathfrak{u}(n)$ with compact support in SM_1^{int} . Let $E : C^\infty(SM, \mathbb{C}^n) \rightarrow C^\infty(SM, \mathbb{C}^n)$ be a Seeley extension operator and define $U_{\mathbb{A}}f = g_1|_{SM}$, where $g_1 : SM_1 \rightarrow \mathbb{C}^n$ is the unique solution to $(X + \mathbb{A}_1)g_1 = Ef$ on SM_1 with $g_1 = 0$ on $\partial_- SM_1$. Then by the previous considerations we have

$$\|U_{\mathbb{A}}f\|_{H^s(SM)} \leq \|g_1\|_{H^s(SM_1)} \lesssim \|Ef\|_{H^s(SM_1)} \lesssim \|f\|_{H^s(SM)},$$

which proves the Lemma. □

Proof of Proposition 3.1. — We first give the proof for a skew-Hermitian pair $\mathbb{A} = (A, \Phi)$. Given $f \in \bigoplus_{k \geq -1} \Omega_k$, we use Lemma 3.4 to obtain a smooth solution u to $(X + \mathbb{A})u = f$ and consider $\tilde{u} = u_0 + u_1 + \dots \in \bigoplus_{k \geq 0} \Omega_k$. Then

$$(X + \mathbb{A})\tilde{u} = f - \mu_+u_{-1} - (\Phi u_{-1} + \mu_+u_{-2}) =: f - q_0 - q_{-1}.$$

By Lemma 3.3 we may solve the equations $\mu_-g_0 = q_{-1}$ and $\mu_+g_{-1} = -q_0$ with $g_m \in \Omega_m$ ($m = -1, 0$) and by Lemma 3.2 there exists a smooth solution v to $(X + \mathbb{A})v = 0$ with $v_{-1} = g_{-1}$ and $v_0 = g_0$. Let $\tilde{v} = v_0 + v_1 + \dots \in \bigoplus_{k \geq 0} \Omega_k$, then

$$(X + \mathbb{A})\tilde{v} = \mu_-v_0 - \mu_+v_{-1} = q_{-1} + q_0.$$

In particular $L_{\mathbb{A}}f := \tilde{u} + \tilde{v} \in \bigoplus_{k \geq 0} \Omega_k$ defines a preimage of $f \in \bigoplus_{k \geq -1} \Omega_k$ under $(X + \mathbb{A})$, which implies surjectivity. For the tame estimates note that

$$\begin{aligned} \tilde{u} &= P_{\geq 0} \circ U_{\mathbb{A}}f, \\ \tilde{v} &= P_{\geq 0} \circ S_{\mathbb{A}, -1} \circ (T_{A, -, 0}, -T_{A, +, -1}) \circ Q \circ U_{\mathbb{A}}f, \end{aligned}$$

where S, T, U are as in the lemmas above, $P_{\geq 0} : C^\infty(SM, \mathbb{C}^{n \times n}) \rightarrow \bigoplus_{k \geq 0} \Omega_k$ is the L^2 -orthogonal projection and $Q : C^\infty(SM, \mathbb{C}^{m \times m}) \rightarrow \Omega_{-1} \oplus \Omega_0$ is defined by

$$Qu = (\Phi u_{-1} + \mu_+u_{-2}) \oplus \mu_+u_{-1}.$$

Each of these linear operators was shown to satisfy a tame estimate $\|\bullet\|_{H^s} \lesssim \|\bullet\|_{H^{s+d}}$ of degree $d \in \mathbb{R}$, which can be read off the preceding lemmas and Lemma 6.1. Combined, we see that $L_{\mathbb{A}}$ is tame of degree 1, as desired.

Next assume that $\mathbb{A} \in \mathcal{U}$ is a general attenuation. By Lemma 5.2 in [33] there exists $F \in \mathbb{G}$ such that $\mathbb{B} = \mathbb{A} \triangleleft F \in \mathcal{U}$ defines a skew-Hermitian pair $\mathbb{B} = (B, \Psi)$. Our previous considerations thus yields a tame right inverse $L_{\mathbb{B}}$ to $X + \mathbb{B}$. It is easy to check that $L_{\mathbb{A}} = FL_{\mathbb{B}}F^{-1}$ gives a right inverse for $X + \mathbb{A}$, which inherits the tameness from $L_{\mathbb{B}}$. This concludes the proof. □

By the same methods we obtain the following variant of Theorem 1.3:

PROPOSITION 3.5. — *Let (M, g) be a simple surface. Then any odd attenuation $\mathbb{A} \in \mathcal{U}$ admits even holomorphic integrating factors.*

Proof. — As for Theorem 1.3, the proposition is equivalent to $\mathbb{G}_{\text{ev}} = \{F \in \mathbb{G} : F \text{ even}\}$ acting transitively on $\mathcal{U}_{\text{odd}} = \{\mathbb{A} \in \mathcal{U} : \mathbb{A} \text{ odd}\}$. Using Nash-Moser's theorem, it suffices to prove that for all $\mathbb{A} \in \mathcal{U}_{\text{odd}}$ also the map $(X + \mathbb{A}) : \bigoplus_{k \geq 0} \Omega_{2k} \rightarrow \bigoplus_{k \geq -1} \Omega_{2k+1}$ has a tame right inverse. In terms of $L_{\mathbb{A}}$ from Proposition 3.1 and the projection $P^{\text{ev}} : \bigoplus_{k \geq 0} \Omega_k \rightarrow \bigoplus_{k \geq 0} \Omega_{2k}$ onto even parts we define $L_{\mathbb{A}}^{\text{ev}} : \bigoplus_{k \geq -1} \Omega_{2k+1} \rightarrow \bigoplus_{k \geq 0} \Omega_{2k}$ by $L_{\mathbb{A}}^{\text{ev}} f = P^{\text{ev}} L_{\mathbb{A}} f$. This is a tame map, as $L_{\mathbb{A}}$ and P^{ev} are tame (see Lemma 6.1) and provides the desired right inverse. \square

4. TWISTOR CORRESPONDENCE

Let (M, g) be an oriented Riemannian surface with smooth, possibly empty boundary ∂M . We construct a *twistor space* Z associated to M , which provides a natural habitat to complexify transport problems on SM .

4.1. A COMPLEX SURFACE. — The twistor space of (M, g) is a degenerate complex surface with underlying smooth manifold $Z = \{(x, v) \in TM : g(v, v) \leq 1\}$. The complex structure on Z is described in terms of a complex distribution

$$\mathcal{D} \subset T_{\mathbb{C}}Z \equiv TZ \otimes \mathbb{C},$$

to be thought of as the $(0, 1)$ -bundle; the structure degenerates on $SM \subset \partial Z$ in the sense that there $\mathcal{D} \cap \overline{\mathcal{D}} \neq 0$. The construction of \mathcal{D} is carried out in the following lemma, precisely in equation (4.1); we then review some standard notions from complex geometry in the degenerate context.

It is most convenient to describe the geometry of Z in terms of the fibration

$$p: SM \times \mathbb{D} \longrightarrow Z, \quad (x, v, \omega) \longmapsto (x, v\omega),$$

where $\mathbb{D} = \{\omega \in \mathbb{C} : |\omega| \leq 1\}$ is the complex unit disk and the product $v\omega \in T_x M$ is explained by the complex structure that g and the orientation induce on M . Note that p is a principal S^1 -bundle for the diagonal action $(x, v, \omega) \triangleleft e^{it} = (x, ve^{it}, e^{-it}\omega)$, which has infinitesimal generator

$$\mathbb{V}(x, v, \omega) = V(x, v) + i(\overline{\omega}\partial_{\overline{\omega}} - \omega\partial_{\omega}), \quad (x, v, \omega) \in SM \times \mathbb{D}.$$

In particular, $SM \times \mathbb{D}/S^1 \cong Z$ as smooth manifolds with corners, with boundary hypersurfaces $p(SM \times \partial\mathbb{D}) \equiv SM$ and $p(\partial SM \times \mathbb{D})$. We equip $T_{\mathbb{C}}(SM \times \mathbb{D})$ with the natural Hermitian structure given in terms of the Sasaki metric on SM and the Euclidean metric on \mathbb{D} .

LEMMA 4.1 (Complex structure on Z)

(i) *The following commutator relations hold on $SM \times \mathbb{D}$:*

$$[\omega^2\eta_+ + \eta_-, \mathbb{V}] = i(\omega^2\eta_+ + \eta_-) \quad \text{and} \quad [\partial_{\overline{\omega}}, \mathbb{V}] = i\partial_{\overline{\omega}}.$$

In particular, the complex distribution $\tilde{\mathcal{D}} = \text{span}_{\mathbb{C}}\{\omega^2\eta_+ + \eta_-, \partial_{\overline{\omega}}\}$ on $SM \times \mathbb{D}$ is S^1 -invariant and descends to a distribution on Z , denoted

$$(4.1) \quad \mathcal{D} = p_*(\tilde{\mathcal{D}}).$$

(ii) The Gram matrix of $\{\omega^2\eta_+ + \eta_-, \bar{\omega}^2\eta_- + \eta_+, \mathbb{V}, \partial_{\bar{\omega}}, \partial_{\omega}\}$ at $(x, v, \omega) \in SM \times \mathbb{D}$, denoted $G(x, v, \omega) \in \mathbb{C}^{5 \times 5}$, satisfies

$$(4.2) \quad \det G(x, v, \omega) = (1 - |\omega|^4)^2/4.$$

(iii) The distribution \mathcal{D} from (4.1) is involutive and satisfies

$$\mathcal{D} \cap \bar{\mathcal{D}} = \begin{cases} 0 & \text{on } Z \setminus SM, \\ \text{span}_{\mathbb{C}} X & \text{on } SM. \end{cases}$$

In particular, Z^{int} has a complex structure for which $\mathcal{D} = T^{0,1}Z^{\text{int}}$.

Proof. — For (i) we use the structure equation $[\eta_{\pm}, V] = \mp i\eta_{\pm}$ to the effect that

$$[\omega^2\eta_+ + \eta_-, \mathbb{V}] = \omega^2[\eta_+, V] + [\eta_-, V] + [\omega^2\eta_+, -i\omega\partial_{\omega}] = -i\omega^2\eta_+ + i\eta_- + 2i\omega^2\eta_+,$$

which gives the first relation; the second one is obvious. To check S^1 -invariance, denote by ξ either of the two vector fields $\omega^2\eta_+ + \eta_-$ or $\partial_{\bar{\omega}}$ and define, for $t \in \mathbb{R}$,

$$\xi_t(x, v, \omega) = d\varphi_{-t}^{\mathbb{V}}(\xi(\varphi_t^{\mathbb{V}}(x, v, \omega))).$$

The Lie derivative of ξ along \mathbb{V} equals $\mathcal{L}_{\mathbb{V}}\xi = -[\xi, \mathbb{V}] = -i\xi$. Hence $(d/dt)\xi_t = -i\xi_t$ for all $t \in \mathbb{R}$, which means that $\xi_t = \exp(-it)\xi_0$ and thus the complex line bundle spanned by ξ is S^1 -invariant.

For (ii) one checks, e.g., that

$$\langle \omega^2\eta_+ + \eta_-, \omega^2\eta_+ + \eta_- \rangle = |\omega|^4 \cdot \langle \eta_+, \eta_+ \rangle + \langle \eta_-, \eta_- \rangle = (|\omega|^4 + 1)/2,$$

where we used that $\sqrt{2}\eta_+$ and $\sqrt{2}\eta_-$ are orthonormal. Proceeding similarly with the other combinations, one sees that $G(x, v, \omega)$ is a block matrix with blocks

$$(4.3) \quad \begin{bmatrix} (1 + |\omega|^4)/2 & \omega^2 \\ \bar{\omega}^2 & (1 + |\omega|^4)/2 \end{bmatrix} \quad \text{and} \quad \begin{bmatrix} 1 + 2|\omega|^2 & i\bar{\omega} & -i\omega \\ -i\omega & 1 & 0 \\ i\bar{\omega} & 0 & 1 \end{bmatrix},$$

and the expression for $\det G(x, v, \omega)$ follows by a simple computation.

For (iii), note that $\omega^2\eta_+ + \eta_-$ and $\partial_{\bar{\omega}}$ commute on $SM \times \mathbb{D}$, hence $\tilde{\mathcal{D}}$ and consequently \mathcal{D} are involutive. On $p^{-1}(Z \setminus SM) = \{|\omega| < 1\}$ we have $\tilde{\mathcal{D}} \cap \bar{\tilde{\mathcal{D}}} = 0$ due to (ii), which implies that $\mathcal{D} \cap \bar{\mathcal{D}} = 0$ away from SM . Further,

$$X = \eta_+ + \eta_- \in \tilde{\mathcal{D}} \cap \bar{\tilde{\mathcal{D}}} \quad \text{on } \{\omega = 1\}$$

and as $p|_{\{\omega=1\}}: SM \times \{1\} \rightarrow SM$ is the identity, this implies $X \in \mathcal{D} \cap \bar{\mathcal{D}}$ on SM . The dimension of $\mathcal{D} \cap \bar{\mathcal{D}}$ at $[(x, v, \omega)]$ equals the deficiency of $G(x, v, \omega)$, which is 1 by (4.3), hence $\mathcal{D} \cap \bar{\mathcal{D}}$ is indeed spanned by X on SM . Finally, note that in the interior of Z we have $T_{\mathbb{C}}Z = \mathcal{D} \oplus \bar{\mathcal{D}}$ and this induces a unique complex structure J as follows: for $(x, v) \in Z^{\text{int}}$ define $J_{(x,v)}: (T_{\mathbb{C}}Z)_{(x,v)} \rightarrow (T_{\mathbb{C}}Z)_{(x,v)}$ by $J_{(x,v)}(w_1 \oplus w_2) = -iw_1 + iw_2$, where $w_1 \oplus w_2$ is the unique decomposition into \mathcal{D} and $\bar{\mathcal{D}}$ -components. It is straightforward to verify that J preserves the real tangent space TZ^{int} and by construction it is an almost complex structure with $T^{0,1}Z^{\text{int}} = \mathcal{D}$. Involutive of \mathcal{D} is equivalent to the formal integrability of J and thus the Newlander-Nirenberg theorem implies that J is a complex structure. \square

The preceding lemma shows that Z^{int} is a complex surface in the classical sense, but with complex structure degenerating at SM . Nevertheless the $\bar{\partial}$ -complex of Z^{int} can be extended to all of Z in a way that is smooth up the boundary. We will build this extended $\bar{\partial}$ -complex from scratch and show a posteriori that it coincides with the standard one in the interior. On an open set $U \subset Z$, we define

$$(4.4) \quad \Omega^0(U) \xrightarrow{\bar{\partial}} \Omega^{0,1}(U) \xrightarrow{\bar{\partial}} \Omega^{0,2}(U)$$

as follows: note that the spaces $C^\infty(U)$ and $C^\infty(p^{-1}(U))$ are well defined also when $U \cap \partial Z \neq \emptyset$, and contain \mathbb{C} -valued functions, smooth up to the boundary. Then

$$\begin{aligned} \Omega^0(U) &:= \{h \in C^\infty(p^{-1}(U)) : \mathbb{V}h = 0\} \cong C^\infty(U), \\ \Omega^{0,1}(U) &:= \{(h_1, h_2) \in C^\infty(p^{-1}(U))^2 : \mathbb{V}h_j + ih_j = 0 \ (j = 1, 2)\}, \\ \Omega^{0,2}(U) &:= \{h \in C^\infty(p^{-1}(U)) : \mathbb{V}h + 2ih = 0\}, \end{aligned}$$

and we define

$$(4.5) \quad \bar{\partial}h := ((\omega^2\eta_+ + \eta_-)h, \partial_{\bar{\omega}}h) \quad \text{and} \quad \bar{\partial}(h_1, h_2) := (\omega^2\eta_+ + \eta_-)h_2 - \partial_{\bar{\omega}}h_1,$$

noting that $\bar{\partial}$ has the mapping properties indicated in (4.4) in view of part (i) of the preceding lemma. See Lemma 4.9 for a description of $\bar{\partial}$ in coordinates. If $U \cap \partial Z = \emptyset$, then we recover the usual $\bar{\partial}$ -complex of the complex surface Z^{int} , via isomorphisms

$$(4.6) \quad \Omega^{0,q}(U) \cong \{\alpha \in \Omega^q(U) : \alpha|_{\bar{\mathcal{D}}} = 0\}, \quad q = 1, 2,$$

exhibited in the following lemma.

LEMMA 4.2 (Comparison with standard $\bar{\partial}$ -complex). — *Let $U \subset Z^{\text{int}}$ be open and consider on $p^{-1}(U) \subset \{|\omega| < 1\}$ the complex 1-forms*

$$\tau = \frac{1}{1 - |\omega|^4} (\eta_-^\vee - \bar{\omega}^2\eta_+^\vee) \quad \text{and} \quad \gamma = d\bar{\omega} - i\bar{\omega}V^\vee,$$

where $\{\eta_+^\vee, \eta_-^\vee, V^\vee\}$ is the coframe on SM that is dual to $\{\eta_+, \eta_-, V\}$. Then:

(i) *The following duality relations hold true:*

$$\begin{aligned} \tau(\omega^2\eta_+ + \eta_-) &= \gamma(\partial_{\bar{\omega}}) \equiv 1, \\ \tau(\partial_{\bar{\omega}}) &= \gamma(\omega^2\eta_+ + \eta_-) \equiv 0, \end{aligned} \quad \text{and} \quad \tau, \gamma = 0 \text{ on } \bar{\mathcal{D}} \oplus \text{span } \mathbb{V}$$

(ii) *For $(h_1, h_2) \in \Omega^{0,1}(U)$ and $h \in \Omega^{0,2}(U)$ the differential forms $h_1\tau + h_2\gamma$ and $h\tau \wedge \gamma$ are S^1 -invariant and the maps*

$$(4.7) \quad (h_1, h_2) \mapsto p_*(h_1\tau + h_2\gamma) \quad \text{and} \quad h \mapsto p_*(h\tau \wedge \gamma)$$

yield isomorphisms as in (4.6) for $q = 1$ and $q = 2$, respectively.

(iii) *The isomorphisms from (ii) intertwine the $\bar{\partial}$ -operators from (4.5) with the standard $\bar{\partial}$ -operators of the complex surface U .*

Proof. — The proof of (i) is a simple computation that we omit. As a consequence, $\tau(\xi)$ is constant for $\xi \in \{\omega^2\eta_+ + \eta_-, \bar{\omega}^2\eta_- + \eta_+, \mathbb{V}, \partial_{\bar{\omega}}, \partial_\omega\}$ and, taking Lie derivatives, we see that

$$0 = \mathcal{L}_{\mathbb{V}}(\tau(\xi)) = \mathcal{L}_{\mathbb{V}}\tau(\xi) + \tau(\mathcal{L}_{\mathbb{V}}\xi) = (\mathcal{L}_{\mathbb{V}}\tau - i\tau)(\xi),$$

where in the last step we used Lemma 4.1(i), noting that while e.g., $\mathcal{L}_{\mathbb{V}}(\partial_{\omega}) = +i\partial_{\omega}$, the equality still holds true as $\tau(\partial_{\omega}) = 0$. By Lemma 4.1(ii) such ξ 's form a frame over $p^{-1}(U)$ and thus $\mathcal{L}_{\mathbb{V}}\tau = i\tau$. Arguing similarly, also $\mathcal{L}_{\mathbb{V}}\gamma = i\gamma$ follows. This implies, e.g., that

$$(4.8) \quad \mathcal{L}_{\mathbb{V}}(h_1\tau) = \mathbb{V}h_1\tau + h_1\mathcal{L}_{\mathbb{V}}\tau = -ih_1\tau + ih_1\tau = 0$$

and overall we obtain the desired S^1 -invariance. Hence $\alpha = p_*(h_1\tau + h_2\gamma)$ and $\alpha' = p_*(h\tau \wedge \gamma)$ are well defined differential forms on U . Using part (i) we see that $\alpha, \alpha' = 0$ on \mathcal{D} such that (4.7) indeed defines a map as in (4.6).

We obtain inverse maps as follows: Given $\alpha \in \Omega^1(U)$, we can express its lift $p^*\alpha$ in terms of the 1-forms $\{\tau, \gamma, \bar{\tau}, \bar{\gamma}, \mathbb{V}^{\vee}\}$ (with \mathbb{V}^{\vee} defined similarly to V^{\vee}), which frame $T_{\mathbb{C}}^*(p^{-1}(U))$ by part (i). If $\alpha|_{\mathcal{D}} = 0$, then only the τ - and γ -coefficients of $p^*\alpha$ are nonzero, which is to say that $p^*\alpha = h_1\tau + h_2\gamma$ for some $h_1, h_2 \in C^{\infty}(p^{-1}(U))$. One computes that

$$0 = \mathcal{L}_{\mathbb{V}}(p^*\alpha) = (\mathbb{V}h_1 + ih_1)\tau + (\mathbb{V}h_2 + ih_2)\gamma,$$

hence $(h_1, h_2) \in \Omega^{0,1}(U)$ and we have found the desired preimage of α . The argument for $q = 2$ is completely analogous.

For (iii) consider $f \in \Omega^0(U)$ with lift $h = p^*f$. Then $\bar{\partial}f \in \Omega^1(U)$ (in the classical sense) is uniquely defined by $\bar{\partial}f = df$ on \mathcal{D} and $\bar{\partial}f = 0$ on $\bar{\mathcal{D}}$, hence

$$(4.9) \quad p^*(\bar{\partial}f) = \begin{cases} p^*(df) & \text{on } \tilde{\mathcal{D}}, \\ 0 & \text{on } \bar{\tilde{\mathcal{D}}} \oplus \text{span } \mathbb{V}. \end{cases}$$

Thus $p^*(\bar{\partial}f) = h_1\tau + h_2\gamma$, where $h_1 = (p^*df)(\omega^2\eta_+ + \eta_-) = (\omega^2\eta_+ + \eta_-)h$ and $h_2 = p^*df(\partial_{\bar{\omega}}) = \partial_{\bar{\omega}}h$ – this gives the desired intertwining property on $\Omega^0(U)$. Similarly one shows that if $\alpha \in \Omega^1(U)$ with $\alpha|_{\mathcal{D}} = 0$ and lift $p^*\alpha = h_1\tau + h_2\gamma$, then $\bar{\partial}\alpha \in \Omega^2(U)$ (in the classical sense) satisfies $p^*(\bar{\partial}\alpha) = h\tau \wedge \gamma$ with

$$h = p^*(d\alpha)(\xi, \partial_{\bar{\omega}}) = d(h_1\tau)(\xi, \partial_{\bar{\omega}}) + d(h_2\gamma)(\xi, \partial_{\bar{\omega}}),$$

where $\xi = \omega^2\eta_+ + \eta_-$. The right hand side is easily computed in view of part (i) and $[\xi, \partial_{\bar{\omega}}] = 0$ and one obtains $h = \xi h_2 - \partial_{\bar{\omega}}h_1$, as desired. \square

DEFINITION 4.3. — A function $f \in \Omega^0(U)$ on an open set $U \subset Z$ is called *holomorphic*, if $\bar{\partial}f = 0 \in \Omega^{0,1}(U)$. We then write $f \in \mathcal{O}(U)$.

We emphasise that holomorphic functions on Z are – by definition – smooth up the boundary. We can now draw the first connection to transport problems on SM .

PROPOSITION 4.4 (Twistor correspondence A). — *There is a one-to-one correspondence between holomorphic functions on Z and fibrewise holomorphic first integrals on SM , implemented by the map*

$$(4.10) \quad \mathcal{O}(Z) \xrightarrow{\sim} \{u \in \bigoplus_{k \geq 0} \Omega_k : Xu = 0\} \subset C^{\infty}(SM), \quad f \mapsto f|_{SM}.$$

Proof. — Given $f \in \mathcal{O}(Z)$, denote $h = p^*f \in C^\infty(SM \times \mathbb{D})$. For fixed $(x, v) \in SM$, the function $h(x, v, \cdot)$ is holomorphic in $\{|\omega| < 1\}$ and thus expands as

$$(4.11) \quad h(x, v, \omega) = \sum_{k \geq 0} \omega^k u_k(x, v),$$

with coefficients $u_k(x, v) \in \mathbb{C}$. By Cauchy's integral formula and the \mathbb{V} -invariance,

$$u_k(x, v) = \frac{1}{2\pi i} \int_{|\zeta|=1} \frac{h(x, v, \zeta)}{\zeta^{k+1}} d\zeta = \frac{1}{2\pi} \int_0^{2\pi} f(x, e^{i\theta}) e^{-ik\theta} d\theta.$$

This shows that $u_k(x, v)$ depends smoothly on $(x, v) \in SM$ and moreover, that it is the k th Fourier mode of the fibrewise holomorphic function $u = f|_{SM} \in C^\infty(SM)$. Using the identity $X = \eta_+ + \eta_-$ and, again, holomorphicity of f we see that

$$Xu = (\omega^2 \eta_+ + \eta_-)h|_{\omega=1} = 0,$$

which shows that the map in (4.10) is well defined. The map is clearly injective.

We construct an inverse map as follows: If $u \in \bigoplus_{k \geq 0} \Omega_k$ is a first integral, then its Fourier modes u_k are easily seen to satisfy $\|u_k\|_{C^m(SM)} = O(k^{-\infty})$ ($m \in \mathbb{N}_0$), such that (4.11) defines a function $h \in C^\infty(SM \times \mathbb{D})$. We compute that

$$(4.12) \quad \mathbb{V}h(x, v, \omega) = \sum_{k \geq 0} \omega^k (Vu_k(x, v) - iku_k(x, v)) = 0,$$

which means that it descends to a function $f = p_*h \in \Omega^0(Z)$ with $f|_{SM} = u$. It remains to show that f is holomorphic, or equivalently that

$$g := (\omega^2 \eta_+ + \eta_-)h = 0 \quad \text{and} \quad \partial_{\bar{\omega}}h = 0.$$

The latter equation is satisfied in view of the expansion (4.11) and we know that $g(x, v, 1) = 0$ for all $(x, v) \in SM$, as $Xu = \eta_+u + \eta_-u = 0$. To see that g indeed vanishes for all $\omega \in \mathbb{D}$ note that

$$\partial_{\bar{\omega}}g = 0 \text{ for } |\omega| \leq 1 \quad \text{and} \quad g = 0 \text{ for } |\omega| = 1,$$

which follows in view of the previous observations from $[\omega^2 \eta_+ + \eta_-, \partial_{\bar{\omega}}] = 0$ and the S^1 -invariance of g , respectively. Thus g vanishes on all of $SM \times \mathbb{D}$ by the maximum modulus principle on \mathbb{D} . This completes the proof. \square

REMARK 4.5. — By the same method of proof, we can associate to any $u \in \bigoplus_{k \geq k_0} \Omega_k$ ($k_0 \in \mathbb{Z}$) a function $h \in C^\infty(SM \times \mathbb{D})$ with $\partial_{\bar{\omega}}h = 0$ by setting

$$(4.13) \quad h(x, v, \omega) = \sum_{k \geq k_0} \omega^{k-k_0} u_k(x, v).$$

This is easily checked to satisfy $(\mathbb{V} - ik_0)h = 0$ such that for $k_0 = 0, -1, -2$ we can generate elements of $\Omega^0(Z)$, $\Omega^{0,1}(Z)$ and $\Omega^{0,2}(Z)$, respectively. Vice versa, if $h \in C^\infty(SM \times \mathbb{D})$ satisfies $\partial_{\bar{\omega}}h = 0$ and $(\mathbb{V} - ik_0)h = 0$, then $u(x, v) = h(x, v, 1)$ defines an element in $\bigoplus_{k \geq k_0} \Omega_k$.

For the next result we consider the embedding $\iota_0: M \rightarrow Z$, $\iota_0(x) = (x, 0)$ as zero section. If we equip M with the complex structure induced by g and the orientation, then ι_0 becomes a holomorphic map.

LEMMA 4.6. — *The embedding $\iota_0: M \rightarrow Z$ as zero section is holomorphic.*

Proof. — We have to show that for all $(x, v) \in SM$,

$$(\iota_0)_*(T_x^{0,1}M) \subset \mathcal{D}_{(x,0)}, \quad \text{where } T_x^{0,1}M = \text{span}_{\mathbb{C}}\{v + iv^\perp\}.$$

To see this, pick a neighbourhood $U \subset M$ of x and let $\tilde{\iota}_0: U \rightarrow SU \times \mathbb{D}$ be a local lift of ι_0 with $\tilde{\iota}_0(x) = (x, v, 0)$. Then modulo $\text{span}_{\mathbb{C}}\mathbb{V}(x, v, 0)$ we have

$$(\tilde{d}\tilde{\iota}_0)_x(v + iv^\perp) \equiv X(x, v) - iX_\perp \equiv \eta_-(x, v) \in \tilde{\mathcal{D}}_{(x,0)},$$

which yields the desired inclusion after push-forward by p . □

As a consequence there is a well-defined map

$$\iota_0^*: \mathcal{O}(Z) \longrightarrow \mathcal{O}(M),$$

where $\mathcal{O}(M)$ denotes the space of holomorphic functions on M that are smooth up to the boundary. Under the identification $C^\infty(M) \cong \Omega_0$ this is also given as

$$\mathcal{O}(M) = \{g \in \Omega_0 : \eta_-g = 0\}.$$

The following result is then a consequence of the characterisation in Proposition 4.4 and a classical result of Pestov and Uhlmann on the surjectivity of the adjoint X-ray transform I_0^* .

COROLLARY 4.7 (Cartan extension – transport version). — *Suppose Z is the twistor space of a simple surface (M, g) . Then the map $\iota_0^*: \mathcal{O}(Z) \rightarrow \mathcal{O}(M)$ is onto.*

Proof. — As explained above, any $g \in \mathcal{O}(M)$ may be viewed as element in Ω_0 with $\eta_-g = 0$. By [40] (see also Theorem 8.2.2 in [36]) there exists a solution $w \in C^\infty(SM)$ to $Xw = 0$ with Fourier mode $w_0 = g$. Then $u = w_0 + w_2 + \dots$ is smooth, fibrewise holomorphic and satisfies $Xu = 0$. By the preceding proposition (and equation (4.11) in the proof) u gives rise to an element $f \in \mathcal{O}(Z)$ with

$$p^*f(x, v, \omega) = w_0(x) + \omega^2w_2(x, v) + \dots,$$

in particular, $\iota_0^*f(x) = p^*f(x, v, 0) = w_0(x) = g(x)$, as desired. □

The preceding result can be viewed as ‘Cartan extension theorem’ and implies in particular that the twistor space Z of a simple surface admits an abundance of holomorphic functions. This is first evidence for Z behaving like a Stein surface, as claimed in Section 1.2. Further evidence is provided by Theorem 4.13 and its corollaries.

4.2. COORDINATES AND EUCLIDEAN CASE. — It is instructive to express the twistor space Z in terms of isothermal coordinates on (M, g) . Suppose that (x_1, x_2) are coordinates on an open subset $O \subset M$, such that $g|_O = e^{2\lambda}dx^2$, with $\lambda \in C^\infty(O, \mathbb{R})$. Viewing O as subset of \mathbb{C} with complex coordinate $z = x_1 + ix_2$, we define

$$(4.14) \quad Z_O = O \times \mathbb{D}, \quad \mathcal{D}_O = \text{span}\{\Xi, \partial_{\bar{\mu}}\} \subset T_{\mathbb{C}}Z_O,$$

where μ is the coordinate of the \mathbb{D} -factor and the vector field Ξ is defined by

$$\Xi = e^{-\lambda}[\mu^2\partial_z + \partial_{\bar{z}} + (\mu^2\partial_z\lambda - \partial_{\bar{z}}\lambda)(\bar{\mu}\partial_{\bar{\mu}} - \mu\partial_{\mu})].$$

On $SM|_O$ we have coordinates (x_1, x_2, θ) , where $\theta \in \mathbb{R}/(2\pi\mathbb{Z})$ is the (oriented) angle of a unit vector with ∂_{x_1} and there is an isomorphism

$$\rho_O : SM|_O \xrightarrow{\sim} O \times S^1, \quad (x_1, x_2, \theta) \mapsto (x_1 + ix_2, e^{i\theta}),$$

which is made implicit below.

The next lemma shows that (Z_O, \mathcal{D}_O) is a (degenerately) complex surface – its proof is independent from the analogous Lemma 4.1 and the two constructions are seen to be equivalent below. We use the following notation:

$$(4.15) \quad \Lambda := e^{-\lambda}(\mu^2\partial_z\lambda - \partial_{\bar{z}}\lambda) \in C^\infty(Z_O, \mathbb{C}).$$

LEMMA 4.8 (Complex structure in coordinates). —

- (i) $[\Xi, \partial_{\bar{\mu}}] = -\Lambda\partial_{\bar{\mu}}$, hence \mathcal{D}_O is involutive.
- (ii) $\mathcal{D}_O \cap \bar{\mathcal{D}}_O = 0$ on $Z_O \setminus (O \times S^1)$.
- (iii) On $O \times S^1 \cong SM|_O$ we have $\Xi = \mu X$.

Proof. — Parts (i) and (ii) follow from simple computations that we omit. For part (iii) note that $e^{i\theta}\partial_z = \frac{1}{2}(\cos\theta\partial_{x_1} + \sin\theta\partial_{x_2}) + \frac{i}{2}(\sin\theta\partial_{x_1} - \cos\theta\partial_{x_2})$, hence the coordinate description of X from (2.1) is equivalent to

$$(4.16) \quad X = e^{-\lambda}(e^{i\theta}\partial_z + e^{-i\theta}\partial_{\bar{z}} + (e^{i\theta}\partial_z\lambda - e^{-i\theta}\partial_{\bar{z}}\lambda)(i\partial_\theta)).$$

Under the isomorphism ρ_O we have $\mu = e^{i\theta}$ such that $\bar{\mu}\partial_{\bar{\mu}} - \mu\partial_{\mu} = i\partial_\theta$ and hence $X = \mu^{-1}\Xi$ for $|\mu| = 1$, as desired. \square

Define a map $\kappa_O : Z_O \rightarrow Z$ by $\kappa_O(x, \mu) = p(x, 0, \mu)$ (where 0 stands for the angle $\theta = 0$) and note that $\kappa_O(Z_O) = \{(x, v) \in Z : x \in O\}$.

LEMMA 4.9 (Comparison with invariant twistor space)

- (i) The map κ_O is a diffeomorphism onto its image and $(\kappa_O)_*(\mathcal{D}_O) = \mathcal{D}$.
- (ii) Let $U \subset \kappa_O(Z_O)$ be open, then pullback by κ_O induces isomorphisms that fit into the commutative diagram

$$(4.17) \quad \begin{array}{ccccc} \Omega^0(U) & \xrightarrow{\bar{\partial}} & \Omega^{0,1}(U) & \xrightarrow{\bar{\partial}} & \Omega^{0,2}(U) \\ \downarrow \wr & & \downarrow \wr & & \downarrow \wr \\ C^\infty(\kappa_O^{-1}(U)) & \xrightarrow{\bar{\partial}} & C^\infty(\kappa_O^{-1}(U))^2 & \xrightarrow{\bar{\partial}} & C^\infty(\kappa_O^{-1}(U)), \end{array}$$

where the $\bar{\partial}$ -operators on the bottom are given (with Λ as in (4.15)) by

$$\bar{\partial}f = (\Xi f, \partial_{\bar{\mu}}f) \quad \text{and} \quad \bar{\partial}(f_1, f_2) = (\Xi + \Lambda)f_2 - \partial_{\bar{\mu}}f_1.$$

Proof. — For (i) define $q : SM|_O \times \mathbb{D} \rightarrow Z_O$ by $q(x, \theta, \omega) = (x, e^{i\theta}\omega)$; then q is smooth, S^1 -invariant and satisfies $\kappa_O(q(x, \theta, \omega)) = [(x, 0, e^{i\theta}\omega)] = p(x, \theta, \omega)$. In particular, q descends to an inverse of κ_O , which is consequently a diffeomorphism. We claim that $q_*(\tilde{\mathcal{D}}) \subset \mathcal{D}_O$ — this will complete the proof of part (i) by comparing ranks. Indeed, one derives, similarly to (4.16), the coordinate expression

$$(4.18) \quad (\omega^2\eta_+ + \eta_-) = e^{-i\theta}e^{-\lambda} (\omega^2e^{2i\theta}\partial_z + \partial_{\bar{z}} + (\omega^2e^{2i\theta}\partial_z\lambda - \partial_{\bar{z}}\lambda)(i\partial_\theta))$$

and employs this to compute the push forwards

$$\begin{aligned} dq_{(x,\theta,\omega)}(\partial_\theta) &= (\partial_t|_{t=0})q(x, \theta + t, \omega) = (\partial_t|_{t=0})q(x, \theta, e^{it}\omega) \\ &= ie^{i\theta}\omega\partial_\mu - ie^{-i\theta}\bar{\omega}\partial_{\bar{\mu}} = -i(\bar{\mu}\partial_{\bar{\mu}} - \mu\partial_\mu), \\ dq_{(x,\theta,\omega)}(\omega^2\eta_+ + \eta_-) &= e^{-i\theta}\Xi(q(x, \theta, \omega)), \\ dq_{(x,\theta,\omega)}(\partial_{\bar{\omega}}) &= e^{-i\theta}\partial_{\bar{\mu}}, \end{aligned}$$

from which the claim follows.

For part (ii) we first note that the vertical arrows in (4.17) are defined as ‘pull-backs’ by κ_O , understood as follows: Given a function $h \in C^\infty(p^{-1}(U))$ (representing an element in $\Omega^0(U), \Omega^{0,1}(U)$ or $\Omega^{0,2}(U)$), we write $\kappa_O^*h(x, \mu) = h(x, 0, \mu)$. Now consider $h \in \Omega^0(U)$, then as $\mathbb{V}h = 0$,

$$(4.19) \quad \partial_\theta h(x, 0, \mu) = \mathbb{V}h(x, 0, \mu) = -i(\bar{\omega}\partial_{\bar{\omega}} - \omega\partial_\omega)h(x, 0, \mu)$$

and together with (4.18) we obtain

$$\kappa_O^* ((\omega^2\eta_+ + \eta_-)h) (x, \mu) = \Xi(\kappa_O^*h)(x, \mu).$$

Similarly, $\kappa_O^*(\partial_{\bar{\omega}})h(x, \mu) = \partial_{\bar{\mu}}(\kappa_O^*h)(x, \mu)$ and thus the left square in (4.17) commutes. Next, if $(h_1, h_2) \in \Omega^{0,1}(U)$, then $\mathbb{V}h_j = -ih_j$ ($j = 1, 2$) and similarly to (4.19) we have

$$\partial_\theta h_2(x, 0, \mu) = \mathbb{V}h_2(x, 0, \mu) = -i(\bar{\omega}\partial_{\bar{\omega}} - \omega\partial_\omega)h_2(x, 0, \mu) - ih_2(x, 0, \mu),$$

such that $\kappa_O^* ((\omega^2\eta_+ + \eta_-)h_2) (x, \mu) = (\Xi + \Lambda)(\kappa_O^*h_2)(x, \mu)$. The computation for $\partial_{\bar{\omega}}$ remains unchanged and thus also the right square in (4.17) commutes. \square

We can gain more insight into the (degenerate) complex surface Z in the case that (M, g) is a Euclidean domain. First suppose that $M = \mathbb{R}^2$, such that $Z = \mathbb{C} \times \mathbb{D}$, with Cauchy-Riemann equations given in terms of

$$\Xi = \mu^2\partial_z + \partial_{\bar{z}} \quad \text{and} \quad \partial_{\bar{\mu}}.$$

Let $W = \mathbb{C} \times \mathbb{D}$ be equipped with the standard complex structure, given in terms of $\partial_{\bar{w}}$ and $\partial_{\bar{\mu}}$ for coordinates $(w, \mu) \in \mathbb{C} \times \mathbb{D}$. Then the map

$$(4.20) \quad \beta : Z \longrightarrow W, \quad (z, \mu) \longmapsto (z - \mu^2\bar{z}, \mu)$$

is holomorphic (in the sense that $\beta_*(\mathcal{D}) \subset \text{span}\{\partial_{\bar{w}}, \partial_{\bar{\mu}}\}$) and maps the interior of Z diffeomorphically onto the interior of W , with inverse given by

$$\beta^{-1}(w, \mu) = \left(\frac{w}{1 + |\mu|^2} + \frac{2\mu\Re(\bar{\mu}w)}{1 - |\mu|^4}, \mu \right), \quad (w, \mu) \in W^{\text{int}}.$$

Thus Z^{int} is biholomorphically equivalent to a polydisk in \mathbb{C}^2 and the degeneracy of the complex structure is encoded in the ‘blow down’ map β . More generally:

LEMMA 4.10. — *Suppose $M \subset \mathbb{R}^2$ is a Euclidean domain. Then the interior of its twistor space Z is a Stein surface that is biholomorphic to a domain in \mathbb{C}^2 .*

Proof. — The restriction of β from (4.20) to $Z^{\text{int}} = M \times \mathbb{D}^{\text{int}}$ gives the desired embedding as domain in \mathbb{C}^2 . In particular, the global holomorphic functions $\beta_1, \beta_2 \in \mathcal{O}(Z^{\text{int}})$ give a global coordinate system and separate points. To show that Z^{int} is a Stein surface, it thus remains to establish holomorphic convexity. To this end, define for $p = (z_*, \mu_*) \in \mathbb{R}^2 \times \mathbb{D} \setminus Z^{\text{int}}$ a function $f_p \in \mathcal{O}(Z^{\text{int}})$ by

$$f_p(z, \mu) = \begin{cases} (\mu - \mu_*)^{-1}, & |\mu_*| = 1, \\ ((z - z_*) - \mu^2(\bar{z} - \bar{z}_*))^{-1}, & |\mu_*| < 1, z_* \in \mathbb{R}^2 \setminus M. \end{cases}$$

Let $K \subset Z^{\text{int}}$ be compact and consider the holomorphic hull $\widehat{K} = \{(z, \mu) \in Z^{\text{int}} : |f(z, \mu)| \leq \sup_K |f| \text{ for all } f \in \mathcal{O}(Z^{\text{int}})\}$. If \widehat{K} was *not* compact, it would contain a sequence (z_n, μ_n) with limit point p as above, which leads to a contradiction, as f_p is unbounded along that sequence. Thus \widehat{K} is compact and, as K was arbitrary, the complex surface Z^{int} is holomorphically convex. \square

In fact, also the twistor space of a simple surface admits a natural, albeit less tractable, holomorphic map $\beta : Z \rightarrow \mathbb{C}^2$ as follows: Passing to global isothermal coordinates and with Λ as in (4.15), we may find a solution $u \in C^\infty(Z, \mathbb{C})$ to

$$\Xi u = \Lambda \quad \text{and} \quad \partial_{\bar{\mu}} u = 0 \quad \text{on } Z.$$

This follows from the existence of *scalar* holomorphic integrating factors on simple surfaces and is also a consequence of the vanishing result $H^1(Z, \mathcal{O}) \equiv H^1_{\bar{\partial}}(Z, [0]) = 0$ from Corollary 4.15 below. Further, by Corollary 4.7 there exists a function $\beta_1 \in \mathcal{O}(Z)$ with $\beta_1(z, 0) = z$ for all $z \in M$ and one checks that

$$\beta(z, \mu) = (\beta_1(z, \mu), e^{u(z, \mu)} \mu) \in \mathbb{C}^2$$

indeed defines a holomorphic map of similar form as (4.20) in the Euclidean case. While it would be interesting to know more about the behaviour of β (e.g., is it also diffeomorphism on the interior of Z ?), our approach to the transport Oka-Grauert principle does not require any such blow-down.

4.3. TRANSPORT OKA-GRAUERT PRINCIPLE. — We now define a ‘moduli space of holomorphic vector bundles’ over an open set $U \subset Z$. Noting that there are natural \mathbb{C}^n and $\mathbb{C}^{n \times n}$ -valued versions of the $\bar{\partial}$ -complex (4.4), we consider *partial connections* $A^{0,1} \in \Omega^{0,1}(U, \mathbb{C}^{n \times n})$ and maps

$$(4.21) \quad \Omega^0(U, \mathbb{C}^n) \xrightarrow{\bar{\partial} + A^{0,1}} \Omega^{0,1}(U, \mathbb{C}^n) \xrightarrow{\bar{\partial} + A^{0,1}} \Omega^{0,2}(Z, \mathbb{C}^n),$$

defined in the obvious way. If $A^{0,1} = (a_1, a_2) \in C^\infty(p^{-1}(U), \mathbb{C}^n)^2$ with $(\mathbb{V} + i)a_j = 0$ ($j = 1, 2$), a computation shows that the *curvature* of (4.21) equals

$$(4.22) \quad (\bar{\partial} + A^{0,1})^2 = (\omega^2 \eta_+ + \eta_-)a_2 - \partial_{\bar{\omega}} a_1 + [a_1, a_2] \in \Omega^{0,2}(U, \mathbb{C}^{n \times n}).$$

DEFINITION 4.11. — For $U \subset Z$ open we define the *moduli space*

$$\mathfrak{M}(U) = \mathfrak{M}_n(U) = \{A^{0,1} \in \Omega^{0,1}(U, \mathbb{C}^{n \times n}) : (\bar{\partial} + A^{0,1})^2 = 0\} / \sim,$$

where $A^{0,1} \sim B^{0,1}$ if and only if there exists $\varphi \in C^\infty(U, \text{GL}(n, \mathbb{C}))$ with

$$B^{0,1} = \varphi^{-1} \bar{\partial} \varphi + \varphi^{-1} A^{0,1} \varphi.$$

If $U \cap \partial Z = \emptyset$, such that U is a classical complex surface, then

$$(4.23) \quad \mathfrak{M}_n(U) \cong \begin{cases} \text{(topologically trivial) holomorphic vector bundles} \\ \text{of rank } n \text{ over } U \text{ up to isomorphism.} \end{cases}$$

Indeed, a representative $A^{0,1}$ of a class in $\mathfrak{M}_n(U)$ equips $U \times \mathbb{C}^n$ with the structure of a holomorphic vector bundle by declaring a local section $f : V \rightarrow \mathbb{C}^n$ (for $V \subset U$ open) to be holomorphic, if $(\bar{\partial} + A^{0,1})f = 0$; equivalent representatives yield isomorphic vector bundles (cf. Chapter 2.1.5 in [8]).

Recall from (1.3) that \mathcal{U} consists of attenuations $\mathbb{A} \in C^\infty(SM, \mathbb{C}^{n \times n})$ with Fourier coefficients $\mathbb{A}_k = 0$ for $k < -1$. The group \mathbb{G} from (1.2) acts on \mathcal{U} by (1.4) and we now establish a correspondence between the orbits of \mathbb{G} and elements in $\mathfrak{M} \equiv \mathfrak{M}_n(Z)$. Define a map $\mathcal{U} \rightarrow \mathfrak{M}$ as follows: For $\mathbb{A} \in \mathcal{U}$ let

$$(4.24) \quad A^{0,1}(x, v, \omega) := (a, 0) \equiv \left(\sum_{k \geq -1} \omega^{k+1} \mathbb{A}_k(x, v), 0 \right) \in \Omega^{0,1}(Z, \mathbb{C}^{n \times n}),$$

noting that $A^{0,1}$ lies in $\Omega^{0,1}(Z, \mathbb{C}^{n \times n})$ and satisfies $(\bar{\partial} + A^{0,1})^2 = 0$ in view of (4.22) and Remark 4.5. We then map \mathbb{A} to the equivalence class $[A^{0,1}] \in \mathfrak{M}$.

PROPOSITION 4.12 (Twistor correspondence B). — *The map $\mathcal{U} \rightarrow \mathfrak{M}$, $\mathbb{A} \mapsto [A^{0,1}]$ is \mathbb{G} -invariant and descends to an injective map $\mathcal{U}/\mathbb{G} \rightarrow \mathfrak{M}$. If (M, g) is diffeomorphic to a disk, the induced map is also surjective, such that*

$$(4.25) \quad \mathcal{U}/\mathbb{G} \cong \mathfrak{M}.$$

In fact, the isomorphism (4.25) holds true for any orientable surface (M, g) . For a proof of this more general statement we refer to Remark 4.4.13 in [4].

Proof. — Suppose that $\mathbb{A}, \mathbb{B} \in \mathcal{U}$, let $A^{0,1} = (a, 0)$ as in (4.24) and define $B^{0,1} = (b, 0)$ analogously. To demonstrate \mathbb{G} -invariance, we assume that $\mathbb{A} \triangleleft F = \mathbb{B}$ for some $F \in \mathbb{G}$ and consider the function

$$(4.26) \quad \phi(x, v, \omega) = \sum_{k \geq 0} \omega^k F_k(x, v),$$

which is smooth and S^1 -invariant by Remark 4.5; further it satisfies $\det \phi(x, v, \omega) \neq 0$ for all $\omega \in \mathbb{D}$. To see this, we define ψ as in (4.26), but with F replaced by F^{-1} . Then $h = \phi\psi - \text{Id}$ satisfies $h(x, v, 1) = 0$ by construction, and, due to S^1 -invariance,

$$h = 0 \text{ on } \{|\omega| = 1\} \quad \text{and} \quad \partial_{\bar{\omega}} h = 0 \text{ on } \{|\omega| \leq 1\}.$$

By the maximum principle, $h \equiv 0$, which means that ψ is an inverse for ϕ .

We claim that $A^{0,1}$ is equivalent to $B^{0,1}$ via the gauge $\varphi = p_*\phi \in C^\infty(Z, \text{GL}(n, \mathbb{C}))$, which is to say that

$$(4.27) \quad \varphi^{-1}(\bar{\partial} + A^{0,1})\varphi \equiv (\phi^{-1}(\omega^2\eta_+ + \eta_- + a)\phi, \phi^{-1}\partial_{\bar{\omega}}\phi) = (b, 0) \in \Omega^{0,1}(Z, \mathbb{C}^{n \times n}).$$

Evidently $\partial_{\bar{\omega}}\phi = 0$, so it remains to show that the function $g = \phi^{-1}(\omega^2\eta_+ + \eta_- + a)\phi - b$ vanishes identically. To see this, note that $g(x, v, 1) = \mathbb{A} \triangleleft F(x, v) - \mathbb{B}(x, v) = 0$. Moreover $(\mathbb{V} + i)g = 0$ which means that g only changes phase along the flow of \mathbb{V} and thus

$$g = 0 \text{ on } \{|\omega| = 1\} \quad \text{and} \quad \partial_{\bar{\omega}}g = 0 \text{ on } \{|\omega| < 1\},$$

where holomorphicity in ω is easily checked. In particular the maximum modulus principle on \mathbb{D} applies to yield $g \equiv 0$.

To show that the induced map $\mathcal{U}/\mathbb{G} \rightarrow \mathfrak{M}$ is injective, we assume that $[A^{0,1}] = [B^{0,1}]$. This means that (4.27) holds true, where now ϕ is defined as $p^*\varphi$ for an appropriate gauge $\varphi \in C^\infty(Z, \text{GL}(n, \mathbb{C}))$. In particular ϕ is holomorphic in ω and thus admits a series expansion as in (4.26) with coefficients $F_k(x, v) \in \mathbb{C}^{n \times n}$ ($(x, v) \in SM$). Similar to the proof of Proposition 4.4 one checks that $F_k \in \Omega_k$ such that

$$F(x, v) = \phi(x, v, 1)$$

defines a smooth, $\text{GL}(n, \mathbb{C})$ -valued map on SM with both F and F^{-1} being fibrewise holomorphic. Further, evaluating (4.27) at $\omega = 1$ yields $\mathbb{A} \triangleleft F = \mathbb{B}$, as desired.

To establish surjectivity of $\mathcal{U}/\mathbb{G} \rightarrow \mathfrak{M}$ we need to show that each class in \mathfrak{M} admits a representative $A^{0,1} = (a_1, a_2) \in \Omega^{0,1}(Z, \mathbb{C}^{n \times n})$ with $a_2 \equiv 0$. Indeed, in that case $a_1 \in C^\infty(SM \times \mathbb{D}, \mathbb{C}^{n \times n})$ satisfies $\partial_{\bar{\omega}}a_1 = 0$ by the curvature condition, hence

$$a_1(x, v, \omega) = \sum_{k \geq -1} \omega^{k+1} \mathbb{A}_k(x, v)$$

for coefficients $\mathbb{A}_k(x, v)$, which can be seen to lie in Ω_k as in the proof of Proposition 4.4; in particular $\mathbb{A}(x, v) = a_1(x, v, 1)$ is a preimage of $[A^{0,1}]$.

We now make use of the fact that M is diffeomorphic to the disk \mathbb{D} , such that global isothermal coordinates become available. Using the description from Section 4.2, the twistor space is then given by $Z = \mathbb{D} \times \mathbb{D}$ and a representative of a class in \mathfrak{M} is a tuple $(b_1, b_2) \in C^\infty(Z, \mathbb{C}^{n \times n})^2$ obeying the curvature condition $(\Xi + \Lambda)b_2 - \partial_{\bar{\mu}}b_1 + [b_1, b_2] = 0$,

where $\Xi, \partial_{\bar{\mu}}$ are as in (4.14) and Λ is as in (4.15). Then by the Oka-Grauert principle on the μ -disk (Lemma 6.2) there exists a solution $\varphi \in C^\infty(Z, \text{GL}(n, \mathbb{C}))$ of

$$\partial_{\bar{\mu}}\varphi - \varphi b_2 = 0,$$

which means that $(a_1, a_2) := \varphi(\Xi + b_1, \partial_{\bar{\mu}} + b_2)\varphi^{-1} \equiv (\varphi\Xi\varphi^{-1} + \varphi b_1\varphi^{-1}, 0)$ defines an equivalent representative with $a_2 \equiv 0$, as desired. \square

In view of the preceding correspondence principle, Theorem 1.3 can be reformulated as:

THEOREM 4.13 (Transport Oka-Grauert principle). — *Suppose Z is the twistor space of a simple surface (M, g) . Then $\mathfrak{M} = \mathfrak{M}_n(Z) = 0$ for all $n \in \mathbb{N}$.* \square

Note that our proof does *not* rely on an integrability theorem as in [8, Th. 2.1.53] – that is, we do not first establish the existence of local holomorphic frames, which are then ‘glued’ by means of a Cartan lemma as in the proof of the standard Oka-Grauert principle. A local integrability theorem – for general twistor spaces – follows a posteriori:

COROLLARY 4.14 (Local integrability). — *Let Z be the twistor space of an arbitrary oriented Riemannian surface. Consider a class $[A^{0,1}] \in \mathfrak{M}_n(Z)$ and a point $p = (x, v) \in Z$ with $x \in M^{\text{int}}$. Then there exists an open neighbourhood U of p and a gauge $\varphi \in C^\infty(U, \text{GL}(n, \mathbb{C}))$ with $(\bar{\partial} + A^{0,1})\varphi = 0$.*

Proof. — There exists a simple surface $M_1 \subset M^{\text{int}}$ containing x in its interior – its twistor space Z_1 is then a subset of Z . By Theorem 4.13, we have $[A^{0,1}|_{Z_1}] = 0 \in \mathfrak{M}_n(Z_1)$ and thus the corollary follows with $U = Z_1|_{M_1^{\text{int}}}$. \square

A further consequence is the following ‘vanishing theorem’ in the spirit of Cartan’s Theorem B. In fact, this is a reformulation of the linear result in Proposition 3.1 (modulo tame estimates) and thus does not require an inverse function theorem.

COROLLARY 4.15. — *Suppose Z is the twistor space of a simple surface (M, g) . Then for $[A^{0,1}] \in \mathfrak{M}$ the 1st cohomology of the twisted $\bar{\partial}$ -complex (4.21) vanishes, i.e.,*

$$(4.28) \quad H_{\bar{\partial}}^1(Z, [A^{0,1}]) \equiv \frac{\ker((\bar{\partial} + A^{0,1})|_{\Omega^{0,1}(Z)})}{\text{im}((\bar{\partial} + A^{0,1})|_{\Omega^0(Z)})} = 0.$$

Proof. — We give a brief sketch: It is straightforward to see that (4.28) is gauge-invariant, so by Proposition 4.12 we may assume that $A^{0,1} = (a, 0)$, where a is as in (4.24) for an attenuation $\mathbb{A} \in \mathcal{U}$. Next, using solvability of the $\partial_{\bar{\omega}}$ -equation, any cohomology class may be represented by a tuple $(h_1, 0) \in \Omega^{0,1}(SM, \mathbb{C}^n)$. Via Remark 4.5 the function h_1 gives rise to a element in $\bigoplus_{k \geq -1} \Omega_k$ and we are in the setting of Proposition 3.1. From this, one deduces that there is a solution $h \in \Omega^0(Z, \mathbb{C}^n)$ to $(\bar{\partial} + A^{0,1})h = (h_1, 0)$ – first for $\omega = 1$, then for all $\omega \in \mathbb{D}$ by invariance and the maximum principle. \square

4.4. DISCUSSION OF RELATED WORK. — The twistor space Z considered in this article is closely related to the more classical twistor notion from [9, 32], used recently, e.g., in the context of projective structures [26, 27]. To explain this relation, we first note that the constructions from Section 4.1 can be carried out in greater generality by substituting the vector fields \mathbb{V} and $\xi = \omega^2\eta_+ + \eta_-$ by

$$\mathbb{V}_n = V + in(\bar{\omega}\partial_{\bar{\omega}} - \omega\partial_{\omega}), \quad \text{and} \quad \xi_k = \begin{cases} \omega^k\eta_+ + \eta_- & k \geq 0, \\ \bar{\omega}^{-k}\eta_+ + \eta_- & k < 0, \end{cases}$$

for $n, k \in \mathbb{Z}$, respectively. A computation similar to Lemma 4.1 shows that $[\xi_k, \mathbb{V}_n] = i\xi_k$ if and only if $nk = 2$, such that we obtain four twistor spaces

$$Z(n) = ((SM \times \mathbb{D})/S^1, \mathcal{D}_n = p_* \text{span}\{\xi_{2/n}, \partial_{\bar{\omega}}\}), \quad n \in \{\pm 1, \pm 2\},$$

where the quotient is taken with respect to the flow of \mathbb{V}_n . Then $(x, v, \omega) \mapsto (x, v, \omega^2)$ induces holomorphic maps $Z(\pm 1) \rightarrow Z(\pm 2)$ and, in particular, $Z \equiv Z(1)$ may be viewed as branched double cover of the space $Z(2)$ (branched double covers were also found useful in [21]).

We claim that the interior of $Z(2)$ is precisely the twistor space considered in the articles mentioned above. Using the description in [27], this means that there is a biholomorphic map

$$(4.29) \quad \mathcal{F} : Z(2)^{\text{int}} \xrightarrow{\sim} P/\text{CO}(2),$$

where P is the oriented frame bundle and $\text{CO}(2)$ is the group of dilations and rotations of \mathbb{R}^2 . Here we assume that the projective class \mathfrak{p} that is used to define the complex structure on $P/\text{CO}(2)$, as explained in [27, §4.1], is given by $\mathfrak{p} = [\nabla^g]$ for the Levi-Civita connection ∇^g of (M, g) . To construct \mathcal{F} , consider

$$SM \times \mathbb{D}^{\text{int}} \xrightarrow{\mathcal{R}} P \times \mathbb{D}^{\text{int}} \xrightarrow{\Upsilon} P,$$

with Υ as in [27, §4.1] and $\mathcal{R}(x, v, \omega) = (x, f_v, -\omega)$, where $f_v = (v, v^\perp) \in P_x$. Here, v^\perp is the rotation of v by $\pi/2$, counterclockwise with respect to the orientation of M . As below equation (4.6) in [27], one checks that

$$\mathcal{R}((x, v, \omega) \triangleleft e^{it}) = ((x, f_v) \triangleleft R_t, R_t^{-1} \triangleright (-\omega)) = \mathcal{R}(x, v, \omega) \triangleleft R_t,$$

where $R_t \in \text{GL}^+(2, \mathbb{R})$ is the rotation matrix corresponding to e^{it} . This shows that \mathcal{R} induces a smooth map between the quotient spaces $Z(2)^{\text{int}} \equiv (SM \times \mathbb{D}^{\text{int}})/S^1$ and $P \times_{\text{GL}^+(2, \mathbb{R})} \mathbb{D}^{\text{int}}$. Also Υ descends to quotient spaces and thus $\mathcal{F}([(x, v, \omega)]) = [\Upsilon \circ \mathcal{R}(x, v, \omega)]$ defines a smooth map \mathcal{F} as in (4.29). To check that \mathcal{F} is also holomorphic one can use the description of $(1, 0)$ -forms on P from [27], in particular the computation of their pull-backs by Υ in (4.3) and (4.5). Pulling these back by \mathcal{R} one obtains (nonzero multiples) of the following 1-forms on $SM \times \mathbb{D}^{\text{int}}$:

$$\eta_+^\vee - \omega\eta_-^\vee \quad \text{and} \quad d\omega + 2i\omega V^\vee.$$

After complex conjugation these equal precisely the 1-forms in Lemma 4.2 (with an additional factor 2 which is due to the choice of $n = 2$ here). This shows that \mathcal{F} is a

holomorphic immersion. It is easily checked that \mathcal{F} is bijective, so overall we obtain an isomorphism as in (4.29).

Next, we briefly discuss the work of Eskin and Ralston [11], who proved a version of Theorem 4.13 in a Euclidean setting. They establish the existence of gauges φ – that is, $\mathrm{GL}(n, \mathbb{C})$ -valued solutions to $(\bar{\partial} + A^{0,1})\varphi = 0$ – that are smooth in Z^{int} and have a *continuous* extension to ∂Z . We give a brief outline of their argument in the language developed above. Recall from Section 4.2 that the twistor space of \mathbb{R}^2 admits a ‘blow down’ map $\beta : Z \rightarrow W$ into a (closed) polydisk. The punchline of [11] is that pull back by β gives a surjective map

$$(4.30) \quad \beta^* : \widetilde{\mathfrak{M}}(W) \longrightarrow \widetilde{\mathfrak{M}}(Z)$$

between appropriate moduli spaces containing ‘holomorphic vector bundles’ which are trivial away from a compact set and have a *continuous* extension to the boundary. The result then follows from the classical Oka-Grauert principle on W – in a version with continuous boundary values (cf. [22, Th. 10.1]) – which implies $\widetilde{\mathfrak{M}}(W) = 0$. Their approach thus parallels the desingularisation by means of a blow down in [20].

In order to establish surjectivity of β^* as in (4.30), the authors prove a local integrability result as in Corollary 4.14 a priori (using the inverse function theorem in a Hölder space, where no loss of derivatives occurs) and then glue local solutions to $(\bar{\partial} + A^{0,1})\varphi = 0$ by means of an appropriate Cartan lemma. The crucial step lies in showing that by such a gluing procedure one can arrange all transition functions to be of the form $h = \beta^*g$ for locally defined functions g on W (cf. equation (12) in [11]) – this is quite delicate and encompasses removing singularities at $\beta(\partial Z)$ using methods from complex analysis.

5. RANGE CHARACTERISATIONS

We now turn to the range characterisation for the non-Abelian X-ray transform, starting with some general considerations that hold on any non-trapping surface (M, g) with strictly convex boundary. Define

$$(5.1) \quad B : C^\infty(\partial SM, \mathrm{GL}(n, \mathbb{C})) \longrightarrow C^\infty(\partial_+ SM, \mathrm{GL}(n, \mathbb{C})), \quad f \longmapsto f(f^{-1} \circ \alpha)|_{\partial_+ SM},$$

where $\alpha : \partial SM \rightarrow \partial SM$ is the scattering relation of (M, g) (see Section 1.3). To motivate the range characterisations in this section, consider the following diagram:

$$(5.2) \quad \begin{array}{ccc} C_{\mathrm{Id}}^\infty(SM, \mathrm{GL}(n, \mathbb{C})) & \xrightarrow{\mathbb{I}^*} & C^\infty(SM, \mathbb{C}^{n \times n}) \\ (\cdot)|_{\partial SM} \downarrow & \searrow & \downarrow \mathbb{A} \mapsto C_{\mathbb{A}} \\ C_{\mathrm{Id}}^\infty(\partial SM, \mathrm{GL}(n, \mathbb{C})) & \xrightarrow{B} & C^\infty(\partial_+ SM, \mathrm{GL}(n, \mathbb{C})). \end{array}$$

Here and below, double-headed arrows stand for surjections; further $C_{\mathrm{Id}}^\infty(\cdot, \mathrm{GL}(n, \mathbb{C}))$ is the space of maps which are homotopic to Id and we define

$$\mathbb{I}^*(R) = -(XR)R^{-1}.$$

The map \mathbb{I}^* , while not being an adjoint in any natural way, serves a similar purpose as I^* in the linear theory. This is illustrated by the following lemma and further substantiated in Section 5.2, where surjectivity results for \mathbb{I}^* in different settings are derived using Theorem 1.3 on simple surfaces.

LEMMA 5.1. — *The diagram (5.2) commutes and the map $\mathbb{I}^* : C_{\text{Id}}^\infty(SM, \text{GL}(n, \mathbb{C})) \rightarrow C^\infty(SM, \mathbb{C}^{n \times n})$ is surjective.*

Proof. — We assume that the diagonal arrow is $B(\cdot|_{\partial SM})$ such that the lower triangle commutes. To check that the upper triangle commutes, let $F \in C_{\text{Id}}^\infty(SM, \text{GL}(n, \mathbb{C}))$ and denote $\mathbb{A} = \mathbb{I}^*(F)$. Let $G : SM \rightarrow \text{GL}(n, \mathbb{C})$ be the unique continuous solution (differentiable along the geodesic flow) of the transport problem

$$XG = 0 \quad \text{and} \quad G = F^{-1} \text{ on } \partial_- SM.$$

Then $R = FG$ satisfies $(X + \mathbb{A})R = 0$ and $R = \text{Id}$ on $\partial_- SM$. In particular, using that $G|_{\partial SM}$ is α -invariant, we have

$$C_{\mathbb{A}} = R|_{\partial_+ SM} = B(R|_{\partial SM}) = FG(G^{-1} \circ \alpha)(F^{-1} \circ \alpha)|_{\partial_+ SM} = B(F|_{\partial SM}).$$

To check that \mathbb{I}^* is onto, we have to show that any $\mathbb{A} \in C^\infty(SM, \mathbb{C}^{n \times n})$ admits a smooth, contractible integrating factor. To this end, embed (M, g) into a closed manifold (N, g) and extend \mathbb{A} smoothly to N . Then there is a smooth cocycle $C : SN \times \mathbb{R} \rightarrow \text{GL}(n, \mathbb{C})$ associated to \mathbb{A} , uniquely defined by

$$\partial_t C(x, v, t) + \mathbb{A}C(x, v, t) = 0 \text{ on } SN \times \mathbb{R} \quad \text{and} \quad C(x, v, 0) = \text{Id on } SN.$$

Let $M_0 \subset N$ be a non-trapping surface with strictly convex boundary, containing M in its interior. Let τ_0 be the exit time of M_0 (which is smooth on SM) and define

$$R_s(x, v) = [C(x, v, s\tau_0(x, v))]^{-1}, \quad 0 \leq s \leq 1, (x, v) \in SM.$$

Then $R_1 \in C^\infty(SM, \text{GL}(n, \mathbb{C}))$ is a smooth integrating factor for \mathbb{A} (cf. Lemma 5.3.2 in [36]) and sending $s \rightarrow 0$ provides a homotopy with Id , as desired. \square

Using the preceding lemma, a simple diagram chase in (5.2) reveals a first range characterisation:

PROPOSITION 5.2. — *An element $q \in C^\infty(\partial_+ SM, \text{GL}(n, \mathbb{C}))$ is given as scattering data $q = C_{\mathbb{A}}$ of a general attenuation $\mathbb{A} \in C^\infty(SM, \mathbb{C}^{n \times n})$ if and only if $q = Bf$ for some $f \in C_{\text{Id}}^\infty(\partial SM, \text{GL}(n, \mathbb{C}))$.* \square

In the remaining section we give similar characterisations, when \mathbb{A} is restricted to certain subclasses $\mathcal{A} \subset C^\infty(SM, \mathbb{C}^{n \times n})$ of attenuations. This involves finding appropriate domains \mathcal{D} and boundary spaces \mathcal{B} for which there is a diagram

$$\begin{array}{ccc} \mathcal{D} & \xrightarrow{\mathbb{I}^*} & \mathcal{A} \\ \downarrow & \searrow & \downarrow \mathbb{A} \mapsto C_{\mathbb{A}} \\ \mathcal{B} & \xrightarrow{P} & C^\infty(\partial_+ SM, \text{GL}(n, \mathbb{C})), \end{array}$$

where P is an appropriate *boundary operator* and the arrows emerging from \mathcal{D} are surjective. If such a diagram commutes (up to gauge), then the range of $\mathcal{A} \ni \mathbb{A} \mapsto C_{\mathbb{A}}$ equals that of P (up to gauge).

5.1. NONLINEAR HILBERT TRANSFORMS. — As building block for the boundary operators considered below we introduce here a nonlinear operator

$$(5.3) \quad \mathcal{H} : C_{\bullet}^{\infty}(\partial SM, \mathrm{GL}(n, \mathbb{C})) \longrightarrow C_{\bullet}^{\infty}(\partial SM, \mathrm{U}(n))$$

which is based on the factorisation theorems discussed in Section 2.1 (see Remark 2.4 for the \bullet -notation) and serves as analogue of the Hilbert transform in the linear theory. Upon choosing a section $\mathbf{1} : M \rightarrow SM$ (or equivalently, fixing a trivialisation of SM) we define $\mathcal{H} \equiv \mathcal{H}_{\mathbf{1}}$ by

$$\mathcal{H}(r) = u^*,$$

where $r = uf$ is a decomposition as in (2.4), normalised such that $u(x, \mathbf{1}(x)) = \mathrm{Id}$. In reference to \mathcal{H} we also define

$$(5.4) \quad \begin{aligned} \mathcal{H}^* : C_{\bullet}^{\infty}(\partial SM, \mathrm{GL}(n, \mathbb{C})) &\longrightarrow C_{\bullet}^{\infty}(\partial SM, \mathrm{U}(n)), & \mathcal{H}^*(r) &= \mathcal{H}(r^{-1}), \\ \mathcal{H}^+ : C^{\infty}(\partial SM, \mathrm{Her}_n^+) &\longrightarrow C_{\mathrm{Id}}^{\infty}(\partial SM, \mathrm{GL}(n, \mathbb{C})), & \mathcal{H}^+(r) &= \mathcal{H}(r^{1/2})r^{1/2}. \end{aligned}$$

Both transforms can be described in terms of suitable decompositions: Indeed, we have $\mathcal{H}^*(r) = u$, where $r = fu$ as in (2.5), normalised such that $u(x, \mathbf{1}(x)) = \mathrm{Id}$ and $\mathcal{H}^+(r) = f$, where $r = f^*f$ is a Birkhoff factorisation as in (2.6), with normalisation inherited from \mathcal{H} .

We introduce two further types of ‘nonlinear Hilbert transforms’, which do not depend on a choice of $\mathbf{1}$, but are only available if r admits a special decomposition. To this end, define spaces

$$(5.5) \quad C_0^{\infty}(\partial SM, \mathrm{GL}(n, \mathbb{C})) \quad \text{and} \quad C_0^{\infty}(\partial SM, \mathrm{Her}_n^+),$$

as follows: An element $r \in C_{\bullet}^{\infty}(\partial SM, \mathrm{GL}(n, \mathbb{C}))$ lies in the left space, if it admits a (necessarily unique) decomposition $r = uf$ as in (2.4) with $f_0 = \mathrm{Id}$ – we then write $\mathcal{H}^0(r) = u^*$. Further, $r \in C^{\infty}(\partial SM, \mathrm{Her}_n^+)$ lies in the right space in (5.5), if $r^{1/2} \in C_0^{\infty}(\partial SM, \mathrm{GL}(n, \mathbb{C}))$ and we set $\mathcal{H}^{+,0}(r) = \mathcal{H}^0(r^{1/2})r^{1/2}$. We obtain transforms:

$$(5.6) \quad \mathcal{H}^0 : C_0^{\infty}(\partial SM, \mathrm{GL}(n, \mathbb{C})) \longrightarrow C_{\bullet}^{\infty}(\partial SM, \mathrm{U}(n)),$$

$$(5.7) \quad \mathcal{H}^{+,0} : C_0^{\infty}(\partial SM, \mathrm{Her}_n^+) \longrightarrow C_{\mathrm{Id}}^{\infty}(\partial SM, \mathrm{GL}(n, \mathbb{C})).$$

Next, we consider the space $C_{\Delta}^{\infty}(\partial SM, \mathrm{GL}(n, \mathbb{C}))$ consisting of those maps

$$r \in C_{\mathrm{Id}}^{\infty}(\partial SM, \mathrm{GL}(n, \mathbb{C}))$$

which factor uniquely as $r = gf$, where $f, g^* \in \mathbb{H}$ (with \mathbb{H} as in Remark 2.4) and $g_0 = \mathrm{Id}$. With respect to this factorisation we define

$$(5.8) \quad \mathcal{H}^{\Delta} : C_{\Delta}^{\infty}(\partial SM, \mathrm{GL}(n, \mathbb{C})) \longrightarrow \mathbb{G}, \quad \mathcal{H}^{\Delta}(r) = f.$$

We discuss the relation of \mathcal{H}^{Δ} to Birkhoff factorisations below Theorem 5.13.

EXAMPLE 5.3. — We consider the ‘nonlinear Hilbert transforms’ from above for $n = 1$. A general element in $C^\infty(\partial SM, \text{GL}(1, \mathbb{C}))$ has the form

$$r = e^{ik\theta} e^{\psi+i\sigma}, \quad \text{where } k \in \mathbb{Z}, \psi, \sigma \in C^\infty(\partial SM, \mathbb{R}).$$

Let $\psi = \psi_{<0} + \psi_0 + \psi_{>0}$ be the decomposition into negative, zero and positive Fourier modes. Then the standard, linear Hilbert transform of ψ is $H\psi = (\psi_{>0} - \psi_{<0})/i$ and thus $\psi = -iH\psi + \psi_0 + 2\psi_{>0}$, which implies that

$$r = (e^{ik\theta} e^{-iH\psi+i\sigma}) \times (e^{\psi_0+2\psi_{>0}}) =: uf$$

is a decomposition as in (2.4) (not necessarily normalised). Then

$$\mathcal{H}(r) = we^{iH\psi} e^{-i\sigma-ik\theta} \quad \text{and} \quad \mathcal{H}^*(r) = we^{-iH\psi} e^{i\sigma+ik\theta},$$

where $w = w_{\mathbf{1}} \in C^\infty(M, \text{U}(n))$ is chosen to achieve the correct normalisation. If r takes values in $\mathbb{R}_{>0} \equiv \text{Her}_1^+$ (such that $k = 0$ and $\sigma = 0$) we see that \mathcal{H} and \mathcal{H}^* are exponentiated linear Hilbert transforms; further

$$\mathcal{H}^+(r) = we^{\frac{1}{2}(\psi+iH\psi)}.$$

Finally, r is in the domain of \mathcal{H}^0 and $\mathcal{H}^{+,0}$ iff $\psi_0 = 0$ and it is in the domain of \mathcal{H}^Δ iff $k = 0$, in which case $\mathcal{H}^\Delta(r) = e^{\psi_0+\psi_{>0}+i\sigma_0+i\sigma_{>0}}$.

5.2. RANGE FOR $\mathfrak{u}(n)$ -ATTENUATIONS. — It is instructive to first consider the non-Abelian X-ray transform on the space \mathcal{U} from (1.3). In terms of the right action of \mathbb{G} on \mathcal{U} (defined in (1.4)), we have $\mathbb{I}^*(F) = 0 \triangleleft F^{-1}$ and hence \mathbb{I}^* fits into an exact sequence (of pointed sets):

$$(5.9) \quad 0 \longrightarrow \mathbb{G}_0 \hookrightarrow \mathbb{G} \xrightarrow{\mathbb{I}^*} \mathcal{U} \longrightarrow \mathfrak{M} \longrightarrow 0.$$

Here \mathbb{G}_0 is the stabiliser of $0 \in \mathcal{U}$ and – for the purpose of this section – we think of \mathfrak{M} as quotient space \mathcal{U}/\mathbb{G} , such that exactness in (5.9) is evident. The identification $\mathcal{U}/\mathbb{G} = \mathfrak{M}$ is justified by Proposition 4.12.

Let us assume now that \mathfrak{M} is trivial – by Theorem 1.3 this holds in particular if (M, g) is simple. Then $\mathbb{I}^* : \mathbb{G} \rightarrow \mathcal{U}$ is surjective and we have a commutative diagram

$$(5.10) \quad \begin{array}{ccc} \mathbb{G} & \xrightarrow{\mathbb{I}^*} & \mathcal{U} \\ (\cdot)|_{\partial SM} \downarrow & \searrow & \downarrow \mathbb{A} \mapsto C_{\mathbb{A}} \\ \mathbb{H} & \xrightarrow{B} & C^\infty(\partial_+ SM, \text{GL}(n, \mathbb{C})), \end{array}$$

where $\mathbb{H} = \{f = F|_{\partial SM} : F \in \mathbb{G}\}$ (see also Remark 2.4). Note that \mathbb{H} is a genuine boundary space in that it has the intrinsic characterisation

$$\mathbb{H} = \{f \in C^\infty_{\text{id}}(\partial SM, \text{GL}(n, \mathbb{C})) : f \text{ is fibrewise holomorphic}\}.$$

We thus obtain the following range characterisation:

PROPOSITION 5.4. — *Suppose $\mathfrak{M} = 0$. Then an element $q \in C^\infty(\partial_+ SM, \text{GL}(n, \mathbb{C}))$ lies in the range of $\mathcal{U} \ni \mathbb{A} \mapsto C_{\mathbb{A}}$ if and only if $q = Bf$ for some $f \in \mathbb{H}$. \square*

Let us now consider attenuations in one of the three classes

$$(5.11) \quad \mathcal{U}(\mathbf{u}(n)) = \{\mathbf{u}(n)\text{-pairs}\}, \quad C^\infty(M, \mathbf{u}(n)) \quad \text{and} \quad \Omega^1(M, \mathbf{u}(n)),$$

all considered as subsets of \mathcal{U} . Note that $\mathcal{U}(\mathbf{u}(n)) = \mathcal{U} \cap C^\infty(SM, \mathbf{u}(n))$, due to the identity $\mathbb{A}_{-k} = -\mathbb{A}_k^*$ ($k \in \mathbb{Z}$) for skew-Hermitian attenuations. The second and third space in (5.11) consist of skew-Hermitian matrix fields Φ and connections A , respectively.

Define

$$\sqrt{\mathbb{G}_0} := \{F \in \mathbb{G} : X(F^*F) = 0\} \subset C^\infty(SM, \mathbb{C}^{n \times n}),$$

and recall that a function F on SM (or ∂SM) is even, if it only has even Fourier modes or equivalently if it obeys the symmetry condition $F(x, v) = F(x, -v)$.

PROPOSITION 5.5. — *Suppose that (M, g) is simple. Then \mathbb{I}^* is well defined and surjective in the following settings:*

- (i) $\mathbb{I}^* : \sqrt{\mathbb{G}_0} \rightarrow \mathcal{U}(\mathbf{u}(n))$,
- (ii) $\mathbb{I}^* : \{F \in \sqrt{\mathbb{G}_0} : F \text{ even}\} \rightarrow \Omega^1(M, \mathbf{u}(n))$,
- (iii) $\mathbb{I}^* : \{F \in \sqrt{\mathbb{G}_0} : F_0 = \text{Id}\} \rightarrow C^\infty(M, \mathbf{u}(n))$.

Part (i) holds in greater generality and is in fact equivalent to the assertion that $\mathfrak{M} = 0$ (assuming that (M, g) is non-trapping and has a strictly convex boundary).

Proof. — For (i) let $F \in \sqrt{\mathbb{G}_0}$, then

$$0 = X(F^*F) = (XF^*)F + F^*(XF) = F^* ((F^{-1})^*(XF^*) - \mathbb{I}^*(F)) F,$$

and hence $(\mathbb{I}^*(F))^* = ((F^{-1})^*(XF^*))^* = -\mathbb{I}^*(F)$, which shows that \mathbb{I}^* indeed maps $\sqrt{\mathbb{G}_0}$ into $\mathcal{U}(\mathbf{u}(n))$. To see that it is onto, let $\mathbb{A} \in \mathcal{U}(\mathbf{u}(n))$ and let $F \in \mathbb{G}$ be an arbitrary HIF for \mathbb{A} . Then, as \mathbb{A} is skew-Hermitian,

$$X(F^*F) = (-\mathbb{A}F)^*F + F^*XF = F^*(\mathbb{A}F + XF) = 0,$$

which means that $F \in \sqrt{\mathbb{G}_0}$.

For (ii) let $F \in \sqrt{\mathbb{G}_0}$ be even. Then $\mathbb{I}^*(F)$ is odd and, being skew-Hermitian, only has Fourier modes in degree ± 1 ; hence $\mathbb{I}^*(F) \in \Omega^1(M, \mathbf{u}(n))$. Conversely, if A is a $\mathbf{u}(n)$ -connection, then by Proposition 3.5 there exists an even HIF $F \in \mathbb{G}$ and we must have $F \in \sqrt{\mathbb{G}_0}$, as above.

Finally, for (iii), let $F \in \sqrt{\mathbb{G}_0}$ with $F_0 = \text{Id}$. Then $\mathbb{A} = \mathbb{I}^*(F)$ satisfies

$$(5.12) \quad \mathbb{A}_{-1} = -(\eta_- F_0)(F^{-1})_0 = 0$$

and, \mathbb{A} being skew-Hermitian, also $\mathbb{A}_1 = -\mathbb{A}_{-1}^* = 0$. Hence $\mathbb{A} \in \Omega_0(SM, \mathbf{u}(n)) \equiv C^\infty(M, \mathbf{u}(n))$. Conversely, given $\Phi \in C^\infty(M, \mathbf{u}(n))$ let $\tilde{F} \in \sqrt{\mathbb{G}_0}$ be any HIF. Then, similar to (5.12) we see that $\eta_- \tilde{F}_0 = 0$ and as \tilde{F}_0 is invertible, also $\eta_- \tilde{F}_0^{-1} = 0$. By Theorem 13.11.6 in [36] there exists $G \in \mathbb{G}_0$ with $G_0 = \tilde{F}_0^{-1}$, hence $F = \tilde{F}G \in \sqrt{\mathbb{G}_0}$ is a new HIF for Φ with $F_0 = \tilde{F}_0 G_0 = \text{Id}$, as desired. \square

The domains in the preceding proposition can be projected onto the following boundary spaces of Her_n^+ -valued functions:

$$(5.13) \quad C_\alpha^\infty(\partial_+ SM, \text{Her}_n^+) = \{w \in C^\infty(\partial_+ SM, \text{Her}_n^+) : A_+ w \text{ smooth on } \partial SM\},$$

$$(5.14) \quad C_{\alpha,1}^\infty(\partial_+ SM, \text{Her}_n^+) = \{w \in C_\alpha^\infty(\partial_+ SM, \text{Her}_n^+) : w \circ \alpha_a = w\},$$

$$(5.15) \quad C_{\alpha,0}^\infty(\partial_+ SM, \text{Her}_n^+) = \left\{ w \in C_\alpha^\infty(\partial_+ SM, \text{Her}_n^+) : \begin{array}{l} w^\sharp = F^*F \text{ for some } \\ F \in \mathbb{G} \text{ with } F_0 = \text{Id} \end{array} \right\}.$$

Here

$$(5.16) \quad A_+ w(x, v) = \begin{cases} w(x, v) & (x, v) \in \partial_+ SM, \\ w \circ \alpha(x, v) & (x, v) \in \partial_- SM, \end{cases} \quad \text{and} \quad \alpha_a(x, v) = \alpha(x, -v),$$

and w^\sharp is the unique solution to $Xw^\sharp = 0$ on SM and $w^\sharp = w$ on $\partial_+ SM$. By a classical result of Pestov and Uhlmann (see [39] or Theorem 5.1.1 in [36]), the first integral w^\sharp is smooth on SM for $w \in C_\alpha^\infty(SM, \text{Her}_n^+)$.

Note that (5.13) and (5.14) define genuine boundary spaces in the sense that membership can be checked on ∂SM only in terms of the scattering relation α . To check whether a function w belongs to $C_{\alpha,0}^\infty$ one first has to find the first integral w^\sharp .

The following result is a consequence of Birkhoff's factorisation theorem for Hermitian matrices (Theorem 2.3) and does *not* require (M, g) to be simple.

PROPOSITION 5.6. — *Let $\sigma(F) = F^*F|_{\partial_+ SM}$, then:*

- (i) $\sigma : \sqrt{\mathbb{G}_0} \rightarrow C_\alpha^\infty(\partial_+ SM, \text{Her}_n^+)$ is surjective,
- (ii) $\sigma : \{F \in \sqrt{\mathbb{G}_0} : F \text{ even}\} \rightarrow C_{\alpha,1}^\infty(\partial_+ SM, \text{Her}_n^+)$ is surjective,
- (iii) $\sigma : \{F \in \sqrt{\mathbb{G}_0} : F_0 = \text{Id}\} \rightarrow C_{\alpha,0}^\infty(\partial_+ SM, \text{Her}_n^+)$ is bijective.

Proof. — For (i) let $w \in C_\alpha^\infty(\partial_+ SM, \text{Her}_n^+)$, then $w^\sharp \in C^\infty(SM, \text{Her}_n^+)$ (it takes values in Her_n^+ , as it is constant along the geodesic flow) and by Theorem 2.3 there exists $F \in \mathbb{G}$ with $w^\sharp = F^*F$. We then automatically have $F \in \sqrt{\mathbb{G}_0}$.

For (ii) note that if $F \in \sqrt{\mathbb{G}_0}$ is even, then $F^*F|_{\partial SM}$ is α -invariant and thus $w = \sigma(F)$ satisfies

$$w(\alpha(x, -v)) = F^*F(x, -v) = F^*F(x, v) = w(x, v).$$

Conversely if $w \in C_\alpha^\infty(\partial_+ SM, \text{Her}_n^+)$ satisfies $w(\alpha(x, -v)) = w(x, v)$, then w^\sharp is even by Lemma 9.4.9 in [36]. Using Theorem 2.3, this factors as $w^\sharp = F^*F$ for an even $F \in \mathbb{G}$. As above, we must have $F \in \sqrt{\mathbb{G}_0}$ and the proof is complete.

For (iii), surjectivity is clear. Here σ is also injective, because $w^\sharp = F^*F = \tilde{F}^*\tilde{F}$ for two different $F, \tilde{F} \in \sqrt{\mathbb{G}_0}$ only if $F = u\tilde{F}$ for $u \in C^\infty(M, \text{U}(n))$, while the requirement $F_0 = \tilde{F}_0 = \text{Id}$ implies $u = \text{Id}$ and thus $F = \tilde{F}$. \square

By the preceding propositions, and if (M, g) is simple, we obtain three diagrams which are analogous to (5.2) and (5.10). The first of these is

$$(5.17) \quad \begin{array}{ccc} \sqrt{\mathbb{G}_0} & \xrightarrow{\mathbb{I}^*} & \mathcal{U}(\mathfrak{u}(n)) \\ \sigma \downarrow & \searrow^{B(\cdot|_{\partial SM})} & \downarrow \mathbb{A} \mapsto C_{\mathbb{A}} \\ C_{\alpha}^{\infty}(\partial_+ SM, \text{Her}_n^+) & \xrightarrow{P} & C^{\infty}(\partial_+ SM, \text{U}(n)), \end{array}$$

with boundary operator P yet to be defined. Here the upper triangle commutes as it arises from restricting (5.2); moreover, any choice of P making the lower triangle commute would have the same range as $\mathcal{U}(\mathfrak{u}(n)) \ni \mathbb{A} \mapsto C_{\mathbb{A}}$. However, commutativity can only be achieved up to unitary gauge, as we now explain.

Note that $C^{\infty}(M, \text{U}(n))$ acts on $\sqrt{\mathbb{G}_0}$ by left-multiplication and $C_{\text{Id}}^{\infty}(\partial M, \text{U}(n))$ (subscript Id means: homotopic to Id) acts on $C^{\infty}(\partial_+ SM, \text{GL}(n, \mathbb{C}))$ by

$$(5.18) \quad h \triangleright q = hq(h^{-1} \circ \alpha).$$

Then the diagonal arrow in (5.17) is equivariant with respect to these group actions, while σ is invariant and in fact injective up the action of $C^{\infty}(M, \text{U}(n))$, i.e., $\sigma(F) = \sigma(F')$ if and only if $F' = UF$ for some $U \in C^{\infty}(M, \text{U}(n))$ (see Theorem 2.3(ii)).

In terms of \mathcal{H}^+ from (5.4) we now define a boundary operator P by

$$(5.19) \quad Pw := B\mathcal{H}^+ A_+ w, \quad w \in C_{\alpha}^{\infty}(\partial_+ SM, \text{Her}_n^+).$$

LEMMA 5.7. — *With P as in (5.19) the lower triangle in (5.17) commutes up to gauge. That is, for any $F \in \sqrt{\mathbb{G}_0}$ there exists $h \in C_{\text{Id}}^{\infty}(\partial M, \text{U}(n))$ such that $B(F|_{\partial SM}) = h \triangleright P(\sigma(F))$.*

Proof. — Let $F \in \sqrt{\mathbb{G}_0}$ and put $w = \sigma(F)$. Then, by definition of \mathcal{H}^+ , we have $A_+ w = f^* f$, where $f = \mathcal{H}^+(A_+ w)$. On the other hand, $A_+ w = F^* F|_{\partial SM}$ is another such decomposition and thus $f = hF$ for some $h \in C_{\text{Id}}^{\infty}(\partial M, \text{U}(n))$. Hence $Pw = B(hF|_{\partial SM}) = h \triangleright B(F|_{\partial SM})$, as desired. \square

THEOREM 5.8 (Range for $\mathfrak{u}(n)$ -pairs). — *Suppose that (M, g) is simple (or more generally, that $\mathfrak{M} = 0$). Then an element $q \in C^{\infty}(\partial_+ SM, \text{U}(n))$ lies in the range of $\mathcal{U}(\mathfrak{u}(n)) \ni (A, \Phi) \rightarrow C_{A, \Phi}$ if and only if*

$$(5.20) \quad q = h \triangleright Pw, \quad \text{for some } (w, h) \in C_{\alpha}^{\infty}(\partial_+ SM, \text{Her}_n^+) \times C_{\text{Id}}^{\infty}(\partial M, \text{U}(n)).$$

Proof. — The proof is essentially a diagram chase in (5.17). First suppose that $q = C_{\mathbb{A}}$ for some $\mathbb{A} \in \mathcal{U}(\mathfrak{u}(n))$. By Proposition 5.5 (valid also if $\mathfrak{M} = 0$) there exists $F \in \sqrt{\mathbb{G}_0}$ with $\mathbb{I}^*(F) = \mathbb{A}$ and consequently, using Lemma 5.7, we have $q = B(F|_{\partial SM}) = h \triangleright P(w)$ for $w = \sigma(F)$ and some $h \in C_{\text{Id}}^{\infty}(\partial M, \text{U}(n))$. For the other direction suppose that $q = h \triangleright P(w)$ for (w, h) as in (5.20). By Proposition 5.6 we have $w = \sigma(F)$ for some $F \in \sqrt{\mathbb{G}_0}$ and by Lemma 5.7, we have $B(F|_{\partial SM}) = h_1 \triangleright Pw$ for some

$h_1 \in C^\infty_{\text{Id}}(\partial M, \text{U}(n))$. We may extend both h and h_1 to functions in $C^\infty(M, \text{U}(n))$ (denoted by the same symbol) and set $\mathbb{A} = \mathbb{I}^*(hh_1^{-1}F)$, such that

$$C_{\mathbb{A}} = (hh_1^{-1}) \triangleright B(F|_{\partial SM}) = h \triangleright Pw = q.$$

This completes the proof. □

REMARK 5.9. — In fact, on simple surfaces the range characterisation in the preceding theorem is *equivalent* to the assertion that $\mathfrak{M} = 0$ in the following sense: If the scattering data of a $\mathfrak{u}(n)$ -pair $\mathbb{A} = (A, \Phi)$ is of the form (5.20), then \mathbb{A} automatically admits holomorphic integrating factors. To see this, let (w, h) be as in (5.20), extend h to a function in $C^\infty(M, \text{U}(n))$ and factor $w^\sharp = G^*G$ for some $G \in \mathbb{G}$. Then also $F := hG$ lies in \mathbb{G} and $\mathbb{B} := -(XF)F^{-1}$ is a $\mathfrak{u}(n)$ -pair with the same scattering data as \mathbb{A} . By Theorem 1.1, the attenuation \mathbb{A} is gauge equivalent to \mathbb{B} and thus it admits holomorphic integrating factors. Granted a characterisation as in the theorem, one thus obtains HIF's for $\mathfrak{u}(n)$ -pairs and as every orbit in $\mathfrak{M} \equiv \mathcal{U}/\mathbb{G}$ contains such a pair [33, Lem. 5.2], it holds that $\mathfrak{M} = 0$.

Similarly, with α_a as in (5.16), we obtain:

THEOREM 5.10 (Range for $\mathfrak{u}(n)$ -connections). — *Suppose that (M, g) is simple. Then an element $q \in C^\infty(\partial_+ SM, \text{U}(n))$ lies in the range of $\Omega^1(M, \mathfrak{u}(n)) \ni A \mapsto C_A$ if and only if one of the following equivalent conditions is satisfied:*

- (i) q satisfies (5.20) and additionally $q \circ \alpha_a = q^{-1}$,
- (ii) q satisfies (5.20) with $w \in C^\infty_{\alpha,1}(\partial_+ SM, \text{Her}_n^+)$ (i.e., $w \circ \alpha_a = w$).

Proof. — To prove the characterisation (i), we first consider an arbitrary attenuation $\mathbb{A} \in C^\infty(SM, \mathbb{C}^{n \times n})$ with integrating factor $R \in C^\infty(SM, \mathbb{C}^{n \times n})$. Then $S(x, v) = R(x, -v)$ defines an integrating factor for $\mathbb{B}(x, v) = -\mathbb{A}(x, -v)$ and we have

$$C_{\mathbb{A}} \circ \alpha_a = (R \circ \alpha_a)(R^{-1} \circ \alpha \circ \alpha_a) = [S(S^{-1} \circ \alpha)]^{-1} = C_{\mathbb{B}}^{-1} \quad \text{on } \partial_+ SM.$$

Now, if $q = C_{\mathbb{A}}$ for \mathbb{A} equal to a connection $A \in \Omega^1(M, \mathfrak{u}(n))$, then also $\mathbb{B} = A$ and thus q has the desired symmetry. Conversely, if $q = C_{\mathbb{A}}$ for a $\mathfrak{u}(n)$ -pair $\mathbb{A} = (A, \Phi)$ and additionally $q \circ \alpha_a = q^{-1}$, then the previous display implies $C_{\mathbb{A}} = C_{\mathbb{B}}$ and by Theorem 1.1 there is a gauge $\varphi \in C^\infty(M, \text{U}(n))$ with $\varphi = \text{Id}$ on ∂M such that

$$\varphi\Phi + \Phi\varphi = 0 \quad \text{and} \quad d\varphi + [A, \varphi] = 0 \quad \text{on } M.$$

By the second equation φ solves an ODE along every curve in M and thus it is determined by its boundary values. It follows that $\varphi \equiv \text{Id}$ and hence $\Phi = 0$.

The characterisation in (ii) follows by the same arguments that lead to Theorem 5.8, replacing diagram (5.17) with the obvious analogue containing the spaces

$$\{F \in \sqrt{\mathbb{G}_0} : F \text{ even}\} \quad \text{and} \quad C^\infty_{\alpha,1}(\partial_+ SM, \text{Her}_n^+).$$

This completes the proof. □

Next, we consider the range of $\Phi \mapsto C_\Phi$ for $\mathfrak{u}(n)$ -valued matrix fields. In this case one defines a boundary operator in terms of the transform $\mathcal{H}^{+,0}$ from (5.7):

$$P_0 w := B\mathcal{H}^{+,0} A_+ w, \quad w \in C_{\alpha,0}^\infty(\partial_+ SM, \text{Her}_n^+).$$

Using Propositions 5.5 and 5.6, one obtains a similar diagram as in (5.17), this time commutative, as σ is bijective in this setting. We obtain the following result, omitting the proof:

THEOREM 5.11 (Range for $\mathfrak{u}(n)$ -matrix fields). — *Suppose that (M, g) is simple. Then an element $q \in C^\infty(\partial_+ SM, \mathfrak{U}(n))$ lies in the range of $C^\infty(M, \mathfrak{u}(n)) \ni \Phi \mapsto C_\Phi$ if and only if $q = P_0 w$ for some $w \in C_{\alpha,0}^\infty(\partial_+ SM, \text{Her}_n^+)$. \square*

REMARK 5.12. — An alternative characterisation for $\mathfrak{u}(n)$ -pairs is obtained as follows: Define $\mathbb{G}_U = \{(U, F) \in C^\infty(SM, \mathfrak{U}(n)) \times \mathbb{G} : X(UF) = 0\}$ and consider the diagram

$$(5.21) \quad \begin{array}{ccc} \mathbb{G}_U & \xrightarrow{(U, F) \mapsto \mathbb{I}^*(F)} & \mathfrak{U}(\mathfrak{u}(n)) \\ \downarrow (U, F) \mapsto UF|_{\partial_+ SM} & \searrow (U, F) \mapsto B(F|_{\partial SM}) & \downarrow \mathbb{A} \mapsto C_{\mathbb{A}} \\ C_\alpha^\infty(\partial_+ SM, \text{GL}(n, \mathbb{C})) & \xrightarrow{B\mathcal{H}A_+} & C^\infty(\partial_+ SM, \mathfrak{U}(n)), \end{array}$$

where \mathcal{H} is as in (5.3). As above, one proves commutativity up to gauge and – assuming that (M, g) is simple – one establishes surjectivity results as indicated by the double headed arrows. This shows that condition (5.20) in Theorem 5.8 can be replaced by

$$q = h \triangleright (B\mathcal{H}A_+ w) \quad \text{for some } (w, h) \in C_\alpha^\infty(\partial_+ SM, \text{GL}(n, \mathbb{C})) \times C_{\text{Id}}^\infty(\partial M, \mathfrak{U}(n)).$$

In order to isolate connections and matrix fields one can use the same ideas that lead to Theorems 5.10 and 5.11; we leave the details to the reader.

5.3. RANGE FOR $\mathfrak{gl}(n, \mathbb{C})$ -ATTENUATIONS. — To obtain range characterisations in the $\mathfrak{gl}(n, \mathbb{C})$ -case, define $\mathbb{G}_G = \{(G, F) \in \mathbb{G}^* \times \mathbb{G} : G_0 = \text{Id}\}$ and consider the diagram

$$(5.22) \quad \begin{array}{ccc} \mathbb{G}_G & \xrightarrow{(G, F) \mapsto \mathbb{I}^*(F)} & \{\mathfrak{gl}(n, \mathbb{C})\text{-pairs}\} \\ \downarrow (G, F) \mapsto GF|_{\partial_+ SM} & \searrow & \downarrow \mathbb{A} \mapsto C_{\mathbb{A}} \\ C_{\alpha,\Delta}^\infty & \xrightarrow{P^\Delta} & C^\infty(\partial_+ SM, \text{GL}(n)), \end{array}$$

where $P^\Delta = B\mathcal{H}^\Delta A_+$ with transform \mathcal{H}^Δ as in (5.8) and with diagonal arrow equal to $(G, F) \mapsto B(F|_{\partial SM})$. Here $C_{\alpha,\Delta}^\infty \subset C_\alpha^\infty(\partial_+ SM, \text{GL}(n, \mathbb{C}))$ is defined – somewhat tautologically – as range of the left vertical map and commutativity holds up to a $\text{GL}(n, \mathbb{C})$ -valued gauge. If (M, g) is simple, the top arrow is surjective and similar to Theorem 5.8, the following result holds true (the proof is omitted):

THEOREM 5.13 (Range for $\mathfrak{gl}(n, \mathbb{C})$ -pairs). — *Suppose that (M, g) is simple (or more generally, that $\mathfrak{M} = 0$). Then an element $q \in C^\infty(\partial_+ SM, \mathrm{GL}(n, \mathbb{C}))$ lies in the range of $\{\mathfrak{gl}(n, \mathbb{C})\text{-pairs}\} \ni (A, \Phi) \mapsto C_{A, \Phi}$ iff*

$$q = h \triangleright P^\Delta w \quad \text{for some } (w, h) \in C_{\alpha, \Delta}^\infty \times C_{\mathrm{Id}}^\infty(\partial M, \mathrm{GL}(n, \mathbb{C})). \quad \square$$

Similar to above one can isolate connections and matrix fields; we omit the details.

The space $C_{\alpha, \Delta}^\infty$ can also be described as follows: Any $w \in C_\alpha^\infty(\partial SM, \mathrm{GL}(n, \mathbb{C}))$ extends to a first integral $w^\sharp : SM \rightarrow \mathrm{GL}(n, \mathbb{C})$ which, by virtue of [41, Th. 8.1.2], admits a Birkhoff factorisation as

$$w^\sharp(x, \cdot) = G(x, \cdot) \Delta(x, \cdot) F(x, \cdot).$$

Here $F(x, \cdot)$ and $G(x, \cdot)^*$ are fibrewise holomorphic and

$$\Delta(x, \theta) = \mathrm{diag}(e^{ia_1(x)\theta}, \dots, e^{ia_n(x)\theta})$$

for not necessarily continuous maps $a_i : M \rightarrow \mathbb{Z}$ ($i = 1, \dots, n$). We then have $w \in C_{\alpha, \Delta}^\infty$ if and only if $\Delta \equiv \mathrm{Id}$, in which case F and G are automatically smooth on SM . To check membership of w in $C_{\alpha, \Delta}^\infty$ it is thus necessary that $\Delta|_{\partial_+ SM} \equiv \mathrm{Id}$; it is an interesting question whether this also sufficient or more generally, whether discontinuities of Δ – also called *jumping lines* – can be detected at the boundary.

5.4. NONTRIVIAL \mathfrak{M} . — Finally, we give a range characterisation if $\mathfrak{M} \neq 0$. In this case the range of $\{\mathfrak{u}(n)\text{-pairs}\} \ni (A, \Phi) \mapsto C_{A, \Phi}$ is parametrised in terms of *both* solitonic and radiative/dispersive degrees of freedom – as discussed below Theorem 1.5 – in the following sense: We construct a ‘boundary space’ $\mathfrak{B}(\partial_+ SM)$ which fits into a short exact sequence (of pointed sets)

$$(5.23) \quad 0 \longrightarrow C_\alpha^\infty(\partial_+ SM, \mathrm{GL}(n, \mathbb{C})) \longrightarrow \mathfrak{B}(\partial_+ SM) \longrightarrow \mathfrak{M} \longrightarrow 0,$$

and is the natural domain of a ‘boundary operator’

$$\mathcal{P} : \mathfrak{B}(\partial_+ SM) \longrightarrow \mathbb{U} \backslash C^\infty(\partial_+ SM, \mathbb{U}(n)).$$

Here the right hand side denotes the quotient by $\mathbb{U} := C_{\mathrm{Id}}^\infty(\partial M, \mathbb{U}(n))$ under the action defined in (5.18).

In order to define the boundary space $\mathfrak{B}(\partial_+ SM)$, denote $\mathfrak{S}_\mathbb{A}^\infty(\partial_+ SM, \mathrm{GL}(n, \mathbb{C}))$ the space of functions $w : \partial_+ SM \rightarrow \mathrm{GL}(n, \mathbb{C})$ for which the extension

$$E_\mathbb{A} w(x, v) = \begin{cases} w(x, v) & (x, v) \in \partial_+ SM, \\ C_\mathbb{A}^{-1} w(\alpha(x, v)) & (x, v) \in \partial_- SM, \end{cases}$$

defines a smooth function on ∂SM . Next, denote with $\mathcal{S}(\partial_+ SM)$ the subset of $\mathbb{U} \times C^\infty(\partial_+ SM, \mathrm{GL}(n, \mathbb{C}))$ consisting of pairs (\mathbb{A}, w) with $w \in \mathfrak{S}_\mathbb{A}^\infty(\partial_+ SM, \mathbb{C}^{n \times n})$. Then \mathbb{G} acts on $\mathcal{S}(\partial_+ SM)$ via $(\mathbb{A}, w) \triangleleft F = (\mathbb{A} \triangleleft F, (F^{-1}|_{\partial_+ SM})w)$ and we define

$$\mathfrak{B}(\partial_+ SM) := \mathcal{S}(\partial_+ SM) / \mathbb{G}.$$

The arrows in (5.23) are given by $w \mapsto [(0, w)]$ and $[(\mathbb{A}, w)] \mapsto [\mathbb{A}]$ respectively (exactness is obvious). The boundary operator \mathcal{P} is defined as concatenation

$$\begin{aligned} \mathfrak{B}(\partial_+ SM) &\longrightarrow \mathbb{G} \backslash C^\infty(\partial SM, \mathrm{GL}(n, \mathbb{C})) \xrightarrow{\mathcal{H}^*} \mathbb{U} \backslash C^\infty(\partial SM, \mathrm{U}(n)), \\ &\mathbb{U} \backslash C^\infty(\partial SM, \mathrm{U}(n)) \xrightarrow{B} \mathbb{U} \backslash C^\infty(\partial_+ SM, \mathrm{U}(n)), \end{aligned}$$

where the first arrow is $[(\mathbb{A}, w)] \mapsto [E_{\mathbb{A}} w]$ and \mathcal{H}^* is the transform defined in (5.4). Both \mathcal{H}^* and B are easily seen to descend to quotient spaces as indicated and we keep denoting them by the same symbols.

THEOREM 5.14. — *An element $q \in C^\infty(\partial_+ SM, \mathrm{U}(n))$ lies in the range of*

$$\{\mathfrak{u}(n)\text{-pairs}\} \ni (A, \Phi) \longmapsto C_{A, \Phi}$$

if and only if

$$[q] = \mathcal{P}(\mathfrak{b}) \quad \text{for some } \mathfrak{b} \in \mathfrak{B}(\partial_+ SM).$$

Proof. — First assume that $q = C_{\mathbb{A}}$ for a $\mathfrak{u}(n)$ -pair $\mathbb{A} = (A, \Phi)$. Let $U \in C^\infty(SM, \mathrm{U}(n))$ be a solution to $(X + \mathbb{A})U = 0$ on SM . Setting $w = U|_{\partial_+ SM}$, we have $\mathfrak{b} := [(\mathbb{A}, w)] \in \mathfrak{B}(\partial_+ SM)$ and $E_{\mathbb{A}} w = U|_{\partial SM}$, which means that $\mathcal{P}(\mathfrak{b}) = [B(U|_{\partial SM})] = [C_{\mathbb{A}}] = [q]$.

Conversely, suppose that $[q] = \mathcal{P}(\mathfrak{b})$ for some $\mathfrak{b} = [(\mathbb{A}, w)] \in \mathfrak{B}(\partial_+ SM)$. Then $w = R|_{\partial_+ SM}$ for a solution $R \in C^\infty(SM, \mathrm{GL}(n, \mathbb{C}))$ to $(X + \mathbb{A})R = 0$. By Theorem 2.3 this may be factored as $R = FU$ for $F \in \mathbb{G}$ and $U \in C^\infty(SM, \mathrm{U}(n))$. Then

$$\mathcal{P}(\mathfrak{b}) = [B(U|_{\partial SM})] = [C_{\mathbb{A} \triangleleft F}],$$

and by Lemma 5.2 in [33], the attenuation $\mathbb{A} \triangleleft F$ is given by a $\mathfrak{u}(n)$ -pair, as desired. \square

6. APPENDIX

6.1. A TAME SETTING. — We discuss the Fréchet structure and tameness of the spaces, Lie groups and actions used in the proofs of Theorem 1.3 and Proposition 3.5. Throughout (M, g) is a compact, oriented Riemannian surface with smooth and possibly empty boundary ∂M .

First recall that $C^\infty(SM, \mathbb{C}^n)$ has a standard Fréchet topology, which can be generated by norms $\|\cdot\|_{H^s}$ of the Sobolev-spaces $H^s(SM, \mathbb{C}^n)$ ($s \in \mathbb{R}$). We view $C^\infty(SM, \mathbb{C}^n)$ as graded Fréchet space, with grading given by

$$(6.1) \quad \|\cdot\|_{L^2} = \|\cdot\|_{H^0} \leq \|\cdot\|_{H^1} \leq \dots$$

Note that while there are several ways to define these norms, a different choice will result in a tamely equivalent grading. As in Section 2 we will tacitly apply the considerations in this section also to $\mathfrak{gl}(n, \mathbb{C})$ -valued functions.

6.1.1. Tame spaces. — The space $\bigoplus_{k \in I} \Omega_k$ ($I \subset \mathbb{Z}$) from (2.2) lies closed in the ambient C^∞ -space and thus inherits a Fréchet topology and a grading. The next lemma implies that $\bigoplus_{k \in I} \Omega_k$ is a tame direct summand and as $C^\infty(SM, \mathbb{C}^n)$ is tame [16, Cor. 1.3.7, §II] the space $\bigoplus_{k \in I} \Omega_k$ must be tame itself [16, Lem. 1.3.3, §II].

In particular, both \mathcal{U} from (1.3) and \mathcal{U}_{odd} from Proposition 3.5 are tame Fréchet spaces.

LEMMA 6.1. — For all $I \subset \mathbb{Z}$, the L^2 -orthogonal projection $P_I : C^\infty(SM, \mathbb{C}^n) \rightarrow \bigoplus_{k \in I} \Omega_k$ satisfies the tame estimate

$$\|P_I u\|_{H^s} \lesssim \|u\|_{H^s}, \quad u \in C^\infty(SM, \mathbb{C}^n), \quad s \geq 0,$$

where \lesssim means up to a constant that may depend on I and s .

Proof. — First note that P_I extends to a bounded map $L^2(SM, \mathbb{C}^n) \rightarrow L^2(SM, \mathbb{C}^n)$, simply because it is a projection. This gives a tame estimate for $s = 0$ and we may proceed by induction. To this end note the Sobolev scale on SM is generated by the operators η_\pm from (2.3) together with the vertical derivative V in the sense that

$$\|\cdot\|_{H^{s+1}} \approx \|\eta_+ \cdot\|_{H^s} + \|\eta_- \cdot\|_{H^s} + \|V \cdot\|_{H^s} + \|\cdot\|_{H^s}, \quad s \geq 0,$$

is an equivalence of norms. Let $I^\pm = I \pm 1 \subset \mathbb{Z}$, then $\eta_\pm P_I = P_{I^\pm} \eta_\pm$ and $[V, P_I] = 0$, which means that for all $u \in C^\infty(SM, \mathbb{C}^n)$

$$\begin{aligned} \|P_I u\|_{H^{s+1}} &\lesssim \|\eta_+ P_I u\|_{H^s} + \|\eta_- P_I u\|_{H^s} + \|V P_I u\|_{H^s} + \|P_I u\|_{H^s} \\ &= \|P_{I^+}(\eta_+ u)\|_{H^s} + \|P_{I^-}(\eta_- u)\|_{H^s} + \|P_I(Vu)\|_{H^s} + \|P_I u\|_{H^s} \\ &\lesssim \|\eta_+ u\|_{H^s} + \|\eta_- u\|_{H^s} + \|Vu\|_{H^s} + \|u\|_{H^s} \lesssim \|u\|_{H^{s+1}}, \end{aligned}$$

where we have used the induction hypothesis. □

6.1.2. *Tame Lie groups.* — The group $\widehat{\mathbb{G}} = C^\infty(SM, \text{GL}(n, \mathbb{C}))$ lies open in the ambient space $C^\infty(SM, \mathbb{C}^{n \times n})$ and thus is a Fréchet manifold. We claim that $\widehat{\mathbb{G}}$ is a tame Lie Group, which means that additionally the maps

$$(6.2) \quad \mathfrak{m} : \widehat{\mathbb{G}} \times \widehat{\mathbb{G}} \longrightarrow \widehat{\mathbb{G}} \quad \text{and} \quad \mathfrak{i} : \widehat{\mathbb{G}} \longrightarrow \widehat{\mathbb{G}},$$

given by multiplication and taking inverses, respectively, are *smooth tame*. Similarly, the subgroups $\mathbb{G} \subset \widehat{\mathbb{G}}$ from (1.2) and $\mathbb{G}_{\text{ev}} \subset \mathbb{G}$ from Theorem 3.5 are tame Lie groups.

To prove tameness of \mathfrak{m} and \mathfrak{i} one may invoke the high-level Theorem 2.2.6 in [16, §II], which states that so called ‘nonlinear vector bundle operators’ are tame. To this end let E and F be trivial vector bundles over SM , with fibres given by $\mathbb{C}^{n \times n}$ and $\mathbb{C}^{n \times n} \times \mathbb{C}^{n \times n}$ respectively. Next, let $U \subset E$ and $V \subset F$ be the open subsets consisting of tuples (x, v, A) and (x, v, A, B) respectively, where $(x, v) \in SM$ and $A, B \in \text{GL}(n, \mathbb{C})$. Then

$$\begin{aligned} p : V &\longrightarrow E, & (x, v, A, B) &\longmapsto (x, v, AB), \\ q : U &\longrightarrow E, & (x, v, A) &\longmapsto (x, v, A^{-1}), \end{aligned}$$

are ‘nonlinear vector bundle maps’ in Hamilton’s sense. Let $\mathcal{V} \subset C^\infty(SM, F)$ be the set of sections f with image in V and denote $Pf = p \circ f$, then the just cited theorem implies that $P : \mathcal{V} \rightarrow C^\infty(SM, E)$ is a tame map. Similarly, q gives rise to a tame map $Q : \mathcal{U} \rightarrow C^\infty(SM, E)$. Identifying \mathcal{U} with $\widehat{\mathbb{G}}$ and \mathcal{V} with $\widehat{\mathbb{G}} \times \widehat{\mathbb{G}}$, we see that P and Q correspond precisely to \mathfrak{m} and \mathfrak{i} , such that we have established tameness of multiplication and inversion on $\widehat{\mathbb{G}}$.

For \mathfrak{m} and \mathfrak{i} to be smooth tame it is required that they be smooth (which is clear) and that all derivatives are tame. However, this is a consequence of the already

obtained tameness, for all derivatives are again given in terms of multiplication and inversion.

6.1.3. *Tame actions.* — Finally, the various Lie group actions defined in the paper are smooth tame. As each of the actions is given in terms of multiplication, inversion and taking adjoints, this can be proved similar to above, by recasting the action map as nonlinear partial differential operator and applying [16, Cor.2.2.7]; we omit the details.

6.2. OKA-GRAUERT PRINCIPLE ON COMPACT DISKS. — Let $\mathbb{D} = \{z \in \mathbb{C} : |z| \leq 1\} \subset \mathbb{C}$.

LEMMA 6.2. — *Let $a \in C^\infty(\mathbb{D}, \mathbb{C}^{n \times n})$ ($n \in \mathbb{N}$). Then there exist $\mathrm{GL}(n, \mathbb{C})$ -valued solutions $f, g \in C^\infty(\mathbb{D}, \mathrm{GL}(n, \mathbb{C}))$ to the equations*

$$(6.3) \quad \partial_{\bar{z}}f + af = 0 \quad \text{and} \quad \partial_{\bar{z}}g - ga = 0 \quad \text{on } \mathbb{D}.$$

Moreover, if $a = a(p, \cdot)$ smoothly depends on a parameter p in some manifold \mathcal{P} — that is, $a \in C^\infty(\mathcal{P} \times \mathbb{D}, \mathbb{C}^{n \times n})$ — then there are corresponding solutions $f = f(p, \cdot)$ and $g = g(p, \cdot)$ in $C^\infty(\mathcal{P} \times \mathbb{D}, \mathrm{GL}(n, \mathbb{C}))$.

This classical result is discussed e.g. in [23] — we include a brief sketch of its proof and refer to the just cited monograph for further background and details.

Proof. — It suffices to solve the second equation in (6.3), the first one is then solved by $f = g^{-1}$. We can extend a to a function $a \in C^\infty(\mathcal{P} \times \mathbb{C}, \mathbb{C}^{n \times n})$ and cover \mathbb{D} by translates of the box $[0, \varepsilon]^2 \subset \mathbb{C}$. For $\varepsilon > 0$ sufficiently small, $\mathrm{GL}(n, \mathbb{C})$ -valued solutions $g_1, \dots, g_{m(\varepsilon)}$, defined in neighbourhoods of the boxes, can be constructed by means of a scaling argument and a Neumann series. By Cartan’s lemma, these local solutions can be patched together and, when restricted to \mathbb{D} , yield the desired solution g . For more details — including smooth parameter dependence in case that \mathcal{P} is an open subset of \mathbb{R}^n — see [23, Th. 1, p. 66]. The passage to \mathcal{P} being a manifold follows from standard arguments. \square

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